

**PROBABILITY, STATISTICS  
AND  
RANDOM PROCESSES**

**with Queueing Theory and  
Queueing Networks**

*Fourth Edition*

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# **PROBABILITY, STATISTICS AND RANDOM PROCESSES**

## **with Queueing Theory and Queueing Networks**

*Fourth Edition*

**T Veerarajan**

*Dean (Retd.), Department of Mathematics  
Velammal College of Engineering and Technology  
Viraganoor, Madurai, Tamil Nadu*



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# Preface to the Fourth Edition

I am extremely gratified by the overwhelming response shown to the first three editions of my book *Probability, Statistics and Random Processes* by the students and teachers throughout Tamil Nadu.

Based on the new Anna University syllabus, as per 2013 regulation, an additional chapter on “Advanced Queue Models and Queue Networks” has been added for the benefit of CSE and IT branches.

The topics ‘Test of Hypotheses’ and ‘Design of Experiments’ have been retained in the revised edition, as they are required for Mechanical Engineering students. Also, a new chapter on ‘Statistical Quality Control’ has been introduced. Solved question papers of 2011–2014 have been added. Thus, this revised edition, will cater to the requirement of all the branches for which this is a core subject.

The salient features of the book are the following:

- Written as per the latest AU syllabus
- Lucid writing style supported by step-by-step solutions
- Application-based problems for better comprehension of concepts
- Solved AU 2011 to 2014 examination question papers
- Rich exam-oriented pedagogy includes
  - ◆ Solved Examples: 379
  - ◆ Exercise Problems: 1184
  - ◆ Figures: 42
  - ◆ Tables: 35

I hope that the book will be received by both the faculty and students as enthusiastically as the previous edition of the book and my other books. Critical evaluation and suggestions for further improvement of the book will be highly appreciated and acknowledged.

**T VEERARAJAN**

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# Preface to the First Edition

This book conforms to the syllabi of the Probability and Queueing Theory paper of computer science, the Random Processes paper of Electronics and Communication and the Probability and Statistics paper of Information Technology streams of engineering at Anna University.

Most engineering students, who are used to a deterministic outlook of Physics and Engineering problems, find the theory of probability unreliable, vague and difficult. This is due to inadequate understanding of the basic concepts of probability theory and the wrong impression that the subject is an advanced branch of Mathematics.

The book is written in such a manner that beginners may develop an interest in the subject and may find it useful to pursue their studies. Basic concepts and proofs of theorems are explained in as lucid a manner as possible. Although the theory of probability is developed rigorously based on measure theory, it is developed in this book by simple set theory approach.

As engineering students find it easier to generalize specific results and examples than to specialize general results, considerable attention is devoted to working of problems. More than 300 problems, including those with applications to communication theory, are worked out in various chapters. Unless the students become personally involved in solving exercises, they cannot really develop an understanding and appreciation of the ideas and a familiarity with the pertinent techniques. Hence, in addition to a large number of short-answer questions under Part-A, over 350 problems have been given under Part-B of the Exercises in various chapters. Answers are provided at the end of every chapter.

Though chapters 7 and 8 are meant for Electrical/Electronics Engineering students, the other chapters that deal with probability theory, random variables, probability distributions and statistics will be useful to the students of other disciplines of engineering as well as those doing MCA and M.Sc courses.

I am sure that the students and the faculty will find this book very useful.

Critical evaluation and suggestions for improvement of the book will be highly appreciated and gratefully acknowledged.

I am extremely grateful to Dr K V Kuppusamy, Chairman, and Mr K Senthil Ganesh, Managing Trustee, RVS Educational Trust, Dindigul, for the support extended to me in this project.

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I have great pleasure in dedicating this book to my beloved students, past and present.

**T VEERARAJAN**

# ROADMAPS TO THE SYLLABUS

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This text is useful for Probability and Random Processes (ECE & Biomedical: MA6451) as taught in Anna University.

## Unit 1 Random Variables

Discrete and continuous random variables; Moments; Moment generating functions; Binomial, Poisson, Geometric, Uniform, Exponential, Gamma and Normal distributions

GO TO

CHAPTER 2 *Random Variables*

CHAPTER 5 *Some Special Probability Distributions*

## Unit 2 Two-Dimensional Random Variables

Joint distributions; Marginal and conditional distributions; Covariance; Correlation and Linear regression; Transformation of random variables

GO TO

CHAPTER 3 *Random Variables*

CHAPTER 4 *Statistical Averages*

## Unit 3 Random Processes

Classification; Stationary process; Markov process; Poisson process; Random telegraph process

GO TO

CHAPTER 7 *Special Random Processes*

## Unit 4 Correlation and Spectral Densities

Auto-correlation functions; Cross-correlation functions; Properties; Power spectral density; Cross spectral density; Properties

GO TO

CHAPTER 6 *Random Processes*

## Unit 5 Linear Systems with Random Inputs

Linear time invariant system; System transfer function; Linear systems with random inputs; Auto-correlation and Cross-correlation functions of input and output

**GO TO**

*CHAPTER 4 Statistical Averages*

*CHAPTER 6 Random Processes*

**This text is useful for Probability and Queueing Theory  
(CSE & IT: MA6453) as taught in Anna University.**

## Unit 1 Random Variables

Discrete and continuous random variables; Moments; Moment generating functions; Binomial, Poisson, Geometric, Uniform, Exponential, Gamma and Normal distributions

**GO TO**

*CHAPTER 2 Random Variables*

## Unit 2 Two-Dimensional Random Variables

Joint distributions; Marginal and conditional distributions; Covariance; Correlation and linear regression; Transformation of random variables

**GO TO**

*CHAPTER 2 Random Variables*

*CHAPTER 4 Statistical Averages*

## Unit 3 Random Processes

Classification; Stationary process; Markov process ; Poisson process; Discrete parameter; Markov chain; Chapman–Kolmogorov equations; Limiting distributions

**GO TO**

*CHAPTER 7 Special Random Processes*

**Unit 4 Queueing Models**

Markovian queues; Birth and death processes; Single and multiple server queueing models; Little's formula; Queues with finite waiting rooms; Queues with impatient customers: Balking and reneging

 GO TOCHAPTER 8 *Queueing Theory***Unit 5 Advanced Queueing Models**

Finite source models; M/G/1 queue; Pollaczek–Khinchin formula; M/D/1 and M/EK/1 as special cases; Series queues; Open Jackson networks

 GO TOCHAPTER 9 *Advanced Queue Models and Queue Networks*



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# Probability Theory

Probability theory had its origin in the analysis of certain games of chance that were popular in the seventeenth century. It has since found applications in many branches of Science and Engineering and this extensive application makes it an important branch of study. Probability theory, as a matter of fact, is a study of random or unpredictable experiments and is helpful in investigating the important features of these random experiments.

## Random Experiment

An experiment whose outcome or result can be predicted with certainty is called a deterministic experiment. For example, if the potential difference  $E$  between the two ends of a conductor and the resistance  $R$  are known, the current  $I$  flowing in the conductor is uniquely determined by Ohm's law,  $I = \frac{E}{R}$ .

Although all possible outcomes of an experiment may be known in advance, the outcome of a particular performance of the experiment cannot be predicted owing to a number of unknown causes. Such an experiment is called a random experiment.

Whenever a fair 6-faced cubic die is rolled, it is known that any of the 6 possible outcomes will occur, but it cannot be predicted what exactly the outcome will be, when the die is rolled at a point of time.

Although the number of telephone calls received in a board in a 5-minutes interval is a non-negative integer, we cannot predict exactly the number of calls received in the next 5-minutes. In such situations, we talk of the chance or the probability of occurrence of a particular outcome, which is taken as a quantitative measure of the likelihood of the occurrence of the outcome.

## Mathematical or Apriori Definition of Probability

Let  $S$  be the sample space (the set of all possible outcomes which are assumed equally likely) and  $A$  be an event (a subset of  $S$  consisting of possible outcomes)

associated with a random experiment. Let  $n(S)$  and  $n(A)$  be the number of elements of  $S$  and  $A$ . Then the probability of event  $A$  occurring, denoted as  $P(A)$ , is defined by

$$P(A) = \frac{n(A)}{n(S)} = \frac{\text{Number of cases favourable to } A}{\text{Exhaustive number of cases in } S}$$

For example, the probability of getting an even number in the die tossing experiment is 0.5, as  $S = \{1, 2, 3, 4, 5, 6\}$ ,  $E = \{2, 4, 6\}$ ,  $n(S) = 6$  and  $n(E) = 3$ .

## Statistical or Aposteriori Definition of Probability

Let a random experiment be repeated  $n$  times and let an event  $A$  occur  $n_A$  times out of the  $n$  trials. The ratio  $\frac{n_A}{n}$  is called the relative frequency of the event  $A$ .

As  $n$  increases,  $\frac{n_A}{n}$  shows a tendency to stabilise and to approach a constant value. This value, denoted by  $P(A)$ , is called the probability of the event  $A$ , i.e.,

$$P(A) = \lim_{n \rightarrow \infty} \frac{n_A}{n}.$$

For example, if we want to find the probability that a spare part produced by a machine is defective, we study the record of defective items produced by the machine over a considerable period of time. If, out of 10,000 items produced, 500 are defective, it is assumed that the probability of a defective item is 0.05.

**Note**

*From both the definitions, it is obvious that  $0 \leq P(A) \leq 1$ . If  $A$  is an impossible event,  $P(A) = 0$ . Conversely, if  $P(A) = 0$ , then  $A$  can occur in a very small percentage of times in the long run. On the other hand, if  $A$  is a certain event,  $P(A) = 1$ . Conversely, if  $P(A) = 1$ , then  $A$  may fail to occur in a very small percentage of times in the long run.*

## Axiomatic Definition of Probability

Let  $S$  be the sample space and  $A$  be an event associated with a random experiment. Then the probability of the event  $A$ , denoted by  $P(A)$ , is defined as a real number satisfying the following axioms.

- (i)  $0 \leq P(A) \leq 1$
- (ii)  $P(S) = 1$
- (iii) If  $A$  and  $B$  are mutually exclusive events,  $P(A \cup B) = P(A) + P(B)$
- (iv) If  $A_1, A_2, \dots, A_n, \dots$  are a set of mutually exclusive events,  $P(A_1 \cup A_2 \cup \dots \cup A_n \dots) = P(A_1) + P(A_2) + \dots + P(A_n) + \dots$

The term *mutually exclusive* used in the above definition can be explained as follows. A set of events is said to be mutually exclusive if the occurrence of any one of them excludes the occurrence of the others. Two events  $A$  and  $B$  are mutually exclusive if  $A$  occurs and  $B$  does not occur and vice versa. In other words,  $A$  and  $B$  cannot occur simultaneously, i.e.,  $P(A \cap B) = 0$ .

In the development of the probability theory, all results are derived directly or indirectly using only the axioms of probability, as can be seen from the following theorems.

### **Theorem 1**

The probability of the impossible event is zero, i.e., if  $\phi$  is the subset (event) containing no sample point,  $P(\phi) = 0$ .

#### **Proof**

The certain event  $S$  and the impossible event  $\phi$  are mutually exclusive.

$$\text{Hence } P(S \cup \phi) = P(S) + P(\phi) \text{ [Axiom (iii)].}$$

$$\text{But } S \cup \phi = S.$$

$$\therefore P(S) = P(S) + P(\phi)$$

$$\therefore P(\phi) = 0$$

### **Theorem 2**

If  $\bar{A}$  is the complementary event of  $A$ ,  $P(\bar{A}) = 1 - P(A) \leq 1$ .

#### **Proof**

$A$  and  $\bar{A}$  are mutually exclusive events, such that  $A \cup \bar{A} = S$ .

$$\begin{aligned} \therefore P(A \cup \bar{A}) &= P(S) \\ &= 1 \text{ [Axiom (ii)]} \end{aligned}$$

$$\text{i.e., } P(A) + P(\bar{A}) = 1 \text{ [Axiom (iii)]}$$

$$\therefore P(\bar{A}) = 1 - P(A)$$

Since  $P(A) \geq 0$ , it follows that  $P(\bar{A}) \leq 1$ .

### **Theorem 3**

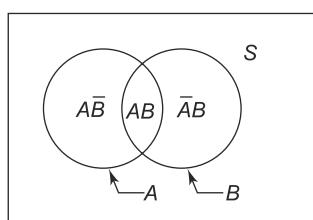
If  $A$  and  $B$  are any 2 events,  $P(A \cup B) = P(A) + P(B) - P(A \cap B) \leq P(A) + P(B)$

#### **Proof**

$A$  is the union of the mutually exclusive events  $A\bar{B}$  and  $AB$  and  $B$  is the union of the mutually exclusive events  $\bar{A}B$  and  $AB$ .

$$\therefore P(A) = P(A\bar{B}) + P(AB) \text{ [Axiom (iii)]}$$

$$\text{And } P(B) = P(\bar{A}B) + P(AB) \text{ [Axiom (iii)]}$$



**Fig. 1.1**

$$\begin{aligned}\therefore P(A) + P(B) &= [P(A \bar{B}) + P(AB) + P(\bar{A}B) + P(AB)] \\ &= P(A \cup B) + P(A \cap B)\end{aligned}$$

The result follows. Clearly,  $P(A) + P(B) - P(A \cap B) \leq P(A) + P(B)$

### Theorem 4

If  $B \subset A$ ,  $P(B) \leq P(A)$ .

### Proof

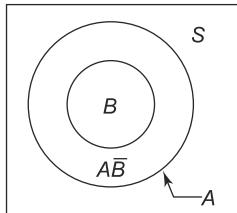


Fig. 1.2

$B$  and  $A \bar{B}$  are mutually exclusive events such that  $B \cup A \bar{B} = A$ .

$$\therefore P(B \cup A \bar{B}) = P(A)$$

i.e.  $P(B) + P(A \bar{B}) = P(A)$  [Axiom (iii)]

$$\therefore P(B) \leq P(A)$$

**Note** In probability theory developed using the classical definition of probability, Theorem 3 above is termed as Addition theorem of probability as applied to any 2 events. The theorem can be extended to any 3 events  $A$ ,  $B$  and  $C$  as follows:

$$\begin{aligned}P(A \cup B \cup C) &= P(\text{at least one of } A, B \text{ and } C \text{ occurs}) \\ &= P(A) + P(B) + P(C) - P(A \cap B) \\ &\quad - P(B \cap C) - P(C \cap A) + P(A \cap B \cap C)\end{aligned}$$

In the classical approach, probability axiom (iii) is termed as addition theorem of probability as applied to 2 mutually exclusive events, which is proved in the following way.

Let the total number of cases (outcomes) be  $n$ , of which  $n_A$  are favourable to the event  $A$  and  $n_B$  are favourable to the event  $B$ .

Therefore the numbers of cases favourable to  $A$  or  $B$ , i.e.  $A \cup B$  is  $(n_A + n_B)$ , since the events  $A$  and  $B$  are disjoint.

$$\therefore P(A \cup B) = \frac{n_A + n_B}{n} = \frac{n_A}{n} + \frac{n_B}{n} = P(A) + P(B)$$

## Conditional Probability

The conditional probability of an event  $B$ , assuming that the event  $A$  has happened, is denoted by  $P(B/A)$  and defined as

$$P(B/A) = \frac{P(A \cap B)}{P(A)}, \text{ provided } P(A) \neq 0$$

For example, when a fair die is tossed, the conditional probability of getting ‘1’, given that an odd number has been obtained, is equal to  $1/3$  as explained below:

$$S = \{1, 2, 3, 4, 5, 6\}; A = \{1, 3, 5\}; B = \{1\}$$

$$\therefore P(B/A) = \frac{n(A \cap B)}{n(A)} = \frac{1}{3}$$

$$\text{As per the definition given above, } P(B/A) = \frac{P(A \cap B)}{P(A)} = \frac{1/6}{1/2} = \frac{1}{3}.$$

Rewriting the definition of conditional probability, we get  $P(A \cap B) = P(A) \times P(B/A)$ . This is sometimes referred to as *product theorem of probability*, which is proved as follows:

Let  $n_A$ ,  $n_{AB}$  be the number of cases favourable to the events  $A$  and  $A \cap B$ , out of the total number  $n$  of cases.

$$\therefore P(A \cap B) = \frac{n_{AB}}{n} = \frac{n_A}{n} \times \frac{n_{AB}}{n_A} = P(A) \times P(B/A)$$

The product theorem can be extended to three events  $A$ ,  $B$  and  $C$  as follows:

$$P(A \cap B \cap C) = P(A) \times P(B/A) \times P(C/A \text{ and } B)$$

The following properties are easily deduced from the definition of conditional probability:

1. If  $A \subset B$ ,  $P(B/A) = 1$ , since  $A \cap B = A$
2. If  $B \subset A$ ,  $P(B/A) \geq P(B)$ , since  $A \cap B = B$ , and  $\frac{P(B)}{P(A)} \geq P(B)$ , as  $P(A) \leq P(S) = 1$
3. If  $A$  and  $B$  are mutually exclusive events,  $P(B/A) = 0$ , since  $P(A \cap B) = 0$
4. If  $P(A) > P(B)$ ,  $P(A/B) > P(B/A)$
5. If  $A_1 \subset A_2$ ,  $P(A_1/B) \leq P(A_2/B)$

## Independent Events

A set of events is said to be independent if the occurrence of any one of them does not depend on the occurrence or non-occurrence of the others.

When two events  $A$  and  $B$  are independent, it is obvious from the definition that  $P(B/A) = P(B)$ . If the events  $A$  and  $B$  are independent, the product theorem takes the form  $P(A \cap B) = P(A) \times P(B)$ . Conversely, If  $P(A \cap B) = P(A) \times P(B)$ , the events  $A$  and  $B$  are said to be independent (pairwise independent). The product theorem can be extended to any number of independent events: If  $A_1, A_2, \dots, A_n$  are  $n$  independent events.

$$P(A_1 \cap A_2 \cap \dots \cap A_n) = P(A_1) \times P(A_2) \times \dots \times P(A_n)$$

When this condition is satisfied, the events  $A_1, A_2, \dots, A_n$  are also said to be *totally independent*. A set of events  $A_1, A_2, \dots, A_n$  is said to be *mutually independent* if the events are totally independent when considered in sets of 2, 3, ...,  $n$  events.

In the case of more than two events, the term ‘independence’ is taken as ‘total independence’ unless specified otherwise.

### Theorem 1

If the events  $A$  and  $B$  are independent, the events  $\bar{A}$  and  $\bar{B}$  (and similarly  $A$  and  $\bar{B}$ ) are also independent.

#### Proof

The events  $A \cap B$  and  $\bar{A} \cap B$  are mutually exclusive such that  $(A \cap B) \cup (\bar{A} \cap B) = B$ .

$$\therefore P(A \cap B) + P(\bar{A} \cap B) = P(B) \text{ (by addition theorem)}$$

$$\begin{aligned}\therefore P(\bar{A} \cap B) &= P(B) - P(A \cap B) \\ &= P(B) - P(A)P(B) \text{ (by product theorem)} \\ &= P(B)[1 - P(A)] \\ &= P(\bar{A})P(B)\end{aligned}$$

### Theorem 2

If the events  $A$  and  $B$  are independent, then so are  $\bar{A}$  and  $\bar{B}$ .

#### Proof

$$P(\bar{A} \cap \bar{B}) = P(\overline{A \cup B}) = 1 - P(A \cup B) \quad (1)$$

$$\begin{aligned}&= 1 - [P(A) + P(B) - P(A \cap B)] \text{ (by addition theorem)} \\ &= 1 - P(A) - P(B) + P(A) \times P(B) \text{ (since } A \text{ and } B \text{ are independent)} \\ &= [1 - P(A)] - P(B)[1 - P(A)] \\ &= P(\bar{A}) \times P(\bar{B})\end{aligned} \quad (2)$$

**Note** From (1) and (2), it follows that when the events  $A$  and  $B$  are independent,  
 $P(A \cup B) = 1 - P(\bar{A}) \times P(\bar{B})$ .

### Worked Example 1(A)

#### Example 1

A fair coin is tossed 4 times. Define the sample space corresponding to this random experiment. Also give the subsets corresponding to the following events and find the respective probabilities:

- (a) More heads than tails are obtained.  
 (b) Tails occur on the even numbered tosses.

$$S = \{\text{HHHH}, \text{HHHT}, \text{HHTH}, \text{HHTT}, \text{HTHH}, \text{HTHT}, \text{HTTH}, \text{HTTT}, \\ \text{THHH}, \text{THHT}, \text{THTH}, \text{THTT}, \text{TTHH}, \text{TTHT}, \text{TTTH}, \text{TTTT}\}$$

- (a) Let  $A$  be the event in which more heads occur than tails.  
 Then  $A = \{\text{HHHH}, \text{HHHT}, \text{HHTH}, \text{HTHH}, \text{THHH}\}$   
 (b) Let  $B$  the event in which tails occur in the second and fourth tosses.  
 Then  $B = \{\text{HTHT}, \text{HTTT}, \text{TTHT}, \text{TTTT}\}$

$$P(A) = \frac{n(A)}{n(S)} = \frac{5}{16}; P(B) = \frac{n(B)}{n(S)} = \frac{1}{4}$$

### Example 2

There are 4 letters and 4 addressed envelopes. If the letters are placed in the envelopes at random, find the probability that (a) none of the letters is in the correct envelope, and (b) at least 1 letter is in the correct envelope, by explicitly writing the sample space and the event spaces.

Let the envelopes be denoted by  $A, B, C$  and  $D$  and the corresponding letters by  $a, b, c$  and  $d$ .

$$S = \{(Aa, Bb, Cc, Dd), (Aa, Bb, Cd, Dc), (Aa, Bc, Cb, Dd), \\ (Aa, Bc, Cd, Db), (Aa, Bd, Cb, Dc), (Aa, Bd, Cc, Db), \\ (Ab, Ba, Cc, Dd), (Ab, Ba, Cd, Dc), (Ab, Bc, Ca, Dd), \\ (Ab, Bc, Cd, Da), (Ab, Bd, Ca, Dc), (Ab, Bd, Cc, Da), \\ (Ac, Ba, Cb, Dd), (Ac, Ba, Cd, Db), (Ac, Bb, Ca, Dd), \\ (Ac, Bb, Cd, Da), (Ac, Bd, Ca, Db), (Ac, Bd, Cb, Da), \\ (Ad, Ba, Cb, Dc), (Ad, Ba, Cc, Db), (Ad, Bb, Ca, Dc), \\ (Ad, Bb, Cc, Da), (Ad, Bc, Ca, Db), (Ad, Bc, Cb, Da)\}$$

where ‘Aa’ means that the letter ‘a’ is placed in the envelope A.

Let  $E_1$  denote the event in which none of the letters is in the correct envelope.

Then  $E_1 = \{(Ab, Ba, Cd, Dc), (Ab, Bc, Cd, Da), (Ab, Bd, Ca, Dc), \\ (Ac, Ba, Cd, Db), (Ac, Bd, Ca, Db), (Ac, Bd, Cb, Da), \\ (Ad, Ba, Cb, Dc), (Ad, Bc, Ca, Db), (Ad, Bc, Cb, Da)\}$

Let  $E_2$  denote the event in which at least one of the letters is in the correct envelope.

We note that  $E_2$  is the complement of  $E_1$ . Therefore,  $E_2$  consists of all elements of  $S$  except those in  $E_1$ .

$$P(E_1) = \frac{n(E_1)}{n(S)} = \frac{9}{24} = \frac{3}{8} \text{ and } P(E_2) = 1 - P(E_1) = \frac{5}{8}$$

---

### Example 3

---

A lot consists of 10 good articles, 4 with minor defects and 2 with major defects. Two articles are chosen from the lot at random (without replacement). Find the probability that (a) both are good, (b) both have major defects, (c) at least 1 is good, (d) at most 1 is good, (e) exactly 1 is good, (f) neither has major defects, and (g) neither is good.

Although the articles may be drawn one after the other, we can consider that both articles are drawn simultaneously, as they are drawn without replacement.

$$(a) \quad P(\text{both are good}) = \frac{\text{No. of ways drawing 2 good articles}}{\text{Total no. of ways of drawing 2 articles}}$$

$$= \frac{10C_2}{16C_2} = \frac{3}{8}$$

$$(b) \quad P(\text{both have major defects})$$

$$= \frac{\text{No. of ways of drawing 2 articles with major defects}}{\text{Total no. of ways}}$$

$$= \frac{2C_2}{16C_2} = \frac{1}{120}$$

$$(c) \quad P(\text{at least 1 is good}) = P(\text{exactly 1 is good or both are good})$$

$$= P(\text{exactly 1 is good and 1 is bad or both are good})$$

$$= \frac{10C_1 \times 6C_1 + 10C_2}{16C_2} = \frac{7}{8}$$

$$(d) \quad P(\text{atmost 1 is good}) = P(\text{none is good or 1 is good and 1 is bad})$$

$$= \frac{10C_0 \times 6C_2 + 10C_1 \times 6C_1}{16C_2} = \frac{5}{8}$$

$$(e) \quad P(\text{exactly 1 is good}) = P(1 \text{ is good and 1 is bad})$$

$$= \frac{10C_1 \times 6C_1}{16C_2} = \frac{1}{2}$$

$$(f) \quad P(\text{neither has major defects})$$

$$= P(\text{both are non-major defective articles})$$

$$= \frac{14C_2}{16C_2} = \frac{91}{120}$$

$$(g) \quad P(\text{neither is good}) = P(\text{both are defective})$$

$$= \frac{6C_2}{16C_2} = \frac{1}{8}$$

---

### Example 4

---

From 6 positive and 8 negative numbers, 4 numbers are chosen at random (without replacement) and multiplied. What is the probability that the product is positive?

If the product is to be positive, all the 4 numbers must be positive or all the 4 must be negative or 2 of them must be positive and the other 2 must be negative.

No. of ways of choosing 4 positive numbers =  $6C_4 = 15$ .

No. of ways of choosing 4 negative numbers =  $8C_4 = 70$ .

No. of ways of choosing 2 positive and 2 negative numbers

$$= 6C_2 \times 8C_2 = 420.$$

Total no. of ways of choosing 4 numbers from all the 14 numbers

$$= 14C_4 = 1001.$$

$P(\text{the product is positive})$

$$= \frac{\text{No. of ways by which the product is positive}}{\text{Total no. of ways}}$$

$$= \frac{15 + 70 + 420}{1001} = \frac{505}{1001}$$

---

### Example 5

---

A box contains tags marked 1, 2, ...,  $n$ . Two tags are chosen at random without replacement. Find the probability that the numbers on the tags will be consecutive integers.

If the numbers on the tags are to be consecutive integers, they must be chosen as a pair from the following pairs.

$$(1, 2); (2, 3); (3, 4); \dots; (n - 1, n)$$

No. of ways of choosing any one pair from the above  $(n - 1)$  pairs =  $(n - 1) C_1 = n - 1$ .

Total no. of ways of choosing 2 tags from the  $n$  tags =  $nC_2$ .

$$\therefore \text{Required probability} = \frac{n-1}{\frac{n(n-1)}{2}} = \frac{2}{n}$$

---

### Example 6

---

If  $n$  biscuits are distributed at random among  $m$  children, what is the probability that a particular child receives  $r$  biscuits, where  $r < n$ ?

The first biscuit can be given to any one of the  $m$  children, i.e. in  $m$  ways.

Similarly, the second biscuit can be given in  $m$  ways.

Therefore, 2 biscuits can be given in  $m^2$  ways.

Extending,  $n$  biscuits can be distributed in  $m^n$  ways. The  $r$  biscuits received by the particular child can be chosen from the  $n$  biscuits in  $nC_r$  ways. If this child has got  $r$  biscuits, the remaining  $(n - r)$  biscuits can be distributed among the remaining  $(m - 1)$  children in  $(m - 1)^{n-r}$  ways.

∴ No. of ways of distributing in the required manner

$$= nC_r(m - 1)^{n-r}$$

$$\therefore \text{Required probability} = \frac{nC_r(m - 1)^{n-r}}{m^n}$$

---

### Example 7

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If  $P(A) = P(B) = P(AB)$ , show that  $P(A\bar{B} + \bar{A}B) = 0$  [ $AB \equiv A \cap B$ ].

By addition theorem,

$$P(A \cup B) = P(A) + P(B) - P(AB) \quad (1)$$

From the Venn diagram on page (1.3), it is clear that

$$A \cup B = A\bar{B} + \bar{A}B + AB$$

$$\therefore P(A \cup B) = P(A\bar{B}) + P(\bar{A}B) + P(AB) \text{ (by probability axiom)} \quad (2)$$

Using the given condition in (1),

$$P(A \cup B) = P(AB) \quad (3)$$

From (2) and (3),  $P(A\bar{B}) + P(\bar{A}B) = 0$

---

### Example 8

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If  $A$ ,  $B$  and  $C$  are any 3 events such that  $P(A) = P(B) = P(C) = 1/4$ ,  $P(A \cap B) = P(B \cap C) = 0$ ;  $P(C \cap A) = 1/8$ . Find the probability that at least 1 of the events  $A$ ,  $B$  and  $C$  occurs.

$$P(\text{at least one of } A, B \text{ and } C \text{ occurs}) = P(A \cup B \cup C)$$

$$\begin{aligned} P(A \cup B \cup C) &= P(A) + P(B) + P(C) - P(A \cap B) \\ &\quad - P(B \cap C) - P(C \cap A) + P(A \cap B \cap C) \end{aligned} \quad (1)$$

Since  $P(A \cap B) = P(B \cap C) = 0$ ,  $P(A \cap B \cap C) = 0$ . Equation (1) becomes

$$P(A \cup B \cup C) = \frac{3}{4} - 0 - 0 - \frac{1}{8} = \frac{5}{8}$$

---

### Example 9

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Solve Example 5, if the tags are chosen at random with replacement.

If the tag with '1' is chosen in the first draw, the tag with '2' must be chosen in the second draw. Probability for each =  $1/n$ .

$\therefore P(\text{'1' in the first draw and '2' in the second draw}) = 1/n^2$  (Product theorem)  
 Similarly,  $P(\text{'n' in the first draw and 'n-1' in the second draw}) = 1/n^2$ .

If the number drawn first is '2', the number drawn second may be '1' or '3'.

Probability of drawing consecutive numbers in this case

$$= \frac{1}{n} \times \frac{2}{n} = \frac{2}{n^2}$$

Similarly, when the first number drawn is '3', '4', ... '(n-1)' probability of drawing consecutive numbers will be  $2/n^2$ .

All the above possibilities are mutually exclusive.

$$\therefore \text{Required probability} = \frac{1}{n^2} + \frac{1}{n^2} + (n-2) \times \frac{2}{n^2} = \frac{2(n-1)}{n^2}$$

---

### Example 10

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A box contains 4 bad and 6 good tubes. Two are drawn out from the box at a time. One of them is tested and found to be good. What is the probability that the other one is also good?

Let  $A$  = one of the tubes drawn is good and  $B$  = the other tube is good.

$$P(A \cap B) = P(\text{both tubes drawn are good})$$

$$= \frac{6C_2}{10C_2} = \frac{1}{3}$$

Knowing that one tube is good, the conditional probability that the other tube is also good is required, i.e.,  $P(B/A)$  is required.

By definition,

$$P(B/A) = \frac{P(A \cap B)}{P(A)} = \frac{1/3}{6/10} = \frac{5}{9}$$

---

### Example 11

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Two defective tubes get mixed up with 2 good ones. *The tubes are tested, one by one, until both defectives are found.* What is the probability that the last defective tube is obtained on (a) the second test, (b) the third test, and (c) the fourth test?

Let  $D$  represent defective and  $N$  represent non-defective tube.

$$(a) P(\text{Second } D \text{ in the II test}) = P(D \text{ in the I test and } D \text{ in the II test})$$

$$= P(D_1 \cap D_2), \text{ say}$$

$$= P(D_1) \times P(D_2) \text{ (by independence)}$$

$$= \frac{2}{4} \times \frac{1}{3} = \frac{1}{6}$$

$$\begin{aligned}
 \text{(b)} \quad P(\text{second } D \text{ in the III test}) &= P(D_1 \cap N_2 \cap D_3 \text{ or } N_1 \cap D_2 \cap D_3) \\
 &= P(D_1 \cap N_2 \cap D_3) + P(N_1 \cap D_2 \cap D_3) \\
 &= \frac{2}{4} \times \frac{2}{3} \times \frac{1}{2} + \frac{2}{4} \times \frac{2}{3} \times \frac{1}{2} \\
 &= \frac{1}{3}
 \end{aligned}$$

$$\begin{aligned}
 \text{(c)} \quad P(\text{second } D \text{ in the IV test}) &= P(D_1 \cap N_2 \cap N_3 \cap D_4) + P(N_1 \cap D_2 \cap N_3 \\
 &\quad \cap D_4) + P(N_1 \cap N_2 \cap D_3 \cap D_4) \\
 &= \frac{2}{4} \times \frac{2}{3} \times \frac{1}{2} \times 1 + \frac{2}{4} \times \frac{2}{3} \times \frac{1}{2} \times 1 + \frac{2}{4} \times \frac{2}{3} \times \frac{1}{2} \times 1 \\
 &= \frac{1}{2}
 \end{aligned}$$

---

### Example 12

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In a shooting test, the probability of hitting the target is  $1/2$  for  $A$ ,  $2/3$  for  $B$  and  $3/4$  for  $C$ . If all of them fire at the target, find the probability that (a) none of them hits the target, and (b) at least one of them hits the target.

Let  $A \equiv$  Event of  $A$  hitting the target, and so on.

$$P(\bar{A}) = \frac{1}{2}, P(\bar{B}) = \frac{1}{3}, P(\bar{C}) = \frac{1}{4}$$

$$\begin{aligned}
 P(\bar{A} \cap \bar{B} \cap \bar{C}) &= P(\bar{A}) \times P(\bar{B}) \times P(\bar{C}) \text{ (by independence)} \\
 &= \frac{1}{2} \times \frac{1}{3} \times \frac{1}{4} = \frac{1}{24}
 \end{aligned}$$

$$\begin{aligned}
 P(\text{at least one hits the target}) &= 1 - P(\text{none hits the target}) \\
 &= 1 - \frac{1}{24} = \frac{23}{24}
 \end{aligned}$$

---

### Example 13

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$A$  and  $B$  alternately throw a pair of dice.  $A$  wins if he throws 6 before  $B$  throws 7 and  $B$  wins if he throws 7 before  $A$  throws 6. If  $A$  begins, show that his chance of winning is  $30/61$ .

Throwing 6 with 2 dice  $\equiv$  Getting 6 as the sum of the numbers shown on the upper faces of the 2 dice.

$$P(\text{throwing 6 with 2 dice}) = \frac{5}{36}$$

$$P(\text{throwing 7 with 2 dice}) = \frac{1}{6}$$

Let  $A \equiv$  Event of  $A$  throwing 6.

Let  $B \equiv$  Event of  $B$  throwing 7.

$A$  plays in the first, third, fifth, ..., trials.

Therefore,  $A$  will win, if he throws 6 in the first trial or third trial or in subsequent (odd) trials.

$$\begin{aligned}\therefore P(A \text{ wins}) &= P(A \text{ or } \bar{A} \bar{B} A \text{ or } \bar{A} \bar{B} \bar{A} \bar{B} A \text{ or } \dots) \\ &= P(A) + P(\bar{A} \bar{B} A) + P(\bar{A} \bar{B} \bar{A} \bar{B} A) + \dots \text{ (Addition theorem)} \\ &= \frac{5}{36} + \left( \frac{31}{36} \times \frac{5}{6} \right) \frac{5}{36} + \left( \frac{31}{36} \times \frac{5}{6} \right)^2 \times \frac{5}{36} + \dots \text{ up to } \infty \\ &= \frac{5/36}{1 - (155/216)} \text{ (since the series is an infinite geometric series)} \\ &= \frac{30}{61}\end{aligned}$$

### Example 14

Show that  $2^n - (n + 1)$  equations are needed to establish the mutual independence of  $n$  events.

$n$  events are mutually independent, if they are totally independent when considered in sets of 2, 3, ...,  $n$  events.

Sets of  $r$  events can be chosen from the  $n$  events in  $nC_r$  ways.

To establish total independence of  $r$  events, say,  $A_1, A_2, \dots, A_r$  chosen in any one of the  $nC_r$  ways, we need one equation, namely,  $P(A_1, A_2, \dots, A_r) = P(A_1) \times P(A_2) \dots \times P(A_r)$ .

Therefore, to establish total independence of all the  $nC_r$  sets, each of  $r$  events, we need  $nC_r$  equations.

Therefore, the number of equations required to establish mutual independence

$$\begin{aligned}&= \sum_{r=2}^n nC_r \\ &= (nC_0 + nC_1 + nC_2 + \dots + nC_n) - (1 + n) \\ &= (1 + 1)^n - (n + 1) \\ &= 2^n - (n + 1)\end{aligned}$$

### Example 15

Two fair dice are thrown independently. Three events  $A$ ,  $B$  and  $C$  are defined as follows:

- (a) Odd face with the first die
- (b) Odd face with the second die

- (c) Sum of the numbers in the 2 dice is odd. Are the events  $A$ ,  $B$  and  $C$  mutually independent?

$$P(A) = \frac{3}{6} = \frac{1}{2}; P(B) = \frac{3}{6} = \frac{1}{2}$$

The outcomes favourable to the event  $C$  are  $(1, 2)$ ,  $(1, 4)$ ,  $(1, 6)$ ,  $(2, 1)$ ,  $(2, 3)$ ,  $(2, 4)$ , and so on.

$$\therefore P(C) = \frac{1}{2}$$

$$P(A \cap B) = P(B \cap C) = P(A \cap C) = \frac{1}{4}$$

$$\therefore P(A \cap B) = P(A) P(B), \text{ and so on}$$

But  $P(A \cap B \cap C) = 0$ , since  $C$  cannot happen when  $A$  and  $B$  occur. Therefore  $P(A \cap B \cap C) \neq P(A) \times P(B) \times P(C)$ .

Therefore, the events are pairwise independent, but not mutually independent.

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### Example 16

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If  $A$ ,  $B$  and  $C$  are random subsets (events) in a sample space and if they are pairwise independent and  $A$  is independent of  $(B \cup C)$ , prove that  $A$ ,  $B$  and  $C$  are mutually independent.

Given  $P(AB) = P(A) \times P(B)$  (1)

$$P(BC) = P(B) \times P(C) \quad (2)$$

$$P(CA) = P(C) \times P(A) \quad (3)$$

$$P[A(B \cup C)] = P(A) \times P(B \cup C) \quad (4)$$

Consider  $P[A(B \cup C)] = P(AB \cup AC)$

$$= P(AB) + P(AC) - P(AB \cap AC) \text{ (by addition theorem)}$$

$$= P(A) \times P(B) + P(A) \times P(C) - P(ABC) \text{ [by (1) and (3)]} \quad (5)$$

Therefore, from (4) and (5), we get

$$\begin{aligned} P(ABC) &= P(A) \times P(B) + P(A) \times P(C) - P(A) \times P(B \cup C) \\ &= P(A) \times [P(B) + P(C) - P(B \cup C)] \\ &= P(A) \times P(B \cap C) \text{ (by addition theorem)} \\ &= P(A) \times P(B) \times P(C) \text{ [by (2)]} \end{aligned} \quad (6)$$

From (1), (2), (3) and (6), the required result follows.

### Exercise 1(A)

#### Part-A (Short-answer Questions)

1. What is a random experiment? Give an example.
2. Give the apriori definition of probability with an example.
3. Give the aposteriori definition of probability with an example.

4. Give the relative frequency definition of probability with an example.
5. Define the sample space and an event associated with a random experiment with an example.
6. Give the axiomatic definition of probability.
7. State the axioms of probability.
8. What do you infer from the statements  $P(A) = 0$  and  $P(A) = 1$ ?
9. Define mutually exclusive events with an example.  
(**Example:** Getting an odd number and getting an even number when a 6-faced dice is tossed are 2 mutually exclusive events.)
10. From a bag containing 3 red and 2 black balls, 2 balls are drawn at random. Find the probability that they are of the same colour.
11. When 2 cards are drawn from a well-shuffled pack of playing cards, what is the probability that they are of the same suit?
12. When  $A$  and  $B$  are 2 mutually exclusive events such that  $P(A) = 1/2$  and  $P(B) = 1/3$ , find  $P(A \cup B)$  and  $P(A \cap B)$ .
13. If  $P(A) = 0.29$ ,  $P(B) = 0.43$ , find  $P(A \cap \bar{B})$ , if  $A$  and  $B$  are mutually exclusive.
14. When  $A$  and  $B$  are 2 mutually exclusive events, are the values  $P(A) = 0.6$  and  $P(A \cap \bar{B}) = 0.5$  consistent? Why?
15. Prove that the probability of an impossible event is zero (or prove) that  $P(\emptyset) = 0$ .
16. Prove that  $P(\bar{A}) = 1 - P(A)$ , where  $\bar{A}$  is the complement of  $A$ .
17. State addition theorem as applied to any 2 events. Extend it to any 3 events.
18. If  $P(A) = 3/4$ ,  $P(B) = 5/8$ , prove that  $P(A \cap B) \geq 3/8$ .
19. A card is drawn from a well-shuffled pack of playing cards. What is the probability that it is either a spade or an ace?
20. The probability that a contractor will get a plumbing contract is  $2/3$  and the probability that he will get an electric contract is  $4/9$ . If the probability of getting at least one contract is  $4/5$ , what is the probability that he will get both?
21. If  $P(A) = 0.4$ ,  $P(B) = 0.7$  and  $P(A \cap B) = 0.3$ , find  $P(\bar{A} \cap \bar{B})$ .
22. If  $P(A) = 0.35$ ,  $P(B) = 0.75$  and  $P(A \cup B) = 0.95$ , find  $P(\bar{A} \cup \bar{B})$ .
23. Prove that  $P(A \cup B) \leq P(A) + P(B)$ . When does the equality hold good?
24. If  $B \subset A$ , prove that  $P(B) \leq P(A)$ .
25. Give the definitions of joint and conditional probabilities with examples.
26. Give the definition of conditional probability and deduce the product theorem of probability.
27. If  $A \subset B$ , prove that  $P(B/A) = 1$ .
28. If  $B \subset A$ , prove that  $P(B/A) \geq P(B)$ .
29. If  $A$  and  $B$  are mutually exclusive events, prove that  $P(B/A) = 0$ .
30. If  $P(A) > P(B)$ , prove that  $P(A/B) > P(B/A)$ .
31. If  $A \subset B$ , prove that  $P(A/C) \leq P(B/C)$ .

32. If  $P(A) = 1/3$ ,  $P(B) = 3/4$  and  $P(A \cup B) = 11/12$ , find  $P(A/B)$  and  $P(B/A)$ .
33. When are 2 events said to be independent? Give an example for 2 independent events.
34. What is the probability of getting atleast 1 head when 2 coins are tossed?
35. When 2 dice are tossed, what is the probability of getting 4 as the sum of the face numbers?
36. If the probability that  $A$  solves a problem is  $1/2$  and that for  $B$  is  $3/4$  and if they aim at solving a problem independently, what is the probability that the problem is solved?
37. If  $P(A) = 0.65$ ,  $P(B) = 0.4$  and  $P(A \cap B) = 0.24$ , can  $A$  and  $B$  independent events?
38. 15% of a firm's employees are BE degree holders, 25% are MBA degree holders and 5% have both the degrees. Find the probability of selecting a BE degree holder, if the selection is confined to MBAs.
39. In a random experiment,  $P(A) = 1/12$ ,  $P(B) = 5/12$  and  $P(B/A) = 1/15$ , find  $P(A \cup B)$ .
40. What is the difference between total independence and mutual independence?
41. Can 2 events be simultaneously independent and mutually exclusive? Explain.
42. If  $A$  and  $B$  are independent events, prove that  $\bar{A}$  and  $B$  are also independent.
43. If  $A$  and  $B$  are independent events, prove that  $A$  and  $\bar{B}$  are also independent.
44. If  $P(A) = 0.5$ ,  $P(B) = 0.3$  and  $P(A \cap B) = 0.15$ , find  $P(A/\bar{B})$ .
45. If  $A$  and  $B$  are independent events, prove that  $\bar{A}$  and  $\bar{B}$  are also independent.
46. If  $A$  and  $B$  are independent events, prove that  

$$P(A \cup B) = 1 - P(\bar{A}) \times P(\bar{B}).$$
47. A and B toss a fair coin alternately with the understanding that the one who obtains the head first wins. If A starts, what is his chance of winning?

**Part-B**

48. Write the sample space associated with the experiment of tossing 3 coins at a time and the event of getting heads from the first 2 coins. Also find the corresponding probability.
49. Items coming off a production line are marked defective ( $D$ ) or non-defective ( $N$ ). Items are observed and their condition listed. This is continued until 2 consecutive defectives are produced or 4 items have been checked, whichever occurs first. Describe a sample space for this experiment.
50. An urn contains 2 white and 4 black balls. Two balls are drawn one by one without replacement. Write the sample space corresponding to this experiment and the subsets corresponding to the following events.

- (a) The first ball drawn is white.  
 (b) Both the balls drawn are black.
- Also find the probabilities of the above events.
51. A box contains three  $10\text{-}\Omega$  resistors labelled  $R_1$ ,  $R_2$  and  $R_3$  and two  $50\text{-}\Omega$  resistors labelled  $R_4$  and  $R_5$ . Two resistors are drawn from this box without replacement. List all the outcomes of this random experiment as pairs of resistors. Also list the outcomes associated with the following events and hence find the corresponding probabilities.
- Both the resistors drawn are  $10\text{-}\Omega$  resistors.
  - One  $10\text{-}\Omega$  resistor and one  $50\text{-}\Omega$  resistor are drawn.
  - One  $10\text{-}\Omega$  resistor is drawn in the first draw and one  $50\text{-}\Omega$  resistor is drawn in the second draw.
52. A box contains 3 white balls and 2 black balls. We remove at random 2 balls in succession. What is the probability that the first removed ball is white and the second is red?
53. An urn contains 3 white balls, 4 red balls and 5 black balls. Two balls are drawn from the urn at random. Find the probability that (i) both of them are of the same colour and (ii) they are of different colours.
54. One integer is chosen at random from the numbers 1, 2, 3, ..., 100. What is the probability that the chosen number is divisible by (i) 6 or 8 and (ii) 6 or 8 or both?
55. If there are 4 persons  $A$ ,  $B$ ,  $C$  and  $D$  and if  $A$  tossed with  $B$ , then  $C$  tossed with  $D$  and then the winners tossed. This process continues till the prize is won. What are the probabilities of each of the 4 to win?
56. Ten chips numbered 1 through 10 are mixed in a bowl. Two chips are drawn from the bowl successively and without replacement. What is the probability that their sum is 10?
57. A bag contains 10 tickets numbered 1, 2, ..., 10. Three tickets are drawn at random and arranged in ascending order of magnitude. What is the probability that the middle number is 5?
58. Two fair dice are thrown independently. Four events  $A$ ,  $B$ ,  $C$  and  $D$  are defined as follows:
- $A$ : Even face with the first dice.
  - $B$ : Even face with the second dice.
  - $C$ : Sum of the points on the 2 dice is odd.
  - $D$ : Product of the points on the 2 dice exceeds 20.
- Find the probabilities of the 4 events.
59. A box contains 4 white, 5 red and 6 black balls. Four balls are drawn at random from the box. Find the probability that among the balls drawn, there is at least 1 ball of each colour.
60. Four persons are chosen at random from a group consisting of 4 men, 3 women and 2 children. Find the chance that the selected group contains at least 1 child.

61. A committee of 6 is to be formed from 5 lecturers and 3 professors. If the members of the committee are chosen at random, what is the probability that there will be a majority of lecturers in the committee?
62. Twelve balls are placed at random in 3 boxes. What is the probability that the first box will contain 3 balls?
63. If  $A$  and  $B$  are any 2 events, show that  $P(A \cap B) \leq P(A) \leq P(A \cup B) \leq P(A) + P(B)$ .
64.  $A$  and  $B$  are 2 events associated with an experiment. If  $P(A) = 0.4$  and  $P(A \cup B) = 0.7$ , find  $P(B)$  if (i)  $A$  and  $B$  are mutually exclusive (ii)  $A$  and  $B$  are independent.
65. If  $P(A + B) = 5/6$ ,  $P(AB) = 1/3$  and  $P(\bar{B}) = 1/2$ , prove that the events  $A$  and  $B$  are independent.
66. If  $A \subset B$ ,  $P(A) = 1/4$  and  $P(B) = 1/3$ , find  $P(A/B)$  and  $P(B/A)$ .
67.  $m$  objects are selected from  $n$  objects ( $m < n$ ). What is the probability that the selection contains a particular object that was present in the  $n$  given objects?
68. What is the probability that there will be 53 sundays in (i) a leap year, and (ii) a non-leap year?
69. If the probability that a communication system has high selectivity is 0.54 and the probability that it will have high fidelity is 0.81 and the probability that it will have both is 0.18, find the probability that (i) a system with high fidelity will also have high selectivity, and (ii) a system with high selectivity will also have high fidelity.
70. An electronic assembly consists of two subsystems  $A$  and  $B$ . From previous testing procedures, the following probabilities are assumed to be known:  $P(A \text{ fails}) = 0.20$ ,  $P(A \text{ and } B \text{ fail}) = 0.15$  and  $P(B \text{ fails alone}) = 0.15$ . Evaluate (i)  $P(A \text{ fails alone})$ , and (ii)  $P(A \text{ fails}/B \text{ has failed})$ .
71. A consignment of 15 tubes contains 4 defectives. The tubes are selected at random, one by one, and examined. Assuming that the tubes tested are not put back, what is the probability that the ninth one examined is the last defective?
72. A card is drawn from a 52-card deck, and without replacing it, a second card is drawn. The first and second cards are not replaced and a third card is drawn.
  - (a) If the first card is a heart, what is the probability of second card being a heart?
  - (b) If the first and second cards are hearts, what is the probability that the third card is the king of clubs?
73. A pair of dice are rolled once. Let  $A$  be the event that the first dice has a 1 on it,  $B$  the event that the second dice has a 6 on it and  $C$  the event that the sum is 7. Are  $A$ ,  $B$  and  $C$  independent?
74. A problem is given to 3 students whose chances of solving it are  $1/2$ ,  $1/3$  and  $1/4$ . What is the probability that (i) only one of them solves the problem and (ii) the problem solved.

75. A and B alternately cut a pack of cards and the pack is shuffled after each cut. If A starts and the game is continued until one cuts a diamond, what is the chance that A wins at his second cut?
76. Players X and Y roll a pair of dice alternately. The player who rolls 11 first wins. If X starts, find his chance of winning.
77. Three persons, A, B and C draw in succession from a bag containing 8 red and 4 white balls until a white ball is drawn. What is the probability that C draws the white ball?

## Theorem of Total Probability

If  $B_1, B_2, \dots, B_n$  be a set of exhaustive and mutually exclusive events, and A is another event associated with (or caused by)  $B_i$ , then

$$P(A) = \sum_{i=1}^n P(B_i)P(A|B_i)$$

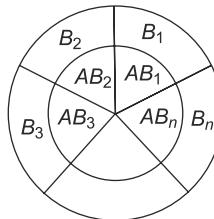


Fig. 1.3

### Proof

The inner circle represents the event A. A can occur along with (or due to)  $B_1, B_2, \dots, B_n$  that are exhaustive and mutually exclusive.

$\therefore AB_1, AB_2, \dots, AB_n$  are also mutually exclusive, such that

$$A = AB_1, AB_2, \dots, AB_n$$

$$\begin{aligned} \therefore P(A) &= P(\sum AB_i) \\ &= \sum P(AB_i) \text{ (since } AB_1, AB_2, \dots, AB_n \text{ are mutually exclusive)} \\ &\quad \text{ (by addition theorem)} \\ &= \sum_{i=1}^n P(B_i) \times P(AB_i) \end{aligned}$$

## Bayes' Theorem of Theorem of Probability of Causes

If  $B_1, B_2, \dots, B_n$  be a set of exhaustive and mutually exclusive events associated with a random experiment and A is another event associated with (or caused by)  $B_i$ , then

$$P(B_i|A) = \frac{P(B_i) \times P(A|B_i)}{\sum_{i=1}^n P(B_i) \times P(A|B_i)}, i = 1, 2, \dots, n$$

**Proof**

$$\begin{aligned} P(B_i \cup A) &= P(B_i) \times P(A|B_i) = P(A) \times P(B_i|A) \\ \therefore P(B_i|A) &= \frac{P(B_i) \times P(A|B_i)}{P(A)} \\ &= \frac{P(B_i) \times P(A|B_i)}{\sum_{i=1}^n P(B_i) \times P(A|B_i)}, i = 1, 2, 3, \dots, n \end{aligned}$$

**Worked Example 1(B)**

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**Example 1**

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A bolt is manufactured by 3 machines  $A$ ,  $B$  and  $C$ .  $A$  turns out twice as many items as  $B$ , and machines  $B$  and  $C$  produce equal number of items. 2% of bolts produced by  $A$  and  $B$  are defective and 4% of bolts produced by  $C$  are defective. All bolts are put into 1 stock pile and 1 is chosen from this pile. What is the probability that it is defective?

Let  $A$  = the event in which the item has been produced by machine  $A$ , and so on.

Let  $D$  = the event of the item being defective.

$$P(A) = \frac{1}{2}, P(B) = P(C) = \frac{1}{4}$$

$P(D/A) = P(\text{an item is defective, given that } A \text{ has produced it})$

$$= \frac{2}{100} = P(D/B)$$

$$P(D/C) = \frac{4}{100}$$

By theorem of total probability,

$$\begin{aligned} P(D) &= P(A) \times P(D/A) + P(B) \times P(D/B) + P(C) \times P(D/C) \\ &= \frac{1}{2} \times \frac{2}{100} + \frac{1}{4} \times \frac{2}{100} + \frac{1}{4} \times \frac{4}{100} \\ &= \frac{1}{40} \end{aligned}$$

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**Example 2**

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An urn contains 10 white and 3 black balls. Another urn contains 3 white and 5 black balls. Two balls are drawn at random from the first urn and placed in the second urn and then 1 ball is taken at random from the latter. What is the probability that it is a white ball?

The two balls transferred may be both white or both black or 1 white and 1 black.

Let  $B_1$  = event of drawing 2 white balls from the first urn,  $B_2$  = event of drawing 2 black balls from it and  $B_3$  = event of drawing 1 white and 1 black ball from it.

Clearly,  $B_1$ ,  $B_2$  and  $B_3$  are exhaustive and mutually exclusive events.

Let  $A$  = event of drawing a white ball from the second urn after transfer.

$$P(B_1) = \frac{10C_2}{13C_2} = \frac{15}{26}; P(B_2) = \frac{3C_2}{13C_2} = \frac{1}{26}; P(B_3) = \frac{10 \times 3}{13C_2} = \frac{10}{26},$$

$$\begin{aligned} P(A/B_1) &= P(\text{drawing a white ball}/2 \text{ white balls have been transferred}) \\ &= P(\text{drawing a white ball}/\text{urn II contains 5 white and 5 black balls}) \\ &= \frac{5}{10} \end{aligned}$$

$$\text{Similarly, } P(A/B_2) = \frac{3}{10} \text{ and } P(A/B_3) = \frac{4}{10}.$$

By theorem of total probability,

$$\begin{aligned} P(A) &= P(B_1) \times P(A/B_1) + P(B_2) \times P(A/B_2) + P(B_3) \times P(A/B_3) \\ &= \frac{15}{26} \times \frac{5}{10} + \frac{1}{26} \times \frac{3}{10} + \frac{10}{26} \times \frac{4}{10} \\ &= \frac{59}{130} \end{aligned}$$

### Example 3

In a coin tossing experiment, if the coin shows head, 1 dice is thrown and the result is recorded. But if the coin shows tail, 2 dice are thrown and their sum is recorded. What is the probability that the recorded number will be 2?

When a single dice is thrown,  $P(2) = \frac{1}{6}$ .

When 2 dice are thrown, the sum will be 2, only if each dice shows 1.

$$\therefore P(\text{getting 2 as sum with 2 dice}) = \frac{1}{6} \times \frac{1}{6} = \frac{1}{36} \text{ (since independence)}$$

By theorem of total probability,

$$\begin{aligned} P(2) &= P(H) \times P(2/H) + P(T) \times P(2/T) \\ &= \frac{1}{2} \times \frac{1}{6} + \frac{1}{2} \times \frac{1}{36} \\ &= \frac{7}{72} \end{aligned}$$

### Example 4

A bag contains 5 balls and it is not known how many of them are white. Two balls are drawn at random from the bag and they are noted to be white. What is the chance that all the balls in the bag are white?

Since 2 white balls have been drawn out, the bag must have contained 2, 3, 4 or 5 white balls.

Let  $B_1$  = Event of the bag containing 2 white balls,  $B_2$  = Events of the bag containing 3 white balls,  $B_3$  = Event of the bag containing 4 white balls and  $B_4$  = Event of the bag containing 5 white balls.

Let  $A$  = Event of drawing 2 white balls.

$$P(A/B_1) = \frac{2C_2}{5C_2} = \frac{1}{10}, P(A/B_2) = \frac{3C_2}{5C_2} = \frac{3}{10}$$

$$P(A/B_3) = \frac{4C_2}{5C_2} = \frac{3}{5}, P(A/B_4) = \frac{5C_2}{5C_2} = 1$$

Since the number of white balls in the bag is not known,  $B_i$ 's are equally likely.

$$\therefore P(B_1) = P(B_2) = P(B_3) = P(B_4) = \frac{1}{4}$$

By Bayes' theorem,

$$P(B_4/A) = \frac{P(B_4) \times P(A/B_4)}{\sum_{i=1}^4 P(B_i) \times P(A/B_i)} = \frac{\frac{1}{4} \times 1}{\frac{1}{4} \times \left( \frac{1}{10} + \frac{3}{10} + \frac{3}{5} + 1 \right)} = \frac{1}{2}$$

### Example 5

There are 3 true coins and 1 false coin with 'head' on both sides. A coin is chosen at random and tossed 4 times. If 'head' occurs all the 4 times, what is the probability that the false coin has been chosen and used?

$$P(T) = P(\text{the coin is a true coin}) = \frac{3}{4}$$

$$P(F) = P(\text{the coin is a false coin}) = \frac{1}{4}$$

Let  $A$  = Event of getting all heads in 4 tosses

$$\text{Then } P(A/T) = \frac{1}{2} \times \frac{1}{2} \times \frac{1}{2} \times \frac{1}{2} = \frac{1}{16} \text{ and } P(A/F) = 1.$$

By Bayes' theorem,

$$\begin{aligned} P(F/A) &= \frac{P(F) \times P(A/F)}{P(F) \times P(A/F) + P(T) \times P(A/T)} \\ &= \frac{\frac{1}{4} \times 1}{\frac{1}{4} \times 1 + \frac{3}{4} \times \frac{1}{16}} = \frac{16}{19} \end{aligned}$$

## Example 6

For a certain binary, communication channel, the probability that a transmitted ‘0’ is received as a ‘0’ is 0.95 and the probability that a transmitted ‘1’ is received as ‘1’ is 0.90. If the probability that a ‘0’ is transmitted is 0.4, find the probability that (i) a ‘1’ is received and (ii) a ‘1’ was transmitted given that a ‘1’ was received.

Let  $A$  = the event of transmitting ‘1’,  $\bar{A}$  = the event of transmitting ‘0’,  $B$  = the event of receiving ‘1’ and,  $\bar{B}$  = the event of receiving ‘0’.

**Given:**  $P(\bar{A}) = 0.4$ ,  $P(B/A) = 0.9$  and  $P(\bar{B}/\bar{A}) = 0.95$

$\therefore P(A) = 0.6$  and  $P(B/\bar{A}) = 0.05$

By the theorem of total probability

$$\begin{aligned} P(B) &= P(A) \times P(B/A) + P(\bar{A}) \times P(B/\bar{A}) \\ &= 0.6 \times 0.9 + 0.4 \times 0.05 \\ &= 0.56 \end{aligned}$$

By Bayes’ theorem,

$$P(A/B) = \frac{P(A) \times P(B/A)}{P(B)} = \frac{0.6 \times 0.9}{0.56} = \frac{27}{28}$$

### Exercise 1(B)

#### Part-A (Short-answer questions)

1. State the theorem of total probability.
2. Bag I contains 2 red and 1 black balls and bag II contains 3 red and 2 black balls. What is the probability that a ball drawn from one of the bags is red?
3. State Bayes’ theorem on inverse probability.
4. Bag I contains 2 white and 3 black balls and bag II contains 4 white and 1 black balls. A ball chosen at random from one of the bags is white. What is the probability that it has come from bag I?
5. Five men out of 100 and 25 women out of 1000 are colour-blind. A colour-blind person is chosen at random. What is the probability that the person is male? (Assume males and females are in equal numbers).

#### Part-B

6. There are 2 bags one of which contains 5 red and 8 black balls and the other 7 red and 10 black balls. A ball is drawn from one or the other of the 2 bags. Find the chance of drawing a red ball.
7. In a bolt factory, machines  $A$ ,  $B$  and  $C$  produce 25, 35 and 40% of the total output, respectively. Of their outputs, 5, 4 and 2%, respectively, are defective bolts. If a bolt is chosen at random from the combined output, what is the probability that it is defective? If a bolt chosen at random is

- found to be defective, what is the probability that it was produced by  $B$  or  $C$ ?
8. A box contains 2000 components of which 5% are defective. A second box contains 500 components of which 40% are defective. Two other boxes contain 1000 components, each with 10% defective components. We select at random one of the above boxes and remove from it at random a single component.
    - (a) What is the probability that the component is defective?
    - (b) Finding that the selected component is defective, what is the probability that it was drawn from box 2?
  9. There are 4 candidates for the office of the highway commissioner; the respective probabilities that they will be selected are 0.3, 0.2, 0.4 and 0.1, and the probabilities for a project's approval are 0.35, 0.85, 0.45 and 0.15, depending on which of the 4 candidates is selected. What is the probability of the project getting approved?
  10. In a binary communication system a '0' or '1' is transmitted. Because of noise in the system, a '0', can be received as a '1' with probability  $p$  and a '1' can be received as a '0' also with probability  $p$ . Assuming that the probability that a '0' is transmitted is  $p_0$ , and that a '1' is transmitted is  $q_0 (= 1 - p_0)$  find the probability that a '1' was transmitted when a '1' is received.
  11. A bag contains 7 red and 3 black marbles, and another bag contains 4 red and 5 black marbles. One marble is transferred from the first bag into the second bag and then a marble is taken out of the second bag at random. If this marble happens to be red, find the probability that a black marble was transferred.
  12. The probability that a student passes a certain exam is 0.9, given that he studied. The probability that he passes the exam without studying is 0.2. Assume that the probability that the student studies for an exam is 0.75. Given that the student passed the exam, what is the probability that he studied?
  13. Urn I has 2 white and 3 black balls, urn II has 4 white and 1 black balls and urn III has 3 white and 4 black balls. An urn is selected at random and a ball drawn at random is found to be white. Find the probability that urn I was selected.
  14. Suppose that coloured balls are distributed in 3 boxes as follows:

	<i>Box 1</i>	<i>Box 2</i>	<i>Box 3</i>
Red	2	4	3
White	3	1	4
Blue	5	3	5

A box is selected at random from which a ball is selected at random and it is observed to be red. What is the probability that box 3 was selected?

15. Three urns contain 3 white, 1 red and 1 black balls; 2 white, 3 red and 4 black balls; 1 white, 3 red and 2 black balls respectively. One urn is chosen at random and from it 2 balls are drawn at random. If they are found to be 1 red and 1 black ball, what is the probability that the first urn was chosen?
16. An urn contains 10 red and 3 black balls. Another urn contains 3 red and 5 black balls. Two balls are transferred from the first urn to the second urn, without noticing their colour. One ball is now drawn from the second urn and it is found to be red. What is the probability that 1 red and 1 black ball were transferred?
17. Box 1 contains 1000 bulbs of which 10% are defective. Box 2 contains 2000 bulbs of which 5% are defective. Two bulbs are drawn (without replacement) from a randomly selected box. (i) Find the probability that both bulbs are defective and (ii) assuming that both are defective, find the probability that they came from box 1.
18. The chance that a doctor  $A$  will diagnose a disease  $x$  correctly is 60%. The chance that a patient will die by his treatment after correct diagnosis is 40% and the chance of death by wrong diagnosis is 70%. A patient of doctor  $A$ , who had disease  $x$ , died. What is the chance that his disease was diagnosed correctly?
19. The chances of  $A$ ,  $B$  and  $C$  becoming the general manager of a certain company are in the ratio 4:2:3. The probabilities that the bonus scheme will be introduced in the company if  $A$ ,  $B$  and  $C$  become general manager are 0.3, 0.7 and 0.8 respectively. If the bonus scheme has been introduced, what is the probability that  $A$  has been appointed as general manager?

## Bernoulli's Trials

Let us consider  $n$  independent repetitions (trials) of a random experiment  $E$ . If  $A$  is an event associated with  $E$  such that  $P(A)$  remains the same for the repetitions, the trials are called *Bernoulli's trials*.

### Theorem

If the probability of occurrence of an event (probability of success) in a single trial of a Bernoulli's experiment is  $p$ , then the probability that the event occurs exactly  $r$  times out of  $n$  independent trials is equal to  $nC_r q^{n-r} p^r$ , where  $q = 1 - p$ , the probability of failure of the event.

### Proof

Getting exactly  $r$  successes means getting  $r$  successes and  $(n - r)$  failures simultaneously.

$$\therefore P(\text{getting } r \text{ successes and } n - r \text{ failures}) = p^r q^{n-r} \quad (\text{since the } n \text{ trials are independent}) \quad (\text{by product theorem}).$$

The trials, from which the successes are obtained, are not specified. There are  $nC_r$  ways of choosing  $r$  trials for successes. Once the  $r$  trials are chosen for successes, the remaining  $(n - r)$  trials should result in failures.

These  $nC_r$  ways are mutually exclusive. In each of these  $nC_r$  ways,  $P(\text{getting exactly } r \text{ success}) = p^r q^{n-r}$ .

Therefore, by the addition theorem the required probability =  $nC_r \times q^{n-r} \times p^r$ .

## De Moivre–Laplace Approximation

A result which is useful when a large number of terms of the form  $nC_r q^{n-r} p^r$  is required to be summed up, is given below without proof.

If the probability of getting exactly  $r$  successes out of  $n$  Bernoulli's trials is denoted by  $P_n(r)$ , then

$$\sum_{r=r_1}^{r_2} P_n(r) = \sum_{r=r_1}^{r_2} nC_r q^{n-r} p^r \text{ is approximately equal to } \int_{r_1 - \frac{1}{2}}^{r_2 + \frac{1}{2}} y dx,$$

where  $y = \frac{1}{\sqrt{2\pi npq}} e^{-(x-np)^2/2npq}$ , which is the density of a normal distribution with mean  $np$  and variance  $npq$ .

As the reader is familiar with normal distribution, it can be easily seen that

$$\sum_{r=r_1}^{r_2} P_n(r) = \int_{t_1}^{t_2} \phi(t) dt$$

where  $\phi(t)$  is the standard normal density and  $t_1 = \frac{r_1 - np - 1/2}{\sqrt{npq}}$  and  $t_2 = \frac{r_2 - np + 1/2}{\sqrt{npq}}$ . Now  $\int_{t_1}^{t_2} \phi(t) dt$  can be computed using the table of areas under normal curve.

## Generalisation of Bernoulli's Theorem Multinomial Distribution

If  $A_1, A_2, \dots, A_k$  are exhaustive and mutually exclusive events associated with a random experiment such that  $P(A_i \text{ occurs}) = p_i$ , where  $p_1 + p_2 + \dots + p_k = 1$ , and if the experiment is repeated  $n$  times, then the probability that  $A_1$  occurs  $r_1$  times,  $A_2$  occurs  $r_2$  times, ...  $A_k$  occurs  $r_k$  times is given by

$$P_n(r_1, r_2, \dots, r_k) = \frac{n!}{r_1! r_2! \dots r_k!} p_1^{r_1} \times p_2^{r_2} \dots p_k^{r_k}$$

where  $r_1 + r_2 + \dots + r_k = n$ .

### Proof

The  $r_1$  trials in which the event  $A_1$  occurs can be chosen from the  $n$  trials in  $nC_{r_1}$  ways. The remaining  $(n - r_1)$  trials are left over for the other events.

The  $r_2$  trials in which the event  $A_2$  occurs can be chosen from the  $(n - r_1)$  trials in  $(n - r_1)C_{r_2}$  ways.

The  $r_3$  trials in which the event  $A_3$  occurs can be chosen from the  $(n - r_1 - r_2)$  trials in  $(n - r_1 - r_2)C_{r_3}$  ways, and so on.

Therefore the number of ways in which the events  $A_1, A_2, \dots, A_k$  can happen

$$= nC_{r_1} \times (n - r_1)C_{r_2} \times (n - r_1 - r_2)C_{r_3} \times \dots \times (n - r_1 - r_2 - \dots - r_{k-1})C_{r_k}$$

$$= \frac{n!}{r_1!r_2!\dots r_k!}.$$

Consider any one of the above ways in which the events  $A_1, A_2, \dots, A_k$  occur. Since the  $n$  trials are independent,  $r_1, r_2, \dots, r_k$  trials are also independent.

$$\therefore P(A_1 \text{ occurs } r_1 \text{ times}) = p_1^{r_1}$$

$$P(A_2 \text{ occurs } r_2 \text{ times}) = p_2^{r_2}, \text{ and so on.}$$

$$\therefore P(A_1 \text{ occurs } r_1 \text{ times}, A_2 \text{ occurs } r_2 \text{ times}, \dots, A_k \text{ occurs } r_k \text{ times}) = p_1^{r_1} \times p_2^{r_2} \times \dots \times p_k^{r_k}$$

Since the ways in which the events happen are mutually exclusive, the required probability is given by

$$P_n(r_1, r_2, \dots, r_k) = \frac{n!}{r_1!r_2!\dots r_k!} p_1^{r_1} \times p_2^{r_2} \times \dots \times p_k^{r_k}$$

### Worked Example 1(C)

#### Example 1

A coin is tossed an infinite number of times. If the probability of a head in a single toss is  $p$ , show that the probability that  $k$ th head is obtained at the  $n$ th tossing, but not earlier is  $(n-1)C_{k-1}p^kq^{n-k}$ , where  $q = 1-p$ .

$k$  heads should be obtained at the  $n$ th tossing, but not earlier.

Therefore  $(k-1)$  heads must be obtained in the first  $(n-1)$  tosses and 1 head must be obtained at the  $n$ th toss.

$$\begin{aligned} \therefore \text{Required probability} &= P[k-1 \text{ heads in } (n-1) \text{ tosses}] \times P(1 \text{ head in } 1 \text{ toss}) \\ &= (n-1) C_{k-1} p^{k-1} q^{n-k} \times p \\ &= (n-1) C_{k-1} p^k q^{n-k} \end{aligned}$$

#### Example 2

Each of two persons  $A$  and  $B$  tosses 3 fair coins. What is the probability that they obtain the same number of heads?

$P(A \text{ and } B \text{ get the same no. of heads})$

$$\begin{aligned} &= P(\text{they get no head each or 1 head each or 2 heads each or 3 heads each}) \\ &= P(\text{each gets 0 head}) + P(\text{each gets 1 head}) + P(\text{each gets 2 heads}) + \\ &\quad P(\text{each gets 3 heads}) \text{ (since the events are mutually exclusive)} \\ &= P(A \text{ gets 0 head}) \times P(B \text{ gets 0 head}) + \dots \text{ (since } A \text{ and } B \text{ toss independently)} \\ &= \left[ 3C_0 \left( \frac{1}{2} \right)^3 \right]^2 + \left[ 3C_1 \left( \frac{1}{2} \right)^3 \right]^2 + \left[ 3C_2 \left( \frac{1}{2} \right)^3 \right]^2 + \left[ 3C_3 \left( \frac{1}{2} \right)^3 \right]^2 \end{aligned}$$

$$= \frac{1}{64}(1+9+9+1) = \frac{5}{16}$$

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### Example 3

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A coin with  $P(\text{head}) = p = 1 - q$  is tossed  $n$  times. Show that the probability that the number of heads obtained is even is  $0.5[1 + (q - p)^n]$ .

$P(\text{even no. of heads are obtained})$

$$\begin{aligned} &= P(0 \text{ head or } 2 \text{ heads or } 4 \text{ heads or } \dots) \\ &= nC_0 q^n p^0 + nC_2 q^{n-2} p^2 + nC_4 q^{n-4} p^4 + \dots \end{aligned} \quad (1)$$

$$\text{Consider } 1 = (q + p)^n = nC_0 q^n p^0 + nC_1 q^{n-1} p^1 + nC_2 q^{n-2} p^2 + \dots \quad (2)$$

$$(q - p)^n = nC_0 q^n p^0 - nC_1 q^{n-1} p^1 + nC_2 q^{n-2} p^2 - \dots \quad (3)$$

Adding (2) and (3), we get

$$1 + (q - p)^n = 2[nC_0 q^n p^0 + nC_2 q^{n-2} p^2 + nC_4 q^{n-4} p^4 + \dots] \quad (4)$$

Using (4) in (1), the required probability =  $0.5[1 + (q - p)^n]$ .

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### Example 4

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If at least 1 child in a family with 2 children is a boy, what is the probability that both children are boys?

$$p = \text{Probability that a child is a boy} = \frac{1}{2}$$

$$\therefore q = \frac{1}{2} \text{ and } n = 2$$

$$P(\text{at least one boy}) = p(\text{exactly 1 boy}) + p(\text{exactly 2 boys})$$

$$= 2C_1 \left(\frac{1}{2}\right)^2 + 2C_2 \left(\frac{1}{2}\right)^2 \quad (\text{by Bernoulli's theorem})$$

$$= \frac{3}{4}$$

$$P(\text{both are boys}/\text{at least one is a boy})$$

$$= \frac{P(\text{both are boys} \cap \text{at least one is a boy})}{P(\text{at least one is a boy})}$$

$$= \frac{P(\text{both are boys})}{P(\text{at least one is a boy})}$$

$$= \frac{\frac{1}{4}}{\frac{3}{4}} = \frac{1}{3}$$

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## Example 5

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Find the probability of getting at least 60 heads, when 100 fair coins are tossed.

Since the coins are fair,  $p = q = \frac{1}{2}$ .

$$n = 100$$

$$\therefore np = 50 \text{ and } \sqrt{npq} = 5$$

$$\begin{aligned}\text{Required probability } P &= \sum_{r=60}^{100} 100C_r \left(\frac{1}{2}\right)^{100} \\ &= \int_{1.9}^{10.1} \phi(t) dt \approx \int_{1.9}^{\infty} \phi(t) dt \quad (\text{by De Moivre–Laplace approximation}) \\ &= 0.5 - \int_0^{1.9} \phi(t) dt = 0.5 - 0.4719 \quad (\text{from the Normal Table}) \\ &= 0.0281\end{aligned}$$

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## Example 6

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A fair dice is rolled 5 times. Find the probability that 1 shows twice, 3 shows twice and 6 shows once.

This problem is a direct application of Bernoulli's generalisation.

Here  $n = 5$ ,  $A_1 \equiv$  (getting),  $A_3 \equiv$  (getting 3)

$A_6 \equiv$  (getting 6),  $k_1 = 2$ ,  $k_3 = 2$  and  $k_6 = 1$ ;  $p_1 = p_3 = p_6 = 1/6$

$$\therefore \text{ Required probability} = \frac{5!}{2!2!1!} \left(\frac{1}{6}\right)^2 \times \left(\frac{1}{6}\right)^2 \times \left(\frac{1}{6}\right)^1 = 0.0039$$

(It is to be noted that  $k_2 = k_4 = k_5 = 0$ , though  $p_2 = p_4 = p_5 = \frac{1}{6}$ .)

### Exercise 1(C)

#### Part-A (Short-answer Questions)

1. State Bernoulli's theorem on independent trials.
2. A fair coin is tossed 4 times. What is the probability of getting more heads than tails?
3. When 12 coins are tossed 256 times, how many times may one expect 8 heads and 4 tails?
4. If war breaks out on the average once in 25 years, find the probability that in 50 years at a stretch, there will be no war.
5. State the generalised form of Bernoulli's theorem on independent trials.
6. State De Moivre–Laplace theorem.

**Part-B**

7. A binary number (composed only of the digits ‘0’ and ‘1’) is made up of  $n$  digits. If the probability of an incorrect digit appearing is  $p$  and that errors in different digits are independent of one another, find the probability of forming an incorrect number.
8. Suppose that twice as many items are produced (per day) by machine 1 as by machine 2. However 4% of the items from machine 1 are defective while machine 2 produces only about 2% defective. Suppose that the daily output of the 2 machines is combined. A random sample of 10 items is taken from the combined output. What is the probability that this sample contains 2 defectives?
9. Binary digits are transmitted over a noisy communication channel in blocks of 16 binary digits. The probability that a received binary digit is in error because of channel noise is 0.1. If errors occur in various digit position within a block independently, find the probability that the number of errors per block is greater than or equal to 5.
10. A company is trying to market a digital transmission system (modem) that has a bit error probability of  $10^{-4}$ , and the bit errors are independent. The buyer will test the modem by sending a known message of  $10^4$  digits and checking the received message. If more than 2 errors occur, the modem will be rejected. Find the probability that the customer will buy the company’s modem.
11. A fair coin is tossed 10,000 times. Find the probability that the number of heads obtained is between 4900 and 5100, using De Moivre–Laplace approximation.
12. Over a period of 12 h, 180 calls are made at random. What is the ‘probability’ that in a 4 h interval the number of calls is between 50 and 70? Use De Moivre–Laplace approximation.

$$\left[ \text{Hint: } P(\text{a particular call occurs in the 4 h interval}) = p = \frac{4}{12} = \frac{1}{3} \right].$$

13. A random experiment can terminate in one of 3 events  $A$ ,  $B$  and  $C$  with probabilities  $1/2$ ,  $1/4$  and  $1/4$  respectively. The experiment is repeated 6 times. Find the probability that the events  $A$ ,  $B$  and  $C$  occur once, twice and thrice respectively.
14.  $A$  throws 3 fair coins and  $B$  throws 4 fair coins. Find the chance that  $A$  will throw more number of heads than would  $B$ .
15. In a large consignment of electric bulbs 10% are defective. A random sample of 20 bulbs is taken for inspection. Find the probability that (i) exactly 3 of them are defective, (ii) at most 3 of them are defective, and (iii) at least 3 of them are defective.
16. A lot contains 1% defective items. What should be the number of items in a random sample so that the probability of finding at least 1 defective in it is at least 0.95?

17. In a precision bombing attack there is a 50% chance that any one bomb will hit the target. Two direct hits are required to destroy the target completely. How many bombs must be dropped to give a 99% or more chance of completely destroying the target?

## **ANSWERS**

### **Exercise 1(A)**

10.  $P(\text{both balls are of the same colour})$   
 $= P(\text{both balls are red or both are black})$   
 $= P(\text{both are red}) + P(\text{both are black})$   
 $= \frac{3C_2}{5C_2} + \frac{2C_2}{5C_2} = \frac{2}{5}$
11. Required probability  $= P(\text{both are spades}) + P(\text{both are clubs}) + P(\text{both are hearts}) + P(\text{both are diamonds})$   
 $= 4 \times \frac{13C_2}{52C_2} = \frac{4}{17}$
12.  $P(A \cup B) = P(A) + P(B) = 5/6$ ;  $P(A \cap B) = 0$ , by definition.
13. When  $A$  and  $B$  are mutually exclusive,  $P(A \cap \bar{B}) = P(A) = 0.29$ .
14.  $P(A) = P(A \cap B) + P(A \cap \bar{B}) = 0 + P(A \cap \bar{B})$  i.e.,  $0.6 = 0.5$ , which is absurd. Hence the given values are inconsistent.
15.  $P(\text{any event}) \leq 1$ ;  $P(A \cup B) \leq 1$ ;  $P(A) + P(B) - P(A \cap B) \leq 1$   
 $\therefore P(A \cap B) \geq \frac{3}{4} + \frac{5}{8} - 1 \left( = \frac{3}{8} \right)$
16.  $P(S \cup A) = P(S) + P(A) - P(S \cap A)$   
 $= \frac{13}{52} + \frac{4}{52} - \frac{1}{52} = \frac{4}{13}$
17.  $P(A \cup B) = P(A) + P(B) - P(A \cap B)$   
 $\therefore P(A \cap B) = \frac{2}{3} + \frac{4}{9} - \frac{4}{5} = \frac{14}{45}$
18.  $P(\bar{A} \cap \bar{B}) = 1 - P(A \cup B)$   
 $= 1 - [P(A) + P(B) - P(A \cap B)]$   
 $= 1 - [0.4 + 0.7 - 0.3] = 0.2$
19.  $P(\bar{A} \cup \bar{B}) = 1 - P(A \cap B)$   
 $= 1 - [P(A) + P(B) - P(A \cup B)] = 0.85$
20. Equality holds good, when  $A$  and  $B$  are mutually exclusive events.
21. The probability for the simultaneous occurrence of two events  $A$  and  $B$  is called the joint probability of  $A$  and  $B$ . Probability of getting 2 heads, when 2 coins are tossed, is an example of joint probability.

32.  $P(A \cup B) = P(A) + P(B) - P(A \cap B)$ ;  $P(A \cap B) = \frac{1}{3} + \frac{3}{4} - \frac{11}{12} = \frac{1}{6}$

$$P(A/B) = \frac{P(A \cap B)}{P(B)} = \frac{2}{9}; P(B/A) = \frac{P(A \cap B)}{P(A)} = \frac{1}{2}$$

34. Required probability =  $1 - P(\text{both tails})$

$$= 1 - \frac{1}{2} \times \frac{1}{2} = \frac{3}{4}$$

35. Required probability =  $P(1 \text{ from I dice and } 3 \text{ from II dice or } 2 \text{ from I and } 2 \text{ from II or } 3 \text{ from I and } 1 \text{ from II})$

$$= \frac{1}{6} \times \frac{1}{6} + \frac{1}{6} \times \frac{1}{6} + \frac{1}{6} \times \frac{1}{6} = \frac{1}{12}$$

36. Required probability =  $1 - P(\text{the problem is not solved})$

$$= 1 - \frac{1}{2} \times \frac{1}{4} = \frac{7}{8}$$

37. No, since  $P(A \cap B) \neq P(A) \times P(B)$

38.  $P(BE/MBA) = \frac{P(BE \cap MBA)}{P(MBA)} = \frac{0.05}{0.25} = 0.2$

39.  $P(A \cap B) = P(A) \times P(B/A) = \frac{1}{12} \times \frac{1}{15} = \frac{1}{180}$

$P(A \cup B) = P(A) + P(B) - P(A \cap B)$

$$= \frac{1}{12} + \frac{5}{12} - \frac{1}{180} = \frac{89}{180}$$

41.  $A$  and  $B$  are independent, if  $P(A \cap B) = P(A) \times P(B)$ . They are mutually exclusive, if  $P(A \cap B) = 0$ .

They are both independent and mutually exclusive if  $P(A) \times P(B) = 0$ , i.e., if  $P(A) = 0$  or  $P(B) = 0$  or  $P(A) = 0$  and  $P(B) = 0$ . The third case is trivial. Hence  $A$  and  $B$  can be both independent and mutually exclusive, provided either of the events is an impossible event.

44.  $P(A \cap B) = P(A) \times P(B)$

$\therefore A$  and  $B$  are independent. Hence,  $A$  and  $\bar{B}$  are also independent.  
 $\therefore P(A/\bar{B}) = P(A) = 0.5$ .

47.  $P = \frac{1}{2} + \frac{1}{2} \left( \frac{1}{2} \right)^2 + \frac{1}{2} \left( \frac{1}{2} \right)^4 + \dots + \infty = \frac{\frac{1}{2}}{1 - \frac{1}{4}} = \frac{2}{3}$

48.  $\frac{1}{4}$

49.  $\{DD, NDD, DNDD, DNDN, DNND, DNNN, NDND, NDNN, NNDD, NNDN, NNND, NNNN\}$

50. (i)  $\frac{1}{3}$       (ii)  $\frac{2}{5}$

51. (i)  $\frac{3}{10}$       (ii)  $\frac{3}{5}$       (iii)  $\frac{3}{10}$

52.  $\frac{6}{25}$  (if the ball is replaced),  $\frac{3}{10}$  (if the ball is not replaced).

53. (i)  $\frac{19}{66}$       (ii)  $\frac{47}{66}$

54. (i)  $\frac{1}{5}$       (ii)  $\frac{6}{25}$

55.  $\frac{1}{4}$

56.  $\frac{4}{45}$

57.  $\frac{1}{6}$

58.  $P(A) = \frac{1}{2} = P(B) = P(C); P(D) = \frac{1}{6}$

59.  $\frac{48}{91}$

60.  $\frac{13}{18}$

61.  $\frac{9}{14}$

62.  $12 C_3 \times 2^9 / 3^{12}$

64. (i) 0.3      (ii) 0.5

66.  $\frac{3}{4}$

67.  $\frac{m}{n}$

68. (i)  $\frac{2}{7}$       (ii)  $\frac{1}{7}$

69. (i)  $\frac{2}{9}$       (ii)  $\frac{1}{3}$

70. (i) 0.05      (ii) 0.50

71.  $\frac{8}{195}$

72. (a)  $\frac{12}{51}$       (b)  $\frac{1}{50}$

73. Pairwise independent, but not totally independent

74. (i)  $\frac{11}{24}$       (ii)  $\frac{3}{4}$

75.  $\frac{9}{64}$

76.  $\frac{18}{35}$

77.  $\frac{7}{33}$

### Exercise 1(B)

2.  $P(\text{red ball}) = P(\text{red ball from bag I}) + P(\text{red ball from bag II})$

$$= \frac{1}{2} \times \frac{2}{3} + \frac{1}{2} \times \frac{3}{5} = \frac{19}{30}$$

4.  $P(B_1/W) = \frac{P(B_1) \times P(W/B_1)}{P(B_1) \times P(W/B_1) + P(B_2) \times P(W/B_2)}$

$$= \frac{\frac{1}{2} \times \frac{2}{5}}{\frac{1}{2} \times \frac{2}{5} + \frac{1}{2} \times \frac{4}{5}} = \frac{1}{3}$$

5.  $P(M) = P(F) = \frac{1}{2}; P(B/M) = \frac{1}{20}; P(B/F) = \frac{1}{40}$

$$\text{By Bayes' theorem, } P(M/B) = \frac{\frac{1}{2} \times \frac{1}{20}}{\frac{1}{2} \times \frac{1}{20} + \frac{1}{2} \times \frac{1}{40}} = \frac{2}{3}$$

6.  $\frac{88}{221}$

7. (i)  $\frac{69}{2000}$       (ii)  $\frac{44}{69}$

8. (i)  $\frac{13}{80}$       (ii)  $\frac{8}{13}$

9. 0.47

10. 
$$\frac{(1-p)q_0}{(1-p)q_0 + pp_0}$$

11.  $\frac{12}{47}$

12.  $\frac{27}{29}$

13.  $\frac{14}{57}$

14.  $\frac{5}{19}$

15.  $\frac{3}{25}$

16.  $\frac{20}{59}$

17. (i) 0.0062    (ii) 0.8005

18.  $\frac{6}{13}$

19.  $\frac{6}{25}$

### Exercise 1(C)

2. Required probability =  $P(\text{getting exactly 3 or 4 heads})$

$$= 4C_3 \left(\frac{1}{2}\right)^4 + 4C_4 \left(\frac{1}{2}\right)^4 = \frac{5}{16}$$

3. Required number =  $256 \times 12 C_8 \left(\frac{1}{2}\right)^{12} \cong 31$

4.  $p = \frac{1}{25}, q = \frac{24}{25}, n = 50; P = 50 C_0 \left(\frac{1}{25}\right)^0 \left(\frac{24}{25}\right)^{50} = \left(\frac{24}{25}\right)^{50}$

7.  $1 - (1 - p)^n$

8. 0.0381

9. 0.017

10. 0.9197

11. 0.9545

12. 0.9876

13. 0.0293

14.  $\frac{29}{128}$

15. (i) 0.1898    (ii) 0.8655    (iii) 0.3243

16. 299

17. 11



# *Chapter* 2

## Random Variables

The outcomes of random experiments may be numerical or non-numerical in nature. For example, the number of telephone calls received in a board in 1 h is numerical in nature, while the result of a coin tossing experiment in which 2 coins are tossed at a time is non-numerical in nature. As it is often useful to describe the outcome of a random experiment by a number, we will assign a number to each non-numerical outcome of the experiment. For example, in the 2 coins tossing experiment we could assign the value 0 to the outcome of getting 2 tails, 1 to the outcome of getting 1 head and 1 tail and 2 to the outcome of getting 2 heads. Thus, in any experimental situation we can assign a real number  $x$  to every element  $s$  of the sample space  $S$ . That is, the function  $X(s) = x$  that maps the elements of the sample space into real numbers is called the random variable associated with the concerned experiment. A formal definition may be given as follows.

**Definition:** A random variable (abbreviatively RV) is a function that assigns a real number  $X(s)$  to every element  $s \in S$ , where  $S$  is the sample space corresponding to a random experiment  $E$ .

**Note** Although we are expected to perform the random experiment  $E$ , we observe the outcome  $s \in S$  and then evaluate  $X(s)$  [i.e., assign a real number  $x$  to  $X(s)$ ], the number  $x = X(s)$  itself can be thought of as the outcome of the experiment and  $R_x$  as the sample space of the experiment. In this sense, we will hereafter talk about a random variable  $X$  taking the value  $x$  and  $P(X = x)$ . Actually,  $P(X = x) = P\{s: X(s) = x\}$ .

Hereafter,  $R_x$  will be referred to as *range space*.

Similarly,  $\{X \leq x\}$  represents the subset  $\{s: X(s) \leq x\}$  and, hence, an event associated with the experiment.

### Discrete Random Variable

If  $X$  is random variable (RV) which can take a finite number or countably infinite number of values,  $X$  is called a discrete RV. When the RV is discrete, the possible

values of  $X$  may be assumed as  $x_1, x_2, \dots, x_n, \dots$ . In the finite case, the list of values terminates and in the countably infinite case, the list goes upto infinity.

For example, the number shown when a dice is thrown and the number of alpha particles emitted by a radioactive source are discrete RVs.

## Probability Function

If  $X$  is a discrete RV which can take the values  $x_1, x_2, \dots, x_n, \dots$  such that  $P(X = x_i) = p_i$  then  $p_i$  is called the *probability function* or *probability mass function* or *point probability function*, provided  $p_i$  ( $i = 1, 2, 3, \dots$ ) satisfy the following conditions:

(i)  $p_i \geq 0$ , for all  $i$ , and

(ii)  $\sum_i p_i = 1$

The collection of pairs  $\{x_i, p_i\}$ ,  $i = 1, 2, 3, \dots$ , is called *the probability distribution of the RV X*, which is sometimes displayed in the form of a table as given below:

$X = x_i$	$P(X = x_i)$
$x_1$	$p_1$
$x_2$	$p_2$
$\vdots$	$\vdots$
$x_r$	$p_r$
$\vdots$	$\vdots$

## Continuous Random Variable

If  $X$  is an RV which can take all values (i.e. *infinite number* of values) in an interval, then  $X$  is called a *continuous* RV.

For example, the length of time during which a vacuum tube installed in a circuit functions is a continuous RV.

## Probability Density Function

If  $X$  is a continuous RV such that

$$P\left\{x - \frac{1}{2}dx \leq X \leq x + \frac{1}{2}dx\right\} = f(x)dx$$

then  $f(x)$  is called the *probability density function* (shortly denoted as pdf) of  $X$ , provided  $f(x)$  satisfies following conditions:

(i)  $f(x) \geq 0$ , for all  $x \in R_x$ , and

(ii)  $\int_{R_X} f(x)dx = 1$

Moreover,  $P(a \leq X \leq b)$  or  $P(a < X < b)$  is defined as

$$P(a \leq X \leq b) = \int_a^b f(x) dx.$$

The curve  $y = f(x)$  is called *the probability curve* of the RV  $X$ .

**Note**

*When  $X$  is a continuous RV*

$$P(X = a) = P(a \leq X \leq a) = \int_a^a f(x) dx = 0$$

This means that it is almost impossible that a continuous RV assumes a specific value. Hence,  $P(a \leq X \leq b) = P(a \leq X < b) = P(a < X \leq b) = P(a < X < b)$ .

## Cumulative Distribution Function (cdf)

If  $X$  is an RV, discrete or continuous then  $P(X \leq x)$  is called the *cumulative distribution function* of  $X$  or *distribution function* of  $X$  and denoted as  $F(x)$ .

If  $X$  is discrete,

$$F(x) = \sum_j P_j \quad \begin{matrix} j \\ x_j \leq x \end{matrix}$$

If  $X$  is continuous,

$$F(x) = P(-\infty < X \leq x) = \int_{-\infty}^x f(x) dx$$

### Properties of the cdf $F(x)$

1.  $F(x)$  is a non-decreasing function of  $x$ , i.e., if  $x_1 < x_2$ , then  $F(x_1) \leq F(x_2)$ .
2.  $F(-\infty) = 0$  and  $F(\infty) = 1$ .
3. If  $X$  is a discrete RV taking values  $x_1, x_2, \dots$ , where  $x_1 < x_2 < x_3 < \dots < x_{i-1} < x_i < \dots$ , then  $P(X = x_i) = F(x_i) - F(x_{i-1})$ .
4. If  $X$  is a continuous R V, then  $\frac{d}{dx} F(x) = f(x)$ , at all points where  $F(x)$  is differentiable.

**Note**

*Although we may talk of probability distribution of a continuous RV, it cannot be represented by a table as in the case of a discrete RV. The probability distribution of a continuous RV is said to be known, if either its pdf or cdf is given.*

## Special Distributions

The probability mass functions of some discrete RVs and the probability density functions of some continuous RVs, which are of frequent applications, are as follows:

## Discrete Distributions

1. If the discrete RV  $X$  can take the values  $0, 1, 2, \dots, n$ , such that  $P(X = i) = nC_i p^i q^{n-i}$ ,  $i = 0, 1, \dots, n$ , where  $p + q = 1$ , then  $X$  is said to follow a binomial distribution with parameters  $n$  and  $p$ , which is denoted a  $B(n, p)$ .
2. If the discrete RV  $X$  can take the value  $0, 1, 2, \dots$ , such that  $P(X = i) = \frac{e^{-\lambda} \lambda^i}{i!}$ ,  $i = 0, 1, 2, \dots$ , then  $X$  is said to follow a Poisson distribution with parameter  $\lambda$ .
3. If the discrete RV  $X$  can take the values  $0, 1, 2, \dots$ , such that  $P(X = i) = (n + i - 1) C_i p^n q^i$ ,  $i = 0, 1, 2, \dots$ , where  $p + q = 1$ , then  $X$  is said to follow a *Pascal (or negative binomial) distribution* with parameter  $n$ .
4. A Pascal distribution with parameter 1 [i.e.,  $P(X = i) = pq^i$ ,  $i = 0, 1, 2, \dots$  and  $p + q = 1$ ] is called a *geometric distribution*.

## Continuous Distributions

5. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{1}{b-a}$  (a constant),  $a \leq x \leq b$ , then  $X$  follows a *uniform distribution* (or *rectangular distribution*).
6. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2}$ ,  $-\infty < x < \infty$ , then  $X$  is said to follow a *normal distribution* (or *Gaussian distribution*) with parameters  $\mu$  and  $\sigma$ , which will be hereafter denoted as  $N(\mu, \sigma)$ .
7. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{1}{(n)} e^{-x} x^{n-1}$ ,  $0 < x < \infty$  and  $n > 0$ , then  $X$  follows a *gamma distribution* with parameter  $n$ . Gamma distribution is a particular case of *Erlang distribution*, the pdf of which is  $f(x) = \frac{c^n}{(n)} x^{n-1} e^{-cx}$ ,  $0 < x < \infty$ ,  $n > 0$ ,  $c > 0$ .
8. An Erlang distribution with  $n = 1$  [i.e.,  $f(x) = ce^{-cx}$ ,  $0 < x < \infty$ ,  $c > 0$ ] is called an *exponential* (or *negative exponential*) *distribution* with parameter  $c$ .
9. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{x}{\alpha^2} e^{-x^2/2\alpha^2}$ ,  $0 < x < \infty$ , then  $X$  follows a *Rayleigh distribution* with parameter  $\alpha$ .
10. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{\sqrt{2}}{\alpha^3 \sqrt{\pi}} x^2 e^{-x^2/2\alpha^2}$ ,  $0 < x < \infty$ , then  $X$  follows a *Maxwell distribution* with parameter  $\alpha$ .

11. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{1}{2\lambda} e^{|x-\mu|/\lambda}$ ,  $-\infty < x < \infty$ ,  $\lambda > 0$ , then  $X$  follows a *Laplace* (or *double exponential*) distribution with parameters  $\lambda$  and  $\mu$ .
12. If the pdf of a continuous RV  $X$  is  $f(x) = \frac{\alpha}{\pi} \times \frac{1}{x^2 + \alpha^2}$ ,  $\alpha > 0$ ,  $-\infty < x < \infty$ , then  $X$  follows a *Cauchy distribution* with parameter  $\alpha$ .

### Worked Example 2(A)

#### **Example 1**

From a lot containing 25 items, 5 of which are defective, 4 items are chosen at random. If  $X$  is the number of defectives found, obtain the probability distribution of  $X$ , when the items are chosen (i) without replacement, and (ii) with replacement.

Since only 4 items are chosen,  $X$  can take the values 0, 1, 2, 3 and 4.  
The lot contains 20 non-defective and 5 defective items.

**Case (i):** When the items are chosen without replacement, we can assume that all the 4 items are chosen simultaneously.

$$\begin{aligned}\therefore P(X = r) &= P(\text{choosing exactly } r \text{ defective items}) \\ &= P(\text{choosing } r \text{ defective and } (4 - r) \text{ good items}) \\ &= \frac{5C_r \times 20C_{4-r}}{25C_4} \quad (r = 0, 1, \dots, 4)\end{aligned}$$

**Case (ii):** When the items are chosen with replacement, we note that the probability of an item being defective remains the same in each draw.

$$\text{i.e., } p = \frac{5}{25} = \frac{1}{5}, q = \frac{4}{5} \text{ and } n = 4$$

The problem is one of performing 4 Bernoulli's trials and finding the probability of exactly  $r$  successes.

$$\therefore P(X = r) = 4C_r \left(\frac{1}{5}\right)^r \left(\frac{4}{5}\right)^{4-r} \quad (r = 0, 1, \dots, 4)$$

#### **Example 2**

A shipment of 6 television sets contains 2 defective sets. A hotel makes a random purchase of 3 of the sets. If  $X$  is the number of defective sets purchased by the hotel, find the probability distribution of  $X$ .

All the 3 sets are purchased simultaneously. Since there are only 2 defective sets in the lot,  $X$  can take the values 0, 1 and 2.

$$\begin{aligned}\therefore P(X = r) &= P(\text{choosing exactly } r \text{ defective sets}) \\ &= P[\text{choosing } r \text{ defective and } (3 - r) \text{ good sets}] \\ &= \frac{2C_r \times 4C_{3-r}}{6C_3} (r = 0, 1, 2)\end{aligned}$$

The required probability distribution is represented in the form of the following table:

$X = r$	$p_r$
0	1/5
1	3/5
2	1/5
Total	1

### Example 3

A random variable  $X$  may assume 4 values with probabilities  $(1 + 3x)/4$ ,  $(1 - x)/4$ ,  $(1 + 2x)/4$  and  $(1 - 4x)/4$ . Find the condition on  $x$  so that these values represent the probability function of  $X$ .

$$\begin{aligned}P(X = x_1) &= p_1 = (1 + 3x)/4; p_2 = (1 - x)/4; \\ p_3 &= (1 + 2x)/4; p_4 = (1 - 4x)/4\end{aligned}$$

If the given probabilities represent a probability function, each  $p_i \geq 0$  and  $\sum_i p_i = 1$ .

In this problem,  $p_1 + p_2 + p_3 + p_4 = 1$ , for any  $x$ .

But  $p_1 \geq 0$ , if  $x \geq -1/3$ ;  $p_2 \geq 0$ , if  $x \leq 1$ ;  $p_3 \geq 0$ , if  $x \geq -1/2$  and  $p_4 \geq 0$ , If  $x \leq 1/4$ .

Therefore, the values of  $x$  for which a probability function is defined lie in the range  $-1/3 \leq x \leq 1/4$ .

### Example 4

If the random variable  $X$  takes the values 1, 2, 3 and 4 such that  $2P(X = 1) = 3P(X = 2) = P(X = 3) = 5P(X = 4)$ , find the probability distribution and cumulative distribution function of  $X$ .

Let  $P(X = 3) = 30K$ . Since  $2P(X = 1) = 30K$ ,  $P(X = 1) = 15K$ .

Similarly,  $P(X = 2) = 10K$  and  $P(X = 4) = 6K$ .

Since  $\sum p_i = 1$ ,  $15K + 10K + 30K + 6K = 1$ .

$$\therefore K = \frac{1}{61}$$

The probability distribution of  $X$  is given in the following table:

$X = i$	1	2	3	4
$P_i$	$\frac{15}{61}$	$\frac{10}{61}$	$\frac{30}{61}$	$\frac{6}{61}$

The cdf  $F(x)$  is defined as  $F(x) = P(X \leq x)$ . Accordingly, the cdf for the above distribution is found out as follows:

When  $x < 1$ ,  $F(x) = 0$

When  $1 \leq x < 2$ ,  $F(x) = P(X = 1) = \frac{15}{61}$

When  $2 \leq x < 3$ ,  $F(x) = P(X = 1) + P(X = 2) = \frac{25}{61}$

When  $3 \leq x < 4$ ,  $F(x) = P(X = 1) + P(X = 2) + P(X = 3) = \frac{55}{61}$

When  $x \geq 4$ ,  $F(x) = P(x = 1) + P(x = 2) + P(x = 3) + P(x = 4) = 1$ .

### Example 5

A random variable  $X$  has the following probability distribution.

$x:$	-2	-1	0	1	2	3
$p(x):$	0.1	$K$	0.2	$2K$	0.3	$3K$

- (a) Find  $K$ . (b) Evaluate  $P(X < 2)$  and  $P(-2 < X < 2)$ . (c) Find the cdf of  $X$ .  
(d) Evaluate the mean of  $X$ .

(a) Since  $\Sigma P(x) = 1$ ,  $6K + 0.6 = 1$

$$\therefore K = \frac{1}{15}$$

$\therefore$  the probability distribution becomes

$x:$	-2	-1	0	1	2	3
$p(x):$	1/10	1/15	1/5	2/15	3/10	1/5

$$\begin{aligned} \text{(b)} \quad P(X < 2) &= P(X = -2, -1, 0 \text{ or } 1) \\ &= P(X = -2) + P(X = -1) + P(X = 0) + P(X = 1) \end{aligned}$$

[since the events  $(X = -2)$ ,  $(X = -1)$ , etc., are mutually exclusive]

$$= \frac{1}{10} + \frac{1}{15} + \frac{1}{5} + \frac{2}{15} = \frac{1}{2}$$

$$\begin{aligned} P(-2 < X < 2) &= P(X = -1, 0 \text{ or } 1) \\ &= P(X = -1) + P(X = 0) + P(X = 1) \\ &= \frac{1}{15} + \frac{1}{5} + \frac{2}{15} = \frac{2}{5} \end{aligned}$$

(c)  $F(x) = 0$ . When  $x < -2$

$$= \frac{1}{10}, \text{ when } -2 \leq x < -1$$

$$= \frac{1}{6}, \text{ when } -1 \leq x < 0$$

$$\begin{aligned}&= \frac{11}{30}, \text{ when } 0 \leq x < 1 \\&= \frac{1}{2}, \text{ when } 1 \leq x < 2 \\&= \frac{4}{5}, \text{ when } 2 \leq x < 3 \\&= 1, \text{ when } 3 \leq x\end{aligned}$$

(d) The mean of  $X$  is defined as  $E(X) = \sum xp(x)$

(refer to Chapter 4)

$$\begin{aligned}\therefore \text{Mean of } X &= \left( -2 \times \frac{1}{10} \right) + \left( -1 \times \frac{1}{15} \right) + \left( 0 \times \frac{1}{5} \right) \\&\quad + \left( 1 \times \frac{2}{15} \right) + \left( 2 \times \frac{3}{10} \right) + \left( 3 \times \frac{1}{5} \right) \\&= -\frac{1}{5} - \frac{1}{15} + \frac{2}{15} + \frac{3}{5} + \frac{3}{5} = \frac{16}{15}\end{aligned}$$

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### Example 6

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The probability function of an infinite discrete distribution is given by  $P(X = j) = 1/2^j$  ( $j = 1, 2, \dots, \infty$ ). Verify that the total probability is 1 and find the mean and variance of the distribution. Find also  $P(X \text{ is even})$ ,  $P(X \geq 5)$  and  $P(X \text{ is divisible by } 3)$ .

Let  $P(X = j) = p_j$

$$\sum_{j=1}^{\infty} p_j = \frac{1}{2} + \frac{1}{2^2} + \frac{1}{2^3} + \dots \infty, \text{ that is a geometric series.}$$

$$\begin{aligned}&= \frac{\frac{1}{2}}{1 - \frac{1}{2}} = 1\end{aligned}$$

The mean of  $X$  is defined as  $E(X) = \sum_{j=1}^{\infty} jp_j$  (refer to Chapter 4).

$$\begin{aligned}\therefore E(x) &= a + 2a^2 + 3a^3 + \dots \infty, \text{ where } a = \frac{1}{2} \\&= a(1 + 2a + 3a^2 + \dots \infty)\end{aligned}$$

$$= a(1-a)^{-2} = \frac{1}{\left(\frac{1}{2}\right)^2} = 2$$

The variance of  $X$  is defined as  $V(X) = E(X^2) - [E(X)]^2$ ,

where  $E(X^2) = \sum_{j=1}^{\infty} j^2 p^j$  (refer to Chapter 4).

$$\begin{aligned} E(X^2) &= \sum_{j=1}^{\infty} j^2 a^j, \text{ where } a = \frac{1}{2} \\ &= \sum_{j=1}^{\infty} [j(j+1) - j] a^j = \sum_{j=1}^{\infty} j(j+1)a^j - \sum_{j=1}^{\infty} ja^j \\ &= a(1.2 + 2.3a + 3.4a^2 + \dots \infty) - a(1 + 2a + 3a^2 + \dots \infty) \\ &= a \times 2(1-a)^3 - a \times (1-a)^{-2} \\ &= \frac{2a}{(1-a)^3} - \frac{a}{(1-a)^2} = 8 - 2 = 6 \\ \therefore \quad V(X) &= E(X^2) - [E(X)]^2 = 6 - 4 = 2 \end{aligned}$$

$P(X \text{ is even}) = P(X = 2 \text{ or } X = 4 \text{ or } X = 6 \text{ or etc.})$

$P(X = 2) + P(X = 4) + \dots + \infty$

(since the events are mutually exclusive)

$$\begin{aligned} &= \left(\frac{1}{2}\right)^2 + \left(\frac{1}{2}\right)^4 + \left(\frac{1}{2}\right)^6 + \dots + \infty \\ &= \frac{\frac{1}{4}}{1 - \frac{1}{4}} = \frac{1}{3} \end{aligned}$$

$P(X \geq 5) = P(X = 5 \text{ or } X = 6 \text{ or } X = 7, \text{ etc.})$

$$= P(X = 5) + P(X = 6) + \dots + \infty$$

$$\begin{aligned} &= \frac{\frac{1}{2^5}}{1 - \frac{1}{2}} = \frac{1}{16} \end{aligned}$$

$$\begin{aligned}
 P(X \text{ is divisible by } 3) &= P(X = 3 \text{ or } X = 6 \text{ or } X = 9, \text{ etc.}) \\
 &= P(X = 3) + P(X = 6) + \dots + \infty \\
 &= \left(\frac{1}{2}\right)^3 + \left(\frac{1}{2}\right)^6 + \left(\frac{1}{2}\right)^9 + \dots + \infty \\
 &= \frac{\frac{1}{8}}{1 - \frac{1}{8}} = \frac{1}{7}
 \end{aligned}$$

---

**Example 7**


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A random variable  $X$  has the following probability distribution.

$x:$	0	1	2	3	4	5	6	7
$p(x):$	0	$K$	$2K$	$2K$	$3K$	$K^2$	$2K^2$	$7K^2 + K$

Find (a) the value of  $K$ , (b)  $P(1.5 < X < 4.5/X > 2)$ , and (c) the smallest value of  $\lambda$  for which  $P(X \leq \lambda) > 1/2$ .

$$(a) \quad \sum p(x) = 1$$

$$\therefore 10K^2 + 9K = 1$$

$$\text{i.e.,} \quad (10K - 1)(K + 1) = 0$$

$$\therefore K = \frac{1}{10} \text{ or } -1.$$

The value  $K = -1$  makes some values of  $p(x)$  negative, which is meaningless.

$$\therefore K = \frac{1}{10}$$

The actual distribution is given below:

$x:$	0	1	2	3	4	5	6	7
$P(x):$	0	$\frac{1}{10}$	$\frac{2}{10}$	$\frac{2}{10}$	$\frac{3}{10}$	$\frac{1}{100}$	$\frac{2}{100}$	$\frac{17}{100}$

$$(b) \quad P(1.5 < X < 4.5/X > 2) = P(A/B), \text{ say}$$

$$= \frac{P(A \cap B)}{P(B)}$$

$$= \frac{P[(1.5 < X < 4.5) \cap (X > 2)]}{P(X > 2)}$$

$$= \frac{P(X = 3) + P(X = 4)}{\sum_{r=3}^7 (X = r)} = \frac{\frac{5}{10}}{\frac{7}{10}} = \frac{5}{7}$$

(c) By trials,  $P(X \leq 0) = 0$ ;  $P(X \leq 1) = \frac{1}{10}$ ;  $P(X \leq 2) = \frac{3}{10}$

$$P(X \leq 3) = \frac{5}{10}; P(X \leq 4) = \frac{8}{10}$$

Therefore, the smallest value of  $\lambda$  satisfying the condition  $P(X \leq \lambda) > 1/2$  is 4.

---

### Example 8

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If  $p(x) = \begin{cases} xe^{-x^2/2} & x \geq 0 \\ 0 & x < 0 \end{cases}$

- (a) show that  $p(x)$  is a pdf (of a continuous RV  $X$ ),
- (b) find its distribution function  $P(x)$ .
- (c) If  $p(x)$  is to be a pdf,  $p(x) \geq 0$  and

$$\int_{R_X} p(x) dx = 1$$

Obviously,  $p(x) = xe^{-x^2/2} \geq 0$ , when  $x \geq 0$

$$\begin{aligned} \text{Now } \int_0^\infty p(x) dx &= \int_0^\infty xe^{-x^2/2} dx = \int_0^\infty e^{-t} dt \text{ (putting } t = x^2/2) \\ &= 1 \end{aligned}$$

$\therefore p(x)$  is a legitimate pdf of a RV  $X$ .

(b)  $F(x) = P(X \leq x) = \int_0^x f(x) dx$   
 $\therefore F(x) = 0$ , when  $x < 0$

and  $F(x) = \int_0^x xe^{-x^2/2} dx = 1 - e^{-x^2/2}$ , when  $x \geq 0$ .

---

### Example 9

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If the density function of a continuous RV  $X$  is given by

$$\begin{aligned} f(x) &= ax, & 0 \leq x \leq 1 \\ &= a, & 1 \leq x \leq 2 \\ &= 3a - ax, & 2 \leq x \leq 3 \\ &= 0, & \text{elsewhere} \end{aligned}$$

- (a) find the value of  $a$ , and
- (b) find the cdf of  $X$ .
- (c) If  $x_1, x_2$  and  $x_3$  are 3 independent observations of  $X$ , what is the probability that exactly one of these 3 is greater than 1.5?

(a) Since  $f(x)$  is a pdf,  $\int_{R_X} f(x) dx = 1$

i.e.,  $\int_0^3 f(x) dx = 1$

i.e.,  $\int_0^1 ax dx + \int_1^2 adx + \int_2^3 (3a - ax) dx = 1$

i.e.,  $2a = 1$

$\therefore a = \frac{1}{2}$

(b)  $F(x) = P(X \leq x) = 0$ , when  $x < 0$

$$F(x) = \int_0^x \frac{x}{2} dx = \frac{x^4}{4}, \text{ when } 0 \leq x \leq 1$$

$$= \int_0^1 \frac{x}{2} dx + \int_1^x \frac{1}{2} dx = \frac{x}{2} - \frac{1}{4} \text{ when } 1 \leq x \leq 2$$

$$= \int_0^2 \frac{x}{2} dx + \int_2^x \frac{1}{2} dx + \int_2^3 \left( \frac{3}{2} - \frac{x}{2} \right) dx = \frac{3}{2}x - \frac{x^2}{4} - \frac{5}{4}, \text{ when } 2 \leq x \leq 3$$

$$= 1, \text{ when } x > 3$$

(c)  $p(X > 1.5) = \int_{1.5}^3 f(x) dx$

$$= \int_{1.5}^2 \frac{1}{2} dx + \int_2^3 \left( \frac{3}{2} - \frac{x}{2} \right) dx = \frac{1}{2}$$

Choosing an  $X$  and observing its value can be considered as a trial and ( $X > 1.5$ ) can be considered a success.

$\therefore p = 1/2, q = 1/2$

As we choose 3 independent observations of  $X, n = 3$ .

By Bernoulli's theorem,

$P(\text{exactly one value} > 1.5)$

$$= p(\text{1 success}) = 3C_1 \times (p)^1 \times (q)^2 = \frac{3}{8}$$

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### Example 10

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A continuous RV  $X$  that can assume any value between  $x = 2$  and  $x = 5$  has a density function given by  $f(x) = k(1 + x)$ . Find  $P(X < 4)$ .

By the property of pdf,

$$\int_{R_x} f(x) dx = 1, \text{ } X \text{ takes values between 2 and 5.}$$

$$\therefore \int_2^5 k(1+x) dx = 1$$

$$\text{i.e., } \frac{27}{2}k = 1$$

$$\therefore k = \frac{2}{27}$$

$$\text{Now } p(X < 4) = p(2 < X < 4) = \int_2^4 k(1+x) dx = \frac{16}{27}$$

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### Example 11

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A continuous RV  $X$  has a pdf  $f(x) = kx^2e^{-x}$ ;  $x \geq 0$ . Find  $k$ , mean and variance.  
By the property of pdf,

$$\int_0^\infty kx^2 e^{-x} dx = 1$$

$$\text{i.e., } 2k = 1$$

$$\therefore k = \frac{1}{2}$$

Mean of  $X$  is defined as

$$E(X) = \int_{R_x} xf(x) dx$$

(refer to Chapter 4)

Variance of  $X$  is defined as

$$V(X) = E(X^2) - \{E(X)\}^2,$$

$$\text{where } E(X)^2 = \int_{R_x} x^2 f(x) dx \text{ (refer to Chapter 4)}$$

$$\begin{aligned} \therefore E(X) &= \frac{1}{2} \int_0^\infty x^3 e^{-x} dx \\ &= \frac{1}{2} [x^3(-e^{-x}) - 3x^2(e^{-x}) + 6x(-e^{-x}) - 6(e^{-x})]_0^\infty \\ &= 3 \end{aligned}$$

$$\begin{aligned}
 E(X)^2 &= \frac{1}{2} \int_0^\infty x^4 e^{-x} dx \\
 &= \frac{1}{2} [x^4(-e^{-x}) - 4x^3(e^{-x}) + 12x^2(-e^{-x}) - 24x(e^{-x}) + 24(-e^{-x})]_0^\infty \\
 &= 12 \\
 \therefore V(X) &= E(X^2) - \{E(X)\}^2 = 3
 \end{aligned}$$

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### Example 12

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The probability that a person will die in the time interval  $(t_1, t_2)$  is given by

$$p(t_1 \leq t \leq t_2) = \int_{t_1}^{t_2} a(t) dt$$

The function  $a(t)$  is determined from long records and can be assumed to be

$$a(t) = \begin{cases} 3 \times 10^{-9} t^2 (100-t)^2 & 0 \leq t \leq 100 \\ 0 & \text{elsewhere} \end{cases}$$

Determine (a) the probability that a person will die between the ages 60 and 70, and (b) the probability that he will die between those ages, assuming he lived up to 60.

$$\begin{aligned}
 (a) \quad P(60 < t < 70) &= \int_{60}^{70} a(t) dt \\
 &= 3 \times 10^{-9} \int_{60}^{70} t^2 (100-t)^2 dt \\
 &= 0.1544
 \end{aligned}$$

$$\begin{aligned}
 (b) \quad P(60 < t < 70 / t \geq 60) &= P(60 < t < 70 / 60 \leq t \leq 100) \\
 &= \frac{P(60 < t < 70)}{P(60 < t < 100)} \\
 &= \int_{60}^{70} a(t) dt / \int_{60}^{100} a(t) dt \\
 &= \frac{0.15436}{0.31744} = 0.4863
 \end{aligned}$$

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### Example 13

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A continuous RV has a pdf  $f(x) = 3x^2$ ,  $0 \leq x \leq 1$ . Find  $a$  and  $b$  such that

- (a)  $P(X \leq a) = P(X > a)$ , and
- (b)  $P(X > b) = 0.05$ .

$$(a) \quad P(X \leq a) = P(X > a)$$

$$\therefore \int_0^a 3x^2 dx = \int_a^1 3x^2 dx$$

$$\text{i.e.,} \quad a^3 = 1 - a^3$$

$$\text{i.e.,} \quad a^3 = \frac{1}{2}$$

$$\therefore \quad a = 0.7937$$

$$(b) \quad P(X > b) = 0.05$$

$$\int_b^1 3x^2 dx = 0.05$$

$$\text{i.e.,} \quad b^3 = 95$$

$$\therefore \quad b = 0.9830$$

### Example 14

The distribution function of a RV  $X$  is given by  $F(x) = 1 - (1+x)e^{-x}$ ,  $x \geq 0$ . Find the density function, mean and variance of  $X$ .

By the property of  $F(x)$ , the pdf  $f(x)$  is given by  $f(x) = F'(x)$  at points of continuity of  $F(x)$ .

The given cdf is continuous for  $x \geq 0$ .

$$\therefore \quad f(x) = (1+x)e^{-x} - e^{-x} = x e^{-x}, x \geq 0$$

$$E(X) = \int_0^\infty x^2 e^{-x} dx = 2$$

$$E(X^2) = \int_0^\infty x^3 e^{-x} dx = 6$$

$$V(X) = E(X^2) - [E(X)]^2 = 2$$

### Example 15

The cdf of a continuous RV  $X$  is given by

$$F(x) = 0, x < 0$$

$$= x^2, 0 \leq x < \frac{1}{2}$$

$$= 1 - \frac{3}{25}(3-x)^2, \frac{1}{2} \leq x < 3$$

$$= 1, x \geq 3$$

Find the pdf of  $X$  and evaluate  $P(|X| \leq 1)$  and  $P\left(\frac{1}{3} \leq X < 4\right)$  using both the pdf and cdf.

The points  $x = 0, 1/2$  and  $3$  are points of continuity.

$$\begin{aligned}\therefore f(x) &= 0, x < 0 \\ &= 2x, \quad 0 \leq x < \frac{1}{2} \\ &= \frac{6}{25}(3-x), \quad \frac{1}{2} \leq x < 3 \\ &= 0, x \geq 3\end{aligned}$$

Although the points  $x = 1/2, 3$  are points of discontinuity for  $f(x)$ , we may assume that  $f\left(\frac{1}{2}\right) = \frac{3}{5}$  and  $f(3) = 0$ .

$$\begin{aligned}P(|X| \leq 1) &= P(-1 \leq x \leq 1) \\ &= \int_{-1}^1 f(x) dx = \int_0^{1/2} 2x dx + \int_{1/2}^{1/2} \frac{6}{25}(3-x) dx \text{ (using property of pdf)} \\ &= \frac{13}{25}\end{aligned}$$

If we use property of cdf,

$$P(|X| \leq 1) = P(-1 \leq x \leq 1) = F(1) - F(-1) = \frac{13}{25}$$

If we use the property of pdf,

$$P(1/3 \leq X < 4) = \int_{1/3}^{1/2} 2x dx + \int_{1/2}^3 \frac{6}{25}(3-x) dx = \frac{8}{9}$$

If we use the property of cdf,

$$\begin{aligned}P(1/3 \leq X < 4) &= F(4) - F\left(\frac{1}{3}\right) \\ &= 1 - \frac{1}{9} = \frac{8}{9}\end{aligned}$$

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### Example 16

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If the RV  $k$  is uniformly distributed over  $(0, 5)$  what is the probability that the roots of the equation  $4x^2 + 4kx + (k+2) = 0$  are real?

The RV  $k$  is  $U(0, 5)$ .

$$\therefore \text{pdf of } k = \frac{1}{5}, 0 < k < 5$$

$P(\text{Roots of } 4x^2 + 4kx + k + 2 = 0 \text{ are real})$

$$\begin{aligned}&= P(\text{Discriminant of the equation} \geq 0) \\ &= P(k^2 - k - 2 \geq 0) = P[(k-2)(k+1) \geq 0] \\ &= P[k \geq -1 \text{ and } k \geq 2 \text{ or } (k \leq 2 \text{ and } k \leq -1)]\end{aligned}$$

$$\begin{aligned}
 &= P(k \geq 2 \text{ or } k \leq -1) = P[(k \geq 2) \text{ [since } k \text{ takes values in } (0, 5)] \\
 &= \int_2^5 f(k) dk = \frac{1}{5}(5-2) \\
 &= \frac{3}{5}
 \end{aligned}$$

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**Example 17**

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A point  $P$  is taken at random on a line  $AB$  of length  $2a$ , all positions of the point being equally likely. Find the probability that the product  $(AP \times PB) > \frac{a^2}{2}$ . Let  $AP = X$ .

$$\therefore PB = (2a - X)$$

Since all positions of the point  $P$  are equally likely,  $X (= AP)$  is uniformly distributed over  $(0, 2a)$ .

$$\therefore \text{pdf of } X = \frac{1}{2a}, 0 < x < 2a$$

$$\begin{aligned}
 P\left[AP \times PB > \frac{a^2}{2}\right] &= P\left[X(2a - X) > \frac{a^2}{2}\right] \\
 &= P(2X^2 - 4aX + a^2 < 0) \\
 &= P\left[\left\{X - \left(1 - \frac{1}{\sqrt{2}}\right)a\right\} \left\{X - \left(1 + \frac{1}{\sqrt{2}}\right)a\right\} < 0\right]
 \end{aligned}$$

[since the factors of  $(2x^2 - 4ax + a^2)$  are  $x - \left(1 - \frac{1}{\sqrt{2}}\right)a$  and  $x - \left(1 + \frac{1}{\sqrt{2}}\right)a$ ]

$$= P\left[\left(1 - \frac{1}{\sqrt{2}}\right)a < X < \left(1 + \frac{1}{\sqrt{2}}\right)a\right]$$

$$\begin{aligned}
 &= \int_{\left(1 - \frac{1}{\sqrt{2}}\right)a}^{\left(1 + \frac{1}{\sqrt{2}}\right)a} f(x) dx = \frac{1}{2a} \int_{\left(1 - \frac{1}{\sqrt{2}}\right)a}^{\left(1 + \frac{1}{\sqrt{2}}\right)a} [x] dx = \frac{\sqrt{2}a}{2a} = \frac{1}{\sqrt{2}}
 \end{aligned}$$

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**Example 18**

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If the continuous RV,  $X$  represents the time of failure of a system that has been put into operation at  $t = 0$ , find the conditional density function of  $X$ , given that the system has survived up to time  $t$ . Deduce the same when  $X$  follows an exponential distribution with parameter  $\lambda$ .

The conditional distribution function of  $X$ , subject to the given condition, is given by

$$\begin{aligned} f(x/X > t) &= \frac{P[X \leq x \text{ and } X > t]}{P(X > t)} \quad [\text{since unconditional } F(x) = P(X \leq x)] \\ &= \frac{P[t < X \leq x]}{P[t < X < \infty]} \\ &= \frac{F(x) - F(t)}{1 - F(t)} \quad \text{for } x > t \\ &= 0 \quad \text{for } x < t \end{aligned}$$

Therefore, the conditional density function  $f(x/X > t)$  is given by

$$\begin{aligned} f(x/X > t) &= \frac{d}{dx} F(x/X > t) \\ &= \frac{f(x)}{1 - F(t)}, \quad x > t \end{aligned}$$

For the exponential distribution with parameter  $\lambda$ ,

$$\begin{aligned} f(x) &= \lambda e^{-\lambda x}, \quad x > 0, \text{ and } F(x) = \int_0^x \lambda e^{-\lambda x} dx = 1 - e^{-\lambda x} \\ \therefore f(x/X > t) &= \frac{\lambda e^{-\lambda X}}{e^{-\lambda t}} = \lambda e^{-\lambda(x-t)} = f(x-t) \end{aligned}$$

### Example 19

If  $f(t)$  is the unconditional density of the time of failure of a system and  $h(t)$  is the hazard rate (or conditional failure rate or conditional density of  $T$ , given  $T > t$ ) find  $f(t)$  in terms of  $h(t)$ . Deduce that  $T$  follows a Rayleigh distribution, when  $h(t) = t$ .

The conditional density of  $T$ , given  $T > t$  or the hazard rate, is given by  $h(t) = f(t/T > t)$

$$\begin{aligned} &= \frac{f(t)}{1 - F(t)} = \frac{F'(t)}{1 - F(t)} \\ \therefore \int_0^t h(t) dt &= \int_0^t \frac{F'(t)}{1 - F(t)} dt \\ &= [-\log\{1 - F(t)\}]_0^t \\ &= -\log\{1 - F(t)\} \quad [\text{since } F(0) = P(T \leq 0) = 0 \text{ as the system was put into operation at } t = 0] \\ \therefore F(t) &= 1 - e^{-\int_0^t h(t) dt} \end{aligned}$$

$$\therefore F(t) = h(t) \times e^{-\int_0^t h(t) dt}$$

When  $h(t) = t$ ,

$$\begin{aligned} f(t) &= t \times e^{-\int_0^t t dt} \\ &= te^{-t^2/2}, \text{ which is the pdf of a Raleigh distribution} \end{aligned}$$

### Example 20

If a continuous RV  $X$  follows  $N(0, 2)$ , find

$$P\{1 \leq X \leq 2\} \text{ and } P\{1 \leq X \leq 2/X \geq 1\}$$

$X$  follows  $N(0, 2)$ , the density function of which is  $f(x) = \frac{1}{2\sqrt{2\pi}} e^{-x^2/8}$ ,  $-\infty < x < \infty$ .

$$\begin{aligned} \therefore P\{1 \leq X \leq 2\} &= \int_1^2 f(x) dx \\ &= \int_{0.5}^1 \phi(t) dt, \text{ putting } t = \frac{x}{2} \end{aligned}$$

where  $\phi(t)$  is the standard normal density.

$$\begin{aligned} &= \int_0^1 \phi(t) dt - \int_0^{0.5} \phi(t) dt \\ &= 0.3413 - 0.1915 \text{ (from the normal tables)} \\ &= 0.1498 \end{aligned}$$

$$\begin{aligned} P\{1 \leq X \leq 2/X \geq 1\} &= \frac{P\{(1 \leq X \leq 2) \text{ and } X \geq 1\}}{P\{X \geq 1\}} \\ &= \frac{P\{1 \leq X \leq 2\}}{P\{1 \leq X < \infty\}} \\ &= \frac{0.1498}{\int_0^\infty \phi(t) dt - \int_0^{0.5} \phi(t) dt} \\ &= \frac{0.1498}{0.5 - 0.1915} = 0.4856 \end{aligned}$$

## Exercise 2(A)

**Part-A (Short-answer Questions)**

1. Define a RV with an example.
2. Define a discrete RV with an example.
3. Define a continuous RV and give an example for the same.
4. Distinguish between a discrete RV and a continuous RV.
5. Define the probability mass function of a discrete RV.
6. Write down the probability distribution of the outcome when 2 fair dice are tossed.
7. Define the pdf of a continuous RV.
8. State the properties of the pdf of a continuous RV.
9. What is the probability curve of a continuous RV? Give an example.
10. Prove that it is almost impossible that a continuous RV assumes a specific value. (OR) If  $X$  is a continuous RV prove that  $P(X = a) = 0$ .
11. If  $X$  represents the total number of heads obtained, when a fair coin is tossed 5 times, find the probability distribution of  $X$ .
12. If the probability distribution of  $X$  is given as:

$x:$	1	2	3	4
$p_x:$	0.4	0.3	0.2	0.1

Find  $P\left(\frac{1}{2} < X < \frac{7}{2} / X > 1\right)$

13. Define the cdf of a RV. Explain how to find it for both kinds of RV.
14. Differentiate between the pdf and cdf of a RV.
15. State the properties of the cdf of a RV.

16. Verify whether  $f(x) = \begin{cases} |x| & \text{in } -1 \leq X \leq 1 \\ 0, & \text{elsewhere} \end{cases}$

can be the pdf of a continuous RV.

17. If  $f(x) = kx^2$ ,  $0 < x < 3$ , is to be a density function, find the value of  $k$ .
18. If the pdf of a RV  $X$  is given by

$$f(x) = \begin{cases} 1/4 & \text{in } -2 < x < 2 \\ 0, & \text{elsewhere} \end{cases}$$

Find  $P\{|X| > 1\}$ .

19. Find the value of  $k$ , if  $f(x) = \begin{cases} kxe^{-x} & x > 0 \\ 0 & \text{elsewhere} \end{cases}$

is the pdf of a RV  $X$ .

20. If the pdf of a RV  $X$  is  $f(x) = \frac{x}{2}$  in  $0 \leq x \leq 2$ , find  $P(X > 1.5/X > 1)$ .

21. The RV  $X$  has the following probability distribution:

$x:$	-2	-1	0	1
$p_x:$	0.4	$k$	0.2	0.3

Find  $k$  and the mean value of  $X$ .

22. If  $X$  represents the outcome of the toss of a 6 faced dice, find  $P(X \leq x)$  as a function of  $x$ .
23. If the pdf of a RV  $X$  is  $f(x) = 2x$ ,  $0 < x < 1$ , find the cdf of  $X$ .
24. If the cdf of a RV  $X$  is given by  $F(x) = 1 - e^{-\lambda x}$ , when  $x \geq 0$  and = 0, when  $x < 0$ , find the pdf of  $X$ .
25. If the cdf of a RV is given by  $F(x) = 0$ , for  $x < 0$ ;  $= x^2/16$  for  $0 \leq x \leq 4$  and = 1, for  $4 \leq x$ , find  $P(X > 1/X < 3)$ .
26. Define binomial distribution. What are its mean and variance?
27. Give the probability law of Poisson distribution and also its mean and variance.
28. Define the exponential distribution.
29. If  $X$  follows an exponential distribution with parameter 1, find  $P(|X| \leq 1)$ .
30. Define Pascal distribution and define geometric distribution as a particular case of Pascal distribution.
31. Write down the pdf's of general normal distribution and standard normal distribution.
32. Define Erlang distribution. Deduce Gamma distribution as a particular case of Erlang distribution.
33. Deduce the pdf of an exponential distribution as a particular case of that of Erlang distribution.
34. Give the pdf of Rayleigh distribution.
35. Define Maxwell distribution.
36. Write down the pdf of Laplace distribution.
37. Define Cauchy distribution.

### Part-B

38. Find the formula for the probability distribution of the number of heads, when a fair coin is tossed 4 times.
39. A coin is known to come up heads 3 times as often as tails. This coin is tossed 3 times. Write down the probability distribution of the number of heads that appear and also the cdf. Make a sketch of both.
40. Consider the experiment of tossing a coin, the 2 events of the space being occurrence of head or tail. Assign probabilities  $p$  and  $q$  for head and tail respectively and define a random variable  $X$  by  $X(h) = 1$  and  $X(t) = 0$ . Determine and plot the probability function  $f(x)$  and the distribution function  $F(x)$ .
41. Two dice are tossed. If  $X$  is the sum of the numbers shown up, find the probability mass function of  $X$ .
42. Consider the experiment of tossing a fair coin 4 times. Define  $X = 0$ , if 0 or 1 head appears;  $X = 1$ , if 2 heads appear;  $X = 2$ , if 3 or 4 heads appear. Find the probability function, mean and variance of  $X$ .

43. A discrete RV  $X$  has the following probability distribution.

$x:$	0	1	2	3	4	5	6	7	8
$p(x):$	$a$	$3a$	$5a$	$7a$	$9a$	$11a$	$13a$	$15a$	$17a$

Find the value of  $a$ ,  $P(X < 3)$ , variance and distribution function of  $X$ .

44. The probability distribution of a RV  $X$  is given below:

$x:$	0	1	2	3
$p(x):$	0.1	0.3	0.5	0.1

If  $Y = X^2 + 2X$ , find the probability distribution, mean and variance of  $Y$ .

45. The probability mass function of a RV  $X$  is defined as

$P(X = 0) = 3C^2$ ,  $P(X = 1) = 4C - 10C^2$  and  $P(X = 2) = 5C - 1$ , where  $C > 0$  and  $P(X = r) = 0$ , if  $r \neq 0, 1, 2$ .

(i) Find the value of  $C$ , (ii) Find  $P\{0 < X < 2/X > 0\}$ , (iii) the distribution function of  $X$ , (iv) the largest value of  $X$  for which  $F(x) < \frac{1}{2}$ , and (v) the smallest value of  $X$  for which  $F(x) > \frac{1}{2}$ .

46. If the probability mass function of a RV  $X$  is given by  $P(X = r) = kr^3$ ;  $r = 1, 2, 3, 4$ , find (i) the value of  $k$  (ii)  $P(1/2 < X \leq 5/2 / X > 1)$ , (iii) the mean and variance of  $X$ , and (iv) the distribution function of  $X$ .

47. Find the values of  $a$  for which  $P(x=j) = (1-a)a^j$ ,  $j = 0, 1, 2, \dots$ , represents a probability mass function. Show also that for any 2 positive integers  $m$  and  $n$

$$P(X > m + n/X > m) = P(X \geq n).$$

48. If a discrete probability distribution is given by  $P(X = r) = k(1-a)^{r-1}$ ,  $0 < a < 1$ , for  $r = 1, 2, \dots, \infty$ , find the value of  $k$  and also the mean and variance of  $X$ .

49. If the probability distribution of a discrete RV  $X$  is given by  $P(X = x) = ke^{-t}(1-e^{-t})^{x-1}$ ,  $x = 1, 2, \dots, \infty$ , find the value of  $k$  and also the mean and variance of  $X$ .

50. In a continuous distribution, the probability density is given by  $f(x) = kx(2-x)$ ,  $0 < x < 2$ . Find  $k$ , mean, variance and the distribution function.

51. The diameter of an electric cable  $X$  is a continuous RV with pdf  $f(x) = kx(1-x)$ ,  $0 \leq x \leq 1$ . Find (i) the value of  $k$ , (ii) cdf of  $X$ , (iii) the value of  $a$  such that  $P(X < a) = 2P(X > a)$ , and (iv)  $P(X \leq 1/2 / 1/3 < X < 2/3)$ .

52.  $X$  is a continuous RV with pdf given by  $f(x) = kx$ , in  $0 \leq x \leq 2$ ;  $= 2k$ , in  $2 \leq x \leq 4$ , and  $= 6k - kx$ , in  $4 \leq x \leq 6$ . Find the value of  $k$  and  $F(x)$ .

53. The continuous RV  $X$  has pdf  $f(x) = \frac{x}{2}$ ,  $0 \leq x \leq 2$ . Two independent determinations of  $X$  are made. What is the probability that both these determinations will be greater than 1? If 3 independent determinations had been made, what is the probability that exactly 2 of these are larger than 1?

54. A continuous RVX that can assume values between  $x = 2$  and  $x = 5$  has a density function given by  $f(x) = 2(1+x)/27$ . Find  $P(3 < X < 4)$ .
55. A continuous RV has the pdf  $f(x) = kx^4$ ,  $-1 < x < 0$ . Find the value of  $k$  and also  $P(X > -1/2 | X < -1/4)$ .
56. Suppose that the life length of a certain radio tube (in hours) is a continuous RV  $X$  with pdf  $f(x) = \frac{100}{x^2}$ ,  $x > 100$  and = 0, elsewhere.
- What is the probability that a tube will last less than 200 h, if it is known that the tube is still functioning after 150 h of service?
  - What is the probability that if 3 such tubes are installed in a set, exactly 1 will have to be replaced after 150 h of service?
  - What is the maximum number of tubes that may be inserted into a set so that there is a probability of 0.1 that after 150 h of service all of them are still functioning?
57. If the cdf of a continuous RV  $X$  is given by  $F(x) = \frac{1}{2}e^{kx}$ ,  $x \leq 0$ , and  $F(x) = 1 - \frac{1}{2}e^{-kx}$ ,  $x > 0$ , find  $P(|x| \leq 1/k)$ . Prove that the density function of  $X$  is  $f(x) = \frac{k}{2}e^{-k|x|}$ ,  $-\infty < x < \infty$ , given that  $k > 0$ .
58. If the distribution function of a continuous RV  $X$  is given by  $F(x) = 0$ , when  $x < 0$ ;  $= x$ , when  $0 \leq x \leq 1$  and  $= 1$ , when  $1 \leq x$ , find the pdf of  $X$ . Also find  $P(1/3 < X < 1/2)$  and  $P(1/2 < X < 2)$  using the cdf of  $X$ .
59. A point is chosen on a line of length  $a$  at random. What is the probability that the ratio of the shorter to the longer segment is less than  $1/4$ ?
60. If the RV  $k$  is uniformly distributed over  $(1, 7)$  what is the probability that the roots of the equation  $x^2 + 2kx + (2k + 3) = 0$  are real?
61. If  $f(t)$  is the unconditional density of time to failure  $T$  of a system and  $h(t)$  is the conditional density of  $T$ , given  $T > t$ , find  $h(t)$  when (i)  $f(t) = \lambda e^{-\lambda t}$ , and (ii)  $f(t) = \lambda^2 t e^{-\lambda t}$ ,  $t > 0$ . Prove also that  $h(t)$  is not a density function.
62. If the continuous RV  $X$  follows  $N(1000, 20)$ , find
  - $P(X < 1024)$ ,
  - $P(X < 1024 | X > 961)$  and
  - $P(31 < \sqrt{X} \leq 32)$ .

## Two-Dimensional Random Variables

So far we have considered only the one-dimensional RV, i.e., we have considered such random experiments, the outcome of which had only one characteristic and hence, was assigned a single real value. In many situations, we will be interested in recording two or more characteristics (numerically) of the outcome of a random experiment. For example, both voltage and current might be of interest in a certain experiment.

**Definitions:** Let  $S$  be the sample space associated with a random experiment  $E$ . Let  $X = X(s)$  and  $Y = Y(s)$  be two functions each assigning a real number to each outcomes  $s \in S$ . Then  $(X, Y)$  is called a *two-dimensional random variable*.

If the possible values of  $(X, Y)$  are finite or countably infinite,  $(X, Y)$  is called a *two-dimensional discrete RV*. When  $(X, Y)$  is a two-dimensional discrete RV the possible values of  $(X, Y)$  may be represented as  $(x_i, y_j)$ ,  $i = 1, 2, \dots, m, \dots; j = 1, 2, \dots, n, \dots$ .

If  $(X, Y)$  can assume all values in a specified region  $R$  in the  $xy$ -plane,  $(X, Y)$  is called a *two-dimensional continuous RV*.

## Probability Function of $(X, Y)$

If  $(X, Y)$  is a two-dimensional discrete RV such that  $P(x = x_i, y = y_j) = p_{ij}$ , then  $p_{ij}$  is called the *probability mass function* or simply the *probability function* of  $(X, Y)$  provided the following conditions are satisfied.

$$(i) \quad p_{ij} \geq 0, \text{ for all } i \text{ and } j$$

$$(ii) \quad \sum_j \sum_i p_{ij} = 1$$

The sets of triples  $\{x_i, y_j, p_{ij}\}$ ,  $i = 1, 2, \dots, m, \dots, j = 1, 2, \dots, n, \dots$ , is called *the joint probability distribution of  $(X, Y)$* .

## Joint Probability Density Function

If  $(X, Y)$  is a two-dimensional continuous RV such that

$$P\left\{x - \frac{dx}{2} \leq X \leq x + \frac{dx}{2} \text{ and } y - \frac{dy}{2} \leq Y \leq y + \frac{dy}{2}\right\} = f(x, y) dx dy, \text{ then } f(x, y) \text{ is}$$

called *the joint pdf of  $(X, Y)$* , provided  $f(x, y)$  satisfies the following conditions.

$$(i) \quad f(x, y) \geq 0, \text{ for all } (x, y) \in R, \text{ where } R \text{ is the range space.}$$

$$(ii) \quad \iint_R f(x, y) dx dy = 1.$$

Moreover, if  $D$  is a subspace of the range space  $R$ ,  $P\{(X, Y) \in D\}$  is defined as

$$P\{(X, Y) \in D\} = \iint_D f(x, y) dx dy. \text{ In particular}$$

$$P\{a \leq X \leq b, c \leq Y \leq d\} = \int_c^d \int_a^b f(x, y) dx dy$$

## Cumulative Distribution Function

If  $(X, Y)$  is a two-dimensional RV (discrete or continuous), then  $F(x, y) = P\{X \leq x \text{ and } Y \leq y\}$  is called *the cdf of  $(X, Y)$* .

In the discrete case,

$$F(x, y) = \sum_j \sum_i p_{ij} \\ y_j \leq y \quad x_i \leq x$$

In the continuous case,

$$F(x, y) = \int_{-\infty}^y \int_{-\infty}^x f(x, y) dx dy$$

### Properties of $F(x, y)$

- (i)  $F(-\infty, y) = 0 = F(x, -\infty)$  and  $F(\infty, \infty) = 1$
- (ii)  $P\{a < X < b, Y \leq y\} = F(b, y) - F(a, y)$
- (iii)  $P\{X \leq x, c < Y < d\} = F(x, d) - F(x, c)$
- (iv)  $P\{a < X < b, c < Y < d\} = F(b, d) - F(a, d) - F(b, c) + F(a, c)$
- (iv) At points of continuity of  $f(x, y)$

$$\frac{\partial^2 F}{\partial x \partial y} = f(x, y)$$

### Marginal Probability Distribution

$$\begin{aligned} P(X = x_i) &= P\{X = x_i \text{ and } Y = y_1\} \text{ or } (X = x_i \text{ and } Y = y_2) \text{ or etc.}\} \\ &= p_{i1} + p_{i2} + \dots = \sum_j p_{ij} \end{aligned}$$

$P(X = x_i) = \sum_j p_{ij}$  is called the *marginal probability function of X*. It is defined for  $X = x_1, x_2, \dots$  and denoted as  $P_{i*}$ . The collection of pairs  $\{x_i, p_{i*}\}, i = 1, 2, 3, \dots$  is called the *marginal probability distribution of X*.

Similarly the collection of pairs  $\{y_j, p_{*j}\}, j = 1, 2, 3, \dots$  is called the *marginal probability distribution of Y*, where  $p_{*j} = \sum_i P_{ij} = P(Y = y_j)$ .

In the continuous case,

$$\begin{aligned} P\left\{x - \frac{1}{2} dx \leq X \leq x + \frac{1}{2} dx, -\infty < Y < \infty\right\} \\ &= \int_{-\infty}^{\infty} \int_{x - \frac{1}{2} dx}^{x + \frac{1}{2} dx} f(x, y) dx dy \\ &= \left[ \int_{-\infty}^{\infty} f(x, y) dy \right] dx \quad [\text{since } f(x, y) \text{ may be treated a constant in } (x - 1/2 dx, x + 1/2 dx)] \\ &= f_X(x) dx, \text{ say} \end{aligned}$$

$$f_X(x) = \int_{-\infty}^{\infty} f(x, y) dy \text{ is called the marginal density of } X.$$

$$\text{Similarly, } f_Y(y) = \int_{-\infty}^{\infty} f(x, y) dx \text{ is called the marginal density of } Y.$$

**Note**

$$\begin{aligned}
 P(a \leq X \leq b) &= P(a \leq X \leq b, -\infty < Y < \infty) \\
 &= \int_{-\infty}^{\infty} \int_a^b f(x, y) dx dy \\
 &= \int_a^b \left[ \int_{-\infty}^{\infty} f(x, y) dy \right] dx = \int_a^b f_X(x) dx \\
 \text{Similarly, } P(c \leq Y \leq d) &= \int_c^d f_Y(y) dy
 \end{aligned}$$

**Conditional Probability Distribution**

$P(X = x_i | Y = y_j) = \frac{P\{X = x_i, Y = y_j\}}{P\{Y = y_j\}} = \frac{p_{ij}}{p_{*j}}$  is called *the conditional probability function of X, given that Y = y<sub>j</sub>*.

The collection of pairs,  $\left\{ x_i, \frac{p_{ij}}{p_{*j}} \right\} i = 1, 2, 3, \dots$ , is called *the conditional probability distribution of X, given Y = y<sub>j</sub>*.

Similarly, the collection of pairs,  $\left\{ y_j, \frac{p_{ij}}{p_{i*}} \right\} j = 1, 2, 3, \dots$ , is called *the conditional probability distribution of Y, given X<sub>i</sub> = x<sub>i</sub>*. In the continuous case,

$$\begin{aligned}
 P\left\{ x - \frac{1}{2} dx \leq X < x + \frac{1}{2} dx | Y = y \right\} \\
 &= P\left\{ x - \frac{1}{2} dx \leq X \leq x + \frac{1}{2} dx | Y - \frac{1}{2} dy \leq Y \leq y + \frac{1}{2} dy \right\} \\
 &= \frac{f(x, y) dx dy}{f_Y(y) dy} = \left\{ \frac{f(x, y)}{f_Y(y)} \right\} dx.
 \end{aligned}$$

$\frac{f(x, y)}{f_Y(y)}$  is called *the conditional density of X, given Y*, and is denoted by  $f(x|y)$ .

Similarly,  $\frac{f(x, y)}{f_X(y)}$  is called *the conditional density of Y, given X*, and is denoted by  $f(y|x)$ .

**Independent RVs**

If  $(X, Y)$  is a two-dimensional discrete RV such that  $P\{X = x_i | Y = y_j\} = P(X = x_i)$ , i.e.,  $\frac{p_{ij}}{p_{*j}} = p_{i*}$ , i.e.,  $p_{ij} = p_{i*} \times p_{*j}$  for all  $i, j$  then  $X$  and  $Y$  are said to be independent RVs.

Similarly, if  $(X, Y)$  is a two-dimensional continuous RV such that  $f(x, y) = f_X(x) \times f_Y(y)$ , then  $X$  and  $Y$  are said to be independent RVs.

## Random Vectors

Sometimes we may have to be concerned with random experiments whose outcomes will have three or more simultaneous numerical characteristics. To study the outcomes of such random experiments we require knowledge of *n-dimensional random variables* or *random vectors*. For example, the location of a space vehicle in a cartesian co-ordinate system is a three-dimensional random vector.

Most of the concepts introduced above for the two-dimensional case can be extended to the *n*-dimensional one.

**Definitions:** A vector  $X: [X_1, X_2, \dots, X_n]$  whose components  $X_i$  are RVs is called a *random vector*.  $(X_1, X_2, \dots, X_n)$  can assume all values in some region  $R_n$  of the *n*-dimensional space.  $R_n$  is called the *range space*.

The joint distribution function of  $(X_1, X_2, \dots, X_n)$  is defined as  $F(x_1, x_2, \dots, x_n) = P[X_1 \leq x_1, X_2 \leq x_2, \dots, X_n \leq x_n]$

The joint pdf of  $[X_1, X_2, \dots, X_n]$  is defined as  $f(x_1, x_2, \dots, x_n)$

$$= \frac{\partial^n F(x_1, x_2 \dots x_n)}{\partial x_1 \cdot \partial x_2 \cdots \partial x_n} \text{ and satisfies the following conditions.}$$

(i)  $f(x_1, x_2, \dots, x_n) \geq 0$ , for all  $(x_1, x_2, \dots, x_n)$

(ii)  $\iint_{R_n} \cdots \int f(x_1, x_2, \dots, x_n) dx_1 dx_2 \dots dx_n = 1$

(iii)  $P[(X_1, X_2, \dots, X_n) \in D] = \iint_D \cdots \int f(x_1, x_2, \dots, x_n) dx_1 dx_2 \dots dx_n$  where  $D$  is a subset of the range space  $R_n$ .

The marginal pdf of any subset of the *n* RVs  $X_1, X_2, \dots, X_n$  is obtained by “integrating out” the variables not in the subset. For example, if  $n = 3$ , then

$$f_{X_1}(x_1) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(x_1, x_2, x_3) dx_2 dx_3 \text{ is the marginal pdf of the one-dimensional}$$

$$\text{RV } X_1 \text{ and } f_{x_1, x_2}(x_1, x_2) = \int_{-\infty}^{\infty} f(x_1, x_2, x_3) dx_3 \text{ is the marginal joint pdf of the two-}$$

dimensional RV  $(X_1, X_2)$ . The concept of independent RVs is also extended in a natural way. The RVs  $(X_1, X_2, \dots, X_n)$  are said to be independent, if  $f(x_1, x_2, \dots, x_n) = f_{X_1}(x_1) \cdot f_{X_2}(x_2) \cdots f_{X_n}(x_n)$ .

The conditional density functions are defined as in the following examples.

If  $n = 3$ ,

$$f(x_1, x_2/x_3) = \frac{f(x_1, x_2, x_3)}{f_{X_3}(x_3)} \text{ and}$$

$$f(x_1/x_2, x_3) = \frac{f(x_1, x_2, x_3)}{f_{X_2, X_3}(x_2, x_3)}$$

### Worked Example 2(B)

#### Example 1

Three balls are drawn at random without replacement from a box containing 2 white, 3 red and 4 black balls. If  $X$  denotes the number of white balls drawn and  $Y$  denotes the number of red balls drawn, find the joint probability distribution of  $(X, Y)$ .

As there are only 2 white balls in the box,  $X$  can take the values 0, 1 and 2 and  $Y$  can take the values 0, 1, 2 and 3.

$$\begin{aligned} P(X = 0, Y = 0) &= P(\text{drawing 3 balls none of which is white or red}) \\ &= P(\text{all the 3 balls drawn are black}) \\ &= 4C_3 / 9C_3 = \frac{1}{21} \end{aligned}$$

$$\begin{aligned} P(X = 0, Y = 1) &= P(\text{drawing 1 red and 2 black balls}) \\ &= \frac{3C_1 \times 4C_2}{9C_3} = \frac{3}{14} \end{aligned}$$

$$\text{Similarly, } P(X = 0, Y = 2) = \frac{3C_2 \times 4C_1}{9C_3} = \frac{1}{7}; P(X = 0, Y = 3) = \frac{1}{84}$$

$$P(X = 1, Y = 0) = \frac{1}{7}; P(X = 1, Y = 1) = \frac{2}{7}; P(X = 1, Y = 2) = \frac{1}{14};$$

$P(X = 1, Y = 3) = 0$  (since only 3 balls are drawn)

$$P(X = 2, Y = 0) = \frac{1}{21}; P(X = 2, Y = 1) = \frac{1}{28}; P(X = 2, Y = 2) = 0;$$

$$P(X = 2, Y = 3) = 0$$

The joint probability distribution of  $(X, Y)$  may be represented in the form of a table as given below:

X	Y			
	0	1	2	3
0	$\frac{1}{21}$	$\frac{3}{14}$	$\frac{1}{7}$	$\frac{1}{84}$
1	$\frac{1}{7}$	$\frac{2}{7}$	$\frac{1}{14}$	0
2	$\frac{1}{21}$	$\frac{1}{28}$	0	0

**Note** Sum of all the cell probabilities = 1.

### Example 2

For the bivariate probability distribution of  $(X, Y)$  given below, find  $P(X \leq 1)$ ,  $P(Y \leq 3)$ ,  $P(X \leq 1, Y \leq 3)$ ,  $P(X \leq 1/Y \leq 3)$ ,  $P(Y \leq 3/X \leq 1)$  and  $P(X + Y \leq 4)$ .

$X \backslash Y$	1	2	3	4	5	6
0	0	0	1/32	2/32	2/32	3/32
1	1/16	1/16	1/8	1/8	1/8	1/8
2	1/32	1/32	1/64	1/64	0	2/64

$$\begin{aligned}
 P(X \leq 1) &= P(X = 0) + P(X = 1) \\
 &= \sum_{j=1}^6 P(X = 0, Y = j) + \sum_{j=1}^6 P(X = 1, Y = j) \\
 &= \left(0 + 0 + \frac{1}{32} + \frac{2}{32} + \frac{2}{32} + \frac{3}{32}\right) + \left(\frac{1}{16} + \frac{1}{16} + \frac{1}{8} + \frac{1}{8} + \frac{1}{8} + \frac{1}{8}\right) \\
 &= \frac{1}{4} + \frac{5}{8} = \frac{7}{8}
 \end{aligned}$$

$$\begin{aligned}
 P(Y \leq 3) &= P(Y = 1) + P(Y = 2) + P(Y = 3) \\
 &= \sum_{i=0}^2 P(X = i, Y = 1) + \sum_{i=0}^2 P(X = i, Y = 2) + \sum_{i=0}^2 P(X = i, Y = 3) \\
 &= \left(0 + \frac{1}{16} + \frac{1}{32}\right) + \left(0 + \frac{1}{16} + \frac{3}{32}\right) + \left(\frac{1}{32} + \frac{1}{8} + \frac{1}{64}\right) \\
 &= \frac{3}{32} + \frac{3}{32} + \frac{11}{64} = \frac{23}{64}
 \end{aligned}$$

$$P(X \leq 1, Y \leq 3) = \sum_{j=1}^3 P(X = 0, Y = j) + \sum_{j=1}^3 P(X = 1, Y = j)$$

$$= \left( 0 + 0 + \frac{1}{32} \right) + \left( \frac{1}{16} + \frac{1}{16} + \frac{1}{8} \right) = \frac{9}{32}$$

$$P(X \leq 1 | Y \leq 3) = \frac{P(X \leq 1, Y \leq 3)}{P(Y \leq 3)} = \frac{9/32}{23/64} = \frac{18}{23}$$

$$P(Y \leq 3 | X \leq 1) = \frac{P(X \leq 1, Y \leq 3)}{P(X \leq 1)} = \frac{9/32}{7/8} = \frac{9}{28}$$

$$\begin{aligned} P(X + Y \leq 4) &= \sum_{j=1}^4 P(X = 0, Y = j) + \sum_{j=1}^3 P(X = 1, Y = j) + \sum_{j=1}^2 P(X = 2, Y = j) \\ &= \frac{3}{32} + \frac{1}{4} + \frac{1}{16} = \frac{13}{32} \end{aligned}$$

### Example 3

The joint probability mass function of  $(X, Y)$  is given by  $p(x, y) = k(2x + 3y)$ ,  $x = 0, 1, 2$ ;  $y = 1, 2, 3$ . Find all the marginal and conditional probability distributions. Also find the probability distribution of  $(X + Y)$ .

The joint probability distribution of  $(X, Y)$  is given below. The relevant probabilities have been computed by using the given law.

X	Y		
	1	2	3
0	$3k$	$6k$	$9k$
1	$5k$	$8k$	$11k$
2	$7k$	$10k$	$13k$

$$\sum_{j=1}^3 \sum_{i=0}^2 p(x_i, y_j) = 1$$

i.e., the sum of all the probabilities in the table is equal to 1.

i.e.,  $72k = 1$ .

$$\therefore k = \frac{1}{72}$$

**Marginal Probability Distribution of  $X$ :  $\{i, p_{i*}\}$** 

$X = i$	$p_{i*} = \sum_{j=1}^3 p_{ij}$
0	$p_{01} + p_{02} + p_{03} = \frac{18}{72}$
1	$p_{11} + p_{12} + p_{13} = \frac{24}{72}$
2	$p_{21} + p_{22} + p_{23} = \frac{30}{72}$
Total = 1	

**Marginal Probability Distribution of  $Y$ :  $\{j, p_{*j}\}$** 

$Y = j$	$p_{*j} = \sum_{i=0}^2 p_{ij}$
1	15/27
2	24/72
3	33/72
Total = 1	

Conditional distribution of  $X$ , given  $Y = 1$ , is given by  $\{i, P(X = i|Y = 1)\} = \{i, P(X = i, Y = 1)/P(Y = 1)\} = \{i, p_{i1}/p_{*1}\} = 0, 1, 2$ .

The tabular representation is given below:

$X = i$	$p_{i1}/p_{*1}$
0	$3k/15k = \frac{1}{5}$
1	$5k/15k = \frac{1}{3}$
2	$7k/15k = \frac{7}{15}$
Total = 1	

The other conditional distributions are given below:

C.P.D. of $X$ , given $Y = 2$	
$X = i$	$p_{i2}/p_{*2}$
0	$\frac{6k}{24k} = \frac{1}{4}$
1	$\frac{8k}{24k} = \frac{1}{3}$
2	$\frac{10k}{24k} = \frac{5}{12}$
Total = 1	

C.P.D. of $X$ , given $Y = 3$	
$X = i$	$p_{i3}/p_{*3}$
0	$\frac{9k}{33k} = \frac{3}{11}$
1	$\frac{11k}{33k} = \frac{1}{3}$
2	$\frac{13k}{33k} = \frac{13}{33}$
Total = 1	

C.P.D. of $Y$ , given $X = 0$	
$Y = j$	$p_{0j}/p_{0*}$
1	$\frac{3k}{18k} = \frac{1}{6}$
2	$\frac{6k}{18k} = \frac{1}{3}$
3	$\frac{9k}{18k} = \frac{1}{2}$
Total = 1	

C.P.D. of $Y$ , given $X = 1$	
$Y = j$	$p_{1j}/p_{1*}$
1	$\frac{5k}{24k} = \frac{5}{24}$
2	$\frac{8k}{24k} = \frac{1}{3}$
3	$\frac{11k}{24k} = \frac{11}{24}$
Total = 1	

C.P.D. of $Y$ , given $X = 2$	
$Y = j$	$p_{2j}/p_{2*}$
1	$\frac{7k}{30k} = \frac{7}{30}$
2	$\frac{10k}{30k} = \frac{1}{3}$
3	$\frac{13k}{30k} = \frac{13}{30}$
Total = 1	

Probability distribution of $(X + Y)$	
$(X + Y)$	$P$
1	$p_{01} = \frac{3}{72}$
2	$p_{02} + p_{11} = \frac{11}{72}$
3	$p_{03} + p_{12} + p_{21} = \frac{24}{72}$
4	$p_{13} + p_{22} = \frac{21}{72}$
5	$p_{23} = \frac{13}{72}$
	Total = 1

### Example 4

A machine is used for a particular job in the forenoon and for a different job in the afternoon. The joint probability distribution of  $(X, Y)$ , where  $X$  and  $Y$  represent the number of times the machine breaks down in the forenoon and in the afternoon respectively, is given in the following table. Examine if  $X$  and  $Y$  are independent RVs.

X	Y		
	0	1	2
0	0.1	0.04	0.06
1	0.2	0.08	0.12
2	0.2	0.08	0.12

$X$  and  $Y$  are independent, if  $P_{i*} \times P_{*j} = P_{ij}$  for all  $i$  and  $j$ . So, let us find  $P_{i*} \times P_{*j}$  for all  $i$  and  $j$ .

$$P_{0*} = 0.1 + 0.04 + 0.06 = 0.2; P_{1*} = 0.4; P_{2*} = 0.4$$

$$P_{*0} = 0.5; P_{*1} = 0.2; P_{*2} = 0.3$$

$$\text{Now, } P_{0*} \times P_{*0} = 0.2 \times 0.5 = 0.1 = P_{00}$$

$$P_{0*} \times P_{*1} = 0.2 \times 0.2 = 0.04 = P_{01}$$

$$P_{0*} \times P_{*2} = 0.2 \times 0.3 = 0.06 = P_{02}$$

Similarly, we can verify that

$$P_{1*} \times P_{*0} = P_{10}; P_{1*} \times P_{*1} = P_{11}; P_{1*} \times P_{*2} = P_{12};$$

$$P_{2*} \times P_{*0} = P_{20}; P_{2*} \times P_{*1} = P_{21}; P_{2*} \times P_{*2} = P_{22}$$

Hence, the RVs  $X$  and  $Y$  are independent.

**Example 5**

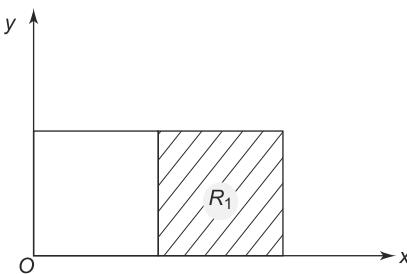
The joint pdf of a two-dimensional RV  $(X, Y)$  is given by  $f(x, y) = xy^2 + \frac{x^2}{8}$ ,  $0 \leq x \leq 2$ ,  $0 \leq y \leq 1$ .

Compute  $P(X < 1)$ ,  $P\left(Y < \frac{1}{2} / X > 1\right)$ ,  $P(X > 1 / Y < 1/2)$

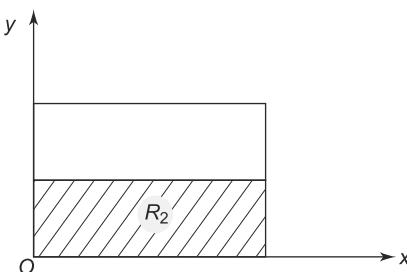
$P\left(Y < \frac{1}{2} / X > 1\right)$ ,  $P(X < Y)$  and  $P(X + Y \leq 1)$ .

Here, the rectangular defined by  $0 \leq x \leq 2$ ,  $0 \leq y \leq 1$  is the range space  $R$ .  $R_1$ ,  $R_2$ , ..., are event spaces.

$$(a) \quad P(X > 1) = \int_{R_1} \int_{(x>1)} f(x, y) dx dy = \int_0^1 \int_{1}^2 \left( xy^2 + \frac{x^2}{8} \right) dx dy = \frac{19}{24}$$

**Fig. 2.1**

$$(b) \quad P(X < 2) = \int_{R_2} \left( xy^2 + \frac{x^2}{8} \right) dx dy \\ = \int_0^{1/2} \int_0^2 \left( xy^2 + \frac{x^2}{8} \right) dx dy = \frac{1}{4}$$

**Fig. 2.2**

$$(c) \quad P(X > 1, Y < 1/2) = \int_{R_3} \left( xy^2 + \frac{x^2}{8} \right) dx dy$$

$\left( x > 1 \& y < \frac{1}{2} \right)$

$$= \int_0^{1/2} \int_1^2 \left( xy^2 + \frac{x^2}{8} \right) dx dy = \frac{5}{24}$$

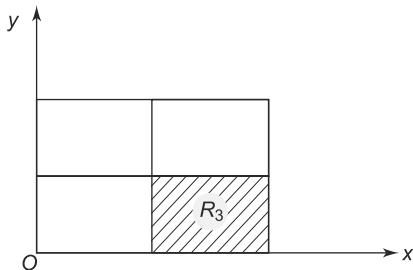


Fig. 2.3

$$(d) \quad P\left(X > 1 / Y < \frac{1}{2}\right) = \frac{P\left(X > 1, Y < \frac{1}{2}\right)}{P\left(Y < \frac{1}{2}\right)} = \frac{5/24}{1/4} = \frac{5}{6}$$

$$(e) \quad P\left(Y < \frac{1}{2} / X > 1\right) = \frac{P\left(X > 1, Y < \frac{1}{2}\right)}{P(X > 1)} = \frac{5/24}{19/24} = \frac{5}{19}$$

$$(f) \quad P(X < Y) = \int_{R_4} \int \left( xy^2 + \frac{x^2}{8} \right) dx dy$$

$(x < y)$

$$= \int_0^1 \int_0^y \left( xy^2 + \frac{x^2}{8} \right) dx dy = \frac{53}{480}$$

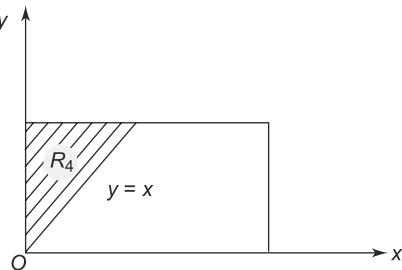


Fig. 2.4

$$(g) \quad P(X + Y \leq 1) = \int_{R_5} \int \left( xy^2 + \frac{x^2}{8} \right) dx dy$$

$$\text{where } R_5 = \{(x, y) | x + y \leq 1\}$$

$$= \int_0^{1-y} \int_0^{1-y} \left( xy^2 + \frac{x^2}{8} \right) dx dy = \frac{13}{480}$$

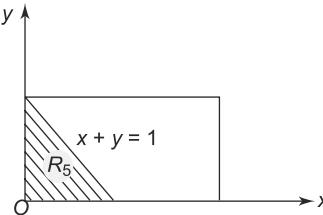


Fig. 2.5

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**Example 6**


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If the joint pdf of the RV  $(X, Y)$  is given by  $f(x, y)$

$$= \frac{1}{2\pi\sigma^2} \exp \{-(x^2 + y^2)/2\sigma^2\}, -\infty < x, y < \infty, \text{ find } P(X^2 + Y^2 \leq a^2).$$

Here, the entire  $xy$ -plane is the range space  $R$  and the event space  $D$  is the interior of the circle  $x^2 + y^2 \leq a^2$ .

$$P(X^2 + Y^2 \leq a^2) = \iint_{x^2+y^2 \leq a^2} f(x, y) dx dy$$

Transform from Cartesian system to polar system, i.e., put  $x = r \cos \theta$  and  $y = r \sin \theta$ .

Then  $dx dy = r dr d\theta$ .

The domain of integration becomes  $r \leq a$ .

$$\begin{aligned} \text{Then } P(X^2 + Y^2 \leq a^2) &= \int_0^{2\pi} \int_0^a \frac{1}{2\pi\sigma^2} e^{-r^2/2\sigma^2} r dr d\theta \\ &= \frac{1}{2\pi} \int_0^{2\pi} (-e^{-r^2/2\sigma^2})_0^a d\theta \\ &= \frac{1}{2\pi} \int_0^{2\pi} (1 - e^{-a^2/2\sigma^2}) d\theta = 1 - e^{-a^2/2\sigma^2} \end{aligned}$$

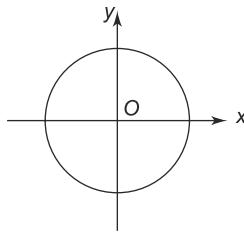


Fig. 2.6

### Example 7

A gun is aimed at a certain point (origin of the coordinate system). Because of the random factors, the actual hit point can be any point  $(X, Y)$  in a circle of radius  $R$  about the origin. Assume that the joint density of  $X$  and  $Y$  is constant in this circle given by

$$\begin{aligned} f_{XY}(x, y) &= c, \text{ for } x^2 + y^2 \leq R^2 \\ &= 0, \text{ otherwise} \end{aligned}$$

(a) Compute  $c$ , and (b) show that

$$\begin{aligned} f_X(x) &= \frac{2}{\pi R} \sqrt{1 - \left(\frac{x}{R}\right)^2}, \text{ for } -R \leq x \leq R \\ &= 0, \text{ otherwise.} \end{aligned}$$

Here, the range space is the interior of the circle  $x^2 + y^2 = R^2$ . By the property of joint pdf.

$$\iint_{x^2+y^2 \leq R} f(x, y) dx dy = 1$$

i.e.,

$$\iint_{x^2+y^2 \leq R} c dx dy = 1$$

Changing over to polar coordinates, we have

$$\int_0^{2\pi} \int_0^R c r dr d\theta = 1$$

$$\therefore c = \frac{1}{\pi R^2}$$

**Note** We have defined earlier that  $f_X(x) = \int_{-\infty}^{\infty} f(x, y) dy$ . This definition holds good if the range space is the entire  $xy$ -plane. If the range space is different from the entire  $xy$ -plane  $f_X(x)$  is given by  $\int_{\text{range}} f(x, y) dy$ , for which the limits are fixed as follows: Draw an arbitrary line parallel to  $y$ -axis (since  $x$  is to be treated as a constant). The  $y$ -coordinates of the ends of the segment of such a line that lies within the range space are the required limits. These limits will be either constants or functions of  $x$ .

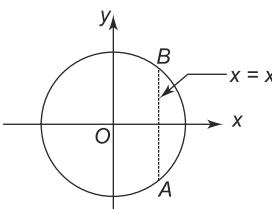


Fig. 2.7

$$\text{The point } A = \{x, -\sqrt{R^2 - x^2}\}$$

$$\text{and the point } B = \{x, \sqrt{R^2 - x^2}\}$$

Now 
$$f_X(x) = \int_{-\sqrt{R^2-x^2}}^{\sqrt{R^2-x^2}} \frac{1}{\pi R^2} dy$$

$$= \frac{2}{\pi R^2} \sqrt{R^2 - x^2} = \frac{2}{\pi R} \sqrt{1 - \left(\frac{x}{R}\right)^2} \quad -R \leq x \leq R$$

**Note** Whenever we are required to find the marginal and conditional density functions, the ranges of the concerned variables should also be specified.

### Example 8

The joint pdf of the RV  $(X, Y)$  is given by  $f(x, y) = kxye^{-(x^2+y^2)}$   $x > 0, y > 0$ . Find the value of  $k$  and prove also that  $X$  and  $Y$  are independent.

Here, the range space is the entire first quadrant of the  $xy$ -plane.

By the property of the joint pdf

$$\iint_{x>0, y>0} kxy e^{-(x^2+y^2)} dx dy = 1$$

i.e., 
$$k \int_0^\infty ye^{-y^2} dy \int_0^\infty xe^{-x^2} dx = 1$$

i.e., 
$$\frac{k}{4} = 1$$
  

$$\therefore k = 4$$

Now, 
$$f_X(x) = \int_0^\infty 4x e^{-x^2} \times ye^{-y^2} dy = 2x e^{-x^2}, x > 0$$

Similarly, 
$$f_Y(y) = 2ye^{-y^2}, y > 0$$

Now 
$$f_X(x) \times f_Y(y) = 4xye^{-(x^2+y^2)} = f(x, y)$$

$\therefore$  the RVs  $x$  and  $y$  are independent.

**Note** If  $f(x, y)$  can be factorised as  $f_1(x) \times f_2(y)$  then  $X$  and  $Y$  will be independent.

### Example 9

Given  $f_{XY}(x, y) = cx(x-y)$ ,  $0 < x < 2$ ,  $-x < y < x$ , and 0 elsewhere, (a) evaluate  $c$ , (b) find  $f_X(x)$ , (c)  $f_{XY}(y/x)$ , and (d)  $f_Y(y)$ .

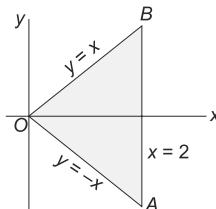


Fig. 2.8

Here, the range space is the area within the triangle  $OAB$  (shown in the figure), defined by  $0 < x < 2$  and  $-x < y < x$ .

(a) By the property of  $j$  pdf,

$$\iint_{\Delta OAB} cx(x-y)dx dy = 1$$

$$\int_0^2 \int_{-x}^x cx(x-y)dy dx = 1$$

$$\text{i.e., } 8c = 1$$

$$\therefore c = \frac{1}{8}$$

$$(b) f_X(x) = \int_{-x}^x \frac{1}{8}x(x-y)dy \\ = \frac{x^3}{4}, \text{ in } 0 < x < 2$$

$$(c) f_X(y/x) = \frac{f(x,y)}{f_X(x)} = \frac{1}{2x^2}(x-y), -x < y < x$$

$$(d) f_Y(y) = \int_{-y}^2 \frac{1}{8}x(x-y)dx, \text{ in } -2 \leq y \leq 0 \\ = \int_y^2 \frac{1}{8}x(x-y)dx, \text{ in } 0 \leq y \leq 2$$

$$\text{i.e., } f_Y(y) = \begin{cases} \frac{1}{3} - \frac{y}{4} + \frac{5}{48}y^3, & \text{in } -2 \leq y \leq 0 \\ \frac{1}{3} - \frac{y}{4} + \frac{1}{48}y^3, & \text{in } 0 \leq y \leq 2 \end{cases}$$

### Example 10

Train  $X$  arrives at a station at random in the time interval  $(0, T)$  and stops for ‘ $a$ ’ min. Train  $Y$  arrives independently in the same interval and stops for ‘ $b$ ’ min.

(a) Find the probability  $P_1$  that  $X$  will arrive before  $Y$ .

(b) Find the probability  $P_2$  that the two trains meet.

(c) Assuming that they meet, find the probability  $P_3$  that  $X$  arrived before  $Y$ .

Let the trains  $X$  and  $Y$  arrive at the station at time instances  $X$  and  $Y$  respectively.

Then the lengths of the intervals  $(0, X)$  and  $(0, Y)$ , namely  $X$  and  $Y$  are continuous RVs. Each of  $X$  and  $Y$  is uniformly distributed in  $(0, T)$  (since the times of arrival are equally likely) with pdf  $\frac{1}{T}$ .

Since the 2 trains arrive independently,  $X$  and  $Y$  are independent RVs.  
 ∴ the joint pdf of  $(X, Y)$  is given by

$$f(x, y) = f_X(x) \times f_Y(y) = \frac{1}{T^2}; 0 \leq x, y \leq T$$

The range space is the square defined by  $0 \leq x \leq T$  and  $0 \leq y \leq T$ .

$$(a) P_1 = P(X < Y) = \iint_{\substack{x < y \\ (\Delta OBC)}} f(x, y) dx dy$$

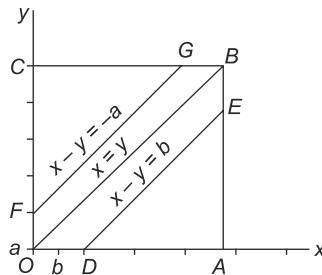


Fig. 2.9

$$= \frac{1}{T^2} \iint_{\Delta OBC} dx dy = \frac{1}{T^2} \times \text{Area of } \Delta OBC$$

$$= \frac{1}{2}$$

- (b) If train  $X$  arrives first, the two trains will meet if  $Y \leq X + a$ .  
 If train  $Y$  arrives first, the two trains will meet if  $X \leq Y + b$ .  
 ∴ for the two trains to meet,  $-a \leq X - Y \leq b$ .

$$\begin{aligned} \therefore P_2 &= P(-a \leq X - Y \leq b) = \iint_{\substack{-a \leq x - y \leq b \\ (\text{ODEBGFO})}} f(x, y) dx dy \\ &= \frac{1}{T^2} \times \text{Area of the figure ODEBGFO} \\ &= \frac{1}{T^2} \times (\text{Area of trapezium ODEB} + \text{that of OBGF}) \\ &= \frac{1}{T^2} \times \left[ \frac{1}{2} \times \frac{b}{\sqrt{2}} \{(T-b)\sqrt{2} + T\sqrt{2}\} + \frac{1}{2} \times \frac{a}{\sqrt{2}} \{(T-a)\sqrt{2} + T\sqrt{2}\} \right] \\ &= \frac{1}{2T^2} \{a(2T-a) + b(2T-b)\} \\ &= \frac{1}{2T^2} \{2(a+b)T - (a^2 - b^2)\} \end{aligned}$$

$$\begin{aligned}
 (c) \quad P_3 &= P\{X < Y \mid -a \leq X - Y \leq b\} \\
 &= \frac{P[X < Y \text{ and } -a \leq X - Y \leq b]}{P(-a \leq X - Y \leq b)} \\
 &= \frac{\frac{1}{T^2} \times \text{Area of trapezium } OBGF}{P_2} \\
 &= \frac{a(2T - a)}{2(a + b)T - (a^2 + b^2)}
 \end{aligned}$$

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**Example 11**

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Two trains arrive at a station at random between 7 a.m. and 7:30 a.m. One train stops for 5 min and the other for  $x$  min. For what value of  $x$ , will be probability that the 2 trains meet be equal to  $\frac{1}{3}$ ?

In the notation of the previous problem,

$$T = 30, a = 5, b = x \text{ and } P_2 = \frac{1}{3}$$

$$\therefore \frac{1}{2T^2} \{2(a+b)T - (a^2 + b^2)\} = \frac{1}{3}$$

$$\text{i.e., } \frac{1}{1800} \{60(x+5) - (x^2 + 25)\} = \frac{1}{3}$$

$$\text{i.e., } x^2 - 60x + 325 = 0$$

Solving,  $x = 53.98$  (or) 6.02

As  $x = 53.98$  is meaningless,  $x = 6$  min (nearly).

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**Example 12**

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The two-dimensional RV  $(X, Y)$  follows a bivariate normal distribution  $N(0, 0; \sigma_x, \sigma_y; r)$ . Find the marginal density function of  $X$  and the conditional density function of  $Y$ , given  $X$ .

The notation  $N(0, 0; \sigma_x, \sigma_y; r)$  refers to a bivariate normal distribution with mean of  $X = \text{mean of } Y = 0$ , variance of  $X = \sigma_x^2$ , variance of  $Y = \sigma_y^2$ , and the coefficient of correlation between  $X$  and  $Y = r$ .

The joint pdf of such a bivariate normal distribution is given by

$$f(x, y) = \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)} \left( \frac{x^2}{\sigma_x^2} - \frac{2rxy}{\sigma_x\sigma_y} + \frac{y^2}{\sigma_y^2} \right) \right\}$$

$-\infty < x, y < \infty$

The marginal density function of  $X$  is given by

$$\begin{aligned}
 f_X(x) &= \int_{-\infty}^{\infty} f(x, y) dy \\
 &= A \int_{-\infty}^{\infty} \exp\left\{-\frac{1}{2(1-r^2)}\left(\frac{y}{\sigma_y} - \frac{rx}{\sigma_x}\right)^2\right\} \times \\
 &\quad \exp\left(\frac{-x^2}{2\sigma_x^2}\right) dy, \text{ where } A = \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \\
 &= A \exp\left(\frac{-x^2}{2\sigma_x^2}\right) \times \sqrt{2}\sigma_y\sqrt{1-r^2} \int_{-\infty}^{\infty} e^{-t^2} dt, \\
 \text{by putting } &\frac{1}{\sqrt{2(1-r^2)}}\left(\frac{y}{\sigma_y} - \frac{rx}{\sigma_x}\right) = t, \\
 &= A \exp\left(\frac{-x^2}{2\sigma_x^2}\right) \sqrt{2} \cdot \sigma_y \sqrt{1-r^2} \left|\left(\frac{1}{2}\right)\right| \\
 &= \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left(\frac{-x^2}{2\sigma_x^2}\right) \sqrt{2} \cdot \sigma_y \sqrt{1-r^2} \sqrt{\pi} \\
 &= \frac{1}{\sigma_x\sqrt{2\pi}} \exp\left(\frac{-x^2}{2\sigma_x^2}\right), -\infty < x < \infty
 \end{aligned}$$

which is the density function of a normal distribution  $N(0, \sigma_x)$ .

The conditional density function of  $Y$  given  $X$  is given by  $f\left(\frac{y}{x}\right) = \frac{f(x, y)}{f_X(x)}$

$$\begin{aligned}
 \therefore f\left(\frac{y}{x}\right) &= \frac{\frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)}\left(\frac{x^2}{\sigma_x^2} - \frac{2rxy}{\sigma_x\sigma_y} + \frac{y^2}{\sigma_y^2}\right)\right\}}{\frac{1}{\sigma_x\sqrt{2\pi}} \exp\frac{-x^2}{2\sigma_x^2}} \\
 &= \frac{1}{\sqrt{2\pi}(\sigma_y\sqrt{1-r^2})} \exp\left\{-\frac{1}{2\sigma_y^2(1-r^2)}\left(y^2 - 2r\frac{\sigma_y}{\sigma_x}xy + \frac{r^2\sigma_y^2}{\sigma_x^2}\right)\right\}
 \end{aligned}$$

$$= \frac{1}{\sqrt{2\pi}(\sigma_y\sqrt{1-r^2})} \exp\left\{-\frac{1}{2(\sigma_y\sqrt{1-r^2})^2}\left(y - \frac{r\sigma_y x}{\sigma_x}\right)^2\right\}$$

which is the density function of a normal distribution

$$N\left(\frac{r\sigma_y}{\sigma_x}x, \sigma_y\sqrt{1-r^2}\right)$$

### Exercise 2(B)

#### Part-A (Short-answer Questions)

1. Define a two-dimensional RV. Give an example for the outcome of a random experiment, that is a two-dimensional RV.
2. Define the joint pmf of a two-dimensional discrete RV.
3. Define the joint pdf of a two-dimensional continuous RV.
4. Write down the joint pdf of a bivariate normal distribution.
5. Define the cdf of a two-dimensional RV and write down the formulas for, finding the cdf of  $(X, Y)$ , when  $(X, Y)$  is (i) a discrete RV, and (ii) a continuous RV.
6. State the properties of the cdf of a two-dimensional RV  $(X, Y)$ .
7. Define the marginal probability distributions of  $X$  and  $Y$ , when  $(X, Y)$  is a discrete RV.
8. Define the marginal probability density functions of  $X$  and  $Y$ , when  $(X, Y)$  is a continuous RV.
9. Define independence of 2 RVs  $X$  and  $Y$ , both in the discrete case and in the continuous case.
10. Define the conditional probability distributions of  $X$  and  $Y$ , given  $Y$  and  $X$  respectively, when  $(X, Y)$  is a discrete RV.
11. Define the conditional probability density functions of  $X$  and  $Y$  given  $Y$  and  $X$  respectively, when  $(X, Y)$  is a continuous RV.
12. Find the probability distribution of  $(X + Y)$  from the bivariate distribution of  $(X, Y)$  given below.

$X$	$Y$	
	1	2
1	0.1	0.2
2	0.3	0.4

13. Find the marginal distributions of  $X$  and  $Y$  from the bivariate distribution of  $(X, Y)$  given in Q.12.
14. Find the conditional distribution of  $X$ , when  $Y = 1$ , from the bivariate distribution of  $(X, Y)$  given in Q.12.
15. Find the value of  $k$ , if  $f(x, y) = k(1-x)(1-y)$ , for  $0 < x, y < 1$ , is to be a joint density function.

16. If  $f(x, y) = k(1 - x - y)$ ,  $0 < x, y < \frac{1}{2}$ , is a joint density function, find  $k$ .
17. If the joint pdf of  $(X, Y)$  is  $f(x, y) = \frac{1}{4}$ ,  $0 \leq x, y \leq 2$ , find  $P(X + Y \leq 1)$ .
18. If the joint pdf of  $(X, Y)$  is  $f(x, y) = 6e^{-2x-3y}$ ,  $x \geq 0, y \geq 0$ , find the marginal density of  $X$  and conditional density of  $Y$  given  $X$ .
19. The joint pdf of  $(X, Y)$  is given by  $f(x, y) = e^{-(x+y)}$ ,  $0 \leq x, y < \infty$ . Are  $X$  and  $Y$  independent? Why?
20. Define a random vector with an example.
21. Define the joint density and distribution functions of an  $n$ -dimensional RV. How are they related?

**Part-B**

22. If  $X$  denotes the number of aces and  $Y$  the number of queens obtained when 2 cards are drawn at random (without replacement) from a deck of cards, obtain the joint probability distribution of  $(X, Y)$ .
23. The joint probability function of two discrete RVs and  $X$  and  $Y$  is given by  $f(x, y) = c(2x + y)$ , where  $x$  and  $y$  can assume all integers such that  $0 \leq x \leq 2$  and  $0 \leq y \leq 3$ , and  $f(x, y) = 0$  otherwise. (i) Find the value of  $c$  and (ii) find  $P(X \geq 1, Y \leq 2)$ .

**Note**

$f(x, y)$  should not be mistaken as pdf, it is used instead of  $p(x_p, y_p)$ .

24. The joint probability distribution of a two-dimensional discrete RV  $(X, Y)$  is given below:

Y	X					
	0	1	2	3	4	5
0	0	0.01	0.03	0.05	0.07	0.09
1	0.01	0.02	0.04	0.05	0.06	0.08
2	0.01	0.03	0.05	0.05	0.05	0.06
3	0.01	0.02	0.04	0.06	0.06	0.05

- (i) Find  $P(X > Y)$  and  $P\{\max(X, Y) = 3\}$ , and  
(ii) Find the probability distribution of the RV  $Z = \min(X, Y)$ .
25. The input to a binary communication system, denoted by a RV  $X$ , takes one of two values 0 or 1 with probabilities  $3/4$  and  $1/4$  respectively. Because of errors caused by noise in the system, the output  $Y$  differs from the input occasionally. The behaviour of the communication system is modeled by the conditional probabilities given below:

$$P(Y = 1/X = 1) = 3/4 \text{ and } P(Y = 0/X = 0) = 7/8$$

Find (i)  $P(Y = 1)$ , (ii)  $P(Y = 0)$ , and (iii)  $P(X = 1/Y = 1)$ .

26. The following table represents the joint probability distribution of the discrete RV  $(X, Y)$ . Find all the marginal and conditional distributions.

Y	X		
	1	2	3
1	1/2	1/6	0
2	0	1/9	1/5
3	1/18	1/4	2/15

27. The joint distribution of  $X_1$  and  $X_2$  is given by  $f(x_1, x_2) = \frac{x_1 + x_2}{21}$ ,  $x_1 = 1, 2$  and  $3$ ;  $x_2 = 1$  and  $2$ . find the marginal distribution of  $X_1$  and  $X_2$ .
28. If the joint pdf of a two-dimensional RV  $(X, Y)$  is given by

$$f(x, y) = x^2 + \frac{xy}{3}; 0 < x < 1, 0 < y < 2 \\ = 0, \text{ elsewhere}$$

$$\text{find (i)} P\left(X > \frac{1}{2}\right), \text{(ii)} P(Y < X), \text{ and (iii)} P\left(Y < \frac{1}{2} / X < \frac{1}{2}\right).$$

29. If the joint pdf of a two-dimensional RV  $(X, Y)$  is given by

$$f(x, y) = k(6 - x - y); 0 < x < 2, 2 < y < 4 \\ = 0, \text{ elsewhere}$$

find (i) the value of  $k$ , (ii)  $P(X < 1, Y < 3)$ , (iii)  $P(X + Y < 3)$ , and (iv)  $P(X < 1 / Y < 3)$ .

30. The joint density function of the RVs  $X$  and  $Y$  is given by

$$f(x, y) = 8xy; 0 < x < 1, 0 < y < x \\ = 0, \text{ elsewhere},$$

$$\text{find } P\left(Y < \frac{1}{8} / X < \frac{1}{2}\right).$$

31. Given that the joint pdf of  $(X, Y)$  is

$$f(x, y) = e^{-y}; x > 0, y > x \\ = 0, \text{ otherwise}$$

Find (i)  $P(X > 1 / Y < 5)$ , and (ii) the marginal distributions of  $X$  and  $Y$ .

32. If the joint pdf of a two-dimensional RV  $(X, Y)$  is given by

$$f(x, y) = 2; 0 < x < 1, 0 < y < x \\ = 0, \text{ otherwise}$$

find the marginal density functions of  $X$  and  $Y$ .

33. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = k$ ,  $0 \leq x < y \leq 2$ , find  $k$  and also the marginal and conditional density functions.

34. The joint density function of a RV  $(X, Y)$  is  $f(x, y) = 8xy$ ,  $0 < x < 1$ ,  $0 < y < x$ , find the conditional density function  $f(y/x)$ .
35. The joint density function of a RV  $(X, Y)$  is given by  $f(x, y) = axy$ ,  $1 \leq x \leq 3$ ,  $2 \leq y \leq 4$ , and  $= 0$ , elsewhere.  
Find (i) the value of  $a$ , (ii) the marginal densities of  $X$  and  $Y$ , and (iii) the conditional densities of  $X$  and  $Y$ , given  $Y$  and  $X$  respectively.
36. Let  $X_1$  and  $X_2$  be two RVs with joint pdf given by  $f(x_1, x_2) = e^{-(x_1+x_2)}$ ;  $x_1, x_2 \geq 0$  and  $= 0$ , otherwise. Find the marginal densities of  $X_1$  and  $X_2$ . Are they independent? Also find  $P[X_1 \leq 1, X_2 \leq 1]$  and  $P(X_1 + X_2 \leq 1)$ .
37. The joint pdf of the RVs  $X$  and  $Y$  is given by  $p(x, y) = xe^{-x(y+1)}$  where  $0 \leq x, y < \infty$ . (i) Find  $p(x)$  and  $p(y)$ , and (ii) Are the RVs independent?
38. If the joint pdf of the RV  $(X, Y)$  is given by  $f(x, y) = k(x^3y + xy^3)$ ,  $0 \leq x \leq 2$ ,  $0 \leq y \leq 2$ , find (i) the value of  $k$ , (ii) the marginal densities of  $X$  and  $Y$ , and (iii) the conditional densities of  $X$  and  $Y$ .
39. If the joint pdf of  $(X, Y)$  is given by

$$f(x, y) = \frac{9(1+x+y)}{2(1+x)^4(1+y)^4} \quad x > 0, y > 0$$

- Find the marginal densities of  $X$  and  $Y$ . Are they independent?
40. Trains  $X$  and  $Y$  arrive at a station at random between 8 A.M. and 8.20 A.M. Train  $A$  stops for 4 min and train  $B$  stops for 5 min. Assuming that the trains arrive independently of each other, find the probability that (i)  $X$  will arrive before  $Y$ , (ii) the trains will meet, and (iii)  $X$  arrived before  $Y$ , assuming that they met.
41. If the two-dimensional RV  $(X, Y)$  follows a bivariate normal distribution  $N(0, 0; \sigma_x, \sigma_y; r)$ , find the marginal density function of  $Y$  and the conditional density function of  $X$ , given  $Y$ .
42. The two-dimensional RV  $(X, Y)$  has the joint density

$$\begin{aligned} f(x, y) &= 8xy, \quad 0 < x < y < 1 \\ &= 0, \text{ otherwise} \end{aligned}$$

- (i) Find  $P(X < 1/2 \cap Y < 1/4)$ ,  
(ii) Find the marginal and conditional distributions, and  
(iii) Are  $X$  and  $Y$  independent? Give reasons for your answer.

## **ANSWERS**

### **Exercise 2(A)**

6. Assign the values 0, 1, 2 to  $X$ , when the outcome consists of 2 tails, 1 tail and 1 head and 2 heads respectively. Then the required probability distribution of  $X$  is

$x:$	0	1	2
$p_x:$	$\frac{1}{4}$	$\frac{1}{2}$	$\frac{1}{4}$

9. **Example:**  $y = \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2}$ ,  $-\infty < x < \infty$ , is the equation of the normal curve, viz., the probability curve of the normal distribution.

11. $X:$	0	1	2	3	4	5
$p_x:$	$\frac{1}{32}$	$\frac{5}{32}$	$\frac{10}{32}$	$\frac{10}{32}$	$\frac{5}{32}$	$\frac{1}{32}$

12. Required probability = 
$$\frac{P\left\{\left(\frac{1}{2} < X < \frac{7}{2}\right) \cap (X > 1)\right\}}{P(X > 1)}$$

$$= \frac{P(X = 2 \text{ or } 3)}{P(X = 2, 3 \text{ or } 4)} = \frac{0.5}{0.6} = \frac{5}{6}$$

16.  $f(x) \geq 0; \int_{-1}^1 |x| dx = 2 \int_0^1 x dx = 1$

$\therefore f(x)$  can be the pdf of continuous RV.

17.  $\int_0^3 kx^2 dx = 1; 9k = 1$

$\therefore k = \frac{1}{9}$

18.  $P\{|X| > 1\} = 1 - P\{|X| < 1\} = 1 - \int_{-1}^1 \frac{1}{4} dx = \frac{1}{2}$

19.  $\int_0^\infty kx e^{-x} dx = 1;$

$$k[x(-e^{-x}) - 1 \times (e^{-x})]_0^\infty = 1$$

$\therefore k = 1$

20. Required probability = 
$$\frac{P(X > 1.5)}{P(X > 1)} = \frac{\left(\frac{x^2}{4}\right)_{1.5}^2}{\left(\frac{x^2}{4}\right)_1^2} = \frac{7}{12}$$

21.  $\sum P_x = 1 \therefore k = 0.1; E(X) = -0.8 - 0.1 + 0 + 0.3 = -0.6$

22. The probability distribution of  $X$  is

$X:$	1	2	3	4	5	6
$P_x:$	$\frac{1}{6}$	$\frac{1}{6}$	$\frac{1}{6}$	$\frac{1}{6}$	$\frac{1}{6}$	$\frac{1}{6}$

$$F(x) = P(X \leq x) = 0, \text{ if } x < 1; = \frac{1}{6}, \text{ if } 1 \leq x \leq 2;$$

$$= \frac{2}{6} \text{ if } 2 \leq x < 3; = \frac{3}{6}, \text{ if } 3 \leq x < 4; = \frac{4}{6}, \text{ if } 4 \leq x < 5;$$

$$= \frac{5}{6}, \text{ if } 5 \leq x < 6 \text{ and } = 1, \text{ if } 6 \leq x$$

23. When  $x < 0, F(x) < 0$ ; when  $0 \leq x < 1, F(x) = x^2$ , when  $1 \leq x, F(x) = 1$ .

24.  $f(x) = \frac{dF}{dx} = \lambda e^{-\lambda x}$ , when  $x > 0$  and  $= 0$  when  $x < 0$ .

25.  $P(X > 1/X < 3)$

$$= \frac{P(1 < X < 3)}{P(0 < X < 3)} = \frac{F(3) - F(1)}{F(3) - F(0)} = \frac{\frac{9}{16} - \frac{1}{16}}{\frac{9}{16} - 0} = \frac{8}{9}$$

26. For the binomial distribution  $B(n, p)$ , mean  $= np$  and variance  $= npq$ .

27. Mean = variance =  $\lambda$  for the Poisson distribution with parameter  $\lambda$ .

29.  $P(|X| < 1) = P(-1 < X < 1)$

$$= \int_0^1 e^{-x} dx = \frac{e-1}{e}$$

### Part-B

38.  $P(X = r) = 4C_r \left(\frac{1}{2}\right)^4, r = 0, 1, 2, 3, 4$

39.  $P(X = 0) = \frac{1}{64}, P(X = 1) = \frac{9}{64}, P(X = 2) = \frac{27}{64}, P(X = 3) = \frac{27}{64}$

40.  $P(X = 0) = q, P(X = 1) = p$ ; when  $x < 0, F(x) = 0$ ; when  $0 \leq x \leq 1, F(x) = q$ ; when  $1 \leq x, F(x) = q + p = 1$

41.  $P(X = r) = p(X = 14 - r) = \frac{(r-1)}{36}, r = 2, 3, 4, 5, 6, 7$

42.  $P(X = 0) = \frac{5}{16}, P(X = 1) = \frac{6}{16}, P(X = 2) = \frac{5}{16}; E(X) = 1; V(X) = \frac{5}{8}$



56. (i)  $\frac{1}{4}$  (ii)  $\frac{4}{9}$  (iii) 5

57.  $P = 1 - \frac{1}{e}; f(x) = \frac{k}{2} e^{-k|x|}$

58.  $f(x) = 1$  in  $0 \leq x < 1$  and = 0, elsewhere;  $\frac{1}{6}; \frac{1}{2}$ .

59.  $\frac{1}{5}$

60.  $\frac{2}{3}$

61. (i)  $\lambda$ ; (ii)  $\frac{\lambda^2 t}{1 + \lambda t}$

62. (i) 0.8849 (ii) 0.8819 (iii) 0.8593

## Exercise 2(B)

4.  $f_{XY}(x, y) = \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}}$   

$$\exp\left[-\frac{1}{2(1-r^2)}\left\{\frac{(x-\mu_x)^2}{\sigma_x^2} - \frac{2r(x-\mu_x)(y-\mu_y)}{\sigma_x\sigma_y} + \frac{(y-\mu_y)^2}{\sigma_y^2}\right\}\right],$$
  

$$-\infty < x, y < \infty$$

12.  $(X + Y): \quad 2 \quad 3 \quad 4$   
 $p: \quad 0.1 \quad 0.5 \quad 0.4$

13.  $X: \quad 1 \quad 2 \quad \text{and} \quad Y: \quad 1 \quad 2$   
 $p_x: \quad 0.3 \quad 0.7 \quad p_y: \quad 0.4 \quad 0.6$

14.  $X: \quad 1 \quad 2$   
 $P_{X/Y=1}: \quad 0.25 \quad 0.75$

15.  $k \int_0^1 \int_0^1 (1-x)(1-y) dx dy = 1;$

$$k \left\{ \frac{(1-x)^2}{-2} \right\}_0^1 \left\{ \frac{(1-y)^2}{-2} \right\}_0^1 = 1$$

i.e.,  $\frac{k}{4} = 1$

$\therefore k = 4$

16.  $\int_0^{1/2} \int_0^{1/2} k(1-x-y) dx dy = 1$

$$k \int_0^{1/2} \left( x - \frac{x^2}{2} - yx \right)_0^{\frac{1}{2}} dy = 1;$$

$$k \left( \frac{3}{8y} - \frac{y^2}{4} \right)_0^{\frac{1}{2}} = 1$$

$$\therefore k = 8$$

17.  $P = \int_0^1 \int_0^{1-y} \frac{1}{4} dx dy = \frac{1}{4} \int_0^1 (1-y) dy = \frac{1}{8}$

18.  $f_X(x) = 6 \int_0^{\infty} e^{-2x} e^{-3y} dy = 2e^{-2x}; x \geq 0$

$$f_{Y/X}(y) = \frac{f(x,y)}{f_X(x)} = 3e^{-3y}; y \geq 0$$

19.  $f_X(x) = e^{-x}$  and  $f_Y(y) = e^{-y}$

$$f(x, y) = f_X(x) f_Y(y)$$

$\therefore X$  and  $Y$  are independent

22.

$X \backslash Y$	0	1	2
0	0.71	0.13	0.01
1	0.13	0.01	0
2	0.01	0	0

23.  $c = \frac{1}{42}; P = \frac{4}{7}$

24. (i) 0.75, 0.21      (ii)

Z	0	1	2	3
P	0.28	0.30	0.25	0.17

25. (i)  $\frac{9}{32}$       (ii)  $\frac{23}{32}$       (iii)  $\frac{2}{3}$

{ $i, p_{i*}$ }		{ $j, p_{*j}$ }		CPD of $X/Y = 1$		CPD of $X/Y = 2$	
$x = i$	$p_{i*}$	$y = j$	$p_{*j}$	$x = i$	$p_{i1}/p_{*1}$	$x = i$	$p_{i2}/p_{*2}$
1	5/36	1	1/4	1	1/3	1	0
2	19/36	2	14/45	2	2/3	2	5/14
3	1/3	3	79/180	3	0	3	9/14

CPD of $X/Y = 3$		CPD of $Y/X = 1$		CPD of $Y/X = 2$		CPD of $Y/X = 3$	
$x = i$	$p_{i3}/p_{*3}$	$x = j$	$p_{1j}/p_{1*}$	$x = j$	$p_{2j}/p_{2*}$	$x = j$	$p_{3j}/p_{3*}$
1	10/79	1	3/5	1	6/19	1	0
2	45/79	2	0	2	4/19	2	3/5
3	24/79	3	2/5	3	9/19	3	2/5

27.

M.D. of $X_1$		M.D. of $X_2$	
$X_1 = i$	$pi^*$	$X_2 = j$	$p^*j$
1	5/21	1	9/21
2	7/21	2	12/21
3	9/21		

28. (i)  $\frac{5}{6}$       (ii)  $\frac{7}{24}$       (iii)  $\frac{5}{32}$

29. (i)  $1/8$       (ii)  $3/8$       (iii)  $5/24$       (iv)  $3/5$

30.  $\frac{31}{256}$

31. (i)  $\frac{e^4 - 5}{e^5 - 6}$  (ii)  $f_X(x) = e^{-x}, x > 0; f_Y(y) = y e^{-y}, y > 0.$

32.  $f_X(x) = 2x$  in  $0 < x < 1; f_Y(y) = 2(1 - y)$  in  $0 < y < 1.$

33.  $k = 1/2; f_X(x) = 1/2 (2 - x), 0 \leq x \leq 2; f_Y(y) = (1/2) y, 0 \leq y \leq 2; f(x/y) = 1/y; 0 < x < y; f(y/x) = \frac{1}{2-x}, x < y < 2.$

34.  $f(y/x) = 2y/x^2, 0 < y < x$

35. (i)  $a = 1/24$

(ii)  $f_X(x) = x/4, 1 \leq x \leq 3; f_Y(y) = y/6, 2 \leq y \leq 4$

(iii)  $f(x/y) = x/4, 1 < x < 3; f(y/x) = y/6, 2 < y < 4$

36.  $f_{X_1}(x_1) = e^{-x_1}; f_{X_2}(x_2) = e^{-x_2}; X_1$  and  $X_2$  are independent;  $(1 - e^{-1})^2; (1 - 2e^{-1})$

37.  $f_X(x) = e^{-x}$ ,  $0 < x < \infty$ ;  $f_Y(y) = (y + 1)^{-2}$ ,  $0 < y < \infty$ ;  $X$  and  $Y$  are not independent.
38.  $k = 1/16$ ;  $f_X(x) = 1/8 (x^3 + 2x)$ ,  $0 \leq x \leq 2$ ;  $f_Y(y) = 1/8 (y^3 + 2y)$ ,  $0 \leq y \leq 2$ ;  
 $f(x/y) = x/2 \frac{(x^2 + y^2)}{(y^2 + 2)}$ ,  $0 < x < 2$ ;  $f(y/x) = y/2 \frac{(x^2 + y^2)}{(x^2 + 2)}$ ,  $0 < y < 2$ .
39.  $f_X(x) = (3/4) \frac{2x+3}{(1+x)^4}$   $x > 0$  and  $f_Y(y) = (3/4) \frac{2x+3}{(1+y)^4}$   $y > 0$ ;  
Not independent.
40. (i) 1/2   (ii) 0.3988   (iii) 0.4514
41.  $f_Y(y) = \frac{1}{\sigma_y \sqrt{2\pi}} \exp\left(-\frac{y^2}{2\sigma_y^2}\right)$ ,  $-\infty < y < \infty$   
 $f(x/y) = N\left(\frac{r\sigma_x}{\sigma_y} \cdot y, \sigma_x \sqrt{1-r^2}\right)$
42. (i) 1/256  
(ii)  $f_X(x) = 4x(1-x^2)$  in  $(0, 1)$ ;  $f_Y(y) = 4y^3$  in  $(0, 1)$ ;  $f(y/x)$   
 $= \left| \frac{2y}{1-x^2} \right| f(x/y) = \frac{2x}{y^2}$   
(iii)  $X$  and  $Y$  are not independent.



# Chapter 3

## Functions of Random Variables

In the analysis of electrical systems, we will be often interested in finding the properties of a signal after it has been subjected to certain processing operations by the system, such as integration, weighted averaging, etc. These signal processing operations may be viewed as transformations of a set of input variables to a set of output variables. If the input is a set of random variables (RVs), then the output will also be a set of RVs. In this chapter, we deal with techniques for obtaining the probability law (distribution) for the set of output RVs when the probability law for the set of input RVs and the nature of transformation are known.

### Function of One Random Variable

Let  $X$  be an RV with the associated sample space  $S_x$  and a known probability distribution. Let  $g$  be a scalar function that maps each  $x \in S_x$  into  $y = g(x)$ . The expression  $Y = g(X)$  defines a new RV  $Y$ . For a given outcome,  $X(s)$  is a number  $x$  and  $g[X(s)]$  is another number specified by  $g(x)$ . This number is the value of the RV  $Y$ , i.e.,  $Y(s) = y = g(x)$ . The sample space  $S_y$  of  $Y$  is the set

$$S_y = \{y = g(x) : x \in S_x\}$$

#### How to find $f_Y(y)$ , when $f_X(x)$ is known

Let us now derive a procedure to find  $f_Y(y)$ , the pdf of  $Y$ , when  $y = g(X)$ , where  $X$  is a continuous RV with pdf  $f_X(x)$  and  $g(x)$  is a strictly monotonic function of  $x$ .

**Case (i):**  $g(x)$  is a strictly increasing function of  $x$ .

$$\begin{aligned}f_Y(y) &= P(Y \leq y), \text{ where } F_Y(y) \text{ is the cdf of } Y \\&= P[g(X) \leq y] \\&= P[X \leq g^{-1}(y)] \\&= F_X(g^{-1}(y))\end{aligned}$$

Differentiating both sides with respect to  $y$ ,

$$f_Y(y) = f_X(x) \frac{dx}{dy}, \text{ where } x = g^{-1}(y) \quad (1)$$

**Case (ii):**  $g(x)$  is a strictly decreasing function of  $x$ .

$$\begin{aligned} f_Y(y) &= P(Y \leq y) \\ &= P[g(X) \leq y] \\ &= P[X \geq g^{-1}(y)] \\ &= 1 - P[X \leq g^{-1}(y)] \\ &= 1 - F_X[g^{-1}(y)] \end{aligned}$$

$$\therefore f_Y(y) = -f_X(x) \frac{dx}{dy} \quad (2)$$

Combining (1) and (2), we get

$$f_Y(y) = f_X(x) \left| \frac{dx}{dy} \right|$$

$$\text{i.e., } f_Y(y) = \frac{f_X(x)}{|g'(x)|}$$

**Note** The above formula for  $f_Y(y)$  can be used only when  $x = g^{-1}(y)$  is single valued.

When  $x = g^{-1}(y)$  takes finitely many values  $x_1, x_2, \dots, x_n$ , i.e., when the roots of the equation  $y = g(x)$  are  $x_1, x_2, \dots, x_n$ , the following extended formula should be used for finding  $f_Y(y)$ :

$$f_Y(y) = \frac{f_X(x_1)}{|g'(x_1)|} + \frac{f_X(x_2)}{|g'(x_2)|} + \dots + \frac{f_X(x_n)}{|g'(x_n)|} \text{ or}$$

$$f_Y(y) = f_X(x_1) \left| \frac{dx_1}{dy} \right| + f_X(x_2) \left| \frac{dx_2}{dy} \right| + \dots + f_X(x_n) \left| \frac{dx_n}{dy} \right|$$

## One Function of Two Random Variables

If a RV  $Z$  is defined as  $Z = g(X, Y)$  where  $X$  and  $Y$  are RVs, we proceed to find  $f_Z(z)$  in the following way.

If  $z$  is a given number, we can find a region  $D_Z$  in the  $xy$ -plane such that all points in  $D_Z$  satisfy the condition  $g(x, y) \leq z$ .

$$\text{i.e., } (Z \leq z) = [g(X, Y) \leq z] = [(X, Y) \in D_Z]$$

$$\text{Now } F_Z(z) = P(Z \leq z) = P[(X, Y) \in D_Z] = \iint_{D_z} f(x, y) dx dy$$

where  $f(x, y)$  is the joint pdf of  $(X, Y)$ . Thus, to find  $F_Z(z)$  it is sufficient to find the region  $D_Z$  for every  $Z$  and to evaluate the above integral.  $f_Z(z)$  is then found out as usual.

### Theorem 1

If two RVs are independent, then the density function of their sum is given by the convolution of their density functions.

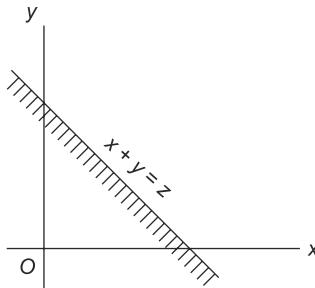


Fig. 3.1

### Proof

Let the joint pdf of  $(X, Y)$  be  $f(x, y)$

$$\text{Let } Z = X + Y$$

$$\begin{aligned} F_Z(z) &= P(X + Y \leq z) \\ &= \iint_{(x+y \leq z)} f(x, y) dx dy \\ &= \int_{-\infty}^{\infty} \int_{-\infty}^{z-y} f(x, y) dx dy \end{aligned}$$

Differentiating both sides with respect to  $z$  (note that the upper limit for the inner integral is a function of  $z$ ),

$$f_Z(z) = \int_{-\infty}^{\infty} f(z - y, y) dy \quad (1)$$

Since  $X, Y$  are independent RVs,

$$\begin{aligned} f(x, y) &= f_X(x) f_Y(y) \\ \therefore f(z - y, y) &= f_X(z - y) f_Y(y) \end{aligned} \quad (2)$$

Using (2) in (1), we get

$$f_Z(z) = \int_{-\infty}^{\infty} f_X(z - y) f_Y(y) dy, \text{ which is the convolution of } f_X(x) \text{ and } f_Y(y).$$

### Corollary

If  $f_X(x) = 0$ , for  $x < 0$ , and  $f_Y(y) = 0$ , for  $y < 0$ , then  $f_X(z - y) f_Y(y) \neq 0$ , only when  $0 < y < z$ .

$$\therefore f_Z(z) = \int_0^z f_X(z-y) f_Y(y) dy, z > 0$$

### Theorem 2

If two RVs  $X$  and  $Y$  are independent, find the pdf of  $Z = XY$  in terms of the density functions of  $X$  and  $Y$ .

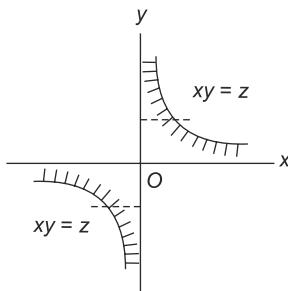


Fig. 3.2

Let the joint pdf of  $(X, Y)$  be  $f(x, y)$

$$f_Z(z) = \iint_{xy \leq z} f(x, y) dx dy$$

#### Note

$xy = z$  is a rectangular hyperbola as shown in the figure.

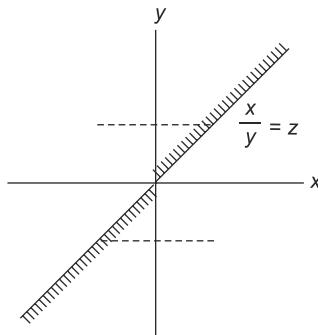
$$f_Z(z) = \int_{-\infty}^0 \int_{z/y}^0 f(x, y) dx dy + \int_0^{\infty} \int_0^{z/y} f(x, y) dx dy$$

Differentiating both sides with respect to  $z$ ,

$$\begin{aligned} f_Z(z) &= \int_{-\infty}^0 -\frac{1}{y} f\left(\frac{z}{y}, y\right) dy + \int_0^{\infty} \frac{1}{y} f\left(\frac{z}{y}, y\right) dy \\ &= \int_{-\infty}^{\infty} \left| \frac{1}{y} \right| f\left(\frac{z}{y}, y\right) dy \\ &= \int_{-\infty}^{\infty} \left| \frac{1}{y} \right| f_X\left(\frac{z}{y}\right) f_Y(y) dy \quad (\text{Since } X \text{ and } Y \text{ are independent}) \end{aligned}$$

### Theorem 3

If two RVs  $X$  and  $Y$  are independent, find the pdf of  $Z = \frac{X}{Y}$  in terms of the density functions of  $X$  and  $Y$ .

**Fig. 3.3**

Let the joint pdf of  $(X, Y)$  be  $f(x, y)$ .

$$F_Z(z) = \iint_{\frac{x}{y} \leq z} f(x, y) dx dy$$

i.e.,  $F_Z(z) = \int_{-\infty}^0 \int_{yz}^{\infty} f(x, y) dx dy + \int_0^{\infty} \int_{-\infty}^{yz} f(x, y) dx dy$

Differentiating both sides with respect to  $z$ ,

$$\begin{aligned} f_Z(z) &= \int_{-\infty}^0 -yf(yz, y) dy + \int_0^{\infty} yf(yz, y) dy \\ &= \int_{-\infty}^{\infty} |y| f(yz, y) dy \\ &= \int_{-\infty}^{\infty} |y| f_X(yz) f_Y(y) dy \quad (\text{since } X \text{ and } Y \text{ are independent}) \end{aligned}$$

## Two Functions of Two Random Variables

### Theorem

If  $(X, Y)$  is a two-dimensional RV with joint pdf  $f_{XY}(x, y)$  and if  $Z = g(X, Y)$  and  $W = h(X, Y)$  are two other RVs, then the joint pdf of  $(Z, W)$  is given by

$$f_{ZW}(z, w) = |J| f_{XY}(x, y), \text{ where } J = \frac{\partial(x, y)}{\partial(z, w)}$$

is called the Jacobian of the transformation and is given by

$$J = \begin{vmatrix} \frac{\partial x}{\partial z} & \frac{\partial x}{\partial w} \\ \frac{\partial y}{\partial z} & \frac{\partial y}{\partial w} \end{vmatrix}$$

**Note** This theorem holds good, only if the equations  $z = g(x, y)$  and  $w = h(x, y)$  when solved, give unique values of  $x$  and  $y$  in terms of  $z$  and  $w$ .

### An alternative method to find the pdf of $Z = g(X, Y)$

Introduce a second RV  $W = h(X, Y)$  and obtain the joint pdf of  $(Z, W)$ , as suggested in the above theorem. Let it be  $f_{ZW}(z, w)$ . The required pdf of  $Z$  is then obtained as the marginal pdf, i.e.,  $f_Z(z)$  is obtained by simply integrating  $f_{ZW}(z, w)$  with respect to  $w$ .

$$\text{i.e., } f_Z(z) = \int_{-\infty}^{\infty} f_{ZW}(z, w) dw$$

### Worked Example 3

#### Example 1

Find the distribution function of the RV  $Y = g(X)$ , in terms of the distribution function of  $X$ , if it is given that

$$g(x) = \begin{cases} x - c & \text{for } x > c \\ 0 & \text{for } |x| \leq c \\ x + c & \text{for } x < -c \end{cases}$$

$$\begin{aligned} \text{If } y < 0, \quad F_Y(y) &= P(Y \leq y) \\ &= P(X + c \leq y) \\ &= P(X \leq y - c) \\ &= F_X(y - c) \end{aligned}$$

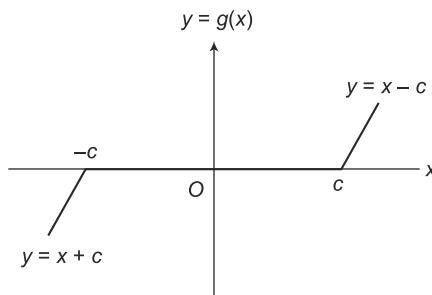


Fig. 3.4

$$\begin{aligned} \text{If } y \geq 0, \quad F_Y(y) &= P(X - c \leq y) \\ &= F_X(y + c) \end{aligned}$$

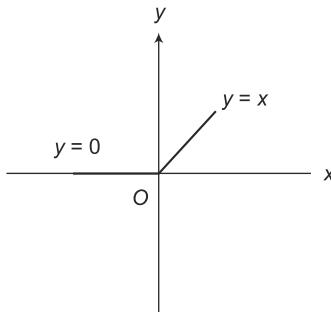
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**Example 2**

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The random variable  $Y$  is defined by  $Y = \frac{1}{2}(X + |X|)$ , where  $X$  is another RV.

Determine the density and distribution function of  $Y$  in terms of those of  $X$ .



**Fig. 3.5**

When  $X \geq 0$ ,  $Y = X$

When  $X < 0$ ,  $Y = 0$

$$\text{If } y < 0, \quad F_Y(y) = P(Y \leq y) = 0 \quad (\text{since there is no } X, \text{ for which } Y \leq y)$$

$$\begin{aligned} \text{If } y \geq 0, \quad F_Y(y) &= P(Y \leq y) \\ &= P(X \leq y/X \geq 0) \\ &= P(0 \leq X \leq y)/P(X \geq 0) \\ &= \frac{F_X(y) - F_X(0)}{1 - F_X(0)} \end{aligned}$$

$$\therefore \text{when } y < 0, f_Y(y) = 0$$

and

$$\text{when } y \geq 0, f_Y(y) = f_X(y)/[1 - F_X(0)]$$

---

**Example 3**

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- (a) Find the density function of  $Y = aX + b$  in terms of the density function of  $X$ .

- (b) Let  $X$  be a continuous RV with pdf

$$\begin{aligned} f(x) &= \frac{x}{12}, \text{ in } 1 < x < 5 \\ &= 0, \text{ elsewhere} \end{aligned}$$

find the probability density function of  $Y = 2X - 3$ .

- (a) (i) Let  $a > 0$

$$F_Y(y) = P(Y \leq y) = P(aX + b \leq y)$$

$$\begin{aligned}
 &= P\left(X \leq \frac{y-b}{a}\right) \quad (\text{since } a > 0) \\
 &= F_x\left(\frac{y-b}{a}\right)
 \end{aligned} \tag{1}$$

(ii) Let  $a < 0$

$$\begin{aligned}
 F_Y(y) &= P(Y \leq y) = P(aX + b \leq y) \\
 &= P(aX \leq y - b) \\
 &= P\left(X \geq \frac{y-b}{a}\right) \quad (\text{since } a < 0) \\
 &= 1 - F_x\left(\frac{y-b}{a}\right)
 \end{aligned} \tag{2}$$

$$\text{From (1), } f_Y(y) = \frac{1}{a} f_x\left(\frac{y-b}{a}\right) \tag{3}$$

$$\text{From (2), } f_Y(y) = -\frac{1}{a} f_x\left(\frac{y-b}{a}\right) \tag{4}$$

Combining (3) and (4),

$$f_Y(y) = \frac{1}{|a|} f_x\left(\frac{y-b}{a}\right)$$

(b)  $y = 2x - 3$ , since  $Y = 2X - 3$

$$\therefore x = \frac{1}{2}(y+3), \text{ i.e., } x \text{ is a single valued function of } y$$

$$\begin{aligned}
 \therefore f_Y(y) &= f_X(x) \left| \frac{dx}{dy} \right| \\
 &= \frac{x}{12} \times \frac{1}{2} \\
 &= \frac{1}{48}(y+3), \text{ in } -1 < y < 7
 \end{aligned}$$

**Note** The range of  $y$  is obtained from that of  $x$  (given in the problem) using the relation between  $x$  and  $y$ .

---

### Example 4

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If  $X$  is a continuous RV with some distribution defined over  $(0, 1)$  such that

$$P(X \leq 0.29) = 0.75, \text{ determine } k \text{ so that}$$

$$P(Y \leq k) = 0.25, \text{ where } Y = 1 - X$$

$$P(Y \leq k) = 0.25$$

$$\text{i.e., } P(1 - X \leq k) = 0.25$$

$$\text{i.e., } P(X \geq 1 - k) = 0.25$$

$$\therefore P(X \leq 1 - k) = 0.75 \quad (1)$$

$$\text{But it is given that } P(X \leq 0.29) = 0.75 \quad (2)$$

Comparing (1) and (2),  $k = 0.71$

### Example 5

If  $Y = X^2$ , where  $X$  is a Gaussian random variable with zero mean and variance  $\sigma^2$ , find the pdf of the random variable  $Y$ .

$$\begin{aligned} F_Y(y) &= P(Y \leq y) = P(X^2 \leq y) \\ &= P(-\sqrt{y} \leq X \leq \sqrt{y}), \text{ if } y \geq 0 \\ &= F_X(\sqrt{y}) - F_X(-\sqrt{y}) \end{aligned} \quad (1)$$

and  $F_Y(y) = 0$ , if  $y < 0$  [since  $X^2 = y$  has no roots, when  $y < 0$ ]

Differentiating (1) with respect to  $y$ ,

$$\begin{aligned} f_Y(y) &= \frac{1}{2\sqrt{y}} \left\{ f_X(\sqrt{y}) + f_X(-\sqrt{y}) \right\}, \text{ if } y \geq 0 \\ &= 0, \text{ if } y < 0 \end{aligned} \quad (2)$$

It is given that  $X$  follows  $N(0, \sigma)$ .

$$\therefore f_X(x) = \frac{1}{\sigma\sqrt{2\pi}} e^{-x^2/2\sigma^2} \quad -\infty < x < \infty$$

Using this value in (2), we get

$$f_Y(y) = \frac{1}{\sigma\sqrt{2\pi y}} e^{-y/2\sigma^2} \quad y > 0$$

### Example 6

If the continuous RV  $X$  has pdf  $f_X(x) = \frac{2}{9}(x+1)$ , in  $-1 < x < 2$  and = 0, elsewhere, find the pdf of  $Y = X^2$ .

The transformation function  $y = x^2$  is not monotonic in  $(-1, 2)$ . So we divide the interval into two parts.

i.e.,  $(-1, 1)$  and  $(1, 2)$

Since  $(-1, 1)$  is a symmetric interval,  $f_Y(y)$  is found out by using the formula (2) of the previous problem.

$\therefore$  when  $-1 < x < 1$ , i.e.,  $0 < y < 1$

$$\begin{aligned} f_Y(y) &= \frac{1}{2\sqrt{y}} \left\{ \frac{2}{9}(1 + \sqrt{y}) + \frac{2}{9}(1 - \sqrt{y}) \right\} \\ &= \frac{2}{9\sqrt{y}} \end{aligned}$$

When  $1 < x < 2$ , i.e.,  $1 < y < 4$ ,  $y = x^2$  is strictly increasing.

$$\begin{aligned}\therefore f_Y(y) &= f_X(x) \left| \frac{dx}{dy} \right| \\ &= \frac{2}{9} (x+1) \times \frac{1}{2x} \\ &= \frac{1}{9} \left( 1 + \frac{1}{\sqrt{y}} \right)\end{aligned}$$

### Example 7

---

According to the Maxwell–Boltzmann law of theoretical physics, the pdf of  $V$ , the velocity of a gas molecule is given by

$$f_V(v) = \begin{cases} kv^2 e^{-av^2}, & v > 0 \\ 0, & \text{elsewhere} \end{cases}$$

where  $a$  is a constant depending on its mass and the absolute temperature and  $k$  is an appropriate constant. Show that the kinetic energy  $Y = \frac{1}{2} mV^2$  is a random variable having Gamma distribution.

By the property of pdf,

$$k \int_0^{\infty} v^2 e^{-av^2} dv = 1$$

i.e.,  $\frac{k}{2a\sqrt{a}} \int_0^{\infty} t^{3/2-1} e^{-t} dt = 1$ , by putting  $t = av^2$

i.e.,  $\frac{k}{2a\sqrt{a}} \left[ \left( \frac{3}{2} \right) \right] = 1$

i.e.,  $\frac{k}{2a\sqrt{a}} \times \frac{1}{2} \left[ \left( \frac{1}{2} \right) \right] = 1$

$\therefore k = \frac{4a\sqrt{a}}{\sqrt{\pi}}$ , since  $\left[ \left( \frac{1}{2} \right) \right] = \sqrt{\pi}$

$$Y = \frac{m}{2} V^2$$

$\therefore v = \pm \sqrt{\frac{2y}{m}}$

since  $v > 0$ ,  $v = \sqrt{\frac{2y}{m}}$  is the only admissible value.

$$\begin{aligned}
 \text{Now, } f_Y(y) &= f_X(v) \left| \frac{dv}{dy} \right| \\
 &= \frac{4a\sqrt{a}}{\sqrt{\pi}} \times \frac{2y}{m} e^{-2ay/m} \frac{1}{\sqrt{2my}} \\
 &= \frac{4\sqrt{2} a \sqrt{a}}{m \sqrt{m} \sqrt{\pi}} \times y^{1/2} e^{-2ay/m} \\
 &= \frac{(2a/m)^{3/2}}{(3/2)} \times y^{3/2-1} e^{-(2a/m)y}, \quad y > 0
 \end{aligned}$$

which is a 2-parameter Gamma distribution or Erlang distribution.

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### Example 8

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Given the RV  $X$  with density function

$$f(x) = \begin{cases} 2x, & 0 < x < 1 \\ 0, & \text{elsewhere} \end{cases}$$

find the pdf of  $Y = 8X^3$ .

Since  $y = 8x^3$  is a strictly increasing function in  $(0, 1)$ ,

$$\begin{aligned}
 f_Y(y) &= f_X(x) \left| \frac{dx}{dy} \right|, \text{ where } x = \frac{1}{2} y^{1/3} \\
 &= y^{1/3} \times \frac{1}{6} y^{-2/3} \\
 &= \frac{1}{6} y^{-1/3} \quad 0 < y < 8
 \end{aligned}$$

---

### Example 9

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If  $X$  is a Gaussian random variable with mean zero and variance  $\sigma^2$ , find the pdf of  $Y = |X|$ .

$$\begin{aligned}
 F_Y(y) &= P(Y \leq y) = P(|X| \leq y) \\
 &= P\{-y \leq X \leq y\} \\
 &= F_X(y) - F_X(-y)
 \end{aligned}$$

Differentiating both sides with respect to  $y$ ,

$$f_Y(y) = f_X(y) + f_X(-y) \quad y > 0 \tag{1}$$

Now  $X$  follows  $N(0, \sigma)$

$$\therefore f_X(x) = \frac{1}{\sigma\sqrt{2\pi}} e^{-x^2/2\sigma^2}, \quad -\infty < x < \infty \tag{2}$$

Using (2) in (1), we get

$$f_Y(y) = \frac{2}{\sigma\sqrt{2\pi}} e^{-y^2/2\sigma^2}, y > 0$$

---

### Example 10

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- (a) If  $X$  is a normal RV with mean zero and variance  $\sigma^2$ , find the pdf of  $Y = e^X$ .
- (b) If  $X$  has an exponential distribution with parameter  $\alpha$ , find the pdf of  $Y = \log X$ .

$$(a) f_Y(y) = f_X(x) \left| \frac{dx}{dy} \right| \quad (\text{since } y = e^x \text{ is monotonic function})$$

$$= \frac{1}{y} f_X(\log y)$$

$$f_Y(y) = \frac{1}{\sigma y \sqrt{2\pi}} e^{-(\log y - \mu)^2/2\sigma^2}, y > 0$$

which is the pdf of a lognormal distribution.

$$(b) f_Y(y) = f_X(x) \left| \frac{dx}{dy} \right|$$

$$= e^y \alpha e^{-\alpha ey}, -\infty < y < \infty$$

---

### Example 11

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- (a) If  $X$  has an exponential distribution with parameter 1, find the pdf of  $Y = \sqrt{X}$ .
- (b) If  $X$  has a Cauchy's distribution with parameter  $\alpha$ , prove that  $Y = 1/X$  has also a Cauchy's distribution with parameter  $1/\alpha$ .

$$(a) f_Y(y) = f_X(x) \left| \frac{dx}{dy} \right| \quad (\text{since } y = \sqrt{x} \text{ is m.i.})$$

$$= 2y e^{-y^2}, \quad y > 0$$

$$(b) f_Y(y) = f_X(x) \left| \frac{dx}{dy} \right| \quad \left( \text{sine } y = \frac{1}{x} \text{ is m.d.} \right)$$

$$= \frac{1}{y^2} \times \frac{\frac{\alpha}{\pi}}{\left( \frac{1}{y} \right)^2 + \alpha^2}$$

$$\left( \text{since pdf of Cauchy's distribution} = \frac{\alpha/\pi}{x^2 + \alpha^2} \right)$$

$$= \frac{1}{\pi \alpha} \frac{1}{y^2 + \frac{1}{\alpha^2}}, \quad -\infty < x < \infty$$

which is a Cauchy's distribution with parameter  $1/\alpha$ .

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### Example 12

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- (a) If the RV  $X$  is uniformly distributed in  $(-\pi, \pi)$  find the pdf of  $Y = a \sin(X + \alpha)$ , where  $a > 0$  and  $\alpha$  are constants.
- (b) The horizontal range of a projectile is given by  $R = \frac{v^2}{g} \sin 2\theta$ . If  $\theta$  is uniformly distributed in  $(0, \pi/2)$  and  $\frac{v^2}{g}$  is a constant, find the pdf of  $R$ .
- (a) If  $|y| > a$ , no solution exists for  $x$ .

$$\therefore f_Y(y) = 0$$

If  $|y| < a$ , there exist only two values for  $x$  in  $(-\pi, \pi)$ .

$$\text{Let them be } x_r = \sin^{-1} \left( \frac{y}{a} \right) - \alpha; \quad r = 1, 2$$

$$\therefore \frac{dx_r}{dy} = \frac{1}{\sqrt{a^2 - y^2}}$$

$$\begin{aligned} \text{Now, } f_Y(y) &= f_X(x_1) \left| \frac{dx_1}{dy} \right| + f_X(x_2) \left| \frac{dx_2}{dy} \right| \\ &= \frac{1}{2\pi} \times \frac{1}{\sqrt{a^2 - y^2}} + \frac{1}{2\pi} \times \frac{1}{\sqrt{a^2 - y^2}} \\ &\quad \left( \text{since the pdf of } X = \frac{1}{2\pi}, \text{ as it is } U(-\pi, \pi) \right) \end{aligned}$$

$$\therefore f_Y(y) = \frac{1}{\pi \sqrt{a^2 - y^2}} \mid y \mid < a$$

- (b)  $\theta$  is uniform in  $(0, \pi/2)$

$\therefore X = 2\theta$  is uniform in  $(0, \pi)$  with pdf  $1/\pi$

$$R = \frac{v^2}{g} \sin X = a \sin X, \text{ say } [R \text{ and } a > 0]$$

Therefore, when  $r < a$ , there exist only two values for  $x$  in  $(0, \pi)$ , given

$$\text{by } x_i = \sin^{-1} \frac{r}{a}, \quad i = 1, 2.$$

$$\therefore \frac{dx_i}{dr} = \frac{1}{\sqrt{a^2 - r^2}}$$

$$\begin{aligned}\text{Now } f_R(r) &= f_X(x_1) \left| \frac{dx_1}{dr} \right| + f_X(x_2) \left| \frac{dx_2}{dr} \right| \\ &= \frac{2}{\pi \sqrt{a^2 - r^2}}, 0 < r < a \\ &= 0, r > a\end{aligned}$$

---

### Example 13

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- (a) If  $X$  is uniformly distributed in  $(-\pi/2, \pi/2)$  find the pdf of  $Y = \tan X$ .
- (b) If  $X$  has the Cauchy's distribution with parameter 1, find the pdf of  $Y = \tan^{-1} X$ .
- (a)  $X$  is  $U(-\pi/2, \pi/2)$

$$\therefore f_X(x) = \frac{1}{\pi}$$

$y = \tan x$

Therefore,  $x = \tan^{-1} y$ , which is single valued in  $(-\pi/2, \pi/2)$ , i.e., for a given value of  $y$ , there exists only one value of  $\tan^{-1} y$  in  $(-\pi/2, \pi/2)$ .

$$\begin{aligned}f_Y(y) &= f_X(x) \left| \frac{dx}{dy} \right| \\ &= \frac{1}{\pi} \frac{1}{1+y^2} \quad -\infty < y < \infty\end{aligned}$$

which is the pdf of a Cauchy's distribution.

- (b)  $y = \tan^{-1} x$  is a monotonic increasing function

$$\begin{aligned}\therefore f_Y(y) &= f_X(x) \left| \frac{dx}{dy} \right| \\ &= \frac{1}{\pi(1+x^2)} \sec^2 y \text{ (since } x = \tan y) \\ &= \frac{1}{\pi}, \quad -\frac{\pi}{2} < y < \frac{\pi}{2}\end{aligned}$$

i.e.,  $Y$  is uniformly distributed in  $(-\pi/2, \pi/2)$ .

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### Example 14

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- (a) If  $X$  has an arbitrary distribution function  $F_X(x)$ , find  $g(x)$  so that the random variable  $Y = g(X)$  may be uniformly distributed in  $(0, 1)$ .

- (b) If  $X$  is uniformly distributed in  $(0, 1)$ , find  $g(x)$ , so that the random variable  $Y = g(X)$  may have an arbitrary distribution with cdf  $F_Y(y)$ .
- (c) If  $X$  is uniformly distributed in  $(-1, 1)$ , find  $g(x)$ , so that the random variable  $Y = g(X)$  may have the density function  $f_Y(y) = 2e^{-2y}$ ,  $y > 0$ .
- (a)  $Y$  is to be uniform in  $(0, 1)$

$$\therefore f_Y(y) = 1 \text{ and } F_Y(y) = \int_0^y f_Y(y) dy = y$$

$$\therefore F_Y[g(x)] = g(x) \quad (1)$$

$$\begin{aligned} \text{Now } F_Y(y) &= P(Y \leq y) = P[g(X) \leq y] \\ &= P[X \leq g^{-1}(y)] \\ &= P(X \leq x) \quad [\text{since } y = g(x) \text{ and hence } x = g^{-1}(y)] \\ &= F_X(x) \end{aligned} \quad (2)$$

$$\text{i.e., } F_Y[g(x)] = F_X(x)$$

$$\text{i.e., } g(x) = F_X(x), \text{ from (1)}$$

- (b)  $X$  is uniform in  $(0, 1)$

$$\therefore f_X(x) = x$$

$$\text{By (2), } F_Y(y) = F_X(x)$$

$$\therefore F_Y[g(x)] = x$$

$$\therefore g(x) = F_Y^{-1}(x)$$

- (c)  $X$  is uniform in  $(-1, 1)$

$$\therefore f_X(x) = \frac{1}{2} \text{ and } F_X(x) = \frac{1}{2}(x + 1)$$

$$f_Y(y) = 2e^{-2y}, y > 0$$

$$\therefore F_Y(y) = \int_0^y 2e^{-2y} dy = 1 - e^{-2y}$$

$$\text{By (2), } 1 - e^{-2y} = \frac{1}{2}(x + 1)$$

$$\text{i.e., } 1 - e^{-2g(x)} = \frac{1}{2}(x + 1)$$

$$\therefore g(x) = \frac{1}{2} \log\left(\frac{2}{1-x}\right)$$

---

### Example 15

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If  $X$  and  $Y$  are independent RVs having density functions,

$$f_1(x) = \begin{cases} 2e^{-2x} & x \geq 0 \\ 0 & , x < 0 \end{cases} \text{ and}$$

$$f_2(y) = \begin{cases} 3e^{-3y} & y \geq 0 \\ 0 & y < 0 \end{cases} \text{ and}$$

find the density function of their sum  $U = X + Y$ .

By corollary under Theorem 1,

$$\begin{aligned} f_U(u) &= \int_0^u f_X(u-y) f_Y(y) dy, \quad u > 0 \\ &= \int_0^u 2e^{-2(u-y)} \times 3e^{-3y} dy \\ &= 6e^{-2u} \int_0^u e^{-y} dy \\ &= 6e^{-2u} (1 - e^{-u}), \quad u > 0 \end{aligned}$$

### Example 16

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If  $X$  and  $Y$  are independent RVs and if  $Y$  is uniformly distributed in  $(0, 1)$ , show that the density of  $Z = X + Y$  is given by  $f_Z(z) = F_X(z) - F_X(z-1)$ .

By Theorem 1,

$$f_Z(z) = \int_{-\infty}^{\infty} f_X(z-y) f_Y(y) dy$$

Now  $f_Y(y) = 1$ , provided  $0 < y < 1$

$$\begin{aligned} f_Z(z) &= \int_0^1 f_X(z-y) dy \\ &= \int_{z-1}^z f_X(t) dt, \text{ by putting } z-y=t \\ &= F_X(z) - F_X(z-1) \end{aligned}$$

### Example 17

---

If  $f_X(x) = ce^{-cx}$   $U(x)$  and  $f_Z(z) = c^2 z e^{-cz}$   $U(z)$ , find  $f_Y(y)$ , if  $Z = X + Y$  and  $X$  and  $Y$  are independent.

**Note**  $U(x)$  is the unit step function defined as

$$\begin{aligned} U(x) &= 1, & \text{if } x \geq 0 \\ &= 0, & \text{if } x < 0 \end{aligned}$$

$$f_X(x) = c e^{-cx}, x \geq 0, \text{ and } f_Z(z) = c^2 z e^{-cz}, z \geq 0$$

By Theorem 1,

$$c^2 z e^{-cz} = \int_0^z c e^{-c(z-y)} f_Y(y) dy \quad (1)$$

Recalling that

$$\frac{\partial}{\partial z} \int_a^{b(z)} \phi(z, y) dy = \int_a^{b(z)} \frac{\partial \phi}{\partial z} dy + \phi[z, b(z)] \frac{\partial b(z)}{\partial z}, \text{ where } a \text{ is a constant, and}$$

differentiating (1) with respect to  $z$  partially, we get,

$$\begin{aligned} c^2 e^{-cz} - c^3 z e^{-cz} &= \int_0^z -c^2 e^{-c(z-y)} f_Y(y) dy + c f_Y(z) \\ &= -c \{ c^2 z e^{-cz} \} + c f_Y(z), \text{ by (1)} \\ \therefore f_Y(z) &= c e^{-cz} z > 0 \\ \therefore f_Y(y) &= c e^{-cy} y > 0 \end{aligned}$$

### Example 18

The current  $I$  and the resistance  $R$  in a circuit are independent continuous RVs with the following density functions.

$$\begin{aligned} f_I(i) &= 2i, & 0 \leq i \leq 1 \\ &= 0 & \text{elsewhere} \\ f_R(r) &= r^2/9 & 0 \leq r \leq 3 \\ &= 0 & \text{elsewhere} \end{aligned}$$

Find the pdf of the voltage  $E$  in the circuit, where  $E = IR$ .

By Theorem 2,

$$f_Z(z) = \int_{-\infty}^{\infty} \frac{1}{|y|} f_X\left(\frac{z}{y}\right) f_Y(y) dy$$

when  $Z = XY$  and  $X$  and  $Y$  are independent RVs.

Taking  $x = i$ ,  $y = r$  and  $z = e$ , we have

$$\begin{aligned} f_E(e) &= \int_{-\infty}^{\infty} \frac{1}{|r|} f_I\left(\frac{e}{r}\right) \times f_R(r) dr \\ &= \int_e^3 \frac{1}{r} \times 2 \frac{e}{r} \times \frac{r^2}{9} dr \end{aligned}$$

#### Note

$$f_I\left(\frac{e}{r}\right) = 2 \frac{e}{r}, \text{ if } 0 \leq \frac{e}{r} \leq 1, \\ \text{i.e., if } r \geq e.$$

$$f_R(r) = \frac{r^2}{9}, \text{ if } 0 \leq r \leq 3$$

Hence, the limits for  $r$  are taken as  $e$  and 3.

$$\therefore f_E(e) = \frac{2e}{9}, (3 - e), \quad 0 \leq e \leq 3$$

### Example 19

---

If  $X$  and  $Y$  are independent RVs each following  $N(0, 2)$  prove that  $Z = \frac{X}{Y}$  follows a Cauchy's distribution.

$$f_X(x) = \frac{1}{2\sqrt{2\pi}} e^{-x^2/8} \quad \text{and} \quad f_Y(y) = \frac{1}{2\sqrt{2\pi}} e^{-y^2/8} \quad -\infty < x, y < \infty$$

By Theorem 3,

$$\begin{aligned} f_Z(z) &= \int_{-\infty}^{\infty} |y| f_X(yz) f_Y(y) dy \\ &= \int_{-\infty}^{\infty} |y| \frac{1}{2\sqrt{2\pi}} e^{-y^2 z^2/8} \times \frac{1}{2\sqrt{2\pi}} e^{-y^2/8} dy \\ &= \frac{1}{4\pi} \int_{-\infty}^{\infty} y e^{-(1+z^2)y^2/8} dy \quad (\text{since the integrand is an even function}) \\ &= \frac{1}{\pi} \times \frac{1}{1+z^2} \quad -\infty < z < \infty \end{aligned}$$

which is the pdf of Cauchy's distribution.

### Example 20

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If  $X$  and  $Y$  are independent RVs each following  $N(0, 2)$ , find the pdf of  $Z = 2X + 3Y$ .

Introduce the auxiliary RV  $W = Y$ .

$$\therefore z = 2x + 3y \text{ and } w = y$$

$$\text{Solving, } x = \frac{1}{2}(z - 3w) \text{ and } y = w$$

$$J = \frac{\partial(x, y)}{\partial(z, w)} = \begin{vmatrix} \frac{\partial x}{\partial z} & \frac{\partial x}{\partial w} \\ \frac{\partial y}{\partial z} & \frac{\partial y}{\partial w} \end{vmatrix} = \begin{vmatrix} \frac{1}{2} & -\frac{3}{2} \\ 0 & 1 \end{vmatrix} = \frac{1}{2}$$

Since  $X$  and  $Y$  are independent normal RVs  $f_{XY}(x, y) = \frac{1}{8\pi} e^{-(x^2 + y^2)/8}$ ,  $-\infty < x, y < \infty$ . The joint pdf of  $(Z, W)$  is given by

$$\begin{aligned} f_{ZW}(z, w) &= |J| f_{XY}(x, y) \\ &= \frac{1}{2} \times \frac{1}{8\pi} e^{-\{(z-3w)^2 + 4w^2\}/32} \quad -\infty < z, w < \infty \end{aligned}$$

The pdf of  $Z$  is the marginal pdf, obtained by integrating  $f_{ZW}(z, w)$  with respect to  $w$  over the range of  $w$ .

$$\begin{aligned} \therefore f_Z(z) &= \frac{1}{16\pi} \int_{-\infty}^{\infty} e^{-(13w^2 - 6zw + z^2)/32} dw \\ &= \frac{1}{16\pi} e^{-z^2/8 \times 13} \int_{-\infty}^{\infty} e^{-13\left(w - \frac{3z}{13}\right)^2/32} dw \\ &= \frac{1}{(2\sqrt{13})\sqrt{2\pi}} e^{-z^2/2(2\sqrt{13})^2} \quad -\infty < z < \infty \end{aligned}$$

which is  $N(0, 2\sqrt{13})$ .

### Example 21

If  $X$  and  $Y$  each follow an exponential distribution with parameter 1 and are independent, find the pdf of  $U = X - Y$ .

$$f_X(x) = e^{-x}, x > 0, \text{ and } f_Y(y) = e^{-y}, y > 0$$

Since  $X$  and  $Y$  are independent,

$$f_{XY}(x, y) = e^{-(x+y)}; x, y > 0$$

Consider the auxiliary RV  $V = Y$  along with

$$U = X - Y.$$

$$\therefore x = u + v \text{ and } y = v$$

$$J = \begin{vmatrix} \frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\ \frac{\partial y}{\partial u} & \frac{\partial y}{\partial v} \end{vmatrix} = \begin{vmatrix} 1 & 1 \\ 0 & 1 \end{vmatrix} = 1$$

The joint pdf of  $(U, V)$  is given by

$$\begin{aligned} f_{UV}(u, v) &= |J| f_{XY}(x, y) \\ &= e^{-(x+y)} = e^{-(u+2v)} \end{aligned}$$

The range space of  $(U, V)$  is found out from the map of the range space of  $(X, Y)$  under the transformations  $x = u + v$  and  $y = v$ .

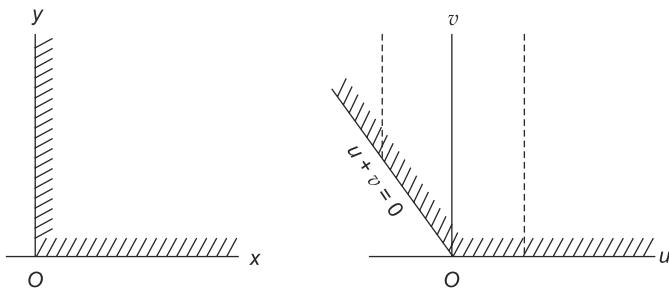


Fig. 3.6 (a) and (b)

Therefore, the range space of  $(U, V)$  is given by  $v > -u$ , when  $u < 0$  and  $v > 0$ , when  $u > 0$

Now the pdf of  $U$  is given by

$$f_U(u) = \int_{-u}^{\infty} e^{-(u+2v)} dv, \quad \text{when } u < 0$$

and  $= \int_0^{\infty} e^{-(u+2v)} dv, \quad \text{when } u > 0$

$$\therefore f_U(u) = \frac{1}{2} e^{-|u|}, \quad \text{when } u < 0$$

and  $= \frac{1}{2} e^{-u}, \quad \text{when } u > 0$

### Example 22

If the joint pdf of  $(X, Y)$  is given by  $f_{XY}(x, y) = x + y$ ;  $0 \leq x, y \leq 1$ , find the pdf of  $U = XY$ .

Introduce the auxiliary RV  $V = Y$ .

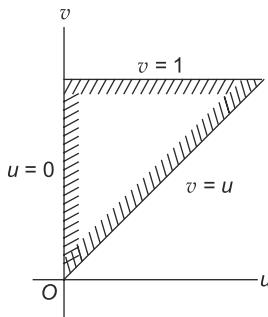
$$\therefore x = \frac{u}{v} \quad \text{and} \quad y = v$$

$$J = \begin{vmatrix} \frac{1}{v} & -\frac{u}{v^2} \\ 0 & 1 \end{vmatrix} = \frac{1}{v}$$

The joint pdf of  $(U, V)$  is given by

$$f_{UV}(u, v) = \frac{1}{|v|} f_{XY}(x, y) = \frac{1}{|v|} \left( \frac{u}{v} + v \right)$$

Range space of  $(X, Y)$  is given by  $0 \leq x \leq 1$  and  $0 \leq y \leq 1$

**Fig. 3.7**

$\therefore$  Range space of  $(U, V)$  is given by

$$0 \leq \frac{u}{v} \leq 1 \text{ and } 0 \leq v \leq 1$$

i.e.  $0 \leq u \leq v$  and  $0 \leq v \leq 1$

The pdf of  $U$  is given by

$$\begin{aligned} f_U(u) &= \int_u^1 f_{UV}(u, v) dv \\ &= \int_u^1 \frac{1}{v} \left( \frac{u}{v} + v \right) dv \\ &= 2(1-u), \quad 0 < u < 1 \end{aligned}$$

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### Example 23

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If  $X$  and  $Y$  are independent RVs with  $f_X(x) = e^{-x} U(x)$  and  $f_Y(y) = 3e^{-3y} U(y)$ , find  $f_Z(z)$ , if  $Z = \frac{X}{Y}$ .

Since  $X$  and  $Y$  are independent,  $f_{XY}(x, y) = 3e^{-(x+3y)}$ ,  $x, y \geq 0$

Introduce the auxiliary RV  $W = Y$ .

$\therefore x = zw$  and  $y = w$

$$J = \begin{vmatrix} w & z \\ 0 & 1 \end{vmatrix} = w$$

The joint pdf of  $(Z, W)$  is given by

$$\begin{aligned} f_{ZW}(z, w) &= |J| f_{XY}(x, y) \\ &= |w| \times 3e^{-(z+3w)}, \quad z, w \geq 0 \end{aligned}$$

The range space is obtained as follows:

Since  $y \geq 0$ ,  $w \geq 0$ . Since  $x \geq 0$ ,  $zw \geq 0$ .

As  $w \geq 0$ ,  $z \geq 0$ .

The pdf of  $z$  is given by

$$\begin{aligned} f_Z(z) &= \int_0^{\infty} 3w e^{-(z+3)w} dw \\ &= \frac{3}{(z+3)^2} \quad z \geq 0 \end{aligned}$$

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### Example 24

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If  $X$  and  $Y$  are independent RVs with pdf's  $e^{-x}$ ,  $x \geq 0$ , and  $e^{-y}$ ,  $y \geq 0$ , respectively, find the density functions of  $U = \frac{X}{X+Y}$  and  $V = X + Y$ . Are  $U$  and  $V$  independent?

Since  $X$  and  $Y$  are independent,  $f_{XY}(x, y) = e^{-(x+y)}$ .

Solving the equations  $u = \frac{x}{x+y}$  and  $v = x+y$ ,

we get  $x = uv$  and  $y = v(1-u)$ .

$$J = \begin{vmatrix} v & u \\ -v & (1-u) \end{vmatrix} = v$$

The joint pdf of  $(U, V)$  is given by

$$\begin{aligned} f_{UV}(u, v) &= |J| e^{-(x+y)} \\ &= |v| e^{-v} \end{aligned}$$

The range space of  $(U, V)$  is obtained as follows:

Since  $x \geq 0$ ,  $uv \geq 0$  and  $v(1-u) \geq 0$ .

Therefore, either  $u \geq 0$ ,  $v \geq 0$  and  $1-u \geq 0$ , i.e.,  $0 \leq u \leq 1$  and  $v \geq 0$

or  $u \leq 0$ ,  $v \leq 0$  and  $1-u \leq 0$ ,

i.e.,  $u \leq 0$  and  $u \geq 1$ , which is absurd.

Therefore, range space of  $(U, V)$  is given by  $0 \leq u \leq 1$  and  $v \geq 0$ .

$$\therefore f_{UV}(u, v) = v e^{-v}; \quad 0 \leq u \leq 1 \text{ and } v \geq 0.$$

$$\text{Pdf of } U \text{ is given by } f_U(u) = \int_0^{\infty} v e^{-v} du = 1, \quad 0 \leq u \leq 1$$

i.e.,  $U$  is uniformly distributed in  $(0, 1)$ .

$$\begin{aligned} \text{Pdf of } V \text{ is given by } f_V(v) &= \int_0^1 v e^{-v} du \\ &= v e^{-v} \quad v \geq 0 \end{aligned}$$

$$\text{Now } f_{UV}(u, v) = f_U(u) \times f_V(v)$$

Therefore,  $U$  and  $V$  are independent RVs.

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### Example 25

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If  $X$  and  $Y$  are independent RVs each normally distributed with mean zero and variance  $\sigma^2$ , find the density functions of  $R = \sqrt{X^2 + Y^2}$  and  $\phi = \tan^{-1} \left( \frac{Y}{X} \right)$ .

Since  $X$  and  $Y$  are independent  $N(0, \sigma)$ ,

$$f_{XY}(x, y) = \frac{1}{2\pi\sigma^2} e^{-(x^2+y^2)/2\sigma^2} \quad -\infty < x, y < \infty$$

$r = \sqrt{x^2 + y^2}$  and  $\theta = \tan^{-1} \left( \frac{y}{x} \right)$  are the transformations from Cartesians to polars.

Therefore, the inverse transformations are given by  $x = r \cos \theta$  and  $y = r \sin \theta$ .

$$J = \begin{vmatrix} \cos \theta & -r \sin \theta \\ \sin \theta & r \cos \theta \end{vmatrix} = r$$

The joint pdf of  $(R, \phi)$  is given by

$$f_{R\phi}(r, \theta) = \frac{|r|}{2\pi\sigma^2} e^{-r^2/2\sigma^2} \quad r \geq 0, 0 \leq \theta \leq 2\pi$$

**Note**  $(-\infty < x, y < \infty)$  and  $(r \geq 0 \text{ and } 0 \leq \theta \leq 2\pi)$  both represent the entire  $xy$ -plane.

The density function of  $R$  is given by

$$f_R(r) = \int_0^{2\pi} f_{R\phi}(r, \theta) d\theta = \frac{r}{\sigma^2} \times e^{-r^2/2\sigma^2} \quad r \geq 0$$

which is a Rayleigh distribution with parameter  $\sigma$ . The density function of  $\phi$  is given by

$$\begin{aligned} f_\phi(\theta) &= \int_0^\infty \frac{r}{2\pi\sigma^2} e^{-r^2/2\sigma^2} dr \\ &= \frac{1}{2\pi} \int_0^\infty e^{-t} dt, \text{ on putting } t = \frac{r^2}{2\sigma^2} \\ &= \frac{1}{2\pi}, 0 \leq \theta \leq 2\pi, \end{aligned}$$

which is a uniform distribution.

### Exercise 3

#### Part-A (Short-answer Questions)

1. If  $X$  and  $Y$  are two RVs where  $Y = g(X)$ , how are the density functions of  $X$  and  $Y$  related?

2. If the pdf of  $X$  is  $f_X(x) = 2x$ ,  $0 < x < 1$ , find the pdf of  $Y = 3X + 1$ .
3. If the pdf of  $X$  is  $f_X(x) = e^{-x}$ ,  $x > 0$ , find the pdf of  $Y = 2X + 1$ .
4. If the RV  $X$  is uniformly distributed in  $(0, 2)$ , find the pdf of  $Y = X^3$ .
5. If the RV  $X$  is uniformly distributed in  $(1, 2)$  find the pdf of  $Y = \frac{1}{X}$ .
6. If  $X$  is uniformly distributed in  $(0, 1)$ , find the pdf of  $Y = \frac{1}{2X+1}$ .
7. If  $X$  is uniformly distributed in  $(1, 2)$  find the pdf of  $Y = e^X$ .
8. If the pdf of a RV  $X$  is  $f_X(x) = 2x$ ,  $0 < x < 1$ , find the pdf of  $Y = e^{-X}$ .
9. If the cdf of a RV  $X$  is  $F(x)$ , show that the RV  $Y = F(x)$  follows a uniform distribution.
10. A RV  $X$  assumes three values  $-1, 0, 1$  with probabilities  $1/3, 1/2, 1/6$  respectively, find the probability distribution of  $Y = 3X + 1$ .
11. If  $X$  and  $Y$  are two RVs such that  $Y = X^2$ , how are the cdf's of  $X$  and  $Y$  related?
12. If  $X$  and  $Y$  are two RVs such that  $Y = X^2$ , how are the pdf's of  $X$  and  $Y$  related?
13. If the RV  $X$  is uniformly distributed in  $(-3, 3)$ , find the pdf of  $Y = X^2$ .
14. If the pdf of  $X$  is  $f(x) = e^{-x}$ ,  $x > 0$ , find the pdf of  $Y = X^2$ .
15. If the RVs  $X$  and  $Y$  are related by  $Y = |X|$ , how are the cdf's of  $X$  and  $Y$  related?
16. If the RVs  $X$  and  $Y$  are related by  $Y = |X|$ , how are the pdf's of  $X$  and  $Y$  related?
17. If the RV  $X$  is uniform distributed in  $(-1, 1)$ , find the pdf of  $Y = |X|$ .
18. If the RVs  $X$  and  $Y$  are related by  $Y = \sqrt{X}$ , how are their pdf's related?
19. If  $X$  is uniformly distributed in  $(0, 1)$ , find the pdf of  $Y = \sqrt{X}$ .
20. If the pdf of a RV  $X$  is  $f_X(x) = 2x$  in  $(0, 1)$ , find the pdf of  $Y = \sqrt{X}$ .
21. If  $X$  is uniformly distributed in  $\left(-\frac{\pi}{2}, \frac{\pi}{2}\right)$ , find the pdf of  $Y = \tan X$ .
22. Write down the formula to find the pdf of  $Z = X + Y$ , if  $X$  and  $Y$  are independent RVs with pdf's  $f_X(x)$  and  $f_Y(y)$  respectively.
23. Write down the formula to find the pdf of  $Z = XY$  in terms of the pdf's of  $X, Y$  if they are independent.
24. Write down the formula for the pdf of  $Z = \frac{X}{Y}$  in terms of the pdf's of  $X, Y$  if they are independent.
25. If  $Z = g(X, Y)$  and  $W = h(X, Y)$ , how are the joint pdf's of  $(X, Y)$  and  $(Z, W)$  related?
26. If  $Z = 2X + 3Y$  and  $W = Y$ , how are the joint pdf's of  $(X, Y)$  and  $(Z, W)$  related?
27. If  $U = XY$  and  $V = Y$ , how are the joint pdf's of  $(X, Y)$  and  $(U, V)$  related?

28. If  $U = \frac{X}{Y}$  and  $V = Y$ , how are the joint pdf's of  $(X, Y)$  and  $(U, V)$  related?
29. If  $U = X + Y$  and  $V = X - Y$ , how are the joint pdf's of  $(X, Y)$  and  $(U, V)$  related?
30. If  $x = R \cos \phi$  and  $y = R \sin \phi$ , how are the joint pdf's of  $(X, Y)$  and  $(R, \phi)$  related?

**Part-B**

31. The RV  $X$  is of the continuous type with distribution function  $F_X(x)$ . If  $g(x)$  is defined as

$$g(x) = \begin{cases} 1 & \text{if } x > 0 \\ -1 & \text{if } x \leq 0 \end{cases}$$

what is the distribution function of  $Y = g(X)$ ?

32. If  $g(x)$  is defined as

$$g(x) = \begin{cases} -b & \text{if } x \leq -b \\ x & \text{if } -b < x < b \\ b & \text{if } x \geq b \end{cases}$$

find the distribution function of  $Y = g(X)$  in terms of that of  $X$ .

33. If  $X$  follows a normal distribution with mean zero and variance  $\sigma^2$ , find the density function of  $Y = \frac{1}{2}(X + |X|)$ .
34. If the density function of a continuous RV  $X$  is given by  $f_X(x) = \frac{2}{9}(x + 1)$ , for  $-1 < x < 2$ , and  $= 0$ , otherwise, find the density function of  $Y = \frac{1}{2}(X + |X|)$ .
35. If  $X$  is a continuous RV with density function  $f_X(x) = e^{-x}$ ,  $x > 0$ , find the density function of  $Y = 2X + 1$ . Hence or otherwise find  $P(Y \geq 5)$ .
36. If the density function of a continuous RV  $X$  is given by  $f_X(x) = 2x$ , for  $0 < x < 1$ , and  $= 0$ , elsewhere, find the density and distribution functions of  $Y = 3X + 1$ .
37. If  $Y = aX^2$  ( $a > 0$ ), where  $X$  is a Gaussian RV with zero mean and variance  $\sigma^2$ , find the pdf of the RV  $Y$ .
38. (i) If the continuous RV  $X$  is uniformly distributed in  $(-3, 3)$ , find the density function of  $Y = 2X^2 - 3$ .  
(ii) If the continuous RV  $X$  is uniformly distributed in  $(-2, 2)$ , find the density function of  $Y = 6 - X^2$ .
39. (i) A fluctuating electric current  $I$  may be considered as a uniformly distributed RV over the interval  $(9, 11)$ . If this current flows through a  $2\Omega$  resistor, find the density function of the power  $P = 2I^2$ .

- (ii) If the voltage  $E$  across a resistor is a RV uniformly distributed between 5 and 10 V, find the density function of the power  $W = \frac{E^2}{r}$ , when  $r = 1000 \Omega$ .
40. (i) If the continuous RV  $X$  has density function  $f_X(x) = 2 e^{-2x}$ ,  $x > 0$ , find the density function of  $Y = X^2$ .  
(ii) If the density function of a continuous RV  $X$  is given by  $f_X(x) = e^{-x}$ ,  $x > 0$ , find the density function of  $Y = 3/(X + 1)^2$ .
41. If the RV  $X$  is uniform of a continuous RV  $X$  is given by  $f_X(x) = e^{-x}$ ,  $x > 0$ , find the density function of  $Y = X^3$ .
42. If the RV  $X$  is uniform in  $(-2\pi, 2\pi)$ , find the density function of RV (i)  $Y = X^3$  and (ii)  $Y = X^4$ .
43. If the radius  $R$  of a sphere is a continuous RV with pdf  $f_R(r) = 6r(1 - r)$ ,  $0 < r < 1$ , find the pdf of (i) the surface area  $S$  of the sphere and (ii) the volume  $V$  of the sphere.
44. If the random variable  $X$  is uniformly distributed over  $(-1, 1)$ , find the density function of  $Y = |X|$ .
45. (i) If the resistance  $R$  follows a uniform distribution between 900 and 1000  $\Omega$ , find the density of the corresponding conductance  $G = \frac{1}{R}$ .  
(ii) If the RV  $X$  is uniformly distributed in  $(0, 1)$ , find the pdf of  $Y = \frac{1}{X+1}$ .
46. (i) If the continuous RV  $X$  is uniformly distributed over  $(1, 3)$ , obtained the pdf of the RV  $Y = e^X$ .  
(ii) If the density function of a random variable  $X$  is given by  $f_X(x) = 2x$ , for  $0 < x < 1$ , and = 0, elsewhere, find the pdf of the RV  $Y = e^{-X}$ .
47. If the RV  $X$  is uniformly distributed over  $(-1, 1)$ , find the density function of (i)  $Y = \sin\left(\frac{\pi X}{2}\right)$  and (ii)  $Y = \cos\left(\frac{\pi X}{2}\right)$ .
48. If  $X$  is an arbitrary RV with continuous distribution function  $F_X(x)$  and if  $Y = F_X(X)$ , show that  $Y$  is uniformly distributed in  $(0, 1)$ .
49. If  $X$  and  $Y$  are independent RVs with density function  $f_X(x) = a e^{-\alpha x} U(x)$  and  $f_Y(y) = \beta e^{-\beta y} U(y)$ , find the density function of  $Z = X + Y$ .
50. If  $X$  and  $Y$  are independent RVs with identical uniform distributions in the interval  $(-1, 1)$  find the density function of  $Z = X + Y$ .
51. If  $X$  and  $Y$  are independent RVs with density functions  $f_X(x) = e^{-x} U(x)$  and  $f_Y(y) = 2e^{-2y} U(y)$ , find the density function of  $Z = X + Y$ .
52. If  $X$  and  $Y$  are independent RVs with density function  $f_X(x) = 1$ , in  $1 \leq x \leq 2$ , and  $f_Y(y) = \frac{y}{6}$  in  $2 \leq y \leq 4$ , find the density function of  $Z = XY$ .

53. If  $X$  and  $Y$  are independent RVs such that  $f_X(x) = \frac{x}{\alpha^2} e^{-x^2/2\alpha^2} U(x)$  and  $f_Y(y) = \frac{y}{\beta^2} e^{-y^2/2\beta^2} U(y)$  prove that the density function of  $Z = \frac{X}{Y}$  is given by

$$f_Z(z) = \frac{2\alpha^2}{\beta^2} \times \frac{z}{\left(z^2 + \frac{\alpha^2}{\beta^2}\right)} \times U(z)$$

54. If  $X$  and  $Y$  are independent RVs with identical uniform distributions in  $(0, a)$ , find the density function of  $Z = |X - Y|$ . (**Hint:** First find the cdf of  $Z$ )  
 55. If  $X$  and  $Y$  are independent RVs with identical uniform distributions in  $(0, 1)$ , find (i) the joint density function of  $(U, V)$ , where  $U = X + Y$  and  $V = X - Y$ , (ii) the density function of  $U$  and (iii) the density function of  $V$ .  
 56. Given the joint density function of  $X$  and  $Y$  as

$$f(x, y) = \begin{cases} \frac{1}{2} x e^{-y}, & 0 < x < 2, \quad y > 0 \\ 0, & \text{elsewhere} \end{cases}$$

find the distribution function of  $(X + Y)$ .

## **ANSWERS**

### **Exercise 3**

- if  $x = g^{-1}(y)$  is single valued, then  $f_Y(y) = \frac{F_X(x)}{|g'(x)|} = f_X(x) \left| \frac{dx}{dy} \right|$ . If  $x = g^{-1}(y)$  takes many values  $x_1, x_2, \dots, x_n$ , then  $f_Y(y) = f_X(x_1) \left| \frac{dx_1}{dy} \right| + f_X(x_2) \left| \frac{dx_2}{dy} \right| + \dots + f_X(x_n) \left| \frac{dx_n}{dy} \right|$ .
- $f_Y(y) = \left| \frac{dx}{dy} \right| f_X(x) = \frac{2}{9} (y-1)$  in  $1 < y < 4$
- $f_Y(y) = \left| \frac{dx}{dy} \right| f_X(x) = \frac{1}{2} e^{-\frac{1}{2}(y-1)}$ ,  $y > 1$
- $f_Y(y) = \left| \frac{1}{3} y^{-2/3} \right| \frac{1}{2} = \frac{1}{6} y^{-2/3}$   $0 < y < 8$
- $f_Y(y) = \left| -\frac{1}{y^2} \right| 1 = \frac{1}{y^2}$  in  $\frac{1}{2} < y < 1$

$$6. \quad f_Y(y) = \left| \frac{-(2x+1)^2}{2} \right| \Big|_{1=\frac{1}{2y^2}} \text{ in } \frac{1}{3} < y < 1$$

$$7. \quad f_Y(y) = \left| \frac{1}{e^x} \right| \Big|_{1=\frac{1}{y}} \text{ in } e < y < e^2$$

$$8. \quad f_Y(y) = \left| -e^x \right| \Big|_{2x=-\frac{2}{y}} \log y \text{ in } \frac{1}{e} < y < 1$$

$$9. \quad f_Y(y) = \left| \frac{1}{f_X(x)} \right| \Big|_{f_X(x)=1} \text{ in } 0 < y < 1$$

$$10. \quad \begin{array}{lcl} Y & : & -2 & 1 & 4 \\ p_Y & : & \frac{1}{3} & \frac{1}{2} & \frac{1}{6} \end{array}$$

$$11. \quad f_Y(y) = F_X(\sqrt{y}) - F_X(-\sqrt{y}) \quad \text{if } y \geq 0 \\ = 0, \quad \text{if } y < 0$$

$$12. \quad f_Y(y) = \frac{1}{2\sqrt{y}} \left[ f_Y(\sqrt{y}) + f_X(-\sqrt{y}) \right] \quad \text{if } y \geq 0 \\ = 0 \quad \text{if } y < 0$$

$$13. \quad f_Y(y) = \frac{1}{2\sqrt{y}} \left\{ \frac{1}{6} + \frac{1}{6} \right\} = \frac{1}{6\sqrt{y}} \text{ in } 0 < y < 9$$

$$14. \quad \text{Since } y = x^2 \text{ is m.i. in } x > 0, f_Y(y) = \left| \frac{1}{2x} \right| e^{-x} = \frac{1}{2\sqrt{y}} e^{-\sqrt{y}} \text{ in } y > 0$$

$$15. \quad f_Y(y) = F_X(y) - F_X(-y)$$

$$16. \quad f_Y(y) = f_X(y) + f_X(-y), y > 0$$

$$17. \quad f_Y(y) = \frac{1}{2} + \frac{1}{2} = 1 \text{ in } 0 < y < 1$$

$$18. \quad f_Y(y) = \left| \frac{dx}{dy} \right| f_X(x) = 2y f_X(y^2)$$

$$19. \quad f_Y(y) = 2y \text{ in } 0 < y < 1$$

$$20. \quad f_Y(y) = 2y \times 2y^2 = 4y^3 \text{ in } 0 < y < 1$$

$$21. \quad f_Y(y) = \left| \frac{dx}{dy} \right| f_X(x) = \frac{1}{\pi} \frac{1}{1+y^2} \text{ in } -\infty < y < \infty$$

$$22. \quad f_Z(z) = \int_{-\infty}^{\infty} f_X(z-y) f_y(y) dy$$

$$23. \quad f_Z(z) = \int_{-\infty}^{\infty} \frac{1}{|y|} f_X\left(\frac{z}{y}\right) f_Y(y) dy$$

24.  $f_Z(z) = \int_{-\infty}^{\infty} |y| f_X(yz) f_Y(y) dy$

25.  $f_{ZW}(z, w) = |J| f_{XY}(x, y)$ , where  $J = \begin{vmatrix} x_z & x_w \\ y_z & y_w \end{vmatrix}$

26.  $f_{ZW}(z, w) = \frac{1}{2} f_{XY}(x, y)$

27.  $f_{UV}(u, v) = \frac{1}{|v|} f_{XY}(x, y)$

28.  $f_{UV}(u, v) = |v| f_{XY}(x, y)$

29.  $f_{UV}(u, v) = \frac{1}{2} f_{XY}(x, y)$

30.  $f_{R\phi}(r, \theta) = |r| f_{XY}(x, y)$

31.  $f_Y(y) = 0$ , if  $y < -1$ ;  $= F_X(0)$ , if  $-1 \leq y < 1$  and  $= 1$ , if  $1 \leq y$

32.  $f_Y(y) = 0$ , if  $y < -b$ ;  $= F_X(y)$ , if  $-b \leq y < b$  and  $= 1$ , if  $b \leq y$

33.  $f_Y(y) = \frac{2}{\sigma \sqrt{2\pi}} e^{-y^2/2\sigma^2}$ ,  $y \geq 0$  and 0, elsewhere

34.  $f_Y(y) = \frac{1}{4} (y + 1)$ ,  $0 < y < 2$  and  $= 0$ , elsewhere

35.  $f_Y(y) = \frac{1}{2} e^{-(y-1)/2}$ ,  $y > 1$ ;  $= 1/e^2$

36.  $f_Y(y) = \frac{2}{9} (y - 1)$  for  $1 < y < 4$  and  $= 0$ , elsewhere;  $F_Y(y) = (y - 1)^2/9$

37.  $f_Y(y) = \frac{1}{\sigma \sqrt{2\pi} ay} e^{-y/2a\sigma^2}$ ,  $y > 0$

38. (i)  $f_Y(y) = \frac{1}{6\sqrt{2} \sqrt{y+3}}$  for  $-3 \leq y \leq 15$

(ii)  $f_Y(y) = \frac{1}{4 \sqrt{6-y}}$  for  $2 \leq y \leq 6$

39. (i)  $f_P(p) = \frac{1}{8} \sqrt{\frac{2}{p}}$  for  $162 \leq p \leq 242$

(ii)  $f_w(w) = \sqrt{\frac{10}{w}}$  for  $\frac{1}{40} \leq w \leq \frac{1}{10}$

40. (i)  $f_Y(y) = \frac{1}{\sqrt{y}} e^{-2\sqrt{y}}$  for  $y > 0$

(ii)  $f_Y(y) = \frac{1}{6} \left( \frac{3}{y} \right)^{3/2} e^{1-\sqrt{3/y}}$  for  $0 \leq y \leq 3$

41.  $f_Y(y) = \frac{1}{3} y^{-2/3} e^{-y^{1/3}}$  for  $y > 0$

42. (i)  $f_Y(y) = \frac{1}{12\pi} y^{-2/3}$  for  $-8\pi^3 \leq y \leq 8\pi^3$

(ii)  $f_Y(y) = \frac{1}{8\pi} y^{-3/4}$  for  $0 \leq y \leq 16\pi^4$

43. (i)  $f_S(s) = \frac{3}{4\pi} \left( 1 - \sqrt{\frac{s}{4\pi}} \right)$  for  $0 \leq s \leq 4\pi$

(ii)  $f_V(v) = \frac{3}{2\pi} \left\{ \left( \frac{3v}{4\pi} \right)^{-\frac{1}{3}} - 1 \right\}$  for  $0 \leq v \leq \frac{4\pi}{3}$

44.  $f_Y(y) = 1$  in  $(0, 1)$

45. (i)  $f_G(g) = \frac{1}{200g^2}$  for  $\frac{1}{1100} \leq g \leq \frac{1}{900}$

(ii)  $f_Y(y) = \frac{1}{y^2}$  for  $\frac{1}{2} \leq y \leq 1$

46. (i)  $f_Y(y) = \frac{1}{2y}$  for  $e \leq y \leq e^3$

(ii)  $f_Y(y) = -\frac{2}{y} \log y$  for  $\frac{1}{e} \leq y \leq 1$

47. (i)  $f_Y(y) = \frac{1}{\pi \sqrt{1-y^2}}$  for  $-1 \leq y \leq 1$

(ii)  $f_Y(y) = \frac{2}{\pi \sqrt{1-y^2}}$  for  $0 \leq y \leq 1$

49.  $f_Z(z) = \frac{\alpha\beta}{\beta-\alpha} (e^{-\alpha z} - e^{-\beta z})$  if  $\beta \neq \alpha$   
 $= \alpha^2 z e^{-\alpha z}$  if  $\beta = \alpha$

50.  $f_Z(z) = \frac{1}{4} (2+z)$  if  $z > 0$

$= \frac{1}{4} (2-z)$  if  $z < 0$

51.  $f_Z(z) = 2[e^{-z} - e^{-2z}] U(z)$

$$52. \quad f_z(z) = \begin{cases} \frac{1}{6}(z-2) & \text{when } 2 \leq z \leq 4 \\ \frac{1}{12}(8-z) & \text{when } 4 \leq z \leq 8 \end{cases}$$

$$54. \quad f_Z(z) = \frac{2}{a^2}(a-z), \text{ for } 0 < z < a$$

$$55. \quad (i) f_{UV}(u, v) = \frac{1}{2} \text{ for } 0 \leq u + v \leq 2 \text{ and } 0 \leq u - v \leq 2$$

$$(ii) \quad f_U(u) = \begin{cases} u & \text{if } 0 \leq u \leq 1 \\ 2-u & \text{if } 1 \leq u \leq 2 \end{cases}$$

$$(iii) \quad f_V(v) = \begin{cases} 1+v & \text{if } -1 \leq v \leq 0 \\ 1-v & \text{if } 0 \leq v \leq 1 \end{cases}$$

$$56. \quad \text{If } Z = X + Y, f_Z(z) = 0, \text{ if } z < 0$$

$$= \frac{1}{2} \left( \frac{z^2}{2} - e^{-z} - z + 1 \right), \text{ if } 0 \leq z < 2$$

$$= \frac{1}{2} [(1 - e^{-2}) + (1 + e^2)(e^{-2} - e^{-z})], \text{ if } 2 \leq z$$



# Chapter 4

## Statistical Averages

A discrete random variable (RV) is no doubt completely described by its probability mass function or probability distribution. Similarly, a continuous RV is completely described by its probability density function. For many purposes, this description is often considered to consist of too many details. It is sometimes simpler and more convenient to describe a RV or to characterise its distribution by a few parameters or summary measures that are representative of the distribution. These parameters or characteristic numbers are the various expected values or statistical averages of the RV.

**Definitions:** If  $X$  is a discrete RV, then *the expected value* or the mean value of  $g(X)$  is defined as

$$E\{g(X)\} = \sum_i g(x_i)p_i,$$

where  $p_i = P(X = x_i)$  is the probability mass function of  $X$ .

If  $X$  is a continuous RV with pdf  $f(x)$ , then

$$E\{g(X)\} = \int_{R_X} g(x)f(x)dx$$

Two expected values which are most commonly used for characterising a RV  $X$  are *its mean*  $\mu_X$  and *variance*  $\sigma_X^2$ , which are defined as follows:

$$\begin{aligned}\mu_X &= E(X) \\ &= \sum_i x_i p_i, \text{ if } X \text{ is discrete} \\ &= \int_{R_X} xf(x)dx, \text{ if } X \text{ is continuous}\end{aligned}$$

$$\begin{aligned}\text{Var}(X) &= \sigma_X^2 = E\{(X - \mu_X)^2\} \\ &= \sum_i (x_i - \mu_X)^2 p_i, \text{ if } X \text{ is discrete}\end{aligned}$$

$$\int_{R_x} (x - \mu_x)^2 f(x) dx, \text{ if } X \text{ is continuous}$$

The square root of variance is called the standard deviation. The mean of an RV is its average value and the variance is a measure of the spread or dispersion of the values of the RV.

**Note**  $Var(X) = E(X^2) - \{E(X)\}^2$

$$\begin{aligned} Var(X) &= E\{X - \mu_x\}^2 \\ &= E\{X^2 - 2\mu_x X + \mu_x^2\} \\ &= E(X^2) - 2\mu_x E(X) + \mu_x^2 \quad (\text{since } \mu_x \text{ is a constant}) \\ &= E(X^2) - \mu_x^2 \quad [\text{since } \mu_x = E(X)] \\ &= E(X^2) - \{E(X)\}^2 \end{aligned}$$

This modified formula for  $Var(X)$  holds good for both discrete and continuous RVs.

**Note** If  $X$  is a discrete RV and a constant, then (i)  $E(aX) = aE(X)$ , (ii)  $Var(aX) = a^2 Var(X)$ .

$$\begin{aligned} (i) \quad E(aX) &= \sum_j a x_i p_i \\ &= a \sum_j x_i p_i \\ &= aE(X) \end{aligned}$$

$$\begin{aligned} (ii) \quad Var(aX) &= E(a^2 X^2) - \{E(aX)\}^2 \quad (\text{by Note 1}) \\ &= a^2 E(X^2) - \{aE(X)\}^2 \\ &= a^2 [E(X^2) - \{E(X)\}^2] \\ &= a^2 Var(X) \end{aligned}$$

This result holds good for a continuous RV also.

**Moments:** If  $X$  is a discrete or continuous RV,  $E(X^n)$  is called  $n$ th order raw moment of  $X$  about the origin and denoted by  $\mu'_n$ .

$E\{(X - \mu_x)^n\}$  is called the  $n$ th order central moment of  $X$  and denoted by  $\mu_n$ .

$E\{|X|^n\}$  and  $E\{|X - \mu_x|^n\}$  are called absolute moments of  $X$ .

$E\{X - a\}^n$  and  $E\{|X - a|^n\}$  are called generalised moments of  $X$ .

## Expected Values of a Two-Dimensional RV

If  $(X, Y)$  is a two-dimensional discrete RV with joint probability mass function  $p_{ij}$ , then  $E\{g(X, Y)\} = \sum_j \sum_i g(x_i, y_i) p_{ij}$ .

If  $(X, Y)$  is a two-dimensional continuous RV with joint pdf  $f(x, y)$ , then

$$E\{g(X, Y)\} = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y) f(x, y) dx dy$$

## Properties of Expected Values

We give below the proofs of the properties for continuous RVs. Students can prove the properties for discrete RVs.

$$(i) \quad E\{g(X)\} = \int_{-\infty}^{\infty} g(x)f_X(x)dx, \text{ where } f_X(x) \text{ is the marginal density of } X.$$

### Proof

$$\begin{aligned} E\{g(X)\} &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x)f(x,y)dx dy \\ &= \int_{-\infty}^{\infty} g(x) \left[ \int_{-\infty}^{\infty} f(x,y)dy \right] dx \\ &= \int_{-\infty}^{\infty} g(x)f_X(x)dx \end{aligned}$$

$$(ii) \quad E\{h(Y)\} = \int_{-\infty}^{\infty} h(y)f_Y(y)dy$$

where  $f_Y(y)$  is the marginal density of  $Y$ .

(Proof is left as an exercise to the student.)

$$(iii) \quad E(X + Y) = E(X) + E(Y)$$

### Proof

$$\begin{aligned} E(X + Y) &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} (x + y)f(x,y)dx dy \\ &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} xf(x,y)dx dy + \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} yf(x,y)dx dy \\ &= E(X) + E(Y) \end{aligned}$$

$$(iv) \quad \text{In general, } E(XY) \neq E(X) \times E(Y), \text{ but if } X \text{ and } Y \text{ are independent RVs, } E(XY) = E(X) \times E(Y).$$

### Proof

$$\begin{aligned} E(XY) &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} xyf(x,y)dx dy \\ &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} xyf_X(x)f_Y(y)dx dy \\ &\quad (\text{since } X \text{ and } Y \text{ are independent}) \end{aligned}$$

$$\begin{aligned}
&= \int_{-\infty}^{\infty} xf_X(x)dx \times \int_{-\infty}^{\infty} yf_Y(y)dy \\
&= E(X) \times E(Y)
\end{aligned}$$

In general, if  $X$  and  $Y$  are independent,

$$E\{g(X) \times h(Y)\} = E\{g(X)\} \times E\{h(Y)\}$$

## Conditional Expected Values

If  $(X, Y)$  is a two-dimensional discrete RV with joint probability mass function  $p_{ij}$ , then the conditional expectations of  $g(X, Y)$  are defined as follows:

$$\begin{aligned}
E\{g(X, Y)/Y = Y_j\} &= \sum_i g(x_i, y_j) \times P(X = x_i/Y = y_j) \\
&= \sum_i g(x_i, y_j) \frac{p_{ij}}{p_{*j}}
\end{aligned}$$

$$\text{and } E\{g(X, Y)/X = x_i\} = \sum_j g(x_i, y_j) p_{ij} / p_{i*}$$

If  $(X, Y)$  is a two-dimensional continuous RV with joint pdf  $f(x, y)$ , then

$$E\{g(X, Y)/Y\} = \int_{-\infty}^{\infty} g(x, y) \times f(x/y) dx \text{ and}$$

$$E\{g(X, Y)/X\} = \int_{-\infty}^{\infty} g(x, y) \times f(y/x) dy$$

In particular, the conditional means are defined as

$$\mu_{Y/X} = E(Y/X) = \int_{-\infty}^{\infty} yf(y/x) dy \text{ and}$$

$$\mu_{X/Y} = E(X/Y) = \int_{-\infty}^{\infty} xf(x/y) dx$$

The conditional variance are defined as

$$\sigma_{Y/X}^2 = E\{(Y - \mu_{Y/X})^2\} = \int_{-\infty}^{\infty} (y - \mu_{Y/X})^2 f(y/x) dy \text{ and}$$

$$\sigma_{X/Y}^2 = E\{(X - \mu_{X/Y})^2\} = \int_{-\infty}^{\infty} (x - \mu_{X/Y})^2 f(x/y) dx$$

## Properties

(1) If  $X$  and  $Y$  are independent RVs, then  $E(Y/X) = E(Y)$  and  $E(X/Y) = E(X)$ .

### Proof

$$\begin{aligned}
 E(Y/X) &= \int_{-\infty}^{\infty} yf(y/x)dy \\
 &= \int_{-\infty}^{\infty} y \frac{f(x,y)}{f_X(x)}dy \\
 &= \int_{-\infty}^{\infty} y \frac{f_X(x) \times f_Y(y)}{f_X(x)}dy \quad (\text{since } X \text{ and } Y \text{ are independent}) \\
 &= \int_{-\infty}^{\infty} yf_Y(y)dy = E(Y)
 \end{aligned}$$

A similar proof can be given for the other result.

(2)  $E[E\{g(X, Y)/X\}] = E\{g(X, Y)\}$

### Proof

$$E\{g(X, Y)/X\} = \int_{-\infty}^{\infty} g(x, y)f(y/x)dy$$

Since  $E\{g(X, Y)/X\}$  is a function of the RV  $X$ ,

$$\begin{aligned}
 E[E\{g(X, Y)/X\}] &= \int_{-\infty}^{\infty} E\{g(X, Y)/X\}f_X(x)dx \\
 &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y)f(y/x)f_X(x)dxdy \quad [\text{from (1)}] \\
 &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y)f(x, y)dxdy \\
 &= E\{g(X, Y)\}
 \end{aligned}$$

In particular,

$$E\{E(Y/X)\} = E(Y) \text{ and similarly}$$

$$E\{E(X/Y)\} = E(X)$$

(3)  $E\{g_1(X) \times g_2(Y)\} = E[g_1(X) \times E\{g_2(Y)/X\}]$

### Proof

$$\begin{aligned}
 E\{g_1(X) \times g_2(Y)\} &= E[E\{g_1(X) \times g_2(Y)/X\}] \quad (\text{by Property (2)}) \\
 &= E[g_1(X) \times E\{g_2(Y)/X\}] \quad (\text{since } X \text{ is given})
 \end{aligned}$$

In particular,

$$\begin{aligned} E(XY) &= E[X \times E(Y/X)] \text{ and} \\ E(X^2Y^2) &= E[X^2 \times E(Y^2/X)] \end{aligned}$$

### Worked Example 4(A)

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#### Example 1

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A lot is known to contain 2 defectives and 8 non-defective items. If these items are inspected at random, one after another, what is the expected number of items that must be chosen in order to remove, both the defective ones?

Let the random variable  $X$  denote the number of items that must be drawn in order to remove both defective items.

Clearly,  $X$  takes the values 2, 3, 4, ..., 10.

$$P(X = r) = P(r \text{ items are to be drawn to remove both defectives})$$

$$\begin{aligned} &= P\{\text{the first } (r - 1) \text{ items drawn should contain 1 defective} \\ &\quad \text{and } r\text{th item drawn should be defective}\} \\ &= \frac{2C_1 \times 8C_{r-2}}{10C_{r-1}} \times \frac{1}{10 - (r - 1)} = \frac{2 \times 8C_{r-2}}{10C_{r-1}(11 - r)} \\ &\quad (r = 2, 3, \dots, 10) \end{aligned}$$

The probability distribution of  $X$  will then be as follows:

$X = r$	2	3	4	5	6	7	8	9	10
$p_r$	1/45	2/45	3/45	4/45	5/45	6/45	7/45	8/45	9/45

$$E(X) = \sum_{r=2}^{10} rp_r = \frac{22}{3}$$

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#### Example 2

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A box contains  $2^n$  tickets of which  $nC_r$  tickets bear the number  $r$  ( $r = 0, 1, 2, \dots, n$ ). Two tickets are drawn from the box. Find the expectation of the sum of their numbers.

Total number of tickets in the box.

$$\begin{aligned} \sum_{r=0}^n nC_r &= nC_0 + nC_1 + \dots + nC_n \\ &= (1 + 1)^n = 2^n, \text{ as given.} \end{aligned}$$

Let the RVs  $X$  and  $Y$  represent the numbers on the first and second tickets respectively.

Then  $E(X + Y) = E(X) + E(Y)$

$X$  can take the values 0, 1, 2, ...,  $n$  with probabilities  $\frac{nC_0}{2^n}, \frac{nC_1}{2^n}, \dots, \frac{nC_n}{2^n}$  respectively.

$$\begin{aligned}\therefore E(X) &= 1 \times \frac{nC_1}{2^n} + 2 \times \frac{nC_2}{2^n} + \cdots + n \times \frac{nC_n}{2^n} \\ &= \frac{n}{2^n} \{(n-1)C_0 + (n-1)C_1 + \cdots + {}^{(n-1)}C_{n-1}\} \\ &= \frac{n}{2^n} (1+1)^{n-1} = \frac{n}{2}\end{aligned}$$

Similarly,  $E(Y) = \frac{n}{2}$

$\therefore E(X + Y) = n$

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### Example 3

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Find the mean and variance of the Pascal's (negative binomial distribution) distribution, given by  $P(X = k) = \binom{n+k-1}{k} p^n q^k, k = 0, 1, 2, \dots$

$$\begin{aligned}E(X) &= \sum_{k=0}^{\infty} k \binom{n+k-1}{k} p^n q^k \\ &= p^n [1 \times nC_1 q^1 + 2(n+1)C_2 q^2 + 3(n+2) C_3 q^3 + \dots] \\ &= np^n q [1 + (n+1)C_1 q + (n+2) C_2 q^2 + \dots] \\ &= np^n q [1 - q]^{-(n+1)} = \frac{nq}{p} \\ \\ E(X^2) &= \sum_{k=0}^{\infty} k^2 \binom{n+k-1}{k} p^n q^k \\ &= p^n [1^2 nC_1 q^1 + 2^2 (n+1)C_2 q^2 + 3^2 (n+2) C_3 q^3 + \dots] \\ &= p^n \left[ nq + \frac{1 \times 2 + 2}{2!} (n+1)nq^2 + \frac{(2 \times 3 + 3)}{3!} (n+2)(n+1)nq^3 + \dots \right] \\ &= np^n q \left[ \left\{ 1 + \frac{(n+1)}{1!} q + \frac{(n+1)(n+2)}{2!} q^2 + \dots \right\} + (n+1)q \times \right. \\ &\quad \left. \left\{ 1 + \frac{(n+2)}{1!} q + \frac{(n+2)(n+3)}{2!} q^2 + \dots \right\} \right] \\ &= np^n q [1 - q]^{-(n+1)} + (n+1)q(1 - q)^{-(n+2)} \\ &= np^n q \left[ \frac{1}{p^{n+1}} + \frac{(n+1)q}{p^{n+2}} \right] \\ &= \frac{nq}{p} + \frac{n(n+1)q^2}{p^2}\end{aligned}$$

$$\begin{aligned}\therefore \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\ &= \frac{nq}{p} + \frac{n(n+1)q^2}{p^2} - \frac{n^2 q^2}{p^2} \\ &= \frac{nq}{p} \left(1 + \frac{q}{p}\right) = \frac{nq}{p^2}\end{aligned}$$

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**Example 4**


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If the continuous RV  $X$  has Rayleigh density  $f(x) = \frac{x}{\alpha^2} e^{-x^2/2\alpha^2} \times U(x)$ , find  $E(X^n)$  and deduce the values of  $E(X)$  and  $\text{var}(X)$ .

By definition,

$$\begin{aligned}E(X^n) &= \int_0^\infty x^n \times \frac{x}{\alpha^2} e^{-x^2/2\alpha^2} dx \\ &= \int_0^\infty (2\alpha^2 t)^{n/2} e^{-t} dt \quad \left(\text{putting } \frac{x^2}{2\alpha^2} = t\right) \\ &= 2^{n/2} \alpha^n \int_0^\infty t^{n/2} e^{-t} dt \\ &= 2^{n/2} \alpha^n [(k+1)] \text{ if } n = 2k \\ &= 2^{n/2} \alpha^n [k = 2^{n/2} \alpha^n | n/2, \text{ if } n \text{ is even}] \\ \\ E(X^n) &= 2^{n/2} \alpha^n \left[ \left( \frac{2k+3}{3} \right) \right] \text{ if } n = 2k+1 \\ &= 2^{n/2} \alpha^n \frac{2k+1}{2} \times \frac{2k-1}{2} \cdots \frac{3}{2} \times \frac{1}{2} \times \left[ \left( \frac{1}{2} \right) \right] \\ &= 2^{n/2} \alpha^n \frac{1 \times 3 \times 5 \times \cdots \times n}{2^{(n+1)/2}} \sqrt{\pi} \\ &= 1 \times 3 \times 5 \times \cdots \times n \alpha^n \sqrt{\pi/2} \text{ if } n \text{ is odd} \\ \therefore E(X) &= \alpha \sqrt{\pi/2}; E(X^2) = 2\alpha^2; \text{var}(X) = \left(2 - \frac{\pi}{2}\right) \alpha^2\end{aligned}$$

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**Example 5**


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A line of length  $a$  units is divided into two parts. If the first part is of length  $X$ , find  $E(X)$ ,  $\text{var}(X)$  and  $E\{X(a-X)\}$ .

Since the positions of the point of division are equally likely,  $X$  is uniformly distributed in  $(0, a)$ .

$$\therefore f(x) = \frac{1}{a}$$

$$E(X) = \int_0^a xf(x)dx = \frac{1}{a} \int_0^a x dx = \frac{a}{2}$$

$$E(X^2) = \int_0^a x^2 f(x)dx = \frac{a^3}{3}$$

$$\therefore \text{Var}(X) = E(X^2) - \{E(X)\}^2 = \frac{a^2}{3} - \frac{a^2}{4} = \frac{a^2}{12}$$

$$E\{X(a-X)\} = a E(X) - E(X^2) = \frac{a^2}{3} - \frac{a^2}{4} = \frac{a^2}{12}$$

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### Example 6

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If  $X$  is a continuous RV, prove that

$$\begin{aligned} E(X) &= \int_0^\infty [1 - F(x)]dx - \int_{-\infty}^0 F(x)dx \\ E(X) &= \int_{-\infty}^\infty xf(x)dx = \int_{-\infty}^0 x dF(x) - \int_0^\infty x d\{1 - F(x)\} \\ &\quad [\text{since } F'(x) = f(x)] \\ &= [xF(x)]_{-\infty}^0 - \int_{-\infty}^0 F(x)dx - [x\{1 - F(x)\}]_0^\infty + \int_0^\infty \{1 - F(x)\}dx \\ &= \int_0^\infty \{1 - F(x)\}dx - \int_{-\infty}^0 F(x)dx \\ &\quad [\text{since } F(-\infty) = 0 \text{ and } F(\infty) = 1] \end{aligned}$$

---

### Example 7

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If the random variable  $X$  follows  $N(0, 2)$  and  $Y = 3X^2$ , find the mean and variance of  $Y$ .

Since  $X$  follows  $N(0, 2)$ ,  $E(X) = 0$  and  $\text{var}(X) = 4$

$$\therefore E(X^2) = \text{var}(X) + \{E(X)\}^2 = 4$$

$$\text{Now } E(Y) = E(3X^2) = 3 \times 4 = 12$$

$$E(Y^2) = E(9X^4) = 9 \times 3 \times 2^4$$

[since for the normal distribution  $N(0, \sigma)$ ,  $E(X^{2r}) = \frac{(2r)! \sigma^{2r}}{2^r r!}$ ]

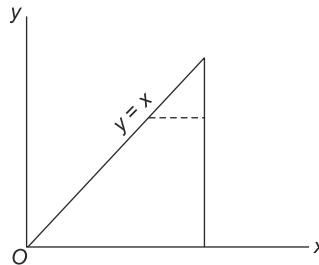
$$\begin{aligned}\text{Var}(Y) &= E(Y^2) - \{E(Y)\}^2 \\ &= 27 \times 2^4 - 12^2 = 288\end{aligned}$$

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### Example 8

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If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 24y(1-x)$ ,  $0 \leq y \leq x \leq 1$ , find  $E(XY)$ .



**Fig. 4.1**

$$\begin{aligned}E(XY) &= \int_0^1 \int_y^1 xyf(x, y) dx dy \\ &= 24 \int_0^1 \int_y^1 xy^2(1-x) dx dy \\ &= 24 \int_0^1 y^2 \left( \frac{1}{6} - \frac{y^2}{2} + \frac{y^3}{3} \right) dy \\ &= \frac{4}{15}\end{aligned}$$

---

### Example 9

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If  $X$  and  $Y$  are two independent RVs with  $f_X(x) = e^{-x}U(x)$  and  $f_Y(y) = e^{-y}U(y)$  and  $Z = (X - Y) U(X - Y)$ , prove that  $E(Z) = 1/2$ .

$$U(X - Y) = \begin{cases} 1 & \text{if } X > Y \\ 0 & \text{if } X < Y \end{cases}$$

$$\therefore Z = \begin{cases} X - Y & \text{if } X > Y \\ 0 & \text{if } X < Y \end{cases}$$

$$\begin{aligned}
 E(Z) &= \int_0^\infty \int_0^\infty z e^{-(x+y)} dx dy \\
 &\quad [\text{since } X \text{ and } Y \text{ are independent, } f(x, y) = f_X(x) \times f_Y(y)] \\
 &= \int_0^\infty \int_y^\infty (x-y) e^{-(x+y)} dx dy \\
 &= \int_0^\infty e^{-y} [(x-y)(-e^{-x}) - e^{-x}]_y^\infty dy \\
 &= \int_0^\infty e^{-2y} dy = \frac{1}{2}
 \end{aligned}$$

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### Example 10

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The joint pdf of  $(X, Y)$  is given by  $f(x, y) = 24xy$ ;  $x > 0, y > 0, x + y \leq 1$ , and  $f(x, y) = 0$ , elsewhere, find the conditional mean and variance of  $Y$ , given  $X$ .

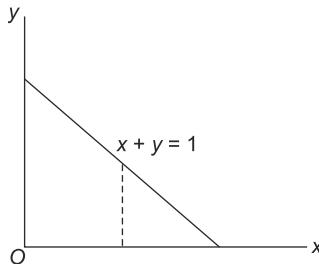


Fig. 4.2

$$\begin{aligned}
 f_X(x) &= \int_0^{1-x} 24xy dy \\
 &= 12x(1-x)^2, \quad 0 < x < 1
 \end{aligned}$$

$$\text{Now, } f(y/x) = \frac{f(x, y)}{f_X(x)} = \frac{2y}{(1-x)^2}, \quad 0 < y < 1-x$$

$$\begin{aligned}
 E(Y/X=x) &= \int_0^{1-x} y f(y/x) dy \\
 &= \int_0^{1-x} \frac{2y^2}{(1-x)^2} dy = \frac{2}{3}(1-x)
 \end{aligned}$$

$$E(Y^2/x) = \int_0^{1-x} y^2 \times f(y/x) dy = \frac{1}{2}(1-x)^2$$

$$\text{Var}(Y^2/x) = E(Y^2/x) - \{E(Y/x)\}^2$$

$$\begin{aligned}
 &= \frac{1}{2}(1-x)^2 - \frac{4}{9}(1-x)^2 \\
 &= \frac{1}{18}(1-x)^2
 \end{aligned}$$

### Example 11

If  $(X, Y)$  is uniformly distributed over the semicircle bounded by  $y = \sqrt{1-x^2}$  and  $y = 0$ , find  $E(X/Y)$  and  $E(Y/X)$ . Also verify the  $E\{E(X/Y)\} = E(X)$  and  $E\{E(Y/X)\} = E(Y)$ .

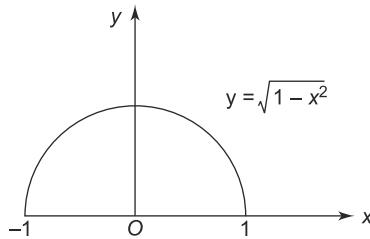


Fig. 4.3

$$f(x, y) = k$$

$$\iint f(x, y) dy dx = 1$$

$$\text{i.e., } \int_{-1}^1 \int_0^{\sqrt{1-x^2}} k dy dx = 1$$

$$\text{i.e., } 2k \int_0^1 \sqrt{1-x^2} dx = 1$$

$$\therefore k = \frac{2}{\pi}$$

$$f_X(x) = \int_0^{\sqrt{1-x^2}} \frac{2}{\pi} dy = \frac{2}{\pi} \sqrt{1-x^2}, -1 \leq x \leq 1$$

$$f_Y(y) = \int_{-\sqrt{1-y^2}}^{\sqrt{1-y^2}} \frac{2}{\pi} dx = \frac{4}{\pi} \sqrt{1-y^2}, 0 \leq y \leq 1$$

$$f(x/y) = \frac{f(x, y)}{f_Y(y)} = \frac{1}{2\sqrt{1-y^2}}, -\sqrt{1-y^2} \leq x \leq \sqrt{1-y^2}$$

$$f(y/x) = \frac{1}{\sqrt{1-x^2}}, 0 \leq y \leq \sqrt{1-x^2}$$

$$E(X) = \int_{-1}^1 xf_X(x) dx = \frac{2}{\pi} \int_1^1 x\sqrt{1-x^2} dx = 0 \quad (\text{since the integrand is odd})$$

$$\begin{aligned} E(X/Y) &= \int_{-\sqrt{1-y^2}}^{\sqrt{1-y^2}} xf(x/y) dx \\ &= \frac{1}{2\sqrt{1-y^2}} \int_{-\sqrt{1-y^2}}^{\sqrt{1-y^2}} dx = 0 \quad (\text{since the integrand is odd}) \end{aligned}$$

$$\therefore E\{E(X/Y)\} = E\{0\} = 0 = E(X)$$

$$E(Y) = \int_0^1 yf_Y(y) dy = \frac{4}{\pi} \int_0^1 y\sqrt{1-y^2} dy = \frac{4}{3\pi}$$

$$E(Y/X) = \int_0^{\sqrt{1-x^2}} yf(y/x) dy = \frac{1}{\sqrt{1-x^2}} \left( \frac{y^2}{2} \right)_0^{\sqrt{1-x^2}} = \frac{1}{2}\sqrt{1-x^2}$$

$$\begin{aligned} \therefore E\{E(Y/X)\} &= E\left\{ \frac{1}{2}\sqrt{1-X^2} \right\} \\ &= \int_{-1}^1 \frac{1}{2}\sqrt{1-x^2} f_X(x) dx \\ &= \frac{2}{\pi} \int_0^1 (1-x^2) dx = \frac{4}{3\pi} \end{aligned}$$

$$\therefore E\{E(Y/X)\} = E(Y)$$

### Example 12

If  $(X, Y)$  follows a bivariate normal distribution  $N(0, 0; \sigma_x, \sigma_y; r)$ , find  $E(Y/X)$ ,  $E(Y^2/X)$ ,  $E(XY)$  and  $E(X^2Y^2)$ .

$$\begin{aligned} f(x, y) &= \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)}\left(\frac{x^2}{\sigma_x^2} - \frac{2rxy}{\sigma_x\sigma_y} + \frac{y^2}{\sigma_y^2}\right)\right\} \\ &= \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)}\left(\frac{y}{\sigma_y} - \frac{rx}{\sigma_x}\right)^2 - \frac{x^2}{2\sigma_x^2}\right\} \\ f_X(x) &= \frac{1}{\sigma_x\sqrt{2\pi}} \exp(-x^2/2\sigma_x^2) \end{aligned}$$

[refer to the worked Example 12 in Chapter 2 on two-dimensional RVs]

$$\therefore f(y/x) = \frac{f(x,y)}{f_X(x)} = \frac{1}{\sigma_y\sqrt{1-r^2}\sqrt{2\pi}}$$

$$\exp\left\{-\frac{1}{2(1-r^2)\sigma_y^2}\left(y - \frac{rx\sigma_y}{\sigma_x}\right)^2\right\}$$

which is a  $N\left(\frac{rx\sigma_y}{\sigma_x}, \sigma_y\sqrt{1-r^2}\right)$

$$\therefore E(Y/X = x) = \int_{-\infty}^{\infty} yf(y/x)dy = \frac{rx\sigma_y}{\sigma_y}$$

$$\text{and } \text{Var}(Y/X = x) = \sigma_y^2(1-r^2)$$

$$\text{Var}(Y/X) = E(Y^2/X) - \{E(Y/X)\}^2$$

$$\therefore E(Y^2/X) = \sigma_y^2(1-r^2) + r^2x^2\sigma_y^2/\sigma_x^2$$

By Property 3 of conditional expected values,

$$\begin{aligned} E(XY) &= E\{XE(Y/X)\} \\ &= E\left\{r\frac{\sigma_y}{\sigma_x}X^2\right\} = r\frac{\sigma_y}{\sigma_x} \times \sigma_x^2 = r\sigma_x\sigma_y \end{aligned}$$

Again, by the same property

$$\begin{aligned} E(X^2Y^2) &= E\{X^2 \times E(Y^2/X)\} \\ &= E\left[\sigma_y^2(1-r^2)X^2 + \frac{r^2\sigma_y^2}{\sigma_x^2}X^4\right] \\ &= \sigma_y^2(1-r^2)E(X^2) + \frac{r^2\sigma_y^2}{\sigma_x^2}E(X^4) \\ &= \sigma_x^2\sigma_y^2(1-r^2) + \frac{r^2\sigma_y^2}{\sigma_x^2} \times 3\sigma_x^4 \\ &= (1+2r^2)\sigma_x^2\sigma_y^2 \end{aligned}$$

### Exercise 4(A)

#### Part-A (Short-answer Questions)

1. Define the expected value of  $g(X)$ , where  $X$  is a RV.
2. Define the mean and variance of a RV.
3. Prove that  $\text{Var}(X) = E(X^2) - E^2(X)$ .
4. If  $X$  is a RV, prove that  $E(X^2) \geq \{E(X)\}^2$ .
5. If  $X$  is a discrete/continuous RV prove that  $E(aX + b) = aE(X) + b$  and  $\text{Var}(aX) = a^2\text{Var}(X)$ .

6. If  $\mu_X$  and  $\sigma_X$  are the mean and SD of the RV  $X$ , find  $\mu_Y$  and  $\sigma_Y$ , where  

$$Y = \frac{1}{\sigma_X}(X - \mu_X).$$
7. Define the raw and central moments of a RV and state the relation between them.
8. The probability distribution of a RV  $X$  is given by  

$X:$	0	1	2	3
$p_X:$	0.1	0.3	0.4	0.2

 find  $E(Y)$ , where  $Y = X^2 + X$ .
9. Find the mean of the RV  $X$ , if its pmf is given by  $P(x=j) = (1-a)a^j, j = 0, 1, 2, \dots, \infty$ .
10. Find the mean of the RV  $X$  if its pdf is  $f(x) = 6x(1-x), 0 \leq x \leq 1$ .
11. Find the mean and variance of the uniform distributed in  $(a, b)$ .
12. Find the mean and variance of a RV  $X$ , that is uniformly distributed in  $(2, 8)$ .
13. If  $X$  is uniformly distributed in  $(1, 2)$  and  $Y = X^3$ , find the mean of  $Y$ .
14. Obtain the mean of the binomial distribution  $B(n; p)$ .
15. Obtain the mean of the Poisson distribution  $P(\lambda)$ .
16. Find the binomial distribution whose mean is 6 and SD is  $\sqrt{2}$ .
17. If  $X$  is a binomial RV with mean 2.4 and variance 1.44, find  $P(X=7)$ .
18. If  $X$  is binomially distributed with  $n=5$  such that  $P(X=1) = 2P(X=2)$ , find  $E(X)$  and  $\text{Var}(X)$ .
19. If  $X$  is binomially distributed with  $n=6$  such that  $P(X=2) = 9 P(X=4)$ , find  $E(X)$  and  $\text{Var}(X)$ .
20.  $X$  is a Poisson RV such that  $P(X=1) = P(X=2)$ , find  $E(X)$  and  $E(X^2)$ .
21. Find the mean of the geometric distribution given by  $P(X=r) = pq^r$  ( $r = 0, 1, 2, \dots$ ). where  $p+q=1$ .
22. On the average, how many times must a dice be thrown until a '6' is obtained?
23. Find the mean and variance of the exponential distribution given by  $f(x) = \lambda e^{-\lambda x}, x > 0$ .
24. If the RV  $X$  follows  $N(0, 2)$ , find  $E(X^2)$ .
25. Define the expected value of  $g(X, Y)$ , where  $(X, Y)$  is a two-dimensional continuous RV with joint pdf  $f(x, y)$ .
26. If  $(X, Y)$  is a two-dimensional continuous RV, express  $E[g(X)]$  and  $E[h(Y)]$  in terms of the marginal densities of  $X$  and  $Y$ .
27. If  $X$  and  $Y$  are independent RVs prove that  $E(XY) = E(X) \times E(Y)$ .
28. If  $X$  and  $Y$  are independent RVs with means 2 and 3 and variances 1 and 2 respectively, find the mean and variance of  $Z = 2X - 5Y$ .
29. If  $(X, Y)$  is a two-dimensional continuous RV, define  $E\{g(X, Y)/X\}$  and  $E\{g(X, Y)/Y\}$ .
30. If  $(X, Y)$  is a two-dimensional continuous RV, define conditional mean and conditional variance of  $X$ , given  $Y$ .

31. If  $X$  and  $Y$  are independent RVs, prove that  $E(Y/X) = E(Y)$ .
32. If  $X$  and  $Y$  are independent RVs, prove that  $E\{E(Y/X)\} = E(Y)$ .
33. If  $X$  and  $Y$  are independent RVs, prove that  $E(X/Y) = E\{X \cdot E(Y/X)\}$ .
34. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 2 - x - y$ , in  $0 \leq x \leq y \leq 1$ , find  $E(X)$  and  $E(Y)$ .
35. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 2$ , in  $0 \leq x \leq y \leq 1$ , find  $E(X)$ .

### Part-B

36. What is the expectation of the number of failures preceding the first success in an infinite series of independent trials with constant probability  $p$  of success?
37. What is the expectation of (a) the sum of the points on  $n$  dice? and (b) the product of the points on  $n$  dice?
38. Three tickets are chosen at random without replacement from 100 tickets, numbered 1, 2, 3, ..., 100. Find the expectation of the sum of the numbers.
39. From an urn containing 3 red and 2 black balls, a man is to draw 2 balls at random without replacement, being promised Rs. 20/- for each red ball he draws and Rs. 10/- for each black ball. Find his expectation.
40. If  $X$  follows a uniform distribution in  $(a, b)$ , find  $E(X)$  and  $\text{Var}(X)$ .
41. Find the mean and variance of the geometric distribution given by  $P(X = r) = pq^r$ ,  $r = 0, 1, 2, \dots$ ;  $p + q = 1$ .
42. Find the mean and variance of the binomial distribution  $B(n; p)$ .
43. Find the mean and variance of the Poisson distribution  $P(\lambda)$ .
44. If the continuous RV  $X$  follows a normal distribution  $N(0, \sigma^2)$ , prove that

$$(i) \quad E(X^n) = \begin{cases} 0 & \text{if } n \text{ is odd} \\ 1 \times 3 \times 5 \cdots (n-1) \sigma^n & \text{if } n \text{ is even} \end{cases}$$

$$(ii) \quad E(|X|^n) = \begin{cases} \sqrt{\frac{2}{\pi}} \times 2^{\frac{n-1}{2}} \left| \frac{(n-1)/2}{\sigma} \right| \sigma^n & \text{if } n \text{ is odd} \\ 1 \times 3 \times 5 \cdots (n-1) \sigma^n & \text{if } n \text{ is even} \end{cases}$$

45. If the continuous RV  $X$  has a Maxwell density, given by

$$f(x) = \sqrt{\frac{2}{\pi}} \frac{x^2}{\alpha^3} e^{-x^2/2\alpha^2} U(x), \text{ prove that}$$

$$E(X^n) = \begin{cases} 1, 3, \dots, (2k+1) \alpha^{2k} & \text{if } n = 2k \\ \sqrt{\frac{\pi}{2}} \times 2^k k \times \alpha^{2k-1} & \text{if } n = 2k-1 \end{cases}$$

Hence, find the mean and variance of the distribution.

46. If  $X$  has a Rayleigh density with parameter  $\alpha$  and  $Y = a + bX^2$ , prove that  $\sigma_y^2 = 4b^2\alpha^4$ .
47. (i) If  $Y = aX + b$ , show that  $\sigma_Y = a\sigma_X$  and  
(ii) If  $Y = (X - \mu_X)/\sigma_X$ , find  $\mu_Y$  and  $\sigma_Y$ .
48. If  $X$  is a RV for which  $E(X) = 10$  and  $\text{Var}(X) = 25$ , for what positive values of  $a$  and  $b$  does  $Y = aX - b$  have expectation 0 and variance 1?
49. If  $X$  is uniformly distributed in (1, 2) and  $Y = X^3$ , find the mean and variance of  $Y$ .
50. If the continuous RV  $X$  has the density function  $f(x) = 2xe^{-x^2}$ ,  $x \geq 0$ , and if  $Y = X^2$ , find the mean and variance  $Y$ .
51. If  $X$  and  $Y$  are independent random variables with density functions  $f_X(x) = \frac{8}{x^3}$ ,  $x > 2$ , and  $f_Y(y) = 2y$ ,  $0 < y < 1$ , respectively, and  $Z = XY$ , find  $E(Z)$ .
52. If each of the independent RVs  $X$  and  $Y$  follows  $N(0, \sigma)$  and  $Z = |X - Y|$ , prove that  $E(Z) = 2\sigma/\sqrt{\pi}$  and  $E(Z)^2 = 2\sigma^2$ .
53. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 2$ ,  $0 \leq x \leq y \leq 1$ , find the conditional mean and conditional variance of  $X$ , given that  $Y = y$ .
54. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 21x^2y^3$ ,  $0 \leq x \leq y \leq 1$ , find the conditional mean and variance of  $X$ , given that  $Y = y$ ,  $0 < y < 1$ .
55. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 3xy(x + y)$ ,  $0 \leq x, y \leq 1$ , verify that  $E\{E(Y/X)\} = E(Y) = \frac{17}{24}$ .

## Linear Correlation

In many situations, the outcome of a random experiment will have two measurable characteristics, viz. will result in two random variables  $X$  and  $Y$ . Often we will be interested in finding whether the two different RV's are related to each other. If they are related, we will try to determine the nature fo relationship and degree of relationship (correlation). Assuming that there is some correlation between  $X$  and  $Y$ , we will then try to find a formula expressing the relationship and use this formula to predict the most likely value of one RV corresponding to any given value of other RV.

To examine whether the two RV's are inter-related, we collect  $n$  pairs of values of  $X$  and  $Y$  corresponding to  $n$  repetitions of the random experiment. Let them be  $(x_1, y_1), (x_2, y_2), \dots, (x_n, y_n)$ . Then we plot the points with co-ordinates  $(x_1, y_1), \dots, (x_n, y_n)$  on a graph paper. The simple figure consisting of the plotted points is called a *scatter diagram*. From the scatter diagram, we can form a fairly good, though vague, idea of the relationship between  $X$  and  $Y$ . If the points are dense or closely packed, we may conclude that  $X$  and  $Y$  are correlated. On the other hand if the points are widely scattered throughout the graph paper, we may conclude that  $X$  and  $Y$  are either not correlated or poorly correlated.

Further if the points in the scatter diagram appear to lie near a straight line, we assume that the RV's have *linear correlation*. If they cluster round a well defined curve other than a straight line, the RV's are assumed to be *non-linear*. In this section we will assume linear correlation between the concerned RV's and discuss how to measure the degree of linear correlation.

## Correlation Coefficient

As the variance  $E\{X - E(X)\}^2$  measures the variations of the RV  $X$  from its mean value  $E(X)$ , the quantity  $E\{[X - E(X)][Y - E(Y)]\}$  measures the simultaneous variation of two RV's  $X$  and  $Y$  from their respective means and hence it is called the *covariance* of  $X$ ,  $Y$  and denoted as  $\text{Cov}(X, Y)$ .

$\text{Cov}(X, Y) = E\{[X - E(X)][Y - E(Y)]\}$  is also called the *product moment* of  $X$  and  $Y$  and also denoted as  $p(X, Y)$ .

Though  $p(X, Y)$  is a useful measure of the degree of correlation between  $X$  and  $Y$ , it is to be expressed in mixed units of  $X$  and  $Y$ . To avoid this difficulty and to express the degree of correlation in absolute units, we divide  $p(X, Y)$  by  $\sigma_x \cdot \sigma_y$ , so that  $\frac{p(x, y)}{\sigma_x \sigma_y}$  is a mere number, free from units of  $X$  and  $Y$ .

$\frac{p(x, y)}{\sigma_x \sigma_y}$  is a measure of intensity of linear relationship between  $X$  and  $Y$  and is called *Karl Pearson's Product Moment Correlation Coefficient* or simply *correlation coefficient* between  $X$  and  $Y$ . It is denoted by  $r(X, Y)$  or  $r_{XY}$  or simply  $r$ .

$$\text{Thus, } r_{XY} = \frac{E\{[X - E(X)][Y - E(Y)]\}}{\sqrt{E\{X - E(X)\}^2 E\{Y - E(Y)\}^2}} \quad (1)$$

since  $\sigma_x$ , the standard deviation of  $X$  is the positive square root of the variance of  $X$ .

$$\begin{aligned} \text{Now, } E[\{X - E(X)\}\{Y - E(Y)\}] \\ &= E[XY] - E(Y) \cdot X - E(X) \cdot Y + E(X) \cdot E(Y) \\ &= E(XY) - E(Y) \cdot E(X) - E(X) \cdot E(Y) + E(X) \cdot E(Y) \\ &\quad [\because E(X) \text{ and } E(Y) \text{ are non-random constants}] \\ &= E(XY) - E(X) \cdot E(Y) \end{aligned} \quad (2)$$

$$\text{Also we know that } E\{X - E(X)\}^2 = E(X^2) - \{E(X)\}^2 \quad (3)$$

$$\text{and } E\{Y - E(Y)\}^2 = E(Y^2) - \{E(Y)\}^2 \quad (4)$$

Using (2), (3) and (4) in (1), we get

$$r_{XY} = \frac{E(XY) - E(X) \cdot E(Y)}{\sqrt{\{E(X^2) - E^2(X)\}\{E(Y^2) - E^2(Y)\}}} \quad (5)$$

where  $E^2(X)$  means  $\{E(X)\}^2$ .

We will mainly deal with linear correlation of discrete RV's  $X$  and  $Y$ .  $X$  will take the values  $x_1, x_2, \dots, x_n$  with frequency 1 each and  $Y$  will simultaneously take the values  $y_1, y_2, \dots, y_n$  with frequency 1 each. Hence,  $E(X) = \frac{1}{n} \sum x_i$ ;  $E(X^2) = \frac{1}{n} \sum x_i^2$ ,  $E(XY) = \frac{1}{n} \sum x_i y_i$  etc. Using these values in (5), the working formula for the computation of  $r_{XY}$  is got as

$$r_{XY} = \frac{\frac{1}{n} \sum x_i y_i - \frac{1}{n} \sum x_i \frac{1}{n} \sum y_i}{\sqrt{\left\{ \frac{1}{n} \sum x_i^2 - \left( \frac{1}{n} \sum x_i \right)^2 \right\} \left\{ \frac{1}{n} \sum y_i^2 - \left( \frac{1}{n} \sum y_i \right)^2 \right\}}} \quad (6)$$

or  $r_{XY} = \frac{n \sum xy - \sum x \cdot \sum y}{\sqrt{\left\{ n \sum x^2 - (\sum x)^2 \right\} \left\{ n \sum y^2 - (\sum y)^2 \right\}}}$  (7)

## Properties of Correlation Coefficient

- $-1 \leq r_{XY} \leq 1$  or  $|\text{Cov}(X, Y)| \leq \sigma_X \cdot \sigma_Y$ .

Let us consider

$$E[a\{X - E(X)\} + \{Y - E(Y)\}]^2 = a^2 \sigma_x^2 + 2a C_{XY} + \sigma_y^2 \quad (1)$$

The R.H.S. expression is a quadratic expression in  $a$ , that is  $a$  real quantity. It is positive, as it is the expected value of a perfect square.

Hence, by the property of quadratic expressions, the discriminant of the R.H.S.  $\leq 0$

$$\text{i.e., } 4C_{XY}^2 - 4\sigma_x^2 \sigma_y^2 \leq 0$$

$$\text{i.e., } C_{XY}^2 \leq \sigma_x^2 \cdot \sigma_y^2 \quad (2)$$

$$\text{i.e., } \frac{C_{XY}^2}{\sigma_x^2 \cdot \sigma_y^2} \leq 1$$

$$\text{i.e., } r_{XY}^2 \leq 1$$

$$\text{i.e., } |r_{XY}| \leq 1 \text{ or } -1 \leq r_{XY} \leq 1$$

From step (2), it is clear that  $|C_{XY}| \leq \sigma_X \cdot \sigma_Y$

### Note

When  $0 < r_{XY} \leq 1$ , the correlation between  $X$  and  $Y$  is said to be positive or direct.  
When  $-1 \leq r_{XY} \leq 0$ , the correlation is said to be negative or inverse.

When  $-1 \leq r_{XY} \leq -0.5$  or  $0.5 \leq r_{XY} \leq 1$ , the correlation is assumed to be high, otherwise the correlation is assumed to be poor.

2. Correlation coefficient is independent of change of origin and scale.

i.e., if  $U = \frac{X-a}{h}$  and  $V = \frac{Y-b}{k}$ , where  $h, k > 0$ , then  $r_{XY} = r_{UV}$ .

By the transformations,  $X = a + hU$  and  $Y = b + kV$

$$\therefore E(X) = a + hE(U) \text{ and } E(Y) = b + kE(V)$$

$$\therefore X - E(X) = h(U - E(U)) \text{ and } Y - E(Y) = k(V - E(V))$$

$$\text{Then } C_{XY} = E[h\{U - E(U)\} \cdot k\{V - E(V)\}] = hk C_{UV}$$

$$\sigma_X^2 = E[h^2\{U - E(U)\}^2] = h^2 \sigma_U^2$$

$$\sigma_Y^2 = E[k^2\{V - E(V)\}^2] = k^2 \sigma_V^2$$

$$\begin{aligned} \therefore r_{XY} &= \frac{\text{Cov}(X, Y)}{\sqrt{\sigma_X^2 \cdot \sigma_Y^2}} = \frac{hk C_{UV}}{\sqrt{h^2 \cdot \sigma_U^2 \cdot k^2 \sigma_V^2}} \\ &= \frac{C_{UV}}{\sigma_U \cdot \sigma_V} = r_{UV} \end{aligned}$$

**Note** If  $X$  and  $Y$  take considerably large values, computation of  $r_{XY}$  will become difficult. In such problems, we may introduce change of origin and scale and compute  $r$  using the above property.

3. Two independent RV's  $X$  and  $Y$  are uncorrelated, but two uncorrelated RV's need not be independent.

When  $X$  and  $Y$  are independent,  $E(XY) = E(X) \cdot E(Y)$ .

$$\therefore C_{XY} = 0 \text{ and hence } r_{XY} = 0$$

viz.,  $X$  and  $Y$  are uncorrelated.

The converse is not true, since  $E(XY) = E(X) \cdot E(Y)$ , when  $r_{XY} = 0$ .

This does not imply that  $X$  and  $Y$  are independent, as  $X$  and  $Y$  are independent only when  $f(x, y) = f_X(x) \cdot f_Y(y)$ .

**Note** When  $E(XY) = 0$ ,  $X$  and  $Y$  are said to be orthogonal RV's.

$$4. r_{XY} = \frac{\sigma_X^2 + \sigma_Y^2 - \sigma_{(X-Y)}^2}{2\sigma_X \sigma_Y}$$

Let  $Z = X - Y$ . Then  $E(Z) = E(X) - E(Y)$

$$\therefore Z - E(Z) = [X - E(X)] - [Y - E(Y)]$$

$$\begin{aligned} \therefore \sigma_Z^2 &= E[Z - E(Z)]^2 = E[(X - E(X)) - (Y - E(Y))]^2 \\ &= E\{X - E(X)\}^2 + E\{Y - E(Y)\}^2 - 2E[(X - E(X))(Y - E(Y))] \\ &= \sigma_X^2 + \sigma_Y^2 - 2C_{XY} \end{aligned}$$

$$\therefore C_{XY} = \frac{\sigma_x^2 + \sigma_y^2 - \sigma_z^2}{2}$$

$$\therefore r_{XY} = \frac{C_{XY}}{\sigma_x \sigma_y} = \frac{\sigma_x^2 + \sigma_y^2 - \sigma_{(X-Y)}^2}{2\sigma_x \sigma_y}$$

Similarly, we can prove that

$$\sigma_{(X+Y)}^2 = \sigma_x^2 + \sigma_y^2 + 2C_{XY}$$

$$\text{and hence, } r_{XY} = \frac{\sigma_{(X+Y)}^2 - \sigma_x^2 - \sigma_y^2}{2\sigma_x \sigma_y}$$

## Rank Correlation Coefficient

Sometimes the actual numerical values of  $X$  and  $Y$  may not be available, but the positions of the actual values arranged in order of merit (ranks) only may be available. The ranks of  $X$  and  $Y$  will in general, be different and hence may be considered as random variables. Let them be denoted by  $U$  and  $V$ . The correlation coefficient between  $U$  and  $V$  is called *the rank correlation coefficient* between (the ranks of)  $X$ ,  $Y$  and denoted by  $\rho_{XY}$ .

Let us now derive a formula for  $\rho_{XY}$  or  $r_{UV}$ . Since  $U$  represents ranks of  $n$  values of  $X$ ,  $U$  takes the values  $1, 2, 3, \dots, n$ .

Similarly,  $V$  takes the same values  $1, 2, 3, \dots, n$  in a different order.

$$E(U) = E(V) = \frac{1}{n}(1 + 2 + \dots + n) = \frac{n+1}{2}$$

$$E(U^2) = E(V^2) = \frac{1}{n}(1^2 + 2^2 + \dots + n^2) = \frac{(n+1)(2n+1)}{6}$$

$$\begin{aligned}\therefore \sigma_U^2 &= \sigma_V^2 = E(U^2) - E^2(U) \\ &= \frac{(n+1)(2n+1)}{6} - \frac{(n+1)^2}{4} \\ &= \frac{(n+1)}{12}\{2(2n+1) - 3(n+1)\} \\ &= \frac{n^2 - 1}{12}\end{aligned}$$

Let  $D = U - V \quad \therefore E(D) = 0$

and  $\sigma_D^2 = E(D^2)$

By Property (4) given above,

$$\rho_{XY} = r_{UV} = \frac{\sigma_U^2 + \sigma_V^2 - \sigma_D^2}{2\sigma_U \sigma_V}, \text{ where } D = U - V$$

$$\begin{aligned}
 &= \frac{\left(\frac{n^2 - 1}{6}\right) - \sigma_D^2}{2\left(\frac{n^2 - 1}{12}\right)} \\
 &= 1 - \frac{6}{n^2 - 1} \sigma_D^2 \text{ or } 1 - \frac{6E(D^2)}{n^2 - 1} \\
 &= 1 - \frac{6\Sigma d^2}{n(n^2 - 1)} \quad \left[ \because E(D^2) = \frac{1}{n} \Sigma d^2 \right]
 \end{aligned}$$

**Note** The formula for the rank correlation coefficient is known as Spearman's formula. The values of  $r_{XY}$  and  $\rho_{XY}$  (or  $r_{UV}$ ) will be, in general, different.

### Worked Example 4(B)

#### Example 1

Compute the coefficient of correlation between  $X$  and  $Y$ , using the following data:

$X$ :	1	3	5	7	8	10
$Y$ :	8	12	15	17	18	20

$x_i$	$y_i$	$x_i^2$	$y_i^2$	$x_i y_i$
1	8	1	64	8
3	12	9	144	36
5	15	25	225	75
7	17	49	289	119
8	18	64	324	144
10	20	100	400	200
34	90	248	1446	352

Thus,  $n = 6$

$$\sum x_i = 34, \sum y_i = 90$$

$$\sum x_i^2 = 248, \sum y_i^2 = 1446$$

$$\sum x_i y_i = 582$$

$$\begin{aligned}
 r_{XY} &= \frac{n \sum xy - \sum x \cdot \sum y}{\sqrt{\{n \sum x^2 - (\sum x)^2\} \{n \sum y^2 - (\sum y)^2\}}} \\
 &= \frac{6 \times 582 - 34 \times 90}{\sqrt{\{6 \times 248 - (34)^2\} \{6 \times 1446 - (90)^2\}}}
 \end{aligned}$$

$$= \frac{432}{\sqrt{332 \times 576}} = 0.9879$$

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### Example 2

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Compute the coefficients of correlation between  $X$  and  $Y$  using the following data:

$X:$	65	67	66	71	67	70	68	69
$Y:$	67	68	68	70	64	67	72	70

We effect change of origin in respect of both  $X$  and  $Y$ . The new origins are chosen at or near the average of extreme values. Thus, we take  $\frac{65+71}{2} = 68$  as the new origin for  $X$  and  $\frac{64+72}{2} = 68$  as the new origin for  $Y$ . viz., we put  $u_i = (x_i - 68)$  and  $v_i = y_i - 68$  and find  $r_{UV}$ .

$X = x_i$	$Y = y_i$	$u_i = x_i - 68$	$v_i = y_i - 68$	$u_i^2$	$v_i^2$	$u_i v_i$
65	67	-3	-1	9	1	3
67	68	-1	0	1	0	0
66	68	-2	0	4	0	0
71	70	3	2	9	4	6
67	64	-1	-4	1	16	4
70	67	2	-1	4	1	-2
68	72	0	4	0	16	0
69	70	1	2	1	1	2
	Total	-1	2	29	39	13

$$\begin{aligned}
 r_{XY} = r_{UV} &= \frac{n \sum uv - \sum u \cdot \sum v}{\sqrt{\{n \sum u^2 - (\sum u)^2\} \{n \sum v^2 - (\sum v)^2\}}} \\
 &= \frac{8 \times 13 - (-1) \times 2}{\sqrt{(8 \times 29 - 1)(8 \times 39 - 4)}} = \frac{106}{\sqrt{231 \times 308}} = 0.3974
 \end{aligned}$$

---

### Example 3

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Find the coefficient of correlation between  $X$  and  $Y$  using the following data:

$X:$	5	10	15	20	25
$Y:$	16	19	23	26	30

As the values of  $X$  are in arithmetic progression, we make the change of origin and scale, by choosing the middle most value 15 as the new origin and the common difference 5 as the new scale.

i.e., we put  $U = \frac{X - 15}{5}$

As the values of  $Y$  are not in A.P., we are content with effecting a change of origin only, i.e., we put  $V = Y - \left(\frac{30+16}{2}\right) = Y - 23$ .

$x$	$y$	$u = \frac{x - 15}{5}$	$v = y - 23$	$u^2$	$v^2$	$uv$
5	16	-2	-7	4	49	14
10	19	-1	-4	1	16	4
15	23	0	0	0	0	0
20	26	1	3	1	9	3
25	30	2	7	4	49	14
	Total	0	-1	10	123	35

$$\begin{aligned}
 r_{XY} = r_{UV} &= \frac{n\sum uv - \sum u \cdot \sum v}{\sqrt{\{n\sum u^2 - (\sum u)^2\} \{n\sum v^2 - (\sum v)^2\}}} \\
 &= \frac{5 \times 35 - 0 \times (-1)}{\sqrt{5 \times 10 - 0)(5 \times 125 - 1)}} \\
 &= \frac{175}{\sqrt{50 \times 624}} = 0.9907
 \end{aligned}$$

#### Example 4

The following table gives the bivariate frequency distribution of marks in an intelligence test obtained by 100 students according to their ages:

Age ( $x$ ) in yrs Marks ( $y$ )	18	19	20	21	Total
10–20	4	2	2	—	8
20–30	5	4	6	4	19
30–40	6	8	10	11	35
40–50	4	4	6	8	22
50–60	—	2	4	4	10
60–70	—	2	3	1	6
Total	19	22	31	28	100

Calculate the coefficient of correlation between age and intelligence.

Since the frequencies of various values of  $x$  and  $y$  are not equal to 1 each the formula for the computation of  $r_{XY}$  is taken with a slight modification as given below:

$$r_{XY} = r_{UV} = \frac{N \sum f_{XY}uv - \sum f_X u \cdot \sum f_Y v}{\sqrt{\{N \sum f_X u^2 - (\sum f_X u)^2\} \{N \sum f_Y v^2 - (\sum f_Y v)^2\}}} \quad (1)$$

where  $u = x - 20$ ,  $v = \frac{y - 35}{10}$ ,  $f_X$  represents frequencies of  $X$ -distribution,  $f_Y$  represents frequencies of  $Y$ -distribution and  $f_{XY}$  are the cell frequencies.

Mid y/mid x	18	19	20	21	$f_Y$	$v$	$f_Y v$	$f_Y v^2$	$f_{XY}uv$
15	4	2	2	—	8	-2	-16	32	20
25	5	4	6	4	19	-1	-19	19	10
35	6	8	10	11	35	0	0	0	0
45	4	4	6	8	22	1	22	22	-4
55	—	2	4	4	10	2	20	40	4
65	—	2	3	1	6	3	18	54	-3
$f_X$	19	22	31	28	100	Total	5	167	27
$u$	-2	-1	0	1	Total				
$f_X u$	-38	-22	0	28					
$f_X u^2$	76	22	0	28					
$\Sigma f_{XY}uv$	18	-6	0	15					

**Note**  $\Sigma f_{XY}uv$  for the first row of the table is computed as follows.

$$\begin{aligned}\Sigma f_{XY}uv &= f_{11} u_1 v_1 + f_{12} u_2 v_1 + f_{13} u_3 v_1 + f_{14} u_4 v_1 \\ &= 4(-2) (-2) + 2(-1) (-2) + 2(0) + 0(1) (-2) \\ &= 20\end{aligned}$$

Similarly, other  $\Sigma f_{XY}uv$  values are computed. Value of  $(\Sigma \Sigma f_{XY}uv)$  obtained as the total of the entries of the last column and as that of the last row must tally.

Using the relevant values obtained in the table in (1), we have

$$\begin{aligned}r_{XY} &= \frac{100 \times 27 - (-32) \times 5}{\sqrt{\{100 \times 126 - (-32)^2\} \{100 \times 167 - 5^2\}}} \\ &= \frac{2860}{\sqrt{13624 \times 16675}} = 0.1897\end{aligned}$$

### Example 5

Calculate the correlation coefficient for the following ages of husbands ( $X$ ) and wives ( $Y$ ), using only standard deviations of  $X$  and  $Y$ :

$$\begin{array}{cccccccccc} X: & 23 & 27 & 28 & 28 & 29 & 30 & 31 & 33 & 35 & 36 \\ Y: & 18 & 20 & 22 & 27 & 21 & 29 & 27 & 29 & 28 & 29 \end{array}$$

$x$	$y$	$u = x - 30$	$v = y - 24$	$u^2$	$v^2$	$d = x - y$	$d^2$
23	18	-7	-6	49	36	5	25
27	20	-3	-4	9	16	7	49
28	22	-2	-2	4	4	6	36
28	27	-2	3	4	9	1	1
29	21	-1	-3	1	9	8	64
30	29	0	5	0	25	1	1
31	27	1	3	1	9	4	16
33	29	3	5	9	25	4	16
35	28	5	4	25	16	7	49
36	29	6	5	36	25	7	49
	Total	0	10	138	174	50	306

$$\sigma_x^2 = \frac{1}{n} \sum u^2 - \left( \frac{1}{n} \sum u \right)^2 = \frac{1}{10} \times 138 = 13.8$$

$$\sigma_y^2 = \frac{1}{n} \sum v^2 - \left( \frac{1}{n} \sum v \right)^2 = \frac{1}{10} \times 174 - \left( \frac{10}{10} \right)^2 = 16.4$$

$$\sigma_{(X-Y)}^2 = \sigma_d^2 = \frac{1}{n} \sum d^2 - \left( \frac{1}{n} \sum d \right)^2 = \frac{1}{10} \times 306 - \left( \frac{50}{10} \right)^2 = 5.6$$

$$r_{XY} = \frac{\sigma_x^2 + \sigma_y^2 - \sigma_d^2}{2\sigma_x\sigma_y} = \frac{13.8 + 16.4 - 5.6}{2 \times \sqrt{13.8} \times \sqrt{16.4}}$$

$$= \frac{24.6}{30.0879} = 0.8176$$

### Example 6

If the independent random variables  $X$  and  $Y$  have the variances 36 and 16 respectively, find the correlation coefficient between  $(X + Y)$  and  $(X - Y)$ .

Let  $U = X + Y$  and  $V = X - Y$

$$E(U) = E(X) + E(Y); E(V) = E(X) - E(Y)$$

$$E(UV) = E(X^2 - Y^2) = E(X^2) - E(Y^2)$$

$$E(U^2) = E\{(X + Y)^2\} = E(X^2) + E(Y^2) + 2E(XY)$$

$$E(V^2) = E(X^2) + E(Y^2) - 2E(XY)$$

$$C_{UV} = E(UV) - E(U) \cdot E(V)$$

$$= E(X^2) - E(Y^2) - \{E^2(X) - E^2(Y)\}$$

$$= [E(X^2) - E^2(X)] - [E(Y^2) - E^2(Y)]$$

$$\begin{aligned}
&= \sigma_x^2 - \sigma_y^2 = 36 - 16 = 20 \\
\sigma_u^2 &= E(U^2) - E^2(U) \\
&= \{E(X^2) + E(Y^2) + 2E(XY)\} - \{E^2(X) + E^2(Y) + 2E(X) \cdot E(Y)\} \\
&= [E(X^2) - E^2(X)] + [E(Y^2) - E^2(Y)] + 2[E(XY) - E(X) \cdot E(Y)] \\
&= 36 + 16 + 2 \times 0 \\
&\quad [\because X \text{ and } Y \text{ are independent and hence uncorrelated}] \\
&= 52
\end{aligned}$$

Similarly,  $\sigma_v^2 = 52$

Now,  $r_{UV} = \frac{C_{UV}}{\sigma_u \cdot \sigma_v} = \frac{20}{52} = \frac{5}{13}$

### Example 7

If  $X$ ,  $Y$  and  $Z$  are uncorrelated RV's with zero means and standard deviations 5, 12 and 9 respectively and if  $U = X + Y$  and  $V = Y + Z$ , find the correlation coefficient between  $U$  and  $V$ .

$$E(X) = E(Y) = E(Z) = 0$$

$$\text{Var}(X) = E(X^2) - E^2(X) = 25 \therefore E(X^2) = 25$$

Similarly,  $E(Y^2) = 144$  and  $E(Z^2) = 81$

$X$  and  $Y$  are uncorrelated.

$$\therefore r_{XY} = 0, \text{ i.e., } E(XY) - E(X) \cdot E(Y) = 0$$

$$\therefore E(XY) = 0. \text{ Similarly } E(YZ) = 0; E(ZX) = 0$$

Now,  $E(U) = E(X + Y) = 0$  and  $E(V) = 0$

$$\begin{aligned}
E(U^2) &= E(X^2 + Y^2 + 2XY) \\
&= 25 + 144 + 2 \times 0 = 169
\end{aligned}$$

$$\begin{aligned}
E(V^2) &= E(Y^2 + Z^2 + 2YZ) \\
&= 144 + 81 + 2 \times 0 = 225
\end{aligned}$$

$$\therefore \sigma_u^2 = E(U^2) - E^2(U) = 169$$

and  $\sigma_v^2 = E(V^2) - E^2(V) = 225$

$$\begin{aligned}
E(UV) &= E\{(X + Y)(Y + Z)\} \\
&= E(XY) + E(XZ) + E(YZ) + E(Y^2)
\end{aligned}$$

$$= 0 + 0 + 0 + 144 = 144$$

$$r_{UV} = \frac{E(UV) - E(U) \cdot E(V)}{\sigma_U \cdot \sigma_V} = \frac{144}{13 \times 15} = \frac{48}{65}$$

### Example 8

If  $X$  and  $Y$  are two RV's with variances  $\sigma_X^2$  and  $\sigma_Y^2$  respectively, find the value of  $k$ , if  $U = X + kY$  and  $V = X + \frac{\sigma_X}{\sigma_Y} \cdot Y$  are uncorrelated.

$U$  and  $V$  are uncorrelated.

$$\therefore \text{Cov}(U, V) = 0$$

$$\text{i.e., } E(UV) - E(U) \cdot E(V) = 0$$

$$\text{i.e., } E\left\{(X + kY)\left(X + \frac{\sigma_X}{\sigma_Y} \cdot Y\right)\right\} - E(X + kY) \cdot E\left(X + \frac{\sigma_X}{\sigma_Y} Y\right) = 0$$

$$\begin{aligned} \text{i.e., } & E\left\{X^2 + k \frac{\sigma_X}{\sigma_Y} \cdot Y^2 + \left(k + \frac{\sigma_X}{\sigma_Y}\right) XY\right\} \\ & - \left[ \{E(X) + kE(Y)\} \left\{E(X) + \frac{\sigma_X}{\sigma_Y} E(Y)\right\} \right] = 0 \end{aligned}$$

$$\begin{aligned} \text{i.e., } & E(X^2) + k \frac{\sigma_X}{\sigma_Y} E(Y^2) + \left(k + \frac{\sigma_X}{\sigma_Y}\right) E(XY) \\ & - \left[ E^2(X) + k \frac{\sigma_X}{\sigma_Y} E^2(Y) + \left(k + \frac{\sigma_X}{\sigma_Y}\right) E(X) \cdot E(Y) \right] = 0 \end{aligned}$$

$$\begin{aligned} \text{i.e., } & \{E(X^2) - E^2(X)\} + \frac{k\sigma_X}{\sigma_Y} \{E(Y^2) - E^2(Y)\} \\ & + \left(k + \frac{\sigma_X}{\sigma_Y}\right) \{E(XY) - E(X) \cdot E(Y)\} = 0 \end{aligned}$$

$$\text{i.e., } \sigma_X^2 + k \frac{\sigma_X}{\sigma_Y} \cdot \sigma_Y^2 + \left(k + \frac{\sigma_X}{\sigma_Y}\right) \text{Cov}(X, Y) = 0$$

Dividing throughout by  $\sigma_X$ ,

$$(\sigma_X + k\sigma_Y) + (\sigma_X + k\sigma_Y) \cdot \frac{\text{cov}(X, Y)}{\sigma_X \sigma_Y} = 0$$

$$\text{i.e. } (\sigma_X + k\sigma_Y) + (1 + r_{XY}) = 0$$

Assuming that  $r_{XY} \neq -1$ , we get

$$\sigma_X + k\sigma_Y = 0$$

$$\therefore k = -\frac{\sigma_X}{\sigma_Y}$$

### Example 9

If  $(X, Y)$  is a two-dimensional RV uniformly distributed over the triangular region  $R$  bounded by  $y = 0$ ,  $x = 3$  and  $y = 4/3 x$ . Find  $f_X(x)$ ,  $f_Y(y)$ ,  $E(X)$ ,  $\text{Var}(X)$ ,  $E(Y)$ ,  $\text{Var}(Y)$  and  $\rho_{XY}$ .

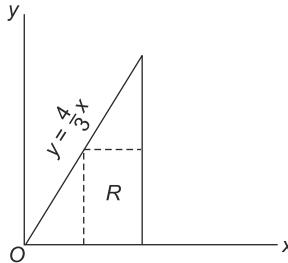


Fig. 4.4

Since  $(X, Y)$  is uniformly distributed,  $f(x, y) = \text{a constant} = k$

$$\iint f(x, y) dx dy = 1$$

$$\text{i.e., } \int_0^4 \int_{3y/4}^3 k dx dy = 1$$

$$\text{i.e., } k \int_0^4 \left( 3 - \frac{3y}{4} \right) dy = 1$$

$$\text{i.e., } 6k = 1$$

$$\therefore k = \frac{1}{6}$$

$$f_Y(y) = \int_{3y/4}^3 \frac{1}{6} dx = \frac{1}{8}(4-y), 0 < y < 4$$

$$f_X(x) = \int_3^{4x/3} \frac{1}{6} dy = \frac{2}{9}x, 0 < x < 3$$

$$E(X) = \int xf_X(x) dx = \int_0^3 \frac{2}{9}x^2 dx = 2$$

$$E(Y) = \int yf_Y(y) dy = \int_0^4 \frac{y}{8} \times (4-y) dy = \frac{4}{3}$$

$$E(X^2) = \int_0^3 \frac{2}{9} \times x^3 \, dx = \frac{9}{2}$$

$$E(Y^2) = \int_0^4 \frac{y^2}{8} \times (4-y) \, dy = \frac{8}{3}$$

$$\text{Var}(X) = E(X^2) - \{E(X)\}^2 = \frac{9}{2} - 4 = \frac{1}{2}$$

$$\text{Var}(Y) = E(Y^2) - \{E(Y)\}^2 = \frac{8}{3} - \frac{16}{9} = \frac{8}{9}$$

$$E(XY) = \int_0^4 \int_{3y/4}^3 \frac{1}{6} xy \, dx \, dy$$

$$= \frac{3}{64} \int_0^4 (16 - y^2) y \, dy = 3$$

$$\rho_{XY} = \frac{E(XY) - E(X) \times E(Y)}{\sigma_x \sigma_y} = \frac{\frac{3}{2} - 2 \times \frac{4}{3}}{\frac{1}{\sqrt{2}} \times 2 \sqrt{\frac{2}{3}}} = \frac{1}{2}$$

### Example 10

Find the correlation co-efficient between  $X$  and  $Y$ , which are jointly normally distributed with

$$f(x, y) = \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)} \left( \frac{x^2}{\sigma_x^2} - \frac{2rxy}{\sigma_x\sigma_y} + \frac{y^2}{\sigma_y^2} \right)\right\}$$

$$\frac{x^2}{\sigma_x^2} - \frac{2rxy}{\sigma_x\sigma_y} + \frac{y^2}{\sigma_y^2} = \left( \frac{x}{\sigma_x} - \frac{ry}{\sigma_y} \right)^2 + (1-r^2) \frac{y^2}{\sigma_y^2}$$

$$= \frac{1}{\sigma_x^2} \left( x - \frac{ry\sigma_x}{\sigma_y} \right)^2 + (1-r^2) \frac{y^2}{\sigma_y^2}$$

$$E(XY) = \frac{1}{2\pi\sigma_x\sigma_y\sqrt{1-r^2}} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \exp \left\{ -\frac{1}{2(1-r^2)} \left[ \frac{1}{\sigma_x^2} \left( x - \frac{ry\sigma_x}{\sigma_y} \right)^2 + \frac{(1-r)^2 y^2}{\sigma_y^2} \right] \right\} xy \, dx \, dy$$

$$\begin{aligned}
 &= \frac{1}{\sigma_y \sqrt{2\pi}} \int_{-\infty}^{\infty} y \exp\left(\frac{-y^2}{2\sigma_y^2}\right) \int_{-\infty}^{\infty} \frac{x}{(\sigma_x \sqrt{1-r^2}) \sqrt{2\pi}} \\
 &\quad \exp\left\{ -\left( x - \frac{ry\sigma_x}{\sigma_y} \right)^2 \right\} dx dy
 \end{aligned} \tag{1}$$

The inner integral is the mean of the normal distribution with mean  $\frac{ry\sigma_x}{\sigma_y}$  and variance  $(1-r^2)\sigma_x^2$ .

$$\therefore \text{the inner integral} = \frac{ry\sigma_x}{\sigma_y}$$

Using this value in (1),

$$\begin{aligned}
 E(XY) &= \left( \frac{r\sigma_x}{\sigma_y} \right) \times \frac{1}{\sigma_y \sqrt{2\pi}} \int_{-\infty}^{\infty} y^2 \exp\left(-\frac{y^2}{2\sigma_y^2}\right) dy \\
 &= \frac{r\sigma_x}{\sigma_y} E(Y^2) \text{ for } N(0, \sigma_y) \\
 &= \frac{r\sigma_x}{\sigma_y} \times \sigma_y^2 \\
 &= r\sigma_X\sigma_Y \\
 \therefore \rho_{XY} &= \frac{E(XY) - E(X)E(Y)}{\sigma_x\sigma_y} = r
 \end{aligned}$$

### Example 11

Ten students got the following percentage of marks in Mathematics and Physical sciences:

Students:	1	2	3	4	5	6	7	8	9	10
Marks in										
Mathematics:	78	36	98	25	75	82	90	62	65	39
Marks in										
Phy. Sciences:	84	51	91	60	68	62	86	58	63	47

Calculate the rank correlation coefficient.

Denoting the ranks in Mathematics and in Phy. Sciences by  $U$  and  $V$ , we have the following values of  $U$  and  $V$ :

$U:$	4	9	1	10	5	3	2	7	6	8
$V:$	3	9	1	7	4	6	2	8	5	10
$D:$	1	0	0	3	1	-3	0	-1	1	-2
$D^2:$	1	0	0	9	1	9	0	1	1	4

$\therefore \Sigma d^2 = 26$

$$\rho_{XY} = r_{UV} = 1 - \frac{6\sum d^2}{n(n^2 - 1)}$$

$$= 1 - \frac{6 \times 26}{10 \times 99} = 0.8424$$

**Example 12**

Ten competitors in a beauty contest were ranked by three judges as follows:

Judges	Competitors									
	1	2	3	4	5	6	7	8	9	10
A:	6	5	3	10	2	4	9	7	8	1
B:	5	8	4	7	10	2	1	6	9	3
C:	4	9	8	1	2	3	10	5	7	6

Discuss which pair of judges have the nearest approach to common taste of beauty.

Rank by A (U)	Rank by B (V)	Rank by C (W)	$d_1 = U - V$	$d_2 = V - W$	$d_3 = U - W$	$d_1^2$	$d_2^2$	$d_3^2$
6	5	4	1	1	2	1	1	4
5	8	9	-3	-1	-4	9	1	16
3	4	8	-1	-4	-5	1	16	25
10	7	1	3	6	9	9	36	81
2	10	2	-8	8	0	64	64	0
4	2	3	2	-1	1	4	1	1
9	1	10	8	-9	-1	64	81	1
7	6	5	1	1	2	1	1	4
8	9	7	-1	2	1	1	4	1
1	3	6	-2	-3	-5	4	9	25
					Total:	158	214	158

$$r_{UV} = 1 - \frac{6\sum d_1^2}{n(n^2 - 1)} = 1 - \frac{6 \times 158}{10 \times 99} = 0.0424$$

$$r_{VW} = 1 - \frac{6\sum d_2^2}{n(n^2 - 1)} = 1 - \frac{6 \times 214}{10 \times 99} = -0.2970$$

$$r_{UW} = 1 - \frac{6\sum d_3^2}{n(n^2 - 1)} = 1 - \frac{6 \times 158}{10 \times 99} = 0.0424$$

Since  $r_{UV}$  is maximum, the judges A and B may be considered to have common taste of beauty to some extent compared to other pairs of judges.

### Exercise 4(B)

**Part-A (Short-answer Questions)**

1. What is the scatter diagram? What is its role in correlation analysis?
2. What do you mean by correlation between two random variables?
3. What is linear correlation? How will you find that two RV's are linearly correlated?
4. Define covariance of  $X, Y$  and coefficient of correlation between  $X$  and  $Y$ .
5. Why is  $r_{XY}$  preferred for measuring the degree of linear correlation to  $\text{Cov}(X, Y)$ ?
6. State the properties of correlation coefficient.
7. State two different formulas used to compute  $r_{XY}$ .
8. Define rank correlation coefficient and write down the formula for computing it.
9. Prove that  $-1 \leq r_{XY} \leq 1$ .
10. Prove that  $\sigma_{(X+Y)}^2 - \sigma_{(X-Y)}^2 = 4 \text{Cov}(X, Y)$
11. If  $C_{XY}$  is the covariance of  $X$  and  $Y$ , prove that  $C_{XY} = E(XY) - E(X) \cdot E(Y)$ .
12. If  $X$  and  $Y$  are independent RV's prove that  $r_{XY} = 0$ . Is the converse true?
13. If  $X$  and  $Y$  are uncorrelated, prove that  $\text{Var}(X + Y) = \text{Var}(X) + \text{Var}(Y)$ .
14. When are two RVs said to be orthogonal?

**Part-B**

15. Ten students got the following marks in Mathematics and Basic Engineering:

Marks in Mathematics } 78    36    98    25    75    82    90    62    65    39

Marks in Basic Engg. } 84    51    91    60    68    62    86    58    53    47

Calculate the coefficient of correlation.

16. Calculate the correlation coefficient between  $X$  and  $Y$  from the following data:

X:	65	66	67	67	68	69	70	72
Y:	67	68	65	68	72	72	69	71

17. Find the coefficient of correlation between  $X$  and  $Y$  using the following data:

X:	5.5	3.6	2.6	3.4	3.1	2.7	3.0	3.1	3.2	3.8
Y:	27	36	39	39	32	35	40	36	44	36

18. Compute the coefficient of correlation between  $X$  and  $Y$  from the following data:

X:	80	45	55	56	58	60	65	68	70	75	85
Y:	82	56	50	48	60	62	64	65	70	74	90

19. Find the coefficient of correlation between  $X$  and  $Y$  from the following data:

$$\begin{array}{llllll} X: & 10 & 14 & 18 & 22 & 26 & 30 \\ Y: & 18 & 12 & 24 & 6 & 30 & 36 \end{array}$$

20. Calculate the coefficient of correlation between  $X$  and  $Y$ , by finding variance only, from the following data:

$$\begin{array}{cccccccccc} X: & 21 & 23 & 30 & 54 & 57 & 58 & 72 & 78 & 87 & 90 \\ Y: & 60 & 71 & 72 & 83 & 110 & 84 & 100 & 92 & 113 & 135 \end{array}$$

21. Calculate  $r_{XY}$  from the following data, where  $X$  represents production (in crore tons) and  $Y$  represents exports (in crore tons), using only the variances.

$$\begin{array}{cccccccc} X: & 55 & 56 & 58 & 59 & 60 & 60 & 62 \\ Y: & 35 & 38 & 38 & 39 & 44 & 43 & 44 \end{array}$$

22. The following table gives the frequency of scores obtained by 65 students in a general knowledge test according to age groups. Measure the degree of linear relationship between age and general knowledge:

Test scores	Age in years			
	19	20	21	22
225	4	4	2	1
275	3	5	4	2
325	2	6	8	5
375	1	4	6	8

23. Compute the value of  $r_{XY}$  between  $X$ , the ages of husbands and  $Y$  the ages of wives from the following data:

$Y \setminus X$	15–25	25–35	35–45	45–55	55–65	65–75	Total
15–25	1	1	—	—	—	—	2
25–35	2	12	1	—	—	—	15
35–45	—	4	10	1	—	—	15
45–55	—	—	3	6	1	—	10
55–65	—	—	—	2	4	2	8
65–75	—	—	—	—	1	2	3
Total	3	17	14	9	6	4	53

24. Find the rank correlation coefficient between the ranks of the variable  $X$  and  $Y$ :

$$\begin{array}{llllllllll} X: & 10 & 15 & 12 & 17 & 13 & 16 & 24 & 14 & 22 \\ Y: & 30 & 42 & 45 & 46 & 33 & 34 & 40 & 35 & 39 \end{array}$$

25. The competitors in a musical contest were ranked by the three judges  $A$ ,  $B$ ,  $C$  in the following order:

Rank by A:1	6	5	10	3	2	4	9	7	8
Rank by B:3	5	8	4	7	10	2	1	6	9
Rank by C:6	4	9	8	1	2	3	10	5	7

- Using rank correlation technique, find which pair of judges have more or less the same taste in music.
26. If  $X, Y, Z$  are uncorrelated RVs having the same variance, find the correlation coefficient between  $(X + Y)$  and  $(Y + Z)$ .
  27. If  $X$  and  $Y$  are two correlated RVs with zero means, prove that  $U = X \cos \alpha + Y \sin \alpha$  and  $V = X \sin \alpha - Y \cos \alpha$  are also uncorrelated.
  28.  $X$  and  $Y$  are independent RV's with means 5 and 10 variances 4 and 9 respectively. Obtain the correlation coefficient between  $U$  and  $V$ , where  $U = 3X + 4Y$  and  $V = 3X - Y$ .
  29. If  $X_1, X_2, X_3$  are three correlated RV's having variances  $v_1, v_2, v_3$  respectively, obtain the coefficient of correlation between  $(X_1 + X_2)$  and  $(X_2 + X_3)$ .
  30. Show that (i)  $E\{aX + bY\} = aE(X) + bE(Y)$  and (ii)  $\text{var}(aX + bY) = a^2 \text{var}(X) + b^2 \text{var}(Y) + 2abC(X, Y)$ , where  $C(X, Y)$  is the covariance of  $(X, Y)$ .
  31. If two RV's are uncorrelated, prove that the variance of their sum is equal to the sum of their variances.
  32. If the joint density function of  $(X, Y)$  is given by  $f(x, y) = 2 - x - y$ ,  $0 \leq x, y \leq 1$ , find  $E(X), E(Y), \text{var}(X), \text{var}(Y)$  and  $r_{XY}$ .
  33. If the two dimensional RV  $(X, Y)$  is uniformly distributed in  $0 \leq x \leq y \leq 1$ , find  $E(X), E(Y), \text{var}(X), \text{var}(Y)$  and  $r_{XY}$ .
  34. If the two dimensional RV  $(X, Y)$  is uniformly distributed over  $R$ , where  $R$  is defined by  $\{(x, y)/x^2 + y^2 \leq 1, y \geq 0\}$ , find  $r_{XY}$ .
  35. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = x + y$ ,  $0 \leq x, y \leq 1$ , find  $r_{XY}$ .
  36. Let  $X$  be a RV with mean value = 3 and variance = 2. Find the second moment of  $X$  about the origin. Another RV  $Y$  is defined by  $Y = -6X + 22$ . Find the mean value of  $Y$  and the correlation of  $X$  and  $Y$ .

## Regression

When the random variables  $X$  and  $Y$  are linearly correlated, the points plotted on the scatter diagram, corresponding to  $n$  pairs of observed values of  $X$  and  $Y$ , will have a tendency to cluster round a straight line. This straight is called the *regression line*. The regression line can be taken as the best fitting straight line for the observed pairs of values of  $X$  and  $Y$  in the least square sense, with which the students are familiar.

When two RV's  $X$  and  $Y$  are linearly correlated, we may not know which variable takes independent values. If we treat  $X$  as the independent variable and hence assume that the values of  $Y$  depend on those of  $X$ , the regression line is called the *regression line of  $Y$  on  $X$* . If we assume that the values of  $X$  depend on those of the independent variable  $Y$ , the *regression line of  $X$  on  $Y$*  is obtained. Thus in situations where the distinction cannot be made between the RVs  $X$  and  $Y$  as to which is the independent variable and which is the dependent variable, there will be two regression lines. However, when the value of  $Y(X)$  is to be predicted corresponding to a specified value of  $X(Y)$ , we should make use of the regression line of  $Y(X)$  on  $X(Y)$ .

## Equation of the Regression Line of $Y$ on $X$

The regression line of  $Y$  on  $X$  is the best-fitting straight line for the observed pairs of values  $(x_1, y_1), (x_2, y_2), \dots, (x_n, y_n)$ , based on the assumption that  $x$  is the independent variable and  $y$  is the dependent variable. Hence, let the equation of the regression line of  $Y$  on  $X$  be assumed as  $y = ax + b$ . (1)

By the principle of least squares, the normal equations which give the values of  $a$  and  $b$ .

$$\text{are} \quad \sum y_i = a \sum x_i + nb \quad (2)$$

$$\text{and} \quad \sum x_i y_i = a \sum x_i^2 + b \sum x_i \quad (3)$$

Dividing the equation (2) by  $n$ , we get

$$\bar{y} = a \bar{x} + b \quad (4)$$

where  $\bar{x} = E(X)$  and  $\bar{y} = E(Y)$ . (1) – (4) gives the required equation as

$$y - \bar{y} = a(x - \bar{x}) \quad (5)$$

Eliminating  $b$  between equations (2) and (3),

$$\begin{aligned} \text{we get} \quad a &= \frac{n \sum x_i y_i - \sum x_i \cdot \sum y_i}{n \sum x_i^2 - (\sum x_i)^2} \\ &= \frac{\frac{1}{n} \sum x_i y_i - \left(\frac{1}{n} \sum x_i\right) \cdot \left(\frac{1}{n} \sum y_i\right)}{\frac{1}{n} \sum x_i^2 - \left(\frac{1}{n} \sum x_i\right)^2} \end{aligned}$$

$$\text{or} \quad a = \frac{E(XY) - E(X) \cdot E(Y)}{E(X^2) - E^2(X)} = \frac{P_{XY}}{\sigma_x^2} \quad (6)$$

Using (6) in (5), we get the equation of the regression line of  $Y$  on  $X$  as

$$y - \bar{y} = \frac{P_{XY}}{\sigma_x^2} (x - \bar{x}) \quad (7)$$

$$\begin{aligned} \text{or} \quad y - \bar{y} &= \frac{r_{XY} \sigma_Y}{\sigma_X} (x - \bar{x}) \quad (8) \\ &\left[ \because r_{XY} = \frac{P_{XY}}{\sigma_X \sigma_Y} \right] \end{aligned}$$

In a similar manner, assuming the equation of the regression line of  $X$  and  $Y$  as  $x = ay + b$  and using the equations

$$\sum x_i = a \sum y_i + nb \text{ and } \sum x_i y_i = a \sum y_i^2 + b \sum y_i,$$

we can get the equation of the regression line of  $X$  on  $Y$  as

$$x - \bar{x} = \frac{p_{XY}}{\sigma_Y^2} (y - \bar{y}) \quad (9)$$

or  $x - \bar{x} = \frac{r_{XY}\sigma_X}{\sigma_Y}(y - \bar{y}) \quad (10)$

### Note

1.  $\frac{p_{XY}}{\sigma_X^2}$  or  $\frac{r_{XY}\sigma_Y}{\sigma_X}$  is called the regression coefficient of  $Y$  on  $X$  and denoted by  $b_1$  or  $b_{YX}$ .

$\frac{p_{XY}}{\sigma_Y^2}$  or  $\frac{r_{XY}\sigma_X}{\sigma_Y}$  is called the regression coefficient of  $X$  on  $Y$  and denoted by  $b_2$  or  $b_{XY}$ .

2. Clearly,  $b_1 b_2 = r_{XY}^2$ , i.e.,  $r_{XY}$  is the geometric mean of  $b_1$  and  $b_2$ .

$$\therefore r_{XY} = \pm \sqrt{b_1 b_2}$$

The sign of  $r_{XY}$  is the same as that of  $b_1$  or  $b_2$ , as  $b_1 = r_{xy} \frac{\sigma_Y}{\sigma_X}$  and  $b_2 = r_{XY} \frac{\sigma_Y}{\sigma_X}$  have the same sign as  $r_{XY}$  ( $\because \sigma_X$  and  $\sigma_Y$  are positive).

$$\text{Also, } \frac{b_1}{b_2} = \frac{\sigma_Y^2}{\sigma_X^2}$$

3. When there is perfect linear correlation between  $X$  and  $Y$ , viz., when  $r_{XY} = \pm 1$ , the two regression lines coincide.

4. The point of intersection of the two regression lines is clearly the point whose co-ordinates are  $(\bar{x}, \bar{y})$ .

5. When there is no linear correlation between  $X$  and  $Y$ , viz., when  $r_{XY} = 0$ , the equations of the regression lines become  $y = \bar{y}$  and  $x = \bar{x}$ , which are at right angles.

## Standard Error of Estimate of $Y$

Although we use the regression line of  $Y$  on  $X$  to predict the value of  $Y$  corresponding to a specified value of  $X$  we may also use it to estimate the value of  $Y$  corresponding to an observed value of  $X = x_i$ , say. The value of  $Y$  estimated in this manner need not, in general, be equal to the corresponding observed value of  $Y$ , namely,  $y_i$ . Hence, the difference between  $Y$  and  $Y_E$  is called the error estimate of  $Y$ . This error will vary from one observed value to the other and a random variable. The standard deviation of this RV ( $Y - Y_E$ ) is called the standard error of estimate of  $Y$  and denoted by  $S_Y$ .

$$\begin{aligned}
 \text{Now, } E\{Y - Y_E\} &= E\left[Y - \left\{\bar{y} + \frac{r_{XY}\sigma_Y}{\sigma_X}(X - \bar{x})\right\}\right] \\
 &= (\bar{y} - \bar{y}) - \frac{r_{XY}\sigma_Y}{\sigma_X}(\bar{x} - \bar{x}) \\
 &= 0 \\
 \sigma_{(Y-Y_E)}^2 &= E\{(Y - Y_E)^2\} - E^2(Y - Y_E) \\
 &= E\left[Y - \left\{\bar{y} + \frac{r_{XY}\sigma_Y}{\sigma_X}(X - \bar{x})\right\}\right]^2 \\
 &= E\left[(Y - \bar{y})^2 + \frac{r_{XY}^2\sigma_Y^2}{\sigma_X^2}(X - \bar{x})^2 - \frac{2r_{XY}\sigma_Y}{\sigma_X}(X - \bar{x})(Y - \bar{y})\right] \\
 \text{(i.e.) } S_Y^2 &= \sigma_Y^2 + \frac{r_{XY}^2\sigma_Y^2}{\sigma_X^2}\sigma_X^2 - \frac{2r_{XY}\sigma_Y}{\sigma_X}\text{Cov}(X, Y) \\
 &= \sigma_Y^2 + r_{XY}^2 \cdot \sigma_Y^2 - 2r_{XY}\sigma_Y^2 \\
 &\quad [\because \text{Cov}(X, Y) = r_{XY}\sigma_X\sigma_Y] \\
 &= (1 - r_{XY}^2)\sigma_Y^2 \text{ or } S_Y = \sqrt{1 - r_{XY}^2}\sigma_Y \tag{1}
 \end{aligned}$$

Similarly, the standard error estimate of  $X$ , denoted by  $S_X$  is given by

$$S_X^2 = (1 - r_{XY}^2)\sigma_X^2 \text{ or } S_X = \sqrt{1 - r_{XY}^2}\sigma_X \tag{2}$$

**Note** We may use (1) or (2) to prove that  $|r_{XY}| \leq 1$ .

$$\text{From (1), } S_Y = \sqrt{1 - r_{XY}^2}\sigma_Y$$

Since  $S_Y$  and  $\sigma_Y$  are positive,  $1 - r_{XY}^2 \geq 0$

$$\therefore r_{XY}^2 \leq 1 \\ \text{i.e., } |r_{XY}| \leq 1 \text{ or } -1 \leq r_{XY} \leq 1]$$

### Worked Example 4(C)

#### Example 1

Obtain the equations of the lines of regression from the following data:

X:	1	2	3	4	5	6	7
Y:	9	8	10	12	11	13	14

$X$	$Y$	$U = X - 4$	$V = Y - 11$	$U^2$	$V^2$	$UV$
1	9	-3	-2	9	4	6
2	8	-2	-3	4	9	6
3	10	-1	-1	1	1	1
4	12	0	1	0	1	0
5	11	1	0	1	0	0
6	13	2	2	4	4	4
7	14	3	3	9	9	9
	Total	0	0	28	28	26

$$\bar{x} = E(X) = 4 + \frac{1}{n} \sum u = 4$$

$$\bar{y} = E(Y) = 11 + \frac{1}{n} \sum v = 11$$

$$\sigma_x^2 = \frac{1}{n} \sum u^2 - \left( \frac{1}{n} \sum u \right)^2 = \frac{1}{7} \times 28 = 4$$

$$\sigma_y^2 = \frac{1}{n} \sum v^2 - \left( \frac{1}{n} \sum v \right)^2 = \frac{1}{7} \times 28 = 4$$

$$C_{XY} = \frac{1}{n} \sum uv - \left( \frac{1}{n} \sum u \right) \cdot \left( \frac{1}{n} \sum v \right) = \frac{1}{7} \times 26 = 3.7$$

The regression line of  $Y$  on  $X$  is

$$y - \bar{y} = \frac{p_{XY}}{\sigma_x^2} (x - \bar{x})$$

$$\text{i.e., } y - 11 = \frac{3.7}{4} (x - 4)$$

$$\text{i.e., } 3.7x - 4y + 29.2 = 0$$

The regression line of  $X$  on  $Y$  is

$$x - \bar{x} = \frac{p_{XY}}{\sigma_x^2} (y - \bar{y})$$

$$\text{i.e., } x - 4 = \frac{3.7}{4} (y - 11)$$

$$\text{i.e., } 4x - 3.7y + 24.7 = 0$$

### Example 2

Obtain the equations of the regression lines from the following data, using the method of least squares. Hence, find the coefficient of correlation between  $X$  and  $Y$ . Also estimate the value of (a)  $Y$ , when  $X = 38$ , and (b)  $X$ , when  $Y = 18$ .

$X:$	22	26	29	30	31	31	34	35
$Y:$	20	20	21	29	27	24	27	31

Put  $U = X - 29$  and  $V = Y - 27$ .

Let the equation of the regression line of  $Y$  on  $X$  be  $y = Ax + B$  or equivalently  $v = au + b$  (1)

The normal equations for finding  $a$  and  $b$  are

$$a \sum u + nb = \sum v \quad (2)$$

$$\text{and } a \sum u^2 + b \sum v = \sum uv \quad (3)$$

$x$	$y$	$u = x - 29$	$v = y - 27$	$u^2$	$v^2$	$uv$
22	20	-7	-7	49	49	49
26	20	-3	-7	9	49	21
29	21	0	-6	0	36	0
30	29	1	2	1	4	2
31	27	2	0	4	0	0
31	24	2	-3	4	9	-6
34	27	5	0	25	0	0
35	31	6	4	36	16	24
	Total	6	-17	128	163	90

Using the relevant values from the table in (2) and (3), we have

$$6a + 8b = -17 \quad (2)'$$

$$128a + 6b = 90 \quad (3)'$$

Solving (2)' and (3)', we get

$$a = 0.83; b = -2.75$$

Hence, the regression line of  $Y$  on  $X$  is

$$y - 27 = 0.83(x - 29) - 2.75$$

$$\text{i.e., } y = 0.83x + 0.18 \quad (4)$$

Let the equation of the regression line of  $X$  on  $Y$  be  $x = Cy + D$  or equivalently  $u = cv + d$  (5)

The normal equations for finding  $c$  and  $d$  are

$$c \sum v + nd = \sum u \quad (6)$$

$$\text{and } c \sum v^2 + d \sum uv = \sum u v \quad (7)$$

Using the relevant values from the table in (6) and (7), we have

$$-17c + 8d = 6 \quad (6)'$$

$$163c - 17d = 90 \quad (7)'$$

Solving (6)' and (7)', we get

$$c = 0.81; d = 2.47$$

Hence, the regression line of  $X$  on  $Y$  is

$$x - 29 = 0.81(y - 27) + 2.47$$

$$\text{i.e.,} \quad x = 0.81y + 9.60 \quad (8)$$

Comparing the equation (4) with

$$y - \bar{y} = r \frac{\sigma_Y}{\sigma_X} (x - \bar{x})$$

$$\text{We get} \quad r \frac{\sigma_Y}{\sigma_X} = 0.83 \quad (9)$$

Comparing the equation (8) with

$$x - \bar{x} = r \frac{\sigma_X}{\sigma_Y} (y - \bar{y})$$

$$\text{We get} \quad r \frac{\sigma_X}{\sigma_Y} = 0.81 \quad (10)$$

From (9) and (10), we get  $r^2 = 0.83 \times 0.81$

$$\therefore r = 0.82 \quad \left( \because b_1 = \frac{r\sigma_Y}{\sigma_X} \text{ and } b_2 = \frac{r\sigma_X}{\sigma_Y} \text{ are both positive} \right)$$

We use the equation (4) to estimate the value of  $Y$  when  $X = 38$ .

$$\therefore Y = 0.83 \times 38 + 0.18 = 31.72$$

Using the equation (8) to estimate the value of  $X$  when  $Y = 18$ , we have

$$X = 0.81 \times 18 + 9.60 = 24.18$$

### Example 3

A study of prices of rice at Chennai and Madurai gave the following data:

	<i>Chennai</i>	<i>Madurai</i>
Mean	19.5	17.75
S.D.	1.75	2.5

Also the coefficient of correlation between the two is 0.8. Estimate the most likely price of rice (a) at Chennai corresponding to the price of 18 at Madurai, and (b) at Madurai corresponding to the price of 17 at Chennai.

Let the prices of rice at Chennai and Madurai be denoted by  $X$  and  $Y$  respectively. Then from the data,

$$\bar{x} = 19.5, \bar{y} = 17.75, \sigma_x = 1.75, \sigma_y = 2.5 \text{ and } r_{XY} = 0.8.$$

Regression line of  $X$  on  $Y$  is

$$x - \bar{x} = \frac{r\sigma_x}{\sigma_y}(y - \bar{y})$$

$$\text{i.e., } x - 19.5 = \frac{0.8 \times 1.75}{2.5}(y - 17.75)$$

$\therefore$  when  $y = 18$ ,

$$\begin{aligned} x &= 19.5 + \frac{0.8 \times 1.75}{2.5} \times 0.25 \\ &= 19.64 \end{aligned}$$

Regression line of  $Y$  on  $X$  is

$$y - \bar{y} = \frac{r\sigma_y}{\sigma_x}(x - \bar{x})$$

$$\text{i.e., } y - 17.75 = \frac{0.8 \times 2.5}{1.75}(x - 19.5)$$

$\therefore$  when  $x = 17$ ,

$$\begin{aligned} y &= 17.75 + \frac{0.8 \times 2.5}{1.75} \times (-2.5) \\ &= 14.89 \end{aligned}$$

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### Example 4

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In a partially destroyed laboratory record of an analysis of correlation data, the following results only are legible: Variable of  $X = 1$ . The regression equations are  $3x + 2y = 26$  and  $6x + y = 31$ . What were (a) the mean values of  $X$  and  $Y$ ? (b) the standard deviation of  $Y$ ? and (c) the correlation coefficient between  $X$  and  $Y$ ?

- Since the lines of regression intersect at  $(\bar{x}, \bar{y})$  we have  $3\bar{x} + 2\bar{y} = 26$  and  $6\bar{x} + \bar{y} = 31$   
Solving these equations, we get  $\bar{x} = 4$  and  $\bar{y} = 7$ .
- Which of the two equations is the regression equation of  $Y$  on  $X$  and which one is the regression equation of  $X$  on  $Y$  are not known.

Let us tentatively assume that the first equation is the regression line of  $X$  on  $Y$  and the second equation is the regression line of  $Y$  on  $X$ . Based on this assumption, the first equation can be re-written as

$$x = -\frac{2}{3}y + \frac{26}{3} \quad (1)$$

and the other as  $y = -6x + 31$  (2)

Then  $b_{XY} = -\frac{2}{3}$  and  $b_{YX} = -6$

$$\therefore r_{XY}^2 = b_{XY} \times b_{YX} = 4$$

$$\therefore r_{XY} = -2, \text{ which is absurd.}$$

Hence, our tentative assumption is wrong.

$\therefore$  the first equation is the regression line of  $Y$  on  $X$  and re-written as

$$y = -\frac{3}{2}x + 13 \quad (3)$$

The second equation is the regression line of  $X$  on  $Y$  and re-written as

$$x = -\frac{1}{6}y + \frac{31}{6} \quad (4)$$

Hence, the correct  $b_{YX} = -\frac{3}{2}$  and the correct  $b_{XY} = -\frac{1}{6}$

$$\therefore r_{XY}^2 = b_{YX} \cdot b_{XY} = \frac{1}{4}$$

$$\therefore r_{XY} = -\frac{1}{2} \quad (\because \text{both } b_{YX} \text{ and } b_{XY} \text{ are negative})$$

$$(c) \quad \text{Now} \quad \frac{\sigma_Y^2}{\sigma_X^2} = \frac{b_{YX}}{b_{XY}} = \frac{-\frac{3}{2}}{-\frac{1}{6}} = 9$$

$$\therefore \sigma_Y^2 = 9 \times \sigma_X^2 = 9$$

$$\therefore \sigma_Y = 3$$

---

### Example 5

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Given that  $x = 4y + 5$  and  $y = kx + 4$  are regression lines of  $X$  on  $Y$  and of  $Y$  on  $X$  respectively, show that  $0 \leq k \leq \frac{1}{4}$ . If  $k = \frac{1}{16}$ , find the means of  $X$  and  $Y$  and  $r_{XY}$ . From the given equations, we note that

$$b_{YX} = k \text{ and } b_{XY} = 4$$

$$r_{XY}^2 = b_{XY} \cdot b_{YX} = 4k$$

Since  $0 \leq r_{XY}^2 \leq 1$ , we get  $0 \leq 4k \leq 1$

$$\therefore 0 \leq k \leq \frac{1}{4}$$

When  $k = \frac{1}{16}$ ,  $r_{XY}^2 = \frac{1}{4}$

$$\therefore r_{XY} = \pm \frac{1}{2}$$

But both  $b_{YX}$  and  $b_{XY}$  are positive.

$$\therefore r_{XY} = \frac{1}{2}$$

When  $k = \frac{1}{16}$ , the regression equations become

$$x = 4y + 5 \quad (1)$$

$$\text{and } y = \frac{1}{16}x + 4 \quad (2)$$

Solving equations (1) and (2), we get

$$x = 28 \text{ and } y = 5.75$$

$$\therefore \bar{x} = 28 \text{ and } \bar{y} = 5.75$$

---

### Example 6

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Find the angle between the two lines of regression. Deduce the condition for the two lines to be (a) at right angles, and (b) coincident.

The equations of the regression lines

$$\text{are } y - \bar{y} = r \frac{\sigma_Y}{\sigma_X} (x - \bar{x}) \quad (1)$$

$$\text{and } x - \bar{x} = r \frac{\sigma_X}{\sigma_Y} (y - \bar{y}) \quad (2)$$

Slope of line (1) =  $r \frac{\sigma_Y}{\sigma_X} = m_1$ , say

Slope of line (2) =  $\frac{\sigma_Y}{r \sigma_X} = m_2$ , say

If  $\theta$  is the acute angle between the two lines, then

$$\tan \theta = \frac{|m_1 - m_2|}{1 + m_1 m_2}$$

$$\begin{aligned}
 &= \frac{\left| r \frac{\sigma_Y}{\sigma_X} - \frac{\sigma_Y}{r \sigma_X} \right|}{1 + \frac{\sigma_Y^2}{\sigma_X^2}} \\
 &= \frac{\left| r - \frac{1}{r} \right| \sigma_X \sigma_Y}{\sigma_X^2 + \sigma_Y^2} \\
 &= \frac{(1 - r^2)}{|r|} \frac{\sigma_X \sigma_Y}{\sigma_X^2 + \sigma_Y^2}
 \end{aligned}$$

The two regression lines are at angles when  $\theta = \frac{\pi}{2}$ , i.e.,  $\tan \theta = \infty$

i.e.,  $r = 0$

$\therefore$  when the linear correlation between  $X$  and  $Y$  is zero, the two lines of regression will be at right angles.

The two regression lines are coincident, when  $\theta = 0$ , i.e., when  $\tan \theta = 0$

i.e., when  $r = \pm 1$ .

$\therefore$  when the correlation between  $X$  and  $Y$  is perfect, the two regression lines will coincide.

### Example 7

For two RV's  $X$  and  $Y$  with the same mean, the two regression equations are  $y = ax + b$  and  $x = cy + d$ . Find the common mean, ratio of the standard deviations and also show that  $\frac{b}{d} = \frac{1-a}{1-c}$ .

If  $\mu$  is the common mean, the point  $(\mu, \mu)$  lies on  $y = ax + b$  and  $x = cy + d$  [ $\because$  they intersect at  $(\bar{x}, \bar{y})$ ]

$$\therefore \mu = a\mu + b \quad (1)$$

$$\mu = c\mu + d \quad (2)$$

$$\text{From (1), } \mu = \frac{b}{1-a}$$

$$\text{From (2), } \mu = \frac{d}{1-c}$$

$$\therefore \frac{b}{1-a} = \frac{d}{1-c}$$

$$\therefore \frac{b}{d} = \frac{1-a}{1-c}$$

$$\text{Now, } \frac{\sigma_Y^2}{\sigma_X^2} = \frac{b_{YX}}{b_{XY}} = \frac{a}{c} \quad \therefore \frac{\sigma_Y}{\sigma_X} = \sqrt{\frac{a}{c}}$$

### Example 8

Find the standard error of estimate of  $Y$  on  $X$  and of  $X$  on  $Y$  from the following data:

$X$ :	1	2	3	4	5
$Y$ :	2	5	9	13	14

$x$	$y$	$x^2$	$y^2$	$xy$
1	2	1	4	2
2	5	4	25	10
3	9	9	81	27
4	13	16	169	52
5	14	25	196	70
15	43	55	475	161

$$\begin{aligned} r_{XY} &= \frac{n\sum xy - \sum x \cdot \sum y}{\sqrt{\{n\sum x^2 - (\sum x)^2\} \{n\sum y^2 - (\sum y)^2\}}} \\ &= \frac{5 \times 161 - 15 \times 43}{\sqrt{\{5 \times 55 - (15)^2\} \{5 \times 475 - (43)^2\}}} \\ &= \frac{160}{\sqrt{50 \times 526}} = 0.9866 \end{aligned}$$

$$\begin{aligned} \sigma_x^2 &= \frac{1}{n} \sum x^2 - \left( \frac{1}{n} \sum x \right)^2 \\ &= \frac{1}{5} \times 55 - \left( \frac{1}{5} \times 15 \right)^2 = 2 \end{aligned}$$

$$\therefore \sigma_x = 1.4142$$

$$\sigma_y^2 = \frac{1}{n} \sum y^2 - \left( \frac{1}{n} \sum y \right)^2$$

$$\begin{aligned}
 &= \frac{1}{5} \times 475 - \left( \frac{1}{5} \times 43 \right)^2 \\
 &= 21.04 \\
 \therefore \sigma_Y &= 4.5869 \\
 S_Y &= \sqrt{1 - r_{XY}^2} \cdot \sigma_Y = \sqrt{1 - (0.9866)^2} \times 4.5869 \\
 &= 0.7484 \\
 S_X &= \sqrt{1 - r_{XY}^2} \cdot \sigma_X = \sqrt{1 - (0.9866)^2} \times 1.4142 \\
 &= 0.2307
 \end{aligned}$$

### Exercise 4(C)

#### **Part-A (Short-answer Questions)**

1. What do you mean by regression line? What is its use?
2. For a given data of  $n$  pairs of values of  $X$  and  $Y$ , why should there be two regression lines?
3. Write down the analytic equations of the regression lines?
4. When will the two regression lines be (i) at right angles, and (ii) coincident?
5. Define regression coefficients.
6. Prove that the correlation coefficient is the geometric mean of the regression coefficients.
7. Find the co-ordinates of the point of intersection of the regression lines.
8. What do you mean by standard error of estimate?
9. Write down the formulas for the standard errors of estimate of  $Y$  and  $X$ .
10. In the usual notation prove that
  - (i)  $r_{XY} \cdot S_X S_Y = (1 - r_{XY}^2) C_{XY}$  and
  - (ii)  $b_1 S_X^2 = b_2 S_Y^2$

#### **Part-B**

11. Find the equations of the regression lines from the following data. Hence, calculate the coefficient of correlation between  $X$  and  $Y$ .
 

$X:$	62	64	65	69	70	71	72	74
$Y:$	126	125	139	145	165	152	180	208
12. Find the equations of the regression lines from the following data. Also estimate the value of  $Y$  when  $X = 71$  and the value of  $X$  when  $Y = 70$ .
 

$X:$	65	66	67	67	68	69	70	72
$Y:$	67	68	65	68	72	72	69	71

13. Find the equation of the regression line of  $Y$  on  $X$  using the method of least squares from the following data. Find the value of  $Y$  corresponding to  $X = 18$ .

$X:$	5	10	15	20	25
$Y:$	16	19	23	26	30

14. Obtain the line of regression of  $X$  on  $Y$  using the method of least squares from the following data. Find the value of  $X$  when  $Y = 45$ .

$X:$	4.7	8.2	12.4	15.8	20.7	24.9	31.9	35.0	39.1	38.8
$Y:$	4.0	8.0	12.5	16.0	20.0	25.0	31.0	36.0	40.0	40.0

15. Find the most likely price in Mumbai corresponding to the price of Rs. 70 at Chennai and that in Chennai corresponding to the price of Rs. 75 at Mumbai from the following:

	Chennai	Mumbai
Mean	65	67
S.D.	2.5	3.5

Coefficient of correlation between the prices in the two cities is 0.8.

16. In a partially destroyed laboratory record of an analysis of correlation data, the following results only are legible.  
 Variance of  $X = 9$ . Regression equations are  $8x - 10y + 66 = 0$  and  $40x - 18y = 214$ . What were (i) the mean values of  $X$  and  $Y$ ?  
 (ii) the correlation coefficient between  $X$  and  $Y$ , and (iii) the standard deviation of  $Y$ ?
17. The equations of two regression lines got in a correlation analysis are  $3x + 12y = 19$  and  $3y + 9x = 46$ . Obtain (i) the correlation coefficient between  $X$  and  $Y$ , (ii) the mean values of  $X$  and  $Y$ , and (iii) the ratio of the coefficient of variation of  $X$  to that of  $Y$ .
18. The equations of lines of regression are given by  $x + 2y - 5 = 0$  and  $2x + 3y - 8 = 0$  and variance of  $X$  is 12. Compute the values of  $\bar{x}$ ,  $\bar{y}$ ,  $\sigma_Y^2$  and  $r_{XY}$ .
19. The regression lines of  $Y$  on  $X$  and of  $X$  on  $Y$  are respectively  $y = a + bx$  and  $x = c + dy$ . Find the values of  $\bar{x}$ ,  $\bar{y}$ , and  $r_{XY}$ . Can you find  $S_X$  and  $S_Y$  from them?
20. If the lines of regression of  $Y$  on  $X$  and  $X$  on  $Y$  are respectively  $a_1x + b_1y + c_1 = 0$  and  $a_2x + b_2y + c_2 = 0$ , prove that  $a_1b_2 \leq a_2b_1$ . Find also the coefficient of correlation between  $X$  and  $Y$  and the ratio of the coefficient of variability of  $Y$  to that of  $X$ .

## Characteristic Function

Although higher order moments of a RV  $X$  may be obtained directly by using the definition of  $E(X^n)$ , it will be easier in many problems to compute them through the characteristic function or equivalently through the moment generating

function of the RV  $X$ . While the characteristic function always exists, the moment generating function need out.

Moment Generating Function (MGF) of a RV  $X$  (discrete or continuous) is defined as  $E(e^{tX})$ , where  $t$  is a real variable and denoted as  $M(t)$ .

If  $X$  is discrete, them  $M(t) = \sum_r e^{tx_r} p_r$ ,

where  $X$  takes the values  $x_1, x_2, x_3, \dots$ , with probabilities  $p_1, p_2, p_3, \dots$

If  $X$  is a continuous RV with density function  $f(x)$ , then

$$M(t) = \int_{-\infty}^{\infty} e^{tx} f(x) dx$$

## Properties of MGF

(Proofs of the properties are omitted, as the proofs of the corresponding properties of characteristic function will be given later.)

$$1. \quad M(t) = \sum_{n=0}^{\infty} t^n E(X^n) / n!$$

i.e.,  $E(X^n) = \mu'_n$  is the co-efficient of  $\frac{t^n}{n!}$  in the expansion of  $M(t)$  in series of powers of  $t$ .

$$2. \quad \mu'_n = E(X^n) = \left[ \frac{d^n}{dt^n} M(t) \right]_{t=0}$$

3. If the MGF of  $X$  is  $M_X(t)$  and if  $Y = aX + b$ , then  $M_Y(t) = e^{bt} M_X(at)$ .

4. If  $X$  and  $Y$  are independent RVs and  $Z = X + Y$ , then  $M_Z(t) = M_X(t) M_Y(t)$ .

**Characteristic function** of a RV  $X$  (discrete or continuous) is defined as  $E(e^{i\omega X})$  and denoted as  $\phi(\omega)$ .

If  $X$  is a discrete RV that can take the values  $x_1, x_2, \dots$ , such that  $P(X = x_r) = p_r$ , then

$$\phi(\omega) = \sum_r e^{i\omega x_r} p_r$$

If  $X$  is a continuous RV with density function  $f(x)$ , then

$$\phi(\omega) = \int_{-\infty}^{\infty} e^{i\omega x} f(x) dx$$

## Properties of Characteristic Function

1.  $\mu'_n = E(X^n) =$  the coefficient of  $\frac{i^n \omega^n}{n!}$  in the expansion of  $\phi(\omega)$  in series of ascending powers of  $i\omega$ .

**Proof**

$$\begin{aligned}
 \phi(\omega) &= E(e^{i\omega X}) \\
 &= E\left(1 + \frac{i\omega X}{1} + \frac{i^2 \omega^2 X^2}{2} + \dots + \frac{i^n \omega^n X^n}{n}\right) \\
 &= 1 + E(X) \times \frac{i\omega}{1} + E(X^2) \times \frac{i^2 \omega^2}{2} + \dots \\
 &= 1 + \mu'_1 \frac{i\omega}{1} + \mu'_2 \frac{i^2 \omega^2}{2} + \dots \\
 \text{i.e., } \phi(\omega) &= \sum_{n=0}^{\infty} \mu'_n \frac{i^n \omega^n}{n} \tag{1}
 \end{aligned}$$

Hence, the result.

$$2. \quad \mu'_n = \frac{1}{i^n} \left[ \frac{d^n}{d\omega^n} \phi(\omega) \right]_{\omega=0}$$

**Proof**

Differentiating both sides of (1) with respect to  $\omega$ ,  $n$  times and then putting  $\omega = 0$ .

$$\begin{aligned}
 [\phi^{(n)}(\omega)]_{\omega=0} &= \mu'_n i^n \\
 \therefore \mu'_n &= \frac{1}{i^n} [\phi^{(n)}(\omega)]_{\omega=0}
 \end{aligned}$$

3. If the characteristic function of a RV  $X$  is  $f_X(\omega)$  and if  $Y = aX + b$ , then

$$\phi_Y(\omega) = e^{ib\omega} \phi_X(a\omega)$$

**Proof**

$$\begin{aligned}
 \phi_Y(\omega) &= E\{e^{i\omega Y}\} \\
 &= E\{e^{i\omega(aX+b)}\} \\
 &= e^{ib\omega} E\{e^{ia(\omega)X}\} \\
 &= e^{ib\omega} \phi_X(a\omega)
 \end{aligned}$$

4. If  $X$  and  $Y$  are independent RVs, then

$$\phi_{X+Y}(\omega) = \phi_X(\omega) \times \phi_Y(\omega)$$

**Proof**

$$\begin{aligned}
 \phi_{X+Y}(\omega) &= E\{e^{i\omega(X+Y)}\} \\
 &= E\{e^{i\omega X} \times e^{i\omega Y}\}
 \end{aligned}$$

$$\begin{aligned}
 &= E(e^{i\omega X}) \times E(e^{i\omega Y}) \quad [\text{since } X \text{ and } Y \text{ are independent}] \\
 &= \phi_X(\omega) \times \phi_Y(\omega)
 \end{aligned}$$

5. If the characteristic function of a continuous RV  $X$  with density function

$$f(x) \text{ is } f(\omega), \text{ then } f(x) = \frac{1}{2\pi} \int_{-\infty}^{\infty} \phi(\omega) e^{-ix\omega} d\omega.$$

### **Proof**

$$\phi(\omega) = \int_{-\infty}^{\infty} f(x) e^{i\omega x} dx$$

RS can be identified as the Fourier transform of  $f(x)$ .

Therefore, by Fourier inversion formula

$$f(x) = \frac{1}{2\pi} \int_{-\infty}^{\infty} \phi(\omega) e^{-i\omega x} d\omega$$

6. If the density function of  $X$  is known, the density function of  $Y = g(X)$  can be found from the CF of  $Y$ , provided  $Y = g(X)$  is one-to-one.

### **Proof**

Let the density function of  $X$  be  $f_X(x)$ .

$$\text{Then } \phi_{g(X)}(\omega) = \int_{-\infty}^{\infty} e^{ig(x)\omega} f_X(x) dx$$

$$\text{Put } g(X) = Y \text{ and hence } g(x) = y$$

$$\text{Then } \phi_Y(\omega) = \int_{-\infty}^{\infty} e^{iy\omega} h(y) dy, \text{ say}$$

Therefore,  $h(y)$  is the density function of  $Y$ .

## **Cumulant Generating Function (CGF)**

If  $M(t)$  is the MGF of a RV  $X$ , then  $\log_e M(t)$  is called the cumulant generating function of  $X$  and denoted by  $K(t)$ .

The coefficient of  $\frac{t^r}{r!}$  in the expansion of  $K(t)$  in ascending powers of  $t$  is called the *rth order cumulant* of  $X$  and denoted by  $\lambda_r$ .

$$\text{i.e., } K(t) = \sum_{r=1}^{\infty} \frac{\lambda_r t^r}{r!}$$

$$\text{Also } \lambda_r = \left\{ \frac{d^r}{dt^r} K(t) \right\}_{t=0}$$

If  $\phi(\omega)$  is the characteristic function of a RV  $X$ , then  $\log_e \phi(\omega)$  is called the *second characteristic function* of  $X$  and denoted by  $\psi(\omega)$ .

The coefficient of  $\frac{i^r \omega^r}{r!}$  in the expansion of  $\psi(\omega)$  in ascending powers of  $\omega$  is the  $r$ th order cumulant of  $X$  and denoted by  $\lambda_r$ .

Thus, 
$$\psi(\omega) = \sum_{r=1}^{\infty} \lambda_r i^r \omega^r / r!$$

Also, 
$$\lambda_r = \frac{1}{i^r} \left\{ \frac{d^r}{d\omega^r} \psi(\omega) \right\}_{\omega=0}$$

## Joint Characteristic Function

If  $(X, Y)$  is a two-dimensional RV, then  $E(e^{i\omega_1 X + i\omega_2 Y})$  is called the joint characteristic function of  $(X, Y)$  and denoted by  $\phi_{XY}(\omega_1, \omega_2)$ .

It is easily verified that

(i)  $\phi_{XY}(0, 0) = 1$

(ii)  $E\{X^m Y^n\} = \frac{1}{i^{m+n}} \left[ \frac{\partial^{m+n}}{\partial \omega_1^m \partial \omega_2^n} \phi_{XY}(\omega_1, \omega_2) \right]_{\omega_1=0, \omega_2=0}$

(iii)  $\phi_X(\omega) = \phi_{XY}(\omega, 0)$  and  $\phi_Y(\omega) = \phi_{XY}(0, \omega)$

(iv) If  $X$  and  $Y$  are independent

$\phi_{XY}(\omega_1, \omega_2) = \phi_X(\omega_1) \times \phi_Y(\omega_2)$  and conversely.

### Worked Example 4(D)

#### Example 1

If  $X$  represents the outcome, when a fair die is tossed, find the MGF of  $X$  and hence find  $E(X)$  and  $\text{Var}(X)$ .

The probability distribution of  $X$  is given by

$$p_i = p(X = i) = \frac{1}{6}, i = 1, 2, \dots, 6$$

$$\begin{aligned} M(t) &= \sum_i e^{tx_i} p_i = \sum_{i=1}^6 e^{ti} p_i \\ &= \frac{1}{6}(e^t + e^{2t} + e^{3t} + e^{4t} + e^{5t} + e^{6t}) \end{aligned}$$

$$E(X) = [M'(t)]_{t=0} = \frac{7}{2}$$

$$\begin{aligned}
 E(X^2) &= [M''(t)]_{t=0} \\
 &= \frac{1}{6}(e^t + 4e^{2t} + 9e^{3t} + 16e^{4t} + 25e^{5t} + 36e^{6t})_{t=0} = \frac{91}{6} \\
 \therefore \text{Var}(X) &= \frac{91}{6} - \frac{49}{4} = \frac{35}{12}
 \end{aligned}$$

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### Example 2

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If a RV  $X$  has the MGF  $M(t) = \frac{3}{3-t}$ , obtain the standard deviation of  $X$ .

$$M(t) = \frac{3}{3\left(1 - \frac{t}{3}\right)} = 1 + t/3 + t^2/9 + \dots + \infty \quad (1)$$

$$E(X) = \text{coefficient of } \frac{t}{[1]} \text{ in (1)} = \frac{1}{3}$$

$$E(X^2) = \text{coefficient of } \frac{t^2}{[2]} \text{ in (1)} = \frac{2}{9}$$

$$\text{Var}(X) = E(X^2) - \{E(X)\}^2 = \frac{1}{9}$$

$$\therefore \sigma_X = \frac{1}{3}$$

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### Example 3

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Find the MGF of the binomial distribution and, hence, find its mean and variance.  
Binomial distribution is given by

$$p_r = P(X = r) = nC_r p^r q^{n-r}, r = 0, 1, 2, \dots, n$$

$$\begin{aligned}
 \therefore M(t) &= \sum_{r=0}^n e^{tr} p_r \\
 &= \sum_{r=0}^n e^{tr} nC_r p^r q^{n-r} \\
 &= \sum_{r=0}^n nC_r (p e^t)^r q^{n-r} \\
 &= (p e^t + q)^n
 \end{aligned}$$

$$M'(t) = n(p e^t + q)^{n-1} \times p e^t$$

$$M''(t) = np[(p e^t + q)^{n-1} \times e^t + (n-1)(p e^t + q)^{n-2} p e^{2t}]$$

$$\begin{aligned}
 E(X) &= M'(0) = np && (\text{since } p + q = 1) \\
 E(X^2) &= M''(0) = np [1 + (n - 1)p] \\
 \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
 &= np - np^2 \\
 &= npq
 \end{aligned}$$

**Example 4**

Find the characteristic function of the Poisson distribution and hence, find the values of the first four central moments.

Poisson distribution is given by

$$\begin{aligned}
 p_r &= P(X = r) = e^{-\lambda} \lambda^r / r! \quad r = 0, 1, 2, \dots, \infty \\
 \therefore \phi(\omega) &= \sum_{r=0}^{\infty} e^{i\omega r} e^{-\lambda} \lambda^r / r! \\
 &= \sum_{r=0}^{\infty} e^{-\lambda} (\lambda e^{i\omega})^r / r! \\
 &= e^{-\lambda} e^{\lambda e^{i\omega}} = e^{-\lambda(1-e^{i\omega})} \\
 \phi^{(1)}(\omega) &= e^{-\lambda} e^{\lambda e^{i\omega}} i \lambda e^{i\omega} \\
 \phi^{(2)}(\omega) &= i \lambda e^{-\lambda} \{e^{\lambda e^{i\omega}} e^{i\omega} i + e^{\lambda e^{i\omega}} i \lambda e^{i2\omega}\} \\
 &= i^2 \lambda e^{-\lambda} \{e^{i\omega} + \lambda e^{i2\omega}\} e^{\lambda e^{i\omega}} \\
 \phi^{(3)}(\omega) &= i^3 \lambda e^{-\lambda} \\
 &\quad [\{i e^{i\omega} + 2i \lambda e^{i2\omega}\} e^{\lambda e^{i\omega}} + \{e^{i\omega} + \lambda e^{i2\omega}\} e^{\lambda e^{i\omega}} i \lambda e^{i\omega}] \\
 &= i^3 \lambda e^{-\lambda} e^{\lambda e^{i\omega}} \{e^{i\omega} + 3\lambda e^{i2\omega} + \lambda^2 e^{i3\omega}\} \\
 \phi^{(4)}(\omega) &= i^3 \lambda e^{-\lambda} [e^{\lambda e^{i\omega}} i \lambda e^{i\omega} \{e^{i\omega} + 3\lambda e^{i2\omega} + \lambda^2 e^{i3\omega}\} \\
 &\quad + e^{\lambda e^{i\omega}} \{e^{i\omega} + 6i \lambda e^{i2\omega} + 3i \lambda^2 e^{i3\omega}\}] \\
 &= i^4 \lambda e^{-\lambda} e^{\lambda e^{i\omega}} \{e^{i\omega} + 7\lambda e^{i2\omega} + 6\lambda^2 e^{i3\omega} + \lambda^3 e^{i4\omega}\} \\
 \therefore E(X) &= \frac{1}{i} \phi^{(1)}(0) = \lambda \\
 E(X^2) &= \frac{1}{i^2} \phi^{(2)}(0) = \lambda(1 + \lambda) \\
 E(X^3) &= \frac{1}{i^3} \phi^{(3)}(0) = \lambda(1 + 3\lambda + \lambda^2)
 \end{aligned}$$

$$E(X^4) = \frac{1}{i^4} \phi^{(4)}(0) = \lambda(1 + 7\lambda + 6\lambda^2 + \lambda^3)$$

The central moments are given by

$$\mu_k = E\{(X - \mu_X)^k\}$$

$$\therefore \mu_1 = E\{X - \lambda\} = 0$$

$$\mu_2 = E\{(X - \lambda)^2\} = \lambda$$

$$\mu_3 = E\{(X - \lambda)^3\} = \lambda$$

$$\mu_4 = E\{(X - \lambda)^4\} = 3\lambda^2 + \lambda$$

### Example 5

Find the characteristic function of the geometric distribution given by  $P(X = r) = q^r p$ ,  $r = 0, 1, 2, \dots, \infty$ ,  $p + q = 1$ .

Hence, find the mean and variance.

$$\begin{aligned}\phi(\omega) &= \sum_{r=0}^{\infty} e^{i\omega r} pq^r \\ &= p \sum_{r=0}^{\infty} (qe^{i\omega})^r = p(1 - qe^{i\omega})^{-1} \\ \phi^{(1)}(\omega) &= p(1 - qe^{i\omega})^{-2} i q e^{i\omega} \\ \phi^{(2)}(\omega) &= i^2 pq[2(1 - qe^{i\omega})^{-3} q e^{i2\omega} + (1 - qe^{i\omega})^{-2} e^{i\omega}] \\ \therefore E(X) &= \frac{1}{i} \phi^{(1)}(0) = \frac{q}{p} \text{ and } E(X^2) = \frac{1}{i^2} \phi^{(2)}(0) = \frac{q}{p^2}(1 + q) \\ \therefore \mu_X &= \frac{q}{p} \text{ and } \sigma_X^2 = \frac{q}{p^2}\end{aligned}$$

### Example 6

Obtain the characteristic function of the normal distribution. Deduce the first four central moments.

Let  $X$  follows  $N(\mu, \sigma)$

Then  $Z = \frac{X - \mu}{\sigma}$  follows  $N(0, 1)$ , i.e., the standard normal distribution whose density function is  $\phi(z) = \frac{1}{\sqrt{2\pi}} e^{-z^2/2}$ .

$$\begin{aligned}\text{Now, } \phi_Z(\omega) &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-z^2/2} e^{i\omega z} dz \\ &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-\frac{1}{2}(z-i\omega)^2 + \frac{i^2\omega^2}{2}} dz\end{aligned}$$

$$\begin{aligned}
 &= e^{-\omega^2/2} \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-t^2} \cdot \sqrt{2} dt, \text{ on putting } \frac{z-i\omega}{\sqrt{2}} = t \\
 &= e^{-\omega^2/2} \left[ \text{since } \int_{-\infty}^{\infty} e^{-t^2} dt = \sqrt{\left(\frac{1}{2}\right)} = \sqrt{\pi} \right] \\
 \therefore \phi_X(\omega) &= f_{\sigma Z + \mu}(\omega) \\
 &= e^{i\mu\omega} \phi_Z(\sigma\omega) \text{ (by Property 3)} \\
 &= e^{i\mu\omega} e^{-\sigma^2\omega^2/2} \\
 \text{Now, } \phi_X(\omega) &= e^{i\omega(\mu+i\sigma^2\omega/2)} \\
 &= 1 + \underbrace{\frac{1}{1} i\omega(\mu + i\sigma^2\omega/2)} + \underbrace{\frac{1}{2} i^2\omega^2(\mu + i\sigma^2\omega/2)^2}_{\dots} \\
 &= 1 + \underbrace{\frac{i\omega\mu}{1} + \frac{i^2\omega^2}{2}(\sigma^2 + \mu^2)} + \underbrace{\frac{i^3\omega^3}{3}(3\mu\sigma^2 + \mu^3)}_{\dots} \\
 &\quad + \underbrace{\frac{i^4\omega^4}{4}(3\sigma^4 + 6\mu^2\sigma^2 + \mu^4)}_{\dots} \tag{1} \\
 \therefore \mu'_1 &= E(X) = \text{coefficient of } \frac{i\omega}{1} \text{ in (1)} = \mu \\
 \text{Similarly, } E(X^2) &= \sigma^2 + \mu^2; \\
 E(X^3) &= 3\mu\sigma^2 + \mu^3 \text{ and,} \\
 E(X^4) &= 3\sigma^4 + 6\mu^2\sigma^2 + \mu^4 \\
 \text{Using the relation } \mu_k &= E\{(X - \mu)^k\}, \text{ we get} \\
 \mu_1 &= 0, \mu_2 = \sigma^2, \mu_3 = 0 \text{ and } \mu_4 = 3\sigma^4
 \end{aligned}$$


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### Example 7

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Find the characteristic function of the Erlang distribution given by  $f(x) = f(x) = \frac{\lambda^n}{|n-1|} x^{n-1} e^{-\lambda x} U(x)$  and, hence, find its mean and variance.

$$\begin{aligned}
 \phi(\omega) &= \frac{\lambda^n}{|n-1|} \int_0^{\infty} x^{n-1} e^{-(\lambda-i\omega)x} dx \\
 &= \frac{\lambda^n}{|n-1|(\lambda-i\omega)^n} \int_0^{\infty} t^{n-1} e^{-t} dt, \text{ on putting } (\lambda - i\omega)x = t \\
 &= \left( \frac{\lambda}{\lambda - i\omega} \right)^n \frac{1}{|n-1|} \overline{(n)} = \frac{1}{\left( 1 - \frac{i\omega}{\lambda} \right)^n} \\
 &\quad [\text{since } \overline{n} = |n-1| \text{ when } n \text{ is a positive integer}]
 \end{aligned}$$

Now  $\phi(\omega) = \left(1 - \frac{i\omega}{\lambda}\right)^{-n}$

$$= 1 + \frac{n}{1} \frac{i\omega}{\lambda} + \frac{n(n+1)}{2} \frac{i^2 \omega^2}{\lambda^2} + \dots + \infty$$

$$\therefore E(X) = \frac{n}{\lambda} \text{ and } E(X^2) = \frac{n(n+1)}{\lambda^2}$$

$$\therefore \text{Var}(X) = \frac{n}{\lambda^2}$$

**Note** If  $n = 1$ , Erlang's distribution becomes the exponential distribution. Hence, the mean and variance of an exponential distribution with parameter  $\lambda$  are  $1/\lambda$  and  $1/\lambda^2$  respectively.

### Example 8

Find the characteristic function of the Laplace distribution with pdf  $f(x) = \frac{\alpha}{2} e^{-\alpha|x|}$ ,  $-\infty < x < \infty$ . Hence, find its mean and variance.

$$\begin{aligned}\phi(\omega) &= \frac{\alpha}{2} \int_{-\infty}^{\infty} e^{-\alpha|x|} e^{i\omega x} dx \\ &= \frac{\alpha}{2} \left[ \int_{-\infty}^0 e^{(\alpha+i\omega)x} dx + \int_0^{\infty} e^{-(\alpha-i\omega)x} dx \right] \\ &= \frac{\alpha}{2} \left\{ \frac{1}{\alpha+i\omega} + \frac{1}{\alpha-i\omega} \right\} = \frac{\alpha^2}{\alpha^2 + \omega^2}\end{aligned}$$

Now,  $\phi(\omega) = \left(1 + \frac{\omega^2}{\alpha^2}\right)^{-1} = 1 - \frac{\omega^2}{\alpha^2} + \frac{\omega^4}{\alpha^4} + \dots + \infty$

$$E(X) = 0 \text{ and } E(X^2) = \text{Var}(X) = \frac{2}{\alpha^2}$$

### Example 9

If  $X_1$  and  $X_2$  are two independent RVs that follow Poisson distribution with parameters  $\lambda_1$  and  $\lambda_2$ , prove that  $(X_1 + X_2)$  also follows a Poisson distribution with parameter  $(\lambda_1 + \lambda_2)$ .

(This property is called the re-productivity property of the Poisson distribution.)

$$\phi_{X_1}(t) = e^{\lambda_1(e^{i\omega}-1)}$$

$$\phi_{X_2}(t) = e^{\lambda_2(e^{i\omega}-1)}$$

Since  $X_1$  and  $X_2$  are independent RVs,

$\phi_{X_1+X_2}(t) = e^{(\lambda_1+\lambda_2)(e^{i\omega}-1)}$ , which is the characteristic function of a Poisson distribution with parameter  $(\lambda_1 + \lambda_2)$ . Hence, the result.

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### Example 10

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Show that the distribution for which the characteristic function is  $e^{-|\omega|}$  has the density function  $f(x) = \frac{1}{\pi} \times \frac{1}{1+x^2}, -\infty < x < \infty$ .

By inversion Property 5,

$$\begin{aligned} f(x) &= \frac{1}{2\pi} \int_{-\infty}^{\infty} \phi(\omega) e^{-ix\omega} d\omega \\ &= \frac{1}{2\pi} \int_{-\infty}^{\infty} e^{-|\omega|} (\cos x\omega - i \sin x\omega) d\omega \\ &= \frac{1}{\pi} \int_0^{\infty} e^{-\omega} \cos x\omega d\omega \quad (\text{by properties of odd and even function}) \\ &= \frac{1}{\pi} \left[ \frac{e^{-\omega}}{1+x^2} (-\cos x\omega + x \sin x\omega) \right]_0^{\infty} \\ &= \frac{1}{\pi} \times \frac{1}{1+x^2}, -\infty < x < \infty \end{aligned}$$

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### Example 11

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Find the density function  $f(x)$  corresponding to the characteristic function defined as

$$\phi(t) = \begin{cases} 1-|t| & \text{for } |t| \leq 1 \\ 0 & \text{for } |t| > 1 \end{cases}$$

**Note** The letter  $t$  is used in the place of  $\omega$ .

By inversion Property 5,

$$f(x) = \frac{1}{2\pi} \int_{-\infty}^{\infty} \phi(t) e^{-ixt} dt$$

$$\begin{aligned}
&= \frac{1}{2\pi} \int_{-1}^1 \{1 - |t|\} e^{-ixt} dt \\
&= \frac{1}{\pi} \int_0^1 (1-t) \cos xt dt \\
&= \frac{1}{\pi} \left\{ (1-t) \frac{\sin xt}{x} - \frac{\cos xt}{x^2} \right\}_0^1 \\
&= \frac{1}{\pi} \frac{(1-\cos x)}{x^2}
\end{aligned}$$

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### Example 12

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Express the first four cumulants in terms of central moments.

By definition,  $K(t) = \log\{M(t)\}$

$$\begin{aligned}
\therefore \sum_{r=1}^{\infty} \frac{\lambda_r}{[r]} &= \log \left\{ \sum_{r=0}^{\infty} \frac{\mu'_r t^r}{[r]} \right\} \\
&= \log \left\{ 1 + \mu'_1 t + \frac{\mu'_2}{2} t^2 + \frac{\mu'_3}{3} t^3 + \dots \right\} \\
&= t \left( \mu'_1 + \frac{\mu'_2}{2} t + \frac{\mu'_3}{6} t^2 + \frac{\mu'_4}{24} t^3 + \dots \right) \\
&\quad - \frac{t^2}{2} \left( \mu'_1 + \frac{\mu'_2 t}{2} + \frac{\mu'_3 t^2}{6} + \frac{\mu'_4 t^3}{24} + \dots \right)^2 + \dots
\end{aligned}$$

Comparing like coefficients, we get,

$$\begin{aligned}
\lambda_1 &= \mu'_1; \lambda'_2 = \mu'_2 - \mu'^2_1 = \mu_2 \\
\lambda_3 &= \mu'_3 - 3\mu'_2 \mu'_1 + 2\mu'^3_1 = \mu_3 \\
\lambda_4 &= \mu'_4 - 4\mu'_3 \mu'_1 - 3\mu'^2_2 + 12\mu'_2 \mu'^2_1 - 6\mu'^4_1 \\
&= \mu_4 - 3\mu_2^2
\end{aligned}$$

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### Example 13

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If  $X$  and  $Y$  are two independent RVs, prove that the cumulant of  $(X + Y)$  of any order is the sum of the cumulants of  $X$  and  $Y$  of the same order.

By Property 4 of characteristic functions,  $\phi_{X+Y}(\omega) = \phi_X(\omega) \times \phi_Y(\omega)$ , when  $X$  and  $Y$  are independent.

$$\therefore \log \phi_{X+Y}(\omega) = \log \phi_X(\omega) + \log \phi_Y(\omega)$$

$$\text{i.e., } \psi_{X+Y}(\omega) = \psi_X(\omega) + \psi_Y(\omega)$$

$$\text{i.e., } \sum_{r=1}^{\infty} \lambda_r(X+Y) \frac{i^r \omega^r}{|r|} = \sum_{r=1}^{\infty} \lambda_r(X) \frac{i^r \omega^r}{|r|} + \sum_{r=1}^{\infty} \lambda_r(Y) \frac{i^r \omega^r}{|r|}$$

$$\therefore \lambda_r(X+Y) = \lambda_r(X) + \lambda_r(Y)$$

### Example 14

If the RV  $X$  follows  $N(0, \sigma)$ , find the density function of  $Y = aX^2$ , using the characteristic function technique.

$$\phi_{aX^2}(\omega) = \int_{-\infty}^{\infty} e^{i\omega aX^2} \times \frac{1}{\sigma\sqrt{2\pi}} e^{-x^2/2\sigma^2} dx \quad (1)$$

Put  $y = ax^2$  in (1)

$$\therefore dx = \frac{dy}{2ax} = \frac{dy}{2\sqrt{ay}}$$

$$\therefore \phi_Y(\omega) = 2 \int_0^{\infty} e^{i\omega y} \times \frac{1}{2\sigma\sqrt{2\pi ay}} \times e^{-y/2a\sigma^2} dy \quad (2)$$

$$\text{But } \phi_Y(\omega) = \int_{-\infty}^{\infty} e^{i\omega y} f_Y(y) dy \quad (3)$$

Comparing (2) and (3), we get

$$f_Y(y) = \frac{1}{\sigma\sqrt{2\pi ay}} e^{-y/2a\sigma^2} U(y)$$

### Example 15

Two RVs  $X$  and  $Y$  have the joint characteristic function  $\phi_{XY}(\omega_1, \omega_2) = \exp(-2\omega_1^2 - 8\omega_2^2)$ . Show that  $X$  and  $Y$  are both zero mean RVs and also that they are uncorrelated.

By the property of joint CF

$$\begin{aligned} E(X) &= \frac{1}{i} \left[ \frac{\partial}{\partial \omega_1} (\exp(-2\omega_1^2 - 8\omega_2^2)) \right]_{\omega_1=0, \omega_2=0} \\ &= \left[ e^{-8\omega_2^2} e^{-2\omega_1^2} 4i\omega_1 \right]_{\omega_1=0, \omega_2=0} \\ &= 0 \end{aligned}$$

$$\begin{aligned}
 E(Y) &= \left[ e^{-2\omega_1^2} e^{-8\omega_2^2} 16i\omega_2 \right]_{\omega_1=0, \omega_2=0} \\
 &= 0 \\
 E(XY) &= \frac{1}{i^2} \left[ \frac{\partial^2}{\partial \omega_1 \partial \omega_2} (e^{-2\omega_1^2 - 8\omega_2^2}) \right]_{\omega_1=0, \omega_2=0} \\
 &= \frac{\partial}{\partial \omega_1} \left[ e^{-2\omega_1^2} e^{-8\omega_2^2} 16\omega_2 \right]_{\omega_1=0, \omega_2=0} \\
 &= \{-64\omega_1\omega_2 e^{-2\omega_1^2 + 8\omega_2^2}\}_{\omega_1=0, \omega_2=0} \\
 &= 0
 \end{aligned}$$

$$C_{XY} = E(XY) - E(X) \times E(Y) = 0$$

Therefore,  $X$  and  $Y$  are uncorrected RVs.

### Exercise 4(D)

#### Part-A (Short-answer Questions)

1. Define the MGF of a RV  $X$ . Why is it called so?
2. State the properties of the MGF of a RV.
3. Derive the relation between the MGFs of  $X$  and  $Y$  when  $Y = aX + b$ .
4. If  $X$  and  $Y$  are independent RVs and  $Z = X + Y$ , prove that  $M_Z(t) = M_X(t) \times M_Y(t)$ .
5. Define the characteristic function of a RV. How does it differ from the MGF?
6. State the properties of the characteristic function of a RV.
7. State the uniqueness theorem of characteristic functions.
8. If  $Y = aX + b$ , find the relation between the characteristic functions of  $X$  and  $Y$ .
9. If  $X$  and  $Y$  are 2 independent RVs prove that  $\phi_{X+Y}(\omega) = \phi_X \times \phi_Y(\omega)$ .
10. If the characteristic function of a continuous RV  $X$  is  $\phi(\omega)$ , express its density function  $f(x)$  in terms of  $\phi(\omega)$ .
11. Define the cumulant generating function/the second characteristic function of a RV  $X$  and what is its use?
12. If the  $r$ th moment of a continuous RV  $X$  about the origin is  $r!$ , find the MGF of  $X$ .
13. If the MGF of a RV  $X$  is  $\frac{2}{2-t}$  find the SD of  $X$ .
14. Find the MGF/CF of a uniform distribution in  $(a, b)$ .
15. Find the MGF/CF of the binomial distribution.
16. Find the MGF/CF of the Poisson distribution.

17. State and prove the reproductive property of the Poisson distribution.
18. Find the CF/MGF of the geometric distribution.
19. Find the CF/MGF of the exponential distribution.
20. If the CF of the standard normal distribution  $N(0, 1)$  is  $e^{-\omega^2/2}$ , find the CF of the general normal distribution  $N(\mu, \sigma)$ .
21. Find the CF of  $X$  whose pdf is given by  $f(x) = \frac{1}{2}e^{-|x|}$ ,  $-\infty < x < \infty$ .

**Part-B**

22. Find the MGF of a RV which is uniformly distributed over  $(-1, 2)$  and hence find its mean and variance.
23. Find the characteristic function of the binomial distribution and, hence, find its mean and variance.
24. Obtain the MGF of the Poisson distribution; deduce the values of the first four central moments.
25. Find the characteristic function of the negative binomial distribution given by  $P(X = r) = (n + r - 1) C_r q^r p^n$ , ( $r = 0, 1, 2, \dots, \infty$ ),  $p + q = 1$ , and hence find its mean and variance.
26. Find the MGF of the two-parameter exponential distribution whose density function is given by  $f(x) = \lambda e^{-\lambda(x-a)}$ ,  $x \geq a$ . Hence, find its mean and variance.
27. If the density function of a continuous RV  $X$  is given by  $f(x) = \frac{1}{2}e^{-|x|}$ ,  $-\infty < x < \infty$ , find the MGF of  $X$ . Hence, find its mean and variance.
28. Find the characteristic function of the Cauchy's distribution given by  $f(x) = \frac{1}{\pi} \times \frac{1}{1+x^2}$   $-\infty < x < \infty$ . Comment about the first two moments.  
**(Hint:** Use contour integration.)
29. Find the characteristic function for the following probability density function:  
$$f_X(x) = \frac{\lambda}{\pi(\lambda^2 + x^2)}$$

30. If  $X$  follows  $N(\mu_X, \sigma_X)$  and  $Y$  follows  $N(\mu_Y, \sigma_Y)$ , prove, by using characteristic functions, that  $(aX + bY)$  follows a normal distribution with mean  $(a\mu_X + b\mu_Y)$  and variance  $(a^2\sigma_X^2 + b^2\sigma_Y^2)$ .
31. Find the density function of the distribution for which the characteristic function is given by  $\phi(t) = e^{-\sigma^2 t^2/2}$ .
32. If a raw moments of a continuous RV  $X$  are given by  $E(X^n) = [n]$ , find the characteristic function of  $X$  and also the density function of  $X$ .  
**(Hint:** Use contour integration to find pdf.)
33. Express the first 4 raw moments about the origin in terms of the cumulants.

34. Prove that the cumulants of all orders are equal for the Poisson distribution.
35. If  $X$  and  $Y$  are two jointly normal RVs whose joint pdf is  $N(\mu_X, \mu_Y; \sigma_X, \sigma_Y; r)$ , find the joint characteristic function of  $(X, Y)$ .
36. If the random variable  $X$  is uniformly distributed over  $(-\pi/2, \pi/2)$ , find the pdf of  $Y = \sin X$ , using characteristic function technique.

## Bounds on Probabilities

If we know the probability distribution of a random variable  $X$  (i.e., the pdf in the continuous case or the pmf in the discrete case), we may compute  $E(X)$  and  $\text{Var}(X)$ . Conversely, if  $E(X)$  and  $\text{Var}(X)$  are known, we cannot construct the probability distribution of  $X$  and hence compute quantities such as  $P\{|X - E(X)| \leq k\}$ . Although we cannot evaluate such probabilities from a knowledge of  $E(X)$  and  $\text{Var}(X)$ , several approximation techniques have been developed to yield upper and/or lower bounds to such probabilities. The most important of such techniques is Tchebycheff inequality.

## Tchebycheff Inequality

If  $X$  is a RV with  $E(X) = \mu$  and  $\text{Var}(X) = \sigma^2$ , then  $P\{|X - \mu| \geq c\} \leq \frac{\sigma^2}{c^2}$  where  $c > 0$ .

### Proof

Let  $X$  be a continuous RV with pdf  $f(x)$ .

$$\begin{aligned} \text{Then } \sigma^2 &= \int_{-\infty}^{\infty} (x - \mu)^2 f(x) dx \\ &\geq \int_{-\infty}^{\mu-c} (x - \mu)^2 f(x) dx + \int_{\mu+c}^{\infty} (x - \mu)^2 f(x) dx \end{aligned}$$

In the first integral,  $x \leq \mu - c$

$$\therefore (x - \mu)^2 \geq c^2$$

In the second integral,  $x \geq \mu + c$

$$\therefore (x - \mu)^2 \geq c^2$$

$$\begin{aligned} \therefore \sigma^2 &\geq c^2 \left\{ \int_{-\infty}^{\mu-c} f(x) dx + \int_{\mu+c}^{\infty} f(x) dx \right\} \quad (1) \end{aligned}$$

$$\begin{aligned} \text{RS of (1)} &= c^2 [1 - P\{\mu - c \leq X \leq \mu + c\}] \\ &= c^2 [1 - P\{-c \leq X - \mu \leq c\}] \\ &= c^2 [1 - P\{|X - \mu| \leq c\}] \\ &= c^2 P\{|X - \mu| \geq c\} \quad (2) \end{aligned}$$

Using (2) in (1), we get

$$P\{|X - \mu| \geq c\} \leq \frac{\sigma^2}{c^2} \quad \text{or} \quad P\{|X - \mu| < c\} \geq 1 - \frac{\sigma^2}{c^2}, c > 0$$

**Alternative forms:** If we put  $c = k\sigma$ , where  $k > 0$ , then Tchebycheff inequality takes the form

$$\begin{aligned} P\left\{\left|\frac{X - \mu}{k}\right| \geq \sigma\right\} &\leq \frac{1}{k^2} \\ \therefore P\left\{\left|\frac{X - \mu}{k}\right| \leq \sigma\right\} &\geq 1 - \frac{1}{k^2} \end{aligned}$$

**Note** Although we have proved the inequality for the continuous case, it holds good for the discrete RV also.

## Bienayme's Inequality

Let us first prove a basic result:

If  $X$  is a RV with  $f(x) = 0$ , when  $x < 0$ , and with  $E(x) = \mu$ , then for any  $\alpha > 0$ ,  $P(X \geq \alpha) \leq \mu/\alpha$ .

### Proof

$$\begin{aligned} \mu = E(X) &= \int_0^\infty xf(x)dx \\ &\geq \int_\alpha^\infty xf(x)dx \quad (\text{since } \alpha > 0) \\ \text{i.e.,} \quad &\geq \alpha \int_\alpha^\infty f(x)dx \quad (\text{since } x \geq \alpha) \\ \text{i.e.,} \quad &\geq \alpha P(X \geq \alpha) \\ \therefore \quad P(X \geq \alpha) &\leq \frac{\mu}{\alpha} \end{aligned} \tag{1}$$

[In particular,  $P(X \geq \sqrt{\mu}) \leq \sqrt{\mu}$  .]

In (1), replace  $X$  by  $|X - a|^n$  and  $\alpha$  by  $c^n$ , where  $X$  is an arbitrary RV and  $c, a, n$  are arbitrary numbers.

$|X - a|^n$  takes only positive values.

$$\begin{aligned} \therefore \quad P\{|X - a|^n \geq c^n\} &\leq \frac{E\{|X - a|^n\}}{c^n} \\ \text{i.e.,} \quad P(\{|X - a| \geq c\}) &\leq \frac{E\{|X - a|^n\}}{c^n} \end{aligned} \tag{2}$$

This inequality is called the Bienayme's inequality.

If we take  $n = 2$  and  $a = \mu$  in (2), we get Tchebycheff's inequality.

## Schwartz Inequality

If  $f(x)$  and  $g(x)$  are real functions of  $x$ , then

$$\left( \int_a^b f(x)g(x) dx \right)^2 \leq \int_a^b [f(x)]^2 dx \times \int_a^b [g(x)]^2 dx$$

### Proof

Consider  $\int_a^b [f(x) - tg(x)]^2 dx$ , where  $t$  is real,

$$= t^2 \int_a^b g^2(x) dx - 2t \int_a^b f(x) g(x) dx + \int_a^b f^2(x) dx$$

For every real  $t$ , LS  $\geq 0$ .

Therefore, RS which is a quadratic expression in  $t \geq 0$ .

Therefore, discriminant  $\leq 0$ .

i.e.,  $\left[ \int_a^b f(x)g(x) dx \right]^2 \leq \int_a^b f^2(x) dx \times \int_a^b g^2(x) dx$

### Note

(i) The equality holds, if  $f(x) = kg(x)$ .

(ii) If  $f$  and  $g$  are complex, then

$$\left| \int f(z)g(z) dz \right|^2 \leq \int |f(z)|^2 dz \int |g(z)|^2 dz$$

## Cauchy-Schwartz Inequality

For any two RVs  $X$  and  $Y$ ,

$$\{E(XY)\}^2 \leq E(X^2)E(Y^2)$$

### Proof

Consider  $E\{(X - tY)^2\} \geq 0$ , where  $t$  is real.

i.e.,  $E(X^2) - 2tE(XY) + t^2E(Y^2) \geq 0$

Since the LS is a quadratic expression in  $t$ , discriminant of LS  $\leq 0$ .

i.e.,  $\{E(XY)\}^2 \leq E(X^2) \times E(Y^2)$

### Worked Example 4(E)

#### Example 1

A RV  $X$  has mean  $\mu = 12$  and variance  $\sigma^2 = 9$  and an unknown probability distribution. Find  $P(6 < X < 18)$ .

Since the probability distribution of  $X$  is not known, we cannot find the value of the required probability. We can find only a lower bound for the probability using Tchebycheff's inequality.

$$P\{|X - \mu| \geq c\} \leq \frac{\sigma^2}{c^2}, c > 0$$

i.e.,  $P\{|X - \mu| \leq c\} \geq 1 - \frac{\sigma^2}{c^2}$

$\therefore P\{\mu - c \leq X \leq \mu + c\} \geq 1 - \frac{\sigma^2}{c^2}$

Taking  $\mu = 12$  and  $\sigma^2 = 9$ , we get,

$$P\{12 - c < X < 12 + c\} \geq 1 - \frac{9}{c^2}$$

Putting  $c = 6$ ,  $P\{6 < X < 18\} \geq 1 - \frac{9}{36}$

i.e.,  $P\{6 < X < 18\} \geq \frac{3}{4}$

## Example 2

If the RV  $X$  is uniformly distributed over  $(-\sqrt{3}, \sqrt{3})$ , compute  $P\left\{|X - \mu| \geq \frac{3\sigma}{2}\right\}$  and compare it with the upper bound obtained by Tchebycheff's inequality. [Refer to Problem 46 in Exercise 4(A)].

$$\mu = 0 \text{ and } \sigma^2 = 1$$

$$\begin{aligned} \therefore P\left\{|X - \mu| \geq \frac{3\sigma}{2}\right\} &= P\{|X| \geq 3/2\} \\ &= 1 - P\{-3/2 < X < 3/2\} \\ &= 1 - \int_{-3/2}^{3/2} \frac{1}{2\sqrt{3}} dx \left( \text{since pdf} = \frac{1}{2\sqrt{3}} \right) \\ &= 1 - \frac{\sqrt{3}}{2} = 0.134 \end{aligned}$$

By Tchebycheff's inequality,

$$P\{|X - \mu| \geq k\sigma\} \leq \frac{1}{k^2}$$

$\therefore P\{|X - \mu| \geq 3/2\sigma\} \leq \frac{4}{9} = 0.444$ ,

which is a poor upper bound.

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### Example 3

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Can we find a RV  $X$  for which  $P\{\mu - 2\sigma \leq X \leq \mu + 2\sigma\} = 0.6?$

$$P\{\mu - 2\sigma \leq X \leq \mu + 2\sigma\} = P\{|X - \mu| \leq 2\sigma\}$$

$$\geq 1 - \frac{1}{2^2} \text{ by Tchebycheff's inequality}$$

$$\text{i.e., } P\{\mu - 2\sigma \leq X \leq \mu + 2\sigma\} \geq 0.75$$

Therefore, there does not exist a RV  $X$  satisfying the given condition.

---

### Example 4

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A discrete RV  $X$  takes the values  $-1, 0, 1$  with probabilities  $\frac{1}{8}, \frac{3}{4}, \frac{1}{8}$  respectively.

Evaluate  $P\{|X - \mu| \geq 2\sigma\}$  and compare it with the upper bound given by Tchebycheff's inequality.

$$\begin{aligned} E(X) &= -1 \times \frac{1}{8} + 0 \times \frac{3}{4} + 1 \times \frac{1}{8} = 0 \\ E(X^2) &= 1 \times \frac{1}{8} + 0 \times \frac{3}{4} + 1 \times \frac{1}{8} = \frac{1}{4}; \text{Var}(X) = \frac{1}{4} \\ \therefore P\{|X - \mu| \geq 2\sigma\} &= P\{|X| \geq 1\} \\ &= P\{X = -1 \text{ or } X = 1\} \\ &= \frac{1}{8} + \frac{1}{8} = \frac{1}{4} \end{aligned} \tag{1}$$

By Tchebycheff's inequality,

$$P\{|X - \mu| \geq 2\sigma\} \leq \frac{1}{2^2} = \frac{1}{4} \tag{2}$$

The two values coincide.

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### Example 5

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If  $X$  denotes the sum of the numbers obtained when 2 dice are thrown, obtain an upper bound for  $P\{|X - 7| \geq 4\}$ . Compare with the exact probability.

Let  $X_1, X_2$  denote the outcomes of the first and second dice respectively.

$$E(X_1) = E(X_2) = \frac{1}{6}(1+2+\dots+6) = \frac{7}{2}$$

$$E(X_1^2) = E(X_2^2) = \frac{1}{6}(1^2 + 2^2 + \dots + 6^2) = \frac{91}{6}$$

$$\therefore \text{Var}(X_1) = \text{Var}(X_2) = \frac{35}{12}$$

$$\therefore E(X) = E(X_1 + X_2) = 7 \text{ and}$$

$$\text{Var}(X) = \text{Var}(X_1 + X_2) = 1^2 \times \frac{35}{12} + 1^2 \times \frac{35}{12} = \frac{35}{6}$$

By Tchebycheff's inequality,

$$\begin{aligned} P\{|X - \mu| \geq c\} &\leq \frac{\sigma^2}{c^2} \\ \therefore P\{|X - 7| \geq 4\} &\leq \frac{35}{96} \end{aligned} \quad (1)$$

$$\text{Now, } P\{|X - 7| \geq 4\} = P\{X = 2, 3, 11 \text{ or } 12\}$$

$$= \frac{1}{36} + \frac{2}{36} + \frac{2}{36} + \frac{1}{36} = \frac{1}{6} \quad (2)$$

There is much difference between the two values.

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### Example 6

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A fair dice is tossed 720 times. Use Tchebycheff's inequality to find a lower bound for the probability of getting 100 to 140 sixes.

Let  $X$  be the number of sixes obtained when a fair dice is tossed 720 times.

$$\begin{aligned} p &= P\{\text{getting '6' in a single toss}\} = \frac{1}{6} \\ q &= \frac{5}{6} \text{ and } n = 720 \end{aligned}$$

$X$  follows a binomial distribution with mean  $np = 120$  and variance  $npq = 100$ , i.e.,  $\mu = 120$  and  $\sigma = 10$ .

By Tchebycheff's inequality,

$$P\{|X - \mu| \leq k\sigma\} \geq 1 - \frac{1}{k^2}$$

$$\text{i.e., } P\{|X - 120| \leq 10k\} \geq 1 - \frac{1}{k^2}$$

$$\text{i.e., } P\{120 - 10k \leq X \leq 120 + 10k\} \geq 1 - \frac{1}{k^2}$$

Taking  $k = 2$ , we get,

$$P\{100 \leq X \leq 140\} \geq \frac{3}{4}$$

Therefore, required lower bound for the probability = 0.75.

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### Example 7

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Use Tchebycheff's inequality to find how many times a fair coin must be tossed in order that the probability that the ratio of the number of heads to the number of tosses will lie between 0.45 and 0.55 will be at least 0.95.

Let  $X$  be the number of heads obtained when a fair coin is tossed  $n$  times. Then  $X$  follows a binomial distribution.

$$B(np, \sqrt{npq}) \text{ where } p = q = \frac{1}{2}$$

$$\therefore \frac{X}{n} \text{ follows } B\left(\frac{1}{2}, \frac{1}{2\sqrt{n}}\right)$$

By Tchebycheff's inequality,

$$P\left\{\left|\frac{X}{n} - \frac{1}{2}\right| \leq c\right\} \geq 1 - \frac{\sigma^2}{c^2}$$

$$\text{i.e., } P\{0.5 - c \leq \frac{X}{n} \leq 0.5 + c\} \geq 1 - \frac{1}{4nc^2}$$

$$\text{Taking } c = 0.05, P\left\{0.45 \leq \frac{X}{n} \leq 0.55\right\} \geq 1 - \frac{100}{n}$$

$$\text{Given that } P\left\{0.45 \leq \frac{X}{n} \leq 0.55\right\} \geq 0.95$$

$$\therefore 1 - \frac{100}{n} = 0.95$$

$$\therefore n = 2000$$

### Example 8

A RV  $X$  is exponentially distributed with parameter 1. Use Tchebycheff's inequality to show that  $P(-1 \leq X \leq 3) \geq 3/4$ . Find the actual probability also.

For an exponential distribution with parameter  $\lambda$ ,  $E(X) = \frac{1}{\lambda}$  and  $\text{Var}(X) = \frac{1}{\lambda^2}$ .

(refer to Worked Example 7 in the previous section of this chapter).

Here,  $\lambda = 1$

$$\therefore \mu = 1 \text{ and } \sigma = 1$$

By Tchebycheff's inequality,

$$P\{|X - 1| \leq 2\} \geq 1 - \frac{1}{4}$$

$$\text{i.e., } P\{-1 \leq X \leq 3\} \geq \frac{3}{4}$$

Density function of the exponential distribution with  $\lambda = 1$  is given by  $f(x) = e^{-x}$ ,  $x > 0$ .

$$\therefore P\{-1 \leq X \leq 3\} = \int_{-1}^3 e^{-x} dx$$

$$= \int_0^3 e^{-x} dx = 1 - e^{-3} = 0.9502$$

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### Example 9

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Use Tchebycheff's inequality to prove that  $P(X = \mu) = 1$ , if  $\text{Var}(X) = 0$   
By Tchebycheff's inequality

$$P\{|X - \mu| \geq c\} \leq \frac{\sigma^2}{c^2}$$

and  $P\{|X - \mu| \leq c\} \geq 1 - \frac{\sigma^2}{c^2}$

Putting  $\sigma^2 = 0$ , these inequalities become

$$P\{|X - \mu| \geq c\} \leq 0$$

and  $P\{|X - \mu| \leq c\} \geq 1$  where  $c > 0$

The above results hold good even for arbitrarily chosen small values of  $c$ . Hence in the limit when  $c \rightarrow 0$ ,

$$P\{|X - \mu| = 0\} = 1$$

i.e.,  $P\{X = \mu\} = 1$

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### Example 10

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A discrete RV  $X$  can assume the values  $x = 1, 2, 3, \dots$  with probability  $2^{-x}$ . Show that  $P\{|X - 2| \geq 2\} \leq 1/2$ , while the actual probability is  $1/8$ .

$$E(X) = \sum_{x=1}^{\infty} \frac{x}{2^x} = 2 \text{ and } \text{Var}(X) = 2$$

(Refer to Worked Example 6 in the first section of Chapter 2)

By Tchebycheff's inequality,

$$P\{|X - \mu| \geq c\} \leq \frac{\sigma^2}{c^2}$$

Putting  $\mu = 2$ ,  $\sigma^2 = 2$ , and  $c = 2$ , we get

$$P\{|X - 2| \geq 2\} \leq \frac{1}{2}$$

Now,  $P\{|X - 2| \geq 2\} = P\{X \geq 4\}$

$$\begin{aligned} &= \frac{1}{2^4} + \frac{1}{2^5} + \frac{1}{2^6} + \dots + \infty \\ &= \frac{\frac{1}{2^2}}{1 - \frac{1}{2}} = \frac{1}{8} \end{aligned}$$

**Exercise 4(E)****Part-A (Short-answer Questions)**

1. State Tchebycheff's inequality.
2. A RV  $X$  has mean 10 and variance 16. Find a lower bound for  $P\{5 < X < 15\}$ .
3. A RV  $X$  has mean 10 and variance 16. Find an upper bound for  $P\{|X - 10| \geq 15\}$ .
4. If Tchebycheff's inequality for a RV  $X$  with mean 12 is  $P\{6 < X < 18\} \geq \frac{3}{4}$ , find the SD of  $X$ .
5. If Tchebycheff's inequality for a RV  $X$  with SD 3 is  $P\{6 < X < 18\} \geq \frac{3}{4}$ , find the mean of  $X$ .
6. If Tchebycheff's inequality for a RV  $X$  is  $P\{-2 < X < 8\} \geq \frac{21}{25}$ , find  $E(X)$  and  $\text{Var}(X)$ .
7. If  $E(X) = 8$  and  $E(X^2) = 68$ , find a lower bound for  $P\{5 < X < 11\}$ , using Tchebycheff's inequality.
8. State Beinayme's inequality. Deduce Tchebycheff's inequality.
9. State Schwarz inequality.
10. State Cauchy–Schwarz inequality.
11. Use Cauchy–Schwarz inequality to prove that  $|\rho_{XY}| \leq 1$ .
12. Does a RV  $X$  with mean  $\mu$  and SD  $\sigma$  satisfying  $P\{\mu - \sqrt{3}\sigma \leq X \leq \mu + \sqrt{3}\sigma\} = 0.5$  exist? Why?

**Part-B**

13. If  $X$  is a RV with  $E(X) = 3$  and  $E(X^2) = 13$ , find the lower bound for  $P(-2 < X < 8)$ , using Tchebycheff's inequality.
14. If the RV  $X$  is uniformly distributed over  $\left(1 - \frac{1}{\sqrt{3}}, 1 + \frac{1}{\sqrt{3}}\right)$  compute  $P\{|X - \mu| \geq 3/2 \sigma\}$  and compare it with the upper bound obtained by Tchebycheff's inequality.
15. A discrete RV  $X$  can take the values  $-a, 0, a$  with probabilities  $1/8, 3/4, 1/8$  respectively. Compute  $P\{|X - 2\sigma\}$  and compare it with Tchebycheff's inequality bound.
16. Two dice are thrown once. If  $X$  is the sum of the numbers showing up, prove that  $P\{|X - 7| \geq 3\} \leq \frac{35}{54}$  and compare this value with the exact probability.
17. A fair dice is tossed 600 times. Use Tchebycheff's inequality to find a lower bound for the probability of getting 80 to 120 sixes.
18. An unbiased coin is tossed 100 times. Show that the probability that the number of heads will be between 30 to 70 is greater than 0.93.

19. If  $X$  is the number obtained in a throw of a fair dice, show that the Tchebycheff's inequality gives  $P\{|X - \mu| > 2.5\} < 0.47$ , while the actual probability is zero.
20. Using Tchebycheff's inequality, find how many times a fair coin must be tossed in order that the probability that the ratio of the number of heads of the number of tosses will lie between 0.4 and 0.6 will be at least 0.9.
21. A random variable  $X$  has pdf  $f(x) = e^{-x}$ ,  $x \geq 0$ . Use Tchebycheff's inequality to show that  $P\{|X - 1| > 1\} < 1/4$  and show also that the actual probability is  $e^{-3}$ .
22. A random variable  $X$  has the pmf  $P(X = 1) = 1/18$ ,  $P(X = 2) = 16/18$ ,  $P(X = 3) = 1/18$ . Show that there is a value of  $c$  such that  $P\{|X - \mu| \geq c\} = \frac{\sigma^2}{c^2}$ , so that, in general, the bound given by Tchebycheff's inequality cannot be improved.

## Convergence Concepts and Central Limit Theorem

Let us consider the sequence of RVs  $X_1, X_2, X_3, \dots, X_n, \dots$  or random sequence. The concept of convergence of random sequences is essential in the study of random signals. A few definitions and criteria that are used for determining the convergence of random sequences are given below.

### 1. Convergence everywhere and almost everywhere

If  $\{X_n\}$  is a sequence of RVs and  $X$  is a RV such that  $\left[ \lim_{n \rightarrow \infty} (X_n) = X \right]$  i.e.,  $X_n \rightarrow X$  as  $n \rightarrow \infty$ , then the sequence  $\{X_n\}$  is said to converge to  $X$  everywhere.

If  $P\{X_n \rightarrow X\} = 1$  as  $n \rightarrow \infty$ , then the sequence  $\{X_n\}$  is said to converge to  $X$  almost everywhere.

### 2. Convergence in probability or stochastic convergence

If  $P\{|X_n - X| > \varepsilon\} \rightarrow 0$  as  $n \rightarrow \infty$ , then the sequence  $\{X_n\}$  is said to converge to  $X$  in probability or stochastically.

As a particular case of this kind of convergence we have the following result, known as *Bernoulli's law of large numbers*.

If  $X$  represents the number of successes out of  $n$  Bernoulli's trials with probability of success  $p$  (in each trial), then  $\{X/n\}$  converges in probability to  $p$ .

$$\text{i.e., } P\left\{\left|\frac{X}{n} - p\right| > \varepsilon\right\} \rightarrow 0 \text{ as } n \rightarrow \infty$$

### 3. Convergence in the mean square sense

If  $E\{|X_n - X|^2\} \rightarrow 0$  as  $n \rightarrow \infty$ , then the sequence  $\{X_n\}$  is said to converge to  $X$  in the mean square sense.

### 4. Convergence in distribution

If  $F_n(x)$  and  $F(x)$  are the distribution functions of  $X_n$  and  $X$  respectively such that  $F_n(x) \rightarrow F(x)$  as  $n \rightarrow \infty$  for every point of continuity of  $F(x)$ , then the sequence  $\{X_n\}$  is said to converge to  $X$  in distribution.

Closely associated to the concept of convergence in distribution is a remarkable result known as central limit theorem, which is given below without proof.

## Central Limit Theorem (Liapounoff's Form)

If  $X_1, X_2, \dots, X_n, \dots$ , be a sequence of independent RVs with  $E(X_i) = \mu_i$  and  $\text{Var}(X_i) = \sigma_i^2$   $i = 1, 2, \dots$ , and if  $S_n = X_1 + X_2 + \dots + X_n$ , then under certain general conditions,

$S_n$  follows a normal distribution with mean  $\mu = \sum_{i=1}^n \mu_i$  and variance  $\sigma^2 = \sum_{i=1}^n \sigma_i^2$  as  $n$  tends to infinity.

## Central Limit Theorem (Lindeberg-Levy's Form)

If  $X_1, X_2, \dots, X_n, \dots$ , be a sequence of independent identically distributed RVs with  $E(X_i) = \mu$  and  $\text{Var}(X_i) = \sigma^2$ ,  $i = 1, 2, \dots$ , and if  $S_n = X_1 + X_2 + \dots + X_n$ , then under certain general conditions,  $S_n$  follows a normal distribution with mean  $n\mu$  and variance  $n\sigma^2$  as  $n$  tends to infinity.

### Corollary

If  $\bar{X} = \frac{1}{n}(X_1 + X_2 + \dots + X_n)$ , then  $E(\bar{X}) = \mu$  and  $\text{Var}(\bar{X}) = \frac{1}{n^2}(n\sigma^2) = \frac{\sigma^2}{n}$

$\therefore \bar{X}$  follows  $N\left(\mu, \frac{\sigma}{\sqrt{n}}\right)$  as  $n \rightarrow \infty$

### Worked Example 4(F)

#### Example 1

The lifetime of a certain brand of an electric bulb may be considered a RV with mean 1200 h and standard deviation 250 h. Find the probability, using central limit theorem, that the average lifetime of 60 bulbs exceeds 1250 h.

Let  $X_i$  represent the lifetime of the bulb.

$$E(X_i) = 1200 \text{ and } \text{Var}(X_i) = 250^2$$

Let  $\bar{X}$  denote the mean lifetime of 60 bulbs.

By corollary of Lindeberg-Levy form of CLT  $\bar{X}$  follows  $N\left(1200, \frac{250}{\sqrt{60}}\right)$

$$P(\bar{X} > 1250) = P\left(\frac{\bar{X} - 1200}{\frac{250}{\sqrt{60}}} > \frac{1250 - 1200}{\frac{250}{\sqrt{60}}}\right)$$

$$\begin{aligned}
 &= P\left(z > \frac{\sqrt{60}}{5}\right) \\
 &= P(z > 1.55), \text{ where } z \text{ is the standard normal variable} \\
 &= 0.0606 \quad (\text{from the table of areas under normal curve})
 \end{aligned}$$

### Example 2

A distribution with unknown mean  $\mu$  has variance equal to 1.5. Use central limit theorem to find how large a sample should be taken from the distribution in order that the probability will be at least 0.95 that the sample mean will be within 0.5 of the population mean.

Let  $n$  be the size of the sample, a typical member of which is  $X_i$ .

**Given:**  $E(X_i) = \mu$  and  $\text{Var}(X_i) = 1.5$

Let  $\bar{X}$  denote the sample mean.

By corollary under CLT,  $\bar{X}$  follows  $N\left(\mu, \frac{\sqrt{1.5}}{\sqrt{n}}\right)$

We have to find  $n$  such that

$$P\{\mu - 0.5 < \bar{X} < \mu + 0.5\} \geq 0.95$$

$$\text{i.e.,} \quad P\{-0.5 < \bar{X} - \mu < 0.5\} \geq 0.95$$

$$\text{i.e.,} \quad P\{|\bar{X} - \mu| < 0.5\} \geq 0.95$$

$$\text{i.e.,} \quad P\left\{\frac{|\bar{X} - \mu|}{\sqrt{\frac{1.5}{n}}} < \frac{0.5}{\sqrt{\frac{1.5}{n}}}\right\} \geq 0.95$$

$$\text{i.e.,} \quad P\{|z| < 0.4082\sqrt{n}\} \geq 0.95$$

where  $z$  is the standard normal variable.

The least value of  $n$  is obtained from

$$P\{|z| < 0.4082\sqrt{n}\} = 0.95$$

From the table of areas under normal curve

$$P\{|z| < 1.96\} = 0.95$$

Therefore, least  $n$  is given by  $0.4082\sqrt{n} = 1.96$ , i.e., least  $n = 24$ .

Therefore, the size of the sample must be at least 24.

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### Example 3

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If  $X_1, X_2, \dots, X_n$  are Poisson variates with parameter  $\lambda = 2$ , use the central limit theorem to estimate  $P(120 \leq S_n \leq 160)$ , where  $S_n = X_1, X_2, \dots, X_n$  and  $n = 75$ .

$$E(X_i) = \lambda = 2 \text{ and } \text{Var}(X_i) = \lambda = 2$$

By CLT,  $S_n$  follows  $N(n\mu, \sigma\sqrt{n})$

i.e.,  $S_n$  follows  $N(150, \sqrt{150})$

$$\begin{aligned} P\{120 \leq S_n \leq 160\} &= P\left\{\frac{-30}{\sqrt{150}} \leq \frac{S_n - 150}{\sqrt{150}} \leq \frac{10}{\sqrt{150}}\right\} \\ &= P\{-2.45 \leq z \leq 0.85\} \end{aligned}$$

where  $z$  is the standard normal variable

$$\begin{aligned} &= 0.4927 + 0.2939, \text{ (from the normal tables)} \\ &= 0.7866 \end{aligned}$$

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### Example 4

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Using the central limit theorem, show that, for large  $n$ ,

$$\frac{c^n}{|n-1|} x^{n-1} e^{-cx} \cong \frac{c}{\sqrt{2\pi n}} e^{-(cx-n)^2/2n}, x > 0$$

Let  $X_1, X_2, \dots, X_n$  be independent RVs each of which is exponentially distributed with parameter  $c$ .

i.e., let the pdf of  $X_i = c e^{-cx}, x > 0$

The characteristic function of  $X_i$  is given by

$$\phi_{X_i}(\omega) = (1 - i\omega/c)^{-1}$$

(refer to the Note under Worked Example 7 of the characteristic function section).

By Property 4 of CFs, since  $X_1, X_2, \dots, X_n$  are independent RVs,

$$\phi_{(X_1+X_2+\dots+X_n)}(\omega) = [\phi_{x_i}(\omega)]^n$$

$$\begin{aligned} &= \left(1 - \frac{i\omega}{c}\right)^{-n} \\ &= \text{CF of Erlang distribution} \end{aligned}$$

(refer to Worked Example 7 of Section 4(B))

Therefore, when  $n$  is finite,  $(X_1 + X_2 + \dots + X_n)$  follows the Erlang distribution whose pdf is given by

$$\frac{c^n}{[n-1]} x^{n-1} e^{-cx}, x > 0 \quad (1)$$

When  $n$  tends to infinity,  $(X_1 + X_2 + \dots + X_n)$  follows a normal distribution with mean  $nE(X_i) = \frac{n}{c}$  and variance  $n\text{Var}(X_i) = \frac{n}{c^2}$  (Central limit theorem), i.e., when  $n \rightarrow \infty$ ,  $(X_1 + X_2 + \dots + X_n)$  follows a normal distribution whose pdf is given by

$$\begin{aligned} & \frac{1}{\sqrt{\frac{n}{c^2} \sqrt{2\pi}}} \exp \left\{ -\left( x - \frac{n}{c} \right)^2 / \frac{2n}{c^2} \right\} \\ &= \frac{c}{\sqrt{2n\pi}} \exp \{ -(cx - n)^2 / 2n \} \end{aligned} \quad (2)$$

From (1) and (2), the required result follows.

### Example 5

Verify central limit theorem for the independent random variables  $X_k$ , where for each  $k$ ,  $P\{X_k = \pm 1\} = \frac{1}{2}$ .

$$E(X_k) = 1 \times 1/2 + (-1) \times 1/2 = 0$$

$$\text{Var}(X_k) = 1^2 \times 1/2 + (-1)^2 \times 1/2 = 1$$

$$\text{Consider } Y_n = \frac{1}{\sqrt{n}}(X_1 + X_2 + \dots + X_n)$$

$$E(Y_n) = 0 \text{ and } \text{Var}(Y_n) = \frac{1}{n} \times n = 1$$

$$\text{Now } \phi_{X_k}(\omega) = E\{e^{i\omega X_k}\}$$

$$\begin{aligned} &= e^{i\omega} \times 1/2 + e^{i\omega(-1)} \times 1/2 \\ &= \cos \omega \end{aligned}$$

$$\begin{aligned} \therefore \phi_{Y_n}(\omega) &= \phi_{\frac{1}{\sqrt{n}}(X_1 + X_2 + \dots + X_n)}(\omega) \\ &= \left\{ \phi_{X_k/\sqrt{n}}(\omega) \right\}^n \quad (\text{since } X_1, X_2, \dots, X_n \text{ are independent}) \\ &= \left[ \cos \left( \frac{\omega}{\sqrt{n}} \right) \right]^n \quad [\text{since } \phi_{aX+b}(\omega) = e^{ib\omega} \phi_X(a\omega)] \end{aligned}$$

$$\begin{aligned}
&= \left[ 1 - \frac{\omega^2}{2n} + \text{terms involving } \frac{1}{n^2} \text{ and higher powers of } \frac{1}{n} \right]^n \\
&= \left( 1 - \frac{\omega^2}{2n} \right)^n + \text{term involving } \frac{1}{n} \text{ and higher powers of } \frac{1}{n} \\
\therefore \quad \lim_{n \rightarrow \infty} \phi_{y_n}(\omega) &= e^{\frac{-1}{2}\omega^2}
\end{aligned}$$

which is the characteristic function of  $N(0, 1)$ .

$$\therefore Y_n \rightarrow N(0, 1), \text{ as } n \rightarrow \infty$$

Therefore, CLT holds good for the sequence  $\{X_k\}$ .

### Example 6

Show that the central limit theorem holds good for the sequence  $\{X_k\}$ , if

$$P\{X_k = \pm k^\alpha\} = \frac{1}{2} \times k^{-2\alpha}, P\{X_k = 0\} = 1 - k^{-2\alpha}, \alpha < \frac{1}{2}.$$

**Note**

Liaponoff's form of CLT holds good for a sequence  $\{X_k\}$ , if

$$\lim_{n \rightarrow \infty} \frac{\sum_{k=1}^n E[|X_k - \mu_k|^3]}{\left\{ \sum_{k=1}^n Var(X_k) \right\}^{\frac{3}{2}}} = 0 \quad \text{Condition is assumed.}$$

We have to verify whether this condition is satisfied by the given  $\{X_k\}$ .

Now

$$\mu_k = E(X_k) = \frac{1}{2} \times k^{-\alpha} - \frac{1}{2} \times k^{-\alpha} = 0$$

$$E(X_k^2) = \frac{1}{2} \times k^{2\alpha} \times k^{-2\alpha} + \frac{1}{2} k^{2\alpha} \times k^{-2\alpha} = 1$$

$\therefore$

$$\text{Var}(X_k) = 1$$

$$\begin{aligned}
E\{|X_k - \mu_k|^3| &= E\{|X_k|^3\} \\
&= k^{3\alpha} \times \frac{1}{2} k^{-2\alpha} + k^{3\alpha} \times \frac{1}{2} k^{-2\alpha} \\
&= k^\alpha
\end{aligned}$$

$$\therefore \sum_{k=1}^n E\{|X_k - \mu_k|^3\} = 1^\alpha + 2^\alpha + \dots + n^\alpha < n \times n^\alpha \quad (\text{since each term } \leq n^\alpha)$$

$$\sum_{k=1}^n \text{Var}(X_k) = n$$

$$\begin{aligned}
 & \sum_{k=1}^n E\{|X_k - \mu_k|^3\} \\
 \therefore \lim_{n \rightarrow \infty} & \frac{\left\{ \sum_{k=1}^n \text{Var}(X_k) \right\}^{3/2}}{\left( \frac{1^\alpha + 2^\alpha + \dots + n^\alpha}{n^{3/2}} \right)} \\
 & = \lim_{n \rightarrow \infty} \frac{1}{\frac{1}{n^{2-\alpha}}} \\
 & = 0 \text{ (since } \alpha < 1/2)
 \end{aligned}$$

i.e., the necessary condition is satisfied. Therefore, CLT holds good for the sequence  $\{X_k\}$ .

### Exercise 4(F)

#### **Part-A (Short-answer Questions)**

1. What is the difference between convergence everywhere and almost everywhere of a random sequence  $\{X_n\}$ ?
2. Define stochastic convergence of a random sequence  $\{X_n\}$ .
3. State Bernoulli's law of large numbers.
4. Define convergence of a random sequence  $\{x_n\}$  in the mean square sense.
5. Define convergence in distribution of a random sequence  $\{X_n\}$ .
6. State the Liapounoff's form of CLT.
7. State the Lindeberg-Levy's form of CLT.
8. What is the importance of CLT?

#### **Part-B**

9. A random sample of size 100 is taken from a population whose mean is 60 and variance is 400. Using CLT, with what probability can we assert that the mean of the sample will not differ from  $\mu = 60$  by more than 4?
10. The guaranteed average life of a certain type of electric light bulb is 1000 h with a standard deviation of 125 h. It is decided to sample the output so as to ensure that 90% of the bulbs do not fall short of the guaranteed average by more than 2.5%. Use CLT to find the minimum sample size.
11. If  $X_i$ ,  $i = 1, 2, \dots, 50$ , are independent RVs, each having a Poisson distribution with parameter  $\lambda = 0.03$  and  $S_n = X_1 + X_2 + \dots + X_n$ , evaluate  $P(S_n \geq 3)$ , using CLT. Compare your answer with the exact value of the probability.
12. If  $V_i$ ,  $i = 1, 2, \dots, 20$ , are independent noise voltages received in an 'adder' and  $V$  is the sum of the voltages received, find the probability that the

- total incoming voltage  $V$  exceeds 105, using CLT. Assume that each of the random variables  $V_i$  is uniformly distributed over  $(0, 10)$ .
13. 30 electronics devices  $D_1, D_1, \dots, D_{30}$  are used in the following manner. As soon as  $D_1$  fails,  $D_2$  becomes operative. When  $D_2$  fails,  $D_3$  becomes operative and so on. If the time to failure of  $D_i$  is an exponentially distributed random variable with parameter  $\lambda = 0.1/\text{h}$  and  $T$  is the total time of operation of all the 30 devices, find the probability that  $T$  exceeds 350 h, using CLT.
  14. Examine if the CLT holds good for the sequence  $\{X_k\}$ , if  $P\{X_k = \pm 2^k\} = 2^{-(2k+1)}$ ,  $P\{X_k = 0\} = 1 - 2^{-2k}$ .
  15. Show that the CLT does not hold if the RV's  $X_i$  have a Cauchy density.

## **ANSWERS**

### **Exercise 4(A)**

4.  $\text{Var}(X) = E\{X - E(X)\}^2 \geq 0$   
 $E(X^2) - \{E(X)\}^2 \geq 0$   
 $E(X^2) \geq \{E(X)\}^2$
6.  $E(Y) = \frac{1}{\sigma_X} \{E(X) - \mu_X\} = 0$   
 $\text{Var}(Y) = \frac{1}{\sigma_X^2} \times \text{Var}(X) = 1$   
 $\text{SD} = 1$
7.  $\mu_k = \mu'_k - kc_1 \mu'_{k-1} D + k c_2 \mu'_{k-2} D^2 + \dots + (-1)^k kc_k D^k$ , where  $D = \mu'_1$
8.  $Y:$       0            2            6            12  
 $P_Y:$      0.1        0.3        0.4        0.2  
 $E(Y) = 0.6 + 2.4 + 2.4 = 5.4$
9.  $E(X) = \sum_{j=0}^{\infty} j(1-a)a^j = (1-a)(a + 2a^2 + 3a^3 + \dots)$   
 $= (1-a) \frac{a}{(1-a)^2} = \frac{a}{1-a}$
10.  $E(X) = 6 \int_0^1 x^2(1-x)dx = 6 \left( \frac{x^3}{3} - \frac{x^4}{4} \right) = \frac{1}{2}$
11.  $E(X) = \int_a^b \frac{1}{b-a} x dx = \frac{b+a}{2}$   
 $E(X^2) = \int_a^b \frac{1}{b-a} x^2 dx = \frac{1}{3}(b^2 + ba + a^2)$

$$\text{Var}(X) = E(X^2) - E^2(X) = \frac{(b-a)^2}{12}$$

$$12. \quad E(X) = \frac{b+a}{2} = 5; \text{Var}(X) = \frac{(b-a)^2}{12} = 3$$

$$13. \quad E(Y) = E(X^3) = \int_1^2 x^3 \, dx = \frac{15}{4}$$

$$14. \quad E(X) = \sum_{r=0}^n r \frac{n!}{r!(n-r)!} p^r q^{n-r}$$

$$= np \sum_{r=1}^n \frac{(n-1)!}{(r-1)!(n-r)!} p^{r-1} q^{n-1} = np$$

$$15. \quad E(X) = \sum_{r=0}^{\infty} r \frac{e^{-\lambda} \lambda^r}{r!} = e^{-\lambda} \lambda \sum_{r=1}^{\infty} \frac{\lambda^{r-1}}{(r-1)!} = \lambda$$

$$16. \quad \frac{npq}{np} = \frac{2}{6} = \frac{1}{3}; p = \frac{2}{3}; \text{and } n = 9$$

$$P(X=r) = 9C_r \left(\frac{2}{3}\right)^r \times \left(\frac{1}{3}\right)^{9-r} \quad r = 0, 1, \dots, 9$$

$$17. \quad \frac{npq}{np} = \frac{1.44}{2.4} = 0.6; p = 0.4; n = 6;$$

Required  $P = 0$

$$18. \quad 5 C_1 p^1 q^4 = 2 \times 5 C_2 P^2 q^3; q = 4p; p = 0.2$$

$$E(X) = 1 \text{ and } \text{Var}(X) = 0.8$$

$$19. \quad 6C_2 p^2 q^4 = 9 \times 6C_4 P^4 q^2; q = 3p; p = 1/4$$

$$E(X) = 1.5 \text{ and } \text{Var}(X) = \frac{9}{8}$$

$$20. \quad \frac{e^{-\lambda} \lambda}{1} = \frac{e^{-\lambda} \lambda^2}{2}; \lambda = 2$$

$$E(X) = 2 = \text{Var}(X)$$

$$E(X^2) = V(X) + E^2(X) = 6$$

$$21. \quad E(X) = \sum_{r=0}^{\infty} r p q^r = pq(1 + 2q + 3q^2 + \dots) = \frac{q}{p}$$

22. If  $X$  represents the number of times, the distribution of  $X$  is

$$x: \quad 1 \quad 2 \quad 3 \quad 4 \quad \dots$$

$$p_x: \quad p \quad qp \quad q^2 p \quad q^3 p \quad \dots, \quad \text{where } p = \frac{1}{6}$$

- $E(X) = p(1 + 2q + 3q^2 + \dots) = \frac{p}{(1-q)^2} = \frac{1}{p} = 6$
23.  $E(X) = \int_0^\infty \lambda x e^{-\lambda x} dx = \frac{1}{\lambda}$   
 $E(X^2) = \int_0^\infty \lambda x^2 e^{-\lambda x} dx = \frac{2}{\lambda^2}$   
 $\text{Var}(X) = E(X^2) - E^2(X) = \frac{1}{\lambda^2}$
24.  $E(X) = 0; \text{Var}(X) = 4; E(X^2) = \text{Var}(X) + E^2(X) = 4$
28.  $E(Z) = 2E(X) - 5E(Y) = 4 - 15 = -11$   
 $\text{Var}(Z) = E(2X - 5Y)^2 - \{E(2X - 5Y)\}^2$   
 $= 4E(X^2) + 25E(Y^2) - 20E(XY)$   
 $- \{4E^2(X) + 25E^2(Y) - 20E(X) \times E(Y)\}$   
 $= 4 \times 1 + 25 \times 2 - 20 \times 0 = 54$
34.  $E(X) = \iint_0^y x(2-x-y) dx dy = \int_0^1 \left( \frac{2}{3} - \frac{y}{2} \right) dy = \frac{5}{12}$   
 $E(Y) = \frac{5}{12}$
35.  $E(X) = \iint_0^y 2x dx dy = \int_0^1 y^2 dy = \frac{1}{3}$
36.  $\frac{1-p}{p}$
37. (i)  $\frac{7n}{2}$       (ii)  $\left(\frac{7}{2}\right)^n$
38. 151.5
39. Rs. 32/-
40.  $\frac{b+a}{2}; \frac{1}{12}(b-a)^2$
41.  $\frac{q}{p}; \frac{2q}{p^2}$
42.  $np; npq$
43.  $\lambda; \lambda$
47. (ii) 0, 1
48.  $a = \frac{1}{5}, b = 2$

49.  $\frac{15}{4}, \frac{457}{112}$

50. 1, 1

51.  $\frac{8}{3}$

53.  $\frac{y}{2}; \frac{y^2}{12}$

54.  $\frac{3y}{4}; \frac{3}{80}y^2$

### Exercise 4(B)

12. No

15. 0.7891

16. 0.6030

17. 0.0390

18. 0.8504

19. 0.6

20. 0.8762

21. 0.9173

22. 0.3875

23. 0.91

24. 0.4

25. A and C

26.  $\frac{1}{2}$

28. 0

29.  $\frac{v_2}{\sqrt{v_1 + v_2}(v_2 + v_3)}$

32.  $\frac{5}{12}; \frac{5}{12}; \frac{1}{16}; \frac{1}{16}; -\frac{1}{9}$

33.  $\frac{1}{3}; \frac{2}{3}; \frac{1}{18}; \frac{1}{18}; -\frac{1}{2}$

34. 0

35.  $-\frac{1}{11}$

36. 11; 4; 0

### Exercise 4(C)

11.  $y = 6x + 114; 20x = 3y + 950; 0.9$

12.  $y = 0.665x + 23.78; x = 0.54y + 30.74; 70.995, 68.54$

13.  $y = 0.7x + 12.7$ ; 25.3  
 14.  $x = 0.936y + 1.388$ ; 43.508  
 15. 72.6; 69.57  
 16. 13; 17; 0.6; 4  
 17.  $-\frac{1}{2\sqrt{3}}; 5; \frac{1}{3}; \frac{2}{15\sqrt{3}}$   
 18. 1; 2; 4; -0.8.66  
 19.  $\bar{x} = \frac{a+bc}{1-bd}$ ;  $\bar{y} = \frac{c+ad}{1-bd}$ ;  $r = \sqrt{bd}$ ; No  
 20.  $\sqrt{\frac{a_1 b_2}{a_2 b_1}}; \left( \frac{b_1 c_2 - b_2 c_1}{c_1 a_2 - c_2 a_1} \right) \sqrt{\frac{a_1 a_2}{b_1 b_2}}$

### Exercise 4(D)

5.  $M(t) = E(e^{tX})$ , whereas  $\phi(\omega) = E(e^{i\omega X})$   
 7. The characteristic function of a RV uniquely determines its pdf. (or) The necessary and sufficient condition for two distributions with pdf's  $f_1(x)$  and  $f_2(x)$  to be identical is that their characteristic functions  $\phi_1(\omega)$  and  $\phi_2(\omega)$  are identical.  
 11. It generates cumulants.
12.  $M(t) = \sum_{r=0}^{\infty} E(X^r) t^r / r! = \sum_{r=0}^{\infty} t^r = \frac{1}{1-t}$   
 13.  $\frac{2}{2-t} = \left(1 - \frac{t}{2}\right)^{-1} = 1 + \frac{t}{2} + \frac{t^2}{4} + \dots$   
 $E(X) = \frac{1}{2}; E(X^2) = \frac{1}{2}; \text{Var}(X) = \frac{1}{4}; \sigma_x = \frac{1}{2}$   
 14.  $M(t) = E(e^{tX}) = \int_a^b \frac{1}{b-a} e^{tx} dx$   
 $= \frac{1}{b-a} \left( \frac{e^{bt} - e^{at}}{t} \right)$   
 $\phi(\omega) = \frac{1}{b-a} \left( \frac{e^{ib\omega} - e^{ia\omega}}{i\omega} \right)$   
 15.  $M(t) = \sum_{r=0}^n nC_r p^r q^{n-r} e^{tr} = (q + pe^t)^n$   
 $\phi(\omega) = (q + pe^{i\omega})^n$

$$16. \quad M(t) = \sum_{r=0}^{\infty} \frac{e^{-\lambda} \lambda^r}{r!} e^{tr} = e^{-\lambda(1-e^t)}$$

$$\phi(\omega) = e^{-\lambda(1-e^{i\omega})}$$

$$18. \quad \text{CF} = \frac{p}{1-qe^{i\omega}}$$

$$\text{MGF} = \frac{p}{1-qe^t}$$

$$19. \quad \phi(\omega) = \int_0^{\infty} \lambda e^{-\lambda x} e^{i\omega x} dx = \lambda \left\{ \frac{e^{-(\lambda-i\omega)x}}{-(\lambda-i\omega)} \right\}_0^{\infty} = \frac{\lambda}{\lambda-i\omega}$$

$$M(t) = \frac{\lambda}{\lambda-t}$$

$$20. \quad \phi(\omega) = e^{i\mu\omega} e^{-\sigma^2\omega^2/2}$$

$$21. \quad \phi(\omega) = \frac{1}{2} \int_{-\infty}^{\infty} e^{-|x|} e^{i\omega x} dx \\ = \frac{1}{2} \left[ \int_{-\infty}^0 e^{(1+i\omega)x} dx + \int_0^{\infty} e^{-(1-i\omega)x} dx \right] \\ = \frac{1}{2} \left[ \frac{1}{1+i\omega} + \frac{1}{1-i\omega} \right] = \frac{1}{1+\omega^2}$$

### Part-B

$$22. \quad \frac{e^{3t}-1}{3te^t}; \frac{1}{2}, \frac{3}{4}$$

$$23. \quad (pe^{i\omega} + q)^n; np; npq$$

$$24. \quad e^\lambda (e^t - 1); 0; \lambda; \lambda; 3\lambda^2 + \lambda$$

$$25. \quad p^n (1-qe^{i\omega})^{-n}; \frac{nq}{p}; \frac{nq}{p^2}$$

$$26. \quad \frac{\lambda e^{at}}{\lambda-t}; \frac{a\lambda+1}{\lambda}; \frac{1}{\lambda^2}$$

$$27. \quad \frac{1}{1-t^2}; 0; 2$$

$$28. \quad e^{-|\omega|}; \mu_X = 0; \sigma_X \text{ does not exist.}$$

29.

31.  $f(x) = \frac{1}{\sigma\sqrt{2\pi}} e^{-x^2/2\sigma^2} \quad -\infty < x < \infty$

32.  $\phi(\omega) = \frac{1}{1-j\omega}; f(x) = e^{-x}, 0 < x < \infty$

33.  $\mu_1 = \lambda_1; \mu_2 = \lambda_2 + \lambda_1^2; \mu_3 = \lambda_3 + 3\lambda_2\lambda_1 + \lambda_1^3$  and

$$\mu_4 = \lambda_4 + 4\lambda_3\lambda_1 + 3\lambda_2^2 + 6\lambda_2\lambda_1^2 + \lambda_1^4$$

35.  $\phi_{XY}(\omega_1, \omega_2) = e^{j(\mu_x\omega_1 + \mu_y\omega_2)} - \frac{1}{2}(\sigma_x^2\omega_1^2 + 2r\sigma_x\sigma_y\omega_1\omega_2 + \sigma_y^2\omega_2^2)$

36.  $f_Y(y) = \frac{1}{\pi} \cdot \frac{1}{\sqrt{1-y^2}}$  in  $-1 \leq y \leq 1$

### Exercise 4(E)

2.  $P\{10 - c < X < 10 + c\} \geq 1 - \frac{16}{c^2}; c = 5$

$$P\{5 < X < 15\} \geq \frac{9}{25}$$

3.  $P\{|x - \mu| \geq k\sigma\} \leq \frac{1}{k^2}$

$$k = \frac{5}{4}$$

$$P\{|X - 10| \geq 5\} \leq \frac{16}{25}$$

4.  $P\{12 - c < X < 12 + c\} \geq 1 - \frac{\sigma^2}{c^2}; c = 6$

$$1 - \frac{\sigma^2}{36} = \frac{3}{4}; \sigma = 3$$

5.  $P\{\mu - c < X < \mu + c\} \geq 1 - \frac{\sigma^2}{c^2}; 1 - \frac{9}{c^2} = \frac{3}{4}; c = 6; \mu = 12$

6.  $P\{\mu - c < X < \mu + c\} \geq 1 - \frac{\sigma^2}{c^2}$

$$\mu = \frac{-2 + 8}{2} = 3; c = 5$$

$$1 - \frac{\sigma^2}{25} = \frac{21}{25}; \sigma^2 = 4$$

7.  $\text{Var}(X) = 68 - 64 = 4$

$$P\{8 - c < X < 8 + c\} \geq 1 - \frac{4}{c^2}; c = 3$$

$$\text{Required Probability} \geq \frac{5}{9}$$

11.  $\{E(XY)^2 \leq E(X^2) E(Y^2)\}$ : Replace  $X$  by  $X - \mu_x$ , and  $Y$  by  $Y - \mu_y$

$$\frac{[E(X - \mu_x)(Y - \mu_y)]^2}{E(X - \mu_x)^2 E(Y - \mu_y)^2} \leq 1; p_{xy}^2 \leq 1; |p_{xy}| \leq 1$$

12. By Tchebycheff's inequality,  $P\{|X - \mu| \leq \sqrt{3}\sigma\} \geq 1 - \frac{1}{3} = \left(\frac{2}{3}\right)$ . Since  $0.5 < \frac{2}{3}$ , such a RV does not exist.

13. 21/25

14. 4/9; 0.134

15. 1/4; 1/4

16. 1/3

17. 19/24

20. 250

22.  $c = 1$

### Exercise 4(F)

8. If  $X_1, X_2, \dots, X_n$  be a sequence of independent and identically distributed RVs with  $E(X_i) = \mu$  and  $\text{Var}(X_i) = \sigma^2$  and if  $\bar{X} = \frac{1}{n} (X_1 + X_2 + \dots + X_n)$ , then  $\bar{X}$  follows  $N\left(\mu, \frac{\sigma}{\sqrt{n}}\right)$  as  $n \rightarrow \infty$ . This result is used in theory of sampling.

9. 0.9544

10. 41

11. 0.1112; 0.1915

12. 0.352

13. 0.1814

14. CLT does not hold good

# *Chapter* 5

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# Some Special Probability Distributions

## Introduction

In Chapter 2, we have just stated the definitions of certain special probability distributions, both discrete and continuous. While constructing probabilistic models for observable phenomena, certain probability distributions arise more frequently than do others. We treat such distributions that play important roles in many engineering applications as special probability distributions. In this chapter we shall discuss a number of discrete as well as continuous (random variables) distributions in considerable detail.

## SPECIAL DISCRETE DISTRIBUTIONS

### 1. Binomial Distribution

**Definition:** Let  $A$  be some event associated with a random experiment  $E$ , such that  $P(A) = p$  and  $P(\bar{A}) = 1 - p = q$ . Assuming that  $p$  remains the same for all repetitions, if we consider  $n$  independent repetitions (or trials) of  $E$  and if the random variable (RV)  $X$  denotes the number of times the event  $A$  has occurred, then  $X$  is called a *binomial random variable* with parameters  $n$  and  $p$  or we say that  $X$  follows a *binomial distribution* with parameters  $n$  and  $p$ , or symbolically  $B(n, p)$ . Obviously the possible values that  $X$  can take, are  $0, 1, 2, \dots, n$ .

By the theorem under Bernoulli's trials in Chapter 1, the probability mass function of a binomial RV is given by

$$P(X = r) = nC_r p^r q^{n-r}; r = 0, 1, 2, \dots, n \quad \text{where } p + q = 1$$

**Note**

1. Binomial distribution is a legitimate probability distribution since

$$\sum_{r=0}^n P(X = r) = \sum_{r=0}^n nC_r q^{n-r} p^r \\ = (q + p)^n = 1$$

2. The name 'binomial distribution' is given since the probabilities  $nC_r q^{n-r} p^r$  ( $r = 0, 1, 2, \dots, n$ ) are the successive terms in the expansion of the binomial expression  $(q + p)^n$ .
3. If we assume that  $n$  trials constitute a set and if we consider  $N$  sets, the frequency function of the binomial distribution is given by  $f(r) = Np(r) = N \cdot nC_r q^{n-r} p^r$ ,  $r = 0, 1, 2, \dots, n$ . In other words, the number of sets in which we get exactly  $r$  successes (the occurrences of the event  $A$ ) =  $N \cdot nC_r q^{n-r} p^r$ ;  $r = 0, 1, 2, \dots, n$ .

### Mean and Variance of the Binomial Distribution

We have already found out  $E(X)$  and  $\text{Var}(X)$  for the binomial distribution  $B(n, p)$  using the moment generating function in Example 3 in Worked Example 4(b). Here, we shall find them directly using the definitions of  $E(X)$  and  $\text{Var}(X)$ .

$$E(X) = \sum_r x_r p_r \\ = \sum_{r=0}^n r \cdot nC_r p^r q^{n-r} \\ = \sum_{r=0}^n r \cdot \frac{n!}{r!(n-r)!} p^r q^{n-r} \quad (1)$$

$$= np \cdot \sum_{r=1}^n r \frac{(n-1)!}{(r-1)!\{(n-1)-(r-1)\}!} p^{r-1} q^{(n-1)-(r-1)} \\ = np \sum_{r=1}^n (n-1) C_{r-1} \cdot p^{r-1} \cdot q^{(n-1)-(r-1)} \\ = np (q+p)^{n-1} \\ = np \quad (2)$$

$$E(X^2) = \sum_r x_r^2 p_r = \sum_0^n r^2 p_r \\ = \sum_{r=0}^n \{r(r-1) + r\} \frac{n!}{r!(n-r)!} p^r q^{n-r} \\ = n(n-1)p^2 \sum_{r=2}^n (n-2)C_{r-2} p^{r-2} q^{n-r} + np, \quad [\text{by (1) and (2)}] \\ = n(n-1)p^2 (q+p)^{n-2} + np$$

$$\begin{aligned}
&= n(n-1)p^2 + np \\
\text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
&= n(n-1)p^2 + np - n^2p^2 \\
&= np(1-p) \\
&= npq
\end{aligned}$$

### **Recurrence Formula for the Central Moments of the Binomial Distribution**

By definition, the  $k$ th order central moment  $\mu_k$  is given by  $\mu_k = E\{X - E(X)\}^k$ .

For the binomial distribution  $B(n, p)$ ,

$$\mu_k = \sum_{r=0}^n (r-np)^k nC_r p^r q^{n-r} \quad (1)$$

Differentiating (1) with respect to  $p$ , we get

$$\begin{aligned}
\frac{d\mu_k}{dp} &= \sum_{r=0}^n nC_r [-nk(r-np)^{k-1} \cdot p^r q^{n-r} + (r-np)^k \\
&\quad \{rp^{r-1} q^{n-r} + (n-r) p^r q^{n-r-1} (-1)\}] \\
&= -nk \mu_{k-1} + \sum_{r=0}^n nC_r (r-np)^k p^{r-1} q^{n-r-1} \{rq - (n-r)p\} \\
&= -nk \mu_{k-1} + \sum_{r=0}^n nC_r (r-np)^k p^{r-1} q^{n-r-1} (r-np) \quad (\because p+q=1) \\
&= -nk \mu_{k-1} + \frac{1}{pq} \sum_{r=0}^n nC_r p^r q^{n-r} (r-np)^{k+1} \\
&= -nk \mu_{k-1} + \frac{1}{pq} \mu_{k+1} \\
\text{i.e., } \mu_{k+1} &= pq \left[ \frac{d\mu_k}{dp} + nk \mu_{k-1} \right] \quad (2)
\end{aligned}$$

Using the recurrence relation (2), we may compute moments of higher order, provided we know moments of lower order.

Putting  $k = 1$  in (2), we get

$$\begin{aligned}
\mu_2 &= pq \left[ \frac{d\mu_1}{dp} + n\mu_0 \right] \\
&= npq \quad (\because \mu_0 = 1 \text{ and } \mu_1 = 0)
\end{aligned}$$

Putting  $k = 2$  in (2), we get

$$\mu_3 = pq \left[ \frac{d}{dp} \mu_2 + 2n \mu_1 \right]$$

$$\begin{aligned}
 &= pq \frac{d}{dp} [np(1-p)] \\
 &= npq [1 - 2p] = npq (q - p)
 \end{aligned}$$

Putting  $k = 3$  in (2), we get

$$\begin{aligned}
 \mu_4 &= pq \left[ \frac{d}{dp} \mu_3 + 3n \mu_2 \right] \\
 &= npq \left[ \frac{d}{dp} \{p(1-p)(1-2p)\} + 3npq \right] \\
 &= npq [1 - 6p + 6p^2 + 3npq] \\
 &= npq [1 - 6pq + 3npq] \\
 &= npq [1 + 3pq(n-2)]
 \end{aligned}$$

**Note**  $\mu_2$  is the variance,  $\mu_3$  is a measure of skewness and  $\mu_4$  is a measure of kurtosis. Sometimes the coefficients  $\beta_1$  and  $\beta_2$  are used to measure skewness and kurtosis respectively,

where  $\mu_1 = \frac{\mu_3^2}{\mu_2^3}$  and  $\mu_2 = \frac{\mu_4}{\mu_2^2}$ .

## 2. Poisson Distribution

**Definition:** If  $X$  is a discrete RV that can assume the values  $0, 1, 2, \dots$ , such that its probability mass function is given by

$$P(X = r) = \frac{e^{-\lambda} \lambda^r}{r!}; \quad r = 0, 1, 2, \dots; \lambda > 0$$

then  $X$  is said to follow a *Poisson distribution* with parameter  $\lambda$  or symbolically  $X$  is said to follow  $P(\lambda)$ .

**Note** Poisson distribution is a legitimate probability distribution, since

$$\begin{aligned}
 \sum_{r=0}^{\infty} P(x = r) &= \sum_{r=0}^{\infty} \frac{e^{-\lambda} \lambda^r}{r!} \\
 &= e^{-\lambda} e^{\lambda} = 1
 \end{aligned}$$

### Poisson Distribution as Limiting form of Binomial Distribution

Poisson distribution is a limiting case of binomial distribution under the following conditions:

- (i)  $n$ , the number of trials is indefinitely large, i.e.,  $n \rightarrow \infty$ .
- (ii)  $p$ , the constant probability of success in each trial is very small, i.e.,  $p \rightarrow 0$ .

(iii)  $np (= \lambda)$  is finite or  $p = \frac{\lambda}{n}$  and  $q = 1 - \frac{\lambda}{n}$ , where  $\lambda$  is a positive real number.

### **Proof**

If  $X$  is a binomially distributed RV with parameters  $n$  and  $p$ , then

$$\begin{aligned}
 P(X = r) &= nC_r p^r q^{n-r}; r = 0, 1, 2, \dots, n. \\
 &= \frac{n(n-1)(n-2)\cdots(n-r+1)}{r!} p^r (1-p)^{n-r} \\
 &= \frac{n(n-1)\cdots(n-r+1)}{r!} \left(\frac{\lambda}{n}\right)^r \left(1 - \frac{\lambda}{n}\right)^{n-r} \quad \left(\text{on putting } p = \frac{\lambda}{n}\right) \\
 &= \frac{\lambda^r}{r!} \left[1\left(1 - \frac{1}{n}\right)\left(1 - \frac{2}{n}\right)\cdots\left(1 - \frac{r-1}{n}\right)\right] \left(1 - \frac{\lambda}{n}\right)^n \cdot \left(1 - \frac{\lambda}{n}\right)^{-r} \\
 \therefore \lim_{\substack{n \rightarrow \infty \\ np = \lambda \text{ finite}}} [P(X = r)] &= \frac{\lambda^r}{r!} \lim_{n \rightarrow \infty} \left[ \left(1 - \frac{1}{n}\right)\left(1 - \frac{2}{n}\right)\cdots\left(1 - \frac{r-1}{n}\right) \right] \\
 \lim_{n \rightarrow \infty} \left(1 - \frac{\lambda}{n}\right)^{-r} \cdot \lim_{n \rightarrow \infty} \left(1 - \frac{\lambda}{n}\right)^n &= \frac{\lambda'}{r!} e^{-\lambda} \\
 \left[ \because \lim_{n \rightarrow \infty} \left(1 - \frac{k}{n}\right) = 1, \text{ when } k \text{ is finite, } \lim_{n \rightarrow \infty} \left(1 - \frac{\lambda}{n}\right)^{-r} = 1 \text{ and } \lim_{n \rightarrow \infty} \left(1 - \frac{\lambda}{n}\right)^n = e^{-\lambda} \right]
 \end{aligned}$$

$$\text{But } P(X = r) = \frac{e^{-\lambda} \cdot \lambda^r}{r!}, \quad r = 0, 1, 2, \dots$$

is the probability mass function of the Poisson random variable.

Thus, the limit of the binomial RV (distribution) is the Poisson RV (distribution).

### **Note**

1. The above result means that we may compute binomial probabilities approximately by using the corresponding Poisson probabilities, whenever  $n$  is large and  $p$  is small.
2. When an event occurs rarely, the number of occurrences of such an event may be assumed to follow a Poisson distribution. The following are some of the examples, which may be analysed using Poisson distribution:
  - (i) the number of alpha particles emitted by a radioactive source in a given time interval
  - (ii) the number of telephone calls received at a telephone exchange in a given time interval
  - (iii) the number of defective articles in a packet of 100
  - (iv) the number of printing errors at each page of a book
  - (v) the number of road accidents reported in a city per day.

## Mean and Variance of Poisson Distribution

We have already found out  $E(X)$  and  $\text{Var}(X)$  for the Poisson distribution  $P(\lambda)$ , using the characteristic function in Example 4 in Worked Example 4(b). Also, since the Poisson distribution is the limit of binomial distribution, the mean and variance of the Poisson distribution may be obtained as the limits of those of binomial distribution when  $n \rightarrow \infty$ ,

i.e., if  $X$  is the Poisson RV

$$E(X) = \lim_{\substack{n \rightarrow \infty \\ np = \lambda}} (np) = \lambda$$

and  $\text{Var}(X) = \lim_{\substack{p \rightarrow 0 \\ np = \lambda}} (npq) = \lim_{p \rightarrow 0} [\lambda(1 - p)] = \lambda$

Now we shall find  $E(X)$  and  $\text{Var}(X)$  for the Poisson distribution directly using the definitions.

$$\begin{aligned} E(X) &= \sum_r x_r p_r \\ &= \sum_{r=0}^{\infty} r \frac{e^{-\lambda} \cdot \lambda^r}{r!} \end{aligned} \tag{1}$$

$$\begin{aligned} &= \lambda e^{-\lambda} \sum_{r=1}^{\infty} \frac{\lambda^{r-1}}{(r-1)!} \\ &= \lambda e^{-\lambda} e^{\lambda} = \lambda \end{aligned} \tag{2}$$

$$\begin{aligned} E(X^2) &= \sum_r x_r^2 p_r \\ &= \sum_{r=0}^{\infty} \{r(r-1) + r\} e^{-\lambda} \frac{\lambda^r}{r!} \\ &= \lambda^2 e^{-\lambda} \sum_{r=2}^{\infty} \frac{\lambda^{r-2}}{(r-2)!} + \lambda \end{aligned} \tag{[by (1) and (2)]}$$

$$= \lambda^2 e^{-\lambda} e^{\lambda} + \lambda = \lambda^2 + \lambda$$

$$\begin{aligned} \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\ &= \lambda^2 + \lambda - \lambda^2 = \lambda \end{aligned}$$

## Recurrence Formula for the Central Moments of the Poisson Distribution

By definition, the  $k$ th order central moment  $\mu_k$  is given by  $\mu_k = E\{X - E(X)\}^k$ .

For the Poisson distribution  $P(\lambda)$ ,

$$\mu_k = \sum_{r=0}^{\infty} (r - \lambda)^k e^{-\lambda} \frac{\lambda^r}{r!} \tag{1}$$

Differentiating (1) with respect to  $\lambda$ , we get

$$\begin{aligned}\frac{d\mu_k}{d\lambda} &= \sum_{r=0}^{\infty} \frac{1}{r!} [-k(r-\lambda)^{k-1} e^{-\lambda} \lambda^r + (re^{-\lambda} \lambda^{r-1} - e^{-\lambda} \lambda^r)(r-\lambda)^k] \\ &= -k \mu_{k-1} + \sum_{r=0}^{\infty} \frac{1}{r!} e^{-\lambda} \lambda^{r-1} (r-\lambda)^{k+1} \\ &= -k \mu_{k-1} + \frac{1}{\lambda} \mu_{k+1}\end{aligned}$$

i.e.,  $\mu_{k+1} = \lambda \left( \frac{d\mu_k}{d\lambda} + k \mu_{k-1} \right)$  (2)

Using recurrence relation (2), we may compute moments of higher order, provided we know moment of lower order.

Putting  $k = 1$  in (2), we get

$$\begin{aligned}\mu_2 &= \lambda \left( \frac{d\mu_1}{d\lambda} + \mu_0 \right) \\ &= 1 (\because \mu_0 = 1 \text{ and } \mu_1 = 0)\end{aligned}$$

Putting  $k = 2$  in (2), we get

$$\mu_3 = \lambda \left( \frac{d\mu_2}{d\lambda} + 2\mu_1 \right) = \lambda$$

Putting  $k = 3$  in (2), we get

$$\mu_4 = \lambda \left( \frac{d\mu_3}{d\lambda} + 3\mu_2 \right) = \lambda(3\lambda + 1)$$

### Note

The interesting property of the Poisson distribution is the equality of its mean, variance and third-order central moment.

## 3. Geometric Distribution

**Definition:** Let the RV  $X$  denote the number of trials of a random experiment required to obtain the first success (occurrence of an event  $A$ ). Obviously,  $X$  can assume the values  $1, 2, 3, \dots$ .

Now,  $X = r$ , if and only if the first  $(r-1)$  trials result in failure (occurrence of  $\bar{A}$ ) and the  $r$ th trial results in success (occurrence  $A$ ). Hence,

$$P(X = r) = q^{r-1} p; r = 1, 2, 3, \dots, \infty$$

where  $P(A) = p$  and  $P(\bar{A}) = q$ .

If  $X$  is a discrete RV that can assume the values  $1, 2, 3, \dots, \infty$  such that its probability mass function is given by

$$P(X = r) = q^{r-1} p; r = 1, 2, \dots, \infty \text{ where } p + q = 1$$

then  $X$  is said to follow a *geometric distribution*.

**Note**

Geometric distribution is a legitimate probability distribution, since

$$\begin{aligned}\sum_{r=1}^{\infty} P(X=r) &= \sum_{r=1}^{\infty} q^{r-1} p \\ &= p(1 + q + q^2 + \dots \infty) \\ &= \frac{p}{1-q} = 1\end{aligned}$$

**Mean and Variance of the Geometric Distribution**

$$\begin{aligned}E(X) &= \sum_r x_r p_r \\ &= \sum_{r=1}^{\infty} r q^{r-1} p \\ &= p[1 + 2q + 3q^2 + \dots + \infty] \\ &= p(1-q)^{-2} = \frac{1}{p}\end{aligned}$$

$$\begin{aligned}E(X^2) &= \sum_r x_r^2 p r \\ &= \sum_{r=1}^{\infty} r^2 q^{r-1} p \\ &= p \sum_{r=1}^{\infty} \{r(r+1) - r\} q^{r-1} \\ &= p[\{1 \times 2 + 2 \times 3q + 3 \times 4q^2 + \dots + \infty\} - \{1 + 2q + 3q^2 + \dots + \infty\}] \\ &= p[2(1-q)^{-3} - (1-q)^{-2}] \\ &= p\left(\frac{2}{p^3} - \frac{1}{p^2}\right) = \frac{1}{p^2}(2-p) = \frac{1}{p^2}(1+q)\end{aligned}$$

$$\begin{aligned}\text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\ &= \frac{1}{p^2}(1+q) - \frac{1}{p^2} = \frac{q}{p^2}.\end{aligned}$$

**Note**

Sometimes the probability mass function of a geometric RV  $X$  is taken as

$$P(X=r) = q^r p; r = 0, 1, 2, \dots, \infty \text{ where } p + q = 1$$

It is this definition that was given in chapter II. If this definition is assumed, then

$$E(X) = \frac{q}{p} \text{ and } \text{Var}(X) = \frac{q}{p^2} \quad [\text{see example (5) in Worked Example 4(b)}]$$

## 4. Hypergeometric Distribution

If  $X$  represents the number of defectives found, when  $n$  items are drawn without replacement from a lot of  $N$  items containing  $k$  defectives and  $(N - k)$  non-defectives, clearly

$$P(X = r) = \frac{kC_r \cdot (N - k) C_{(n-r)}}{NC_n}; r = 0, 1, 2, \dots, \min(n, k)$$

**Note**

If  $n > k$ , then the maximum value of  $X$  is  $k$ ;

If  $n < k$ , then the maximum value of  $X$  is  $n$ ,

i.e., the maximum value of  $X$  is  $\min(n, k)$ ,

i.e.,  $r$  can take the value  $0, 1, 2, \dots, \min(n, k)$ .

**Definition:** If  $X$  is a discrete RV that can assume, non-negative values  $0, 1, 2, \dots$ , such that its probability mass function is given by

$$P(X = r) = \frac{kC_r \cdot (N - k) C_{(n-r)}}{NC_n}; r = 0, 1, 2, \dots, \min(n, k)$$

then  $X$  is said to follow *hypergeometric distribution* with the parameters  $N, k$  and  $n$ .

**Note**

(1) In the probability mass function of  $X$ ,  $r$  can be assumed to take the values  $0, 1, 2, \dots, n$ , which is true when  $n < k$ . But when  $n > k$ ,  $r$  can take the values  $0, 1, 2, \dots, k$ . In other words,

$$P(X = r) = 0, \text{ when } r = k + 1, k + 2, \dots, n.$$

This value (namely, zero) of the probability is provided by the probability mass function formula itself, since  $kC_r = 0$ , for  $r = k + 1, k + 2, \dots, n$ .

Thus, in the value of  $P(X = r)$ ,  $\min(n, k)$  can be replaced by  $n$ .

(2) Hypergeometric distribution is a legitimate probability distribution, since

$$\begin{aligned} \sum_{r=0}^n P(X = r) &= \sum_{r=0}^n \frac{kC_r \cdot (N - k) C_{(n-r)}}{N C_n} \\ &= \frac{1}{N C_n} N C_n = 1 \text{ since} \end{aligned} \tag{1}$$

$$\begin{aligned} \sum_{r=0}^n kC_r \cdot (N - k) C_{n-r} &= \text{coefficient of } x^n \text{ in } (1 + x)^k (1 + x)^{N-k} \\ &= \text{coefficient of } x^n \text{ in } (1 + x)^N \\ &= NC_n \end{aligned}$$

### Mean and Variance of the Hypergeometric Distribution

$$E(X) = \sum_r x_r p_r$$

$$\begin{aligned}
 &= \sum_{r=0}^n r k C_r \cdot (N-k) C_{(n-r)} / N C_n \\
 &= \frac{k}{N C_n} \sum_{r=1}^n (k-1) C_{(r-1)} (N-k) C_{n-r} \\
 &= \frac{k}{N C_n} \sum_{r'=0}^{n-1} k' C_{r'} (N-1-k') C_{n-1-r'}, \\
 &\quad \text{(on putting } k' = k-1 \text{ and } r' = r-1\text{)} \\
 &= \frac{k}{N C_n} \cdot (N-1) C_{n-1}, \quad [\text{by step (1) in note (2) given above}] \\
 &= \frac{n k}{N}.
 \end{aligned}$$

$$\begin{aligned}
 E(X^2) &= E\{X(X-1) + X\} \\
 &= E\{X(X-1)\} + \frac{n k}{N} \\
 &= \frac{n k}{N} + \sum_{r=0}^n r(r-1) k C_r \cdot (N-k) C_{(n-r)} / N C_n \\
 &= \frac{n k}{N} + \frac{k(k-1)}{N C_n} \sum_{r=2}^n (k-2) C_{(r-2)} \cdot (N-k) C_{(n-r)} \\
 &= \frac{n k}{N} + \frac{k(k-1)}{N C_n} \sum_{r'=0}^{n-2} k' C_{r'} \cdot (N-2-k') C_{(n-2-r')}, \\
 &\quad \text{[on putting } k' = k-2 \text{ and } r' = r-2\text{]} \\
 &= \frac{n k}{N} + \frac{k(k-1)}{N C_n} \cdot (N-2) C_{n-2},
 \end{aligned}$$

$$\begin{aligned}
 &\quad \text{[by step (1) in note (2) given above]} \\
 &= \frac{n k}{N} + \frac{k(k-1) n(n-1)}{N(N-1)}
 \end{aligned}$$

$$\begin{aligned}
 \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
 &= \frac{n k}{N} + \frac{k(k-1) \cdot n(n-1)}{N(N-1)} - \frac{n^2 k^2}{N^2} \\
 &= \frac{n k}{N^2 (N-1)} [N(N-1) + N(k-1)(n-1) - (N-1)nk] \\
 &= \frac{n k}{N^2 (N-1)} [N^2 - Nn - Nk + nk] \\
 &= \frac{n k (N-k) (N-n)}{N^2 (N-1)}
 \end{aligned}$$

**Note** If we denote the proportion of defective items in the lot as  $p$ , i.e.,  $p = \frac{k}{N}$  and  $q = 1 - p$ , then  $E(X) = np$  and  $Var(X) = npq\left(\frac{N-n}{N-1}\right)$ .

### Binomial Distribution as a Limiting Form of the Hypergeometric Distribution

Hypergeometric distribution tends to binomial distribution as  $N \rightarrow \infty$  and  $\frac{k}{N} = p$ .

#### Proof

If  $X$  follows a hypergeometric distribution with parameters  $N, k$  and  $n$ , then

$$\begin{aligned} P\{X=r\} &= \frac{kC_r \cdot (N-k) C_{n-r}}{N C_n}, r=0,1,2,\dots,n \\ &= \frac{k(k-1)\cdots(k-r+1)}{r!} \cdot \frac{(N-k)(N-k-1)\cdots(N-k-n+r+1)}{(n-r)!} \\ &\quad \times \frac{n!}{N(N-1)\cdots(N-n+1)} \\ &= \frac{n!}{r!(n-r)!} \\ &\times \frac{\left[\left(\frac{k}{N}\right)\left(\frac{k}{N}-\frac{1}{N}\right)\cdots\left\{\frac{k}{N}-\frac{(r-1)}{N}\right\}\right]\left[\left(1-\frac{k}{N}\right)\left\{1-\frac{k+1}{N}\right\}\cdots\left\{1-\frac{k+n-r-1}{N}\right\}\right]}{1\cdot\left(1-\frac{1}{N}\right)\left(1-\frac{2}{N}\right)\cdots\left\{1-\frac{(n-1)}{N}\right\}} \end{aligned}$$

(by dividing each factor in the numerator and denominator by  $N$ .)

Putting  $\frac{k}{N} = p$  and proceeding to the limit as  $N \rightarrow \infty$ , we get

$$\begin{aligned} \lim_{N \rightarrow \infty} P\{X=r\} &= nC_r \cdot p^r (1-p)^{n-r} \\ &\quad \frac{k}{N} = p \\ &= nC_r p^r q^{n-r}, r=0,1,2,\dots,n \end{aligned}$$

Thus, the limit of a hypergeometric distribution is a binomial distribution.

**Note** We know that the binomial distribution holds good when we draw samples with replacement (since the probability of getting a defective item has to remain constant), while the hypergeometric distribution holds good when we draw samples without replacement. If the lot size  $N$  is very large, there is not much difference in the proportions of defective items in the lot whether the item drawn is replaced or not. The previous result is simply a mathematical statement of this fact.

## 5. Negative Binomial Distribution

**Definition:** If  $X$  denotes the number of failures preceding the  $n$ th success in a sequence of independent Bernoulli's trials, then  $X$  is said to follow a *negative binomial distribution* with parameter  $n$ .

If  $r$  failures have to occur preceding the  $n$ th success,  $(n + r)$  trials are required, as the first  $(n + r - 1)$  trials should result in  $r$  failures and  $(n - 1)$  successes and  $(n + r)^{\text{th}}$  trial should result in a success, where  $r = 0, 1, 2, \dots$

Hence,  $P\{X = r\} = P[\text{getting } (n - 1) \text{ successes and } r \text{ failures in } (n + r - 1) \text{ trials}] \times P[\text{getting success in the } (n + r)^{\text{th}} \text{ trial}]$

$$= (n + r - 1) C_{(n-1)} p^{n-1} q^r \times p, \text{ where}$$

$p$  and  $q$  have the usual meaning in Bernouilli's trials

$$\text{i.e., } P\{X = r\} = (n + r - 1) C_r p^n q^r (r = 0, 1, 2, \dots)$$

### Note

The negative binomial distribution is a legitimate probability distribution, since

$$\begin{aligned} \sum_{r=0}^{\infty} P(X = r) &= p^n \sum_{r=0}^{\infty} (n + r - 1) C_r q^r \\ &= p^n [1 + n C_1 q + (n + 1) C_2 q^2 + (n + 2) C_3 q^3 + \dots] \\ &= p^n \left[ 1 + \frac{n}{1!} q + \frac{n(n+1)}{2!} q^2 + \frac{n(n+1)(n+2)}{3!} q^3 + \dots \right] \\ &= p^n (1 - q)^{-n} = 1, \text{ since } p + q = 1 \end{aligned}$$

### Note

When  $n = 1$ , the negative binomial distribution reduces to the geometric distribution whose probability law is given by  $P\{X = r\} = q^r p (r = 0, 1, 2, \dots)$

### MGF of the Negative Binomial Distribution

The moment generating function  $M(t)$  of the negative binomial distribution is given by

$$\begin{aligned} M(t) &= E\{e^{tX}\} \\ &= \sum_{r=0}^{\infty} e^{tr} \times P(X = r) \\ &= \sum_{r=0}^{\infty} (n + r - 1) C_r p^n q^r \cdot e^{tr} \\ &= \sum_{r=0}^{\infty} (n + r - 1) C_r p^n (q e^t)^r \end{aligned}$$

$$\begin{aligned}
 &= p^n (1 - q e^t)^{-n} \\
 &= \left( \frac{1 - q e^t}{p} \right)^{-n}
 \end{aligned}$$

Let us now use the MGF to find the mean and variance of the negative binomial distribution

**Note** We have already found out the mean and variance directly in example (3) of section 4(A)]

$$\begin{aligned}
 M(t) &= \left\{ \frac{1 - q \left( 1 + \frac{t}{1!} + \frac{t^2}{2!} + \dots \right)}{p} \right\}^{-n} \\
 &= \left\{ 1 - \frac{qt}{p} \left( 1 + \frac{t}{2} + \dots \right) \right\}^{-n} \\
 &= 1 + \frac{nq}{1! \cdot p} t \left( 1 + \frac{t}{2} + \dots \right) + \frac{n(n+1)}{2!} \frac{q^2 t^2}{p^2} \left( 1 + \frac{t}{2} + \dots \right)^2 + \dots
 \end{aligned}$$

$E(X)$  = Coefficient of  $\frac{t}{1!}$  in the expansion of  $M(t)$ .

$$= \frac{nq}{p}$$

$E(X^2)$  = Coefficient of  $\frac{t^2}{2!}$  in the expansion of  $M(t)$

$$= \frac{nq}{p} + \frac{n(n+1)q^2}{p^2}$$

$$\text{Var}(X) = E(X^2) - E^2(X)$$

$$= \frac{nq}{p} + \frac{n(n+1)q^2}{p^2} - \frac{n^2 q^2}{p^2}$$

$$= \frac{nq}{p} \left( 1 + \frac{q}{p} \right) = \frac{nq}{p^2}$$

### Worked Example 5(A)

#### Example 1

Out of 800 families with 4 children each, how many families would be expected to have (a) 2 boys and 2 girls, (b) at least 1 boy, (c) at most 2 girls, and (d) children of both sexes. Assume equal probabilities for boys and girls.

Considering each child as a trial,  $n = 4$ . Assuming that birth of a boy is a success,  $p = \frac{1}{2}$  and  $q = \frac{1}{2}$ . Let  $X$  denote the number of successes (boys).

$$(a) \quad P(2 \text{ boys and } 2 \text{ girls}) = P(X = 2)$$

$$\begin{aligned} &= 4C_2 \cdot \left(\frac{1}{2}\right)^2 \cdot \left(\frac{1}{2}\right)^{4-2} \\ &= 6 \times \left(\frac{1}{2}\right)^4 = \frac{3}{8} \end{aligned}$$

$\therefore$  No. of families having 2 boys and 2 girls

$$= N \cdot P(X = 2) \text{ (where } N \text{ is the total no. of families considered)}$$

$$\begin{aligned} &= 800 \times \frac{3}{8} \\ &= 300. \end{aligned}$$

$$(b) \quad P(\text{at least } 1 \text{ boy}) = P(X \geq 1)$$

$$\begin{aligned} &= P(X = 1) + P(X = 2) + P(X = 3) + P(X = 4) \\ &= 1 - P(X = 0) \\ &= 1 - 4C_0 \cdot \left(\frac{1}{2}\right)^0 \cdot \left(\frac{1}{2}\right)^4 \\ &= 1 - \frac{1}{16} = \frac{15}{16} \end{aligned}$$

$\therefore$  No. of families having at least 1 boy

$$= 800 \times \frac{15}{16} = 750.$$

$$(c) \quad P(\text{at most } 2 \text{ girls}) = P(\text{exactly } 0 \text{ girl, } 1 \text{ girl or } 2 \text{ girls})$$

$$\begin{aligned} &= P(X = 4, X = 3 \text{ or } X = 2) \\ &= 1 - \{P(X = 0) + P(X = 1)\} \end{aligned}$$

$$= 1 - \left\{ 4C_0 \cdot \left(\frac{1}{2}\right)^4 + 4C_1 \cdot \left(\frac{1}{2}\right)^4 \right\}$$

$$= \frac{11}{16}$$

$\therefore$  No. of families having at most 2 girls

$$= 800 \times \frac{11}{16} = 550$$

(d)  $P(\text{children of both sexes})$

$$\begin{aligned} &= 1 - P(\text{children of the same sex}) \\ &= 1 - \{P(\text{all are boys}) + P(\text{all are girls})\} \\ &= 1 - \{P(X = 4) + P(X = 0)\} \\ &= 1 - \left\{ 4C_4 \cdot \left(\frac{1}{2}\right)^4 + 4C_0 \cdot \left(\frac{1}{2}\right)^4 \right\} \\ &= 1 - \frac{1}{8} = \frac{7}{8} \end{aligned}$$

$\therefore$  No. of families having children of both sexes

$$= 800 \times \frac{7}{8} = 700.$$

---

### Example 2

---

An irregular 6-faced die is such that the probability that it gives 3 even numbers in 5 throws is twice the probability that it gives 2 even numbers in 5 throws. How many sets of exactly 5 trials can be expected to give no even number out of 2500 sets?

Let the probability of getting an event number with the unfair die be  $p$ .

Let  $X$  denote the number of even numbers obtained in 5 trials (throws).

**Given:**  $P(X = 3) = 2 \times P(X = 2)$

$$\text{i.e., } 5C_3 p^3 q^2 = 2 \times 5C_2 p^2 q^3$$

$$\text{i.e., } p = 2q = 2(1-p)$$

$$\therefore 3p = 2 \text{ or } p = \frac{2}{3} \text{ and } q = \frac{1}{3}$$

Now,  $P(\text{getting no even number})$

$$= P(X = 0)$$

$$= 5C_0 \cdot p^0 \cdot q^5 = \left(\frac{1}{3}\right)^5 = \frac{1}{243}$$

$\therefore$  Number of sets having no success (even number) out of  $N$  sets  $= N \times P(X = 0)$

$$\begin{aligned} \therefore \text{Required number of sets} &= 2500 \times \frac{1}{243} \\ &= 10, \text{ nearly} \end{aligned}$$

---

### Example 3

---

The probability of a successful rocket launching is  $p$ . If launching attempts are made until 3 successful launching have occurred, what is the probability that exactly 5 attempts will be necessary? What is the probability that fewer than 5 attempts will be necessary?

If launching attempts are made until 3 consecutive successful launching occur, what are the probabilities? [See Example (1) in Section 1 (c)]

- (a) Exactly 5 attempts will be required to get 3 successes (successful launching of rockets), if 2 successes occur in the first 4 attempts and third success occurs in the fifth attempt.

$$\therefore P(\text{exactly 5 attempts are required})$$

$$= P(2 \text{ successes in 4 attempts})$$

$$\times P(\text{success in the single 5th attempt})$$

$$= 4C_2 p^2 q^2 \times p$$

[ $\because$  The no. of successes in the 4 independent attempts follow  $B(4, p)$ ]

$$= 6 p^3 q^3$$

- (b)  $P\{\text{fewer than 5 attempts are required}\}$

$$= P\{\text{exactly 3 or 4 attempts are required}\}$$

$$= [P\{2 \text{ successes in the first 2 attempts}\}$$

$$\times P(\text{success in the 3rd attempt})]$$

$$+ [P\{2 \text{ successes in the first 3 attempts}\}$$

$$\times P(\text{success in the 4th attempt})]$$

$$= 2C_2 p^2 q^0 \times p + 3C_2 p^2 q^1 \times p$$

$$= p^3 + 3p^3 q = p^3 (1 + 3q)$$

- (c) Five attempts will be required to get 3 consecutive successes, if the first 2 attempts result in failures and the last 3 attempts result in successes.

$$\therefore \text{Required probability} = q \cdot q \cdot p \cdot p \cdot p = p^3 q^2$$

- (d) Three attempts will be required to get 3 consecutive successes, if each attempt results in a success.

$$\therefore \text{Probability for this} = p^3.$$

Four attempts will be required to get 3 consecutive successes, if the first attempt results in a failure and the remaining attempts result in a success each.

$$\therefore \text{Probability for this} = qp^3$$

$$\therefore P\{\text{fewer than 5 attempts are required}\}$$

$$= p^3 + qp^3 = p^3 (1 + q).$$

---

### Example 4

---

A communication system consists of  $n$  components, each of which will independently function with probability  $p$ . The total system will be able to operate effectively if at least one-half of its components function. For what values of  $p$

is a 5-component system more likely to operate effectively than a 3-component system?

Since the probability  $p$  of functioning of every component is a constant and the  $n$  components function independently, the number of components  $X$  that function follow a binomial distribution with parameters  $n$  and  $p$ .

$$\therefore P(X = r) = nC_r p^r q^{n-r}; r = 0, 1, 2, \dots, n.$$

$$P(5\text{-component system functions effectively})$$

$$= P(X = 3 \text{ or } 4 \text{ or } 5)$$

$$= \sum_{r=3}^5 5C_r p^r q^{5-r} \quad (\because n = 5)$$

$$P(3\text{-component system functions effectively})$$

$$= P(X = 2 \text{ or } 3)$$

$$= \sum_{r=2}^3 3C_r p^r q^{3-r} \quad (\because n = 3)$$

5-component system will function more effectively than the 3-component system,

$$\text{if } \sum_{r=3}^5 5C_r p^r q^{5-r} \geq \sum_{r=2}^3 3C_r p^r q^{3-r}$$

$$\text{i.e., } 10p^3q^2 + 5p^4q + p^5 \geq 3p^2q + p^3$$

$$\text{i.e., } 10p^3(1 - 2p + p^2) + 5p^4(1 - p) + p^5 \geq 3p^2(1 - p) + p^3$$

$$\text{i.e., } 3p^2(2p^3 - 5p^2 + 4p - 1) \geq 0$$

$$\text{i.e., } 3p^2(p - 1)^2(2p - 1) \geq 0$$

$$\text{i.e., } (2p - 1) \geq 0, \quad [\text{since } 3p^2(p - 1)^2 \geq 0]$$

$$\text{i.e., } p \geq \frac{1}{2}.$$

### Example 5

If the probability that a child is a boy is  $p$ , where  $0 < p < 1$ , find the expected number of boys in a family with  $n$  children, given that there is at least one boy.

Let  $X$  be the number of boys (successes) out of  $n$  children (trials)

Then  $X$  follows a  $B(n, p)$ .

Required to find  $E\{X/X \geq 1\}$ .

$$E\{X/X \geq 1\} = \sum_r r \cdot P(X = r / X \geq 1)$$

$$= \sum_{r=1}^n r \cdot \frac{P(X = r)}{P(X \geq 1)}$$

$$\begin{aligned}&= \sum_{r=1}^n \frac{r \cdot nC_r p^r q^{n-r}}{1 - P(X = 0)} \\&= \frac{1}{1 - q^n} \sum_{r=0}^n r \cdot nC_r p^r q^{n-r} \\&= \frac{n p}{1 - q^n}\end{aligned}$$

---

### Example 6

---

Two dice are thrown 120 times. Find the average number of times in which the number of the first dice exceeds the number on the second dice.

The number on the first dice exceeds that on the second dice, in the following combinations:

(2, 1); (3, 1), (3, 2); (4, 1), (4, 2), (4, 3); (5, 1), (5, 1), (5, 2), (5, 3), (5, 4); (6, 1), (6, 2), (6, 3), (6, 4), (6, 5),

where the numbers in the parentheses represent the numbers in the first and second dice respectively.

∴  $P(\text{success}) = P(\text{number in the first dice exceeds the number in the second dice})$

$$= \frac{15}{36} = \frac{5}{12}$$

This probability remains the same in all the throws that are independent.

If  $X$  is the number of successes, then  $X$  follows a binomial distribution with parameters  $n (= 120)$  and  $p \left(= \frac{5}{12}\right)$ .

$$\therefore E(X) = np = 120 \times \frac{5}{12} = 50$$

---

### Example 7

---

Fit a binomial distribution for the following data:

x:	0	1	2	3	4	5	6	Total
f.:	5	18	28	12	7	6	4	80

Fitting a binomial distribution means assuming that the given distribution is approximately binomial and hence finding the probability mass function and then finding the theoretical frequencies.

To find the binomial frequency distribution  $N(q + p)^n$ , which fits the given data, we required  $N$ ,  $n$  and  $p$ . We assume  $N = \text{total frequency} = 80$  and  $n = \text{no. of trials} = 6$  from the given data.

To find  $p$ , we compute the mean of the given frequency distribution and equate it to  $np$  (means of the binomial distribution).

$x:$	0	1	2	3	4	5	6	Total
$f:$	5	18	28	12	7	6	4	80
$fx:$	0	18	56	36	28	30	24	192

$$\bar{x} = \frac{\sum fx}{\sum f} = \frac{192}{80} = 2.4$$

i.e.,  $np = 2.4$  or  $6p = 2.4$

$\therefore p = 0.4$  and  $q = 0.6$

If the given distribution is nearly binomial, the theoretical frequencies are given by the successive terms in the expansion of  $80(0.6 + 0.4)^6$ . Thus we get,

$x:$	0	1	2	3	4	5	6
Theoretical $f:$	3.73	14.93	24.88	22.12	11.06	2.95	0.33

Converting these values into whole numbers consistent with the condition that the total frequency is 80, the corresponding binomial frequency distribution is as follows:

$x:$	0	1	2	3	4	5	6	Total
$f:$	4	15	25	22	11	3	0	80

### Example 8

The number of monthly breakdowns of a computer is a RV having a Poisson distribution with mean equal to 1.8. Find the probability that this computer will function for a month

- (a) without a breakdown,
- (b) with only one breakdown, and
- (c) with at least one breakdown.

Let  $X$  denote the number of breakdowns of the computer in a month.  $X$  follows a Poisson distribution with mean (parameter)  $\lambda = 1.8$ .

$$\therefore P\{X = r\} = \frac{e^{-\lambda} \cdot \lambda^r}{r!} = \frac{e^{-1.8} \cdot (1.8)^r}{r!}$$

- (a)  $P(X = 0) = e^{-1.8} = 0.1653$
- (b)  $P(X = 1) = e^{-1.8} (1.8) = 0.2975$
- (c)  $P(X \geq 1) = 1 - P(X = 0) = 0.8347$

### Example 9

It is known that the probability of an item produced by a certain machine will be defective is 0.05. If the produced items are sent to the market in packets of 20, find the number of packets containing at least, exactly and at most 2 defective items in a consignment of 1000 packets using (a) binomial distribution, and (b) Poisson approximation to binomial distribution.

### **Use of Binomial Distributions**

$p$  = probability that an item is defective = 0.05,  $q = 0.95$  and  $n$  = No. of independent items (trials) considered = 20.

Let  $X$  denote the number of defectives in the  $n$  items considered.

$$\begin{aligned} P(X = r) &= nC_r p^r q^{n-r} \\ \text{(a)} \quad \therefore P(X = 2) &= 20 C_2 (0.05)^2 (0.95)^{18} \\ &= 0.1887 \end{aligned}$$

If  $N$  is the number of sets (packets), each set (packet) containing 20 trials (items), then the number of sets containing exactly 2 successes (defectives) is given by

$$\begin{aligned} N(X = 2) &= N \times P(X = 2) \\ &= 1000 \times 0.1887 = 189, \text{ nearly} \end{aligned}$$

$$\begin{aligned} \text{(b)} \quad P(\text{at least 2 defectives}) &= P(X \geq 2) \\ &= 1 - \{P(X = 0) + P(X = 1)\} \\ &= 1 - [20C_0 (0.05)^0 (0.95)^{20} + 20 C_1 (0.05)^1 (0.95)^{19}] \\ &= 1 - [0.3585 + 0.3774] \\ &= 0.2641 \end{aligned}$$

$$\begin{aligned} \therefore N(X \geq 2) &= N \times P(X \geq 2) \\ &= 1000 \times 0.2641 = 264, \text{ nearly} \end{aligned}$$

$$\begin{aligned} \text{(c)} \quad P(\text{at most 2 defectives}) &= P(X \leq 2) \\ &= P(X = 0) + P(X = 1) + P(X = 2) \\ &= \sum_{r=0}^2 20 C_r (0.05)^r (0.95)^{20-r} \\ &= 0.3585 + 0.3774 + 0.1887 \\ &= 0.9246 \end{aligned}$$

$$\begin{aligned} \therefore N(X \leq 2) &= N \times P(X \leq 2) \\ &= 1000 \times 0.9246 = 925, \text{ nearly} \end{aligned}$$

### **Use of Poisson Distributions**

Since  $p = 0.05$  is very small and  $n = 20$  is sufficiently large, binomial distribution may be approximated by Poisson distribution with parameter  $\lambda = np = 1$ .

$$\therefore P(X = r) = \frac{e^{-\lambda} \cdot \lambda^r}{r!} = \frac{e^{-1}}{r!}$$

$$\text{(a)} \quad P(X = 2) = \frac{e^{-1}}{2!} = 0.1839$$

$$\therefore N(X = 2) = 1000 \times 0.1839 = 184, \text{ nearly}$$

$$(b) P(X \geq 2) = 1 - \{P(X = 0) + P(X = 1)\} \\ = 1 - \{e^{-1} + e^{-1}\} = 0.2642$$

$\therefore N(X \geq 2) = 1000 \times 0.2642 = 264$ , nearly.

$$(c) P(X \leq 2) = \sum_{r=0}^2 P(X = r) = \sum_{r=0}^2 \frac{e^{-1}}{r!} \\ = 0.9197$$

$\therefore N(X \leq 2) = 920$ , nearly.

### Example 10

Prove the reproductive property of independent Poisson RVs. Hence, find the probability of 5 or more telephone calls arriving in a 9-min period in a college switch-board, if the telephone calls that are received at the rate of 2 every 3 minutes follow a Poisson distribution.

Let  $X_1$  and  $X_2$  be independent RVs that follow Poisson distributions with parameters  $\lambda_1$  and  $\lambda_2$  respectively.

Let  $X = X_1 + X_2$

$$P\{X = n\} = P\{X_1 + X_2 = n\} \\ = \sum_{r=0}^n P\{X_1 = r\} \cdot P\{X_2 = n - r\}, \\ \text{(since } X_1 \text{ and } X_2 \text{ are independent)}$$

$$\begin{aligned} &= \sum_{r=0}^n \frac{e^{-\lambda_1} \cdot \lambda_1^r}{r!} \frac{e^{-\lambda_2} \cdot \lambda_2^{n-r}}{(n-r)!} \\ &= \frac{e^{-(\lambda_1 + \lambda_2)}}{n!} \sum_{r=0}^n \frac{n!}{r!(n-r)!} \lambda_1^r \cdot \lambda_2^{n-r} \\ &= \frac{e^{-(\lambda_1 + \lambda_2)}}{n!} \sum_{r=0}^n nC_r \cdot \lambda_1^r \cdot \lambda_2^{n-r} \\ &= \frac{e^{-(\lambda_1 + \lambda_2)} \cdot (\lambda_1 + \lambda_2)^n}{n!} \end{aligned}$$

Thus, the sum of 2 independent Poisson variables with parameters  $\lambda_1$  and  $\lambda_2$  is also a Poisson variable with parameter  $(\lambda_1 + \lambda_2)$ .

This property, which can be extended to any finite number of independent Poisson variables is known as *the Reproductive Property* of Poisson RVs. [For an alternative proof, see example (9) in section 4(b)]

Let  $X_1, X_2, X_3$  denote the number of telephone calls received in three consecutive 3-min periods.

Each of  $X_1, X_2, X_3$  follows a Poisson distribution with parameter (mean)  $\lambda = 2$ .

∴  $X = X_1 + X_2 + X_3$  follow a Poisson distribution with parameter 6.

Clearly,  $X$  represents the number of calls received in a 9-min period.

Now,  $P(X \geq 5) = 1 - P\{X \leq 4\}$

$$\begin{aligned} &= 1 - \sum_{r=0}^4 \frac{e^{-6} \cdot 6^r}{r!} \\ &= 1 - (0.0025 + 0.0149 + 0.0446 + 0.0892 + 0.1339) \\ &= 1 - 0.2851 = 0.7149 \end{aligned}$$

---

### Example 11

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If the number  $X$  of particles emitted during a 1 h period from a radioactive source has a Poisson distribution with parameter  $\lambda = 4$  and that the probability that any emitted particle is recorded is  $p = 0.9$ , find the probability distribution of the number  $Y$  of the particles recorded in a 1 h period and hence the probability that no particle is recorded.

$$\begin{aligned} P\{Y = n\} &= \sum_{r=0}^{\infty} P[X = n+r \text{ and } n \text{ of them are recorded}] \\ &= \sum_{r=0}^{\infty} \frac{e^{-\lambda} \cdot \lambda^{n+r}}{(n+r)!} (n+r) C_n \cdot p^n q^r \\ &= \sum_{r=0}^{\infty} \frac{e^{-\lambda} \cdot (\lambda p)^n}{(n+r)!} \frac{(n+r)!}{n! r!} (\lambda q)^r \\ &= \frac{e^{-\lambda} \cdot (\lambda p)^n}{n!} \sum_{r=0}^{\infty} \frac{1}{r!} (\lambda q)^r \\ &= \frac{e^{-\lambda} \cdot (\lambda p)^n}{n!} e^{\lambda q} \\ &= \frac{e^{-\lambda(1-q)} \cdot (\lambda p)^n}{n!} = \frac{e^{-\lambda p} \cdot (\lambda p)^n}{n!}, \quad n = 0, 1, 2, \dots, \infty. \end{aligned}$$

Therefore,  $Y$ , the number of recorded particles, follows a Poisson distribution with parameter  $\lambda p$ .

$$\begin{aligned} \text{Hence, } P(Y = 0) &= e^{-\lambda p} \\ &= e^{-4 \times 0.9} = e^{-3.6} \\ &= 0.0273 \end{aligned}$$

---

### Example 12

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If  $X$  and  $Y$  are independent Poisson RVs, show that the conditional distribution of  $X$ , given the value of  $X + Y$ , is a binomial distribution.

Let  $X$  and  $Y$  follow Poisson distributions with parameters  $\lambda_1$  and  $\lambda_2$  respectively.

$$\begin{aligned}
 & \text{Now, } P\{X = r/(X + Y) = n\} \\
 &= \frac{P\{X = r \text{ and } (X + Y) = n\}}{P\{(X + Y) = n\}} = \frac{P\{X = r; Y = n - r\}}{P[X + Y = n]} \\
 &= \frac{P\{X = r\} \cdot P\{X = n - r\}}{P\{(X + Y) = n\}} \quad (\text{by independence of } X \text{ and } Y) \\
 &= \frac{\{e^{-\lambda_1} \cdot \lambda_1^r / r!\} \{e^{-\lambda_2} \cdot \lambda_2^{n-r} / (n-r)!\}}{e^{-(\lambda_1 + \lambda_2)} \cdot (\lambda_1 + \lambda_2)^n / n!} \\
 &\qquad\qquad\qquad (\text{by the reproductive property}) \\
 &= \frac{n!}{r!(n-r)!} \left( \frac{\lambda_1}{\lambda_1 + \lambda_2} \right)^r \cdot \left( \frac{\lambda_2}{\lambda_1 + \lambda_2} \right)^{n-r} \\
 &= nC_r p^r q^{n-r}, \text{ where } p = \frac{\lambda_1}{\lambda_1 + \lambda_2} \text{ and } q = \frac{\lambda_2}{\lambda_1 + \lambda_2}
 \end{aligned}$$

Hence, the required result.

### Example 13

Fit a Poisson distribution for the following distribution:

x:	0	1	2	3	4	5	Total
f:	142	156	69	27	5	1	400

Fitting a Poisson distribution for a given distribution means assuming that the given distribution is approximately Poisson and, hence, finding the probability mass function and then finding the theoretical frequencies.

To find the probability mass function

$$P\{X = r\} = \frac{e^{-\lambda} \cdot \lambda^r}{r!} \quad r = 0, 1, 2, \dots, \infty$$

of the approximate Poisson distribution, we require  $\lambda$ , which is the mean of the Poisson distribution.

We find the mean of the given distribution and assume it as  $\lambda$ .

x :	0	1	2	3	4	5	Total
f:	142	156	69	27	5	1	400
fx :	0	156	138	81	20	5	400

$$\bar{x} = \frac{\sum f x}{\sum f} = \frac{400}{400} = 1 = \lambda$$

The theoretical frequencies are given by

$$\begin{aligned}
 & \frac{N e^{-\lambda} \cdot \lambda^r}{r!} \quad \text{where } N = 400, \text{ obtained from the given distribution.} \\
 &= \frac{400 e^{-1}}{r!}, r = 0, 1, 2, \dots, \infty
 \end{aligned}$$

Thus, we get

$x:$	0	1	2	3	4	5
Theoretical $f:$	147.15	147.15	73.58	24.53	6.13	1.23

The theoretical frequencies for  $x = 6, 7, 8, \dots$  are very small and hence neglected.

Converting the theoretical frequencies into whole numbers consistent with the condition that the total frequency = 400, we get the following Poisson frequency distribution which fits the given distribution:

$x:$	0	1	2	3	4	5
Theoretical $f:$	147	147	74	25	6	1

### Example 14

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If the probability that an applicant for a driver's license will pass the road test on any given trial is 0.8, what is the probability that he will finally pass the test (a) on the fourth trial, and (b) in fewer than 4 trials?

Let  $X$  denote the number of trials required to achieve the first success. Then  $X$  follows a geometric distribution given by

$$P(X = r) = q^{r-1} p; \quad r = 1, 2, 3, \dots, \infty$$

Here,  $p = 0.8$  and  $q = 0.2$

$$\begin{aligned} \text{(a)} \quad P(X = 4) &= 0.8 \times (0.2)^{4-1} \\ &= 0.8 \times 0.008 = 0.0064 \end{aligned}$$

$$\begin{aligned} \text{(b)} \quad P(X < 4) &= \sum_{r=1}^3 0.8 \times (0.2)^{r-1} \\ &= 0.8[(0.2)^0 + (0.2)^1 + (0.2)^2 + (0.2)^3] \\ &= 0.9984. \end{aligned}$$

### Example 15

---

$A$  and  $B$  shoot independently until each has hit his own target. The probabilities of their hitting the target at each shot are  $\frac{3}{5}$  and  $\frac{5}{7}$  respectively. Find the probability that  $B$  will require more shots than  $A$ .

Let  $X$  denote the number of trials required by  $A$  to get his first success. Then  $X$  follows a geometric distribution given by

$$P(X = r) = p_1 q_1^{r-1} = \frac{3}{5} \left( \frac{2}{5} \right)^{r-1}; \quad r = 1, 2, \dots, \infty$$

Let  $Y$  denote the number of trials required by  $B$  to get his first success. Then  $Y$  follows a geometric distribution given by

$$P(Y = r) = p_2 \cdot q_2^{r-1} = \frac{5}{7} \cdot \left(\frac{2}{7}\right)^{r-1}; \quad r = 1, 2, \dots, \infty$$

$P\{B$  requires more trials to get his first success than  $A$  requires to get his first success $\}$

$$\begin{aligned} &= \sum_{r=1}^{\infty} P\{X = r \text{ and } Y = r + 1 \text{ or } r + 2, \dots, \infty\} \\ &= \sum_{r=1}^{\infty} [P\{X = r\} \cdot P\{Y = r + 1 \text{ or } r + 2, \dots, \infty\}] \text{ (by independence)} \\ &= \sum_{r=1}^{\infty} \frac{3}{5} \cdot \left(\frac{2}{5}\right)^{r-1} \cdot \sum_{k=1}^{\infty} \frac{5}{7} \cdot \left(\frac{2}{7}\right)^{r+k-1} \\ &= \sum_{r=1}^{\infty} \frac{3}{7} \cdot \left(\frac{4}{35}\right)^{r-1} \sum_{k=1}^{\infty} \left(\frac{2}{7}\right)^k \\ &= \sum_{r=1}^{\infty} \frac{3}{7} \cdot \left(\frac{4}{35}\right)^{r-1} \left( \frac{\frac{2}{7}}{1 - \frac{2}{7}} \right) \\ &= \frac{6}{35} \sum_{r=1}^{\infty} \left(\frac{4}{35}\right)^{r-1} = \frac{6}{35} \cdot \frac{1}{1 - \frac{4}{35}} = \frac{6}{31}. \end{aligned}$$

### Example 16

A coin is tossed until the first head occurs. Assuming that the tosses are independent and the probability of a head occurring is  $p$ , find the value of  $p$  so that the probability that an odd number of tosses is required is equal to 0.6.

Can you find a value of  $p$  so that the probability is 0.5 that an odd number of tosses is required?

Let  $X$  denote the number of tosses required to get the first head (success). Then  $X$  follows a geometric distribution given by

$$\begin{aligned} P(X = r) &= pq^{r-1}; \quad r = 1, 2, \dots, \\ \therefore P(X = \text{an odd number}) &= P(X = 1 \text{ or } 3 \text{ or } 5, \dots) \end{aligned}$$

$$\begin{aligned} &= \sum_{r=1}^{\infty} P(X = 2r - 1) = \sum_{r=1}^{\infty} p q^{2r-2} \\ &= \frac{p}{q^2} (q^2 + q^4 + q^6 + \dots) \end{aligned}$$

$$= \frac{p}{q^2} \cdot \frac{q^2}{1-q^2} = \frac{1}{1+q} \quad (\text{since } p+q=1) \quad (1)$$

Now,  $\frac{1}{1+q} = 0.6$ , if  $\frac{1}{2-p} = 0.6$

i.e.,  $0.6p = 0.2$

i.e.,  $p = \frac{1}{3}$

Now,  $\frac{1}{1+q} = 0.5$ , if  $\frac{1}{2-p} = 0.5$

i.e.,  $1 - \frac{p}{2} = 1$

i.e.,  $p = 0$

though we get  $p = 0$ , it is meaningless, because

$$P(X = \text{an odd number})$$

$$\begin{aligned} &= \sum_{r=1}^{\infty} pq^{2r-2} \\ &= 0, \text{ when } p = 0 \end{aligned}$$

Hence, the value of  $p$  cannot be found out.

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### Example 17

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Establish the *memoryless property* of geometric distribution, that is, if  $X$  is a discrete RV following a geometric distribution, then  $P\{X > m + n | X > m\} = P\{X > n\}$ , where  $m$  and  $n$  are any two positive integers. Prove the converse also, if it is true.

Since  $X$  follows a geometric distribution,

$$P(X = r) = pq^{r-1}; r = 1, 2, \dots, \infty; p + q = 1$$

$$\begin{aligned} \therefore P\{X > k\} &= \sum_{r=k+1}^{\infty} pq^{r-1} = p(q^k + q^{k+1} + q^{k+2} + \dots + \infty) = \frac{pq^k}{1-q} \\ &= q \end{aligned}$$

$$\text{Now, } P\{X > m + n | X > m\} = \frac{P\{X > m + n \text{ and } X > m\}}{P\{X > m\}}$$

$$= \frac{P\{X > m + n\}}{P\{X > m\}}$$

$$= \frac{q^{m+n}}{q^m} = q^n = P(X > n)$$

The converse of the above result is also true, i.e., if  $P\{X > m + n/X > m\} = P\{X > n\}$ , where  $m$  and  $n$  are any two positive integers, then  $X$  follows a geometric distribution.

Since  $X$  takes the values  $1, 2, 3, \dots$ ,  $P(X \geq 1) = 1$ .

Let  $P\{X > 1\} = q$

$$\text{Now, } P\{X = (r+1)\} = P\{X > r\} - P\{X > (r+1)\} \quad (1)$$

$$\begin{aligned} \therefore \frac{P\{X = (r+1)\}}{P\{X > r\}} &= 1 - \frac{P\{X > (r+1)\}}{P\{X > r\}} \\ &= 1 - P\{X > (r+1)/X > r\} \\ &= 1 - P(X > 1) \text{ (by the data)} \\ &= 1 - q \\ \therefore P\{X = r+1\} &= (1-q) P(X > r) \quad (2) \\ &= (1-q) [P\{X > (r-1)\} - P\{X = r\}] \\ &\quad [\text{from (1), on changing } r \text{ to } (r-1)] \\ &= (1-q) [P\{X > (r-1)\} - (1-q) P(X > (r-1))] \\ &\quad [\text{from (2) on changing } r \text{ to } (r-1)] \\ &= (1-q) q P\{X > (r-1)\} \\ &= (1-q) q^2 P\{X > (r-2)\} \\ &= (1-q) q^{r-1} P(X > 1) \\ &= (1-q) q^r \\ \therefore P\{X = r\} &= pq^{r-1}, \text{ where } p = 1 - q \text{ and } r = 1, 2, \dots \end{aligned}$$

That is,  $X$  follows a geometric distribution.

### Example 18

If two independent RVs  $X$  and  $Y$  have the same geometric distribution, prove that the conditional distribution of  $X$ , given that  $X + Y = k$  is a discrete uniform distribution.

**Given:**  $P(X = r) = P(Y = r) = pq^{r-1}; r = 1, 2, \dots$

$$\begin{aligned} \text{Now, } P\{X = r/X + Y = k\} &= \frac{P\{X = r \text{ and } X + Y = k\}}{P\{X + Y = k\}} \\ &= \frac{P\{X = r \text{ and } Y = k - r\}}{P\{X + Y = k\}} \\ &= \frac{P\{X = r\} \cdot P\{Y = k - r\}}{\sum_{r=1}^{k-1} P\{X = r\} \cdot P\{Y = k - r\}} \quad (\text{by independence}) \\ &= \frac{pq^{r-1} \cdot pq^{k-r-1}}{\sum_{r=1}^{k-1} pq^{r-1} \cdot pq^{k-r-1}} \end{aligned}$$

$$= \frac{q^{k-2}}{\sum_{r=1}^{k-1} q^{k-2}} = \frac{1}{k-1}; \quad r = 1, 2, 3, \dots, (k-1)$$

**Note** When  $P\{X = r\} = \text{constant}$ , the discrete RV  $X$  is said to follow a discrete uniform distribution.

Thus, the conditional distribution of  $X$ , given that  $X + Y = k$ , is a discrete uniform distribution.

### Example 19

A taxicab company has 12 Ambassadors and 8 Fiats. If 5 of these taxi cabs are in the workshop for repairs and an Ambassador is as likely to be in for repairs as a Fiat, what is the probability that

- (a) 3 of them are Ambassadors and 2 are Fiats,
- (b) at least 3 of them are Ambassadors, and
- (c) all the 5 are of the same make?

Let  $X$  denote the number of Ambassadors in the workshop out of the 5 taxicabs.

We note that  $N = 20$ ,  $k = 12$ , and  $n = 5$  and  $X$  follows a hypergeometric distribution given by

$$P(X = r) = \frac{12 C_r \cdot 8 C_{5-r}}{20 C_5} \quad r = 0, 1, \dots, 5$$

$$\begin{aligned} \text{(a)} \quad P(\text{3 Ambassadors and 2 Fiats}) &= \frac{12 C_3 \cdot 8 C_2}{20 C_5} \\ &= \frac{220 \times 28}{15,504} = \frac{385}{969} \end{aligned}$$

$$\begin{aligned} \text{(b)} \quad P(\text{at least 3 Ambassadors}) &= P(X \geq 3) \\ &= P(X = 3) + P(X = 4) + P(X = 5) \\ &= \sum_{r=3}^5 \frac{12 C_r \cdot 8 C_{5-r}}{20 C_5} \\ &= \frac{1}{15,504} \{220 \times 28 + 495 \times 8 + 792 \times 1\} = \frac{682}{969} \end{aligned}$$

$$\begin{aligned} \text{(c)} \quad P(\text{all the 5 are of the same make}) &= P(\text{all are Ambassadors or all are Fiats}) \\ &= P(X = 5 \text{ or } X = 0) \\ &= P(X = 5) + P(X = 0) \\ &= \frac{12 C_5 \cdot 8 C_0}{20 C_5} + \frac{12 C_0 \cdot 8 C_5}{20 C_5} \\ &= \frac{1}{15,504} (792 + 56) = \frac{53}{969} \end{aligned}$$

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**Example 20**

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A panel of 7 judges is to decide which of the 2 final contestants  $A$  and  $B$  will be declared the winner; a simple majority of the judges will determine the winner. Assume that 4 of the judges will vote for  $A$  and the other 3 will vote for  $B$ . If we randomly select 3 of the judges and seek their verdict, what is the probability that a majority of them will favour  $A$ ?

Let  $X$  denote the number of judges favouring  $A$ . We note that  $N = 7$ ,  $k = 4$  and  $n = 3$  and  $X$  follows a hypergeometric distribution, given by

$$P(X = r) = \frac{4C_r \cdot 3C_{3-r}}{7C_3} \quad r = 0, 1, 2, 3$$

$$P(\text{a majority of 3 chosen judges will favour } A) = P(X \geq 2)$$

$$= P(X = 2) + P(X = 3)$$

$$= \frac{1}{7C_3} \{4C_2 \cdot 3C_1 + 4C_3 \cdot 3C_0\}$$

$$= \frac{1}{35} (6 \times 3 + 4 \times 1) = \frac{22}{35}$$

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**Example 21**

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If a boy is throwing stones at a target, what is the probability that his 10th throw is his 5th hit, if the probability of hitting the target at any trial is  $\frac{1}{2}$ ?

Since the 10th throw should result in the 5th success, viz., hit, the first 9 throws ought to have resulted in 4 successes and 5 failures. Hence, in the usual notation,  $n = 5$ ,  $r = 5$ ,  $p = \frac{1}{2} = q$ .

$$\therefore \text{Required probability} = P(X = 5) = (5 + 5 - 1) C_5 \cdot \left(\frac{1}{2}\right)^5 \cdot \left(\frac{1}{2}\right)^5 \\ = 9C_4 \times \frac{1}{2^{10}} = 0.123$$

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**Example 22**

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An item is produced in large numbers. The machine is known to produce 2% defectives. A quality control inspector is examining the items by taking them one by one at random. What is the probability that at least 4 items are to be examined in order to get 2 defectives?

Success  $\equiv$  defective and failure  $\equiv$  non-defective

$$p = 0.02 \text{ and } q = 0.98$$

If at least 4 items are to be examined (viz. 4 trials are required) to give 2 defectives (success), 4 or 5 or 6 or ... items are to be examined i.e., the first 3 or 4 or 5 or ... trials must result in 1 success and the next trial in a success.

∴ Required probability

$$\begin{aligned}
 &= 3C_1 p^2 q^2 + 4C_1 p^2 q^3 + 5C_1 p^2 q^4 + \dots \\
 &= p^2 \{3q^2 + 4q^3 + 5q^4 + \dots\} \\
 &= p^2 [(1-q)^{-2} - 1 - 2q] \\
 &= 1 - p^2 - 2p^2 q \\
 &= 1 - (0.02)^2 - 2 \times (0.02)^2 \times 0.98 \\
 &= 0.9988
 \end{aligned}$$

### Example 23

Find the probability that a person tossing 3 fair coins get either all heads or all tails for the second time on the fifth trial.

$$p = P(3 \text{ heads or 3 tails in tossing 3 coins})$$

$$= \frac{1}{8} + \frac{1}{8} = \frac{1}{4}$$

and  $q = \frac{3}{4}$ .

5th trial must result in the 2nd success.

∴ the first 4 trials must have resulted in 1 success and 3 failures.

$$\therefore \text{Required probability} = 3C_1 p^1 q^3 \times p$$

$$\begin{aligned}
 &= 4 \times \left(\frac{1}{4}\right)^2 \times \left(\frac{3}{4}\right)^3 \\
 &= \frac{27}{256}
 \end{aligned}$$

### Exercise 5(A)

#### Part-A (Short-answer Questions)

- The mean and variance of a binomial distribution are 4 and  $\frac{4}{3}$  respectively. Find  $P(X \geq 1)$ , if  $n = 6$ .
- If the recurrence relation for the central moments of the binomial distribution is  $\mu_{r+1} = pq \left( nr \mu_{r-1} + \frac{d\mu_r}{dp} \right)$ , find the value of  $\beta_1$ .
- In 256 sets of 8 tosses of a coin, in how many sets one may expect heads and tails in equal numbers?
- An experiment succeeds twice as often as it fails. Find the chance that in the next 4 trials, there shall be at least one success.
- In a family of 4 children, what is the probability that there will be at least 1 boy and at least 1 girl, assuming equal probability for boys and girls.

6. If  $X$  has the distribution  $B\left(25, \frac{1}{5}\right)$ , find  $P(X < \mu - 2\sigma)$  where  $\mu$  and  $\sigma^2$  are the mean and variance of the distribution.
7. Show that the largest value of the variance of a binomial distribution is  $\frac{n}{4}$ .
8. Find the mean and SD of the distribution whose moment generating function is  $0.4 e^t + 0.6$ .
9. When will the sum of 2 binomial variates having distributions  $B(n_1, p_1)$  and  $B(n_2, p_2)$  be also a binomial variate?
10. If  $X$  follows  $B\left(3, \frac{1}{3}\right)$  and  $Y$  follows  $B\left(5, \frac{1}{3}\right)$ , find  $P(X + Y \geq 1)$ .
11. Write down the pmf of the Poisson distribution which is approximately equivalent to  $B(100, 0.02)$ .
12. If  $X$  is a Poisson variate such that  $2P(X = 0) + P(X = 2) = 2P(X = 1)$ , find  $E(X)$ .
13. If  $X$  is a Poisson variate such that  $E(X^2) = 6$ , find  $E(X)$ .
14. If  $X$  is a Poisson variate such that  $P(X = 0) = 0.5$ , find  $\text{Var}(X)$ .
15. If  $X$  is a Poisson variate with parameter  $\lambda > 0$ , prove that  $E(X^2) = \lambda E(X + 1)$ .
16. If  $X$  is a Poisson variate with parameter  $\lambda$ , prove that
- $$E(X \text{ is even}) = \frac{1}{2} (1 + e^{-2\lambda})$$
17. If the MGF of a discrete RV  $X$  is  $e^{4(e^t - 1)}$ , find  $P(X = \mu + \sigma)$ , where  $\mu$  and  $\sigma$  are the mean and SD of  $X$ .
18. If  $X$  and  $Y$  are independent identical Poisson variates with mean 1, find  $P(X + Y = 2)$ .
19. Find the mean and variance of the discrete probability distribution given by  $P(X = r) = e^{-t} (1 - e^{-t})^{r-1}$   $r = 1, 2, 3, \dots, \infty$ .
20. If  $X$  is a geometric variate, taking values  $1, 2, 3, \dots, \infty$  find  $P(X \text{ is odd})$ .
21. Find the mean and variance of the distribution given by
- $$P(X = r) = \frac{2}{3^r}, r = 1, 2, \dots, \infty$$
22. For the geometric distribution of  $X$ , which represents the number of Bernoulli's trials required to get the first success  $\text{Var}(X) = 2E(X)$ . Find the pmf of the distribution.
23. Find the MGF of the geometric distribution, given by  $P(X = r) = q^{r-1} p$ ,  $r = 1, 2, \dots, \infty$ .
24. If the MGF of a discrete RV  $X$ , taking values  $1, 2, \dots, \infty$ , is  $e^t (5 - 4e^t)^{-1}$ , find the mean and variance of  $X$ .
25. Define hypergeometric distribution and give an example for the situation where it arises.
26. Write down the mean and variance of the hypergeometric distribution given by
- $$P(X = r) = k C_r (N - k) C_{n-r} / N C_n, r = 0, 1, 2, \dots$$

27. State the conditions under which the hypergeometric distribution tends to the binomial distribution. Hence, deduce the mean and variance of the binomial distribution from those of the hypergeometric distribution.

**Part-B**

28. It is known that diskettes produced by a certain company are defective with a probability 0.01 independently of each other. The company markets diskettes in packages of 10 and offers a money-back guarantee that atmost 1 of the 10 diskettes is defective. What proportion of diskettes are returned? If someone buys 3 diskettes, what is the probability that he will return exactly one of them?
29. Assuming that half the population is vegetarian and that 100 investigators each take 10 individuals to see whether they are vegetarians, how many would you expect to report that 3 people or less were vegetarians?
30. Show that, if 2 symmetric binomial distributions of degree  $n$  are so superposed that the  $r$ th term of the one coincides with the  $(r + 1)$ th term of the other, the distribution formed by adding superposed terms is a binomial distribution of degree  $(n + 1)$ .

[Hint:  $nC_{r-1} + nC_r = (n + 1)C_r$ ]

31. A factory produces 10 articles daily. It may be assumed that there is a constant probability  $p = 0.1$  of producing a defective article. Before these articles are stored, they are inspected and the defective ones are set aside. Suppose that there is a constant probability  $r = 0.1$  that a defective article is misclassified. If  $X$  denotes the number of articles classified as defective at the end of a production day, find (a)  $P(X = 3)$  and  $P(X > 3)$ .  
[Hint:  $P(\text{a defective article is classified as defective}) = P(\text{an article produced is defective}) \times P(\text{it is classified as defective}) = 0.1 \times 0.9 = 0.09$ ]
32. A fair coin is tossed 4 times. If  $X$  denotes the number of heads obtained and  $Y$  denotes the excess of the number of heads over the number of tails, obtained the probability mass function of  $Y$ .
33. An irregular 6-faced dice is thrown and the expectation that in 10 throws it will give 5 even numbers is twice the expectation that it will give 4 even numbers. How many times in 10,000 sets of 10 throws would you expect to give no even number?
34. If  $m$  things are distributed among  $a$  men and  $b$  women, show that the probability that the number of things received by men is odd is
- $$\frac{1}{2} \left[ \frac{(b+a)^m - (b-a)^m}{(b+a)^m} \right].$$
- [Hint:  $P(\text{a thing is received by men}) = p = \frac{a}{a+b}$  and  $q = \frac{b}{a+b}$ ]
35. At least one half of an airplane's engines are required to function in order for it to operate. If each engine independently functions with probability

- $p$ , for what values of  $p$  is a 4-engine plane to be preferred for operation to a 2-engine plane?
36. At least one half of an airplane's engines are required to function in order for it to operate. If each engine functions independently with probability of failure  $q$ , for what values of  $q$  is a 2-engine plane to be preferred for operation to a 4-engine plane?
37. If a fair coin is flipped an even number  $2n$  times, show that the probability of getting more heads than tails is  $\frac{1}{2} \left[ 1 - 2nC_n \left( \frac{1}{2} \right)^{2n} \right]$ ,
- [**Hint:**  $P(\text{more heads than tails}) = P(\text{less heads than tails}) = \frac{1}{2} [1 - P(\text{equal number of heads and tails})]$ ]
38. If a fair coin is tossed at random 5 independent times, find the conditional probability of 5 heads relative to the hypothesis that there are at least 4 heads.
39. A factory has 10 machines which may need adjustment from time to time during the day. Three of these machines are old, each having a probability of  $\frac{1}{11}$  needing adjustment during the day and 7 are new, having the corresponding probability of  $\frac{1}{21}$ . Assuming that no machine needs adjustment twice on the same day, find the probabilities that on a particular day
- (i) just 2 old and no new machine need adjustment and
  - (ii) just 2 machines that need adjustment are of the same type.
40. The probability of a man hitting a target is  $\frac{1}{4}$ . (i) If he fires 7 times, what is the probability of his hitting the target at least twice? and (ii) How many times must he fire so that the probability of his hitting the target at least once is greater than  $\frac{2}{3}$ ?
41. A set of 6 similar coins are tossed 640 times with the following results:
- |                  |   |    |     |     |     |    |    |
|------------------|---|----|-----|-----|-----|----|----|
| Number of heads: | 0 | 1  | 2   | 3   | 4   | 5  | 6  |
| Frequency :      | 7 | 64 | 140 | 210 | 132 | 75 | 12 |
- Calculate the binomial frequencies on the assumption that the coins are symmetrical.
42. Fit a binomial distribution for the following data and hence find the theoretical frequencies:
- |      |   |    |    |    |   |
|------|---|----|----|----|---|
| $x:$ | 0 | 1  | 2  | 3  | 4 |
| $f:$ | 5 | 29 | 36 | 25 | 5 |
43. A car hire firm has 2 cars which it hires out day by day. The number of demands for a car on each day follows a Poisson distribution with mean

- 1.5. Calculate the proportion of days on which (i) neither car is used and (ii) some demand is not fulfilled.
44. The proofs of a 500-page book contains 500 misprints. Find the probability that there are at least 4 misprints in a randomly chosen page.
45. If the average number of claims handled daily by an insurance company is 5, what proportion of days will have less than 3 claims? What is the probability that there will be 4 claims in exactly 3 of the next 5 days. Assume that the number of claims on different days are independent.
46. In a certain factory producing razor blades, there is a small change  $\frac{1}{500}$  for any blade to be defective. The blades are supplied in packets of 10. Use Poisson distribution to calculate the approximate number of packets containing (i) no defective blade, (ii) at least 1 defective blade and (iii) at most 1 defective blade in a consignment of 10,000 packets.
47. An insurance company has discovered that only about 0.1% of the population is involved in a certain type of accident each year. If its 10,000 policy holders were randomly selected from the population, what is the probability that not more than 5 of its clients are involved in such an accident next year?
48. In a given city, 4% of all licensed drivers will be involved in at least 1 road accident in any given year. Determine the probability that among 150 licensed drivers randomly chosen in this city  
(i) only 5 will be involved in at least 1 accident in any given year, and  
(ii) at most 3 will be involved in at least 1 accident in any given year.
49. A radioactive source emits on the average 2.5 particles per second. Find the probability that 3 or more particles will be emitted in an interval of 4s.
50. It has been established that the number of defective stereos produced daily at a certain plant is Poisson distributed with a mean of 4. Over a 2-day span, what is the probability that the number of defective stereos does not exceed 3?
51. In an industrial complex, the average number of fatal accidents per month is one-half. The number of accidents per month is adequately described by a Poisson distribution. What is the probability that 6 months will pass without a fatal accident?
52. If the numbers of telephone calls coming into a telephone exchange between 9 am and 10 am and between 10 am and 11 am are independent and follow Poisson distributions with parameters 2 and 6 respectively, what is the probability that more than 5 calls come between 9 am and 11 am?
53. Patients arrive randomly and independently at a doctor's consulting room from 5 pm at an average rate of one in 5 min. The waiting room can hold 12 persons. What is the probability that the room will be full, when the doctor arrives at 6 pm?

54. The number of blackflies on a broad bean leaf follows a Poisson distribution with mean 2. A plant inspector, however, records the number of flies on a leaf only if at least 1 fly is present. What is the probability that he records 1 or 2 flies on a randomly chosen leaf? What is the expected number of flies recorded per leaf?  
**[Hint:** If  $X$  is the number of flies on a leaf, we have to find  $P\{X = r/X \geq 1\}$ ,  $r = 1, 2$ , and add them.]
55. A radioactive source emits particles at a rate of 10 per minute in accordance with Poisson law. Each particle emitted has a probability of  $\frac{2}{5}$  being recorded. Find the probability that at least 4 particles are recorded in a 2-min period.
56. Fit a Poisson distribution for the following distribution and hence find the expected frequencies.
- |       |     |     |     |    |    |   |   |
|-------|-----|-----|-----|----|----|---|---|
| $x :$ | 0   | 1   | 2   | 3  | 4  | 5 | 6 |
| $f :$ | 314 | 335 | 204 | 86 | 29 | 9 | 3 |
57. If the probability that a certain test yields a positive reaction equals 0.4, what is the probability that fewer than 5 negative reactions occur before the first positive one?
58. In a test a light switch is turned on and off until it fails. If the probability that the switch will fail any time it is turned on or off is 0.001, what is the probability that the switch will not fail during the first 800 times it is turned on or off?
59. An item is inspected at the end of each day to see whether it is still functioning properly. If it is found to fail at the 10th inspection and not earlier, what is the maximum value of the probability of its failure on any day?
60. If  $X$  and  $Y$  are 2 independent RVs, each representing the number of failures preceding the first success in a sequence of Bernoulli's trials with  $p$  as the probability of success in a single trial, show that  $P(X = Y) = \frac{p}{1 + q}$ , where  $p + q = 1$ .
61. A throws 2 dices until he gets 6 and  $B$  throws independently 2 other dice until he gets 7. Find the probability that  $B$  will require more throws than  $A$ .
62. If 2 independent RVs  $X$  and  $Y$  have identical geometric distributions with parameter  $p$ , find the probability mass function of  $(X + Y)$  and hence the expected value of  $(X + Y)$ .
63. As part of an air-pollution survey, an inspector decides to examine the exhaust of 6 of a company's 24 trucks. If four of the company's trucks emit excessive amounts of pollutants, what is the probability that none of them will be included in the inspector's sample?
64. Among the 120 applicants for a job, only 80 are actually qualified. If 5 of the applicants are randomly selected for an in-depth interview, find

the probability that only 2 of them will be qualified for the job using (i) hypergeometric probability, and (ii) binomial approximation to hypergeometric probability.

65. If an auditor selects 5 returns from among 15 returns of which 9 contains illegitimate deductions, what is the probability that a majority of the selected returns contains illegitimate deductions?
66. In sampling a lot of 100 items, the sampling plan calls for inspections of 20 pieces. Find the probability of accepting a lot with 5 defectives, if we allow 1 defective in the sample.
67. If  $X$  and  $Y$  are independent binomial random variables having respective parameters  $(n, p)$  and  $(m, p)$ , prove that the conditional probability mass function of  $X$ , given that  $X + Y = k$ , is that of a hypergeometric RV.  
[Hint: The joint probability mass function of  $X$  and  $Y$  is that of a binomial distribution with parameters  $(m + n, p)$ .]
68. A student takes a 5-answer multiple choice test orally. He continues to answer questions until he gets 5 correct answers. What is the probability that he gets them on the 12<sup>th</sup> question, if he guesses at each answer?
69. A consignment of 15 tubes contains 4 defectives. The tubes are selected at random, one by one, and examined. Assuming that the tubes tested are not put back, what is the probability that the ninth one examined is the last defective?
70. A machine is known to produce 3% defective items. What is the probability that at least 5 items are to be examined in order to get 2 defective items?

## SPECIAL CONTINUOUS DISTRIBUTIONS

### 1. Uniform or Rectangular Distribution

**Definition:** A continuous RV  $X$  is said to follow a *uniform* or *rectangular distribution* in any finite interval, if its probability density function is a constant in that interval.

If  $X$  follows a uniform distribution in  $a < x < b$ , then  $f(x) = \frac{1}{b-a}$  in  $a < x < b$ , as explained below:

When  $X$  follows a uniform distribution in  $(a, b)$ ,

$$f(x) = k.$$

By the basic property of a probability density function,

$$\int_{R_X} f(x) dx = 1$$

$$\therefore \int_a^b k dx = 1$$

$$\therefore k = \frac{1}{b-a}.$$

Thus, for a uniform distribution in  $(a, b)$ ,

$$f(x) = \frac{1}{b-a}.$$

When  $X$  follows a uniform distribution in  $(a, b)$ , it is symbolically written as:  $X$  follows  $U(a, b)$ .

### **Moments of the Uniform Distribution $U(a, b)$**

Raw moments  $\mu_r$  of the uniform distribution  $U(a, b)$  about the origin are given by

$$\mu'_r = E\{X^r\}, \text{ where } X \text{ follows } U(a, b)$$

$$\begin{aligned} &= \int_a^b x^r \frac{1}{b-a} dx \\ &= \frac{b^{r+1} - a^{r+1}}{(r+1)(b-a)} \end{aligned} \tag{1}$$

$$\therefore E(X) = \text{Mean of } U(a, b) = \mu'_1 = \frac{1}{2}(b+a) \tag{2}$$

Central moments  $\mu_r$  of the uniform distribution  $U(a, b)$  are given by

$$\begin{aligned} \mu_r &= E[\{X - E(X)\}^r] \\ &= E[\{X - \frac{1}{2}(b+a)\}^r] \\ &= \int_a^b \frac{\left\{x - \frac{1}{2}(b+a)\right\}^r}{b-a} dx \\ &= \frac{1}{b-a} \int_{-c}^c t^r dt, \text{ on putting } t = x - \frac{1}{2}(b+a) \text{ and } c = \frac{1}{2}(b-a) \\ &= \begin{cases} 0 & \text{if } r \text{ is odd} \\ \frac{1}{r+1} \cdot \left(\frac{b-a}{2}\right)^r & \text{if } r \text{ is even} \end{cases} \end{aligned}$$

$$\text{Thus, } \mu_{2n-1} = 0 \text{ and } \mu_{2n} = \frac{1}{2n+1} \left(\frac{b-a}{2}\right)^{2n} \text{ for } n = 1, 2, 3, \dots \tag{3}$$

$$\text{In particular, } \mu_2 = \text{variance of } U(a, b) = \frac{1}{12}(b-a)^2 \tag{4}$$

$$\mu_3 = 0 \text{ and } \mu_4 = \frac{1}{80} (b-a)^4$$

The absolute central moments  $v_r$  of the uniform distribution  $U(a, b)$  are given by

$$\begin{aligned}
v_r &= E\{|X - E(X)|^r\} \\
&= \int_a^b \frac{\left|x - \frac{1}{2}(b+a)\right|^r}{b-a} dx \\
&= \frac{1}{b-a} \int_{-c}^c |t|^r dt, \text{ on putting } t = x - \frac{1}{2}(b+a) \text{ and } c = \frac{1}{2}(b-a) \\
&= \frac{2}{b-a} \int_0^c t^r dt, \quad (\text{since the integrand is an even function of } t) \\
&= \frac{1}{r+1} \cdot \left(\frac{b-a}{2}\right)^r
\end{aligned} \tag{5}$$

**Definition:**  $E\{|X - E(X)|\}$  is called the *mean deviation* (MD) about the mean of the RV  $X$  or of the corresponding distribution.

Thus, the MD about the mean of the distribution  $U(a, b)$  is given by

$$v_1 = \frac{1}{4} (b-a)$$

## 2. Exponential Distribution

**Definition:** A continuous RV  $X$  is said to follow an *exponential distribution* or *negative exponential distribution* with parameter  $\lambda > 0$ , if its probability density function is given by

$$f(x) = \begin{cases} \lambda e^{-\lambda x} & x \geq 0 \\ 0 & \text{otherwise} \end{cases}$$

We note that  $\int_0^\infty f(x) dx = \int_0^\infty \lambda e^{-\lambda x} dx = 1$  and hence  $f(x)$  is a legitimate density function.

### Mean and Variance of the Exponential Distribution

Raw moments  $\mu'_r$  about the origin of the exponential distribution are given by

$$\begin{aligned}
\mu'_r &= E(X^r) \\
&= \int_0^\infty x^r \cdot \lambda e^{-\lambda x} dx \\
&= \frac{1}{\lambda^r} \int_0^\infty y^r e^{-y} dy, \text{ (on putting } y = \lambda x)
\end{aligned}$$

$$\begin{aligned}
 &= \frac{1}{\lambda^r} \sqrt{(r+1)} \\
 &= \frac{r!}{\lambda^r}
 \end{aligned} \tag{1}$$

$\therefore E(X) = \text{Mean of the exponential distribution}$

$$= \mu_1' = \frac{1}{\lambda}, [\text{from (1)}]$$

Putting  $r = 2$  in (1), we get

$$\begin{aligned}
 \mu_2' &= \frac{2}{\lambda^2} \\
 \therefore \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
 &= \frac{2}{\lambda^2} - \frac{1}{\lambda^2} = \frac{1}{\lambda^2}
 \end{aligned}$$

### **Memoryless Property of the Exponential Distribution**

If  $X$  is exponentially distributed, then

$$P(X > s + t | X > s) = P(X > t), \text{ for any } s, t > 0$$

$$\begin{aligned}
 P(X > k) &= \int_k^{\infty} \lambda e^{-\lambda x} dx \\
 &= (-e^{-\lambda x})_k^{\infty} = e^{-\lambda k}
 \end{aligned} \tag{1}$$

Now  $P(X > s + t | X > s) = \frac{P\{X > s + t \text{ and } X > s\}}{P\{X > s\}}$

$$\begin{aligned}
 &= \frac{P\{X > s + t\}}{P\{X > s\}} \\
 &= \frac{e^{-\lambda(s+t)}}{e^{-\lambda s}}, [\text{by (1)}] \\
 &= e^{-\lambda t} = P(X > t).
 \end{aligned}$$

**Note** The converse of this result is also true. That is, if  $P(X > s + t | X > s) = P(X > t)$ , then  $X$  follows an exponential distribution. See Example (8) in Worked Example 5(b).]

### **3. Erlang Distribution, or General Gamma Distribution**

**Definition:** A continuous RV  $X$  is said to follow an *Erlang distribution* or *General Gamma distribution* with parameters  $\lambda > 0$  and  $k > 0$ , if its probability density function is given by

$$f(x) = \begin{cases} \frac{\lambda^k x^{k-1} e^{-\lambda x}}{\Gamma(k)}, & \text{for } x \geq 0 \\ 0, & \text{otherwise} \end{cases}$$

$$\begin{aligned} \text{We note that } \int_0^{\infty} f(x) dx &= \frac{\lambda^k}{\Gamma(k)} \int_0^{\infty} x^{k-1} e^{-\lambda x} dx \\ &= \frac{1}{\Gamma(k)} \int_0^{\infty} t^{k-1} e^{-t} dt, \text{ [on putting } \lambda x = t] \end{aligned}$$

Hence,  $f(x)$  is a legitimate density function.

- Note**
- When  $\lambda = 1$ , the Erlang distribution is called Gamma distribution or simple Gamma distribution with parameter  $k$  whose density function is  $f(x) = \frac{1}{\Gamma(k)} x^{k-1} e^{-x}$ ,  $x \geq 0$ ;  $k > 0$ .
  - When  $k = 1$ , the Erlang distribution reduces to the exponential distribution with parameter  $\lambda > 0$ .
  - Sometimes, the Erlang distribution itself is called Gamma distribution.

### Mean and Variance of the Erlang Distribution

The raw moments  $\mu'_r$  about the origin of the Erlang distribution are given by

$$\begin{aligned} \mu'_r &= E(X^r) \\ &= \int_0^{\infty} \frac{\lambda^k}{\Gamma(k)} x^{k+r-1} e^{-\lambda x} dx \\ &= \frac{\lambda^k}{\Gamma(k)} \cdot \frac{1}{\lambda^{k+r}} \int_0^{\infty} t^{k+r-1} e^{-t} dt, \text{ (on putting } \lambda x = t) \\ &= \frac{1}{\lambda^r} \frac{\Gamma(k+1)}{\Gamma(k)} \end{aligned}$$

$$\text{Mean} = E(X) = \frac{1}{\lambda} \cdot \frac{\Gamma(k+1)}{\Gamma(k)} = \frac{k}{\lambda}$$

$$\begin{aligned} \text{Var}(X) &= E(X^2) - [E(X)]^2 \\ &= \frac{1}{\lambda^2} \cdot \frac{\Gamma(k+2)}{\Gamma(k)} - \left( \frac{k}{\lambda} \right)^2 \\ &= \frac{1}{\lambda^2} \{k(k+1) - k^2\} = \frac{k}{\lambda^2} \end{aligned}$$

### Reproductive Property of the Gamma Distribution

The sum of a finite number of independent Erlang variables is also an Erlang variable. That is, if  $X_1, X_2, \dots, X_n$  are independent Erlang variables with parameters  $(\lambda, k_1), (\lambda, k_2), \dots, (\lambda, k_n)$ , then  $X_1 + X_2 + \dots + X_n$  is also an Erlang variable with parameter  $(\lambda, k_1 + k_2 + \dots + k_n)$ .

Let us first find the moment generating function of the Erlang variable  $X$  with parameters  $\lambda$  and  $k$  and use it to prove this property. MGF of  $X$  is given by

$$\begin{aligned}
 M_X(t) &= E\{e^{tx}\} \\
 &= \int_0^{\infty} \frac{\lambda^k}{(k)} x^{k-1} e^{-\lambda x} e^{-tx} dx \\
 &= \frac{\lambda^k}{(k)} \int_0^{\infty} x^{k-1} e^{-(\lambda-t)x} dx \\
 &= \frac{\lambda^k}{(k)} \cdot \frac{1}{(\lambda-t)^k} \int_0^{\infty} y^{k-1} e^{-y} dy, \text{ (on putting } \lambda-t=y) \\
 &= \left( \frac{\lambda}{\lambda-t} \right)^k [\because \text{the integral} = \frac{1}{(k)}] \\
 &= \left( 1 - \frac{t}{\lambda} \right)^{-k}
 \end{aligned}$$

Now,  $M_{(X_1 + X_2 + \dots + X_n)}(t) = M_{X_1}(t)M_{X_2}(t)\dots M_{X_n}(t)$  (since  $X_1, X_2, \dots, X_n$  are independent)

[Refer to Property (4) of MGF given in Section 4(b) of Chapter 4.]

$$\begin{aligned}
 &= \left( 1 - \frac{t}{\lambda} \right)^{-k_1} \left( 1 - \frac{t}{\lambda} \right)^{-k_2} \dots \left( 1 - \frac{t}{\lambda} \right)^{-k_n} \\
 &= \left( 1 - \frac{t}{\lambda} \right)^{-(k_1 + k_2 + \dots + k_n)}
 \end{aligned}$$

which is the MGF of an Erlang variable with parameters  $(\lambda, k_1 + k_2 + \dots + k_n)$ . Hence the reproductive property.

### **Relation Between the Distribution Functions (cdf) of the Erlang Distribution with $\lambda = 1$ (or Simple Gamma Distribution) and Poisson Distribution**

If  $X$  is a Poisson random variable with mean  $\lambda$ ,

$$\text{then } P(X \leq K) = \sum_{r=0}^k \frac{e^{-\lambda} \lambda^r}{r!} \quad (1)$$

Differentiating both sides with respect to  $\lambda$ , we get

$$\begin{aligned}
 \frac{d}{d\lambda} P(X \leq k) &= \sum_{r=0}^k \frac{1}{r!} \{e^{-\lambda} \cdot r \lambda^{r-1} - e^{-\lambda} \cdot \lambda^r\} \\
 &= e^{-\lambda} \sum_{r=0}^k \left[ \frac{\lambda^{r-1}}{(r-1)!} - \frac{\lambda^r}{r!} \right]
 \end{aligned}$$

$$\begin{aligned}
&= e^{-\lambda} \left[ -1 + \left( 1 - \frac{\lambda}{1!} \right) + \left( \frac{\lambda}{1!} - \frac{\lambda^2}{2!} \right) + \dots + \left\{ \frac{\lambda^{k-1}}{(k-1)!} - \frac{\lambda^k}{k!} \right\} \right] \\
&= \frac{e^{-\lambda} \lambda^k}{k!}
\end{aligned} \tag{2}$$

Integrating both sides of (2) with respect to  $\lambda$  from  $\lambda$  to  $\infty$ , we get

$$\begin{aligned}
\left[ \sum_{r=0}^K \frac{e^{-\lambda} \lambda^r}{r!} \right]_{\lambda}^{\infty} &= - \int_{\lambda}^{\infty} \frac{1}{k!} e^{-\lambda} \lambda^k d\lambda \\
\text{i.e., } \sum_{r=0}^K \frac{e^{-\lambda} \lambda^r}{r!} &= \int_{\lambda}^{\infty} \frac{1}{k!} e^{-y} y^k dy \\
\text{i.e., } P(X \leq k) &= P(Y \geq \lambda), \\
&\quad [\text{where } Y \text{ is the Erlang variable with parameters 1 and } (k+1)] \\
\text{or } P(X \leq k) &= 1 - P(Y \leq \lambda)
\end{aligned}$$

**Note** The above relationship is valid only when the parameter  $k$  is a positive integer.

## 4. Weibull Distribution

**Definition:** A continuous RV  $X$  is said to follow a *Weibull distribution* with parameters  $\alpha, \beta > 0$ , if the RV  $Y = \alpha X^\beta$  follows the exponential distribution with density function  $f_Y(y) = e^{-y}$ ,  $y > 0$ .

### Density Function of the Weibull Distribution

Since  $Y = \alpha \cdot X^\beta$ , we have  $y = \alpha \cdot x^\beta$ .

By the transformation rule, derived in Chapter 3, we have  $f_X(x) = f_Y(y) \left| \frac{dy}{dx} \right|$ , where  $f_X(x)$  and  $f_Y(y)$  are the density functions of  $X$  and  $Y$  respectively.

$$\begin{aligned}
\therefore f_X(x) &= e^{-y} \alpha \beta x^{\beta-1} \\
&= \alpha \beta x^{\beta-1} e^{-\alpha x^\beta}; x > 0 \quad [\because y > 0]
\end{aligned}$$

**Note** When  $\beta = 1$ , Weibull distribution reduces to the exponential distribution with parameter  $\alpha$ .

### Mean and Variance of the Weibull Distribution

The raw moments  $\mu'_r$  about the origin of the Weibull distribution are given by

$$\begin{aligned}
\mu'_r &= E(X^r) \\
&= \alpha \beta \int_0^{\infty} x^{r+\beta-1} e^{-\alpha x^\beta} dx \\
&= \int_0^{\infty} \left( \frac{y}{\alpha} \right)^{\frac{r}{\beta} + 1 - \frac{1}{\beta}} e^{-y} \left( \frac{y}{\alpha} \right)^{\frac{1}{\beta} - 1} dy
\end{aligned}$$

on putting  $y = \alpha x^\beta$  or  $x = \left(\frac{y}{\alpha}\right)^{\frac{1}{\beta}}$

$$= \alpha^{-r/\beta} \int_0^{\infty} y^{r/\beta} e^{-y} dy$$

$$= \alpha^{-r/\beta} \sqrt{\left(\frac{r}{\beta} + 1\right)}$$

$$\therefore \text{Mean} = E(X) = \mu_1' = \alpha^{-\frac{1}{\beta}} \sqrt{\left(\frac{1}{\beta} + 1\right)}$$

$$\text{Var}(X) = E(X^2) - \{E(X)\}^2$$

$$= \alpha^{-2/\beta} \left[ \left( \frac{2}{\beta} + 1 \right) - \left\{ \left( \frac{1}{\beta} + 1 \right) \right\}^2 \right]$$

**Note** Weibull distribution finds frequent applications in Reliability Theory. It is assumed as the probability distribution of the time to failure (or length of life) of a component in a system. Other distributions used to describe the failure law are the exponential and normal distributions. see Example (19) in Worked Example 5(b).

## 5. Normal (or Gaussian) Distribution

**Definition:** A continuous RV  $X$  is said to follow a *normal distribution* or *Gaussian distribution* with parameters  $\mu$  and  $\sigma$ , if its probability density function is given by

$$f(x) = \frac{1}{\sigma \sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2}; -\infty < x < \infty$$

$$-\infty < \mu < \infty \quad \sigma > 0 \quad (1)$$

Symbolically, ' $X$  follows  $N(\mu, \sigma)$ '. Sometimes it is also given as  $N(\mu, \sigma^2)$ . We shall use only the notation  $N(\mu, \sigma)$  as in the earlier chapters.

$f(x)$  is a legitimate density function, as

$$\begin{aligned} \int_{-\infty}^{\infty} f(x) dx &= \frac{1}{\sigma \sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-(x-\mu)^2/2\sigma^2} dx \\ &= \frac{1}{\sigma \sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-t^2} \sigma \sqrt{2} dt, \quad \left( \text{on putting } t = \frac{x-\mu}{\sigma \sqrt{2}} \right) \\ &= \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} e^{-t^2} dt \\ &= \frac{1}{\sqrt{\pi}} 2 \int_0^{\infty} e^{-t^2} dt = \frac{1}{\sqrt{\pi}} \cdot \sqrt{\left(\frac{1}{2}\right)} = \frac{1}{\sqrt{\pi}} \cdot \sqrt{\pi} = 1 \end{aligned}$$

**(a) Standard Normal Distribution**

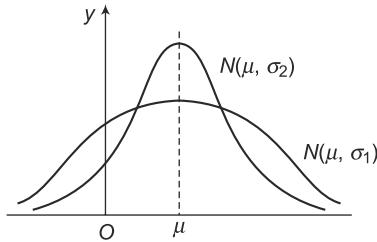
The normal distribution  $N(0, 1)$  is called the standardised or simply the standard normal distribution, whose density function is given by

$$\phi(z) = \frac{1}{\sqrt{2\pi}} e^{-z^2/2} \quad -\infty < z < \infty$$

This is obtained by putting  $\mu = 0$  and  $\sigma = 1$  and by changing  $x$  and  $f$  respectively into  $z$  and  $\phi$ . If  $X$  has distribution  $N(\mu, \sigma)$  and if  $Z = \frac{X - \mu}{\sigma}$ , then we can prove that  $Z$  has distribution  $N(0, 1)$ .

[See the corollary under the property (6) of normal distribution]

The importance of  $N(0, 1)$  is due to the fact that the values of  $\phi(z)$  and  $\int_0^z \phi(z) dz$  are tabulated.

**(b) Normal Probability Curve**

**Fig. 5.1**

The graph of  $y = f(x)$ , that is given above for  $\sigma = \sigma_1$  and  $\sigma_2$ , is a well-known bell-shaped curve and is called the normal probability curve (Fig. 5.1).

The curve is symmetrical about the ordinate at  $x = \mu$ . The ordinate  $f(x)$  decreases rapidly as  $x$  increases numerically, the maximum (occurring at  $x = \mu$ ) given by  $\frac{1}{\sigma\sqrt{2\pi}}$ . The curve extends up to infinity on either side of  $x = \mu$  and the  $x$ -axis is an asymptote to the curve.

The graph is concave downward at  $x = \mu$  and it is concave upward for large numerical values of  $x$ . The points at which the concavity changes are called the points of inflexion of the curve. They are found by solving the equation  $y'' = 0$  [i.e.,  $f''(x) = 0$ ]. We can prove that the points of inflexion of the normal probability curve occur at  $x = \mu \pm \sigma$ , that is, at points which are at a distance of  $\sigma$  on either side of  $x = \mu$ . Thus, if  $\sigma$  is relatively large, the curve tends to be flat, while if  $\sigma$  is small, the curve tends to be peaked. Hence, the steepness of the curve is determined by  $\sigma$ . The two curves given in the figure relate to  $N(\mu, \sigma_1)$  and  $N(\mu, \sigma_2)$ , where  $\sigma_1 > \sigma_2$ .

**(i) Properties of the Normal Distribution  $N(\mu, \sigma)$** 

If  $X$  follows  $N(\mu, \sigma)$ , then  $E(X) = \mu$  and  $\text{Var}(X) = \sigma^2$

$$\begin{aligned}
 E(X) &= \int_{-\infty}^{\infty} x f(x) dx \\
 &= \frac{1}{\sigma \sqrt{2\pi}} \int_{-\infty}^{\infty} x e^{-(x-\mu)^2/2\sigma^2} dx \\
 &= \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} (\mu + \sqrt{2} \sigma t) e^{-t^2} dt, \quad \left( \text{on putting } t = \frac{x-\mu}{\sigma \sqrt{2}} \right) \\
 &= \frac{\mu}{\sqrt{\pi}} \int_{-\infty}^{\infty} e^{-t^2} dt + \sqrt{\frac{2}{\pi}} \sigma \int_{-\infty}^{\infty} t e^{-t^2} dt \\
 &= \frac{\mu}{\sqrt{\pi}} \cdot \sqrt{\pi} = \mu.
 \end{aligned}$$

(since the integrand in the second integral is an odd function of  $t$ )

$$\begin{aligned}
 E(X^2) &= \frac{1}{\sigma \sqrt{2}} \int_{-\infty}^{\infty} x^2 e^{-(x-\mu)^2/2\sigma^2} dx \\
 &= \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} (\mu + \sqrt{2} \sigma t)^2 e^{-t^2} dt, \quad \left( \text{on putting } t = \frac{x-\mu}{\sigma \sqrt{2}} \right) \\
 &= \frac{1}{\sqrt{\pi}} \left[ \int_{-\infty}^{\infty} \mu^2 e^{-t^2} dt + 2\sqrt{2} \mu \sigma \int_{-\infty}^{\infty} t e^{-t^2} dt + 2\sigma^2 \int_{-\infty}^{\infty} t^2 e^{-t^2} dt \right] \\
 &= \mu^2 + 0 + \frac{2\sigma^2}{\sqrt{\pi}} \int_0^{\infty} t e^{-t^2} 2t dt, \quad (\because t^2 e^{-t^2} \text{ is even}) \\
 &= \mu^2 + \frac{2\sigma^2}{\sqrt{\pi}} \int_0^{\infty} u^{\frac{1}{2}} e^{-u} du, \quad (\text{on putting } u = t^2) \\
 &= \mu^2 + \frac{2\sigma^2}{\sqrt{\pi}} \left[ \left( \frac{3}{2} \right) \right] \\
 &= \mu^2 + 2 \frac{\sigma^2}{\sqrt{\pi}} \cdot \frac{1}{2} \left[ \left( \frac{1}{2} \right) \right] \\
 &= \mu^2 + \sigma^2 \quad \left( \because \left[ \left( \frac{1}{2} \right) \right] = \sqrt{\pi} \right) \\
 \therefore \text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
 &= \sigma^2
 \end{aligned}$$

**(ii) Median and Mode of the Normal Distribution  $N(\mu, \sigma)$** 

**Definition:** If  $X$  is a continuous RV with density function  $f(x)$ , then  $M$  is called the median value of  $X$ , provided that

$$\int_{-\infty}^M f(x) dx = \int_M^\infty f(x) dx = \frac{1}{2}$$

For the normal distribution  $N(\mu, \sigma)$ , the median  $M$  is given by

$$\int_M^\infty \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx = \frac{1}{2}$$

$$\text{i.e., } \int_M^\mu \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx + \int_\mu^\infty \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx = \frac{1}{2}$$

$$\text{i.e., } \int_M^\mu \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx + \frac{1}{2} = \frac{1}{2},$$

$$\left. \begin{array}{l} \text{since } \int_{-\infty}^\infty \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx = 1 \text{ and the normal curve is} \\ \text{symmetrical about } x = \mu \end{array} \right\}$$

$$\text{i.e., } \int_M^\mu f(x) dx = 0$$

$$\therefore M = \mu$$

**Definition:** Mode of a continuous RV  $X$  is defined as the value of  $x$  for which the density function  $f(x)$  is maximum.

For the normal distribution  $N(\mu, \sigma)$ ,

$$f(x) = \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} \quad -\infty < x < \infty$$

$$\therefore \log f(x) = k - \frac{1}{2\sigma^2} (x - \mu)^2$$

Differentiating with respect to  $x$ ,

$$\frac{f'(x)}{f(x)} = -\frac{1}{\sigma^2} (x - \mu)$$

$$\text{i.e., } f'(x) = -\frac{1}{\sigma^2} (x - \mu) f(x) \\ = 0, \text{ when } x = \mu$$

$$f''(x) = -\frac{1}{\sigma^2} \{(x - \mu) f'(x) + f(x)\}$$

$$\therefore [f''(x)]_{x=\mu} = -\frac{1}{\sigma^2} f(\mu) < 0$$

Therefore,  $f(x)$  is maximum at  $x = \mu$ . That is,

Mode of the distribution  $N(\mu, \sigma) = \mu$ .

**Note** For the normal distribution, mean, median and mode are equal.

**(iii) Central Moments of the Normal Distribution  $N(\mu, \sigma)$** 

Central moments  $\mu_r$  of  $N(\mu, \sigma)$  are given by  $\mu_r = E(x - \mu)^r$

$$\begin{aligned} &= \frac{1}{\sigma \sqrt{2\pi}} \int_{-\infty}^{\infty} (x - \mu)^r e^{-(x - \mu)^2/2\sigma^2} dx \\ &= \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} (\sqrt{2} \sigma t)^r e^{-t^2} dt \\ &= \frac{2^{r/2} \sigma^r}{\sqrt{\pi}} \int_{-\infty}^{\infty} t^r e^{-t^2} dt \end{aligned}$$

**Case (A):**  $r$  is an odd integer, that is,  $r = 2n + 1$ .

$$\begin{aligned} \therefore \mu_{2n+1} &= \frac{2^{(n+1)/2} \sigma^{2n+1}}{\sqrt{\pi}} \int_{-\infty}^{\infty} t^{2n+1} e^{-t^2} dt \\ &= 0, \text{ (since the integrand is an odd function of } t) \end{aligned}$$

**Case (B):**  $r$  is an even integer, that is,  $r = 2n$

$$\begin{aligned} \therefore \mu_{2n} &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \int_{-\infty}^{\infty} t^{2n} e^{-t^2} dt \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \cdot 2 \int_0^{\infty} t^{2n} e^{-t^2} dt \\ &\quad (\because \text{the integrand is an even function of } t) \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \int_{-\infty}^{\infty} u^{n-\frac{1}{2}} e^{-u} du. \quad [\text{on putting } u = t^2] \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \left[ \left( n + \frac{1}{2} \right) \right] \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \left( \frac{2n-1}{2} \right) \left[ \left( \frac{2n-1}{2} \right) \right] \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \left( \frac{2n-1}{2} \right) \cdot \left( \frac{2n-3}{2} \right) \left[ \left( \frac{2n-3}{2} \right) \right] \\ &= \frac{2^n \sigma^{2n}}{\sqrt{\pi}} \frac{2n-1}{2} \cdot \frac{2n-3}{2} \cdots \frac{1}{2} \left[ \left( \frac{1}{2} \right) \right] \\ &= 1.3.5 \dots (2n-1) \sigma^{2n} \end{aligned} \tag{1}$$

From (1), we get,

$$\mu_{2n-2} = \frac{2^{n-1} \sigma^{2n-2}}{\sqrt{\pi}} \cdot \left[ \left( n - \frac{1}{2} \right) \right] \tag{2}$$

From (1) and (2), we get

$$\frac{\mu_{2n}}{\mu_{2n-2}} = 2\sigma^2(n - 1/2)$$

i.e.,

$$\mu_{2n} = (2n - 1)\sigma^2 \mu_{2n-2} \quad (3)$$

(3) gives a recurrence relation for the even order central moments of the normal distribution  $N(\mu, \sigma)$ .]

#### (iv) Mean Deviation about the Mean of the Normal Distribution $N(\mu, \sigma)$

**Definition:** The absolute (central) moment of the first order of a RV  $X$  is called the mean deviation about the mean of  $X$ .

i.e., MD about the mean =  $E\{|x - E(X)|\}$

For the normal distribution  $N(\mu, \sigma)$ ,

$$\begin{aligned} \text{The MD about the mean} &= \int_{-\infty}^{\infty} |x - \mu| \frac{1}{\sigma \sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx \\ &= \frac{1}{\sigma \sqrt{2\pi}} \int_{-\infty}^{\infty} |\sqrt{2} \sigma t| e^{-t^2} \sqrt{2} \sigma dt \\ &= \sqrt{\frac{2}{\pi}} \sigma \int_{-\infty}^{\infty} |t| e^{-t^2} dt \\ &= 2 \sqrt{\frac{2}{\pi}} \sigma \int_{-\infty}^{\infty} t e^{-t^2} dt, \\ &\quad (\text{since the integrand is an even function of } t) \\ &= \sqrt{\frac{2}{\pi}} \sigma (-e^{-t^2})_0^\infty \\ &= \sqrt{\frac{2}{\pi}} \sigma = \frac{4}{5} \sigma \quad (\text{approximately}). \end{aligned}$$

#### (v) Quartile Deviation of the Normal Distribution $N(\mu, \sigma)$

**Definition:** The first quartile  $Q_1$  and the third quartile  $Q_3$  of  $N(\mu, \sigma)$  (or of any continuous random variable) are defined by the equations

$$\int_{-\infty}^{Q_1} f(x) dx = \frac{1}{4} \text{ and } \int_{-\infty}^{Q_3} f(x) dx = \frac{3}{4}$$

or equivalently,

$$\int_{Q_1}^{\mu} f(x) dx = \frac{1}{4} \text{ and } \int_{\mu}^{Q_3} f(x) dx = \frac{1}{4},$$

[if the curve  $y = f(x)$  is symmetrical about  $x = \mu$ ]

Then the quartile deviation (QD) is defined as

$$\text{QD} = \frac{1}{2}(Q_3 - Q_1)$$

For the normal distribution  $N(\mu, \sigma)$ ,  $Q_1$  is given by

$$\int_{Q_1}^{\mu} \frac{1}{\sigma \sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2} dx = 0.25$$

i.e.,  $\int_{(Q_1-\mu)/\sigma}^0 \frac{1}{\sqrt{2\pi}} e^{-z^2/2} dz = 0.25 \quad \left( \text{on putting } z = \frac{x-\mu}{\sigma} \right)$

i.e.,  $\int_0^{(\mu-Q_1)/\sigma} \phi(z) dz = 0.25,$

$\left( \text{by symmetry of the normal curve and since } \frac{Q_1 - \mu}{\sigma} < 0. \right)$

From the table of normal areas (areas under standard normal curve), we get

$$\int_0^{0.674} \phi(z) dz = 0.25,$$

$$\therefore \frac{\mu - Q_1}{\sigma} = 0.674$$

i.e.,  $Q_1 = \mu - 0.674 \sigma$

By symmetry,  $Q_3 = \mu + 0.674 \sigma$

$$\therefore \text{QD} = \frac{1}{2}(Q_3 - Q_1) = 0.674 \sigma = \frac{2}{3} \sigma \text{ (approximately)}$$

#### (vi) Moment Generating Function of $N(0, 1)$ and $N(\mu, \sigma)$

The moment generating function of  $N(0, 1)$  is given by

$$\begin{aligned} M_Z(t) &= E(e^{tZ}) \\ &= \int_{-\infty}^{\infty} e^{tz} \phi(z) dz, \quad [\text{where } \phi(z) \text{ is the density function of } N(0, 1)] \\ &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-z^2/2} e^{tz} dz \\ &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-(z^2 - 2tz)/2} dz \\ &= \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-(z-t)^2/2} dz \end{aligned}$$

$$\begin{aligned}
 &= \frac{1}{\sqrt{2\pi}} e^{t^2/2} \sqrt{2} \int_{-\infty}^{\infty} e^{-u^2} du, \quad \left( \text{on putting } u = \frac{z-t}{\sqrt{2}} \right) \\
 &= \frac{1}{\sqrt{\pi}} e^{t^2/2} \left[ \left( \frac{1}{2} \right) \right] = e^{t^2/2} \left( \because \left[ \left( \frac{1}{2} \right) \right] = \sqrt{\pi} \right)
 \end{aligned}$$

The moment generating function of  $N(\mu, \sigma)$  is given by

$$\begin{aligned}
 M_X(t) &= M_{\sigma Z + \mu}(t), \quad \left( \text{since } Z = \frac{X - \mu}{\sigma} \right) \\
 &= e^{\mu t} \cdot \mu_Z(\sigma t), \quad (\text{by the property of MGF}) \\
 &= e^{\mu t} e^{\sigma^2 t^2/2} \\
 &= e^{t \left( \mu + \frac{\sigma^2 t}{2} \right)}
 \end{aligned}$$

$$\begin{aligned}
 \text{Now, } M_X(t) &= 1 + \frac{t}{1!} \left( \mu + \frac{\sigma^2 t}{2} \right) + \frac{t^2}{2!} \left( \mu + \frac{\sigma^2 t}{2} \right)^2 \\
 &\quad + \frac{t^3}{3!} \left( \mu + \frac{\sigma^2 t}{2} \right)^3 + \frac{t^4}{4!} \left( \mu + \frac{\sigma^2 t}{2} \right)^4 + \dots + \infty
 \end{aligned}$$

$$\therefore E(X) = \text{Coefficient of } \frac{t}{1!} = \mu$$

$$E(X^2) = \text{Coefficient of } \frac{t^2}{2!} = \sigma^2 + \mu^2$$

$$E(X^3) = \text{Coefficient of } \frac{t^3}{3!} = 3\mu\sigma^2 + \mu^3 \text{ and}$$

$$E(X^4) = \text{Coefficient of } \frac{t^4}{4!} = 3\sigma^4 + 6\mu^2\sigma^2 + \mu^4.$$

Using the relation  $\mu_k = k\text{th order central moment} = E\{(X - \mu)^k\}$ , we get

$$\mu_1 = 0, \mu_2 = \sigma^2, \mu_3 = 0, \mu_4 = 3\sigma^4$$

We could have got these values from the formulas  $\mu_{2n+1} = 0$  and  $\mu_{2n} = 1.35 \dots (2n-1)\sigma^{2n}$ , which we have derived already.

### Corollary

If  $X$  and the distribution  $N(\mu, \sigma)$  then  $Y = aX + b$  has the distribution  $N(a\mu + b, a\sigma)$

$$\begin{aligned}
 M_X(t) &= e^{\mu t + \sigma^2 t^2/2} \\
 \therefore M_Y(t) &= M_{aX+b}(t) \\
 &= e^{bt} M_X(at) \\
 &= e^{bt} e^{a\mu t + a^2 \sigma^2 t^2/2} \\
 &= e^{(a\mu+b)t + (a^2\sigma^2)t^2/2}
 \end{aligned}$$

which is the MGF of  $N(a\mu + b, a\sigma)$ .

In particular, if  $X$  has the distribution  $N(\mu, \sigma)$ , then  $Z = \frac{X - \mu}{\sigma}$  has the distribution  $N\left(\frac{1}{\sigma}\mu - \frac{\mu}{\sigma}, \frac{1}{\sigma} \cdot \sigma\right)$  that is,  $N(0, 1)$ .

### (vii) Additive Property of the Normal Distribution

If  $X_i$  ( $i = 1, 2, \dots, n$ ) be  $n$  independent normal RVs with mean  $\mu_i$  and variance  $\sigma_i^2$ , then  $\sum_{i=1}^n a_i X_i$  is also a normal RV with mean  $\sum_{i=1}^n a_i \mu_i$  and variance  $\sum_{i=1}^n a_i^2 \sigma_i^2$ .

$$\begin{aligned} M_{\left(\sum_{i=1}^n a_i X_i\right)}(t) &= M_{a_1 X_1}(t) \cdot M_{a_2 X_2}(t) \cdots M_{a_n X_n}(t), \text{ (by independence)} \\ &= e^{a_1 \mu_1 t + a_1^2 \sigma_1^2 t^2/2} \times e^{a_2 \mu_2 t + a_2^2 \sigma_2^2 t^2/2} \cdots \times e^{a_n \mu_n t + a_n^2 \sigma_n^2 t^2/2} \\ &= e^{(\sum a_i \mu_i)t + \sum a_i^2 \sigma_i^2 t^2/2} \end{aligned}$$

which is the MGF of a normal RV with mean  $\sum a_i \mu_i$  and variance  $\sum a_i^2 \sigma_i^2$ . Hence the property.

### Deductions

- Putting  $a_1 = a_2 = 1$  and  $a_3 = a_4 = \dots = a_n = 0$ , we get the following result, in particular:

If  $X_1$  is  $N(\mu_1, \sigma_1)$  and  $X_2$  is  $N(\mu_2, \sigma_2)$ , then  $X_1 + X_2$  is  $N(\mu_1 + \mu_2, \sqrt{\sigma_1^2 + \sigma_2^2})$ .

Similarly,  $X_1 - X_2$  is  $N(\mu_1 - \mu_2, \sqrt{\sigma_1^2 + \sigma_2^2})$ .

- Putting  $a_1 = a_2 = \dots = a_n = \frac{1}{n}$  and assuming that each  $X_i$  is  $N(\mu, \sigma)$ , then

$\bar{X} = \frac{1}{n} \sum_{i=1}^n X_i$  has a normal distribution

$$N\left\{ \frac{1}{n} \sum_{i=1}^n \mu, \sqrt{\sum_{i=1}^n \frac{1}{n^2} \sigma^2} \right\}$$

i.e.,  $N\{\mu, \sigma/\sqrt{n}\}$ .

Thus, if  $X_i$  ( $i = 1, 2, \dots, n$ ) are independent and identically distributed normal variables with mean  $\mu$  and standard deviation  $\sigma$ , then their mean  $\bar{X}$  is  $N\{\mu, \sigma/\sqrt{n}\}$ .

### (viii) Normal Distribution as a Limiting Form of the Binomial Distribution

When  $n$  is very large and neither  $p$  nor  $q$  is very small, the standard normal distribution can be regarded as the limiting form of the standardised binomial distribution.

When  $X$  follows the binomial distribution  $B(n, p)$ , the standardised binomial variable  $Z$  is given by  $Z = \frac{X - np}{\sqrt{npq}}$ . As  $X$  varies from 0 to  $n$  with step size 1,

$Z$  varies from  $\frac{-np}{\sqrt{npq}}$  to  $\frac{np}{\sqrt{npq}}$  with step size  $\frac{1}{\sqrt{npq}}$ . When neither  $p$  nor  $q$

is very small and  $n$  is very large,  $Z$  varies from  $-\infty$  to  $\infty$  with infinitesimally small step size. Hence, in the limit, the distribution of  $Z$  may be expected to be a continuous distribution extending from  $-\infty$  to  $\infty$  and having mean 0 and standard deviation 1. In fact the limiting form of the distribution of  $Z$  is standard normal distribution as seen belows:

If  $X$  follows  $B(n, p)$ , then the MGF of  $X$  is given by  $M_X(t) = (q + p e^t)^n$ .

If  $Z = \frac{X - np}{\sqrt{npq}}$ , then

$$\begin{aligned}
 M_Z(t) &= M_{\frac{1}{\sqrt{npq}}X - \frac{np}{\sqrt{npq}}} (t) = e^{\frac{-np}{\sqrt{npq}}} \{q + p e^{t/\sqrt{npq}}\}^n \\
 \therefore \log M_Z(t) &= -\frac{np}{\sqrt{npq}} + n \log \{q + p e^{t/\sqrt{npq}}\} \\
 &= -\frac{np}{\sqrt{npq}} + \\
 &\quad n \log \left[ q + p \left\{ 1 + \frac{t}{\sqrt{npq}} + \frac{t^2}{2\sqrt{npq}} + \frac{t^3}{6(npq)^{3/2}} + \dots \right\} \right] \\
 &= -\frac{np}{\sqrt{npq}} + \\
 &\quad n \log \left[ 1 + \left\{ \frac{pt}{\sqrt{npq}} + \frac{pt^2}{2\sqrt{npq}} + \frac{pt^3}{6(npq)^{3/2}} + \dots \right\} \right] \\
 &= -\frac{np}{\sqrt{npq}} + n \left[ \frac{pt}{\sqrt{npq}} \left\{ 1 + \frac{t}{2\sqrt{npq}} + \frac{t^2}{6n^2 p^2 q^2} + \dots \right\} \right. \\
 &\quad \left. - \frac{1}{2} \cdot \frac{p^2 t^2}{npq} \left\{ 1 + \frac{t}{2\sqrt{npq}} + \frac{t}{6n^2 p^2 q^2} + \dots \right\}^2 + \dots \right] \\
 &= \frac{t^2}{2} + \text{terms containing } \frac{1}{\sqrt{n}} \text{ and lower powers of } n
 \end{aligned}$$

$$\therefore \lim_{n \rightarrow \infty} \log M_Z(t) = \frac{t^2}{2}$$

$$\text{i.e., } \log_e \left[ \lim_{n \rightarrow \infty} M_z(t) \right] = \frac{t^2}{2}$$

$$\therefore \lim_{n \rightarrow \infty} M_Z(t) = e^{t^2/2}$$

which is the MGF of the standard normal distribution. Hence, the limit of the standardised binomial distribution, as  $n$  tends to  $\infty$ , is the standard normal distribution.

**Note** We recall De Moivre–Laplace approximation for the sum of a large number of terms of the form  $nC_r p^r q^{n-r}$  in terms of the integral of standard normal density function, which was discussed in section 1(c). It was stated that

$$\sum_{r=r_1}^{r_2} nC_r p^r q^{n-r} = \int_{z_1}^{z_2} \phi(z) dz$$

where  $z_1 = \frac{r_1 - np - \frac{1}{2}}{\sqrt{n pq}}$  and  $z_2 = \frac{r_2 - np + \frac{1}{2}}{\sqrt{n pq}}$  and  $\phi(z)$  is the density function of the standard normal distribution.

### (c) Importance of Normal Distributions

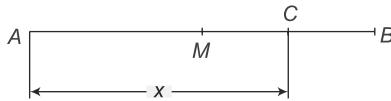
Normal distributions play a very important role in statistical theory because of the following reasons:

- (i) A large number of RVs, such as binomial and Poisson, occurring in many applications have a distribution closely resembling the normal distribution.
- (ii) Many of the distributions of sample statistics, such as sample mean and sample variance, tend to normality for samples of large size. In particular, the sampling distributions like Student's  $t$ , Snedecor's  $F$  and Chi-square distributions tend to normality when the size of the sample is large.
- (iii) Tests of significance for small samples are based on the assumption that samples have been drawn from normal populations.
- (iv) Even if a variable is not normally distributed, it can sometimes be converted into a normal variable by simple transformation of the variable.
- (v) Normal distribution is applied to a large extent in Statistical Quality Control in industry.

### Worked Example 5(B)

#### Example 1

If a string, 1 m long, is cut into 2 pieces at a random point along its length, what is the probability that the longer piece is at least twice the length of the shorter?



**Fig. 5.2**

Let  $C$  be the point of cut on  $AB$  such that  $AC = X$ . Since all positions of  $C$  are equally likely,  $X$  is uniformly distributed over  $(0, 1)$  [ $\because AB = 1$ ]

$$\therefore f(x) = 1$$

If  $X$  represents the length of the longer piece, then  $C$  lies in  $MB$ .

$$P\{\text{longer piece length} \geq 2 \times \text{shorter piece length}\}$$

$$\begin{aligned} &= P\{X \geq 2(1-X)\} = P\left(X \geq \frac{2}{3}\right) \\ &= \int_{2/3}^1 1 dx = \frac{1}{3} \end{aligned}$$

If  $X$  represents the length of the shorter piece, then  $C$  lies in  $AM$ .

In this case also, required probability

$$\begin{aligned} &= P\{(1-X) \geq 2X\} = P\left(X \leq \frac{1}{3}\right) \\ &= \int_0^{1/3} 1 dx = \frac{1}{3} \end{aligned}$$

Hence, the required probability =  $\frac{1}{3}$  [Also see Example (17) in Worked

Example 2(A) and Problem (59) in Exercise 2(A)].

#### Example 2

Buses arrive at a specified stop at 15 min. intervals starting at 7 am that is, they arrive at 7, 7:15, 7:30, 7:45, and so on. If a passenger arrives at the stop at a random time that is uniformly distributed between 7 and 7:30 am, find the probability that he waits.

- (a) less than 5 min for a bus, and
- (b) at least 12 min for a bus.

Let  $X$  denote the time in minutes past 7 am, when the passenger arrives at the stop.

Then  $X$  is uniformly distributed over  $(0, 30)$ , i.e.,  $f(x) = \frac{1}{30}$ ,  $0 < x < 30$

- (a) The passenger will have to wait less than 5 min. if he arrives at the stop between 7:10 and 7:15 or 7:25 and 7:30.

$$\therefore \text{Required probability} = P(10 < x < 15) + P(25 < x < 30)$$

$$\begin{aligned} &= \int_{10}^{15} \frac{1}{30} dx + \int_{25}^{30} \frac{1}{30} dx \\ &= \frac{1}{3} \end{aligned}$$

- (b) The passenger will have to wait at least 12 min. if he arrives at the stop between 7:00 and 7:03 or 7:15 and 7:18.

$$\therefore \text{Required probability} = P(0 < x < 3) + P(15 < x < 18)$$

$$\begin{aligned} &= \int_0^3 \frac{1}{30} dx + \int_{15}^{18} \frac{1}{30} dx \\ &= \frac{1}{5} \end{aligned}$$

### Example 3

If the roots of the quadratic equation  $x^2 - ax + b = 0$  are real and  $b$  is positive but otherwise unknown, what are the expected values of the roots of the equation. Assume that  $b$  has a uniform distribution in the permissible range.

The roots of the equation  $x^2 - ax + b = 0$  are given by

$$x = \frac{1}{2}(a \pm \sqrt{a^2 - 4b})$$

Since the roots are real,  $a^2 - 4b > 0$

$$\text{i.e., } 0 < b < \frac{a^2}{4} \quad (\because b > 0)$$

Therefore,  $b$  is a random variable, uniformly distributed in  $\left(0, \frac{a^2}{4}\right)$ .

Therefore, its density function  $f(b) = \frac{4}{a^2}$

$$\begin{aligned} E\{\text{the roots}\} &= E\left\{\frac{1}{2}(a \pm \sqrt{a^2 - 4b})\right\} \\ &= \frac{1}{2} \int_0^{a^2/4} (a \pm \sqrt{a^2 - 4b}) \cdot \frac{4}{a^2} db \end{aligned}$$

$$\begin{aligned}
 &= \frac{2}{a^2} \left[ ab \pm \frac{(a^2 - 4b)^{\frac{3}{2}}}{\frac{3}{2}(-4)} \right]_0^{a^2/4} \\
 &= \frac{2}{a^2} \left[ \frac{a^2}{4} \mp \frac{1}{6}(0 - a^3) \right] \\
 &= a \left( \frac{1}{2} \pm \frac{1}{3} \right) \\
 &= \frac{5a}{6} \text{ and } \frac{a}{6}
 \end{aligned}$$

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**Example 4**

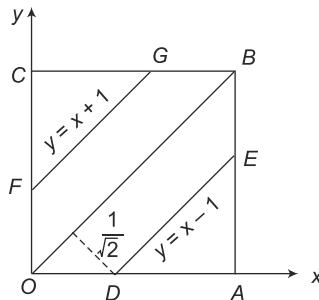
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Two points are taken at random on a given straight line of length 2 units. Prove that the probability of the distance between them exceeds 1 unit is  $\frac{1}{4}$ .

[See Example (10) in Worked Example 2(B)]

Let  $X$  and  $Y$  be the distance of the two points  $P, Q$ , taken on the line  $AB$  from  $A$ . Each of  $X$  and  $Y$  follows a uniform distribution in  $(0, 2)$ .

Therefore, the joint density function of  $(X, Y) = f(x, y) = \frac{1}{4}$ ,  $0 < x < 2; 0 < y < 2$ .



**Fig. 5.3**

$$\begin{aligned}
 \text{Now, } P\{PQ < 1\} &= P\{|X - Y| < 1\} \\
 &= P\{-1 < X - Y < 1\} \\
 &= \iint_{-1 < x - y < 1} f(x, y) dx dy \\
 &= \frac{1}{4} \iint_{-1 < x - y < 1} dx dy = \frac{1}{4}. \text{ Area of } ODEBGFO \\
 &= \frac{1}{2}. \text{ Area of trapezium } ODEB, \text{ by symmetry}
 \end{aligned}$$

$$\begin{aligned}
 &= \frac{1}{2} \cdot \frac{1}{2} \cdot \frac{1}{\sqrt{2}} (\sqrt{2} + 2\sqrt{2}) \\
 &= \frac{3}{4} \\
 \therefore P\{PQ > 1\} &= \frac{1}{4}
 \end{aligned}$$

---

### Example 5

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If  $X$  is a RV with a continuous distribution function  $F(x)$ , prove that  $Y = F(X)$  has a uniform distribution in  $(0, 1)$ . Further,

if  $f(x) = \begin{cases} \frac{1}{2}(x-1) & 1 \leq x \leq 3 \\ 0, & \text{otherwise} \end{cases}$

find the range of  $Y$  corresponding to the range  $1.1 \leq X \leq 2.9$ .

The distribution function of  $Y$  is given by

$$\begin{aligned}
 G_Y(y) &= P(Y \leq y) \\
 &= P\{F(X) \leq y\} \\
 &= P\{X \leq F^{-1}(y)\} && [\text{The inverse exists, as } F(x) \text{ is non-decreasing and continuous}] \\
 &= F[F^{-1}(y)] && [\because P\{X \leq x\} = F(x)] \\
 &= y
 \end{aligned}$$

Therefore, the density function of  $Y$  is given by

$$g_Y(y) = \frac{d}{dy}[G_Y(y)] = 1$$

Also the range of  $Y$  is  $0 \leq y \leq 1$ , since the range of  $F(x)$  is  $(0, 1)$ .

Therefore,  $Y$  follows a uniform distribution in  $(0, 1)$ .

**Note** The converse of this problem has been worked in Example (14) of Worked Example (3).

When  $f(x) = \begin{cases} \frac{1}{2}(x-1) & 1 \leq x \leq 3 \\ 0 & \text{otherwise} \end{cases}$

$$F(x) = \int_1^x \frac{1}{2}(x-1) dx = \frac{1}{4}(x-1)^2$$

Since  $Y = F(X)$ ,  $Y = \frac{1}{4}(X-1)^2$

$$\therefore \text{when } 1.1 \leq X \leq 2.9, \frac{1}{4}(1.1 - 1)^2 \leq Y \leq \frac{1}{4}(2.9 - 1)^2$$

i.e., the required range of  $Y$  is

$$0.0025 \leq Y \leq 0.9025$$

---

### Example 6

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The mileage which car owners get with a certain kind of radial tire is a RV having an exponential distribution with mean 40,000 km. Find the probabilities that one of these tires will last (a) at least 20,000 km, and (b) at most 30,000 km.

Let  $X$  denote the mileage obtained with the tire

$$f(x) = \frac{1}{40,000} e^{-x/40,000} \quad x > 0$$

$$\begin{aligned} \text{(a)} \quad P(X \geq 20,000) &= \int_{20,000}^{\infty} \frac{1}{40,000} e^{-x/40,000} dx \\ &= [-e^{-x/40,000}]_{20,000}^{\infty} \\ &= e^{-0.5} = 0.6065 \end{aligned}$$

$$\begin{aligned} \text{(b)} \quad P(X \leq 30,000) &= \int_0^{30,000} \frac{1}{40,000} e^{-x/40,000} dx \\ &= [-e^{-x/40,000}]_0^{30,000} \\ &= 1 - e^{-0.5} = 0.5270 \end{aligned}$$

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### Example 7

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If the time  $T$  to failure of a component is exponentially distributed with parameter  $\lambda$  and if  $n$  such components are installed, what is the probability that one-half or more of these components are still functioning at the end of  $t$  hours?

The density function of  $T$  is given by

$$f(t) = \lambda e^{-\lambda t}, \quad t \geq 0$$

$P$ (a component functions at the end of or after  $t$  hours)

$$= P(T \geq t) = \int_t^{\infty} \lambda e^{-\lambda t} dt = e^{-\lambda t}$$

If we consider a component functioning at the end or after  $t$  hours as a success in a single trial, we have  $p = e^{-\lambda t}$  and  $q = 1 - e^{-\lambda t}$ .

Then the number  $X$  of successes in  $n$  independent trials follows a binomial distribution with parameters  $n$  and  $p$ .

$$\therefore P(X = r) = nC_r p^r q^{n-r}; \quad r = 0, 1, 2, \dots, n$$

If  $n$  is even the required probability is given by

$$\sum_{r=\frac{n}{2}}^n P(X=r) = \sum_{r=\frac{n}{2}}^n nC_r e^{-\lambda rt} (1-e^{-\lambda t})^{n-r}$$

If  $n$  is odd, the required probability is given by

$$\sum_{r=\frac{n+1}{2}}^n P(X=r) = \sum_{r=\frac{n+1}{2}}^n nC_r e^{-\lambda rt} (1-e^{-\lambda t})^{n-r}$$

### Example 8

If a continuous RV  $X(>0)$  possesses memoryless property, that is  $P(X>x+h) = P(X>x)$ .  $P(X>h)$ , then  $X$  follows an exponential distribution.

Let  $G(x) = P(X>x)$

$\therefore$  the given condition means that

$$G(x+h) = G(x) G(h)$$

$$\begin{aligned} \therefore \frac{G(x+h) - G(x)}{h} &= G(x) \left\{ \frac{G(h) - 1}{h} \right\} \\ &= G(x) \frac{[G(h) - G(0)]}{h} [\because G(0) = P(X>0) = 1, \text{ as } x>0] \end{aligned}$$

Taking limits on both sides as  $h \rightarrow 0$ , we have

$$\begin{aligned} G'(x) &= G(x) \cdot G'(0) \\ &= -\lambda \cdot G(x), \quad [\text{on putting } \lambda = -G'(0)] \end{aligned} \tag{1}$$

Solving the differential equation (i), we get

$$\log G(x) = -\lambda x + \log C$$

$$\text{i.e., } G(x) = C e^{-\lambda x} \tag{2}$$

Using the fact that  $G(0) = 1$  in (2), we get  $C = 1$

$$\text{Thus, } G(x) = P(X>x) = e^{-\lambda x}$$

Now, the distribution function  $F(x)$  of  $X$  is given by  $F(x) = P(X \leq x)$

$$\begin{aligned} &= 1 - P(X>x) = \{1 - G(x)\} \\ &= 1 - e^{-\lambda x} \end{aligned}$$

Therefore, the density function  $f(x)$  of  $X$  is given by

$$f(x) = F'(x) = \lambda e^{-\lambda x}, x>0$$

i.e.,  $X$  follows an exponential distribution with parameter

$$\lambda = -G'(0) = F'(0) > 0$$

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### Example 9

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The time (in hours) required to repair a machine is exponentially distributed with parameter  $\lambda = 1/2$ .

- What is the probability that the repair time exceeds 2 h?
- What is the conditional probability that a repair takes at least 10 h given that its duration exceeds 9 h?

If  $X$  represents the time to repair the machine, the density function of  $X$  is given by

$$f(x) = \lambda e^{-\lambda x} = \frac{1}{2} e^{-\frac{x}{2}}, x > 0$$

$$\begin{aligned}(a) \quad P(X > 2) &= \int_2^{\infty} \frac{1}{2} e^{-\frac{x}{2}} dx \\ &= \left(-e^{-\frac{x}{2}}\right)_2^{\infty} = e^{-1} = 0.3679\end{aligned}$$

(b)  $P\{X \geq 10 | X > 9\} = P\{X > 1\}$ , (by the memoryless property)

$$\begin{aligned}&= \int_1^{\infty} \frac{1}{2} e^{-\frac{x}{2}} dx \\ &= \left(-e^{-\frac{x}{2}}\right)_0^{\infty} = e^{-0.5} = 0.6065\end{aligned}$$

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### Example 10

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The life length  $X$  of an electronic component follows an exponential distribution. There are 2 processes by which the component may be manufactured. The expected life length of the component is 100 h. if the process I is used to manufacture, while it is 150 h if the process II is used. The cost of manufacturing a single component by process I is Rs. 10, while it is Rs. 20 for process II. Moreover if the component lasts less than the guaranteed life of 200 h, a loss of Rs. 50 is to be borne by the manufacturer. Which process is advantageous to the manufacturer?

If the process I is used, the density function of  $X$  is given by

$$f(x) = \frac{1}{100} e^{-x/100}, x > 0.$$

$$\begin{aligned}\therefore \quad P(X \geq 200) &= \int_{200}^{\infty} \frac{1}{100} e^{-x/100} dx \\ &= (e^{-x/100})_{200}^{\infty} = e^{-2}\end{aligned}$$

$$\therefore \quad P(X < 200) = 1 - e^{-2}$$

Similarly, if the process II is used,

$$P(X \geq 200) = e^{-4/3} \text{ and } P(X < 200) = 1 - e^{-4/3}$$

Let  $C_1$  and  $C_2$  be the costs per component corresponding to the processes I and II respectively.

Then 
$$C_1 = \begin{cases} 10, & X \geq 200 \\ 60, & X < 200 \end{cases}$$

$$\begin{aligned}\therefore E(C_1) &= 10 \times P(X \geq 200) + 60 \cdot P(X < 200) \\ &= 10 e^{-2} + 60 (1 - e^{-2}) \\ &= 60 - 50 e^{-2} = 53.235\end{aligned}$$

Now 
$$C_2 = \begin{cases} 20, & X \geq 200 \\ 70, & X < 200 \end{cases}$$

$$\begin{aligned}\therefore E(C_2) &= 20 \times P(X \geq 200) + 70 \times P(X < 200) \\ &= 20 e^{-4/3} + 70 (1 - e^{-4/3}) \\ &= 70 - 50 e^{-4/3} = 56.765\end{aligned}$$

Since  $E(C_1) < E(C_2)$ , process I is advantageous to the manufacturer.

### Example 11

If the density function of a continuous RV  $X$  is  $f(x) = c e^{-b(x-a)}$ ,  $a \leq x$ , where  $a$ ,  $b$ ,  $c$  are constants. Show that  $b = c = \frac{1}{\sigma}$  and  $a = \mu - \sigma$ , where  $\mu = E(X)$  and  $\sigma^2 = \text{Var}(X)$ .

Since  $f(x)$  is a density function,  $\int_a^{\infty} f(x) dx = 1$ .

i.e.,  $\int_a^{\infty} c e^{-b(x-a)} dx = 1$

i.e.,  $c \left\{ \frac{e^{-b(x-a)}}{-b} \right\}_a^{\infty} = 1$

i.e.,  $\frac{c}{b} = 1 \text{ or } b = c$  (1)

Now,  $\mu = E(X) = \int_a^{\infty} bx e^{-b(x-a)} dx$

$$= b e^{ab} \left[ x \cdot \left( \frac{e^{-bx}}{-b} \right) - \frac{e^{-bx}}{b^2} \right]_a^{\infty}$$

$$\begin{aligned}
&= b e^{ab} \left[ \frac{a}{b} e^{-ab} + \frac{1}{b^2} e^{-ab} \right] \\
&= a + \frac{1}{b}
\end{aligned} \tag{2}$$

$$\begin{aligned}
E(X^2) &= \int_a^\infty b x^2 e^{-b(x-a)} dx \\
&= b e^{ab} \left[ x^2 \left( \frac{e^{-bx}}{-b} \right) - 2x \left( \frac{e^{-bx}}{-b^2} \right) + 2 \left( \frac{e^{-bx}}{-b^3} \right) \right]_a^\infty \\
&= b \left[ \frac{a^2}{b} + \frac{2a}{b^2} + \frac{2}{b^3} \right] \\
&= \frac{1}{b^2} (a^2 b^2 + 2ab + 2)
\end{aligned}$$

$$\begin{aligned}
\text{Var}(X) &= E(X^2) - \{E(X)\}^2 \\
&= \frac{1}{b^2} (a^2 b^2 + 2ab + 2) - \left( a^2 + \frac{2a}{b} + \frac{1}{b^2} \right)
\end{aligned}$$

i.e.,  $\sigma^2 = \frac{1}{b^2}$  or  $\sigma = \frac{1}{b}$  (3)

From (1) and (3), we get  $b = c = \frac{1}{\sigma}$ .

From (2) and (3),  $\mu - \sigma = a$ .

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### Example 12

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In a certain city, the daily consumption of electric power in millions of kilowatt-hours can be treated as a RV having an Erlang distribution with parameters  $\lambda = \frac{1}{2}$  and  $k = 3$ . If the power plant of this city has a daily capacity of 12 millions kilowatt-hours, what is the probability that this power supply will be inadequate on any given day.

Let  $X$  represent the daily consumption of electric power (in millions of kilowatt-hours). Then the density function of  $X$  is given as

$$f(x) = \frac{\left(\frac{1}{2}\right)^3}{\Gamma(3)} x^2 e^{-x/2}, x > 0$$

$P\{\text{the power supply is inadequate}\}$

$$\begin{aligned}
 &= P(X > 12) = \int_{12}^{\infty} f(x) dx \quad [\because \text{the daily capacity is only } 12] \\
 &= \int_{12}^{\infty} \frac{1}{(3)} \cdot \frac{1}{2^3} x^2 e^{-x/2} dx \\
 &= \frac{1}{16} \left[ x^2 \left( \frac{e^{-x/2}}{-\frac{1}{2}} \right) - 2x \left( \frac{e^{-x/2}}{\frac{1}{4}} \right) + 2 \left( \frac{e^{-x/2}}{-\frac{1}{8}} \right) \right]_{12}^{\infty} \\
 &= \frac{1}{16} e^{-6} (288 + 96 + 16) \\
 &= 25 e^{-6} = 0.0625
 \end{aligned}$$

---

### Example 13

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If a company employs  $n$  sales persons, its gross sales in thousands of rupees may be regarded as a RV having an Erlang distribution with  $\lambda = \frac{1}{2}$  and  $k = 80\sqrt{n}$ . If the sales cost is Rs. 8000 per salesperson, how many salespersons should the company employ to maximise the expected profit?

Let  $X$  represent the gross sales (in Rupees) by  $n$  salespersons.

$X$  follows the Erlang distribution with parameters  $\lambda = \frac{1}{2}$  and  $k = 80,000\sqrt{n}$ .

$$\therefore E(X) = \frac{k}{\lambda} = 1,60,000\sqrt{n}$$

If  $y$  denotes the total expected profit of the company, then

$$y = \text{total expected sales} - \text{total sales cost}$$

$$= 1,60,000\sqrt{n} - 8000n$$

$$\frac{dy}{dn} = \frac{80,000}{\sqrt{n}} - 8000$$

$$= 0, \text{ when } \sqrt{n} = 10 \text{ or } n = 100$$

$$\frac{d^2y}{dn^2} = -\frac{40,000}{n^{3/2}} < 0, \text{ when } n = 100.$$

Therefore,  $y$  is maximum, when  $n = 100$ .

That is the company should employ 100 salespersons in order to maximise the total expected profit.

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**Example 14**

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Consumer demand for milk in a certain locality, per month, is known to be a general Gamma (Erlang) RV. If the average demand is  $a$  litres and the most likely demand is  $b$  litres ( $b < a$ ), what is the variance of the demand?

Let  $X$  represent the monthly consumer demand of milk.

Average demand is the value of  $E(X)$ .

Most likely demand is the value of the mode of  $X$  or the value of  $X$  for which its density function is maximum.

If  $f(x)$  is the density function of  $X$ , then

$$\begin{aligned}f(x) &= \frac{\lambda^k}{(k)} x^{k-1} e^{-\lambda x} \quad x > 0 \\f'(x) &= \frac{\lambda^k}{(k)} [(k-1)x^{k-2} e^{-\lambda x} - \lambda x^{k-1} e^{-\lambda x}] \\&= \frac{\lambda^k}{(k)} x^{k-2} e^{-\lambda x} \{(k-1) - \lambda x\} \\&= 0, \text{ when } x = 0, \quad x = \frac{k-1}{\lambda} \\f''(x) &= \frac{\lambda^k}{(k)} [-\lambda x^{k-2} e^{-\lambda x} + \{(k-1) - \lambda x\} \frac{d}{dx} \{x^{k-2} e^{-\lambda x}\}] \\&< 0, \text{ when } x = \frac{k-1}{\lambda}\end{aligned}$$

Therefore,  $f(x)$  is maximum, when  $x = \frac{k-1}{\lambda}$ .

$$\text{i.e., most likely demand} = \frac{k-1}{\lambda} = b \quad (1)$$

$$\text{and} \quad E(X) = \frac{k}{\lambda} = a \quad (2)$$

$$\begin{aligned}\text{Now,} \quad \text{Var}(X) &= \frac{k}{\lambda^2} = \frac{k}{\lambda} \cdot \frac{1}{\lambda} \\&= a(a-b), \quad [\text{from (1) and (2)}]\end{aligned}$$

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**Example 15**

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A random sample of size  $n$  is taken from a general Gamma (Erlang) distribution with parameters  $\lambda$  and  $k$ . Show that the mean  $\bar{X}$  of the sample also follows a Gamma distribution with parameters  $n\lambda$  and  $nk$ .

If  $X$  follows Erlang distribution with parameters  $\lambda$  and  $k$ , then the MGF of  $X$  is given by

$$M_X(t) = \left(1 - \frac{t}{\lambda}\right)^{-k}$$

If  $X_1, X_2, \dots, X_n$  are the members of the sample drawn, then each  $X_i$  follows Erlang distribution with MGF equal to  $\left(1 - \frac{t}{\lambda}\right)^{-k}$  and also they are independent.

Therefore, by the reproductive property,

$$M_{X_1 + X_2 + \dots + X_n}(t) = \left(1 - \frac{t}{\lambda}\right)^{-nk}$$

$$\begin{aligned} \therefore M_{\bar{X}}(t) &= M_{\frac{1}{n}(X_1 + X_2 + \dots + X_n)(t)} \\ &= M_{X_1 + X_2 + \dots + X_n}\left(\frac{t}{n}\right) && [\because M_{aX}(t) = M_X(at)] \\ &= \left(1 - \frac{t}{n\lambda}\right)^{-nk} \end{aligned}$$

which is the MGF of an Erlang distribution with parameters  $n\lambda$  and  $nk$ . Therefore,  $\bar{X}$  also follows an Erlang distribution with density function

$$\frac{(n\lambda)^{nk}}{(nk)!} \cdot x^{nk-1} \cdot e^{-n\lambda x}, x > 0.$$

---

### Example 16

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If the conditional distribution of  $Y$ , given  $X = x$ , is an exponential distribution with parameter  $x$  and if the unconditional distribution of  $X$  is an Erlang distribution with parameters  $\lambda (> 0)$  and  $k (> 2)$ , prove that the conditional distribution of  $X$ , given  $Y = y$ , is an Erlang distribution with parameters  $\lambda + y$  and  $k + 1$ .

**Given:**  $F_{Y|X}(y) = x e^{-xy}$ ,  $y > 0$  and  $x > 0$

and  $f_X(x) = \frac{\lambda^k}{(k)!} x^{k-1} e^{-\lambda x}$ ,  $x > 0$

If  $f(x, y)$  denotes the joint density function of  $(X, Y)$ , then  $f_{Y|X}(y) = \frac{f(x, y)}{f_X(x)}$

$$\therefore f(x, y) = \frac{\lambda^k}{(k)!} x^k e^{-(\lambda+y)x}, x > 0, y > 0$$

Now,  $f_Y(y)$  = the marginal density function of  $Y$

$$\begin{aligned}
 &= \int_0^{\infty} f(x, y) dx \\
 &= \frac{\lambda^k}{\Gamma(k)} \cdot \int_0^{\infty} x^k e^{-(\lambda+y)x} dx \\
 &= \frac{\lambda^k}{\Gamma(k)} \frac{1}{(\lambda+y)^{k+1}} \int_0^{\infty} t^k e^{-t} dt \quad [\text{on putting } (\lambda+y)x=t] \\
 &= \frac{\lambda^k}{\Gamma(k)} \frac{1}{(\lambda+y)^{k+1}} \Gamma(k+1) = \frac{k \lambda^k}{(\lambda+y)^{k+1}}, y > 0
 \end{aligned}$$

Now,  $f_{X/Y}(x) = \frac{f(x, y)}{f_Y(y)}$

$$\begin{aligned}
 &= \frac{\frac{\lambda^k}{\Gamma(k)} x^k e^{-(\lambda+y)x}}{\frac{k \lambda^k}{(\lambda+y)^{k+1}}} , x > 0 \text{ and } y > 0 \\
 &= \frac{(\lambda+y)^{k+1}}{\Gamma(k+1)} \cdot x^k e^{-(\lambda+y)x}, x > 0
 \end{aligned}$$

This is the density function of an Erlang distribution with parameters  $\lambda + y$  and  $k + 1$ .

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### Example 17

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Each of the 6 tubes of a radio set has a life length (in years) which may be considered as a RV that follows a Weibull distribution with parameters  $\alpha = 25$  and  $\beta = 2$ . If the tubes function independently of one another, what is the probability that no tube will have to be replaced during the first 2 months of service?

If  $X$  represents the life length of each tube, then its density function ( $f(x)$ ) is given by

$$f(x) = \alpha \beta x^{\beta-1} e^{-\alpha x \beta} \quad x > 0$$

i.e.,  $f(x) = 50x e^{-25x} > 0$

Now,  $P(\text{a tube is not to be replaced during the first 2 months})$

$$= P\left(X > \frac{1}{6}\right)$$

$$\begin{aligned}
 &= \int_{\frac{1}{6}}^{\infty} 50x e^{-25x^2} dx \\
 &= \left(-e^{-25x^2}\right)_{\frac{1}{6}}^{\infty} = e^{-25/36}
 \end{aligned}$$

$$\begin{aligned}
 \therefore P(\text{all the 6 tubes are not to be replaced during the first 2 months}) \\
 &= (e^{-25/36})^6 \quad (\text{by independence}) \\
 &= e^{-25/6} \\
 &= 0.0155
 \end{aligned}$$

### Example 18

If the life  $X$  (in years) of a certain type of car has a Weibull distribution with the parameter  $\beta = 2$ , find the value of the parameter  $\alpha$ , given that probability that the life of the car exceeds 5 years is  $e^{-0.25}$ . For these values of  $\alpha$  and  $\beta$ , find the mean and variance of  $X$ .

The density function of  $X$  is given by

$$f(x) = 2\alpha x e^{-\alpha x^2}, x > 0 \quad [\because \beta = 2]$$

$$\begin{aligned}
 \text{Now, } P(X > 5) &= \int_5^{\infty} 2\alpha x e^{-\alpha x^2} dx \\
 &= \left(-e^{-\alpha x^2}\right)_5^{\infty} \\
 &= e^{-25\alpha}
 \end{aligned}$$

$$\text{Given that } P(X > 5) = e^{-0.25}$$

$$\therefore e^{-25\alpha} = e^{-0.25}$$

$$\therefore \alpha = \frac{1}{100}$$

For the Weibull distribution with parameters  $\alpha$  and  $\beta$ ,  $E(X) = \alpha^{-1/\beta} \sqrt{\left(\frac{1}{\beta} + 1\right)}$

$$\begin{aligned}
 \therefore \text{Required mean} &= \left(\frac{1}{100}\right)^{-\frac{1}{2}} \cdot \sqrt{\left(\frac{3}{2}\right)} \\
 &= 10 \times \frac{1}{2} \sqrt{\left(\frac{1}{2}\right)} \\
 &= 5 \sqrt{\pi}.
 \end{aligned}$$

$$\text{Var}(X) = \alpha^{-\frac{2}{\beta}} \left[ \left( \frac{2}{\beta} + 1 \right) - \left\{ \left( \frac{1}{\beta} + 1 \right) \right\}^2 \right]$$

$$\begin{aligned}&= \left( \frac{1}{100} \right)^{-1} \left[ \overline{(2)} - \left\{ \left( \frac{3}{2} \right)^2 \right\} \right] \\&= 100 \left[ 1 - \left( \frac{1}{2} \sqrt{\pi} \right)^2 \right] \\&= 100 \left( 1 - \frac{\pi}{4} \right)\end{aligned}$$

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### Example 19

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If the time  $T$  to failure of a component follows a Weibull distribution with parameters  $\alpha$  and  $\beta$ , find the hazard rate or conditional failure rate at time  $t$  of the component.

Refer to Example (19) in Worked Example 2(A).

If  $f(t)$  is the density function of  $T$  and  $h(t)$  is the hazard rate at time  $t$ , then

$$h(t) = \frac{f(t)}{1 - F(t)}$$

where  $F(t)$  is the distribution function of  $T$ .

Now, 
$$f(t) = \alpha\beta \times t^{\beta-1} e^{-\alpha t^\beta} \quad t > 0$$

$$\therefore F(t) = P(T \leq t)$$

$$\begin{aligned}&= \int_0^t \alpha\beta t^{\beta-1} \cdot e^{-\alpha t^\beta} dt \\&= \left[ -e^{-\alpha t^\beta} \right]_0^t \\&= 1 - e^{-\alpha t^\beta}\end{aligned}$$

$$\begin{aligned}\therefore h(t) &= \frac{\alpha\beta t^{\beta-1} e^{-\alpha t^\beta}}{e^{-\alpha t^\beta}} \\&= \alpha\beta t^{\beta-1}\end{aligned}$$

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### Example 20

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If  $Y$  is the smallest item of 3 independent observations  $X_1, X_2, X_3$  from a Weibull distribution with parameters  $\alpha$  and  $\beta$ , show that  $Y$  also has a Weibull distribution. What are its parameters?

Each of  $X_1, X_2, X_3$  follows the Weibull distribution whose density function is given by

$$f(x) = \alpha\beta \times x^{\beta-1} e^{\alpha x^\beta}, x > 0$$

$$\begin{aligned}
 \text{Now, } P(Y > y) &= P[\min(X_1, X_2, X_3) > y] \\
 &= P(X_1 > y) \times P(X_2 > y) \times P(X_3 > y), \\
 &\quad (\text{since } X_1, X_2, X_3 \text{ are independent}) \\
 &= \{P(X_i > y)\}^3
 \end{aligned} \tag{1}$$

$$\begin{aligned}
 \text{Now, } P(X_i > y) &= \int_y^{\infty} \alpha \beta x^{\beta-1} e^{-\alpha x^{\beta}} dx \\
 &= (-e^{-\alpha x^{\beta}})_y^{\infty} \\
 &= e^{-\alpha y^{\beta}}
 \end{aligned} \tag{2}$$

Using (2) in (1), we have

$$P(Y > y) = (e^{-\alpha y^{\beta}})^3 = e^{3\alpha y^{\beta}}$$

Therefore,  $Y$  also has Weibull distribution with parameters  $3\alpha$  and  $\beta$ .

**Note**

The result can be extended to  $n$  independent observations.

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### Example 21

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There are 400 students in the first-year class of an engineering college. The probability that any student requires a copy of a particular Mathematics book from the college library on any day is 0.1. How many copies of the book should be kept in the library so that the probability may be greater than 0.95 that none of the students requiring a copy from the library has to come back disappointed? (Use normal approximation to the binomial distribution).

$$p = P(\text{a student requires the book}) = 0.1 \text{ and } q = 0.9$$

$$n = \text{number of students} = 400$$

If  $X$  represents the number of students requiring the book, then  $X$  follows a binomial distribution with mean  $= np = 40$  and  $\text{SD} = \sqrt{npq} = 6$ .

As given in the problem, we may assume that  $X$  follows the distribution  $N(40, 6)$ .

Let  $m$  be the required number of books, satisfying the given condition.

$$\text{i.e., } P(X < m) > 0.95$$

$$\text{i.e., } P\left(-\infty < \frac{X - 40}{6} < \frac{m - 40}{6}\right) > 0.95$$

$$\text{i.e., } P\left(0 < Z < \frac{m - 40}{6}\right) > 0.45$$

where  $Z$  is the standard normal variate.

From the table of areas under normal curve, we find that

$$P\{0 < Z < 1.65\} > 0.45$$

$$\therefore \frac{m - 40}{6} = 1.65$$

i.e.,  $m = 49.9$

Therefore, at least 50 copies of the book should be kept in the library.

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### Example 22

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The marks obtained by a number of students in a certain subject are approximately normally distributed with mean 65 and standard deviation 5. If 3 students are selected at random from this group, what is the probability that at least 1 of them would have scored above 75?

If  $X$  represents the marks obtained by the students.  $X$  follows the distribution  $N(65, 5)$ .

$$P(\text{a student scores above } 75)$$

$$\begin{aligned} &= P(X > 75) = P\left(\frac{75 - 65}{5} < \frac{X - 65}{5} < \infty\right) \\ &= P(2 < Z < \infty), (\text{where } Z \text{ is the standard normal variate}) \\ &= 0.5 - P(0 < Z < 2) \\ &= 0.5 - 0.4772, (\text{from the table of areas}) \\ &= 0.0228 \end{aligned}$$

Let  $p = P(\text{a student scores above } 75) = 0.0228$  then  $q = 0.9772$  and  $n = 3$ .

Since  $p$  is the same for all the students, the number  $Y$ , of (successes) students scoring above 75, follows a binomial distribution.

$$P(\text{at least 1 student scores above } 75)$$

$$\begin{aligned} &= P(\text{at least 1 success}) \\ &= P(Y \geq 1) = 1 - P(Y = 0) \\ &= 1 - nC_0 \times p^0 q^n \\ &= 1 - 3C_0 (0.9772)^3 \\ &= 1 - 0.9333 \\ &= 0.0667 \end{aligned}$$

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### Example 23

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If the actual amount of instant coffee which a filling machine puts into '6-ounce' jars is a RV having a normal distribution with  $SD = 0.05$  ounce and if only 3% of the jars are to contain less than 6 ounces of coffee, what must be the mean fill of these jars?

Let  $X$  be the actual amount of coffee put into the jars.

Then  $X$  follows  $N(\mu, 0.05)$

**Given:**  $P(X < 6) = 0.03$

$$\therefore P\left(-\infty < \frac{X - \mu}{0.05} < \frac{6 - \mu}{0.05}\right) = 0.03$$

$$\text{i.e., } P\left\{-\infty < Z < \frac{6 - \mu}{0.05}\right\} = 0.03$$

$$\therefore P\left\{0 < Z < \frac{\mu - 6}{0.05}\right\} = 0.47, \quad (\text{by symmetry})$$

From the table of areas, we have

$$P\{0 < Z < 1.808\} = 0.47$$

$$\therefore \frac{\mu - 6}{0.05} = 1.808$$

$$\therefore \mu = 6.094 \text{ ounces.}$$

### Example 24

In an engineering examination, a student is considered to have failed, secured second class, first class and distinction, according as he scores less than 45%, between 45% and 60%, between 60% and 75% and above 75% respectively. In a particular year 10% of the students failed in the examination and 5% of the students got distinction. Find the percentages of students who have got first class and second class. (Assume normal distribution of marks).

Let  $X$  represent the percentage of marks scored by the students in the examination.

Let  $X$  follow the distribution  $N(\mu, \sigma)$ .

**Given:**  $P(X < 45) = 0.10$  and  $P(X > 75) = 0.05$

$$\text{i.e., } P\left(-\infty < \frac{X - \mu}{\sigma} < \frac{45 - \mu}{\sigma}\right) = 0.10 \text{ and}$$

$$P\left(\frac{75 - \mu}{\sigma} < \frac{X - \mu}{\sigma} < \infty\right) = 0.05$$

$$\text{i.e., } P\left(-\infty < Z < \frac{45 - \mu}{\sigma}\right) = 0.10 \text{ and}$$

$$P\left(\frac{75 - \mu}{\sigma} < Z < \infty\right) = 0.05$$

$$\therefore P\left(0 < Z < \frac{\mu - 45}{\sigma}\right) = 0.40 \text{ and}$$

$$P\left(0 < Z < \frac{75 - \mu}{\sigma}\right) = 0.45$$

From the table of areas, we get

$$\frac{\mu - 45}{\sigma} = 1.28 \text{ and } \frac{75 - \mu}{\sigma} = 1.64 \quad (1)$$

$$\text{i.e., } \mu - 1.28 \sigma = 45 \quad (1)$$

$$\text{and } \mu + 1.64 \sigma = 75 \quad (2)$$

Solving equations (1) and (2), we get

$$\mu = 58.15 \text{ and } \sigma = 10.28$$

Now,  $P$  (a student gets first class)

$$\begin{aligned} &= P(60 < X < 75) \\ &= P\left\{\frac{60 - 58.15}{10.28} < Z < \frac{75 - 58.15}{10.28}\right\} \\ &= P\{0.18 < Z < 1.64\} \\ &= P\{0 < Z < 1.64\} - P\{0 < Z < 0.18\} \\ &= 0.4495 - 0.0714 = 0.3781 \end{aligned}$$

$\therefore$  Percentage of students getting first class = 38 (approximately)

Now, percentage of students getting second class

$$\begin{aligned} &= 100 - (\text{sum of the percentages of students who have failed,} \\ &\quad \text{got first class and got distinction}) \\ &= 100 - (10 + 38 + 5), \text{ approximately.} \\ &= 47 \text{ (approximately)} \end{aligned}$$

### Example 25

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The percentage  $X$  of a particular compound contained in a rocket fuel follows the distribution  $N(33, 3)$ , though the specification for  $X$  is that it should lie between 30 and 35. The manufacturer will get a net profit (per unit of the fuel) of Rs. 100, if  $30 < X < 35$ , Rs. 50, if  $25 < X \leq 30$  or  $35 \leq X < 40$  and incur a loss of Rs. 60 per unit of the fuel otherwise. Find the expected profit of the manufacturer. If he wants to increase his expected profit by 50% by increasing the net profit on that category of the fuel that meets the specification, what should be the new net profit per unit of the fuel of this category?

$$\begin{aligned} P(30 < X < 35) &= P\left\{\frac{30 - 33}{3} < \frac{X - 33}{3} < \frac{35 - 33}{3}\right\} \\ &= P\{-1 < Z < 0.67\} \\ &= P\{0 < Z < 1\} + P\{0 < Z < 0.67\} \\ &= 0.3413 + 0.2486, \quad (\text{using the table of areas}) \\ &= 0.5899 \end{aligned}$$

$$\begin{aligned} P(25 < X \leq 30) &= P\left\{\frac{25 - 33}{3} < Z < \frac{30 - 33}{3}\right\} \\ &= P(-2.67 < Z < -1) \\ &= P(1 < Z < 2.67), \quad (\text{by symmetry}) \\ &= P(0 < Z < 2.67) - P(0 < Z < 1) \\ &= 0.4962 - 0.3413 \\ &= 0.1549 \end{aligned}$$

$$P(35 \leq X < 40) = P\left\{\frac{35 - 33}{3} < Z < \frac{40 - 33}{3}\right\}$$

$$\begin{aligned}
 &= P\{0.67 < Z < 2.33\} \\
 &= P\{0 < Z < 2.33\} - P\{0 < Z < 0.67\} \\
 &= 0.4901 - 0.2486 \\
 &= 0.2415 \\
 \therefore P\{(25 < X \leq 30) \text{ or } (35 \leq X < 40)\} &= P(25 < X \leq 30) + P(35 \leq X < 40), \\
 &\quad (\text{since the 2 ranges are mutually exclusive}) \\
 &= 0.1549 + 0.2415 = 0.3964
 \end{aligned}$$

$$\begin{aligned}
 P\{X < 25 \text{ or } X > 40\} &= 1 - (0.5899 + 0.3964) \\
 &= 0.0137
 \end{aligned}$$

<i>Profit/unit</i>	<i>Probability</i>
Rs. 100	0.5899
Rs. 50	0.3964
Rs. -60	0.0137

$$\begin{aligned}
 \therefore E(\text{Profit per unit}) &= \text{Rs.}(100 \times 0.5899 + 50 \times 0.3964 - 60 \times 0.0137) \\
 &= \text{Rs.} 79, \text{ approximately}
 \end{aligned}$$

Let the revised net profit per unit of the first category fuel be  $k$ .

The  $E(\text{Revised profit per unit})$

$$\begin{aligned}
 &= \text{Rs.}(k \times 0.5899 + 50 \times 0.3964 - 60 \times 0.0137) \\
 &= \text{Rs.}(0.5899 k + 18.998)
 \end{aligned}$$

$E(\text{Revised profit per unit}) = \text{Rs.} 79 + \text{Rs.} 39.5$ , as per the manufacturer's wish

$$\therefore 0.5899 k + 18.998 = 118.5$$

$$\begin{aligned}
 k &= \frac{118.5 - 18.998}{0.5899} \\
 &= 168.68 \approx \text{Rs.} 169 \text{ nearly.}
 \end{aligned}$$

### Example 26

The marks obtained by the students in Mathematics, Physics and Chemistry in an examination are normally distributed with the means 52, 50 and 48 and with standard deviations 10, 8 and 6 respectively. Find the probability that a student selected at random has secured a total of (a) 180 or above, and (b) 135 or less.

Let  $X$ ,  $Y$ ,  $Z$  denote the marks obtained by students in Mathematics, Physics and Chemistry respectively.

**Given:**  $X$  flows  $N(52, 10)$ ,  $Y$  follows  $N(50, 8)$  and  $Z$  follows  $N(48, 6)$ .

By the additive property of normal distribution,  $T = X + Y + Z$  follows the distribution

$$\begin{aligned}
 &N\{52 + 50 + 48, \sqrt{10^2 + 8^2 + 6^2}\} \\
 \text{i.e.,} \quad &N(150, 14.14)
 \end{aligned}$$

$$\begin{aligned}
 \text{(a)} \quad P(T \geq 180) &= P\left\{\frac{180 - 150}{14.14} < \frac{T - 150}{14.14} < \infty\right\} \\
 &= P\{2.12 < Z < \infty\}
 \end{aligned}$$

$$\begin{aligned}
 &= 0.5 - P\{0 < Z < 2.12\} \\
 &= 0.5 - 0.4830 && \text{(from the table of areas)} \\
 &= 0.0170
 \end{aligned}$$

$$\begin{aligned}
 \text{(b)} \quad P(T \leq 135) &= P\left\{\frac{T - 150}{14.14} < \frac{135 - 150}{14.14}\right\} \\
 &= P(-\infty < Z < -1.06) \\
 &= P\{1.06 < Z < \infty\}, && \text{(by symmetry)} \\
 &= 0.5 - P\{0 < Z < 1.06\} \\
 &= 0.5 - 0.3554 \\
 &= 0.1446
 \end{aligned}$$

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### Example 27

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The independent RVs  $X$  and  $Y$  have distributions  $N(45, 2)$  and  $N(44, 1.5)$  respectively. What is the probability that randomly chosen values of  $X$  and  $Y$  differ by 1.5 or more?

$X$  is  $N(45, 2)$  and  $Y$  is  $N(44, 1.5)$

$\therefore$  by the additive property,

$$U = X - Y \text{ follows the distribution } N(1, \sqrt{4 + 2.25}) \\ \text{i.e., } N(1, 2.5)$$

Now  $P\{X \text{ and } Y \text{ differ by 1.5 or more}\}$

$$\begin{aligned}
 &= P\{|X - Y| \geq 1.5\} \\
 &= P\{|U| \geq 1.5\} \\
 &= 1 - P(|U| \leq 1.5) \\
 &= 1 - P\{-1.5 \leq U \leq 1.5\} \\
 &= 1 - P\left\{\frac{-1.5 - 1}{2.5} \leq \frac{U - 1}{2.5} \leq \frac{1.5 - 1}{2.5}\right\} \\
 &= 1 - P\{-1 \leq Z \leq 0.2\} \\
 &= 1 - [P(0 \leq Z \leq 1) + P(0 \leq Z \leq 0.2)] \\
 &= 1 - (0.3413 + 0.0793), \text{ (from the table of areas)} \\
 &= 0.5794
 \end{aligned}$$

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### Example 28

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If  $X$  and  $Y$  are independent RVs, each following  $N(0, 3)$ , what is the probability that the point  $(X, Y)$  lies between the lines  $3X + 4Y = 5$  and  $3X + 4Y = 10$ ?

$X$  follows  $N(0, 3)$  and  $Y$  follows  $N(0, 3)$ .

Therefore, by the additive property of normal distribution,

$$U = 3X + 4Y \text{ follows } N[3 \times 0 + 4 \times 0, \sqrt{9 \times 9 + 16 \times 9}]$$

i.e.,  $N(0, 15)$

Now,  $P\{\text{the point } (X, Y) \text{ lies between the lines } 3X + 4Y = 5 \text{ and } 3X + 4Y = 10\}$

$$\begin{aligned} &= P\{5 < 3X + 4Y < 10\} \\ &= P\{5 < U < 10\} \\ &= P\left\{\frac{5-0}{15} < \frac{U-0}{15} < \frac{10-0}{15}\right\} \\ &= P\{0.33 < Z < 0.67\}, \end{aligned}$$

where  $Z$  is the standard normal variable

$$\begin{aligned} &= P(0 < Z < 0.67) - P(0 < Z < 0.33) \\ &= 0.2486 - 0.1293, \text{ (from the table of areas)} \\ &= 0.1193. \end{aligned}$$

### Example 29

If  $X$  and  $Y$  are independent RVs following  $N(8, 2)$  and  $N(12, 4\sqrt{3})$  respectively, find the value of  $\lambda$  such that

$$P(2X - Y \leq 2\lambda) = P(X + 2Y \geq \lambda)$$

By the additive property of normal distribution,

$$U = 2X - Y \text{ follows } N\{2 \times 8 - 12, \sqrt{4 \times 4 + 1 \times 48}\}$$

i.e.,  $N(4, 8)$

and

$$V = X + 2Y \text{ follows } N\{8 \times 2 - 12, \sqrt{4 \times 4 \times 48}\}$$

i.e.,  $N(32, 14)$

$$\text{Now, } P(2X - Y \leq 2\lambda) = P(X + 2Y \geq \lambda)$$

$$\text{i.e., } P(U \leq 2\lambda) = P(V \geq \lambda)$$

$$\text{i.e., } P\left(\frac{U-4}{8} \leq \frac{2\lambda-4}{8}\right)$$

$$= P\left(\frac{V-32}{14} \geq \frac{\lambda-32}{14}\right)$$

$$\text{i.e., } P\left(Z \leq \frac{2\lambda-4}{8}\right) = P\left(Z \geq \frac{\lambda-32}{14}\right),$$

where  $Z$  is the standard normal variable.

$$\therefore \frac{2\lambda-4}{8} = -\left(\frac{\lambda-32}{14}\right)$$

$$\text{i.e., } 28\lambda - 56 = 256 - 8\lambda$$

$$\therefore \lambda = \frac{26}{3}$$

**Example 30**

Fit a normal distribution to the following frequency distribution and hence find the theoretical frequencies:

$x$ :	125, 135, 145, 155, 165, 175, 185, 195, 205	Total
$f$ :	1, 1, 14, 22, 25, 19, 13, 3, 2	100

To fit a normal distribution for the given data, we require the density function of the normal distribution which involves the mean and SD. Let us now compute the mean  $\bar{x}$  and SDs of the given distribution and assume them as  $\mu$  and  $\sigma$  of the approximate normal distribution.

$x$	$f$	$d = \frac{x - 165}{10}$	$fd$	$fd^2$
125	1	-4	-4	16
135	1	-3	-3	9
145	14	-2	-28	56
155	22	-1	-22	22
165	25	0	0	0
175	19	1	19	19
185	13	2	26	52
195	3	3	9	27
205	2	4	8	32
Total:	100	-	5	233

$$\bar{x} = A + \frac{c}{N} \sum fd = 165 + \frac{10}{100} \times 5 = 165.5$$

$$\begin{aligned} s^2 &= c^2 \left\{ \frac{1}{N} \sum fd^2 - \left( \frac{1}{N} \sum f d \right)^2 \right\} \\ &= 10^2 (2.33 - 0.0025) \\ &= 232.75 \end{aligned}$$

$$\therefore s = 15.26$$

Therefore, the density function of the approximate normal distribution that fits the given distribution is

$$f(x) = \frac{1}{15.26 \sqrt{2\pi}} e^{-(x - 165.5)^2 / 465.5} \quad -\infty < x < \infty$$

To find the theoretical frequency of the class  $120 \leq X \leq 130$ , whose mid-value is 125, we first get  $P(120 \leq X \leq 130) = P\left\{\frac{120 - 165.5}{15.26} \leq Z \leq \frac{130 - 165.5}{15.26}\right\}$  and multiply this probability by the total frequency. Proceeding likewise, we get all the theoretical frequencies. The computations are shown in the table given in the next page.

Class mid-value ( $x$ )	Class		$P(Z_L \leq Z \leq Z_R)$	Theoretical frequency	Corrected frequency
	$(X_L \leq X \leq X_R)$	$Z_L \leq Z \leq Z_R$			
125	120 $\leq X \leq$ 130	$-2.98 \leq Z \leq -2.33$	0.0085	0.85	1
135	130 $\leq X \leq$ 140	$-2.33 \leq Z \leq -1.67$	0.0376	3.76	4
145	140 $\leq X \leq$ 150	$-1.67 \leq Z \leq -1.02$	0.1064	10.64	11
155	150 $\leq X \leq$ 160	$-1.02 \leq Z \leq -0.36$	0.2055	20.55	21
165	160 $\leq X \leq$ 170	$-0.36 \leq Z \leq 0.29$	0.2547	25.47	25
175	170 $\leq X \leq$ 180	$0.29 \leq Z \leq 0.95$	0.2148	21.48	21
185	180 $\leq X \leq$ 190	$0.95 \leq Z \leq 1.61$	0.1174	11.74	12
195	190 $\leq X \leq$ 200	$1.61 \leq Z \leq 2.26$	0.0418	4.18	4
205	200 $\leq X \leq$ 210	$2.26 \leq Z \leq 2.92$	0.0102	1.02	1
					Total: 100

### Exercise 5(B)

**Part-A (Short-answer Questions)**

1. If  $X$  has uniform distribution in  $(-3, 3)$ , find  $P(|X - 2| < 2)$ .
2. If  $X$  has uniform distribution in  $(-a, a)$ ,  $a > 0$ , find ‘ $a$ ’ such that  $P(|X| < 1) = P(|X| > 1)$ .
3. If the MGF of a continuous RV  $X$  is  $\frac{1}{t}(e^{5t} - e^{4t})$ ,  $t \neq 0$ , what is the distribution of  $X$ ? What are its mean and variance?
4. A continuous RV  $X$  has the density function  $c e^{-x/5}$ ;  $x > 0$ . Find  $c$ ,  $E(X)$  and  $\text{Var}(X)$ .
5. What do you mean by memoryless property of the exponential distribution?
6. If  $X$  and  $Y$  are independent identically distributed RVs, each with density function  $e^{-x}$ ,  $x > 0$ , find the density function of  $(X + Y)$ .
7. Define Erlang distribution and also give its mean and variance.
8. Write down the MGF of simple Gamma distribution and hence find its mean and variance.
9. Give the values of  $\beta_1$  and  $\beta_2$  coefficients of the Erlang distribution with parameters  $(1, k)$ .
10. Find where the maximum occurs for the Erlang density function.
11. If  $X$  has uniform distribution in  $(0, 2)$  and  $Y$  has exponential distribution with parameter  $\lambda$ , find  $\lambda$  such that  $P(X < 1) = P(Y < 1)$ .
12. If  $X$  has uniform distribution in  $(-1, 3)$  and  $Y$  has exponential distribution with parameter  $\lambda$ , find  $\lambda$  such that  $\text{Var}(X) = \text{Var}(Y)$ .
13. Define Weibull distribution and also give its mean and variance.
14. Find the value of  $k$ , mean and variance of the normal distribution whose density function is  $k \cdot 2^{-x^2}$   $-\infty < x < \infty$ .
15. If  $X$  follows  $N(30, 5)$  and  $Y$  follows  $N(15, 10)$  show that  $P(26 \leq X \leq 40) = P(7 \leq Y \leq 35)$ .
16. If  $X$  follows  $N(3, 2)$ , find the value of  $k$  such that  $P(|X - 3| > k) = 0.05$ .
17. If  $\log_{10} X$  follows  $N(4, 2)$ , find  $P(1.202 < X < 83180000)$ , given that  $\log_{10}(1202) = 3.08$  and  $\log_{10}(8318) = 3.92$ .
18. For a certain normal distribution, the first moment about 10 is 40 and the fourth moment about 50 is 48. What are its mean and SD?
19. Show that, for a normal distribution, the quartile deviation, the mean deviation and the standard deviation are in the ratio  $10 : 12 : 15$ .
20. If 2 normal universes  $A$  and  $B$  have the same total frequency, but the SD of  $A$  is  $k$  times the SD of  $B$ , prove that the maximum frequency of  $A$  is  $\frac{1}{k}$  times that of  $B$ .
21. State the reproductive property of normal distribution.
22. If  $X$  and  $Y$  are independent RVs having  $N(1, 2)$  and  $N(2, 2)$  respectively find the density function of  $(X + 2Y)$ .

23. Why is normal distribution considered an important distribution?

**Part-B**

24.  $X$  is uniformly distributed with mean 1 and variance  $\frac{4}{3}$ . If 3 independent observations of  $X$  are made, what is the probability that all of them are negative?
25. A point  $D$  is chosen on the line  $AB$  whose length is 1 and whose mid-point is  $C$ . If the distance  $X$  from  $D$  to  $A$  is a RV having a uniform distribution in  $(0, a)$ , what is 'the probability that  $AD$ ,  $BD$  and  $AC$  will form a triangle?
26. A passenger arrives at a bus stop at 10 am, knowing that the bus will arrive at some time uniformly distributed between 10 am and 10.30 am. What is the probability that he will have to wait longer than 10 min? If at 10.15 am the bus has not yet arrived, what is the probability that he will have to wait at least 10 additional minutes?
27. A man and a woman agree to meet at a certain place between 10 am. and 11 am. They agree that the one arriving first will wait 15 min for the other to arrive. Assuming that the arrival times are independent and uniformly distributed, find the probability that they meet.
28. The RVs  $a$  and  $b$  are independently and uniformly distributed in the intervals  $(0, 6)$  and  $(0, 9)$  respectively. Find the probability that the roots of the equation  $x^2 - ax + b = 0$  are real.
29. If  $a, b, c$  are randomly chosen between 0 and 1, find the probability that the quadratic equation  $ax^2 + bx + c = 0$  has real roots.
30.  $X, Y, Z$  are independent RVs, each following a uniform distribution in  $(0,1)$ . If  $U = \text{maximum of } X, Y, Z$  and  $V = \text{minimum of } X, Y, Z$ , find  
 (i)  $P(U \leq \frac{1}{2})$ ,   (ii)  $P(V \geq \frac{1}{3})$ , and   (iii)  $P\{ U \leq \frac{1}{2} \text{ and } V \geq \frac{1}{3} \}$ .
31. If the number of kilometres that a car can run before its battery wears out is exponentially distributed with an average value of 10,000 km and if the owner desires to take a 5000 km trip, what is the probability that he will be able to complete his trip without having to replace the car battery. Assume that the car has been used for some time. What is the probability, when the distribution is not exponential?
32. If  $X$  is exponentially distributed, prove that the probability that  $X$  exceeds its expected value is less than 0.5.
33. The amount of time that a watch will run without having to be reset is a RV having an exponential distribution with mean 120 days. Find the probability that such a watch will  
 (i) have to be set in less than 24 days, and  
 (ii) not have to be reset in at least 180 days.
34. The daily consumption of milk in excess of 20,000 litres in a town is approximately exponentially distributed with parameter  $1/3000$ . The town

has a daily stock of 35,000L. What is the probability that of 2 days selected at random, the stock is insufficient for both days?

[Hint: If  $Y$  denotes the daily consumption of milk, then  $X = Y - 20,000$  follows the exponential distribution.]

35. The length of the shower on a tropical island during rainy season has an exponential distribution with parameter 2, time being measured in minutes. What is the probability that a shower will last more than 3 min? If a shower has already lasted for 2 min, what is the probability that it will last for at least one more minute?
36. If  $X$  is exponentially distributed with parameter  $\lambda$ , find the value of  $k$  such that

$$P(X > k)/P(X \leq K) = a.$$

37. If  $X$  is exponentially distributed with parameter  $\lambda$ , prove that the RV  $Y = e^{-\lambda X}$  is uniformly distributed in  $(0, 1)$ .
38. If  $X_1, X_2, X_3$  are independent RVs having exponential distributions with parameters  $\lambda_1, \lambda_2, \lambda_3$  respectively, prove that  $Y = \min(X_1, X_2, X_3)$  follows an exponential distribution with parameter  $(\lambda_1 + \lambda_2 + \lambda_3)$ .  
[Hint: Find the distribution function of  $Y = F(y) = 1 - P\{\min(X_1, X_2, X_3) > y\}$ .]
39. The daily consumption of milk in a town in excess of 20,000L is approximately distributed as an Erlang variate with parameters

$\lambda = \frac{1}{10,000}$  and  $k = 2$ . The town has a daily stock of 30,000L. What is the probability that the stock is insufficient on a particular day?

40. Find the probabilities that the value of a RV will exceed 4, if it has an Erlang distribution with

(i)  $\lambda = \frac{1}{3}$  and  $k = 2$ , and (ii)  $\lambda = \frac{1}{4}$  and  $k = 3$ .

41. Show that for the Erlang distribution with parameters  $\lambda, k$ , (mean-mode)/

$$\text{SD} = \frac{1}{\sqrt{k}}$$

[Hint: If  $f(x)$  is the Erlang density function, the mode is the value of  $x$  for which  $f(x)$  is maximum.]

42. If  $X$  follows the Erlang distribution with parameters  $\lambda$  and  $k$ , prove that the expected value of the positive square root of  $X$  is  $\sqrt{\left(k + \frac{1}{2}\right)} / \sqrt{\lambda}$ .

43. If  $X_1, X_2, \dots, X_n$  are independent RVs, each following the same exponential distribution with parameter  $\lambda$ , prove that  $X_1 + X_2 + \dots + X_n$  follows an Erlang distribution with parameters  $\lambda$  and  $n$ .

[Hint: Use moment generating function. Also see Example (7) in Section 4(b).]

44. A random sample of size  $n$  is taken from a population which is exponentially distributed with parameter  $\lambda$ . If  $\bar{X}$  is the sample mean, show that  $n\lambda \bar{X}$  follows a simple Gamma distribution with parameter  $n$ .  
**[Hint:** Use moment generating function.]
45. If the service life, in hours, of a semiconductor is a RV having a Weibull distribution with the parameters  $\alpha = 0.025$  and  $\beta = 0.5$ ,
- How long can such a semiconductor be expected to last?
  - What is the probability that such a semiconductor will still be in operating condition after 4000 h?
46. Find the mode of the Weibull distribution with parameters  $\alpha$  and  $\beta$  when  $\alpha > 1$ .
47. If the hazard rate at time  $t$  of a system is given by  $h(t) = \alpha \beta t^{\beta-1}$ , prove that the time to failure of the system follows a Weibull distribution with parameters  $\alpha$  and  $\beta$ .
48. If the RV  $X$  follows an exponential distribution with parameter 2, prove that  $Y = X^3$  follows a Weibull distribution with parameters 2 and  $\frac{1}{3}$ .
49. Find the probability of failure-free performance of roller-bearings over a period of  $10^4$  h if the life expectancy of the bearings is defined by Weibull distribution with parameters  $\alpha = 10^{-7}$  and  $\beta = 1.5$ .  
**[Hint:**  $P(\text{failure-free performance over a period } t) = P(\text{the component does not fail in } (0, t)] = P(T \geq t)$ , where  $T$  is the life expectancy or time to failure of the component.]
50. The time when a country bus passes a certain point is distributed normally with a mean 9.25 am and a SD of 3 min. What is the least time one could arrive at this point and still have a probability of 0.99 of catching the bus?  
**[Hint:** If  $T$  is the time in minutes past 9 am, then  $T$  follows  $N(25, 3)$ .]
51. The marks obtained by a number of students in a certain subject are assumed to be approximately normally distributed with mean 55 and a SD of 5. If 5 students are taken at random from this set, what is the probability that 3 of them would have scored marks above 60?
52. The life lengths in hours of 2 electronic devices  $A$  and  $B$  have distributions  $N(40, 6)$  and  $N(45, 3)$  respectively. If the electronic device is to be used for a 45-h period, which device is to be preferred? If it is to be used for a 48-h period, which device is to be preferred?
53. The mean and SD of a certain group of 1000 high school grades, that are normally distributed are 78% and 11% respectively.
- Find how many grades were above 90%?
  - What was the highest grade of the lowest 10?
  - What was the semi-interquartile range (Quartile deviation)?
  - Within what limits did the middle 900 lie?
54. The local authorities in a certain city instal 10,000 electric lamps in the streets of the city. If these lamps have an average life of 1,000 burning

hours with a standard deviation of 200 h, how many lamps might be expected to fail (i) in the first 800 burning hours? (ii) between 800 and 1200 burning hours? After how many burning hours would you expect (iii) 10% of the lamps to fail? (iv) 10% of the lamps to be still burning? Assume that the life of lamps is normally distributed.

55. In a normal population with mean 15 and SD 3.5, it is found that 647 observations exceed 16.25. What is the total number of observations in the population?
56. A RV has a normal distribution with SD 10. If the probability that the RV will take on a value less than 82.5 is 0.8212, what is the probability that it will take on a value greater than 58.3?
57. In a normal distribution, 7% of the items are under 35 and 89% are under 63. What are the mean and standard deviation of the distribution? What percentage of items are under 49?
58. A normal population has coefficient of variation equal to 2% and 8% of the population lies above 120 cm. What percentage of the population lies below 115 cm?
59. The breaking strength  $X$  of a certain kind of rope (in kg) has distribution  $N(45, 1.8)$ . Each 50 metre coil of rope brings a profit of Rs 1000, provided  $X > 43$ . If  $X \leq 43$ , the rope may be used for a different purpose and a profit of Rs. 400 per coil is realised. Find the expected profit per coil.
60. The mean and standard deviation of marks in Mathematics are 45 and 10 respectively. The corresponding values for computer science are 50 and 15 respectively. Assuming that the marks in the two subjects are independent normal variates, find the probability that a student scores a total of marks lying between 100 and 120 in the 2 subjects.
61. If  $\log_{10} X$  has the distribution  $N(7, \sqrt{3})$  and  $\log_{10} Y$  has the distribution  $N(3, 1)$ , find  $P\left\{1.202 < \frac{X}{Y} < 8318 \times 10^4\right\}$ , given that  $X$  and  $Y$  are independent.  
**[Hint:** Find  $P\{\log(1.202) < (\log X - \log Y) < \log(10^4 \times 8318)\}]$
62. If  $X$  and  $Y$  are independent RVs having normal distributions with a common mean  $\mu$ , but with variances 4 and 48 respectively, such that  $P(X + 2Y \leq 3) = P(2X - Y \geq 4)$ , determine  $\mu$ .
63. Fit a normal distribution to the following distribution and hence find the theoretical frequencies:

Class	:	60–65	65–70	70–75	75–80	80–85
Frequency :		3	21	150	335	326
		85–90	90–95	95–100	Total	
		135	26	4		1000

## ANSWERS

### **Exercise 5(A)**

1.  $p = \frac{2}{3}, q = \frac{1}{3}; n = 6; P(X \geq 1) = 1 - nC_0 p^0 q^n = \frac{728}{729}$

2.  $\mu_3 = npq(q-p); \beta_1 = \frac{\mu_3^2}{\mu_2^3} = \frac{(q-p)^2}{npq}$

3. Reqd. no. =  $256 \times 8C_4 \left(\frac{1}{2}\right)^8 = 70$

4.  $p = \frac{2}{3}, q = \frac{1}{3}; P(X \geq 1) = 1 - 4C_0 p^0 q^4 = \frac{80}{81}$

5. Reqd. Probability =  $\sum_{r=1}^3 4C_r \left(\frac{1}{2}\right)^4 = \frac{7}{8}$

6.  $\mu = 5, \sigma^2 = 4; P(X < \mu - 2\sigma) = P(X < 1) = \left(\frac{4}{5}\right)^{25}$

7.  $V = np(1-p); \frac{dV}{dp} = n(1-2p) = 0, \text{ when } p = \frac{1}{2} \text{ and } \frac{d^2V}{dp^2} < 0$

8. The MGF of  $B(n, p)$  is  $(q + p e^t)^n$ . The given MGF is that of  $B(6, 0.4)$ .  
Hence, mean = 2.4 and SD = 1.2

9. When  $p_1 = p_2$ , the sum is also a binomial variate

10.  $(X + Y)$  follows  $B\left(8, \frac{1}{3}\right); P(X + Y \geq 1) = 1 - \left(\frac{2}{3}\right)^8 = \frac{6305}{6561}$

11.  $P(X = r) = e^{-2} 2^r/r!$

12.  $2e^{-\lambda} + e^{-\lambda} \frac{\lambda^2}{2} = 2 e^{-\lambda} \lambda. \therefore (\lambda - 2)^2 = 0 \text{ or } \lambda = 2$

13. If  $\lambda$  is the parameter of the Poisson distribution,  $\text{Var}(X) = E(X^2) - E^2(X)$ .  
i.e.,  $\lambda = 6 - \lambda^2$

$$\therefore \lambda^2 + \lambda - 6 = 0 \quad \therefore \lambda = E(X) = 2, \text{ since } \lambda > 0$$

14. If  $\lambda$  is the parameter,  $e^{-\lambda} = 0.5 \therefore \text{Var}(X) = \lambda = \log 2$

15.  $E(X^2) = \lambda^2 + \lambda = \lambda(\lambda + 1) = \lambda E(X + 1)$

16.  $P(X \text{ is even}) = e^{-\lambda} \left\{1 + \frac{\lambda^2}{2!} + \frac{\lambda^4}{4!} + \dots\right\} = e^{-\lambda} \cosh \lambda = \frac{1}{\lambda} e^{-\lambda} (e^\lambda + e^{-\lambda})$   
 $= \frac{1}{2} (1 + e^{-2\lambda})$

17.  $e^{4(e^t - 1)}$  is the MGF of a Poisson distribution with parameter 4.  $\therefore \mu = 4$  and  $\sigma = 2$ .  
 $\therefore P(X = 6) = e^{-4} \cdot 4^6/6!$
18.  $(X + Y)$  is a Poisson variate with parameter 2.  
 $\therefore P(X + Y = 2) = e^{-2} 2^2/2! = \frac{2}{e}$
19. The given distribution is a geometric distribution with  $p = e^{-t}$   
 $\therefore E(X) = \frac{q}{p} = e^t - 1$  and  $\text{Var}(X) = \frac{q}{p^2} = e^t(e^t - 1)$ .
20.  $P(X = r) = pq^{r-1}$   
 $\therefore P(X \text{ is odd}) = p + pq^2 + pq^4 + \dots = \frac{p}{1-q^2} = \frac{1}{1+q}$ .
21. The given distribution is a geometric distribution with  $p = \frac{2}{3}$  and  $q = \frac{1}{3}$ .  
 $\therefore E(X) = \frac{q}{p} = \frac{1}{2}$  and  $\text{Var}(X) = \frac{q}{p^2} = \frac{3}{4}$
22.  $\frac{q}{p^2} = 2 \frac{q}{p} \therefore p = \frac{1}{2} \therefore P(X = r) = 2^{-r}; r = 1, 2, \dots, \infty$
23.  $M(t) = \sum_{r=1}^{\infty} e^{tr} q^{r-1} p = p e^t \sum_{r=1}^{\infty} (q e^t)^{r-1} = \frac{pe^t}{1-qe^t}$
24. The given MGF  $\frac{1}{5} e^t / \left(1 - \frac{4}{5} e^t\right)$  is that of a geometric distribution with  
 $p = \frac{1}{5}$  and  $q = \frac{4}{5}$   
 $\therefore E(X) = 4$  and  $\text{Var}(X) = 20$ .
26. Mean =  $\frac{nk}{N}$  and Variance =  $\frac{nk(N-k)(N-n)}{N^2(N-1)}$
27. Conditions:  $N \rightarrow \infty$  and  $\frac{k}{N} = p$ ;  $\lim \left[ n \left( \frac{k}{N} \right) \right] = np$ ;  
 $\lim \left[ n \left( \frac{k}{N} \right) \left( 1 - \frac{k}{N} \right) \left( \frac{1 - \frac{n}{N}}{1 - \frac{1}{N}} \right) \right] = npq$
28.  $P(X > 1) = .005; 0.015$

29. 17
31.  $P(X = 3) = 0.0452; P(X > 3) = 0.0099$
32.  $P(Y = k) = \frac{1}{16} \cdot 4C_{\left(\frac{k}{2} + 2\right)}, k = -4, -2, 0, 2, 4$
33. & 34.  $\frac{2}{3} < p < 1$
35.  $\frac{1}{3} < q < 1$
38.  $\frac{1}{6}$
39. 0.016; 0.044
40. (i) 0.5551, (ii) 4
41. 10, 60, 150, 200, 150, 60, 10
42. 7, 26, 37, 34, 6
43. (i) 0.2231, (ii) 0.1913
44. 0.019
45. 0.1247; 0.0367
46. (i) 9802, (ii) 198, (iii) 9998
47. 0.067
48. (i) 0.1606, (ii) 0.1512
49. 0.0028
50. 0.0424
51. 0.0498
52. 0.8088
53. 0.1144
54. 0.6244, 2.3
55. 0.9577
56. 301, 362, 217, 87, 26, 6, 1
57. 0.92
58. 0.4529
59.  $\frac{1}{10}$
61.  $\frac{25}{61}$
62.  $P(X + Y = k) = (k - 1) p^2 q^{k-2}, k = 2, 3, \dots, \infty; \frac{2}{p}$ .
63. 0.2880
64. (i) 0.164, (ii) 0.165
65. 0.7134

- 66. 0.8121
- 68. 0.0221
- 69. 0.0410
- 70. 0.9922

### Exercise 5(B)

1.  $f(x) = \frac{1}{6}; P(|X - 2| < 2) = P(0 < X < 4) = P(0 < X < 3) = \frac{1}{2}$
2.  $a > 1; P(|X| < 1) = \int_{-1}^1 \frac{1}{2a} dx = \frac{1}{a}; P(|X| > 1) = 1 - \frac{1}{a} \therefore a = 2$
3. Comparing the given MGF with that of uniform distribution, i.e.,  $\frac{1}{b-a} \cdot \frac{e^{bt} - e^{at}}{a}$ ,  $X$  follows  $U(4, 5)$   $E(X) = \frac{9}{2}$  and  $\text{Var}(X) = \frac{1}{12}$ .
4.  $c \int_0^\infty e^{-x/5} dx \therefore c = \frac{1}{5}$ ; Exponential distribution with parameter  $\lambda = \frac{1}{5}$ ;  $E(X) = \frac{1}{\lambda} = 5$ ;  $\text{Var}(X) = \frac{1}{\lambda^2} = 25$
5. If  $X$  is exponentially distributed, then  $P(X > s + t | X > s) = P(X > t)$ , for any  $s, t > 0$
6.  $X$  and  $Y$  follows exponential distribution with parameter 1.  $\therefore M_X(t) = M_Y(t) = \left(1 - \frac{t}{\lambda}\right)^{-1}$   
 $\therefore M_{X+Y}(t) = M_X(t) M_Y(t) = \left(1 - \frac{t}{\lambda}\right)^{-2} =$  MGF of Erlang distribution  $(\lambda, 2)$ . Its density function is  $\lambda^2 x e^{-\lambda x}, x > 0$
8.  $M(t) = (1-t)^{-k} = 1 + kt + \frac{k(k+1)}{2!} t^2 + \dots + \infty$   
 $E(X) = k; E(X^2) = k(k+1); \text{Var}(X) = k$
9. Using  $M(t) = (1-t)^{-k}$ ,  $\mu_3 = 2k; \mu_4 = 3k^2 + 6k$   
 $\therefore \beta_1 = \mu_3^2 / \mu_2^3 = \frac{4}{k}$  and  $\beta_2 = \mu_4 / \mu_2^2 = 3 + \frac{6}{k}$
10. Maximum of  $f(x) = \frac{\lambda^k}{\lfloor k \rfloor!} |x|^{k-1} e^{-\lambda x}$  occurs at  $x = \frac{\lambda - 1}{k}$
11.  $P(X < 1) = \int_0^1 \frac{1}{2} dx = \frac{1}{2}; P(Y < 1) = \int_0^1 \lambda e^{-\lambda x} dx = 1 - e^{-\lambda} \cdot 1 - e^{-\lambda} = \frac{1}{2};$

$$e^{-\lambda} = \frac{1}{2} \therefore \lambda = \log 2$$

12.  $\text{Var}(X) = \frac{(3+1)^2}{12} = \frac{4}{3}$ ;  $\text{Var}(Y) = \frac{1}{\lambda^2}$   
 $\therefore \frac{1}{\lambda^2} = \frac{4}{3}$  gives  $\lambda = \sqrt{3}/2$

14.  $f(x) = k e^{-(\log 2)x^2}$ , Mean = 0; variance is given by  $2\sigma^2 = \frac{1}{\log 2}$   
 $\therefore \sigma = \frac{1}{\sqrt{\log 4}}$ ;  $k = \frac{1}{\sigma \sqrt{2\pi}} = \sqrt{\frac{\log 2}{\pi}}$

15.  $P(26 \leq X \leq 40) = P(-0.8 \leq Z \leq 2)$  and  $P(7 \leq Y \leq 35) = P(-0.8 \leq Z \leq 2)$ .  
Hence, equality.  
16.  $P(|X - 3| < k) = 0.95 \therefore P(-k < X - 3 < k) = 0.95$   
i.e.,  $P\left(-\frac{k}{2} < Z < \frac{k}{2}\right) = 0.95$ ;  $\therefore \frac{k}{2} = 1.96$   
 $\therefore k = 3.92$

17. Req'd. probability =  $P\{\log_{10}(1.202) < \log X < \log_{10}83180000\} = P\{0.08 < \log X < 7.92\}$   
=  $P\{-1.96 < Z < 1.96\} = 0.95$

18.  $E(X - 10) = 40$ .  $\therefore E(X) = \mu = 50$ .  $E(X - 50)^4 = 48$   
 $\therefore \mu_4 = 48$ , i.e.,  $3\sigma^2 = 48 \therefore \sigma = 2$

19. QD : MD : SD =  $\frac{2}{3}\sigma$  :  $\frac{4}{5}\sigma$  :  $\sigma = 10 : 12 : 15$

20. Max. ordinate of A and B are  $\frac{1}{\sigma_A \sqrt{2\pi}}$  and  $\frac{1}{\sigma_B \sqrt{2\pi}}$

$$\text{Max } (f_A) = C \frac{N}{\sigma_A \sqrt{2\pi}} \text{ and } \text{Max } (f_B) = C \frac{N}{\sigma_B \sqrt{2\pi}}$$

$$\therefore \frac{\text{Max } (f_A)}{\text{Max } (f_B)} = \frac{\sigma_B}{\sigma_A} = \frac{1}{k}$$

21. See Property (7) of normal distribution.  
22. By additive property,  $(X + 2Y)$  follows  $N(5, 2\sqrt{5})$ . Its density function  
is  $f(x) = \frac{1}{2\sqrt{10\pi}} e^{-(x-5)^2/40}$ ,  $-\infty < x < \infty$

24.  $\frac{1}{64}$

25.  $\frac{9}{16}$

26.  $\frac{2}{3}; \frac{1}{3}$

27.  $\frac{7}{16}$

28.  $\frac{1}{3}$

29.  $\frac{1}{9}$

30. (i)  $\frac{1}{8}$ , (ii)  $\frac{8}{27}$ , (iii)  $\frac{1}{216}$

31. (i) 0.6065;  
(ii)  $\frac{1 - F(t + 5000)}{1 - F(t)}$ , where  $F(t)$  is the distribution function of the life of the car.

33. (i) 0.1813. (ii) 0.2231

34.  $e^{-10}$

35. (i) 0.0025, (ii) 0.1353

36.  $\lambda^{-1} \log(1 + a^{-1})$

39. 0.7378

40. (i) 0.0879, (ii) 0.9197

45. (i) 3200h, (ii) 0.2057

46.  $\beta(\alpha - 1)$

49. 0.9048

50. 9.18 am

51. 0.0283

52. (i) B, (ii) B

53. (i) 138, (ii) 52%, (iii) 7.3%, (iv) 60% and 96.1%

54. (i) 1587, (ii) 6826, (iii) 744h, (iv) 1256h

55. 1800

56. 0.4332

57. 50.3, 10.33, 45%

58. 116.72, 2.33, 23%

59. Rs. 920 nearly

60. 0.3074

71. 0.95

62. 2.1

63. 3, 31, 148, 322, 319, 144, 30, 3

# Chapter 6

## Random Processes

In electrical systems, voltage or current waveforms are used as signals for collecting, transmitting or processing information, as well as for controlling and providing power to a variety of devices. These signals (voltage or current waveforms) are functions of time and are of two classes—deterministic and random. *Deterministic signals* can be described by the usual mathematical functions with time  $t$  as the independent variable. But a random signal always has some element of uncertainty associated with it and hence it is not possible to determine its value exactly at any given point of time. However, we may be able to describe the random signal in terms of its average properties such as the average power in the random signal, its spectral distribution and the probability that the signal amplitude exceeds a given value. The probabilistic model used for characterising a random signal is called a *random process* or *stochastic process*.

A random variable (RV) is a rule (or function) that assigns a real number to every outcome of a random experiment, while a random process is a rule (or function) that assigns a time function to every outcome of a random experiment. For example, consider the random experiment of tossing a dice at  $t = 0$  and observing the number on the top face. The sample space of this experiment consists of the outcomes  $\{1, 2, 3, \dots, 6\}$ . For each outcome of the experiment, let us arbitrary assign a function of time  $t$  ( $0 \leq t < \infty$ ) in the following manner.

<i>Outcome:</i>	1	2	3	4	5	6
<i>Function of time</i>	$x_1(t)$	$x_2(t)$	$x_3(t)$	$x_4(t)$	$x_5(t)$	$x_6(t)$
	= -4	= -2	= 2	= 4	= $-t/2$	= $t/2$

The set of functions  $\{x_1(t), x_2(t), \dots, x_6(t)\}$  represents a random process.

**Definition:** A random process is a collection (or ensemble) of RVs  $\{X(s, t)\}$  that are functions of a real variable, namely time  $t$  where  $s \in S$  (sample space) and  $t \in T$  (parameter set or index set).

The set of possible values of any individual member of the random process is called *state space*. Any individual member itself is called a *sample function* or a realisation of the process.

**Note**

- (i) If  $s$  and  $t$  are fixed,  $\{X(s, t)\}$  is a number.
- (ii) If  $t$  is fixed  $\{X(s, t)\}$  is a RV.
- (iii) If  $s$  is fixed,  $\{X(s, t)\}$  is a single time function.
- (iv) If  $s$  and  $t$  are variables,  $\{X(s, t)\}$  is a collection of RVs that are time functions.

**Notation:** As the dependence of a random process on  $s$  is obvious,  $s$  will be omitted hereafter in the notation of a random process. If the parameter set  $T$  is discrete, the random process will be noted by  $\{X(n)\}$  or  $\{X_n\}$ .

If the parameter set  $T$  is continuous, the process will be denoted by  $\{X(t)\}$ .

## Classification of Random Processes

Depending on the continuous or discrete nature of the state space  $S$  and parameter set  $T$ , a random process can be classified into four types:

- (i) If both  $T$  and  $S$  are discrete, the random process is called a *discrete random sequence*. For example, if  $X_n$  represents the outcome of the  $n$ th toss of a fair dice, then  $\{X_n, n \geq 1\}$  is a discrete random sequence, since  $T = \{1, 2, 3, \dots\}$  and  $S = \{1, 2, 3, 4, 5, 6\}$ .
- (ii) If  $T$  is discrete and  $S$  is continuous, the random process is called a *continuous random sequence*.  
For example, if  $X_n$  represents the temperature at the end of the  $n$ th hour of a day, then  $\{X_n, 1 \leq n \leq 24\}$  is a continuous random sequence, since temperature can take any value in an interval and hence continuous.
- (iii) If  $T$  is continuous and  $S$  is discrete, the random process is called a *discrete random process*.  
For example, if  $X(t)$  represents the number of telephone calls received in the interval  $(0, t)$ , then  $X(t)$  is a discrete random process, since  $S = \{0, 1, 2, 3, \dots\}$ .
- (iv) If both  $T$  and  $S$  are continuous, the random process is called *continuous random process*. For example, if  $X(t)$  represents the maximum temperature at a place in the interval  $(0, t)$ ,  $\{X(t)\}$  is a continuous random process.  
In the names given above, the word ‘discrete’ or ‘continuous’ is used to refer to the nature of  $S$  and the word ‘sequence’ or ‘process’ is used to refer to the nature of  $T$ .

## Methods of Description of a Random Process

Since a random process is an indexed set of RVs, we can obviously use the joint probability distribution functions to describe a random process.

For a specific  $t$ ,  $X(t)$  is a RV as was observed earlier.

$F(x, t) = P\{X(t) \leq x\}$  is called the first-order distribution of the process  $\{X(t)\}$

and  $f(x, t) = \frac{\partial}{\partial x} F(x, t)$  is called the first-order density of  $\{X(t)\}$ .

$F(x_1, x_2, t_1, t_2) = P\{X(t_1) \leq x_1; X(t_2) \leq x_2\}$  is the joint distribution of the RVs  $X(t_1)$  and  $X(t_2)$  and is called the second-order distribution of the process  $\{X(t)\}$  and  $f(x_1, x_2, t_1, t_2) = \frac{\partial^2}{\partial x_1 \cdot \partial x_2} F(x_1, x_2, t_1, t_2)$  is called the second-order density of  $\{X(t)\}$ .

Similarly the  $n$ th order distribution  $\{X(t)\}$  is the joint distribution  $F(x_1, x_2, \dots, x_n; t_1, t_2, \dots, t_n)$  of the RVs  $X(t_1), X(t_2), \dots, X(t_n)$ .

The first-order distribution function describes the instantaneous *amplitude* distribution of the process and the second-order distribution function tells us something about the structure of the signal in the time domain and, hence, the spectral content of the signal. Although the higher-order distributions describe the process in a more detailed manner, the first and second-order distribution functions are primarily used to describe the process.

## Special Classes of Random Processes

The important feature of a random process is the relationship among the members of the family. Usually, the nature of relationship is understood by the joint distribution function of the member RVs. A random process is said to be specified only when the parameter set, the state space and the nature of dependence relationship existing among the members of the family are specified.

Based on the dependence relationship among the members of the process, random processes are classified broadly into a few special types such as the ones explained below.

### (i) Markov Process

If, for  $t_1 < t_2 < t_3 < \dots < t_n < t$ ,  $P\{X(t) \leq x/X(t_1) = x_1, X(t_2) = x_2, \dots, X(t_n) = x_n\} = P\{X(t) \leq x/X(t_n) = x_n\}$ , then the process  $\{X(t)\}$  is called a Markov process.

In other words, if the future behaviour of a process depends only on the present state, but not on the past, the process is a Markov process.

A discrete parameter Markov process is called a Markov chain.

### (ii) Process with Independent Increments

If, for all choices of  $t_1, t_2, \dots, t_n$  such that  $t_1 < t_2 < t_3 < \dots < t_n$ , the random variables  $X(t_2) - X(t_1), X(t_3) - X(t_2), \dots, X(t_n) - X(t_{n-1})$  are independent, then the process  $\{X(t)\}$  is said to be a random process with independent increments.

If  $T = \{0, 1, 2, \dots\}$  is the parameter set for  $\{X_n\}$ , then  $\{Z_n\}$ , where  $Z_0 = X_0$  and  $Z_n = X_n - X_{n-1}$ , is a random sequence with independent increments if the RVs  $Z_0, Z_1, Z_2, \dots$  are independent.

Two processes with independent increments play an important role in the theory of random processes. One is the Poisson process that has a Poisson distribution for the increments and the other is the Wiener process with a Gaussian distribution for the increments. We will take up the study of Poisson and Gaussian processes in Chapter 7.

### (iii) Stationary Processes

If certain probability distribution or averages do not depend on  $t$ , then the random process  $\{X(t)\}$  is called stationary. A rigorous definition and detailed study of stationary processes will be taken up in the following articles.

## Average Values of Random Processes

As in the case of RVs random processes can be described in terms of averages or expected values, mostly derived from the first and second-order distributions of  $\{X(t)\}$ . *Mean* of the process  $\{X(t)\}$  is the expected value of a typical member  $X(t)$  of the process.

$$\text{i.e., } \mu(t) = E\{X(t)\}$$

*Autocorrelation* of the process  $\{X(t)\}$ , denoted by  $R_{xx}(t_1, t_2)$  or  $R_x(t_1, t_2)$  or  $R(t_1, t_2)$ , the expected value of the product of any two members  $X(t_1)$  and  $X(t_2)$  of the process.

$$\text{i.e., } R(t_1, t_2) = E\{X(t_1) \times X(t_2)\}$$

*Autocovariance* of the process  $\{X(t)\}$ , denoted by  $C_{xx}(t_1, t_2)$  or  $C_x(t_1, t_2)$  or  $C(t_1, t_2)$ , is defined as

$$\begin{aligned} C(t_1, t_2) &= E[\{X(t_1) - \mu(t_1)\} \{X(t_2) - \mu(t_2)\}] \\ &= R(t_1, t_2) - \mu(t_1) \times \mu(t_2) \end{aligned}$$

*Correlation coefficient* of the process  $\{X(t)\}$ , denoted by  $\rho_{xx}(t_1, t_2)$  or  $\rho(t_1, t_2)$ , is defined as

$$\rho(t_1, t_2) = \frac{C(t_1, t_2)}{\sqrt{C(t_1, t_1) \times C(t_2, t_2)}}$$

where  $C(t_1, t_2)$  is the variance of  $X(t_1)$ .

When we deal with 2 or more random processes, we can use joint distribution functions or averages to describe the relationship between them.

*Cross-correlation* of 2 processes  $\{X(t)\}$  and  $\{Y(t)\}$  is defined as

$$R_{xy}(t_1, t_2) = E\{X(t_1) \times Y(t_2)\}$$

*Cross-covariance* of 2 processes  $\{X(t)\}$  and  $\{Y(t)\}$  is defined as

$$C_{xy}(t_1, t_2) = R_{xy}(t_1, t_2) - \mu_x(t_1) \times \mu_y(t_2)$$

*Cross correlation coefficient* of 2 processes  $\{X(t)\}$  and  $\{Y(t)\}$  is defined as

$$\rho_{xy}(t_1, t_2) = \frac{C_{xy}(t_1, t_2)}{\sqrt{C_{xx}(t_1, t_1) \times C_{yy}(t_2, t_2)}}$$

## Stationarity

A random process is called a *strongly stationary process* or *strict sense stationary process* (abbreviated as SSS process), if all its finite dimensional distributions are invariant under translation of time parameter.

That is, if the joint distribution (and hence the joint density) of  $X(t_1), X(t_2), \dots, X(t_n)$  is the same as that of  $X(t_1 + h), X(t_2 + h), \dots, X(t_n + h)$  for all  $t_1, t_2, \dots, t_n$  and  $h > 0$  and for all  $n \geq 1$ , then the random process  $\{X(t)\}$  is called a SSS process. If the definition given above holds good for  $n = 1, 2, \dots, k$  only and not for  $n > k$ , then the process is called  $k$ th order stationary.

**Note** If a random process is a SSS process, as per the definition, its first-order densities must be invariant under translation of time, i.e., the densities of  $X(t)$  and  $X(t + h)$  are the same, i.e.,  $f(x, t) = f(x, t + h)$ .

This is possible only if  $f(x, t)$  is independent of  $t$ .

Therefore, first-order densities (and hence distribution function) of a SSS process are independent of time.

As a consequence,  $E\{X(t)\}$  is also independent of  $t$ .

i.e.,  $E\{X(t)\} = \mu = \text{a constant}$

Also the second-order densities must be invariant under translation of time, i.e., the joint pdf of  $X(t_1), X(t_2)$  is the same as that of  $X(t_1 + h), X(t_2 + h)$ .

i.e.,  $f(x_1, x_2, t_1, t_2) = f(x_1, x_2, t_1 + h, t_2 + h)$

This is possible only if  $f(x_1, x_2, t_1, t_2)$  is function of  $t = t_1 - t_2$ .

Therefore, second-order densities (and hence distribution functions) of a SSS process are functions of  $\tau = t_1 - t_2$ .

As a consequence,  $R(t_1, t_2) = E\{X(t_1) \times X(t_2)\}$  is also a function of  $\tau = t_1 - t_2$ .

It is pointed out that if  $E\{X(t)\}$  is a constant and  $R(t_1, t_2)$ , is a function of  $(t_1 - t_2)$  the random process  $\{X(t)\}$  need not be a SSS process.

The definition of strict sense stationarity can be extended as follows.

Two real-valued random processes  $\{X(t)\}$  and  $\{Y(t)\}$  are said to be jointly stationary in the strict sense, if the joint distribution of  $X(t)$  and  $Y(t)$  are invariant under translation of time.

The complex random process  $\{Z(t)\}$ , where  $Z(t) = X(t) + iY(t)$ , is said to be a SSS process if  $\{X(t)\}$  and  $\{Y(t)\}$  are jointly stationary in the strict sense.

**Wide-sense stationarity:** A random process  $\{X(t)\}$  with finite first- and second-order moments is called a weakly stationary process or covariance stationary process or wide-sense stationary process (abbreviated as WSS process) if its mean is a constant and the autocorrelation depends only on the time difference.

i.e., if  $E\{X(t)\} = \mu$  and

$E\{X(t)\} \times X(t - \tau)\} = R(\tau)$

**Note** From the definitions given above, it is clear that a SSS process with finite first- and second-order moments is a WSS process, while a WSS process need not be a SSS process.

A random process that is not stationary in any sense is called an evolutionary process.

Two random processes  $\{X(t)\}$  and  $\{Y(t)\}$  are said to be jointly stationary in the wide sense, if each process is individually a WSS process and  $R_{xy}(t_1, t_2)$  is a function of  $(t_1, t_2)$  only.

## Example of a SSS Process

Let  $X_n$  denote the presence or absence of a pulse at the  $n$ th time instant in a digital communication system or digital data processing system.

If  $P\{X_n = 1\} = p$  and  $P\{X_n = 0\} = 1 - p = q$ , then the random process (sequence)  $\{X_n, n \geq 1\}$ , called the Bernoulli's process, is a SSS process, for, its first-order distribution is given by

$X_n = r$	1	0
$P(X_n = r)$	$p$	$q$

This distribution is the same for any  $X_n$ , i.e., for  $X_m$  and  $X_{m+p}$ .

Consider the second-order distribution of the process, i.e., the joint distribution of  $X_r$  and  $X_s$ .

$X_r$	$X_s$	
	1	0
1	$p^2$	$pq$
0	$pq$	$q^2$

This joint distribution is the same for the pair of members  $X_r$  and  $X_s$  and for the pair  $X_{r+p}$  and  $X_{s+p}$  of the process.

Consider the third-order distribution of the process, i.e., the joint distribution of  $X_r$ ,  $X_s$  and  $X_t$  that is given below.

$$P\{X_r = 0, X_s = 0, X_t = 0\} = q^3$$

$$P\{X_r = 0, X_s = 0, X_t = 1\} = pq^2$$

$$P\{X_r = 0, X_s = 1, X_t = 0\} = pq^2$$

$$P\{X_r = 0, X_s = 1, X_t = 1\} = p^2q$$

$$P\{X_r = 1, X_s = 0, X_t = 0\} = pq^2$$

$$P\{X_r = 1, X_s = 0, X_t = 1\} = p^2q$$

$$P\{X_r = 1, X_s = 1, X_t = 0\} = p^2q$$

$$P\{X_r = 1, X_s = 1, X_t = 1\} = p^3$$

This joint distribution is the same for the triple of members  $X_r$ ,  $X_s$ ,  $X_t$  and for  $X_{r+p}$ ,  $X_{s+p}$ ,  $X_{t+p}$  of the process, and so on, i.e., distributions of all orders are invariant under translation of time.

**Note** If  $Y_n = \sum_{n=1}^n X_n$  = the total number of pulses from time instant 1 through  $n$ , then the random process  $\{Y_n, n \geq 1\}$ , called the Binomial process, is not a SSS process, for  $P\{Y_n = i\} = nC_i p^i q^{n-i}$  ( $i = 0, 1, 2, \dots, n$ ) depends on  $n$ , i.e., the distributions of  $Y_m$  and  $Y_{m+p}$  are not the same).

## Analytical Representation of a Random Process

Deterministic signals are usually expressed in simple analytical forms such as  $X(t) = e^{-t^2}$  and  $Y(t) = 20 \sin 10 t$ . It is sometimes possible to express a random

process in an analytical form using one or more RVs. For example, consider an FM station that is broadcasting a ‘tone’,  $X(t) = 100 \cos(10^8 t)$ , to a large number of receivers distributed randomly in a metropolitan area. The amplitude and phase of the waveform received by any receiver will depend on the distance between the transmitter and the receiver. Since there are a large number of receivers distributed randomly over an area, the distance can be considered as a continuous RV. Since the amplitude and the phase are functions of distance, they are also RVs. So we can represent the ensemble (collection) of received waveforms by a random process  $\{X(t)\}$  of the form

$$X(t) = A \cos(10^8 t + \theta)$$

where  $A$  and  $\theta$  are RVs representing the amplitude and phase of the received waveforms.

Such representation of a random process in terms of one or more RVs whose probability law is known is used in several applications in communication systems.

### Worked Example 6(A)

#### Example 1

Examine whether the Poisson process  $\{X(t)\}$ , given by the probability law  $P\{X(t) = r\} = e^{-\lambda t} (\lambda t)^r / r!$ ,  $\{r = 0, 1, 2, \dots\}$ , is covariance stationary.

The probability distribution of  $X(t)$  is a Poisson distribution with parameter  $\lambda t$ .

$$\therefore E\{X(t)\} = \lambda t \neq \text{a constant}$$

Therefore, the Poisson process is not covariance stationary.

#### Example 2

The process  $\{X(t)\}$ , whose probability distribution under certain conditions is given by

$$\begin{aligned} P\{X(t) = n\} &= \frac{(at)^{n-1}}{(1+at)^{n+1}}, n = 1, 2, \dots \\ &= \frac{at}{1+at}, n = 0 \end{aligned}$$

Show that it is not stationary.

The probability distribution of  $X(t)$  is

$X(t) = n:$	0	1	2	3	...
$p_n:$	$\frac{at}{1+at}$	$\frac{1}{(1+at)^2}$	$\frac{at}{(1+at)^3}$	$\frac{(at)^2}{(1+at)^4}$	...

$$\begin{aligned}
 E\{X(t)\} &= \sum_{n=0}^{\infty} np_n \\
 &= \frac{1}{(1+at)^2} + \frac{2at}{(1+at)^3} + \frac{3(at)^2}{(1+at)^4} + \dots \\
 &= \frac{1}{(1+at)^2} \{1 + 2\alpha + 3\alpha^2 + \dots\}, \text{ where } \alpha = \frac{at}{1+at} \\
 &= \frac{1}{(1+at)^2} (1-\alpha)^{-2} = \frac{1}{(1+at)^2} (1+at)^2 = 1 \\
 E\{X^2(t)\} &= \sum_{n=0}^{\infty} n^2 p_n = \sum_{n=1}^{\infty} n^2 \frac{(at)^{n-1}}{(1+at)^{n+1}} \\
 &= \frac{1}{(1+at)^2} \left[ \sum_{n=1}^{\infty} n(n+1) \left( \frac{at}{1+at} \right)^{n-1} - \sum_{n=1}^{\infty} n \left( \frac{at}{1+at} \right)^{n-1} \right] \\
 &= \frac{1}{(1+at)^2} \left[ \frac{2}{\left( 1 - \frac{at}{1+at} \right)^3} - \frac{1}{\left( 1 - \frac{at}{1+at} \right)^2} \right] \\
 &= 1 + 2at
 \end{aligned}$$

$$\therefore \text{Var}\{X(t)\} = 2at$$

If  $\{X(t)\}$  is a stationary process,  $E\{X(t)\}$  and  $\text{Var}\{X(t)\}$  are constants.  
Since  $\text{Var}\{X(t)\}$  is a function of  $t$ , the given process is not stationary.

<b>Note</b>	When $\{X(t)\}$ is a stationary process, $R\{t_1, t_2\} = E\{X(t_1)X(t_2)\}$ is function of $(t_1 - t_2)$ .
$\therefore$	$E\{X^2(t)\}$ is a constant.
Also	$E\{X(t)\}$ is a constant.
$\therefore$	$\text{Var}\{X(t)\}$ is a constant.

### Example 3

Show that the random process  $X(t) = A \cos(\omega_0 t + \theta)$  is wide-sense stationary, if  $A$  and  $\omega_0$  are constants and  $\theta$  is a uniformly distributed RV in  $(0, 2\pi)$ .

Since  $\theta$  is uniformly distributed in  $(0, 2\pi)$

$$f_{\theta}(\theta) = \frac{1}{2\pi}, 0 < \theta < 2\pi$$

$$E\{X(t)\} = E\{A \cos(\omega_0 t + \theta)\}$$

$$\begin{aligned}
 &= A \int_0^{2\pi} \frac{1}{2\pi} \cos(\omega_0 t + \theta) d\theta \\
 &\quad [\text{since } E\{g(\theta)\} = \int g(\theta) f_{\theta}(\theta) d\theta]
 \end{aligned}$$

$$\begin{aligned}
&= \frac{A}{2\pi} \{ \sin(2\pi + \omega_0 t) - \sin \omega_0 t \} \\
&= 0 = \text{a constant} \\
E\{X(t_1) X(t_2)\} &= E\{A^2 \cos(\omega_0 t_1 + \theta) \times \cos(\omega_0 t_2 + \theta)\} \\
&= \frac{A^2}{2} E\{\cos[(t_1 + t_2)\omega_0 + 2\theta] + \cos[\omega_0(t_1 - t_2)]\} \\
&= \frac{A^2}{2} \int_0^{2\pi} \frac{1}{2\pi} [\cos\{(t_1 + t_2)\omega_0 + 2\theta\} + \cos\{(t_1 - t_2)\omega_0\}] d\theta \\
&= \frac{A^2}{2} \cos \omega_0(t_1 - t_2)
\end{aligned}$$

i.e.,  $R(t_1, t_2)$  = a function of  $(t_1 - t_2)$

Therefore,  $\{X(t)\}$  is a WSS process.

### Example 4

Given a RV  $Y$  with characteristic function

$$\begin{aligned}
\phi(\omega) &= E\{e^{i\omega Y}\} \\
&= E\{\cos \omega Y + i \sin \omega Y\}
\end{aligned}$$

and a random process defined by  $X(t) = \cos(\lambda t + Y)$ , show that  $\{X(t)\}$  is stationary in the wide sense

if  $\phi(1) = \phi(2) = 0$

$$\begin{aligned}
E\{X(t)\} &= E\{\cos(\lambda t + Y)\} \\
&= \cos \lambda t \times E(\cos Y) - \sin \lambda t \times E(\sin Y) \tag{1}
\end{aligned}$$

Given  $\phi(1) = 0$

i.e.,  $E\{\cos Y + i \sin Y\} = 0$

$$\therefore E(\cos Y) = 0 = E(\sin Y) \tag{2}$$

Using (2) in (1), we get  $E\{X(t)\} = 0$  (3)

$$\begin{aligned}
E\{X(t_1) \times X(t_2)\} &= E\{\cos(\lambda t_1 + Y) \times \cos(\lambda t_2 + Y)\} \\
&= \cos \lambda t_1 \cos \lambda t_2 E(\cos^2 Y) + \sin \lambda t_1 \sin \lambda t_2 E(\sin^2 Y) \\
&\quad - \sin \lambda(t_1 + t_2) E(\sin Y \cos Y) \\
&= \cos \lambda t_1 \cos \lambda t_2 \left( \frac{1}{2} + \frac{1}{2} \cos 2Y \right) + \sin \lambda t_1 \sin \lambda t_2 \\
&\quad \cdot E\left( \frac{1}{2} - \frac{1}{2} \cos 2Y \right) - \frac{1}{2} \sin \lambda(t_1 + t_2) E(\sin 2Y) \tag{4}
\end{aligned}$$

Given:  $\phi(2) = 0$

i.e.,  $E\{\cos 2Y + i \sin 2Y\} = 0$

$$\therefore E(\cos 2Y) = 0 = E(\sin 2Y) \tag{5}$$

Using (5) in (4), we get

$$\begin{aligned} R(t_1, t_2) &= E\{X(t_1) \times X(t_2)\} = \frac{1}{2}\{\cos \lambda t_1 \cos \lambda t_2 + \sin \lambda t_1 \sin \lambda t_2\} \\ &= \frac{1}{2}\cos \lambda(t_1 - t_2) \end{aligned} \quad (6)$$

From (3) and (6), it follows that  $\{X(t)\}$  is a WSS process.

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### Example 5

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In the fair coin experiment, we define the process  $\{X(t)\}$  as follows:

$$\begin{aligned} X(t) &= \sin \pi t, \text{ if head shows, and} \\ &= 2t, \text{ if tail shows.} \end{aligned}$$

(a) Find  $E\{X(t)\}$  and (b) find  $F(x, t)$  for  $t = 0.25$

(a) The probability distribution of  $X(t)$  is given by

$$P\{X(t) = \sin \pi t\} = \frac{1}{2} \text{ and } P\{X(t) = 2t\} = \frac{1}{2}$$

$$\therefore E\{X(t)\} = \frac{1}{2} \sin \pi t + t$$

$$(b) \text{ When } t = 0.25, P\left\{X(t) = \frac{1}{2}\right\} = \frac{1}{2} \text{ and } P\left\{X(t) = \frac{1}{\sqrt{2}}\right\} = \frac{1}{2}$$

$\therefore F(x, 0.25)$  is given by

$$\begin{aligned} F(x, 0.25) &= 0, \text{ if } x < \frac{1}{2} \\ &= \frac{1}{2}, \text{ if } \frac{1}{2} \leq x < \frac{1}{\sqrt{2}} \\ &= 1, \text{ if } \frac{1}{\sqrt{2}} \leq x \end{aligned}$$

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### Example 6

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If  $\{X(t)\}$  is a wide-sense stationary process with autocorrelation  $R(\tau) = Ae^{-\alpha|\tau|}$ , determine the second-order moment of the RV  $X(8) - X(5)$ .

Second moment of  $X(8) - X(5)$  is given by

$$E[\{X(8) - X(5)\}^2] = E\{X^2(8)\} + E\{X^2(5)\} - 2E\{X(8)X(5)\} \quad (1)$$

Given:  $R(\tau) = Ae^{-\alpha|\tau|}$

$$\text{i.e., } R(t_1, t_2) = Ae^{-\alpha|t_1 - t_2|}$$

$$\therefore E\{X^2(t)\} = R(t, t) = A$$

$$\therefore E\{X^2(8)\} = E\{X^2(5)\} = A \quad (2)$$

$$\text{Also, } E\{X(8) \times X(5)\} = R(8, 5) = Ae^{-3\alpha} \quad (3)$$

Using (2) and (3) in (1), we get

$$E[\{X(8) - X(5)\}^2] = 2A(1 - e^{-3\alpha})$$

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### Example 7

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Show that the process  $X(t) = A \cos \lambda t + B \sin \lambda t$  (where  $A$  and  $B$  are RVs) is wide-sense stationary, if

- (a)  $E(A) = E(B) = 0$ ,
- (b)  $E(A^2) = E(B^2)$ , and
- (c)  $E(AB) = 0$ .

$$E\{X(t)\} = \cos \lambda t \times E(A) + \sin \lambda t \times E(B) \quad (1)$$

If  $\{X(t)\}$  is to be a WSS process,  $E\{X(t)\}$  must be a constant (i.e., independent of  $t$ ).

In (1), if  $E(A)$  and  $E(B)$  are any constants other than zero,  $E\{X(t)\}$  will be a function of  $t$ .

$$\therefore E(A) = E(B) = 0$$

$$\begin{aligned} R(t_1, t_2) &= E\{X(t_1) \times X(t_2)\} \\ &= E\{(A \cos \lambda t_1 + B \sin \lambda t_1)(A \cos \lambda t_2 + B \sin \lambda t_2)\} \\ &= E(A^2) \cos \lambda t_1 \cos \lambda t_2 + E(B^2) \sin \lambda t_1 \sin \lambda t_2 \\ &\quad + E(AB) \sin \lambda(t_1 + t_2) \end{aligned} \quad (2)$$

If  $\{X(t)\}$  is to be a WSS process,  $R(t_1, t_2)$  must be a function of  $(t_1 - t_2)$ .

$$\therefore \text{In (2), } E(AB) = 0 \text{ and } E(A^2) = E(B^2) = k$$

$$\text{Then } R(t_1, t_2) = k \cos \lambda(t_1 - t_2)$$

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### Example 8

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If the  $2n$  RVs  $A_r$  and  $B_r$  are correlated with zero mean and  $E(A_r^2) = E(B_r^2) = \sigma_r^2$ , show that the process  $X(t) = \sum_{r=1}^n (A_r \cos \omega_r t + B_r \sin \omega_r t)$  is wide-sense stationary.

What are the mean and autocorrelation of  $X(t)$ ?

$$E\{X(t)\} = \sum_{r=1}^n E(A_r) \cos \omega_r t + E(B_r) \sin \omega_r t = 0$$

$$E\{X(t_1) \times X(t_2)\} = E \left\{ \sum_{r=1}^n \sum_{s=1}^n (A_r \cos \omega_r t_1 + B_r \sin \omega_r t_1) \times (A_s \cos \omega_s t_2 + B_s \sin \omega_s t_2) \right\}$$

Since  $E\{A_r A_s\}$ ,  $E\{B_r B_s\}$ ,  $E\{A_r B_s\}$  and  $E\{B_r A_s\}$  are all zero, for  $r \neq s$ , we have

$$E\{X(t_1) \times X(t_2)\} = \sum_{r=1}^n [E(A_r^2) \cos \omega_r t_1 \cos \omega_r t_2 + E(B_r^2) \sin \omega_r t_1 \sin \omega_r t_2]$$

$$= \sum_{r=1}^n \sigma_r^2 \cos \omega_r (t_1 - t_2)$$

Therefore,  $\{X(t)\}$  is a WSS process.

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### Example 8

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Given a RV  $\Omega$  with density  $f(\omega)$  and another RV  $\phi$  uniformly distributed in  $(-\pi, \pi)$  and independent of  $\Omega$  and  $X(t) = a \cos (\Omega t + \phi)$ , prove that  $\{X(t)\}$  is a WSS process.

Recall that  $E\{g(X, Y)\} = E[E\{g(X, Y)/X\}]$

$$\begin{aligned} \therefore E\{X(t)\} &= E\{a \cos (\Omega t + \phi)\} \\ &= aE[\cos (\Omega t + \phi)/\Omega] \\ &= aE[\cos \Omega t \times E(\cos \phi) - \sin \Omega t \times E(\sin \phi)] \\ &= aE\left[\cos \Omega t \int_{-\pi}^{\pi} \frac{1}{2\pi} \cos \phi d\phi - \sin \Omega t \int_{-\pi}^{\pi} \frac{1}{2\pi} \sin \phi d\phi\right] \\ &= 0 \quad [\text{since } \phi \text{ is uniform in } (-\pi, \pi)] \end{aligned}$$

$$\begin{aligned} R(t_1, t_2) &= E\{X(t_1) \times X(t_2)\} \\ &= E\{a^2 \cos (\Omega t_1 + \phi) \cos (\Omega t_2 + \phi)\} \\ &= a^2 E[\cos \Omega t_1 \cos \Omega t_2 \cos^2 \phi + \sin \Omega t_1 \sin \Omega t_2 \sin^2 \phi \\ &\quad - (\sin \Omega t_1 \cos \Omega t_2 + \cos \Omega t_1 \sin \Omega t_2) \sin \phi \cos \phi / \Omega] \\ &= a^2 E\left[\cos \Omega t_1 \cos \Omega t_2 \int_{-\pi}^{\pi} \cos^2 \phi d\phi + \sin \Omega t_1 \sin \Omega t_2 \times \right. \\ &\quad \left. \int_{-\pi}^{\pi} \sin^2 \phi d\phi - \sin \Omega(t_1 + t_2) \int_{-\pi}^{\pi} \sin \phi \cos \phi d\phi\right] \end{aligned}$$

$$= \frac{1}{2} a^2 E\{\cos \Omega t_1 \cos \Omega t_2 + \sin \Omega t_1 \sin \Omega t_2\}$$

$$= \frac{1}{2} a^2 E\{\cos \Omega(t_1 - t_2)\}$$

= a function of  $(t_1 - t_2)$ , whatever be the value of  $f(\omega)$ .

Therefore,  $\{X(t)\}$  is a WSS process.

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### Example 10

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Verify whether the sine wave process  $\{X(t)\}$ , where  $X(t) = Y \cos \omega t$ , where  $Y$  is uniformly distributed in  $(0, 1)$  is a SSS process.

$$P\{X(t) \leq x\} = P\{Y \cos \omega t \leq x\}$$

$$\begin{aligned}
 &= \begin{cases} P\left\{Y \leq \frac{x}{\cos \omega t}\right\} & \text{if } \cos \omega t > 0 \\ P\left\{Y \geq \frac{x}{\cos \omega t}\right\} & \text{if } \cos \omega t < 0 \end{cases} \\
 \text{i.e., } F_{X(t)}(x) &= \begin{cases} F_Y\left(\frac{x}{\cos \omega t}\right) & \text{if } \cos \omega t > 0 \\ 1 - F_Y\left(\frac{x}{\cos \omega t}\right) & \text{if } \cos \omega t < 0 \end{cases} \\
 \therefore f_{X(t)}(x) &= \frac{1}{|\cos \omega t|} \times f_Y\left(\frac{x}{\cos \omega t}\right) \\
 &= \text{a function of } t
 \end{aligned}$$

If  $\{X(t)\}$  is to be a SSS process, its first-order density must be independent of  $t$ . Therefore,  $\{X(t)\}$  is not a SSS process.

### Example 11

If  $X(t) = Y \cos \omega t + Z \sin \omega t$ , where  $Y$  and  $Z$  are two independent normal RVs with  $E(Y) = E(Z) = 0$ ,  $E(Y^2) = E(Z^2) = \sigma^2$  and  $\omega$  is a constant, prove that  $\{X(t)\}$  is a SSS process of order 2.

Since  $X(t)$  is a linear combination of  $Y$  and  $Z$ , that are independent,  $X(t)$  follows a normal distribution with

$$E\{X(t)\} = \cos \omega t E(Y) + \sin \omega t E(Z) = 0$$

$$\begin{aligned}
 \text{and } \text{Var}\{X(t)\} &= \cos^2 \omega t \times E(Y^2) + \sin^2 \omega t \times E(Z^2) \\
 &= \sigma^2
 \end{aligned}$$

Since  $X(t_1)$  and  $X(t_2)$  are each  $N(0, \sigma)$ ,  $X(t_1)$  and  $X(t_2)$  are jointly normal with the joint pdf given by

$$f(x_1, x_2, t_1, t_2) = \frac{1}{2\pi\sigma^2\sqrt{1-r^2}} \exp\left\{\frac{-(x_1^2 - 2rx_1 x_2 + x_2^2)}{2(1-r^2)\sigma^2}\right\}; -\infty < x_1, x_2 < \infty \quad (1)$$

In (1),  $r = \text{correlation co-efficient between } X(t_1) \text{ and } X(t_2)$

$$\begin{aligned}
 r &= \frac{C(t_1, t_2)}{\sqrt{\text{Var}\{X(t_1)\} \times \text{Var}\{X(t_2)\}}} \\
 &= \frac{1}{\sigma^2} E\{X(t_1) \times X(t_2)\} \\
 &= \frac{1}{\sigma^2} E[Y \cos \omega t_1 + Z \sin \omega t_1] (Y \cos \omega t_2 + Z \sin \omega t_2)
 \end{aligned}$$

$$\begin{aligned}
 &= \frac{1}{\sigma^2} [E(Y^2) \cos \omega t_1 \cos \omega t_2 + E(Z^2) \sin \omega t_1 \sin \omega t_2] \\
 &\quad [\text{since } E(YZ) = 0, \text{ as } Y \text{ and } Z \text{ are independent}] \\
 &= \cos \omega(t_1 - t_2)
 \end{aligned}$$

Now, the joint pdf of  $X(t_1 + h)$  and  $X(t_2 + h)$  is given by a similar expression as in (1), where

$$\begin{aligned}
 r &= \cos \omega \{(t_1 + h) - (t_2 + h)\} \\
 &= \cos \omega(t_1 - t_2)
 \end{aligned}$$

Thus, the joint pdf's of  $\{X(t_1), X(t_2)\}$  and  $\{X(t_1 + h), X(t_2 + h)\}$  are the same. Therefore,  $\{X(t)\}$  is a SSS process of order 2.

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### Example 12

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Two random processes  $X(t)$  and  $Y(t)$  are defined by  $X(t) = A \cos \omega_0 t + B \sin \omega_0 t$  and  $Y(t) = B \cos \omega_0 t - A \sin \omega_0 t$ . Show that  $X(t)$  and  $Y(t)$  are jointly wide-sense stationary, if  $A$  and  $B$  are uncorrelated RVs with zero means and the same variances and  $\omega_0$  is a constant.

$$\begin{aligned}
 E(A) &= E(B) = 0; \text{Var}(A) = \text{Var}(B) \\
 \therefore E(A^2) &= E(B^2)
 \end{aligned}$$

Since  $A$  and  $B$  are uncorrelated,  $E(AB) = 0$ .

Therefore, by Example 7,  $\{X(t)\}$  and  $\{Y(t)\}$  are individually WSS processes.

$$\begin{aligned}
 \text{Now, } R_{xy}(t_1, t_2) &= E\{X(t_1) \times Y(t_2)\} \\
 &= E\{(A \cos \omega_0 t_1 + B \sin \omega_0 t_1)(B \cos \omega_0 t_2 - A \sin \omega_0 t_2)\} \\
 &= E(B^2) \sin \omega_0 t_1 \cos \omega_0 t_2 - E(A^2) \cos \omega_0 t_1 \sin \omega_0 t_2 \\
 &= \sigma^2 \sin \omega_0(t_1, t_2) \quad [\text{assuming } E(A^2) = E(B^2) = \sigma^2] \\
 &= \text{a function of } (t_1 - t_2)
 \end{aligned}$$

Therefore,  $\{X(t)\}$  and  $\{Y(t)\}$  are jointly WSS processes.

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### Example 13

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(a) Define Random Walk and prove that the limiting form of random walk is the Wiener process, and (b) find the covariance function of the Wiener process.

**(a) Definition of Random Walk**

Suppose we toss a fair coin every  $T$  seconds and instantly after each toss we move a distance  $d$  to the right if heads show and to the left if tails show. If the process starts at  $t = 0$ , our position at time  $t = nT$  is a random sequence  $X(nT)$  that may be simply denoted as  $X(n)$  also. The process  $\{X(nT)\}$  is called a random walk.

Suppose that  $r$  heads and  $(n - r)$  tails have occurred in the first  $n$  tosses of the coin. Then the random walk consists of  $r$  steps to the right and  $(n - r)$  steps to the left.

$$\begin{aligned}\therefore X(nT) &= rd - (n-r)d \\ &= (2r-n)d = md, \text{ say}\end{aligned}$$

Note that  $X(nT)$  is a RV, taking the values  $md$ , where  $m = -n, -n+2, \dots, n-2, n$ .

Now,  $P\{X(nT) = md\} = P\{\text{getting } r \text{ heads in } n \text{ tosses}\}$

$$= nC_r \left(\frac{1}{2}\right)^r, \text{ where } r = \frac{m+n}{2}, \text{ since } 2r-n=m$$

Now,  $X(nT)$  can also be expressed as a sum as given below,

$$X(nT) = X_1 + X_2 + \dots + X_n$$

where  $X_i$  represents the distance moved in the  $i$ th step.

The RVs  $X_i$  are independent, taking the values  $\pm d$  with equal probability.

$$\therefore E\{X(nT)\} = \sum_{i=1}^n E(X_i) = 0$$

$$\begin{aligned}\text{and } E\{X^2(nT)\} &= \sum_{i=1}^n E(X_i^2) \\ &= \sum_{i=1}^n \left( \frac{1}{2} \times d^2 + \frac{1}{2} \times d^2 \right) = nd^2\end{aligned}$$

We know that limiting form of the binomial distribution with mean  $np$  and variance  $npq$  as  $n \rightarrow \infty$  is the normal distribution  $N(np, \sqrt{npq})$ .

$$\begin{aligned}\text{i.e. } nC_r p^r q^{n-r} &\cong \frac{1}{\sqrt{2\pi npq}} e^{-(r-np)^2/2npq} \\ \therefore P\{X(nT) = md\} &\quad \frac{1}{\sqrt{2\pi \frac{n}{4}}} e^{-m^2/2n} \left( \text{since } r = \frac{m+n}{2} \right) \quad (1)\end{aligned}$$

### Wiener Process as Limiting form of Random Walk

In (1), put  $nT = t$ ,  $md = x$  and  $d^2 = \alpha T$  and take limits as  $T \rightarrow 0$  and  $n \rightarrow \infty$ .

In the limit,  $\{X(t)\}$  becomes a continuous process.

$$\text{Now, } \frac{m}{\sqrt{n}} = \frac{x/d}{\sqrt{t/T}} = \frac{x}{\sqrt{d^2 t/T}} = \frac{x}{\sqrt{\alpha t}} \quad (2)$$

$$\begin{aligned}\text{Also, } \frac{n}{4} &= E\{X^2(nT)\} \\ &= nd^2 \\ &= n\alpha t \\ &= \alpha t\end{aligned} \quad (3)$$

If we use (2) and (3) in (1), when we proceed to limits,

$$P\{x \leq X(t) \leq x + dx\} = \frac{1}{\sqrt{2\pi\alpha t}} e^{-x^2/2\alpha t} dx, -\infty < x < \infty$$

i.e., the pdf of Wiener process  $\{X(t)\}$  is

$$f_{X(t)}(x) = \frac{1}{\sqrt{2\pi\alpha t}} e^{-x^2/2\alpha t}$$

which is  $N(0, \sqrt{\alpha t})$ .

(b) Obviously, the random walk  $\{X(nT)\}$  is a process with independent increments.

i.e.,  $\{X(n_2 T) - X(n_1 T)\}$  and  $\{X(n_1 T) - X(0)\}$  are independent.

Since Wiener process  $\{X(t)\}$  is the limiting form of random walk,  $\{X(t_2) - X(t_1)\}$  and  $X(t_1)$  are independent.

Let  $t_1 < t_2$

Then  $E[\{X(t_2) - X(t_1)\} \times X(t_1)]$

$$\begin{aligned} &= E[\{X(t_2) - X(t_1)\}] \times E[X(t_1)] \\ &= 0 \end{aligned} \quad [\text{since } E[X(t)] = 0]$$

i.e.,  $E\{X(t_1) \times X(t_2)\} = E\{X^2(t_1)\}$

$$= \alpha t_1 \quad [\text{since } \text{Var } \{X(t)\} = \alpha t]$$

i.e.,  $R(t_1, t_2) = \alpha t_1$ .

Similarly, when  $t_2 < t$ ,  $R(t_1, t_2) = \alpha t_2$ .

$\therefore R(t_1, t_2) = \alpha \min(t_1, t_2)$

$$\begin{aligned} C(t_1, t_2) &= R(t_1, t_2) - \mu(t_1) \times \mu(t_2) \\ &= \alpha \min(t_1, t_2) \end{aligned} \quad [\text{since } \mu(t) = 0]$$

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### Example 14

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If  $X(t)$  with  $X(0) = 0$  and  $\mu = 0$  is a Wiener process, show that  $Y(t) = \sigma X\left(\frac{t}{\sigma^2}\right)$

is also a *Wiener process*. Find its covariance function.

The pdf of Wiener process  $\{X(t)\}$  is given by

$$f(x) = \frac{1}{\sqrt{2\pi\alpha t}} e^{-x^2/2\alpha t}$$

The standard deviation of  $X(t)$  is given by

$$\sigma = \sqrt{\alpha t}$$

Therefore, the pdf of the RV  $\left\{X\left(\frac{t}{\sigma^2}\right)\right\} = \left\{X\left(\frac{1}{\alpha}\right)\right\}$  is

$$f_1(x) = \frac{1}{\sqrt{2\pi}} e^{-x^2/2}$$

To find the pdf of  $Y(t) = \sigma X\left(\frac{t}{\sigma^2}\right)$ , let us use the transformation rule.

$$\begin{aligned} f_{Y(t)}(y) &= f_1(x) \left| \frac{dx}{dy} \right| \text{ where } y = \sigma x \\ &= \frac{1}{\sigma} f_1\left(\frac{y}{\sigma}\right) \\ &= \frac{1}{\sqrt{\sigma t}} f_1\left(\frac{y}{\sqrt{\sigma t}}\right) \\ &= \frac{1}{\sqrt{2\pi\sigma t}} e^{-y^2/2\sigma t} \end{aligned}$$

which is the pdf of the Wiener process.

The covariance function of the Wiener process was found out in the previous example.

### Example 15

Define semi-random telegraph signal process and random telegraph signal process and prove also that the former is evolutionary and the latter is wide-sense stationary.

If  $N(t)$  represents the number of occurrences of a specified event in  $(0, t)$  and  $X(t) = (-1)^{N(t)}$ , then  $\{X(t)\}$  is called a *semi-random telegraph signal process*.

If  $\{X(t)\}$  is a semi-random telegraph signal process,  $\alpha$  is a RV which is independent of  $X(t)$  and which assumes the values +1 and -1 with equal probability and  $Y(t) = \alpha X(t)$ , then  $\{Y(t)\}$  is called a *random telegraph signal process*.

It will be proved in Chapter 7 that the distribution of  $N(t)$  is Poisson with mean  $\lambda t$ , where the probability of exactly one occurrence in a small interval of length  $h = \lambda h$ .

In other words, the process  $\{N(t)\}$  is a Poisson process with the probability law

$$P\{N(t) = r\} = e^{-\lambda t} (\lambda t)^r / r! ; r = 0, 1, 2, \dots, \infty$$

If  $\{X(t)\}$  is the semi-random telegraph signal process, then as per the definition given above,  $X(t)$  can take the values +1 and -1 only.

$$\begin{aligned} P\{X(t) = 1\} &= P\{N(t) \text{ is even}\}. \\ &= P\{N(t) = 0\} + P\{N(t) = 2\} + P\{N(t) = 4\} \\ &\quad + \dots + \infty \text{ (since the events are mutually exclusive)} \\ &= e^{-\lambda t} \left\{ 1 + \frac{(\lambda t)^2}{2} + \frac{(\lambda t)^4}{4} + \dots + \infty \right\} \end{aligned}$$

$$\begin{aligned}
 &= e^{-\lambda t} \cosh \lambda t \\
 P\{X(t) = -1\} &= P\{N(t) \text{ is odd}\} \\
 &= P\{N(t) = 1\} + P\{N(t) = 3\} + \dots + \infty \\
 &\quad (\text{since the events are mutually exclusive}) \\
 &= e^{-\lambda t} \left\{ \frac{\lambda t}{1} + \frac{(\lambda t)^3}{3} + \dots + \infty \right\} \\
 &= e^{-\lambda t} \sinh \lambda t \\
 \therefore E\{X(t)\} &= 1 \times e^{-\lambda t} \cosh \lambda t + (-1) \times e^{-\lambda t} \sinh \lambda t \\
 &= e^{-\lambda t} e^{-\lambda t} = e^{-2\lambda t}
 \end{aligned}$$

To find  $E\{X(t_1) \times X(t_2)\}$ , we require the joint probability distribution of  $\{X(t_1), X(t_2)\}$ .

$$\begin{aligned}
 \text{Now, } P\{X(t_1) = 1, X(t_2) = 1\} &= P\{X(t_1) = 1/X(t_2) = 1\} \times P\{X(t_2) = 1\} \\
 &= P\{\text{an even number of occurrences of the event in } (t_1 - t_2)\} \\
 &\quad \times P\{X(t_2) = 1\} \\
 &= e^{-\lambda\tau} \cosh \lambda\tau \times e^{-\lambda t_2} \cosh \lambda t_2; \text{ where } \tau = t_1 - t_2
 \end{aligned}$$

$$\text{Similarly, } P\{X(t_1) = -1, X(t_2) = -1\}$$

$$= e^{-\lambda\tau} \cosh \lambda\tau e^{-\lambda/2} \sinh \lambda t_2$$

$$P\{X(t_1) = -1, X(t_2) = 1\} = e^{-\lambda\tau} \sinh \lambda\tau e^{-\lambda t_2} \sinh \lambda t_2$$

$$\text{and } P\{X(t_1) = -1, X(t_2) = 1\} = e^{-\lambda\tau} \sinh \lambda\tau e^{-\lambda t_2} \cosh \lambda t_2$$

$$\text{Now, } X(t_1) \times X(t_2) = 1, \text{ if } \{X(t_1) = 1 \text{ and } X(t_2) = 1\} \text{ or } \{X(t_1) = -1, X(t_2) = -1\}$$

$$\begin{aligned}
 \therefore P\{X(t_1) \times X(t_2) = 1\} &= e^{-\lambda(\tau+t_2)} \cosh \lambda\tau (\cosh \lambda t_2 + \sinh \lambda t_2) \\
 &= e^{-\lambda t} \cosh \lambda\tau
 \end{aligned}$$

$$\text{and } P\{X(t_1) \times X(t_2) = -1\}$$

$$\begin{aligned}
 &= e^{-\lambda(\tau+t_2)} \sin \lambda\tau (\cosh \lambda t_2 + \sinh \lambda t_2) \\
 &= e^{-\lambda\tau} \sinh \lambda\tau
 \end{aligned}$$

$$\begin{aligned}
 \therefore R(t_1, t_2) &= E\{X(t_1) X(t_2)\} \\
 &= 1 \times e^{-\lambda\tau} \cosh \lambda\tau - 1 \times e^{-\lambda\tau} \sinh \lambda\tau \\
 &= e^{-2\lambda\tau} \\
 &= e^{-2\lambda(t_1-t_2)}
 \end{aligned}$$

Although  $R(t_1, t_2)$  is a function of  $(t_1 - t_2)$ ,  $E\{X(t)\}$  is not a constant.

Therefore,  $\{X(t)\}$  is *evolutionary*.

Let us now consider the random telegraph signal process  $\{Y(t)\}$ , where  $Y(t) = \alpha X(t)$ .

By definition,  $P\{\alpha = 1\} = \frac{1}{2}$  and  $P(\alpha = -1) = \frac{1}{2}$

$$\therefore E(\alpha) = 0 \text{ and } E(\alpha^2) = 1$$

Now,  $E\{Y(t)\} = E(\alpha) \times E\{X(t)\}$

$= 0$  [since  $\alpha$  and  $X(t)$  are independent]

$$\begin{aligned} E\{Y(t_1) \times Y(t_2)\} &= E\{\alpha^2 X(t_1) \times X(t_2)\} \\ &= E\{\alpha^2\} \times E\{X(t_1) \times X(t_2)\} \quad (\text{by independence}) \\ &= 1 \times e^{-2\lambda(t_1-t_2)} \end{aligned}$$

i.e.,  $R_{yy}(t_1, t_2) = \text{a function of } (t_1 - t_2)$

Therefore,  $\{Y(t)\}$  is a wide-sense stationary process.

### Exercise 6(A)

#### Part-A (Short-answer Questions)

1. What is the difference between a RV and a random process?
2. Define a random process and give an example of a random process.
3. Explain the terms ‘state space’ and ‘parameter set’ associated with a random process.
4. If  $\{X(s, t)\}$  is a random process, what is the nature of  $X(s, t)$  when (i)  $s$  is fixed and (ii)  $t$  is fixed?
5. What is the difference between a random sequence and random process?
6. What is a discrete random sequence? Give an example.
7. What is a continuous random sequence? Give an example.
8. What is a discrete random process? Give an example.
9. What is a continuous random process? Give an example.
10. How is a random process described mathematically?
11. Name 3 classes of RP's into which RP's are generally divided.
12. Name 2 important RP's with independent increments.
13. What do you mean by the mean and variance of a random process?
14. Define the autocorrelation of a RP  $\{X(t)\}$ .
15. Define the autocovariance of a RP  $\{X(t)\}$ .
16. Define the correlation coefficient of a RP  $\{X(t)\}$ .
17. Is the autocorrelation of a RP the same as the correlation coefficient of the process? Why?
18. Define the cross-correlation of 2 random processes.
19. When are 2 random processes said to be orthogonal?
20. Define the cross-covariance of 2 random processes.
21. Define the cross-correlation coefficient of 2 random processes.
22. Define a strict-sense stationary process and give an example.
23. Define a  $k$ th-order stationary process. When will it becomes a SSS process?

24. Prove that the first-order density function of a SSS process  $\{X(t)\}$  is independent of  $t$ .
25. If  $\{X(t)\}$  is a SSS process, prove that  $E\{X(t)\}$  is a constant.
26. If  $\{X(t)\}$  is a SSS process, prove that the joint pdf of  $X(t_1)$  and  $X(t_2)$  is a function of  $(t_1 - t_2)$ .
27. Prove that the autocorrelation of a SSS process  $\{X(t)\}$  is a function of  $(t_1 - t_2)$ .
28. When are  $\{X(t)\}$  and  $\{Y(t)\}$  said to be jointly stationary in the strict sense?
29. When is a complex random process  $\{Z(t)\}$ , where  $Z(t) = X(t) + iY(t)$ , said to be a SSS process?
30. Define wide-sense stationary process. Give an example.
31. What is the difference between a SSS process and a WSS process?
32. When is a random process said to be evolutionary? Give an example of an evolutionary process.
33. When are the processes  $\{X(t)\}$  and  $\{Y(t)\}$  said to be jointly stationary in the wide sense?
34. If  $\{X(t)\}$  is a stationary process in any sense, prove that  $\text{Var}\{X(t)\}$  is a constant.
35. Give the one-dimensional density function of Wiener process. What are its mean and variance?
36. Define a semi-random telegraph signal process. Is it stationary?
37. Define a random telegraph signal process. Is it stationary?

#### Part-B

38. If  $X(t) = P + Qt$ , where  $P$  and  $Q$  are independent RVs with  $E(P) = p$ ,  $E(Q) = q$ ,  $\text{Var}(P) = \sigma_1^2$  and  $\text{Var}(Q) = \sigma_2^2$ , find  $E\{X(t)\}$ ,  $R(t_1, t_2)$  and  $C(t_1, t_2)$ . Is the process  $\{X(t)\}$  stationary?
39. If  $X(t) = \sin(\omega t + Y)$ , where  $Y$  is uniformly distributed in  $(0, 2\pi)$ , prove that  $\{X(t)\}$  is a wide-sense stationary process.
40. If  $X(t) = Y \cos t + Z \sin t$  for all  $t$  where  $Y$  and  $Z$  are independent binary RVs, each of which assumes the values  $-1$  and  $+2$  with probabilities  $2/3$  and  $1/3$  respectively, prove that  $\{X(t)\}$  is wide-sense stationary.
41. Calculate the autocorrelation function of the process  $X(t) = A \sin(\omega_0 t + \phi)$ , where  $A$  and  $\omega_0$  are constants and  $\phi$  is a uniformly distributed RV in  $(0, 2\pi)$ .
42. Consider the random process  $V(t) = \cos(\omega t + \theta)$ , where  $\theta$  is a RV with probability density

$$P(\theta) = \begin{cases} \frac{1}{2\pi} & -\pi \leq \theta \leq \pi \\ 0 & \text{elsewhere} \end{cases}$$

- (i) Show that the first and second moments of  $V(t)$  are independent of time.
- (ii) If  $\theta = \text{constant}$ , will the ensemble mean of  $V(t)$  be time-independent?

43. If  $X(t) = R \cos(\omega t + \phi)$ , where  $R$  and  $\phi$  are independent RVs and  $\phi$  is uniformly distributed in  $(-\pi, \pi)$ , prove that  $R(t_1, t_2) = \frac{1}{2} E(R^2) \times \cos \omega(t_1 - t_2)$ .
44. In the fair-coin experiment, we define the process  $X(t)$  as follows:
- $$X(t) = \sin \pi t, \text{ if head shows and}$$
- $$X(t) = 2t, \text{ if tail shows.}$$
- Find (i)  $E\{X(t)\}$  and (ii)  $F(x, t)$  for  $t = 0.25, 0.5, 1$ .
45. A stochastic process is described by  $X(t) = A \sin t + B \cos t$ , where  $A$  and  $B$  are independent RVs with zero means and equal standard deviations. Show that process is stationary of the second order.  
**[Hint:** Without knowing the distributions of  $A$  and  $B$ , it is not possible to prove strict sense stationarity of the second order. We can prove that  $\{X(t)\}$  is WSS.]
46. Consider a random process  $Z(t) = X_1 \cos \omega_0 t - X_2 \sin \omega_0 t$ , where  $X_1$  and  $X_2$  are independent Gaussian RVs with zero mean and variance  $\sigma^2$ , find  $E\{z\}$  and  $E\{z^2\}$ .
47. Suppose that  $X(t)$  is a process with mean  $\mu(t) = 3$  and autocorrelation  $R(t_1, t_2) = 9 + 4 e^{-0.2|t_1 - t_2|}$ . Determine the mean, variance and the covariance of the RVs  $Z = X(5)$  and  $W = X(8)$ .
48. If the RVs  $A_i$  are uncorrelated with zero mean and  $E\{|A_i|^2\} = \sigma_i^2$ , prove that the process  $X(t) = \sum_{i=1}^n A_i e^{j\omega_i t}$  is wide-sense stationary with zero mean. Show also that for  $X(t)$ ,  $R(z) = \sum_{i=1}^n \sigma_i^2 e^{j\omega_i z}$ .
- [Hint:** For a complex-valued random process  $\{X(t)\}$ , the autocorrelation is defined as  $R(t_1, t_2) = E\{X(t_1) X^*(t_2)\}$ , where  $X^*(t_2)$  is the complex conjugate of  $X(t_2)$ .]
49. If  $U(t) = X \cos t + Y \sin t$  and  $V(t) = Y \cos t + X \sin t$ , where  $X$  and  $Y$  are independent RVs such that  $E(X) = 0 = E(Y)$ ,  $E(X^2) = E(Y^2) = 1$ , show that  $\{U(t)\}$  and  $\{V(t)\}$  are individually stationary in the wide sense, but they are not jointly wide-sense stationary.
50. If  $X(t) = 5 \cos(10t + \theta)$  and  $Y(t) = 20 \sin(10t + \theta)$ , where  $\theta$  is a RV uniformly distributed in  $(0, 2\pi)$ , prove that the processes  $\{X(t)\}$  and  $\{Y(t)\}$  are jointly wide-sense stationary.
51. If  $X(t) = A \sin(\omega t + \theta)$  where  $A$  and  $\omega$  are constants and  $\theta$  is RV, uniformly distributed over  $(-\pi, \pi)$ , find the autocorrelation of  $\{Y(t)\}$ , where  $Y(t) = X^2(t)$ .
52. Find the first-order characteristic function of a Wiener process.

53. Consider the process  $W(t) = X(t) \cos \omega t + Y(t) \sin \omega t$ , where  $X(t)$  and  $Y(t)$  are 2 real jointly stationary processes. What are the conditions for  $W(t)$  to be WSS? In case  $W(t)$  is wide-sense stationary, what is its autocorrelation of  $X(t)$  and  $Y(t)$ ?
54. Show that, if the process  $X(t) = a \cos \omega t + b \sin \omega t$  is SSS, where  $a$  and  $b$  are independent RVs, then they are normal.
55. The RVs  $A$  and  $B$  are independent  $N(0, \sigma)$  and  $p$  is the probability that the process  $\{X(t)\}$ , where  $X(t) = A - Bt$ , crosses the  $t$ -axis in the interval  $(0, T)$ . Show that  $\pi p = \tan^{-1} T$ .

[Hint:  $Z = A/B$  follows the Cauchy distribution with  $f(z) = \frac{1}{\pi} \times \frac{1}{1+z^2}$  and  $p = P\{0 \leq A/B \leq T\}$ ].

## Autocorrelation Function and its Properties

**Definition:** If the process  $\{X(t)\}$  is stationary either in the strict sense or in the wide sense, then  $E\{X(t) \times X(t - \tau)\}$  is a function of  $\tau$ , denoted by  $R_{xx}(\tau)$  or  $R(\tau)$  or  $R_x(\tau)$ . This function  $R(\tau)$  is called the autocorrelation function of the process  $\{X(t)\}$ .

## Properties of $R(\tau)$

1.  $R(t)$  is an even function of  $\tau$

### Proof

$$\begin{aligned} R(\tau) &= E\{X(t) \times X(t - \tau)\} \\ \therefore R(-\tau) &= E\{X(t) \times X(t + \tau)\} \\ &= E\{X(t + \tau) \times X(t)\} \\ &= R(\tau) \end{aligned}$$

Therefore,  $R(\tau)$  is an even function of  $\tau$ .

2.  $R(\tau)$  is maximum at  $\tau = 0$

i.e.,  $|R(\tau)| \leq R(0)$

### Proof

The Cauchy–Schwarz inequality is

$$\{E(XY)\}^2 \leq E(X^2) \times E(Y^2)$$

Put  $X = X(t)$  and  $Y = X(t - \tau)$

Then  $[E\{X(t) \times X(t - \tau)\}]^2 \leq E\{X^2(t)\} \times E\{X^2(t - \tau)\}$

i.e.,  $\{R(\tau)\}^2 \leq [E\{X^2(t)\}]^2$

[since  $E\{X(t)\}$  and  $\text{Var}\{X(t)\}$  are constant for a stationary process]

i.e.,  $\{R(\tau)\}^2 \leq \{R(0)\}^2$

Taking square-root on both sides

$$|R(\tau)| \leq R(0) \quad [\text{since } R(0) = E\{X^2(t)\} \text{ is positive}]$$

**3. If the autocorrelation function  $R(t)$  of a real stationary process  $\{X(t)\}$  is continuous at  $\tau = 0$ , it is continuous at every other point.**

### Proof

Consider

$$\begin{aligned} E[\{X(t) - X(t - \tau)\}^2] &= E[\{X^2(t)\} + E\{X^2(t - \tau)\} - 2E\{X(t) \times X(t - \tau)\}] \\ &= R(0) + R(0) - 2R(\tau) \\ &= 2[R(0) - R(\tau)] \end{aligned} \quad (1)$$

Since  $R(\tau)$  is continuous at  $\tau = 0$ ,  $\lim_{\tau \rightarrow 0} R(\tau) = R(0)$

i.e.,  $\lim_{\tau \rightarrow 0} \{\text{R.S. of (1)}\} = 0$

$\therefore \lim_{\tau \rightarrow 0} \{\text{L.S. of (1)}\} = 0$

$\therefore \lim_{\tau \rightarrow 0} \{X(t - \tau)\} = X(t)$

i.e.,  $X(t)$  is continuous for all  $t$  (2)

Consider  $R(\tau + h) - R(\tau)$

$$\begin{aligned} &= E[\{X(t)\} \times X\{t - (\tau + h)\}] - E\{X(t) \times X(t - \tau)\} \\ &= E[X(t) \{X(t - \tau - h) - X(t - \tau)\}] \end{aligned} \quad (3)$$

Now,  $\lim_{h \rightarrow 0} [X(t - \tau - h) - X(t - \tau)] = 0$ , by (2)

$\therefore \lim_{h \rightarrow 0} \{\text{R.S. of (3)}\} = 0$

$\therefore \lim_{h \rightarrow 0} \{\text{L.S. of (3)}\} = 0$

i.e.,  $\lim_{h \rightarrow 0} \{R(\tau + h)\} = R(\tau)$

i.e.,  $R(\tau)$  is continuous for all  $\tau$

**4. If  $R(\tau)$  is the autocorrelation function of a stationary process  $\{X(t)\}$  with no periodic component, then  $\lim_{\tau \rightarrow \infty} R(\tau) = \mu_x^2$ , provided the limit exists.**

### Proof

$$R(\tau) = E\{X(t)\} \times X(t - \tau)\}$$

When  $\tau$  is very large,  $X(t)$  and  $X(t - \tau)$  are two sample functions (members) of the process  $\{X(t)\}$  observed at a very long interval of time.

Therefore,  $X(t)$  and  $X(t - \tau)$  tend to become independent [ $X(t)$  and  $X(t - \tau)$  may be dependent, when  $X(t)$  contains a periodic component, which is not true].

$$\therefore \lim_{\tau \rightarrow \infty} \{R(\tau)\} = E\{X(t)\} \times E\{X(t - \tau)\}$$
$$= \mu_x^2 \quad [\text{since } E\{X(t)\} \text{ is a constant}]$$

i.e.,  $\mu_x = \sqrt{\lim_{\tau \rightarrow \infty} R(\tau)}$

## Cross-Correlation Function and its Properties

**Definition:** If the processes  $\{X(t)\}$  and  $\{Y(t)\}$  are jointly wide-sense stationary, then  $E\{X(t) \times Y(t - \tau)\}$  is a function of  $\tau$ , denoted by  $R_{xy}(\tau)$ . This function  $R_{xy}(\tau)$  is called the cross-correlation function of the processes  $\{X(t)\}$  and  $\{Y(t)\}$ .

We give below the properties of  $R_{xy}(\tau)$  without proof. Proofs of these properties are left as exercises to the reader.

### Properties

1.  $R_{yx}(\tau) = R_{xy}(-\tau)$
2.  $|R_{xy}(\tau)| \leq \sqrt{R_{xx}(0) \times R_{yy}(0)}$   
This means that the maximum of  $R_{xy}(\tau)$  can occur anywhere, but it cannot exceed  $\sqrt{R_{xx}(0) \times R_{yy}(0)}$ .
3.  $|R_{xy}(\tau)| \leq 1/2 \{R_{xx}(0) + R_{yy}(0)\}$
4. If the processes  $\{X(t)\}$  and  $\{Y(t)\}$  are orthogonal, then  $R_{xy}(\tau) = 0$
5. If the processes  $\{X(t)\}$  and  $\{Y(t)\}$  are independent, then  $R_{xy}(\tau) = \mu_x \times \mu_y$

## Ergodicity

When we wish to take a measurement of a variable quantity in the laboratory, we usually obtain multiple measurements of the variable and average them to reduce measurement errors. If the value of the variable being measured is constant and errors are due to disturbances (noise) or due to the instability of the measuring instrument, then averaging is, in fact, a valid and useful technique. ‘Time averaging’ is an extension of this concept, which is used in the estimation of various statistics of random processes.

We normally use ensemble averages (or statistical averages) such as the mean and autocorrelation function for characterising random processes. To estimate ensemble averages, one has to compute a weighted average over all the member functions of the random process.

For example, the ensemble mean of a discrete random process  $\{X(t)\}$  is computed by the formula  $\mu_x = \sum_i x_i p_i$ . If we have access only to a single sample function of the process, then we use its *time-average* to estimate the ensemble averages of the process.

**Definition:** If  $\{X(t)\}$  is a random process, then  $\frac{1}{2T} \int_{-T}^T X(t) dt$  is called the *time-average* of  $\{X(t)\}$  over  $(-T, T)$  and denoted by  $\bar{X}_T$ .

In general, ensemble averages and time averages are not equal except for a very special class of random processes called *ergodic processes*. The concept of ergodicity deals with the equality of time averages and ensemble averages.

**Definition:** A random process  $\{X(t)\}$  is said to be ergodic, if its ensemble averages are equal to appropriate time averages.

This definition implies that, with probability 1, any ensemble average of  $\{X(t)\}$  can be determined from a single sample function of  $\{X(t)\}$ .

**Note** Ergodicity is a stronger condition than stationarity and hence all random processes that are stationary are not ergodic. Moreover, ergodicity is usually defined with respect to one or more ensemble averages (such as mean and autocorrelation function) as discussed below and a process may be ergodic with respect to one ensemble average but not others.

## Mean-Ergodic Process

If the random process  $\{X(t)\}$  has a constant mean  $E\{X(t)\} = \mu$  and if  $\bar{X}_T =$

$\frac{1}{2T} \int_{-T}^T X(t) dt \rightarrow \mu$ , as  $T \rightarrow \infty$ , then  $\{X(t)\}$  is said to be mean-ergodic.

## Mean-Ergodic Theorem

If  $\{X(t)\}$  is a random process with constant mean  $\mu$  and if  $\bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt$ , then

$\{X(t)\}$  is mean-ergodic (or ergodic in the mean), provided  $\lim_{T \rightarrow \infty} \{\text{Var } \bar{X}_T\} = 0$ .

### Proof

$$\begin{aligned} \bar{X}_T &= \frac{1}{2T} \int_{-T}^T X(t) dt \\ \therefore E(\bar{X}_T) &= \frac{1}{2T} \int_{-T}^T E\{X(t)\} dt \\ &= \mu \end{aligned} \tag{1}$$

By Tchebycheff's inequality,

$$P\{|\bar{X}_T - E(\bar{X}_T)| \leq \epsilon\} \geq 1 - \frac{\text{Var}(\bar{X}_T)}{\epsilon^2} \tag{2}$$

Taking limits as  $T \rightarrow \infty$  and using (1) we get

$$P\left\{\left|\lim_{T \rightarrow \infty}(\bar{X}_T) - \mu\right| \leq \epsilon\right\} \geq 1 - \frac{\lim_{T \rightarrow \infty} \text{Var}(\bar{X}_T)}{\epsilon^2}$$

$\therefore$  when  $\lim_{T \rightarrow \infty} \text{Var}(\bar{X}_T) = 0$ , (2) becomes

$$P\left\{\left|\lim_{T \rightarrow \infty}(\bar{X}_T) - \mu\right| \leq \epsilon\right\} \geq 1$$

i.e.,  $\lim_{T \rightarrow \infty}(\bar{X}_T) = E\{X(t)\}$  with probability 1.

**Note** This theorem provides a sufficient condition for the mean-ergodicity of a random process. That is, to prove the mean-ergodicity of  $\{X(t)\}$ , it is enough to prove  $\lim_{T \rightarrow \infty} \text{Var}(\bar{X}_T) = 0$ .

## Correlation Ergodic Process

The stationary process  $\{X(t)\}$  is said to be correlation ergodic (or ergodic in the correlation), if the process  $\{Y(t)\}$  is mean-ergodic, where  $Y(t) = X(t + \lambda) \times X(t)$ . That is, the stationary process  $\{X(t)\}$  is correlation ergodic, if  $\bar{Y}_T = \frac{1}{2T} \int_{-T}^T X(t + \lambda) X(t) dt$  tends to

$$E\{X(t + \lambda) X(t)\} = R(\lambda) \text{ as } T \rightarrow \infty.$$

## Distribution Ergodic Process

If  $\{X(t)\}$  is a stationary process and if  $\{Y(t)\}$  is another process such that

$$Y(t) = \begin{cases} 1 & \text{if } X(t) \leq x \\ 0 & \text{if } X(t) > x \end{cases}$$

then  $\{X(t)\}$  is said to be distribution-ergodic, if  $\{Y(t)\}$  is mean-ergodic. That is, the stationary process  $\{X(t)\}$  is distribution ergodic, if  $\bar{Y}_T = \frac{1}{2T} \int_{-T}^T Y(t) dt \rightarrow E\{Y(t)\}$  as  $T \rightarrow \infty$ .

We note that

$$\begin{aligned} E\{Y(t)\} &= 1 \times P\{X(t) \leq x\} + 0 \times P\{X(t) > x\} \\ &= F_X(x) \end{aligned}$$

Thus, the stationary process  $\{X(t)\}$  is distribution-ergodic,

$$\text{if } \frac{1}{2T} \int_{-T}^T Y(t) dt \rightarrow F_X(x) \text{ as } T \rightarrow \infty$$

### Worked Example 6(B)

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#### Example 1

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Given that the autocorrelation function for a stationary ergodic process with no periodic components is

$$R_{xx}(\tau) = 25 + \frac{4}{1+6\tau^2}$$

Find the mean value and variance of the process  $\{X(t)\}$ .

By the property of autocorrelation function,

$$\begin{aligned} \mu_x^2 &= \lim_{\tau \rightarrow \infty} R_{xx}(\tau) \\ &= 25 \\ \therefore \mu_x &= 5 \\ E\{X^2(t)\} &= R_{xx}(0) \\ &= 25 + 4 = 29 \\ \therefore \text{Var}\{X(t)\} &= E\{X^2(t)\} - E^2\{X(t)\} \\ &= 29 - 25 = 4 \end{aligned}$$

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#### Example 2

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Express the autocorrelation function of the process  $\{X'(t)\}$  in terms of the autocorrelation function of the process  $\{X(t)\}$ .

Consider

$$\begin{aligned} R_{xx'}(t_1, t_2) &= E\{X(t_1) X'(t_2)\} \\ &= E\left[X(t_1) \left\{ \frac{X(t_2 + h) - X(t_2)}{h} \right\} \right] \text{ as } h \rightarrow 0 \\ &= \lim_{h \rightarrow 0} \left[ \frac{R_{xx}(t_1, t_2 + h) - R_{xx}(t_1, t_2)}{h} \right] \\ &= \frac{\partial}{\partial t_2} R_{xx}(t_1, t_2) \end{aligned} \tag{1}$$

Similarly,

$$R_{x'x'}(t_1, t_2) = \frac{\partial}{\partial t_1} R_{xx'}(t_1, t_2) \tag{2}$$

Using (1) in (2),

$$R_{x'x'}(t_1, t_2) = \frac{\partial^2}{\partial t_1 \partial t_2} R_{xx}(t_1, t_2) \quad (3)$$

If  $\{X(t)\}$  is a stationary process, we put  $t_1 - t_2 = \tau$ . From (1), (2) and (3), we then get

$$R_{xx'}(\tau) = -\frac{\partial}{\partial \tau} R_{xx}(\tau)$$

$$R_{x'x'}(\tau) = \frac{\partial}{\partial \tau} R_{xx'}(\tau) \text{ and}$$

$$R_{x'x'}(\tau) = \frac{\partial^2}{\partial \tau^2} R_{xx}(\tau)$$

### Example 3

Prove that the random process  $\{X(t)\}$  with constant mean is mean-ergodic,

$$\text{if } \lim_{T \rightarrow \infty} \left[ \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1, t_2) dt_1 dt_2 \right] = 0$$

As per mean-ergodic theorem, the condition for the mean-ergodicity of the process  $\{X(t)\}$  is

$$\lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} = 0, \text{ where}$$

$$\bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt \text{ and } E(\bar{X}_T) = E\{X(t)\}$$

$$\text{Now, } \bar{X}_T^2 = \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T X(t_1) X(t_2) dt_1 dt_2$$

$$\therefore E\{\bar{X}_T^2\} = \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T R(t_1, t_2) dt_1 dt_2$$

$$\begin{aligned} \therefore \text{Var}(\bar{X}_T) &= E\{\bar{X}_T^2\} - E^2(\bar{X}_T) \\ &= \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T [R(t_1, t_2) - E\{X(t_1)\} E\{X(t_2)\}] dt_1 dt_2 \\ &= \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1, t_2) dt_1 dt_2 \end{aligned} \quad (1)$$

Therefore, the condition  $\lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} = 0$  is equivalent to the condition

$$\lim_{T \rightarrow \infty} \left[ \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1, t_2) dt_1 dt_2 \right] = 0$$

Hence, the result.

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### Example 4

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If  $\bar{X}_T$  is the time-average of a stationary random process  $\{X(t)\}$  over  $\{-T, T\}$ , prove that  $\text{Var}(\bar{X}_T) = \frac{1}{T} \int_0^{2T} C(\tau) \left[ 1 - \frac{|\tau|}{2T} \right] d\tau$  and hence prove that the sufficient condition for the mean-ergodicity of the process  $\{X(t)\}$  is

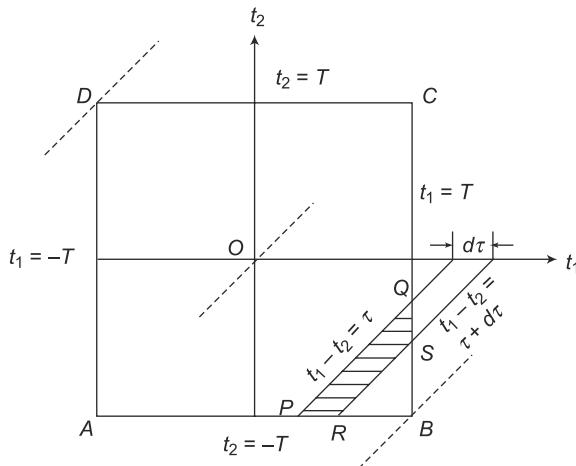
$$(a) \lim_{T \rightarrow \infty} \frac{1}{T} \int_0^{2T} C(\tau) \left[ 1 - \frac{|\tau|}{2T} \right] d\tau = 0 \quad \text{and} \quad (b) \int_{-\infty}^{\infty} |C(\tau)| d\tau < \infty.$$

(a) Step (1) of the previous example gives

$$\text{Var}(\bar{X}_T) = \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1, t_2) dt_1 dt_2 \quad (1)$$

We shall convert the double integral (1) into a single definite integral with respect to the variable  $\tau = t_1 - t_2$  as explained below:

The double integral (1) is evaluated over the area of the square bounded by  $T_1 = -T, T$  and  $t_2 = -T, T$  as shown in the figure.



**Fig. 6.1**

We divide the area of the square  $ABCD$  into a number of strips parallel to the line  $t_1 - t_2 = 0$ . Let a typical strip be  $PQRS$ , where  $PQ$  is given by  $t_1 - t_2 = \tau$  and  $RS$  is given by  $t_1 - t_2 = \tau + d\tau$ .

When  $PQRS$  is at the initial position  $D$ ,  $t_1 - t_2 = -2T$ , i.e., the initial value of  $\tau = -2T$ .

When  $PQRS$  is at the final position  $B$ ,  $t_1 - t_2 = 2T$ , i.e., final value of  $\tau = 2T$ .

Hence, to cover the given area  $ABCD$ ,  $\tau$  has to vary from  $-2T$  to  $2T$ .

Since  $d\tau$  is very small,  $C(t_1 - t_2) = C(\tau)$  can be assumed to be a constant in the strip  $PQRS$ .

$$\text{Now, } dt_1 dt_2 = \text{elemental area in the } t_1 t_2\text{-plane} \\ = \text{area of the small strip } PQRS \quad (2)$$

$t_1$  co-ordinate of  $P$  is obtained by solving the equations  $t_1 - t_2 = \tau$  and  $t_2 = -T$ .

Thus,  $(t_1)_P = \tau - T$ .

$$\therefore PB (= BQ) = T - (\tau - T) = 2T - \tau \text{ if } \tau > 0 \\ = 2T + \tau \text{ if } \tau < 0$$

When  $\tau > 0$

$$\begin{aligned} \text{Area of } PQRS &= \text{Area of } \Delta PBQ - \text{Area of } \Delta RSB \\ &= \frac{1}{2}(2T - \tau)^2 - \frac{1}{2}(2T - \tau - d\tau)^2 \\ &= (2T - \tau) d\tau, \text{ omitting } (d\tau)^2 \end{aligned} \quad (3)$$

From (2) and (3),

$$dt_1 dt_2 = \{2T - |\tau|\} d\tau \quad (4)$$

Using (4) in (1),

$$\text{Var}(\bar{X}_T) = \frac{1}{2T} \int_{-2T}^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau$$

$$\text{i.e., } \text{Var}(\bar{X}_T) = \frac{1}{T} \int_0^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \text{ (since the integrand is even)}$$

Therefore, sufficient condition for mean-ergodicity of a stationary process  $\{X(t)\}$  can also be stated as

$$\lim_{T \rightarrow \infty} \left[ \frac{1}{T} \int_0^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \right] = 0$$

(b) The sufficient condition for mean ergodicity of  $\{X(t)\}$  can also given as

$$\lim_{T \rightarrow \infty} \left[ \frac{1}{T} \int_{-2T}^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \right] = 0 \quad (5)$$

Since  $\tau$  varies from  $-2T$  to  $2T$ ,  $|\tau| \leq 2T$ .

$$\therefore 1 - \frac{|\tau|}{2T} \leq 1$$

$$\therefore \frac{1}{2T} \int_{-2T}^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \leq \frac{1}{2T} \int_{-2T}^{2T} |C(\tau)| d\tau$$

$$(5) \text{ will be true, only if } \lim_{T \rightarrow \infty} \frac{1}{2T} \int_{-2T}^{2T} |C(\tau)| d\tau = 0$$

i.e., if  $\int_{-\infty}^{\infty} |C(t)| d\tau$  is finite

i.e., if  $\int_{-\infty}^{\infty} |C(\tau)| d\tau < \infty$

Therefore, sufficient condition for mean ergodicity of the stationary process  $\{X(t)\}$  can also be stated as  $\int_{-\infty}^{\infty} |C(\tau)| d\tau < \infty$ .

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### Example 5

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The random binary transmission process  $\{X(t)\}$  is a WSS process with zero mean and autocorrelation function  $R(\tau) = 1 - \frac{|\tau|}{T}$ , where  $T$  is a constant. Find the mean and variance of the time average of  $\{X(t)\}$  over  $(0, T)$ . Is  $\{X(t)\}$  mean ergodic?

$$\bar{X}_T = \frac{1}{T} \int_0^T X(t) dt$$

$$E(\bar{X}_T) = E\{X(t)\} = 0$$

$$\begin{aligned} \text{Var}(\bar{X}_T) &= \frac{1}{T} \int_{-T}^T \left\{ 1 - \frac{|\tau|}{T} \right\} C(t) d\tau \quad [\text{see (5) in the previous problem}] \\ &= \frac{1}{T} \int_{-T}^T \left\{ 1 - \frac{|\tau|}{T} \right\}^2 d\tau \\ &= \frac{2}{T} \int_0^T \left\{ 1 - \frac{\tau}{T} \right\}^2 d\tau = \frac{2}{3} \end{aligned}$$

$$\lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} = \frac{2}{3} \neq 0$$

i.e., the condition for mean-ergodicity of  $\{X(t)\}$  is not satisfied. Therefore,  $\{X(t)\}$  is not mean-ergodic.

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### Example 6

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If  $\{X(t)\}$  is a WSS process with mean  $\mu$  and autocovariance function

$$C_{xx}(\tau) = \begin{cases} \sigma_x^2 \left( 1 - \frac{|\tau|}{\tau_0} \right) & \text{for } 0 \leq |\tau| \leq \tau_0 \\ 0 & \text{for } |\tau| \geq \tau_0 \end{cases}$$

find the variance of the time average of  $\{X(t)\}$  over  $(0, T)$ . Also examine if the process  $\{X(t)\}$  is mean-ergodic.

$$\begin{aligned}
\text{Var}(\bar{X}_T) &= \frac{1}{T} \int_{-T}^T \left\{ 1 - \frac{|\tau|}{T} \right\} C(\tau) d\tau \\
&= \frac{2\sigma_x^2}{T} \int_0^T \left( 1 - \frac{\tau}{T} \right) \left( 1 - \frac{\tau}{\tau_0} \right) d\tau, \text{ if } 0 < T \leq \tau_0 \\
&= \frac{2\sigma_x^2}{T} \left( T - \frac{T}{2} - \frac{T^2}{2\tau_0} + \frac{T^2}{3\tau_0} \right) \\
&= \sigma_x^2 \left( 1 - \frac{T}{3\tau_0} \right), \text{ if } 0 < T \leq \tau_0
\end{aligned} \tag{1}$$

and

$$\begin{aligned}
\text{Var}(\bar{X}_T) &= \frac{2\sigma_x^2}{T} \int_0^{\tau_0} \left( 1 - \frac{\tau}{T} \right) \left( 1 - \frac{\tau}{\tau_0} \right) d\tau, \text{ if } T \geq \tau_0 \\
&= \frac{\sigma_x^2 \tau_0}{T} \left( 1 - \frac{\tau_0}{3T} \right), \text{ if } T \geq \tau_0
\end{aligned} \tag{2}$$

When  $T$  is sufficiently large, (2) holds good.

$$\therefore \lim_{T \rightarrow \infty} \text{Var}(\bar{X}_T) = \lim_{T \rightarrow \infty} \left\{ \frac{\sigma_x^2 \tau_0}{T} \left( 1 - \frac{\tau_0}{3T} \right) \right\} = 0$$

Therefore,  $\{X(t)\}$  is mean-ergodic.

### Example 7

If  $X(t) = \mu + N(t)$ , where  $N(t)$  is white noise with  $C(t_1, t_2) = \phi(t_1)\delta(t_1 - t_2)$ , where  $\phi(t)$  is a bounded function of  $t$  and  $\delta$  is the unit impulse function, prove that  $\{X(t)\}$  is a mean-ergodic process.

Let us consider the time-average of  $\{X(t)\}$  over  $(-T, T)$

$$\begin{aligned}
\bar{X}_T &= \frac{1}{2T} \int_{-T}^T X(t) dt \\
E\{\bar{X}_T\} &= \frac{1}{2T} \int_{-T}^T E\{\mu + N(t)\} dt \\
&= \mu \quad [\text{since } E\{N(t)\} = 0, \text{ as } N(t) \text{ is white noise}]
\end{aligned}$$

**Note**

The unit impulse function  $\delta(t - a)$  is defined as

$$\delta(t - a) = \begin{cases} \frac{1}{\epsilon} & \text{if } a - \frac{\epsilon}{2} \leq t \leq a + \frac{\epsilon}{2} \\ 0 & \text{otherwise, where } \epsilon \rightarrow 0 \end{cases}$$

$$\begin{aligned} \int_{-\infty}^{\infty} \phi(t) \delta(t-a) dt &= \int_{a-\frac{\epsilon}{2}}^{a+\frac{\epsilon}{2}} \phi(t) \frac{1}{\epsilon} dt \\ &= \phi(a) \quad [\text{since } \phi(t) \text{ can be considered a constant in the small interval of length } \epsilon] \end{aligned}$$

Now,  $\text{Var}(\bar{X}_T) = \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1 t_2) dt_1 dt_2$  [refer to Example (3)]

$$\begin{aligned} &= \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T \phi(t_1) \delta(t_1 - t_2) dt_1 dt_2 \\ &= \frac{1}{4T^2} \int_{-T}^T \phi(t_2) dt_2 \quad (\text{by the note above}) \end{aligned}$$

$\therefore \lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} = 0$ , since  $\phi(t)$  is bounded

$\therefore (\bar{X}_T) \rightarrow \mu$ , as  $T \rightarrow \infty$

i.e.,  $\{X(t)\}$  is a mean-ergodic process.

### Example 8

If the WSS prove  $\{X(t)\}$  is given by  $X(t) = 10 \cos(100t + \theta)$ , where  $\theta$  is uniformly distributed over  $(-\pi, \pi)$ , prove that  $\{X(t)\}$  is correlation ergodic.

$$\begin{aligned} R(\tau) &= E[10 \cos(100t + 100\tau + \theta) \times 10 \cos(100t + \theta)] \\ &= 50 \cos(100\tau) \end{aligned}$$

Consider  $\bar{Z}_T = \frac{1}{2T} \int_{-T}^T X(t + \tau) X(t) dt$

$$\begin{aligned} &= \frac{1}{2T} \int_{-T}^T 100 \cos(100t + 100\tau + \theta) \cos(100t + \theta) dt \\ &= \frac{25}{T} \int_{-T}^T \cos(100\tau) dt + \frac{25}{T} \int_{-T}^T \cos(200t + 100\tau + 2\theta) dt \\ &= 50 \cos(100\tau) + \frac{25}{T} \int_{-T}^T \cos(200t + 100\tau + 2\theta) dt \end{aligned}$$

Now,  $\lim_{T \rightarrow \infty} (Z_T) = 50 \cos(100\tau) = R(\tau)$

Therefore,  $\{X(t)\}$  is correlation-ergodic.

## Exercise 6(B)

**Part-A (Short-answer Questions)**

1. Define autocorrelation function (ACF) of a stationary process.
2. Prove that the ACF  $R(\tau)$  of a real process is an even function of  $\tau$ .
3. If  $R(\tau)$  is the ACF of a complex process, prove that  $R^*(\tau) = R(-\tau)$ .
4. If  $R(\tau)$  is the ACF of a stationary process, prove that  $|R(\tau)| \leq R(0)$ .
5. If  $R(\tau)$  is the ACF of a stationary process  $\{X(t)\}$ , prove that  $\lim_{\tau \rightarrow \infty} [R(\tau)] = \mu_x^2$ . Is it true for any stationary process?
6. Find the mean of the stationary process  $\{X(t)\}$ , whose ACF is given by  

$$R(\tau) = \frac{25\tau^2 + 36}{6.25\tau^2 + 4}.$$
7. Find the variance of the stationary process  $\{X(t)\}$ , whose ACF is given by  

$$R(\tau) = \frac{25\tau^2 + 36}{6.25\tau^2 + 4}.$$
8. Find the variance of the stationary process  $\{X(t)\}$ , whose ACF is given by  

$$R(\tau) = 2 + 4e^{-2|\tau|}.$$
9. Find the variance of the stationary process  $\{X(t)\}$ , whose ACF is given by  

$$R(\tau) = 16 + \frac{9}{1 + 6\tau^2}.$$
10. Define the cross-correlation function and state any 2 of its properties.
11. Find the cross-correlation function of 2 stationary processes that are orthogonal.
12. When are the jointly stationary processes  $\{X(t)\}$  and  $\{Y(t)\}$  independent, prove that  $R_{xy}(\tau) = \mu_x \mu_y$ .
13. Define ensemble average and time average of a random process  $\{X(t)\}$ .
14. What is the difference between ensemble average and time average of a stochastic process  $\{X(t)\}$ ?
15. When is a random process said to be ergodic? Give an example for an ergodic process.
16. Distinguish between stationarity and ergodicity.
17. What do you mean by mean-ergodicity of a RP?
18. State mean ergodic theorem.
19. State the sufficient conditions for the mean ergodicity of a RP  $\{X(t)\}$ .
20. State 2 different sufficient conditions for  $\{X(t)\}$  with constant mean to be mean-ergodic.
21. Examine if the process  $\{X(t)\}$ , where  $X(t) = X$ , a random variable is mean-ergodic.
22. Give an example of a WSS process which is not mean-ergodic.
23. If  $\bar{X}_T$  is the time-average of a stationary random process  $\{X(t)\}$  over  $(-T, T)$ , express  $\text{Var}(\bar{X}_T)$  in terms of the autocovariance function of

- $\{X(t)\}$  and hence state the sufficient condition for the mean-ergodicity of  $\{X(t)\}$ .
24. When is a random process said to be correlation ergodic?
  25. When is a random process said to be distribution ergodic?

**Part-B**

26. A stationary process has an autocorrelation function given by  $R(\tau) = \frac{25\tau^2 + 36}{6.25\tau^2 + 4}$ . Find the mean value, mean-square value and variance of the process.
  27. If the autocorrelation of a process  $\{X(t)\}$  is  $R(t_1, t_2)$  and if  $Y(t) = X(t+a) - X(a)$ , where  $a$  is a constant, express  $R_{yy}(t_1, t_2)$  in terms of  $R$ 's.
  28. If  $\{X(t)\}$  is a WSS process with autocorrelation function  $R_{xx}(\tau)$  and if  $Y(t) = X(t+a) - X(t-a)$  show that  $R_{yy}(\tau) = 2R_{xx}(\tau) - R_{xx}(\tau+2a) - R_{xx}(\tau-2a)$ .
  29. If  $\{X(t)\}$  and  $\{Y(t)\}$  are independent WSS processes with zero means, find the autocorrelation function of  $\{Z(t)\}$ , when
    - (i)  $Z(t) = a + bX(t) + cY(t)$
    - (ii)  $Z(t) = aX(t)Y(t)$
  30. If  $X(t) = A$ , where  $A$  is random variable, prove that  $\{X(t)\}$  is not mean-ergodic.
  31. If  $S = \int_0^{10} X(t)dt$ , show that  $E(S^2) = \int_{-10}^{10} (10 - |\tau|)R_{xx}(\tau)dt$ . Find also the mean and variance of  $S$ , if  $E\{X(t)\} = 8$  and  $R_{xx}(\tau) = 64 + 10e^{-2|\tau|}$ .
  32. A stationary zero mean random process  $\{X(t)\}$  has the autocorrelation function  $R_{xx}(\tau) = 10e^{-0.1\tau^2}$ . Find the mean and variance of  $\bar{X}_T = \frac{1}{5} \int_0^5 X(t)dt$ .
  33. If  $\{X(t)\}$  is a WSS process with  $E\{X(t)\} = 2$  and  $R_{xx}(\tau) = 4 + e^{-|\tau|/10}$ , find the mean and variance of  $S = \int_0^1 X(t)dt$ .
  34.  $\{X(t)\}$  is the random telegraph signal process with  $E\{X(t)\} = 0$  and  $R(\tau) = e^{-2\lambda|\tau|}$ . Find the mean and variance of the time average of  $\{X(t)\}$  over  $(-T, T)$ . Is it mean-ergodic?
  35. The random process  $\{X(t)\}$  is stationary with  $E\{X(t)\} = 1$  and  $R(\tau) = 1 + e^{-2|\tau|}$ . Find the mean and variance of  $S = \int_0^1 X(t)dt$ .
  36. If the autocovariance function of a stationary process  $\{X(t)\}$  is given by  $C(\tau) = Ae^{-c\tau}$ , prove that  $\{X(t)\}$  is mean-ergodic. Also find  $\text{Var}(\bar{X}_T)$ , where  $\bar{X}_T$  is the time average of  $\{X(t)\}$  over  $(-T, T)$ .
- [Hint:  $c(\tau) \rightarrow 0$  as  $\tau \rightarrow \infty$ ]

37. If the autocorrelation function of a WSS process  $\{X(t)\}$  is  $R(\tau)$ , show that  $P\{|X(t + \tau) - X(t)| \geq a\} \leq 2[R(0) - R(\tau)]/a^2$ .  
**[Hint:** Use Tchebycheff's inequality.]
38. Show that, if  $\{X(t)\}$  is normal with  $\mu_x = 0$  and  $R_{xx}(\tau) = 0$  for  $|\tau| > a$ , then it is correlation-ergodic.  
**[Hint:**  $Z(t) = X(t + \lambda)X(t)$ ;  $R_{zz}(\tau) = R^2(\lambda) + R^2(\tau) + R(\lambda + \tau)R(\lambda - \tau)$ . Since  $R_{zz}(\tau) \rightarrow 0$  as  $\tau \rightarrow \infty$ ,  $C_{zz}(\tau) \rightarrow 0$  as  $\tau \rightarrow \infty$ ]
39. If the autocorrelation function of a stationary Gaussian process  $\{X(t)\}$  is  $R(\tau) = 10e^{-|\tau|}$ , prove that  $\{X(t)\}$  is ergodic both in mean and correlation.
40. Give an example of a WSS process which is not ergodic in mean.

## Power Spectral Density Function

The autocorrelation function  $R(\tau)$  tells us something about how rapidly we can expect the random signal  $X(t)$  to change as a function of time. If the autocorrelation function decays rapidly (slowly), it indicates that the process can be expected to change rapidly (slowly). Moreover if the autocorrelation function has periodic components, then the corresponding process also will have periodic components. Hence, we may conclude that the autocorrelation function contains information about the expected frequency content of the random process.

For example, if we assume that  $X(t)$  is a voltage waveform across a  $1 \Omega$  resistance, then the ensemble average value of  $X^2(t)$  is the average value of power delivered to the  $1 \Omega$  resistance by  $X(t)$ .

$$\begin{aligned}\text{i.e., Average power of } X(t) &= E\{X^2(t)\} \\ &= R(0) \geq 0\end{aligned}$$

Now, if  $R(0)$  can be expressed as

$$R(0) = \int_{-\infty}^{\infty} S(f) df, \text{ since } R(0) \text{ represents power,}$$

$S(f)$  will be expressed in units of power per Hertz. That is,  $S(f)$  gives the distribution of power of  $\{X(t)\}$  as a function of frequency and hence is called the *power spectral density function* or *simply spectral density* or *power spectrum* of the stationary process  $\{X(t)\}$ . We shall now give the mathematical definition of power spectral density function of a stationary process.

**Definition:** If  $\{X(t)\}$  is a stationary process (either in the strict sense or wide sense) with autocorrelation function  $R(\tau)$ , then the Fourier transform of  $R(\tau)$  is called the power spectral density function of  $\{X(t)\}$  and denoted as  $S_{xx}(\omega)$  or  $S_x(\omega)$  or  $S(\omega)$ .

$$\text{Thus, } S(\omega) = \int_{-\infty}^{\infty} R(\tau) e^{-i\omega\tau} d\tau \quad (1)$$

Sometimes  $\omega$  is replaced by  $2\pi f$ , where  $f$  is the frequency variable, in which case the power spectral density function will be a function of  $f$ , denoted by  $S(f)$ .

Then 
$$S(f) = \int_{-\infty}^{\infty} R(\tau) e^{-i2\pi f \tau} d\tau \quad (2)$$

**Note** Equation (1) or (2) is sometimes called the Wiener-Kinchine relation. We shall mostly follow the definition (1) and denote the power spectral density as a function of  $\omega$  only.

Given the power spectral density functions  $S(\omega)$ , the autocorrelation function  $R(\tau)$  is given by the Fourier inverse transform of  $S(\omega)$ .

i.e., 
$$R(\tau) = \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) e^{i\omega\tau} d\omega \quad (3)$$

(or) 
$$R(\tau) = \int_{-\infty}^{\infty} S(f) e^{i2\pi f \tau} df \quad (4)$$

If  $\{X(t)\}$  and  $\{Y(t)\}$  are two jointly stationary random processes with cross-correlation function  $R_{xy}(\tau)$ , then the Fourier transform of  $R_{xy}(\tau)$  is called the *cross-power spectral density* of  $\{X(t)\}$  and  $\{Y(t)\}$  and denoted as  $S_{xy}(\omega)$ .

i.e., 
$$S_{xy}(\omega) = \int_{-\infty}^{\infty} R_{xy}(\tau) e^{-i\omega\tau} d\tau$$

## Properties of Power Spectral Density Function

1. The value of the spectral density function at zero frequency is equal to the total area under the graph of the autocorrelation function. By putting  $\omega = 0$  in (1) or  $f = 0$  in (2), we get

$$S(0) = \int_{-\infty}^{\infty} R(\tau) d\tau, \text{ which is the given property.}$$

2. The mean square value of a wide-sense stationary process is equal to the total area under the graph of the spectral density. By putting  $f = 0$  in (4), we get

$$E\{X^2(t)\} = R(0) = \int_{-\infty}^{\infty} S(f) df, \text{ which is the given property.}$$

3. The spectral density function of a real random process is an even function.

### Proof

$$S(\omega) = \int_{-\infty}^{\infty} R(\tau) e^{-i\omega\tau} d\tau, \text{ by definition}$$

$$\therefore S(-\omega) = \int_{-\infty}^{\infty} R(\tau) e^{i\omega\tau} d\tau$$

Putting  $\tau = -u$ ,

$$\begin{aligned} S(-\omega) &= \int_{-\infty}^{\infty} R(-u) e^{-i\omega u} du \\ &= \int_{-\infty}^{\infty} R(u) e^{-i\omega u} du \quad [\text{since } R(\tau) \text{ is an even function of } \tau] \\ &= S(\omega) \end{aligned}$$

Therefore,  $S(\omega)$  is an even function of  $\omega$ .

4. The spectral density of a process  $\{X(t)\}$ , real or complex, is a real function of  $\omega$  and non-negative.

### Proof

$$\begin{aligned} R(\tau) &= E\{X(t)X^*(t-\tau)\} \\ \therefore R(-\tau) &= E\{X(t)X^*(t+\tau)\} \\ \therefore R^*(-\tau) &= E\{X(t+\tau)X^*(t)\} \\ &= R(\tau) \\ (\text{or}) \quad R^*(\tau) &= R(-\tau) \end{aligned} \tag{1}$$

$$\begin{aligned} \text{Now,} \quad S(\omega) &= \int_{-\infty}^{\infty} R(\tau) e^{-i\omega\tau} d\tau \\ \therefore S^*(\omega) &= \int_{-\infty}^{\infty} R^*(\tau) e^{i\omega\tau} d\tau \\ &= \int_{-\infty}^{\infty} R(-\tau) e^{i\omega\tau} d\tau \quad \text{by (1)} \\ &= \int_{-\infty}^{\infty} R(u) e^{-i\omega u} du, \text{ by putting } u = -\tau \\ &= S(\omega) \end{aligned}$$

Hence,  $S(\omega)$  is a real function of  $\omega$ .

**Note** It will be proved that  $S(\omega) \geq 0$ , in Worked Example 14).

5. The spectral density and the autocorrelation function of a real WSS process form a Fourier cosine transform pair.

### Proof

$$\begin{aligned} S(\omega) &= \int_{-\infty}^{\infty} R(\tau) \{ \cos \omega \tau - i \sin \omega \tau \} d\tau \\ &= 2 \int_{-\infty}^{\infty} R(t) \cos \omega \tau d\tau \quad [\text{since } R(\tau) \text{ is even}] \end{aligned}$$

= Fourier cosine transform of  $[2R(\tau)]$

$$\begin{aligned} R(\tau) &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) (\cos \tau\omega + i \sin \tau\omega) d\omega \\ &= \frac{1}{\pi} \int_0^{\infty} S(\omega) \cos \tau\omega d\omega \quad [\text{since } S(\omega) \text{ is even}] \\ &= \text{Fourier inverse consine transform of } \left[ \frac{1}{2} S(\omega) \right] \end{aligned}$$

## 6. Wiener-Khinchine Theorem

If  $X_T(\omega)$  is the Fourier transform of the truncated random process defined as

$$X_T(t) = \begin{cases} X(t) & \text{for } |t| \leq T \\ 0 & \text{for } |t| > T \end{cases}$$

where  $\{X(t)\}$  is a real WSS process with power spectral density function  $S(\omega)$ , then

$$S(\omega) = \lim_{T \rightarrow \infty} \left[ \frac{1}{2T} E\{|X_T(\omega)|^2\} \right]$$

### Proof

$$\begin{aligned} \text{Given: } X_T(\omega) &= \int_{-\infty}^{\infty} X_T(t) e^{-i\omega t} dt \\ &= \int_{-T}^{T} X(t) e^{-i\omega t} dt \end{aligned}$$

Since  $\{X(t)\}$  is real

$$\begin{aligned} |X_T(\omega)|^2 &= X_T(\omega) X_T(-\omega) \\ &= \int_{-T}^{T} X(t_1) e^{-i\omega t_1} dt_1 \int_{-T}^{T} X(t_2) e^{i\omega t_2} dt_2 \\ &= \int_{-T}^{T} \int_{-T}^{T} X(t_1) X(t_2) e^{-i\omega(t_1-t_2)} dt_1 dt_2 \\ \therefore E|X_T(\omega)|^2 &= \int_{-T}^{T} \int_{-T}^{T} E\{X(t_1) X(t_2)\} e^{-i\omega(t_1-t_2)} dt_1 dt_2 \\ &= \int_{-T}^{T} \int_{-T}^{T} R(t_1 - t_2) e^{-i\omega(t_1-t_2)} dt_1 dt_2 \quad [\text{since } \{X(t)\} \text{ is WSS}] \end{aligned}$$

$$= \int_{-T}^T \int_{-T}^T \phi(t_1 - t_2) dt_1 dt_2, \text{ say} \quad (1)$$

The double integral (1) is evaluated over the area of the square  $ABCD$  bounded by  $t_1 = -T, T$  and  $t_2 = -T, T$  as shown in the figure.

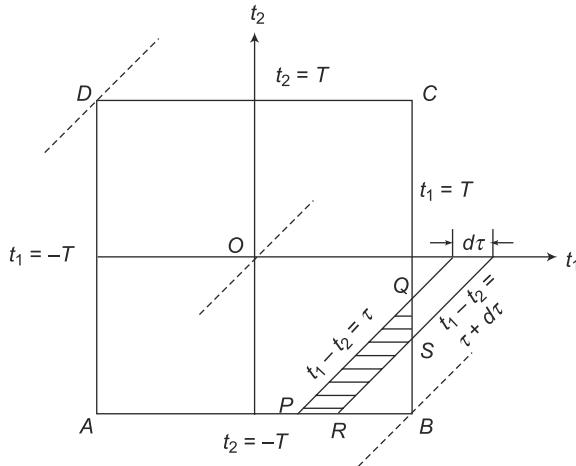


Fig. 6.2

We divide the area of the square into a number of strips like  $PQRS$ , where  $PQ$  is given by  $t_1 - t_2 = \tau$  and  $RS$  is given by  $t_1 - t_2 = \tau + d\tau$ .

When  $PQRS$  is at the initial position  $D$ ,  $t_1 - t_2 = -2T$ . When it is at the final position  $B$ ,  $t_1 - t_2 = 2T$ . Hence when  $\tau$  varies from  $-2T$  to  $2T$ , the area  $ABCD$  is covered.

Now,

$$\begin{aligned} dt_1 dt_2 &= \text{elemental area of the } t_1 t_2\text{-plane} \\ &= \text{area of } PQRS \\ (t_1)_P &= t - T \text{ and } PB (= BQ) = T - (\tau - T) = 2T - \tau, \text{ if } \tau > 0 \\ &= 2T + \tau, \text{ if } \tau < 0 \end{aligned} \quad (2)$$

When  $\tau > 0$ ,

$$\begin{aligned} \text{Area of } PQRS &= \Delta PBQ - \Delta RBS \\ &= \frac{1}{2}(2T - \tau)^2 - \frac{1}{2}(2T - \tau - d\tau)^2 \\ &= (2T - \tau) d\tau, \text{ omitting } (d\tau)^2 \end{aligned} \quad (3)$$

From (2) and (3),

$$dt_1 dt_2 = (2T - |\tau|) d\tau \quad (4)$$

Using (4) in (1), we get,

$$E\{|X_T(\omega)|^2\} = \int_{-2T}^{2T} \phi(\tau)(2T - |\tau|) d\tau$$

$$\begin{aligned}
 \therefore \frac{1}{2T} E\{|X_T(\omega)|^2\} &= \int_{-2T}^{2T} \phi(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \\
 \therefore \lim_{T \rightarrow \infty} \frac{1}{2T} E\{|X_T(\omega)|^2\} &= \lim_{T \rightarrow \infty} \int_{-2T}^{2T} \phi(\tau) d\tau - \lim_{T \rightarrow \infty} \frac{1}{2T} \int_{-2T}^{2T} |\tau| \phi(\tau) d\tau \\
 &= \int_{-\infty}^{\infty} \phi(\tau) d\tau \quad [\text{assuming that } \int_{-\infty}^{\infty} |\tau| \phi(\tau) d\tau \text{ is bounded}] \\
 &= \int_{-\infty}^{\infty} R(\tau) e^{-i\omega\tau} d\tau \\
 &\quad [\text{provided } \int_{-\infty}^{\infty} |\tau| R(\tau) e^{-i\omega\tau} d\tau \text{ is bounded}] \\
 &= S(\omega), \text{ by definition of } S(\omega)
 \end{aligned}$$

**Note** This theorem provides an alternative method for finding  $S(\omega)$  for a WSS process.

**Linear systems with random inputs:** Mathematically, a *system* is a functional relationship between the input  $x(t)$  and the output  $y(t)$ . Usually, this relationship is written as  $y(t) = f[x(t)]$ ,  $-\infty < t < \infty$ .

If we assume that  $x(t)$  represents a sample function of a random process  $\{X(t)\}$ , the system produces an output or response  $y(t)$  and the ensemble of the output functions forms a random process  $\{Y(t)\}$ . The process  $\{Y(t)\}$  can be considered as the output of the system or transformation ' $f$ ' with  $\{X(t)\}$  as the input, the system is completely specified by the operator ' $f$ '.

We recall that  $X(t)$  actually means  $X(s, t)$ , where  $s \in S$  (sample space). If the system operates only on the variable  $t$  treating  $s$  as a parameter, it is called a *deterministic system*. If the system operates on both  $t$  and  $s$ , it is called *stochastic*. We shall consider only deterministic systems in our study.

**Definitions:** If  $f[a_1 X_1(t) \pm a_2 X_2(t)] = a_1 f[X_1(t)] \pm a_2 f[X_2(t)]$ , then  $f$  is called a linear system.

If  $Y(t+h) = f[X(t+h)]$ , where  $Y(t) = f[X(t)]$ ,  $f$  is called a *time-invariant system* or  $X(t)$  and  $Y(t)$  are said to form a time-invariant system.

If the output  $Y(t_1)$  at a given time  $t = t_1$  depends only on  $X(t_1)$  and not on any other past or future values of  $X(t)$ , then the system  $f$  is called a *memoryless system*.

If the value of the output  $Y(t)$  at  $t = t_1$  depends only on the past values of the input  $X(t)$ ,  $t \leq t_1$ , i.e.,  $Y(t_1) = f[X(t); t \leq t_1]$ , then the system is called a *causal system*.

## System in the Form of Convolution

Very often in electrical systems, the output  $Y(t)$  is expressed as a convolution of the input  $X(t)$  with a *system weighting function*  $h(t)$ , i.e., the input–output relationship will be of form

$$Y(t) = \int_{-\infty}^{\infty} h(u)X(t-u)du \quad (1)$$

## Unit Impulse Response of the System

Refer to the Note in Worked Example 7 of the previous section, in which we have proved that

$$\int_{-\infty}^{\infty} \phi(t)\delta(t-a)dt = \phi(a)$$

where  $\delta(t-a)$  is the unit impulse function at  $a$ .

If we take  $a = 0$ , we get

$$\int_{-\infty}^{\infty} \phi(t)\delta(t)dt = \phi(0) \quad (2)$$

Put  $X(t) = \delta(t)$  in (1)

$$\begin{aligned} \text{Then } Y(t) &= \int_{-\infty}^{\infty} h(u)\delta(t-u)du \\ &= \int_{-\infty}^{\infty} h(t-u)\delta(u)du \quad (\text{by the property of the convolution}) \\ &= h(t-0), \text{ by (2)} \\ &= h(t) \end{aligned}$$

Thus, if the input of the system is the unit impulse function, then the output or response is the system weighting function. Hence, the system weighting function  $h(t)$  will be hereafter called unit impulse response function.

### Properties

1. If a system is such that its input  $X(t)$  and its output  $Y(t)$  are related by a convolution integral, i.e., if  $Y(t) = \int_{-\infty}^{\infty} h(u)X(t-u)du$ , then the system is a linear time-invariant system.

### Proof

Let  $X(t) = a_1X_1(t) + a_2X_2(t)$

Then

$$\begin{aligned} Y(t) &= \int_{-\infty}^{\infty} h(u)[a_1 X_1(t-u) + a_2 X_2(t-u)] du \\ &= a_1 Y_1(t) + a_2 Y_2(t) \end{aligned}$$

Therefore, the system is linear. If  $X(t)$  is replaced by  $X(t+h)$ , then

$$\int_{-\infty}^{\infty} h(u) X(t+h-u) du = Y(t+h)$$

Therefore, the system is time-invariant.

**Note** If  $h(t)$  is absolutely integrable, viz.,  $\int_{-\infty}^{\infty} |h(t)| dt < \infty$ , then the system is said to be stable in the sense that every bounded input gives a bounded output.

In addition, if  $h(t) = 0$ , when  $t < 0$ , the system is said to be causal.

2. If the input to a time-invariant, stable linear system is a WSS process, the output will also be a WSS process.

### Proof

$$\begin{aligned} Y(t) &= \int_{-\infty}^{\infty} h(u) X(t-u) du \\ \therefore E\{Y(t)\} &= \int_{-\infty}^{\infty} h(u) E\{X(t-u)\} du \\ &= \mu_x \int_{-\infty}^{\infty} h(u) du \quad [\text{since } \{X(t)\} \text{ is WSS}] \\ &= \text{a finite constant, independent of } t \end{aligned}$$

(since the system is stable)

$$\begin{aligned} R_{yy}(t_1, t_2) &= E\{Y(t_1) Y(t_2)\} \\ &= E \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} h(u_1) h(u_2) X(t_1 - u_1) X(t_2 - u_2) du_1 du_2 \\ &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} h(u_1) h(u_2) R_{xx}(t_1 - u_1, t_2 - u_2) du_1 du_2 \\ &= \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} h(u_1) h(u_2) R_{xx}[(t_1 - t_2) - (u_1 - u_2)] du_1 du_2 \end{aligned}$$

[since  $\{X(t)\}$  is WSS]

Since the RS is a function of  $(t_1 - t_2)$ , so will be the LS. Therefore,  $R_{yy}(t_1, t_2)$  will be a function of  $(t_1 - t_2)$ . Therefore,  $\{Y(t)\}$  is a WSS process.

3. If  $\{X(t)\}$  is a WSS process and if

$$Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du, \text{ then}$$

- (i)  $R_{xy}(\tau) = R_{xx}(\tau)^* h(-\tau)$  and
- (ii)  $R_{yy}(\tau) = R_{xy}(\tau)^* h(\tau)$ , where \* denotes convolution. Also
- (iii)  $S_{xy}(\omega) = S_{xx}(\omega)H^*(\omega)$  and
- (iv)  $S_{yy}(\omega) = S_{xx}(\omega)|H(\omega)|^2$

### Proof

$$\begin{aligned} Y(t) &= \int_{-\infty}^{\infty} X(t-\alpha) h(\alpha) d\alpha \\ \therefore X(t+\tau)Y(t) &= \int_{-\infty}^{\infty} X(t+\tau) X(t-\alpha) h(\alpha) d\alpha \\ \therefore E\{X(t+\tau)Y(t)\} &= \int_{-\infty}^{\infty} E\{X(t+\tau)X(t-\alpha)\} h(\alpha) d\alpha \\ &= \int_{-\infty}^{\infty} R_{xx}(\tau+\alpha) h(\alpha) d\alpha && [\text{since } \{X(t)\} \text{ is WSS}] \\ &= \int_{-\infty}^{\infty} R_{xx}(\tau-\beta) h(-\beta) d\beta && (\text{putting } \beta = -\alpha) \end{aligned}$$

$$\text{i.e., } R_{xy}(\tau) = R_{xx}(\tau)^* h(-\tau) \quad (1)$$

$$\text{Similarly, } R_{yx}(\tau) = R_{xx}(\tau)^* h(\tau) \quad (1a)$$

$$\begin{aligned} \text{Now, } Y(t) Y(t-\tau) &= \int_{-\infty}^{\infty} X(t-\alpha) Y(t-\tau) h(\alpha) d\alpha \\ \therefore E\{Y(t) Y(t-\tau)\} &= \int_{-\infty}^{\infty} R_{xy}(\tau-\alpha) h(\alpha) d\alpha \end{aligned}$$

Assuming that  $\{X(t)\}$  and  $\{Y(t)\}$  are jointly WSS

$$\text{i.e., } R_{yy}(\tau) = R_{xy}(\tau)^* h(\tau) \quad (2)$$

Taking Fourier transforms of (1) and (2), we get

$$S_{xy}(\omega) = S_{xx}(\omega)H^*(\omega) \quad (3)$$

$$\text{where } H^*(\omega) \text{ is the conjugate of } H(\omega) \text{ and } S_{yy}(\omega) = S_{xy}(\omega)H(\omega) \quad (4)$$

### Note

Inserting (1) in (2),  $R_{yy}(\tau) = R_{xx}(\tau)^* h(\tau)^* h(-\tau)$

$$(5)$$

Inserting (3) in (4),  $S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$  (6)

However, alternative proofs for (5) and (6) are given in the following properties.

4. If  $\{X(t)\}$  is a WSS process and if

$$Y(t) = \int_{-\infty}^{\infty} h(u)X(t-u)du, \text{ then}$$

$$R_{yy}(\tau) = R_{xx}(\tau)^* K(\tau),$$

where  $K(t) = h(t)^* h(-t) = \int_{-\infty}^{\infty} h(u)h(t+u)du$

### **Proof**

$$\begin{aligned} Y(t) &= \int_{-\infty}^{\infty} X(u)h(t-u)du \\ \therefore Y(t) Y(t-\tau) &= \int_{-\infty}^{\infty} X(u)h(t-u)du \int_{-\infty}^{\infty} X(v)h(t-\tau-v)dv \\ &= \int_{-\infty}^{\infty} \left[ \int_{-\infty}^{\infty} X(u)X(v)h(t-u)du \right] h(t-\tau-v)dv \\ \therefore R_{yy}(\tau) &= \int_{-\infty}^{\infty} \left[ \int_{-\infty}^{\infty} R_{xx}(u-v)h(t-u)du \right] h(t-\tau-v)dv \end{aligned}$$

[taking expectations on both sides]

In the inner integral, put  $u - v = -w$ , treating  $v$  and  $t$  as parameters and also change the order of integration.

$$\therefore R_{yy}(\tau) = \int_{-\infty}^{\infty} R_{xx}(-w) \left\{ \int_{-\infty}^{\infty} h(t-v+w)h(t-\tau-v) \right\} dw$$

In the inner integral, put  $t - \tau - v = \alpha$ , treating  $w$ ,  $t$  and  $\tau$  as parameters.

$$\begin{aligned} \therefore R_{yy}(\tau) &= \int_{-\infty}^{\infty} R_{xx}(-w) \left[ \int_{-\infty}^{\infty} h(\alpha)h(\tau+w+\alpha)d\alpha \right] dw \\ &= \int_{-\infty}^{\infty} R_{xx}(-w)K(\tau+w)dw \end{aligned}$$

[by definition of  $K(t)$ ]

$$\begin{aligned}
 &= \int_{-\infty}^{\infty} R_{xx}(\beta)K(\tau - \beta) d\beta \\
 &= R_{xx}(\tau)^* K(\tau)
 \end{aligned}
 \quad (\text{putting } \beta = -w)$$

5. The power spectral densities of the input processes in the system are connected by the relation  $S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$ , where  $H(\omega)$  is the Fourier transform of unit impulse response function  $h(t)$ .

### **Proof**

By the previous property,

$$R_{yy}(\tau) = R_{xx}(\tau)^* K(\tau),$$

where  $K(t) = h(t)^* h(-t)$

Taking Fourier transform on both sides,

$$S_{yy}(\omega) = F\{K(\tau)\} S_{xx}(\omega) \quad (1)$$

$$\text{Let } H(\omega) = F\{h(t)\}$$

$$= \int_{-\infty}^{\infty} h(t) e^{-i\omega t} dt \quad (2)$$

$$\begin{aligned}
 \therefore H^*(\omega) &= \int_{-\infty}^{\infty} h(t) e^{i\omega t} dt \\
 &= \int_{-\infty}^{\infty} h(-s) e^{-i\omega s} ds \quad (\text{putting } s = -t) \\
 &= F\{h(-t)\} \quad (3)
 \end{aligned}$$

$$\text{Now } K(t) = h(t)^* h(-t)$$

$$\begin{aligned}
 \therefore F\{K(t)\} &= F\{h(t)\} F\{h(-t)\} \quad (\text{by convolution theorem}) \\
 &= H(\omega) H^*(\omega) \quad [\text{by (2) and (3)}] \\
 &= |H(\omega)|^2 \quad (4)
 \end{aligned}$$

Inserting (4) in (1), we get

$$S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

**Definition:**  $H(\omega)$  is called the *system function* or the *power transfer function*. It is the Fourier transform of the unit impulse response function of the system.

## Worked Example 6(C)

### Example 1

The autocorrelation function of the random telegraph signal process is given by  $R(\tau) = a^2 e^{-2\gamma|\tau|}$ . Determine the power density spectrum of the random telegraph signal.

$$\begin{aligned}
 S(\omega) &= \int_{-\infty}^{\infty} R(\tau) e^{-i\omega\tau} d\tau \\
 &= a^2 \int_{-\infty}^{\infty} e^{-2\gamma|\tau|} (\cos \omega\tau - i \sin \omega\tau) d\tau \\
 &= 2a^2 \int_0^{\infty} e^{-2\gamma\tau} \cos \omega\tau d\tau \\
 &= \left[ \frac{2a^2 e^{-2\pi}}{4\gamma^2 + \omega^2} (-2\gamma \cos \omega\tau + \omega \sin \omega\tau) \right]_0^{\infty} \\
 &= \frac{4a^2 \gamma}{4\gamma^2 + \omega^2}
 \end{aligned}$$

### Example 2

The autocorrelation function of the Poisson increment process is given by

$$R(\tau) = \begin{cases} \lambda^2 & \text{for } |\tau| > \epsilon \\ \lambda^2 + \frac{\lambda}{\epsilon} \left(1 - \frac{|\tau|}{\epsilon}\right) & \text{for } |\tau| \leq \epsilon \end{cases}$$

Prove that its spectral density is given by

$$\begin{aligned}
 S(\omega) &= 2\pi\lambda^2 \delta(\omega) + \frac{4\lambda \sin^2(\omega\epsilon/2)}{\epsilon^2 \omega^2} \\
 S(\omega) &= \int_{-\epsilon}^{\epsilon} \left\{ \lambda^2 + \frac{\lambda}{\epsilon} \left(1 - \frac{|\tau|}{\epsilon}\right) \right\} e^{-i\omega\tau} d\tau + \int_{-\infty}^{-\epsilon} \lambda^2 e^{-i\omega\tau} d\tau \\
 &\quad + \int_{\epsilon}^{\infty} \lambda^2 e^{-i\omega\tau} d\tau \\
 &= \frac{\lambda}{\epsilon} \int_{-\epsilon}^{\epsilon} \left(1 - \frac{|\tau|}{\epsilon}\right) e^{-i\omega\tau} d\tau + \int_{-\infty}^{\infty} \lambda^2 e^{-i\omega\tau} d\tau \\
 &= \frac{2\lambda}{\epsilon} \int_0^{\epsilon} \left(1 - \frac{\tau}{\epsilon}\right) \cos \omega\tau d\tau + F\{\lambda^2\}
 \end{aligned}$$

where  $F(\lambda^2)$  is the Fourier transform of  $\lambda^2$ .

$$\begin{aligned}
 &= \frac{2\lambda}{\epsilon} \left[ \left( 1 - \frac{\tau}{\epsilon} \right) \frac{\sin \omega \tau}{\omega} + \frac{1}{\epsilon} \left( \frac{-\cos \omega \tau}{\omega^2} \right) \right]_0^\epsilon + F(\lambda^2) \\
 &= \frac{2\lambda}{\epsilon^2 \omega^2} (1 - \cos \omega \epsilon) + F(\lambda^2) \\
 &= \frac{4\lambda \sin^2(\omega \epsilon / 2)}{\epsilon^2 \omega^2} + F(\lambda^2)
 \end{aligned} \tag{1}$$

The Fourier inverse transform of  $S(\omega)$  is given by

$$\begin{aligned}
 R(\tau) &= F^{-1}\{S(\omega)\} \\
 &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) e^{i\tau\omega} d\omega
 \end{aligned}$$

Let us now find  $R(\tau)$  corresponding to  $S(\omega) = 2\pi\lambda^2\delta(\omega)$ , where  $\delta(\omega)$  is the unit impulse function.

$$\text{i.e., } R(\tau) = F^{-1}\{2\pi\lambda^2\delta(\omega)\}$$

$$\begin{aligned}
 &= \frac{2\pi\lambda^2}{2\pi} \int_{-\infty}^{\infty} \delta(\omega) e^{i\tau\omega} d\omega \\
 &= \lambda^2 \left[ \text{since } \int_{-\infty}^{\infty} \phi(t) \delta(t) dt = \phi(0) \right]
 \end{aligned}$$

$$\therefore F(\lambda^2) = 2\pi \lambda^2 \delta(\omega) \tag{2}$$

Inserting (2) in (1) the required result is obtained.

### Example 3

Find the power spectral density of a WSS process with autocorrelation function

$$\begin{aligned}
 R(\tau) &= e^{-\alpha\tau^2} \\
 S(\omega) &= \int_{-\infty}^{\infty} e^{-\alpha\tau^2} e^{-i\omega\tau} d\tau \\
 &= e^{-\omega^2/4\alpha} \int_{-\infty}^{\infty} e^{-\alpha\left(\tau + \frac{i\omega}{2\alpha}\right)^2} d\tau \\
 &= \frac{1}{\sqrt{\alpha}} e^{-\omega^2/4\alpha} \int_{-\infty}^{\infty} e^{-x^2} dx, \text{ putting } \sqrt{\alpha}\left(\tau + \frac{i\omega}{2\alpha}\right) = x \\
 &= \sqrt{\frac{\pi}{\alpha}} e^{-\omega^2/4\alpha} \quad \left[ \text{since } \int_{-\infty}^{\infty} e^{-x^2} dx = \left[ \left( \frac{1}{2} \right) \right] = \sqrt{\pi} \right]
 \end{aligned}$$

---

### Example 4

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A random process  $\{X(t)\}$  is given by  $X(t) = A \cos pt + B \sin pt$ , where  $A$  and  $B$  are independent RVs such that  $E(A) = E(B) = 0$  and  $E(A^2) = E(B^2) = \sigma^2$ . Find the power spectral density of the process [refer to Problem 7 in Exercise 6(a)].

The autocorrelation function of the given process can be found as

$$R(\tau) = \sigma^2 \cos p\tau$$

$$S(\omega) = \int_{-\infty}^{\infty} \sigma^2 \cos p\tau e^{-i\omega\tau} d\tau \quad (1)$$

Consider  $F^{-1}\{\pi\sigma^2[\delta(\omega+p) + \delta(\omega-p)]\}$

$$= \frac{1}{2\pi} \pi\sigma^2 \int_{-\infty}^{\infty} [\delta(\omega+p) + \delta(\omega-p)] e^{i\tau\omega} d\omega$$

$$= \frac{\sigma^2}{2} [e^{-i\tau p} + e^{i\tau p}] \quad \left\{ \text{since } \int_{-\infty}^{\infty} \phi(t)\delta(t-a) dt = \phi(a) \right\}$$

$$= \sigma^2 \cos p\tau$$

$$\therefore F(\sigma^2 \cos p\tau) = \pi\sigma^2[\delta(\omega+p) + \delta(\omega-p)] \quad (2)$$

Using (2) in (1), we get,

$$S(\omega) = \pi\sigma^2[\delta(\omega+p) + \delta(\omega-p)]$$

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### Example 5

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If  $Y(t) = X(t+a) - X(t-a)$ , prove that  $R_{yy}(\tau) = 2 R_{xx}(\tau) - R_{xx}(\tau+2a) - R_{xx}(\tau-2a)$ . Hence, prove that  $S_{yy}(\omega) = 4 \sin^2 a\omega S_{xx}(\omega)$

[refer to Problem 28 in Exercise 6(b)].

$$R_{yy}(\tau) = 2 R_{xx}(\tau) - R_{xx}(\tau+2a) - R_{xx}(\tau-2a)$$

Taking Fourier transforms on both sides,

$$S_{yy}(\omega) = 2S_{xx}(\omega) - \int_{-\infty}^{\infty} R_{xx}(\tau+2a) e^{-i\omega\tau} d\tau - \int_{-\infty}^{\infty} R_{xx}(\tau-2a) e^{-i\omega\tau} d\tau$$

$$= 2S_{xx}(\omega) - e^{i2a\omega} \int_{-\infty}^{\infty} R_{xx}(u) e^{-i\omega u} du - e^{-i2a\omega} \int_{-\infty}^{\infty} R_{xx}(v) e^{-i\omega v} dv$$

(putting  $\tau+2a=u$  in the first integral and  $\tau-2a=v$  in the second integral)

i.e.,

$$S_{yy}(\omega) = 2 S_{xx}(\omega) - \{e^{i2a\omega} + e^{-i2a\omega}\} S_{xx}(\omega)$$

$$= 2(1 - \cos 2a\omega) S_{xx}(\omega)$$

$$= 4 \sin^2 a\omega S_{xx}(\omega)$$

---

### Example 6

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If the process  $\{X(t)\}$  is defined as  $X(t) = Y(t) Z(t)$ , where  $\{Y(t)\}$  and  $\{Z(t)\}$  are independent WSS processes, prove that

$$(a) \quad R_{xx}(\tau) = R_{yy}(\tau)R_{zz}(\tau) \text{ and}$$

$$(b) \quad S_{xx}(\omega) = \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{yy}(\alpha)S_{zz}(\omega - \alpha) d\alpha$$

[refer to Problem 29(ii) in Exercise 6(b) for part (i)].

$$S_{xx}(\omega) = F\{R_{xx}(\tau)\} = F\{R_{yy}(\tau)R_{zz}(\tau)\} \quad (1)$$

$$\text{Consider } F^{-1}\left[\int_{-\infty}^{\infty} S_{yy}(\alpha)S_{zz}(\omega - \alpha) d\alpha\right]$$

$$= \frac{1}{2\pi} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} S_{yy}(\alpha)S_{zz}(\omega - \alpha) e^{i\omega\tau} d\alpha d\omega$$

Putting  $\alpha = y$  and  $\omega - \alpha = z$ , we get (from calculus)

$$d\alpha d\omega = \begin{vmatrix} \alpha_y & \alpha_z \\ \omega_y & \omega_z \end{vmatrix} dy dz = \begin{vmatrix} 1 & 0 \\ 1 & 1 \end{vmatrix} dy dz$$

$$\begin{aligned} \therefore F^{-1}\left[\int_{-\infty}^{\infty} S_{yy}(\alpha)S_{zz}(\omega - \alpha) d\alpha\right] \\ &= \frac{1}{2\pi} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} S_{yy}(y)S_{zz}(z) e^{i(y+z)\tau} dy dz \\ &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{yy}(y) e^{iy\tau} dy \int_{-\infty}^{\infty} S_{zz}(z) e^{iz\tau} dz \\ &= F^{-1}\{S_{yy}(\omega)\} \times 2\pi F^{-1}\{S_{zz}(\omega)\} \\ &= 2\pi R_{yy}(\tau) R_{zz}(\tau) \end{aligned}$$

$$\therefore F\{R_{yy}(\tau)R_{zz}(\tau)\} = \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{yy}(\alpha)S_{zz}(\omega - \alpha) d\alpha \quad (2)$$

Using (2) in (1), we get  $S_{xx}(\omega)$  in the required form.

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### Example 7

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If the power spectral density of a WSS process is given by

$$S(\omega) = \begin{cases} \frac{b}{a}(a - |\omega|), & |\omega| \leq a \\ 0, & |\omega| > a \end{cases}$$

find the autocorrelation function of the process.

The autocorrelation function  $R(\tau)$  is given by

$$\begin{aligned}
 R(\tau) &= F^{-1}\{S(\omega)\} \\
 &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) e^{i\tau\omega} d\omega \\
 &= \frac{1}{2\pi} \int_{-a}^a \frac{b}{a}(a - |\omega|) e^{i\tau\omega} d\omega \\
 &= \frac{1}{\pi} \int_0^a \frac{b}{a}(a - \omega) \cos \tau\omega d\omega \\
 &= \frac{b}{\pi a} \left\{ (a - \omega) \frac{\sin \tau\omega}{\tau} - \frac{\cos \tau\omega}{\tau^2} \right\}_0^a \\
 &= \frac{b}{\pi a \tau^2} (1 - \cos a\tau) \\
 &= \frac{ab}{2\pi} \left( \frac{\sin a\frac{\tau}{2}}{a\frac{\tau}{2}} \right)^2
 \end{aligned}$$

---

### Example 8

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The power spectrum of a WSS process  $\{X(t)\}$  is given by  $S(\omega) = \frac{1}{(1+\omega^2)^2}$ .

Find its autocorrelation function  $R(\tau)$  and average power.

$$R(\tau) = \frac{1}{2\pi} \int_{-\infty}^{\infty} \frac{1}{(1+\omega^2)^2} e^{i\tau\omega} d\omega \quad (1)$$

The integral in (1) is evaluated by contour integration technique as given below.

Consider  $\int_C \frac{e^{iaz}}{(1+z^2)} dz$ , where  $C$  is the closed contour consisting of the real axis from  $-R$  to  $+R$  and the upper half of the circle  $|z| = R$ .

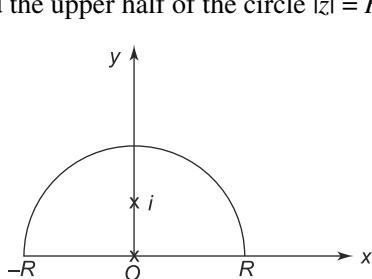


Fig. 6.3

The only singularity of the integrand lying within  $C$  is the double pole  $z = i$ .

$$(\text{Residue})_{z=i} = \frac{1}{[1]} \left[ \frac{d}{dz} \left\{ \frac{e^{iaz}}{(z+i)^2} \right\} \right]_{z=i} = \frac{-i}{4}(1+a)e^{-a}$$

Using Cauchy's residue theorem, taking limits as  $R \rightarrow \infty$  and using Jorden's lemma, we get

$$\int_{-\infty}^{\infty} \frac{e^{iax}}{(1+x^2)^2} dx = \frac{\pi}{2}(1+a)e^{-a} \quad (2)$$

Using (2) in (1),

$$R(\tau) = \frac{1}{4}(1+\tau)e^{-\tau}$$

Average power of

$$\{X(\tau)\} = R(0) = 0.25.$$

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### Example 9

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The power spectral density function of a zero mean WSS process  $\{X(t)\}$  is given by

$$S(\omega) = \begin{cases} 1, & |\omega| < \omega_0 \\ 0, & \text{elsewhere} \end{cases}$$

Find  $R(\tau)$  and show also that  $X(t)$  and  $X\left(t + \frac{\tau}{\omega_0}\right)$  are uncorrelated.

$$R(\tau) = \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) e^{i\tau\omega} d\omega$$

$$= \frac{1}{2\pi} \int_{-\omega_0}^{\omega_0} e^{i\tau\omega} d\omega$$

$$= \frac{1}{2\pi} \left( \frac{e^{i\tau\omega}}{i\tau} \right) \Big|_{-\omega_0}^{\omega_0}$$

$$= \frac{1}{\pi\tau} \sin \omega_0 \tau$$

$$\therefore E \left\{ X \left( t + \frac{\pi}{\omega_0} \right) X(t) \right\} = R \left( \frac{\pi}{\omega_0} \right) = \frac{1}{\pi\tau} \sin \pi$$

Since the mean of the process is zero,

$$C \left\{ X \left( t + \frac{\pi}{\omega_0} \right) X(t) \right\} = E \left\{ X \left( t + \frac{\pi}{\omega_0} \right) X(t) \right\} = 0$$

Therefore,  $X(t)$  and  $X\left(t + \frac{\pi}{\omega_0}\right)$  are uncorrelated.

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### Example 10

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The short-time moving average of a process  $\{X(t)\}$  is defined as  $Y(t) = \frac{1}{T} \int_{t-T}^t X(s) ds$ .

Prove that  $X(t)$  and  $Y(t)$  are related by means of a convolution type integral. Find the unit impulse response of the system also.

$$Y(t) = \frac{1}{T} \int_{t-T}^t X(s) ds \quad (1)$$

Putting  $s = t - u$  and treating  $t$  as a parameter, (1) becomes

$$Y(t) = \frac{1}{T} \int_0^T X(t-u) du \quad (2)$$

Let us define the unit impulse response of the system as follows:

$$h(t) = \begin{cases} \frac{1}{T}, & \text{for } 0 \leq t \leq T \\ 0, & \text{otherwise} \end{cases}$$

Then (2) can be expressed as

$$Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$$

which is a convolution type integral.

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### Example 11

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If the input  $x(t)$  and the output  $y(t)$  are connected by the differential equation  $T \frac{dy(t)}{dt} + y(t) = x(t)$ , prove that they can be related by means of a convolution type integral. Assume that  $x(t)$  and  $y(t)$  are zero for  $t \leq 0$ .

The given differential equation  $y'(t) + \frac{1}{T} y(t) = \frac{1}{T} x(t)$  is a linear equation. Its solution is

$$y(t)e^{t/T} = \int \frac{1}{T} x(u) e^{u/T} du + c$$

i.e.,  $y(t) = \frac{1}{T} \int x(u) e^{-(t-u)/T} du + c$

Since  $y(0) = 0$ ,

$$y(t) = \frac{1}{T} \int_0^t x(u) e^{-(t-u)/T} du$$

$$(or) \quad y(t) = \frac{1}{T} \int_0^t x(t-u) e^{-u/T} du \quad (1)$$

**Given:**

$$x(t) = 0, \text{ for } t < 0$$

$$\therefore x(t-u) = 0, \text{ for } t < u$$

$\therefore$  (1) can be written as

$$y(t) = \frac{1}{T} \int_0^\infty x(t-u) e^{-u/T} du \quad (2)$$

Now, if we define

$$h(t) = \begin{cases} \frac{1}{T} e^{-t/T}, & \text{for } t \geq 0 \\ 0, & \text{otherwise} \end{cases}$$

(2) can be rewritten as

$$y(t) = \int_{-\infty}^{\infty} h(t) x(t-u) du$$

Hence, the result.

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### Example 12

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$X(t)$  is the input voltage to a circuit (system) and  $Y(t)$  is the output voltage.  $\{X(t)\}$  is a stationary random process with  $\mu_x = 0$  and  $R_{xx}(\tau) = e^{-\alpha|\tau|}$ . Find  $\mu_y$ ,  $S_{yy}(\omega)$  and  $R_{yy}(\tau)$ , if the power transfer function is

$$\begin{aligned} H(\omega) &= \frac{R}{R + iL\omega} \\ Y(t) &= \int_{-\infty}^{\infty} h(\alpha) X(t-\alpha) d\alpha \\ \therefore E\{Y(t)\} &= \int_{-\infty}^{\infty} h(\alpha) E\{X(t-\alpha)\} d\alpha \\ &= 0 \quad [\text{since } E\{X(t-\alpha)\} = \mu_x = 0] \\ S_{xx}(\omega) &= \int_{-\infty}^{\infty} R_{xx}(\tau) e^{-i\omega\tau} d\tau \\ &= \int_{-\infty}^0 e^{\alpha\tau} e^{-i\omega\tau} d\tau + \int_0^{\infty} e^{-\alpha\tau} e^{-i\omega\tau} d\tau \\ &= \left\{ \frac{e^{(\alpha-i\omega)t}}{\alpha-i\omega} \right\}_{-\infty}^0 + \left\{ \frac{e^{-(\alpha+i\omega)\tau}}{-(\alpha+i\omega)} \right\}_0^{\infty} \end{aligned}$$

$$= \frac{1}{\alpha - i\omega} + \frac{1}{\alpha + i\omega} = \frac{2\alpha}{\alpha^2 + \omega^2}$$

Now,  $S_{yy}(\omega) = S_{xx}(\omega) = |H(\omega)|^2$

$$\begin{aligned} &= \frac{2\alpha}{\alpha^2 + \omega^2} \times \frac{R^2}{R^2 + L^2 \omega^2} \\ &= \frac{\{(2\alpha R^2 / (R^2 - L^2 \alpha^2))\}}{\alpha^2 + \omega^2} + \frac{\{2\alpha R^2 / (\alpha^2 - R^2 / L^2)\}}{R^2 + L^2 \omega^2} \\ &= \frac{2\alpha \left(\frac{R}{L}\right)^2}{\left(\frac{R}{L}\right)^2 - \alpha^2} \times \frac{1}{\alpha^2 + \omega^2} + \frac{2\alpha R^2 / L^2}{\alpha^2 - \left(\frac{R}{L}\right)^2} \times \frac{1}{\left(\frac{R}{L}\right)^2 + \omega^2} \\ &= \lambda \frac{1}{\alpha^2 + \omega^2} + \mu \frac{1}{\left(\frac{R}{L}\right)^2 + \omega^2}, \text{ say} \end{aligned}$$

$$\therefore R_{yy}(\tau) = \frac{\lambda}{2\pi} \int_{-\infty}^{\infty} \frac{e^{i\tau\omega}}{\alpha^2 + \omega^2} d\omega + \frac{\mu}{2\pi} \int_{-\infty}^{\infty} \frac{e^{i\tau\omega}}{\left(\frac{R}{L}\right)^2 + \omega^2} d\omega \quad (1)$$

We can prove that, by contour integration technique,

$$\int_{-\infty}^{\infty} \frac{e^{iaz}}{z^2 + b^2} dz = \frac{\pi}{b} e^{-ab}; a > 0 \quad (2)$$

Using (2) in (1),

$$R_{yy}(\tau) = \frac{\left(\frac{R}{L}\right)^2}{\left(\frac{R}{L}\right)^2 - \alpha^2} e^{-\alpha|\tau|} + \frac{\left(\frac{R}{L}\right)^2 \alpha}{\alpha^2 - \left(\frac{R}{L}\right)^2} e^{-\left(\frac{R}{L}\right)|\tau|}$$

### Example 13

Given that  $Y(t) = \frac{1}{2\epsilon} \int_{t-\epsilon}^{t+\epsilon} X(\alpha) d\alpha$ , where  $\{X(t)\}$  is a WSS process, prove that

$$S_{yy}(\omega) = \frac{\sin^2 \epsilon \omega}{\epsilon^2 \omega^2} S_{xx}(\omega).$$

Hence, find the relation between  $R_{xx}(\tau)$  and  $R_{yy}(\tau)$ .

Putting  $\alpha = t - u$ , we get  $Y(t) = \frac{1}{2\epsilon} \int_{-\epsilon}^{\epsilon} X(t-u) du$ .

If we define  $h(t)$  as follows:

$$h(t) = \begin{cases} \frac{1}{2\epsilon}, & \text{for } |t| \leq \epsilon \\ 0, & \text{for } |t| > \epsilon \end{cases}$$

then

$$Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$$

$$\therefore S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega), \text{ where}$$

$$H(\omega) = F\{h(t)\}$$

$$= \int_{-\epsilon}^{\epsilon} \frac{1}{2\epsilon} e^{-i\omega\tau} dt = \frac{\sin \epsilon \omega}{\epsilon \omega}$$

i.e.,

$$S_{yy}(\omega) = \frac{\sin^2 \epsilon \omega}{(\epsilon \omega)^2} S_{xx}(\omega)$$

$\therefore$

$$\begin{aligned} R_{yy}(\tau) &= F^{-1} \left\{ \frac{\sin^2 \epsilon \omega}{\epsilon^2 \omega^2} S_{xx}(\omega) \right\} \\ &= F^{-1} \left\{ \frac{\sin^2 \epsilon \omega}{\epsilon^2 \omega^2} \right\} * R_{xx}(\tau) \end{aligned} \quad (1)$$

We can prove that

$$\text{if } R(\tau) = \begin{cases} 1 - \frac{|\tau|}{2\epsilon} & \text{if } |\tau| \leq 2\epsilon \\ 0 & \text{if } |\tau| > 2\epsilon \end{cases}$$

then

$$S(\omega) = 2\epsilon \times \frac{\sin^2 \epsilon \omega}{\epsilon^2 \omega^2} \quad [\text{refer to Problem 35 in Exercise 6(c)}]$$

$$\therefore F^{-1} \left( \frac{\sin^2 \epsilon \omega}{\epsilon^2 \omega^2} \right) = \begin{cases} \frac{1}{2\epsilon} \left( 1 - \frac{|\tau|}{2\epsilon} \right) & \text{if } |\tau| \leq 2\epsilon \\ 0 & \text{if } |\tau| > 2\epsilon \end{cases} \quad (2)$$

Using (2) in (1),

$$R_{yy}(\tau) = \frac{1}{2\epsilon} \int_{-2\epsilon}^{2\epsilon} \left( 1 - \frac{|u|}{2\epsilon} \right) R_{xx}(\tau-u) du$$

### Example 14

**Property (4) of power spectral density.** Prove that the spectral density of any WSS process is non-negative, i.e.  $S(\omega) \geq 0$ .

If possible, let  $S(\omega) < 0$  at  $\omega = \omega_0$ . That is, let  $S(\omega) < 0$  in  $\omega_0 - \frac{\epsilon}{2} < \omega < \omega_0 + \frac{\epsilon}{2}$ ,

where  $\epsilon$  is very small. Let us assume that the system function of the convolution type linear system is

$$H(\omega) = \begin{cases} 1, & \omega_0 - \frac{\epsilon}{2} < \omega < \omega_0 + \frac{\epsilon}{2} \\ 0, & \text{elsewhere} \end{cases}$$

**Note** In this case, the system is called a narrow band filter.

$$\text{Now, } S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

$$= \begin{cases} S_{xx}(\omega), & \omega_0 - \frac{\epsilon}{2} < \omega < \omega_0 + \frac{\epsilon}{2} \\ 0, & \text{elsewhere} \end{cases}$$

$$E\{Y^2(t)\} = R_{yy}(0)$$

$$= \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{yy}(\omega) d\omega$$

$$= \frac{1}{2\pi} \int_{\omega_0 - \frac{\epsilon}{2}}^{\omega_0 + \frac{\epsilon}{2}} S_{xx}(\omega) d\omega$$

$$= \frac{\epsilon}{2\pi} S_{xx}(\omega_0)$$

[since  $S_{xx}(\omega)$  can be considered a constant  $S_{xx}(\omega_0)$ , as the band is narrow]

Since  $E\{Y^2(t)\} \geq 0$ ,  $S_{xx}(\omega) \geq 0$ , which is contrary to our initial assumption. Therefore,  $S_{xx}(\omega) \geq 0$ , everywhere, since  $\omega = \omega_0$  is arbitrary.

### Exercise 6(C)

#### Part-A (Short-answer Questions)

1. Define the power spectral density (PSD) function of a stationary process.
2. Express each of ACF and PSD of a stationary process in terms of the other.
3. Write down the Wiener-Kinchine relations.
4. Define the cross PSD of the random processes  $\{X(t)\}$  and  $\{Y(t)\}$ .
5. State any 2 properties of the PSD function of a stationary process.
6. What is average power of a WSS process  $\{X(t)\}$  and express it in terms of the PSD of the process.
7. Find the mean-square value (or the average power) of the process  $\{X(t)\}$ , if its ACF is given by  $R(\tau) = e^{-\tau^2/2}$ .
8. Prove that the PSD of a real stationary process is an even function.
9. Prove that the PSD of a real or complex stationary process is a real function of  $\omega$ .
10. Prove that the PSD of a real WSS process is twice the Fourier cosine transform of its ACF.

11. Prove that the ACF of a real WSS process is half the Fourier inverse cosine transform of its PSDF.

12. Find the PSD function of a stationary process whose ACF is  $e^{-|t|}$ .

13. Find the ACF of a stationary process whose PSDF is given by

$$S(\omega) = \begin{cases} \omega^2, & \text{for } |\omega| \leq 1 \\ 0, & \text{for } |\omega| > 1 \end{cases}$$

14. State Wiener-Khinchine theorem.

15. What is the use of Wiener-Khinchine theorem?

16. What do you mean by a system? When is it called (i) a deterministic system, and (ii) a stochastic system?

17. Define a system. When is it called a linear system?

18. Define a system. When is it called a time-invariant system?

19. Define a system. When is it called a memoryless system?

20. Define a system. When is it called a causal system?

21. Define system weighting function.

22. If a system is defined as  $Y(t) = \frac{1}{T} \int_0^T X(t-u) du$ , find its weighting function.

23. What is unit impulse response of a system? Why is it called so?

24. If the input  $X(t)$  of the system  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$  is the unit impulse function, prove that  $Y(t) = h(t)$ .

25. Express  $Y(t) = \frac{1}{2c} \int_{t-c}^{t+c} X(\alpha) d\alpha$  as a convolution type of linear system and hence find the unit impulse response of the system.

26. If a system is defined as  $Y(t) = \frac{1}{T} \int_{-\infty}^{\infty} X(t-u) e^{-u/T} du$ , find its unit impulse response.

27. Prove that the system  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$  is a linear time-invariant system.

28. When is the system  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$  said to be stable?

29. If  $\{X(t)\}$  and  $\{Y(t)\}$  in the system  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$  are WSS processes, how are their ACF's related?

30. If the input and output of the system  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$  are WSS processes, how are their PSDF's related?
31. Define the power transfer function (or system function) of the system 
$$Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du.$$
32. If the system function of a convolution type of linear system is given by 
$$h(t) = \begin{cases} \frac{1}{2c}, & \text{for } |t| \leq c \\ 0, & \text{for } |t| > c \end{cases}$$

Find the relation between PSDF's of the input and output processes.

### Part-B

33. Calculate the power spectral density of a stationary random process for which the autocorrelation is  $R_{xx}(\tau) = e^{-\alpha|\tau|}$ .
34. If the autocorrelation function of a WSS process is  $R(\tau) = \rho e^{-\rho|\tau|}$ , show that its spectral density is given by  $S(\omega) = \frac{2}{1 + \left(\frac{\omega}{\rho}\right)^2}$ .
35. Find the power spectral density of the random binary transmission process whose autocorrelation function is 
$$R(\tau) = \begin{cases} 1 - \frac{|\tau|}{T}, & \text{for } |\tau| \leq T \\ 0, & \text{for elsewhere} \end{cases}$$
36. Find the power spectral density of the random process  $\{X(t)\}$ , if  $E\{X(t)\} = 1$  and  $R_{xx}(\tau) = 1 + e^{-\alpha|\tau|}$ .
37. If  $\{X(t)\}$  is a constant random process with  $R(\tau) = m^2$  for all  $\tau$ , where  $m$  is a constant, show that the spectral density of the process is  $S(\omega) = 2\pi m^2 \delta(\omega)$ .
38. Find the power spectral density of the random process whose autocorrelation function is  $R(\tau) = e^{-\alpha\tau^2} \cos \omega_0 \tau$ .
39. Find the power spectral density of the random process if its autocorrelation function is  $R(\tau) = e^{-\alpha|\tau|} \cos \beta \tau$ .
40. For the process  $\{X(t)\}$ , where  $X(t) = \sum_{i=1}^n (A_i \cos p_i t + B_i \sin p_i t)$ , where  $A_i$  and  $B_i$  are uncorrelated RVs with mean zero and variance  $\sigma_i^2$ , show that the autocorrelation function is given by  $R(\tau) = \sum_{i=1}^n \sigma_i^2 \cos p_i \tau$ . Prove also, that the power spectrum for this process is given by

$$S(\omega) = \pi \sum_{i=1}^n \sigma_i^2 [\delta(\omega - p_i) + \delta(\omega + p_i)]$$

**(Hint:** Refer to Worked Example 4 in Section 6(c)).

41. (i) For the process  $\{X(t)\}$ , where  $X(t) = a \cos(bt + Y)$ , where  $Y$  is uniformly distributed over  $(-\pi, \pi)$  find the autocorrelation function and the spectral density.  
(ii) For the process  $\{X(t)\}$ , where  $X(t) = a \sin(bt + Y)$ , where  $Y$  is uniformly distributed over  $(0, 2\pi)$ , find the autocorrelation function and the spectral density.
42. Find the autocorrelation function of the process  $\{X(t)\}$ , for which the power spectral density is given by

$$S(\omega) = \begin{cases} 1 + \omega^2, & \text{for } |\omega| \leq 1 \\ 0, & \text{for } |\omega| > 1 \end{cases}$$

43. Find the average power of the random process  $\{X(t)\}$ , if its power spectral density is given by

$$S(\omega) = \frac{10\omega^2 + 35}{(\omega^2 + 4)(\omega^2 + 9)}$$

- (i) using  $S(\omega)$  directly, and  
(ii) using the autocorrelation function  $R(\tau)$ .

**[Hint:** Average power =  $\frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) d\omega$  or  $R(0)$ .]

44. Find the autocorrelation function of the random process  $\{X(t)\}$  for which the power spectral density is given by  
(i)  $S(\omega) = 2\pi\delta(\omega) + \frac{2\alpha}{\alpha^2 + \omega^2}$ , and (ii)  $S(\omega) = \frac{1}{\omega^2 + 4}$
45.  $\{X(t)\}$  is a stationary random process with spectral density function  $S_{xx}(\omega)$  and  $\{Y(t)\}$  is another independent random process where  $Y(t) = A \cos(\omega_c t + \theta)$ , where  $\theta$  is a RV uniformly distributed over  $(-\pi, \pi)$ . Find the spectral density function of  $\{Z(t)\}$ , where  $Z(t) = X(t) Y(t)$ .
46. If  $\{Y(t)\}$  is the moving time average of  $\{X(t)\}$  over  $\{t - T, t + T\}$ , express  $S_{yy}(\omega)$  in terms of  $S_{xx}(\omega)$ . Hence, find the autocorrelation function of  $\{Y(t)\}$  in terms of that of  $\{X(t)\}$ .
47.  $X(t)$  is the input voltage to a circuit and  $Y(t)$  is the output voltage.  $\{X(t)\}$  is a stationary random process with  $\mu_x = 0$  and  $R_{xx}(\tau) = e^{-2|\tau|}$ . Find  $\mu_y$ ,  $S_{yy}(\omega)$  and  $R_{yy}(\tau)$ , if the system function is given by  $H(\omega) = \frac{1}{\omega + i^2}$ .

## ANSWERS

### **Exercise 6(A)**

4. (i) When  $s$  is fixed,  $X(s, t)$  is a time function. (ii) When  $t$  is fixed,  $X(s, t)$  is a RV.
10. The jpdf's and the joint distribution functions of the member functions of the random process are used to describe it. The first- and second-order distribution functions, namely,  $F(x, t) = P\{X(t) \leq x\}$  and  $F(x_1, x_2; t_1, t_2) = P\{X(t_1) \leq x_1; X(t_2) \leq x_2\}$  are primarily used to describe the random process  $\{X(t)\}$ .
11. (i) Markov process (ii) Process with independent increments and (iii) Stationary process.
12. Poisson process and Wiener process.
13. If  $X(t)$  is a representative member function of the random process  $\{X(t)\}$ ,  $E\{X(t)\}$  and  $\text{Var}\{X(t)\}$  are called the mean and variance of the process.
17. No. Their definitions are entirely different.
19. Two processes  $\{X(t)\}$  and  $\{Y(t)\}$  are said to be orthogonal, if  $E\{X(t_1)Y(t_2)\} = 0$ .
22. Bernoulli's process.
23. A  $k^{\text{th}}$ -order stationary process becomes a SSS process when  $k \rightarrow \infty$ .
32. Poisson process.
35.  $f_{X(t)}(x) = \frac{1}{\sqrt{2\pi\alpha t}} e^{-x^2/2\alpha t}; E\{X(t)\} = 0;$   
 $\text{Var}\{X(t)\} = \alpha t.$
36. No
37. Yes, it is a WSS process.
38.  $p + qt; \sigma_1^2 + p^2 + (t_1 + t_2)pq + t_1t_2(\sigma_2^2 + q^2); \sigma_1^2 + t_1t_2\sigma_2^2; \text{No.}$
39.  $E\{X(t)\} = 0; R(t_1, t_2) = 1/2 \cos \omega(t_1 - t_2)$
40.  $E\{X(t)\} = 0; R(t_1, t_2) = 2 \cos (t_1 - t_2)$
41.  $R(t_1, t_2) = (A^2/2) \cos \omega(t_1 - t_2)$
42.  $E\{X(t)\} = 0; E\{X^2(t)\} = 1/2; \text{No}$
44.  $E\{X(t)\} = t + 1/2 \sin \pi t; \text{when } t = 0.25, F(x) = 0, \text{if } x < \frac{1}{2}, = \frac{1}{2} \text{ if } \frac{1}{2} \leq x < \frac{1}{\sqrt{2}}$   
 $\text{and} = 1, \text{if } \frac{1}{\sqrt{2}} \leq x; \text{when } t = 0.5,$   
 $F(x) = 0, \text{if } x < 1 \text{ and} = 1, \text{if } x \geq 1; \text{when } t = 1,$   
 $F(x) = 0, \text{if } x < 0, = 1/2, \text{if } 0 \leq x < 2 \text{ and} = 1, \text{if } x \geq 2.$
45.  $E\{X(t)\} = 0; R(t_1, t_2) = \sigma^2 \cos (t_1 - t_2)$
46.  $E(z) = 0; E(z^2) = \sigma^2$
47.  $E\{Z\} = E\{w\} = 3; \text{Var}(z) = \text{Var}(w) = 13; \text{cov}(z, w) = 2.195.$
49.  $R_{xy}(t_1, t_2) = \sin (t_1 + t_2)$
50.  $R_{xy}(t_1, t_2) = 50 \sin 10 (t_1 - t_2)$

51.  $R(t_1, t_2) = \frac{A^4}{8} \{2 + \cos 2\omega(t_1 - t_2)\}$

52.  $\phi(\omega) = e^{-\alpha t \omega^2 / 2}$

53.  $R_{xx}(\tau) = R_{yy}(\tau)$  and  $R_{xy}(-\tau) = -R_{xy}(\tau); R_{ww}(\tau) = R_{xx}(\tau) \cos \omega\tau + R_{yx}(\tau) \sin \omega\tau$  (or)  $R_{xx}(\tau) \cos \omega\tau - R_{xy}(\tau) \sin \omega\tau$

### Exercise 6(B)

5. No, it is true only when the stationary process does not contain periodic components.

6.  $\mu_x^2 = \lim_{\tau \rightarrow \infty} \left( \frac{25 + 36/\tau^2}{6.25 + 4/\tau^2} \right) = 4; \mu_x = 2$

7.  $E\{X^2(t)\} = R(0) = 9;$   
 $\text{Var}\{X(t)\} = E\{X^2(t)\} - E^2\{X(t)\} = 5$

8.  $E\{X(t)\} = \sqrt{2}; E\{X^2(t)\} = 6; \text{Var}\{X(t)\} = 4$

9.  $E\{X(t)\} = 4; E\{X^2(t)\} = 25; \text{Var}\{X(t)\} = 9$

11.  $R_{xy}(\tau) = E[X(t)]Y(t - \tau) = 0$

12.  $R_{xy}(\tau) = E[X(t)]Y(t - \tau) = E[X(t)]E[Y(t - \tau)] = \mu_x \mu_y$

13.  $E\{X(t)\} = \sum x_i p_i$  or  $\int_{R_X} xf(x) dx; \bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt$

14. The ensemble average is given by  $E\{X(t)\} = \sum x_i p_i$  or  $\int_{R_X} xf(x) dx$ ; the

time average is given by  $\bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt$ . To compute  $E\{X(t)\}$ , we should know the probability distribution or density function of  $X(t)$ ; To compute  $\bar{X}_T$ , it is enough we know a single sample function of the process.

16. Stationary of a random process is the property of the process by which certain probability distributions or averages do not depend on  $t$ .

Ergodicity of a random process is the property by which almost every member of the process exhibits the same statistical behaviour as the whole process. Ergodicity is a stronger condition than stationarity and hence all stationary processes are not ergodic.

20. (i)  $\lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} = 0$  where  $\bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt$

(ii)  $\lim_{T \rightarrow \infty} \left[ \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T C(t_1, t_2) dt_1 dt_2 \right] = 0$

21. No;  $\bar{X}_T = \frac{1}{2T} \int_{-T}^T X \, dt = X$ ;  $\text{Var}(X_T) = \text{Var}(X) \neq 0$ ;  $\lim_{T \rightarrow \infty} \{\text{Var}(\bar{X}_T)\} \neq 0$
22. Random binary transmission process.
23.  $\text{Var}(\bar{X}_T) = \frac{1}{T} \int_0^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau$ ;  

$$\lim_{T \rightarrow \infty} \left[ \frac{1}{T} \int_0^{2T} C(\tau) \left\{ 1 - \frac{|\tau|}{2T} \right\} d\tau \right] = 0$$
26.  $\pm 2, 9, 5$
27.  $R_{yy}(t_1, t_2) = R(t_1 + a, t_2 + a) - R(t_1 + a, t_2)$   
 $- R(t_1, t_2 + a) + R(t_1, t_2)$
29. (i)  $a^2 + b^2 R_{xx}(\tau) + c^2 R_{yy}(\tau)$   
(ii)  $a^2 R_{xx}(\tau) \cdot R_{yy}(\tau)$
31. 80;  $5(e^{-20} + 19)$
32.  $\sqrt{10}; 4 \int_4^5 e^{-x^2/10} \, dx + 4e^{-2.5} - 14$
33. 2;  $20(10e^{-0.1} - 9)$
34. 0;  $\frac{1}{2\lambda T} - \frac{1}{8\lambda^2 T^2} (1 - e^{-4\lambda T})$ ; Yes
35. 1;  $1/2 (1 + e^{-2})$
36.  $\frac{A}{\alpha T} \left\{ 1 - \frac{(1 - e^{-2\alpha T})}{2\alpha T} \right\}$
40. Random binary transmission process.

### Exercise 6(C)

6. Average power of  $\{X(t)\} = E\{X^2(t)\} = R(0)$ ;

$$E\{X^2(t)\} = \int_{-\infty}^{\infty} S(f) \, df \text{ or } \frac{1}{2\pi} \int_{-\infty}^{\infty} S(\omega) \, d\omega$$

7.  $E\{X^2(t)\} = R(0) = 1$

$$12. \quad S(\omega) = \int_{-\infty}^{\infty} e^{-|t|} e^{i\omega t} \, dt = 2 \int_0^{\infty} e^{-t} \cos \omega t \, dt = \frac{2}{1 + \omega^2}$$

$$13. \quad R(\tau) = \frac{1}{2\pi} \int_{-1}^1 \omega^2 \cos \tau \omega \, d\omega = \frac{1}{\pi} \left( \frac{\sin \tau}{\tau} + \frac{2\cos \tau}{\tau^2} - \frac{2\sin \tau}{\tau^3} \right)$$

21. If the output  $Y(t)$  of a system is expressed as the convolution of the input  $X(t)$  and a function  $h(t)$ , i.e.,  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$ , then  $h(t)$  is called the system weighting function.

$$22. h(t) = \begin{cases} \frac{1}{T}, & 0 \leq t \leq T \\ 0, & \text{elsewhere} \end{cases}$$

23. If a system is of the form  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$ , then the system weighting function  $h(t)$  is also called unit impulse response of the system. It is called so because the response (output)  $Y(t)$  will be  $h(t)$ , when the input  $X(t) =$  the unit impulse function  $\delta(t)$ .

25. Put  $\alpha = t - u$ . Then

$$\begin{aligned} Y(t) &= \int_{-c}^c \frac{1}{2c} X(t-u) du \\ &= \int_{-\infty}^{\infty} h(u) X(t-u) du \\ \therefore h(t) &= \begin{cases} \frac{1}{2c}, & \text{in } |t| \leq c \\ 0, & \text{in } |t| > c \end{cases} \end{aligned}$$

$$26. h(t) = \begin{cases} \frac{1}{T} e^{-t/T}, & \text{for } t \geq 0 \\ 0, & \text{elsewhere} \end{cases}$$

$$29. R_{yy}(\tau) = R_{xy}(\tau)*h(\tau),$$

$$\text{where } R_{xy}(\tau) = R_{xx}(\tau)*h(-\tau)$$

$$(\text{or}) \quad R_{yy}(\tau) = R_{xx}(\tau)*h(\tau)*h(-\tau)$$

$$30. S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega), \text{ where } H(\omega) \text{ is the Fourier transform of } h(t)$$

$$32. H(\omega) = \int_{-c}^c \frac{1}{2c} e^{-i\omega t} dt = \frac{\sin c\omega}{c\omega}; S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

$$S_{YY}(\omega) = \frac{\sin^2 c\omega}{c^2 \omega^2} S_{xx}(\omega)$$

$$33. \frac{2\alpha}{\alpha^2 + \omega^2}$$

35. 
$$4 \sin^2\left(\frac{\omega T}{2}\right) / T \omega^2$$

36. 
$$2\pi\delta(\omega) + \frac{2\alpha}{\omega^2 + \alpha^2}$$

38. 
$$\sqrt{\frac{\pi}{4\alpha}} \{e^{-(\omega+\omega_0)^2/4\alpha} + e^{-(\omega-\omega_0)^2/4\alpha}\}$$

39. 
$$\frac{\alpha}{\alpha^2 + (\omega - \beta)^2} + \frac{\alpha}{\alpha^2 + (\omega + \beta)^2}$$

41. (i) and (ii)  $\frac{a^2}{2} \cos b\tau; \frac{\pi\alpha^2}{2} \cdot [\delta(\omega - b) + (\omega + b)]$

42. 
$$\frac{1}{\pi\tau^3} [\tau^2 \sin \tau + \tau \cos \tau - \sin \tau]$$

43. 
$$R(\tau) = \frac{11}{6}e^{-3|\tau|} - \frac{1}{4}e^{-2|\tau|}; \text{ Average power} = \frac{19}{12}$$

44. (i)  $1 + e^{-\alpha|\tau|}$

(ii)  $\frac{1}{8}e^{-|\tau|}(\cos \tau + \sin |\tau|)$

45. 
$$\frac{A^2}{4} \{S_{xx}(\omega - \omega_c) + S_{xx}(\omega + \omega_c)\}$$

46. 
$$S_{yy}(\omega) = \frac{\sin^2 \omega T}{\omega^2 T^2} S_{xx}(\omega);$$

$$R_{yy}(\tau) = \frac{1}{2T} \int_{-2T}^{2T} \left(1 - \frac{|\alpha|}{2T}\right) R_{xx}(\tau - \alpha) d\alpha$$

47. 
$$0; \frac{4}{(\omega^2 + 4)^2}; \frac{1}{8}(1 + 2|\tau|)e^{-2|\tau|}$$



# *Chapter* 7

## Special Random Processes

Many random phenomena in physical problems including ‘noise’ are well approximated by a special class of random process, namely Gaussian random process. A number of processes such as the Wiener process and the shot-noise process can be approximated, as per central limit theorem, by a Gaussian process. Moreover the output of a linear system in which the input is a weighted sum of a large number of independent samples of a random process tends to approach a Gaussian process. Gaussian processes play an important role in the theory and analysis of random phenomena, because they are good approximations to the observations, and multivariate Gaussian distributions are analytically simple.

One of the most important uses of the Gaussian process is to model and analyse the effects of thermal noise in electronic circuits used in communication systems. Individual circuits contain resistors, inductors and capacitors as well as semiconductor devices. The resistors and semiconductor elements contain charged particles (free electrons) subjected to random motion due to thermal energy. The random motion of charged particles causes fluctuations in the current waveforms or information bearing signals that flow through these components. These fluctuations are called thermal noise, which are of sufficient strength to disturb a weak signal and to make the recognition of signals a difficult task. Models of thermal noise are used to identify and minimise the effects of noise in signal recognition.

### **Definition of a Gaussian Process**

A real-valued random process  $\{X(t)\}$  is called a Gaussian process or normal process, if the random variables  $X(t_1), X(t_2), \dots, X(t_n)$  are jointly normal for every  $n = 1, 2, \dots$  and for any set of  $t_i$ 's.

The  $n$ th order density of a Gaussian process is given by  $f(x_1, x_2, \dots, x_n; t_1, t_2, \dots, t_n)$

$$= \frac{1}{(2\pi)^{n/2} |\Lambda|^{1/2}} \exp \left[ -\frac{1}{2|\Lambda|} \sum_{i=1}^n \sum_{j=1}^n |\Lambda|_{ij} (x_i - \mu_i)(x_j - \mu_j) \right]$$

where  $\mu_i = E\{X_i(t_i)\}$  and  $\Lambda$  is the  $n$ th order square matrix ( $\lambda_{ij}$ ), where  $\lambda_{ij} = C\{X(t_i), X(t_j)\}$  and  $|\Lambda|_{ij}$  = cofactor of  $\lambda_{ij}$  in  $|\Lambda|$  (1)

**Note** Gaussian process is completely specified by the first- and second-order moments, viz., mean and co-variances (variances).

**Note** When we consider the first-order density of a Gaussian process,

$$\begin{aligned} \Lambda &= (\lambda_{11}) = [Cov\{X(t_1), X(t_1)\}] \\ &= [Var\{X(t_1)\}] = (\sigma_1^2) \\ \therefore |\Lambda| &= \sigma_1^2 \text{ and } |\Lambda|_{11} = 1 \\ \therefore f(x_1, t_1) &= \frac{1}{\sigma_1(2\pi)} \exp \{-(x_1 - \mu_1)^2 / 2\sigma_1^2\} \end{aligned}$$

**Note** When we consider the second-order density of a Gaussian process,

$$\begin{aligned} \Lambda &= \begin{pmatrix} \lambda_{11} & \lambda_{12} \\ \lambda_{21} & \lambda_{22} \end{pmatrix} = \begin{pmatrix} \sigma_1^2 & r_{12} \sigma_1 \sigma_2 \\ r_{21} \sigma_1 \sigma_2 & \sigma_2^2 \end{pmatrix} \\ \therefore |\Lambda| &= \sigma_1^2 \sigma_2^2 (1 - r^2), \text{ where } r_{12} = r_{21} = r \\ |\Lambda|_{11} &= \sigma_1^2, |\Lambda|_{12} = -r \sigma_1 \sigma_2, |\Lambda|_{21} = -r \sigma_1 \sigma_2, |\Lambda|_{22} = \sigma_2^2 \\ \therefore f(x_1, x_2, t_1, t_2) &= \frac{1}{2\pi\sigma_1\sigma_2\sqrt{1-r^2}} \exp \left[ -\frac{1}{2\sigma_1^2\sigma_2^2(1-r^2)} \{ \sigma_2^2 (x_1 - \mu_1)^2 \right. \\ &\quad \left. - 2r\sigma_1\sigma_2(x_1 - \mu_1)(x_2 - \mu_2) + \sigma_1^2 (x_2 - \mu_2)^2 \} \right] \end{aligned}$$

$$\begin{aligned} \text{i.e., } f(x_1, x_2, t_1, t_2) &= \frac{1}{2\pi\sigma_1\sigma_2\sqrt{1-r^2}} \exp \\ &\quad \left[ -\frac{1}{2(1-r^2)} \left\{ \frac{(x_1 - \mu_1)^2}{\sigma_1^2} - \frac{2r(x_1 - \mu_1)(x_2 - \mu_2)}{\sigma_1 \sigma_2} + \frac{(x_2 - \mu_2)^2}{\sigma_2^2} \right\} \right] \end{aligned}$$

which we have made use of in many problems earlier.

## Properties

1. If a Gaussian process is wide-sense stationary, it is also strict-sense stationary.

### Proof

The  $n$ th order density of a Gaussian process is given by  $f(x_1, x_2, \dots, x_n; t_1, t_2, \dots, t_n)$

$$= \frac{1}{(2\pi)^{n/2} |\Lambda|^{1/2}} \exp \left[ -\frac{1}{2|\Lambda|} \sum_{i=1}^n \sum_{j=1}^n |\Lambda|_{ij} (x_i - \mu_i)(x_j - \mu_j) \right]$$

where  $\mu_i = E(X(t_i))$  and  $\Lambda$  is the  $n$ th order square matrix  $(\lambda_{ij})$ ,  $i, j = 1, 2, \dots, n$ , where  $\lambda_{ij} = C\{X(t_i), X(t_j)\}$  and  $|\Lambda|_{ij}$  is the cofactor of  $\lambda_{ij}$  in  $|\Lambda|$ .

If the Gaussian process is WSS, then  $\lambda_{ij} = C\{X(t_i), X(t_j)\}$  is a function of  $(t_i - t_j)$  for all  $i$  and  $j$ .

Therefore, the  $n$ th order densities of  $\{X(t_1), X(t_2), \dots, X(t_n)\}$  and  $\{X(t_1 + h), X(t_2 + h), \dots, X(t_n + h)\}$  are identical.

Therefore, the Gaussian process is a SSS process.

2. If the member functions of a Gaussian process are uncorrelated, then they are independent.

### Proof

Consider  $n$  member functions  $X(t_1), X(t_2), \dots, X(t_n)$  of the Gaussian process  $\{X(t)\}$ .

Since the members are uncorrelated,

$$\begin{aligned} C\{X(t_i), X(t_j)\} &= \lambda_{ij} = 0, & \text{for } i \neq j \\ &= \sigma_i^2, & \text{for } i = j \end{aligned}$$

Therefore,  $[\Lambda]$  is a diagonal matrix with terms in the principal diagonal equal to  $\sigma_i^2$ ,

$$\therefore |\Lambda| = \sigma_1^2 \sigma_2^2 \dots \sigma_n^2; |\Lambda|_{ij} = \sigma_1^2 \sigma_2^2 \dots \sigma_{i-1}^2 \sigma_{i+1}^2 \dots \sigma_n^2$$

$$|\Lambda|_{ij} = 0, \text{ for } i \neq j$$

Hence, the  $n$ th order density function of the Gaussian process becomes  $f(x_1, x_2, \dots, x_n; t_1, t_2, \dots, t_n)$ .

$$\begin{aligned} &= \frac{1}{(2\pi)^{n/2} \sigma_1 \sigma_2 \dots \sigma_n} \exp \left[ -\frac{1}{2} \sum_{j=1}^n (x_j - \mu_j)^2 / \sigma_j^2 \right] \\ &= \left\{ \frac{1}{\sigma_1 \sqrt{2\pi}} e^{-(x_1 - \mu_1)^2 / 2\sigma_1^2} \right\} \left\{ \frac{1}{\sigma_2 \sqrt{2\pi}} e^{-(x_2 - \mu_2)^2 / 2\sigma_2^2} \right\} \dots \\ &\quad \left\{ \frac{1}{\sigma_n \sqrt{2\pi}} e^{-(x_n - \mu_n)^2 / 2\sigma_n^2} \right\} \\ &= f(x_1, t_1) f(x_2, t_2) \dots f(x_n, t_n) \end{aligned}$$

i.e.,  $X(t_1), X(t_2), \dots, X(t_n)$  are independent.

3. If the input  $\{X(t)\}$  of a linear system is a Gaussian process, the output will also be a Gaussian process.

### Proof

Let  $\{Y(t)\}$  be the output process.

$$\begin{aligned} \text{Then } Y(t) &= \int_{-\infty}^{\infty} X(t-u) h(u) du \text{ (or)} \int_{-\infty}^{\infty} h(t-u) X(u) du \\ &= \sum_{j=1}^n \{h(t-u_j) \Delta u\} X(u_j), \text{ as } n \rightarrow \infty \end{aligned}$$

Consider the  $n$  sample functions of the input process  $\{X(t)\}$ . Let them be  $X(t_1), X(t_2), \dots, X(t_n)$ .

Let the corresponding output sample functions be  $Y(t_1), Y(t_2), \dots, Y(t_n)$ .

$$\therefore Y(t_i) = \sum_{j=1}^n \{h(t_i - u_j) \Delta u\} X(u_j) \text{ in the limit}$$

$$\text{i.e., } Y_i = \sum_{j=1}^n h_{ij} X_j, \text{ say, where } i = 1, 2, \dots, n \quad (1)$$

$$\therefore y_i = \sum_{j=1}^n h_{ij} x_j, \quad i = 1, 2, \dots, n \quad (1a)$$

$$\text{i.e., } \begin{pmatrix} y_1 \\ y_2 \\ \vdots \\ y_n \end{pmatrix} = \begin{pmatrix} h_{11} & h_{12} & \cdots & h_{1n} \\ h_{21} & h_{22} & \cdots & h_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ h_{n1} & h_{n2} & \cdots & h_{nn} \end{pmatrix} \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{pmatrix}$$

$$\text{i.e., } \bar{Y} = H \bar{X}, \text{ say}$$

$$\therefore E(\bar{Y}) = H E(\bar{X}) \text{ (or)} \mu_{\bar{Y}} = H \mu_{\bar{X}}$$

$$\therefore f_Y(y_1, y_2, \dots, y_n) = |J(x_1, x_2, \dots, x_n)| f_X(x_1, x_2, \dots, x_n)$$

$$\text{where } |J(x_1, x_2, \dots, x_n)| = |J(y_1, y_2, \dots, y_n)|^{-1}$$

$$\text{Now, } |J(y_1, y_2 \dots y_n)|$$

$$= \begin{vmatrix} \frac{\partial y_1}{\partial x_1} & \frac{\partial y_1}{\partial x_2} & \cdots & \frac{\partial y_1}{\partial x_n} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\partial y_n}{\partial x_1} & \frac{\partial y_n}{\partial x_2} & \cdots & \frac{\partial y_n}{\partial x_n} \end{vmatrix} = \begin{vmatrix} h_{11} & h_{12} & \cdots & h_{1n} \\ h_{21} & h_{22} & \cdots & h_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ h_{n1} & h_{n2} & \cdots & h_{nn} \end{vmatrix} = |H|$$

Since  $\frac{\partial y_1}{\partial x_j} = h_{ij}$  ( $i, j = 1, 2, \dots, n$ ) [from (1)]

$$\begin{aligned}\therefore f_Y(y_1, y_2, \dots, y_n) &= \frac{1}{|H|} f_X(x_1, x_2, \dots, x_n) \\ &= \frac{1}{|H|(2\pi)^{n/2} |\Lambda_x|^{1/2}} \exp \left[ -\frac{1}{2|\Lambda_x|} \sum_{i=1}^n \sum_{j=1}^n |\Lambda_x|_{ij} (x_i - \mu_i)(x_j - \mu_j) \right] \quad (2)\end{aligned}$$

where  $\Lambda_x$  is the square matrix whose elements are the covariance functions of  $\{X(t)\}$ , viz.

$$\Lambda_x = \begin{pmatrix} c_{11} & c_{12} & \cdots & c_{1n} \\ c_{21} & c_{22} & \cdots & c_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ c_{n1} & c_{n2} & \cdots & c_{nn} \end{pmatrix}$$

where  $c_{ij} = C\{X(t_i), X(t_j)\}$ .

If we express the RS of (2) in matrix form, then  $f_Y(y_1, y_2, \dots, y_n)$

$$\begin{aligned}&= \frac{1}{(2\pi)^{n/2} |H| |\Lambda_x|^{1/2}} \exp \left[ -\frac{1}{2} (\bar{X} - \mu_{\bar{X}})^T \Lambda_x^{-1} (\bar{X} - \mu_{\bar{X}}) \right] \\ &= \frac{1}{(2\pi)^{n/2} |H| |\Lambda_x|^{1/2}} \exp \left[ -\frac{1}{2} (H^{-1} \bar{Y} - H^{-1} \mu_{\bar{Y}})^T \Lambda_x^{-1} (H^{-1} \bar{Y} - H^{-1} \mu_{\bar{Y}}) \right] \\ &= \frac{1}{(2\pi)^{n/2} |H| |\Lambda_x|^{1/2}} \exp \left[ -\frac{1}{2} \{H^{-1}(\bar{Y} - \mu_{\bar{Y}})\}^T \Lambda_x^{-1} \{H^{-1}(\bar{Y} - \mu_{\bar{Y}})\} \right] \\ &= \frac{1}{(2\pi)^{n/2} |H| |\Lambda_x|^{1/2}} \exp \\ &\quad \left[ -\frac{1}{2} (\bar{Y} - \mu_{\bar{Y}})^T \{(H^T)^{-1}\} \Lambda_x^{-1} H^{-1} \} (\bar{Y} - \mu_{\bar{Y}}) \right] \quad (3)\end{aligned}$$

$$\begin{aligned}\text{Now, let } C'_{ij} &= C\{Y(t_i), Y(t_j)\} \\ &= C\{Y_i, Y_j\} \\ &= E[\{Y_i - E(Y_i)\} \{Y_j - E(Y_j)\}] \\ &= E \left[ \sum_{r=1}^n h_{ir} (X_r - E(X_r)) \sum_{s=1}^n h_{js} (X_s - E(X_s)) \right] \quad [\text{by (1)}] \\ &= \sum_{r=1}^n \sum_{s=1}^n h_{ir} h_{is} c_{rs}\end{aligned}$$

Therefore,  $\Lambda_y = H \Lambda_x H^T$ , where  $\Lambda_y$  is the covariance matrix of the  $\{Y(t)\}$  process.

$$\therefore \Lambda_y^{-1} = (H^T)^{-1} \Lambda_x^{-1} H^{-1} \quad (4)$$

$$\text{and } |\Lambda_y| = |H|^2 |\Lambda_x| \quad (5)$$

Using (4) and (5) in (3), we get

$$f_Y(y_1, y_2, \dots, y_n) = \frac{1}{(2\pi)^{n/2} |\Lambda|^{1/2}} \exp \left[ -\frac{1}{2} (\bar{Y} - \mu_{\bar{Y}})^T \Lambda_y^{-1} (\bar{Y} - \mu_{\bar{Y}}) \right]$$

which is the  $n$ th order density function of a Gaussian process.

Therefore, the output process  $\{Y(t)\}$  is a Gaussian process.

**Note** This property can be interpreted as follows:

Since  $Y(t) = \sum_{j=t}^n \{b(t - u_j) \Delta u\} X(u_j)$ ,  $Y(t)$  is a linear combination of  $X(u_j)$ 's.

Since each  $X(u_j)$  follows a normal distribution,  $Y(t)$  also follows a normal distribution.

## Processes Depending on a Stationary Gaussian Process

### 1. Square-Law Detector Process

If  $\{X(t)\}$  is a zero mean stationary Gaussian process and if  $Y(t) = X^2(t)$ , then  $\{Y(t)\}$  is called a square-law detector process.

$$E\{Y(t)\} = E\{X^2(t)\} = \text{Var}\{X(t)\} = R_{xx}(0)$$

$$\begin{aligned} R_{yy}(t_1, t_2) &= E\{Y(t_1) Y(t_2)\} \\ &= E\{X^2(t_1) X^2(t_2)\} \\ &= E\{X^2(t_1)\} E\{X^2(t_2)\} + 2E\{X(t_1) X(t_2)\} \\ &\quad [\text{since when } X \text{ and } Y \text{ are joint normal}] \\ E(X^2 Y^2) &= E(X^2) E(Y^2) + 2E^2(XY) \quad (\text{refer to Chapter 4}) \\ &= R_{xx}^2(0) + 2R_{xx}^2(\tau) \quad [\text{since } X(t) \text{ is stationary}] \end{aligned}$$

Since the RS is function of  $\tau$ , LS is also a function of  $\tau = t_1 - t_2$ .

$$\text{i.e., } R_{yy}(\tau) = R_{xx}^2(0) + 2R_{xx}^2(\tau)$$

Therefore,  $\{Y(t)\}$  is also a stationary process (at least in the wide-sense).

$$\text{We note that } E\{Y^2(t)\} = R_{yy}(0) = 3R_{xx}^2(0)$$

$$\text{Var}\{Y(t)\} = 2R_{xx}^2(0) \text{ and } C_{yy}(\tau) = 2R_{xx}^2(\tau)$$

Power spectral density of  $\{Y(t)\}$  is given by

$$\begin{aligned} S_{yy}(\omega) &= \int_{-\infty}^{\infty} \{R_{xx}^2(0) + 2R_{xx}^2(\tau)\} e^{-i\tau\omega} d\tau \\ &= 2\pi R_{xx}^2(0) \delta(\omega) + 2F\{R_{xx}(\tau) R_{xx}(\tau)\} \\ &\quad [\text{since } F^{-1}\{2\pi m^2 \delta(\omega)\} = m^2] \end{aligned} \quad (1)$$

Consider  $F^{-1}\{s(\omega) * s(\omega)\}$

$$= \frac{1}{2\pi} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} s(\omega - \alpha) s(\alpha) e^{i\tau\omega} d\alpha d\omega \quad (2)$$

Put  $\omega - \alpha = \beta$  and  $\alpha = \gamma$

i.e.,  $\omega = \beta + \gamma$        $\alpha = \gamma$

Then from calculus,

$$d\omega d\alpha = \begin{vmatrix} \frac{\partial \omega}{\partial \beta} & \frac{\partial \omega}{\partial \gamma} \\ \frac{\partial \alpha}{\partial \beta} & \frac{\partial \alpha}{\partial \gamma} \end{vmatrix} d\beta d\gamma = d\beta d\gamma \quad (3)$$

Using (3) in (2),

$$\begin{aligned} F^{-1}\{s(\omega) * s(\omega)\} &= \frac{1}{2\pi} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} s(\beta) s(\gamma) e^{i\tau(\beta+\gamma)} d\beta d\gamma \\ &= 2\pi \frac{1}{2\pi} \int_{-\infty}^{\infty} s(\beta) e^{i\tau\beta} d\beta \frac{1}{2\pi} \int_{-\infty}^{\infty} s(\gamma) e^{i\tau\gamma} d\gamma \\ &= 2\pi R(\tau)R(\tau) \end{aligned} \quad (4)$$

Using (4) in (1),

$$S_{yy}(\omega) = 2\pi R_{xx}^2(0) \delta(\omega) + \frac{1}{\pi} s_{xx}(\omega)^* s_{xx}(\omega)$$

## Two Important Results

We now consider two important results which will be used in the discussion of other processes depending on stationary Gaussian process, that will follow.

- (i) If  $X$  and  $Y$  are two normal RVs with zero means, variances  $\sigma_1^2$  and  $\sigma_2^2$  and correlation coefficient  $r$ , then the probability that they are of the same sign  $= \frac{1}{2} + \frac{1}{\pi} \sin^{-1}(r)$  and the probability that they are of opposite signs  $= \frac{1}{2} - \frac{1}{\pi} \sin^{-1}(r)$ .

The joint density function of  $(X, Y)$  is

$$f_{XY}(x, y) = \frac{1}{2\pi\sigma_1\sigma_2\sqrt{1-r^2}} \exp\left[-\frac{1}{2(1-r^2)} \left\{ \frac{x^2}{\sigma_1^2} - \frac{2xy}{\sigma_1\sigma_2} + \frac{y^2}{\sigma_2^2} \right\}\right]$$

If we put  $U = \frac{X}{Y}$  and  $V = Y$ , then by the usual procedure, we can find that the joint pdf of  $(U, V)$  is

$$\begin{aligned} f_{UV}(u, v) &= \frac{1}{2\pi\sigma_1\sigma_2\sqrt{1-r^2}} \exp \\ &\quad \left[ -\frac{1}{2(1-r^2)} \left\{ \frac{u^2}{\sigma_1^2} - \frac{2ru}{\sigma_1\sigma_2} + \frac{1}{\sigma_2^2} \right\} v^2 \right] |v| \quad (-\infty < u, v < \infty) \end{aligned}$$

The pdf of  $U$  is simply the marginal pdf, given by

$$\begin{aligned} f_U(u) &= \frac{1}{\pi \sigma_1 \sigma_2 \sqrt{1-r^2}} \int_0^\infty \exp \left[ -\frac{1}{2(1-r^2)} \left\{ \frac{u^2}{\sigma_1^2} - \frac{2ru}{\sigma_1 \sigma_2} + \frac{1}{\sigma_2^2} \right\} v^2 \right] v dv \\ &= \frac{\sqrt{1-r^2}}{\pi \sigma_1 \sigma_2} \frac{1}{\left( \frac{u}{\sigma_1} - \frac{r}{\sigma_2} \right)^2 + \frac{1-r^2}{\sigma_2^2}}, -\infty < u < \infty \end{aligned}$$

which is a Cauchy's distribution.

Now,  $P\{X \text{ and } Y \text{ are of the same sign}\}$

$$\begin{aligned} &= P\left(\frac{X}{Y} > 0\right) = P(U > 0) \\ &= \frac{\sqrt{1-r^2}}{\pi \sigma_1 \sigma_2} \int_0^\infty \frac{du}{\left( \frac{u}{\sigma_1} - \frac{r}{\sigma_2} \right)^2 + \frac{1-r^2}{\sigma_2^2}} \\ &= \frac{\sqrt{1-r^2}}{\pi \sigma_2} \int_{-\frac{r}{\sigma_2}}^\infty \frac{dz}{z^2 + \left( \sqrt{\frac{1-r^2}{\sigma_2^2}} \right)^2}, \text{ putting } z = \frac{u}{\sigma_1} - \frac{r}{\sigma_2} \\ &= \frac{1}{\pi} \left( \tan^{-1} \frac{\sigma_2 z}{\sqrt{1-r^2}} \right) \Big|_{-\frac{r}{\sigma_2}}^\infty \\ &= \frac{1}{\pi} \left\{ \frac{\pi}{2} + \tan^{-1} \frac{r}{\sqrt{1-r^2}} \right\} \\ &= \frac{1}{2} + \frac{1}{\pi} \sin^{-1}(r) \end{aligned}$$

Similarly,  $P\{X \text{ and } Y \text{ are of opposite signs}\}$

$$\begin{aligned} &= P\left(\frac{X}{Y} < 0\right) = P(U < 0) \\ &= \frac{1}{2} - \frac{1}{\pi} \sin^{-1}(r) \end{aligned}$$

- (ii) If  $X$  and  $Y$  are two normal random variables with zero means, variances  $\sigma_1^2$  and  $\sigma_2^2$  and correlation coefficient  $r$ , then  $E\{|XY|\} = \frac{2}{\pi} \sigma_1 \sigma_2 (\cos \alpha + \alpha \sin \alpha)$ , where  $\sin \alpha = r$ .

We make use of a theorem, called *Price's theorem*, which we state next without proof:

If  $X$  and  $Y$  are two normal RVs as described above, then

$$\frac{\partial^n}{\partial C^n} E[g(X, Y)] = E \left[ \frac{\partial^{2n} g(X, Y)}{\partial X^n \partial Y^n} \right], \text{ where } C = C(X, Y)$$

is the covariance of  $X$  and  $Y$ .

If we take  $g(X, Y) = |XY|$  and  $n = 1$ , then Price's theorem gives

$$\frac{\partial}{\partial C} E\{|XY|\} = E \left\{ \frac{\partial^2 |XY|}{\partial X \partial Y} \right\} = E \left\{ \frac{d|X|}{dX} \frac{d|Y|}{dY} \right\} \quad (1)$$

$$\text{Now, } \frac{d|X|}{dX} = \begin{cases} 1 & \text{if } X > 0 \\ -1 & \text{if } X < 0 \end{cases}$$

$$\text{and } \frac{d|Y|}{dY} = \begin{cases} 1 & \text{if } Y > 0 \\ -1 & \text{if } Y < 0 \end{cases}$$

$$\therefore \frac{d}{dX} |X| \frac{d}{dY} |Y| = 1, \text{ if } \frac{X}{Y} > 0 \\ = -1, \text{ if } \frac{X}{Y} < 0$$

$$\begin{aligned} \therefore E \left\{ \frac{d}{dX} |X| \frac{d}{dY} |Y| \right\} &= 1 \times P \left\{ \frac{X}{Y} > 0 \right\} - 1 \times P \left\{ \frac{X}{Y} < 0 \right\} \\ &= \left( \frac{1}{2} + \frac{1}{\pi} \sin^{-1} r \right) - \left( \frac{1}{2} - \frac{1}{\pi} \sin^{-1} r \right) \\ &\quad [\text{by the previous result (i)}] \\ &= \frac{2}{\pi} \sin^{-1} (r) \\ &= \frac{2}{\pi} \sin^{-1} \left( \frac{C}{\sigma_1 \sigma_2} \right) \end{aligned} \quad (2)$$

Using (2) in (1), we get,

$$\frac{\partial}{\partial C} E\{|XY|\} = \frac{2}{\pi} \sin^{-1} \left( \frac{C}{\sigma_1 \sigma_2} \right)$$

Integrating both sides with respect to  $C$ , between 0 and  $C$ ,

$$E\{|XY|\} = E\{|XY|\}_{C=0} + \frac{2}{\pi} \int_0^C \sin^{-1} \left( \frac{C}{\sigma_1 \sigma_2} \right) dC \quad (3)$$

When  $C = 0$ ,  $X$  and  $Y$  are uncorrelated.

$$\therefore E\{|XY|\} = E\{|X|\} E\{|Y|\}$$

$$= \int_{-\infty}^{\infty} |x| \frac{1}{\sigma_1 \sqrt{2\pi}} e^{-x^2/2\sigma_1^2} dx \int_{-\infty}^{\infty} |y| \frac{1}{\sigma_2 \sqrt{2\pi}} e^{-y^2/2\sigma_2^2} dy$$

$$= \sqrt{\frac{2}{\pi}} \sigma_1 \sqrt{\frac{2}{\pi}} \sigma_2$$

i.e.,  $E\{|XY|\}_{C=0} = \frac{2}{\pi} \sigma_1 \sigma_2$  (4)

Using (4) in (3), we get,

$$\begin{aligned} E\{|XY|\} &= \frac{2}{\pi} \sigma_1 \sigma_2 + \frac{2}{\pi} \left[ C \sin^{-1} \frac{C}{\sigma_1 \sigma_2} - \int_0^C \frac{1}{\sqrt{1 - \frac{C^2}{\sigma_1^2 \sigma_2^2}}} \frac{1}{\sigma_1 \sigma_2} dC \right] \\ &= \frac{2}{\pi} \sigma_1 \sigma_2 + \frac{2C}{\pi} \sin^{-1}(r) + \frac{2}{\pi} \left( \sqrt{\sigma_1^2 \sigma_2^2 - C^2} \right)_0^C \\ &= \frac{2}{\pi} \sigma_1 \sigma_2 \left( r \sin^{-1} r + \sqrt{1 - r^2} \right) \\ &= \frac{2}{\pi} \sigma_1 \sigma_2 (\cos \alpha + \alpha \sin \alpha), \text{ putting } r = \sin \alpha \end{aligned}$$

## 2. Full-Wave Linear Detector Process

If  $\{X(t)\}$  is a zero mean stationary Gaussian process and if  $Y(t) = |X(t)|$ , then  $\{Y(t)\}$  is called a full-wave linear detector process.

$$\begin{aligned} E\{Y(t)\} &= E\{|X(t)|\} \\ &= \int_{-\infty}^{\infty} |X| \frac{1}{\sigma \sqrt{2\pi}} e^{-x^2/2\sigma^2} dx \\ &= \frac{1}{\sigma} \sqrt{\frac{2}{\pi}} \int_0^{\infty} x e^{-x^2/2\sigma^2} dx \\ &= \sqrt{\frac{2}{\pi}} \sigma \int_0^{\infty} e^{-t} dt, \text{ putting } \frac{x^2}{2\sigma^2} = t \\ &= \sqrt{\frac{2}{\pi}} \sigma = \sqrt{\frac{2 R_{xx}(0)}{\pi}} \end{aligned}$$

$$\begin{aligned} R_{yy}(t_1, t_2) &= E\{|X(t_1) X(t_2)|\} \\ &= \sqrt{\frac{2}{\pi}} \sigma^2 (\cos \alpha + \alpha \sin \alpha) \quad (\text{by the previous result}) \end{aligned}$$

where  $\sin \alpha = r = \frac{C\{X(t_1), X(t_2)\}}{\sigma^2}$

$$= \frac{E\{X(t_1), X(t_2)\}}{\sigma^2}$$

$$= \frac{R_{xx}(t_1 - t_2)}{\sigma^2} \quad [\text{since } \{X(t)\} \text{ is stationary}]$$

Therefore,  $\{Y(t)\}$  is wide-sense stationary, with  $R_{yy}(\tau) = \frac{2}{\pi} R_{xx}(0) (\cos \alpha + \alpha \sin \alpha)$ , where

$$\sin \alpha = \frac{R_{xx}(\tau)}{R_{xx}(0)}$$

Now,  $E\{Y^2(t)\} = R_{yy}(0) = \frac{2}{\pi} R_{xx}(0) \left\{ 0 + \frac{\pi}{2} 1 \right\},$

since  $\sin \alpha = \frac{R_{xx}(0)}{R_{xx}(0)} = 1$  and  $\alpha = \frac{\pi}{2}$

$\therefore E\{Y^2(t)\} = R_{xx}(0)$  and  $\text{Var}\{Y(t)\} = \left(1 - \frac{2}{\pi}\right) R_{xx}(0)$

### 3. Half-wave Linear Detector Process

If  $\{X(t)\}$  is a zero mean stationary Gaussian process and if

$$Z(t) = \begin{cases} X(t), & \text{for } X(t) \geq 0 \\ 0, & \text{for } X(t) < 0 \end{cases}$$

then  $\{Z(t)\}$  is called a half-wave linear detector process.

$$Z(t) \text{ can be rewritten as } Z(t) = \frac{1}{2} \{X(t) + |X(t)|\}$$

$$\begin{aligned} \therefore E\{Z(t)\} &= \frac{1}{2} [E\{X(t)\} + [E|X(t)|]] \\ &= \frac{1}{2} \left[ 0 + \sqrt{\frac{2}{\pi}} R_{xx}(0) \right] \quad (\text{refer to the previous process}) \\ &= \sqrt{\frac{R_{xx}(0)}{2\pi}} \end{aligned}$$

$$E\{Z(t) Z(t - \tau)\} = E[E\{Z(t) Z(t - \tau)/X(t) X(t - \tau)\}] \quad (1)$$

$$\text{Now, } Z(t) Z(t - \tau)/X(t) X(t - \tau) = \frac{1}{2} \{X(t) X(t - \tau) + |X(t) X(t - \tau)|\} \text{ (or)} = 0$$

The first value is assumed, when  $X(t) X(t - \tau) > 0$ , i.e., when  $X(t)$  and  $X(t - \tau)$  are both positive or both negative.

$$\therefore P\{\text{The first value of assumed}\} = \frac{1}{4} + \frac{1}{4} = \frac{1}{2}$$

$$\text{Similarly, } P\{\text{the second value is assumed}\} = \frac{1}{2}$$

$$\therefore E\{Z(t) Z(t - \tau)/X(t) X(t - \tau)\} = \frac{1}{4} \{X(t) X(t - \tau) + |X(t) X(t - \tau)|\} \quad (2)$$

Using (2) in (1), we get,

$$\begin{aligned} E\{Z(t) Z(t - \tau)\} &= \frac{1}{4} [E\{X(t) X(t - \tau) + E\{|X(t) X(t - \tau)|\}\}] \\ &= \frac{1}{4} [R_{xx}(\tau) + R_{yy}(\tau)] \quad [\text{where } \{Y(t)\} \text{ is the full-wave linear detector process}] \end{aligned}$$

i.e.,  $R_{zz}(\tau) = \frac{1}{4} [R_{xx}(\tau) + \frac{2}{\pi} R_{xx}(0) (\cos \alpha + \alpha \sin \alpha)],$

where  $\sin \alpha = \frac{R_{xx}(\tau)}{R_{xx}(0)}$

Therefore, the process  $\{Z(t)\}$  is wide-sense stationary.

Now  $E\{Z^2(t)\} = R_{zz}(0) \frac{1}{2} R_{xx}(0)$

$$\begin{aligned} \therefore \text{Var}\{Z(t)\} &= \frac{1}{2} R_{xx}(0) - \frac{1}{2\pi} R_{xx}(0) \\ &= \frac{1}{2} \left(1 - \frac{1}{\pi}\right) R_{xx}(0) \end{aligned}$$

#### 4. Hard Limiter Process

If  $\{X(t)\}$  is a zero mean stationary Gaussian process and if

$$Y(t) = \begin{cases} +1 & \text{for } X(t) \geq 0 \\ -1 & \text{for } X(t) < 0 \end{cases}$$

the  $\{Y(t)\}$  is called a hard limiter process or ideal limiter process.

$$\begin{aligned} E\{Y(t)\} &= 1 \times P\{X(t) \geq 0\} - 1 \times P\{X(t) < 0\} \\ &= 0 \end{aligned}$$

Now  $Y(t) Y(t - \tau) = 1, \text{ if } X(t) X(t - \tau) \geq 0$   
 $= -1, \text{ if } X(t) X(t - \tau) < 0$

i.e.,  $P\{Y(t) Y(t - \tau) = 1\} = P\{X(t) X(t - \tau) \geq 0\} \quad (i)$

$$\begin{aligned} &= \frac{1}{2} + \frac{1}{\pi} \sin^{-1} (r_{xx}) \quad [\text{by result (i) above}] \\ &= \frac{1}{2} + \frac{1}{\pi} \sin^{-1} \left\{ \frac{R_{xx}(\tau)}{R_{xx}(0)} \right\} \end{aligned}$$

and  $P\{Y(t) Y(t - \tau) = -1\} = P\{X(t) X(t - \tau) < 0\}$

$$= \frac{1}{2} - \frac{1}{\pi} \sin^{-1} \left\{ \frac{R_{xx}(\tau)}{R_{xx}(0)} \right\} \quad [\text{by result (1)}]$$

$\therefore E\{Y(t) Y(t - \tau)\} = \frac{2}{\pi} \sin^{-1} \left\{ \frac{R_{xx}(\tau)}{R_{xx}(0)} \right\}$

i.e.,  $R_{yy}(\tau) = \frac{2}{\pi} \sin^{-1} \left\{ \frac{R_{xx}(\tau)}{R_{xx}(0)} \right\}$  (1)

[(1) is called the *arcsine law*]

Therefore  $\{Y(t)\}$  is wide-sense stationary.

Also  $E\{Y^2(t)\} = 1$  and  $\text{Var}\{Y(t)\} = 1$

## Band Pass Process (Signal)

If the power spectrum of a random process  $\{X(t)\}$  is zero outside a certain band (an interval in the  $\omega$ -axis),

i.e.,  $S_{xx}(\omega) \neq 0$ , in  $|\omega - \omega_0| \leq \frac{\omega_B}{2}$  and in  $|\omega + \omega_0| \leq \frac{\omega_B}{2}$   
 $= 0$ , in  $|\omega - \omega_0| > \frac{\omega_B}{2}$  and in  $|\omega + \omega_0| > \frac{\omega_B}{2}$

the  $\{X(t)\}$  is called a *band pass process*.

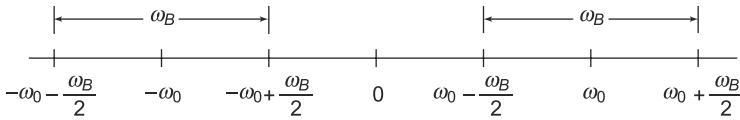


Fig 7.1

If  $S_{xx}(\omega) \neq 0$ , in  $|\omega| \leq \omega_B$   
 $= 0$ , in  $|\omega| > \omega_B$

the  $\{X(t)\}$  is called a *low pass process* or *ideal low pass process*.

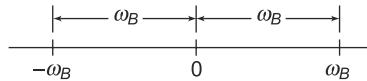


Fig. 7.2

If the bandwidth  $\omega_B$  of a band pass process is small compared with the centre frequency  $\omega_0$ , the process is called *narrow-band process* or *quasimonochromatic*.

If the power spectrum  $S_{xx}(\omega)$  of a bandpass process  $\{X(t)\}$  is an impulse function, then the process is called *monochromatic*.

## Narrow-Band Gaussian Process

In communication system, information bearing signals are often narrow-band Gaussian processes. When such signals are viewed on an oscilloscope, they appear like a sine wave with slowly varying amplitude and phase. Hence, a narrow-band Gaussian process  $\{X(t)\}$  is often represented as

$$X(t) = R_X(t) \cos [\omega_0 t + \theta_X(t)] \quad (1)$$

$R_X(t)$  and  $\theta_X(t)$ , which are low pass processes, are called the *envelope* and *phase* of the process  $\{X(t)\}$  respectively. (1) can be rewritten as

$$X(t) = [R_X(t) \cos \theta_X(t)] \cos \omega_0 t \mp [R_X(t) \sin \theta_X(t)] \sin \omega_0 t \quad (2)$$

$R_X(t) \cos \theta_X(t)$  is called the *inphase component* of the process  $\{X(t)\}$  and denoted as  $X_c(t)$  or  $I(t)$ .  $R_X(t) \sin \theta_X(t)$  is called the *quadrature component* of  $\{X(t)\}$  and denoted as  $X_s(t)$  or  $Q(t)$ .

Both  $X_c(t)$  and  $X_s(t)$  are low pass processes.

## Property

The envelope of a narrow-band Gaussian process follows a Rayleigh distribution and the phase follows a uniform distribution in  $(0, 2\pi)$ .

We note that

$$\sqrt{X_c^2(t) + X_s^2(t)} = R_X(t) \text{ and } \tan^{-1} \left\{ \frac{X_c(t)}{X_s(t)} \right\} = \theta_x(t)$$

Refer to Worked Example 25 in Chapter 3, in which we have proved the following result:

If  $X$  and  $Y$  are two independent  $N(0, \sigma)$  then  $R = \sqrt{X^2 + Y^2}$  follows a Rayleigh distribution and  $\phi = \tan^{-1} \frac{Y}{X}$  follows a uniform distribution in  $(0, 2\pi)$ .

According to this problem, the required property follows.

## Quadrature Representation of a WSS Process

In order to represent a process in the quadrature form, it need not be a narrowband Gaussian process. Any arbitrary zero mean WSS process  $\{X(t)\}$  can be represented in the quadrature form as proved below.

Let  $\{A(t)\}$  and  $\{B(t)\}$  be any two zero mean, jointly WSS processes,  $\alpha$  be a constant and  $\{Y(t)\}$  be the ‘dual’ process of the given process  $\{X(t)\}$ , defined below:

$Y(t)$  is so chosen that

$$X(t) + iY(t) = \{A(t) + iB(t)\} e^{+i\alpha t} \quad (1)$$

$$\text{or, } A(t) + iB(t) = \{X(t) + iY(t)\} e^{-i\alpha t} \quad (2)$$

$$\text{Then } A(t) = X(t) \cos \alpha t + Y(t) \sin \alpha t \quad (3)$$

$$\text{and } B(t) = Y(t) \cos \alpha t - X(t) \sin \alpha t \quad (4)$$

It is easily verified, using (3) and (4), that  $E\{A(t) A(t - \tau)\}$ ,  $E\{B(t) B(t - \tau)\}$  and  $E\{A(t) B(t - \tau)\}$  are independent of  $t$ , i.e., functions of  $\tau$ .

$$\text{only if } R_{xx}(\tau) = R_{yy}(\tau) \quad (5)$$

$$\text{and } R_{xy}(\tau) = -R_{yx}(\tau) \quad (6)$$

Since  $\{A(t)\}$  and  $\{B(t)\}$  are jointly WSS, conditions (5) and (6) must be satisfied by  $\{X(t)\}$  and its dual  $\{Y(t)\}$ .

When (5) and (6) are true, from (1) we get,

$$X(t) = A(t) \cos \alpha t - B(t) \sin \alpha t \quad (7)$$

$$Y(t) = A(t) \sin \alpha t + B(t) \cos \alpha t \quad (8)$$

Thus, if we can get a suitable  $Y(t)$  and  $\alpha$ , satisfying (5) and (6),  $X(t)$  can be expressed in form (7), which is the quadrature form.

Consider the linear time-invariant system

$$Z(t) = \int_{-\infty}^{\infty} X(\alpha) h(t - \alpha) d\alpha \quad (9)$$

$$\text{where } h(t) = \frac{1}{\pi t} \quad (10)$$

$$\text{We note that } F\{h(t)\} = \frac{1}{\pi} \int_{-\infty}^{\infty} \frac{e^{-i\omega t}}{t} dt$$

$$= \frac{1}{\pi} \left[ \int_{-\infty}^{\infty} \frac{\cos \omega t}{t} dt - i \int_{-\infty}^{\infty} \frac{\sin \omega t}{t} dt \right]$$

$$= \frac{1}{\pi} (0 - i\pi)$$

(using well-known results in contour integration if  $\omega > 0$ )

$$= -i, \text{ if } \omega > 0$$

$$\text{and } F[h(t)] = i, \text{ if } \omega < 0$$

$$\text{i.e., } H(\omega) = \begin{cases} i, & \text{if } \omega < 0 \\ -i, & \text{if } \omega > 0 \end{cases}$$

Using (10) in (9), we get

$$Z(t) = \frac{1}{\pi} \int_{-\infty}^{\infty} \frac{X(\alpha)}{t - \alpha} d\alpha$$

which is called the *Hilbert transform* of  $X(t)$  and denoted as  $\bar{X}(t)$ .

$$\therefore \bar{X}(t) = X(t) * \frac{1}{\pi t}$$

Therefore, Property 3 of the linear time-invariant system,

$$\begin{aligned} S_{X\bar{X}}(\omega) &= H^*(\omega) S_{xx}(\omega) \\ &= \begin{cases} -iS_{xx}(\omega), & \text{if } \omega < 0 \\ +iS_{xx}(\omega), & \text{if } \omega > 0 \end{cases} \end{aligned} \quad (11)$$

$$\text{and } S_{\bar{X}\bar{X}}(\omega) = H(\omega) S_{xx}(\omega)$$

$$= \begin{cases} iS_{xx}(\omega), & \text{if } \omega < 0 \\ -iS_{xx}(\omega), & \text{if } \omega > 0 \end{cases} \quad (12)$$

$$\text{From (11) and (12), } S_{X\bar{X}}(\omega) = -S_{\bar{X}\bar{X}}(\omega) \quad (13)$$

$$\text{Also, } S_{\bar{X}\bar{X}}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

$$\text{i.e., } S_{\bar{X}\bar{X}}(\omega) = S_{xx}(\omega) \quad (14)$$

From (13) and (14), it follows that

$$R_{X\bar{X}}(\tau) = -R_{\bar{X}X}(\tau) \quad (15)$$

and  $R_{\bar{X}\bar{X}}(\tau) = R_{xx}(\tau) \quad (16)$

Thus, conditions (5) and (6) are satisfied by  $Y(t) = \bar{X}(t)$ .

Therefore, we can use  $\bar{X}(t)$  in the place of  $Y(t)$  for the quadrature representation of  $X(t)$ .

Now, let  $X(t) + i\bar{X}(t) = \{I(t) + iQ(t)\} e^{i\omega_0 t}$  (17)

Then, replacing  $A(t)$ ,  $B(t)$ ,  $\alpha$ , by  $I(t)$ ,  $Q(t)$  and  $\omega_0$  respectively in (7), we get  $X(t) = I(t) \cos \omega_0 t - Q(t) \sin \omega_0 t$ , which is the required quadrature representation of  $\{X(t)\}$ .

**Note**

1.  $I(t)$  and  $Q(t)$  are called the inphase and quadrature components of  $\{X(t)\}$ .
2. The quadrature representation of  $\{X(t)\}$  is not unique.
3. The quadrature representation is **useful**, only when  $\{X(t)\}$  is a zero mean WSS bandpass process.

## Noise in Communication Systems

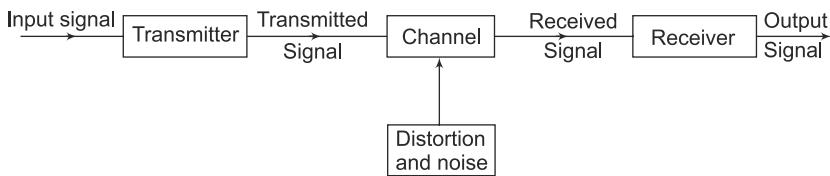


Fig. 7.3

In communication systems, the message to be transmitted to a far-off location is first converted into an electrical waveform called input signal, before being sent into the transmitter. The transmitter processes and modifies the input signal for efficient transmission. The transmitter output is then sent through the channel which is just a medium such as wire, coaxial cable or optical fibre. The channel output or the received signal is then reprocessed by the receiver which sends out the output signal. The output signal is converted to its original form, namely the message.

When the message is communicated in this manner, the signal is not only distorted by the channel but also contaminated along the path by undesirable signals that are generally referred to by the term *noise*. The noise can come from many external and internal sources and take many forms.

External noise includes interfering signals from nearby sources, man-made noise generated by faulty contact switches for the electrical equipment, by ignition radiation, fluorescent lights, natural noise from lighting and extraterrestrial radiation, etc. Internal noise results from thermal motion of electrons in

conductors, random emission and diffusion or recombination of charged carriers in electronic devices. By careful engineering techniques, the effects of many unwanted signals can be eliminated or minimised. But there always remain certain inescapable random signals that set a limit to system performance, i.e., on the efficiency of communication.

One of the main reasons for introducing probability theory in the study of ‘Signal Analysis’ is the random nature of noise. Because of this randomness, it is usual to describe noise as a random process and hence in terms of a probabilistic model. Such a model describes the noise amplitude or any other parameter by means of a probability density function  $f(x)$  [ $x$  represents voltage]. For many important types of noise, the density function can be determined theoretically and for others it has been estimated empirically.

Certain properties of noise, such as mean value, mean square value and the root-mean square value can be found by using the probability density function.

However the probability density function does not describe a noise waveform sufficiently so as to determine its effect on the performance of a communication system. To achieve this, it is necessary to know how the noise change with time. This information is provided by a mean-square voltage spectrum, called the power spectrum or spectral density, that represents the distribution of signal power as a function of frequency.

## Thermal Noise

Thermal noise is the noise because of the random motion of free electrons in conducting media such as a resistor. Thermal noise generated in resistors and semiconductors is assumed to be a zero mean, stationary Gaussian random process  $\{N(t)\}$  with a power spectral density that is flat over a very wide range of frequencies, i.e., the graph of  $S_{NN}(\omega)$  is a straight line parallel to the  $\omega$ -axis. Since  $S_{NN}(\omega)$  contains all frequencies in equal amount, the noise is also called *white Gaussian noise* or simply *white noise* in analogy to white light which consists of all colours.

It is customary to denote the constant spectral density of white noise by  $\frac{N_0}{2}$  (or  $\frac{\eta}{2}$ ).

$$\text{i.e., } S_{NN}(\omega) = \frac{N_0}{2}$$

The autocorrelation function of the white noise is given by

$$R_{NN}(\tau) = \frac{N_0}{2} \delta(\tau) \text{ since}$$

$$\int_{-\infty}^{\infty} \frac{N_0}{2} \delta(\tau) e^{-i\omega\tau} d\tau = \frac{N_0}{2}$$

The average power of the white noise  $\{N(t)\}$  is given by

$$R_{NN}(0) = \int_{-\infty}^{\infty} S_{NN}(\omega) d\omega = \int_{-\infty}^{\infty} \frac{N_0}{2} d\omega \rightarrow \infty$$

Therefore, the spectral density of  $\{N(t)\}$  is not physically realisable. However, since the bandwidths of real processes are always finite and since

$$\int_{-\omega_B}^{\omega_B} S_{NN}(\omega) d\omega = N_0 \omega_B < \infty$$

for any finite bandwidth, the spectral density  $S_{NN}(\omega)$  can be used over finite bandwidths.

**Definition:** Noise having a nonzero and constant spectral density over a finite frequency band and zero elsewhere is called *band-limited white noise*.

i.e., if  $\{N(t)\}$  is a band-limited white noise then

$$S_{NN}(\omega) = \begin{cases} \frac{N_0}{2}, & |\omega| \leq \omega_0 \\ 0, & \text{elsewhere} \end{cases}$$

We give below a few properties of the band-limited white noise which can be easily verified by the reader.

$$1. \quad E\{N^2(t)\} = \frac{N_0 \omega_B}{2\pi}$$

$$2. \quad R_{NN}(\tau) = \frac{N_0 \omega_B}{2\pi} \left( \frac{\sin \omega_B \tau}{\omega_B \tau} \right)$$

$$3. \quad N(t) \text{ and } N\left(t + \frac{k\pi}{\omega_B}\right) \text{ are independent, where } k \text{ is a nonzero integer.}$$

## Filters

Filtering is commonly used in electrical systems to reject undesirable signals and noise and to select the desired signal. A simple example of filtering occurs when we ‘tune’ in a particular radio to ‘select’ one of many signals.

Filtering actually means selecting carefully the transfer function  $H(\omega)$  in a stable, linear, time-invariant system, so as to modify the spectral components of the input signal. The system function  $H(\omega)$  or the linear system itself is referred to as filter, when it does the filtering.

The commonly used filters are narrow-band filters, i.e., band pass and low pass filters.

If the system function  $H(\omega)$  is defined as

$$\begin{aligned} H(\omega) &\neq 0, \text{ for } \omega_0 - \varepsilon/2 < \omega < \omega_0 + \varepsilon/2 \\ &= 0, \text{ otherwise} \end{aligned}$$

then the filter is called a *band pass filter*.

$$\text{If } H(\omega) \neq 0, \text{ for } -\varepsilon/2 < \omega < \varepsilon/2 \\ = 0, \text{ otherwise}$$

then the filter is called a *low pass filter*.

The equation  $S_{yy}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$  shows that the spectral properties of a signal can be modified by passing it through a linear time-invariant system with the appropriate transfer function. By carefully choosing  $H(\omega)$ , we can remove or filter out certain spectral components in the input. For example, let the input  $X(t) = S(t) + N(t)$ , where  $S(t)$  is the signal of interest and  $N(t)$  is an unwanted noise process. If the spectral densities of  $\{S(t)\}$  and  $\{N(t)\}$  are non-overlapping in the frequency domain, the noise  $N(t)$  can be removed by passing  $X(t)$  through a filter  $H(\omega)$  that has a response of 1 for the range of frequencies occupied by the signal and a response of 0 for the range of frequencies occupied by the noise. But in most practical situations there is spectral overlap and the design of optimum filters to separate signal and noise is somewhat difficult. The discussion of this problem and the various optimum filters in common use such as matched filter and Wiener filter may be found in textbooks on Random Signal Analysis. It is beyond the scope of this book.

### Worked Example 7(A)

#### Example 1

If  $\{X(t)\}$  is a Gaussian process with  $\mu(t) = 10$  and  $C(t_1, t_2) = 16e^{-|t_1-t_2|}$  find the probability that (i)  $X(10) \leq 8$  and (ii)  $|X(10) - X(6)| \leq 4$ .

Since  $\{X(t)\}$  is a Gaussian process, any member of the process is a normal RV.

Therefore,  $X(10)$  is a normal RV with mean  $\mu(10) = 10$  and variance  $C(10, 10) = 16$ .

$$\begin{aligned} P\{X(10) \leq 8\} &= P\left\{\frac{X(10) - 10}{4} \leq -0.5\right\} \\ &= P\{Z \leq -0.5\} \quad (\text{where } Z \text{ is the standard normal RV}) \\ &= 0.5 - P\{0 \leq Z \leq 0.5\} \\ &= 0.5 - 0.1915 \quad (\text{from normal tables}) \\ &= 0.3085 \end{aligned}$$

$X(10) - X(6)$  is also a normal RV with mean  $\mu(10) - \mu(6) = 10 - 10 = 0$ .

$$\begin{aligned} \text{Var}\{X(10) - X(6)\} &= \text{Var}\{X(10)\} + \text{Var}\{X(6)\} - 2 \text{Covar}\{X(10), X(6)\} \\ &= C(10, 10) + C(6, 6) - 2C(10, 6) \\ &= 16 + 16 - 2 \times 16e^{-4} \\ &= 31.4139 \end{aligned}$$

$$\text{Now, } P\{|X(10) - X(6)| \leq 4\} = P\left\{\frac{|X(10) - X(6)|}{5.6048} \leq \frac{4}{5.6048}\right\}$$

$$\begin{aligned}
 &= P\{|Z| \leq 0.7137\} \\
 &= 2 \times 0.2611 \\
 &= 0.5222
 \end{aligned}$$

---

### Example 2

---

The process  $\{X(t)\}$  is normal with  $\mu_t = 0$  and  $R_x(\tau) = 4e^{-3|\tau|}$ . Find a memoryless system  $g(x)$  such that the first order density  $f_Y(y)$  of the resulting output  $Y(t) = g[X(t)]$  is uniform in the interval (6, 9).

Since  $\{X(t)\}$  is a normal process, a sample function  $X(t)$  follows a normal distribution with mean 0 and variance given by  $R_x(0) = 4$ .

$$\therefore f_X(x) = \frac{1}{2\sqrt{2\pi}} e^{-x^2/8}, -\infty < x < \infty$$

Now,  $Y(t)$  is to be uniform in (6, 9)

$$\therefore f_Y(y) = \frac{1}{3}, 6 < y < 9$$

Therefore, the distribution function of  $Y$  is given by

$$\begin{aligned}
 F_Y(y) &= \int_0^y f_Y(y) dy \\
 &= \frac{1}{3} (y - 6)
 \end{aligned} \tag{1}$$

Now,

$$\begin{aligned}
 F_Y(y) &= P\{Y(t) \leq y\} \\
 &= P\{g[X(t)] \leq y\} \\
 &= P\{X(t) \leq g^{-1}(y)\} \\
 &= P\{X(t) \leq x\} \quad [\text{since } y = g(x)] \\
 &= F_X(x)
 \end{aligned}$$

$$\text{But, from (1), } F_Y\{g(x)\} = \frac{1}{3}\{g(x) - 6\}$$

$$\therefore \frac{1}{3}\{g(x) - 6\} = F_X(x)$$

$$\begin{aligned}
 \therefore g(x) &= 6 + 3F_X(x) \\
 &= 6 + 3 \int_{-\infty}^{\infty} \frac{1}{2\sqrt{2\pi}} e^{-x^2/8} dx
 \end{aligned}$$

---

### Example 3

---

Let  $Z$  and  $\theta$  be independent RVs such that  $Z$  has a density function

$$f(z) = \begin{cases} 0, & \text{in } z < 0 \\ ze^{-z^2/2}, & \text{in } z > 0 \end{cases}$$

and  $\theta$  is uniformly distributed in  $(0, 2\pi)$ . Show that  $\{X_t; -\infty < t < \infty\}$  is a Gaussian process; if  $X_t = Z \cos(2\pi t + \theta)$ .

Let us first find the density function of

$$W = \cos(2\pi t + \theta), \text{ where } f_\theta(\theta) = \frac{1}{2\pi}$$

Since  $w = \cos(2\pi t + \theta)$ ,  $\theta = \cos^{-1}(w) - 2\pi t$ .

There are only two values of  $\theta$  in  $(0, 2\pi)$  for a given value of  $w$ . Let them be  $\theta_1$  and  $\theta_2$ .

By the transformation rule (refer Chapter 3)

$$\begin{aligned} f_W(w) &= f_\theta(\theta_1) \left| \frac{d\theta_1}{dw} \right| + f_\theta(\theta_2) \left| \frac{d\theta_2}{dw} \right| \\ &= 2 \times \frac{1}{2\pi} \times \left| -\frac{1}{1-w^2} \right| \\ &= \frac{1}{\pi\sqrt{1-w^2}}, |w| < 1 \end{aligned}$$

Let us now find the first-order density of  $X = ZW$ , where  $Z$  has been taken as  $X$ .

Introduce the auxiliary variable  $Y = W$ , so that we may find the joint pdf of  $(X, Y)$

$$x = zw \text{ and } y = w$$

$$\text{i.e., } z = \frac{x}{y} \text{ and } w = y$$

$$\therefore f_{XY}(x, y) = |J| f_{ZW}(z, w)$$

$$\text{where } J = \begin{vmatrix} \frac{\partial z}{\partial x} & \frac{\partial z}{\partial y} \\ \frac{\partial w}{\partial x} & \frac{\partial w}{\partial y} \end{vmatrix} = \begin{vmatrix} \frac{1}{y} & -\frac{x}{y^2} \\ 0 & 1 \end{vmatrix} = \frac{1}{|y|}$$

$\therefore f_X(x)$  is the marginal density function of  $X$ .

$$\text{i.e., } f_X(x) = \int_{-1}^1 \frac{1}{|y|} f_Z(z) f_W(w) dy, \text{ where } z = \frac{x}{y} \text{ and } w = y$$

(by independence of  $Z$  and  $W$ )

$$= \int_{-1}^1 \frac{1}{|y|} \frac{x}{y} e^{-x^2/2y^2} \frac{1}{\pi\sqrt{1-y^2}} dy, \quad \frac{x}{y} > 0$$

$$= \begin{cases} \frac{1}{\pi} \int_{-1}^0 -\frac{x}{y^2} e^{-x^2/2y^2} \frac{1}{\sqrt{1-y^2}} dy, & \text{where } x < 0 \text{ and } y < 0 \quad (1) \end{cases}$$

$$= \begin{cases} \frac{1}{\pi} \int_0^1 \frac{x}{y^2} e^{-x^2/2y^2} \frac{1}{\sqrt{1-y^2}} dy, & \text{where } x > 0 \text{ and } y > 0 \quad (2) \end{cases}$$

Changing  $y$  to  $-y'$ , we note that the integral in (1) becomes

$$\frac{1}{\pi} \int_0^1 \frac{x}{y'^2} e^{-x^2/2y'^2} \frac{1}{\sqrt{1-y'^2}} dy', \text{ which is the same as integral (2).}$$

$$\therefore f_X(x) = \frac{1}{\pi} \int_0^1 \frac{x}{y^2} e^{-x^2/2y^2} \frac{1}{\sqrt{1-y^2}} dy, -\infty < x < \infty \quad (3)$$

Put  $\frac{x^2}{2y^2} = t$  in (3), treating  $x$  as a parameter,

$$\text{Then } f_X(x) = \frac{1}{\pi} \int_{x^2/2}^{\infty} \frac{1}{\sqrt{2t-x^2}} e^{-t} dt \quad (4)$$

Put  $t - \frac{x^2}{2} = u$  in (4), treating  $x$  as a parameter.

$$\begin{aligned} \text{Then } f_X(x) &= \frac{1}{\pi\sqrt{2}} \int_0^{\infty} e^{-x^2/2} u^{-1/2} e^{-u} du \\ &= \frac{1}{\pi\sqrt{2}} e^{-x^2/2} \left[ \left( \frac{1}{2} \right) \right] \\ &= \frac{1}{\sqrt{2\pi}} e^{-x^2/2}, \infty < x < \infty \end{aligned}$$

Thus, each member of the process  $\{X_t\}$  follows a normal distribution with mean zero and variance 1.

Therefore,  $a_1 X_{t_1} + a_2 X_{t_2} + \dots + a_n X_{t_n}$  also follows a normal distribution, for any set of  $a_1, a_2, \dots, a_n$ .

Therefore,  $\{X_{t_1}, X_{t_2}, \dots, X_{t_n}\}$  are jointly normal for any  $n$  [see the Problem 2 given in Exercise 8 (A)].

Therefore, the process  $\{X_t\}$  is Gaussian.

#### Example 4

It is given that  $R_x(\tau) = e^{-|\tau|}$  for a certain stationary Gaussian random process  $\{X(t)\}$ . Find the joint pdf of the RVs  $X(t), X(t+1), X(t+2)$ .

Let us denote the given RVs by  $X(t_1), X(t_2), X(t_3)$ .

The joint pdf of  $\{X(t_1), X(t_2), X(t_3)\}$  is given by

$$F(x_1, x_2, x_3, t_1, t_2, t_3) =$$

$$\frac{1}{(2\pi)^{3/2} |\Lambda|^{1/2}} \exp \left[ -\frac{1}{2|\Lambda|} \sum_{i=1}^3 \sum_{j=1}^3 |\Lambda|_{ij} (x_i - \mu_i)(x_j - \mu_j) \right]$$

where  $\mu_i = E\{X(t_i)\}$  and  $\Lambda$  is the third order square matrix  $(\lambda_{ij})$ , where  $\lambda_{ij} = C\{X(t_i), X(t_j)\}$  and  $|\Lambda|_{ij}$  = cofactor of  $\lambda_{ij}$  in  $|\Lambda|$ .

$$\begin{aligned}
 E\{X(t)\} &= \lim_{\tau \rightarrow \infty} R_x(\tau) = \lim_{\tau \rightarrow \infty} e^{-|\tau|} = 0 \\
 \therefore \lambda_{ij} &= C\{X(t_i) X(t_j)\} = R(t_i - t_j) \\
 \therefore \lambda_{11} &= R(t_1, t_1) = R(t, t) = R(0) = 1 \\
 \lambda_{12} &= R(t, t+1) = R(1) = e^{-1} \text{ etc.} \\
 \therefore \Lambda &= \begin{pmatrix} 1 & \frac{1}{e} & \frac{1}{e^2} \\ \frac{1}{e} & 1 & \frac{1}{e} \\ \frac{1}{e^2} & \frac{1}{e} & 1 \end{pmatrix} \text{ and } |\Lambda| = \left(1 - \frac{1}{e^2}\right)^2 \\
 |\Lambda|_{11} &= 1 - \frac{1}{e^2} \quad |\Lambda|_{12} = -\frac{1}{e} + \frac{1}{e^3} \quad |\Lambda|_{13} = 0 \text{ etc.}
 \end{aligned}$$

Therefore, the required joint pdf is given by

$$\begin{aligned}
 f(x_1, x_2, x_3) &= \frac{1}{(2\pi)^{3/2} \left(1 - \frac{1}{e^2}\right)} \exp \\
 &\left[ -\frac{1}{2 \left(1 - \frac{1}{e^2}\right)^2} \left\{ \left(1 - \frac{1}{e^2}\right) x_1^2 - \frac{2}{e} \left(1 - \frac{1}{e^2}\right) x_1 x_2 + \left(1 - \frac{1}{e^4}\right) x_2^2 \right. \right. \\
 &\quad \left. \left. - \frac{2}{e} \left(1 - \frac{1}{e^2}\right) x_2 x_3 + \left(1 - \frac{1}{e^2}\right) x_3^2 \right\} \right]
 \end{aligned}$$

$$\begin{aligned}
 \text{i.e., } f(x_1, x_2, x_3) &= \frac{1}{(2\pi)^{3/2} \left(1 - \frac{1}{e^2}\right)} \exp \\
 &\left[ -\frac{1}{2 \left(1 - \frac{1}{e^2}\right)} \left\{ x_1^2 - \frac{2}{e} x_1 x_2 + \left(1 + \frac{1}{e^2}\right) x_2^2 - \frac{2}{e} x_2 x_3 + x_3^2 \right\} \right]
 \end{aligned}$$

### Example 5

If  $\{Y(t)\}$  is the square-law detector process and if  $Z(t) = Y(t) - E\{Y(t)\}$ , show that the spectral density of  $\{Z(t)\}$  is given by  $S_{zz}(\omega) = \frac{1}{\pi} \int_{-\infty}^{\infty} S_{xx}(\alpha) S_{xx}(\omega - \alpha) d\alpha$ , where  $S_{xx}(\omega)$  is the input spectral density.

**Note**

$Z(t)$  is called the fluctuation of the square law detector.

$$\begin{aligned} E\{Z(t) Z(t-\tau)\} &= E[\{(Y(t) - E[Y(t)])\} \{Y(t-\tau) - E[Y(t-\tau)]\}] \\ &= E\{Y(t) Y(t-\tau)\} - E\{Y(t)\} E\{Y(t-\tau)\} \end{aligned}$$

$$\begin{aligned} \text{i.e., } R_{zz}(\tau) &= R_{yy}(\tau) - E\{Y(t)\} E\{Y(t-\tau)\} \\ &= R_{xx}^2(0) + 2R_{xx}^2(\tau) - R_{xx}^2(0) \text{ [see 'square law detector process']} \\ &= 2R_{xx}^2(\tau) \end{aligned}$$

Taking Fourier transforms,

$$S_{zz}(\omega) = \frac{1}{\pi} \int_{-\infty}^{\infty} S_{xx}(\alpha) S_{xx}(\omega - \alpha) d\alpha$$

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**Example 6**


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If  $\{X(t)\}$  is a band pass process, prove that its inphase and quadrature components are low pass processes (refer to the book-work on quadrature representation of a WSS process).

If we take the input of a system as  $X(t)$  and the impulse response function

$$h(t) = \delta(t) + i \frac{1}{\pi t}$$

Then the output of the system is given by

$$\begin{aligned} Z(t) &= \int_{-\infty}^{\infty} x(\alpha) \left[ \delta(t-\alpha) + i \frac{1}{\pi(t-\alpha)} \right] d\alpha \\ &= \int_{-\infty}^{\infty} X(t-\alpha) \delta(\alpha) d\alpha + \frac{i}{\pi} \int_{-\infty}^{\infty} \frac{X(\alpha)}{t-\alpha} d\alpha \\ &= X(t) + i\check{X}(t), \end{aligned}$$

[by the property of  $\delta(t)$  and by the definition of Hilbert transform]

$$\therefore S_{zz}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

$$\text{where } H(\omega) = F\{h(t)\}$$

$$\begin{aligned} &= F[\delta(t)] + iF\left[\frac{1}{\pi t}\right] \\ &= \begin{cases} 1 + i(i), & \text{if } \omega < 0 \\ 1 + i(-i), & \text{if } \omega > 0 \end{cases} \\ &= \begin{cases} 0, & \text{if } \omega < 0 \\ 2, & \text{if } \omega > 0 \end{cases} \\ &= 2U(\omega), \end{aligned}$$

where  $U(\omega)$  is the unit step function.

$$\therefore S_{zz}(\omega) = 4S_{xx}(\omega)U(\omega) \quad (1)$$

Also,  $Z(t) Z^*(t - \tau) = \{X(t) + i\bar{X}(t)\} \{X(t - \tau) - i\bar{X}(t - \tau)\}$   
 $= X(t)X(t - \tau) + \bar{X}(t)\bar{X}(t - \tau) + i\{\bar{X}(t)X(t - \tau) - X(t)\bar{X}(t - \tau)\}$

$$\therefore R_{zz}(\tau) = R_{xx}(\tau) + R_{\bar{x}\bar{x}}(\tau) + i\{R_{\bar{x}x}(\tau) - R_{x\bar{x}}(\tau)\}$$
  
 $= 2\{R_{xx}(\tau) + iR_{\bar{x}x}(\tau)\}, \quad (2)$

by steps (15) and (16) of the book-work

$$\therefore S_{zz}(\omega) = 2\{S_{xx}(\omega) + iS_{\bar{x}x}(\omega)\} \quad (3)$$

From step (17) of the book-work

$$I(t) + iQ(t) = \{X(t) + i\bar{X}(t)\} e^{-i\omega_0 t}$$
  
 $\therefore I(t) = X(t)\cos \omega_0 t + \bar{X}(t)\sin \omega_0 t \quad (4)$

$$Q(t) = \bar{X}(t)\cos \omega_0 t - X(t)\sin \omega_0 t \quad (5)$$

From (4) and (5),

$$R_{II}(\tau) = R_{xx}(\tau)\cos \omega_0 \tau + R_{\bar{x}x}(\tau)\sin \omega_0 \tau$$
  
 $R_{QI}(\tau) = -R_{xx}(\tau)\sin \omega_0 \tau + R_{\bar{x}x}(\tau)\cos \omega_0 \tau$   
 $\therefore R_{II}(\tau) + iR_{QI}(\tau) = [R_{xx}(\tau) + iR_{\bar{x}x}(\tau)]e^{-i\omega_0 \tau}$   
 $= \frac{1}{2}R_{zz}(\tau)e^{i\omega_0 \tau}$   
 $\therefore R_{II}(\tau) = \frac{1}{2}Rl\{R_{zz}(\tau)e^{-i\omega_0 \tau}\}$   
 $\therefore S_{II}(\omega) = \frac{1}{2}Rl\{S_{zz}(\omega + \omega_0)\}$   
 $= S_{xx}(\omega + \omega_0) \quad (6)$

Changing  $\tau$  into  $-\tau$  in (2),

$$R_{zz}(-\tau) = 2\{R_{xx}(-\tau) + iR_{\bar{x}x}(-\tau)\}$$
  
 $= 2\{R_{xx}(\tau) - iR_{\bar{x}x}(\tau)\} \quad (7)$

From (2) and (7),

$$R_{xx}(\tau) = \frac{1}{4}\{R_{zz}(\tau) + R_{zz}(-\tau)\}$$
  
 $\therefore S_{xx}(\omega) = \frac{1}{4}\{S_{zz}(\omega) + S_{zz}(-\omega)\} \quad (8)$

Using (8) in (6),

$$S_{II}(\omega) = \frac{1}{4}\{S_{zz}(\omega + \omega_0) + S_{zz}(-\omega + \omega_0)\}$$
  
 $= S_{xx}(\omega + \omega_0) U(\omega + \omega_0) + S_{xx}(-\omega + \omega_0)$   
 $U(-\omega + \omega_0) [\text{by (1)}] \quad (9)$

Since  $R_{QQ}(\tau) = R_{II}(\tau)$ ,  $S_{QQ}(\omega) = S_{II}(\omega)$  (10)

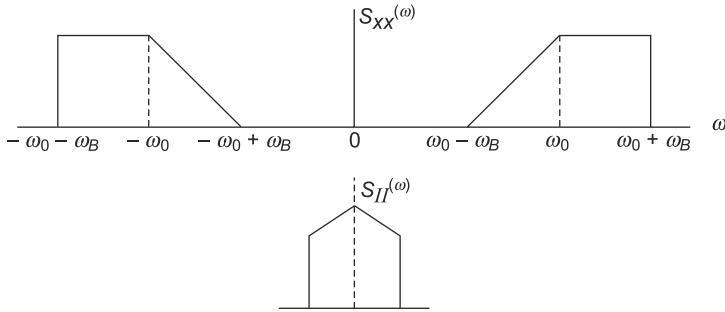


Fig. 7.4

Therefore, when  $\{X(t)\}$  is a band pass process, both  $\{I(t)\}$  and  $\{Q(t)\}$  are low pass processes as given in the above figures.

### Example 7

Obtain the autocorrelation for an ideal low pass stochastic process.

Let the spectral density function of the low pass process  $\{X(t)\}$  be  $S_{xx}(\omega)$ , in  $|\omega| < \omega_B$ .

Let the complex form of Fourier series of  $S_{xx}(\omega)$  in  $(-\omega_B, \omega_B)$  be

$$S_{xx}(\omega) = \sum_{n=-\infty}^{\infty} c_n e^{in\pi\omega/\omega_B} \quad (1)$$

where  $c_n$  is given by

$$c_n = \frac{1}{2\omega_B} \int_{-\omega_B}^{\omega_B} S_{xx}(\omega) e^{-in\pi\omega/\omega_B} d\omega \quad (2)$$

Taking the inverse Fourier transform of (1),

$$\begin{aligned} R_{xx}(\tau) &= \frac{1}{2\pi} \int_{-\infty}^{\infty} \sum c_n e^{in\pi\omega'/\omega_B} e^{i\tau\omega'} d\omega' \\ &= \sum \frac{1}{2\pi} \int_{-\omega_B}^{\omega_B} \frac{1}{2\omega_B} \int_{-\omega_B}^{\omega_B} S_{xx}(\omega) e^{-in\pi\omega/\omega_B} d\omega e^{i\left(\frac{n\pi}{\omega_B} + \tau\right)\omega'} d\omega' \end{aligned}$$

[since  $\{X(t)\}$  is low pass]

$$\begin{aligned} &= \sum \frac{1}{2\omega_B} \int_{-\omega_B}^{\omega_B} R_{xx}\left(-\frac{n\pi}{\omega_B}\right) e^{i\left(\frac{n\pi}{\omega_B}\right)\omega'} \omega' d\omega' \\ &= \sum R_{xx}\left(-\frac{n\pi}{\omega_B}\right) \frac{1}{\omega_B} \left\{ \frac{\sin\left(\frac{n\pi}{\omega_B} + \tau\right)\omega'}{\left(\frac{n\pi}{\omega_B} + \tau\right)} \right\}_0^{\omega_B} \end{aligned}$$

$$\begin{aligned}
&= \sum_{n=-\infty}^{\infty} R_{xx}\left(-\frac{n\pi}{\omega_B}\right) \frac{\sin\left(\frac{n\pi}{\omega_B} + \tau\right)\omega_B}{\left(\frac{n\pi}{\omega_B} + \tau\right)\omega_B} \\
&= \sum_{n=-\infty}^{\infty} R_{xx}\left(\frac{n\pi}{\omega_B}\right) \frac{\sin\left(\tau - \frac{n\pi}{\omega_B}\right)\omega_B}{\left(\tau - \frac{n\pi}{\omega_B}\right)\omega_B} \quad (\text{Changing } n \text{ to } -n)
\end{aligned}$$

Let us assume that the values of  $X(t)$  at  $t = nT$ ,  $n = \dots, -3, -2, -1, 0, 1, 2, 3, \dots$  are given, using which we can construct  $X(t)$ , where  $T = \frac{\pi}{\omega_B}$ .

$$\therefore R_{xx}(\tau) = \sum_{n=-\infty}^{\infty} R_{xx}(nT) \frac{\sin \omega_B(\tau - nT)}{\omega_B(\tau - nT)}$$

Thus, when  $\{X(t)\}$  is a low pass process, its autocorrelation is found out by summation.

### Example 8

If  $\{X(t)\}$  is a band limited process such that  $S_{xx}(\omega) = 0$ , when  $|\omega| > \sigma$ , prove that  $2[R_{xx}(0) - R_{xx}(\tau)] \leq \sigma^2 \tau^2 R_{xx}(0)$ .

$$\begin{aligned}
R_{xx}(\tau) &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{xx}(\omega) e^{i\tau\omega} d\omega \\
&= \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{xx}(\omega) \cos \tau\omega d\omega \quad [\text{since } S_{xx}(\omega) \text{ is even}] \\
R_{xx}(0) - R_{xx}(\tau) &= \frac{1}{2\pi} \int_{-\sigma}^{\sigma} S_{xx}(\omega) (1 - \cos \tau\omega) d\omega \quad [\text{since } \{X(t)\} \text{ is band limited}] \\
&= \frac{1}{2\pi} \int_{-\sigma}^{\sigma} S_{xx}(\omega) \times 2 \sin^2\left(\frac{\tau\omega}{2}\right) d\omega
\end{aligned} \tag{1}$$

From trigonometry,  $|\sin \theta| \leq \theta$

$$\therefore \sin^2 \theta \leq \theta^2$$

$$\therefore 2 \sin^2\left(\frac{\tau\omega}{2}\right) \leq \frac{\tau^2 \omega^2}{2} \tag{2}$$

Inserting (2) in (1),

$$R_{xx}(0) - R_{xx}(\tau) \leq \frac{1}{2\pi} \int_{-\sigma}^{\sigma} S_{xx}(\omega) \frac{\tau^2 \omega^2}{2} d\omega$$

$$\begin{aligned} &\leq \frac{\sigma^2 \tau^2}{4\pi} \int_{-\sigma}^{\sigma} S_{xx}(\omega) d\omega \\ &\leq \frac{\sigma^2 \tau^2}{4\pi} \int_{-\infty}^{\infty} S_{xx}(\omega) d\omega \\ \text{i.e., } &\leq \frac{\sigma^2 \tau^2}{2} R_{xx}(0) \end{aligned}$$

### Example 9

Consider a white Gaussian noise of zero mean and power spectral density  $N_0/2$  applied to a low pass RC filter whose transfer function is  $H(f) = \frac{1}{1 + i2\pi fRC}$ .

Find the autocorrelation function of the output random process.

The simple RC circuit for which the transfer function is given is a linear time-invariant system.

The power spectral densities of the input  $\{X(t)\}$  and the output  $\{Y(t)\}$  of a linear system are connected by

$$S_{YY}(\omega) = |H(\omega)|^2 S_{xx}(\omega)$$

In the problem, the transfer function is expressed in terms of the frequency  $f$ . Therefore, the above relation is

$$\begin{aligned} S_{yy}(f) &= |H(f)|^2 S_{xx}(f) \\ &= \frac{1}{1 + 4\pi^2 f^2 R^2 C^2} \frac{N_0}{2} \quad (\text{since the input is a white noise}) \\ \therefore R_{yy}(\tau) &= \frac{N_0}{2} \int_{-\infty}^{\infty} \frac{e^{i2\pi\tau f}}{1 + 4\pi^2 f^2 R^2 C^2} df \\ &= \frac{N_0}{8\pi^2 R^2 C^2} \int_{-\infty}^{\infty} \frac{e^{i(2\pi\tau)f} df}{\left(\frac{1}{2\pi RC}\right)^2 + f^2} \end{aligned} \tag{1}$$

Compare the integral in (1) with  $\int_{-\infty}^{\infty} \frac{e^{imx} dx}{a^2 + x^2}$ , which can be evaluated by contour integration technique [refer to the Worked Example 12 in Section 6(c) ‘power spectral density’].

$$\int_{-\infty}^{\infty} \frac{e^{imx} dx}{a^2 + x^2} = \frac{\pi}{a} e^{-|m|a} \tag{2}$$

Using (2) in (1),

$$\begin{aligned} R_{yy}(\tau) &= \frac{N_0}{8\pi^2 R^2 C^2} \pi \times 2\pi RC e^{-|2\pi\tau|/2\pi RC} \\ &= \frac{N_0}{4RC} e^{-|\tau|/RC} \end{aligned}$$

The mean square value of  $\{Y(t)\}$  is given by

$$E\{Y^2(t)\} = R_{yy}(0) = \frac{N_0}{4RC}$$

### Example 10

If  $Y(t) = A \cos(\omega_0 t + \theta) + N(t)$ , where  $A$  is a constant,  $\theta$  is a random variable with a uniform distribution in  $(-\pi, \pi)$  and  $\{N(t)\}$  is a band-limited Gaussian white noise with a power spectral density

$$S_{NN}(\omega) = \begin{cases} \frac{N_0}{2}, & \text{for } |\omega - \omega_0| < \omega_B \\ 0, & \text{elsewhere} \end{cases}$$

find the power spectral density of  $\{Y(t)\}$ . Assume that  $N(t)$  and  $\theta$  are independent.

$$\begin{aligned} Y(t_1) Y(t_2) &= \{A \cos(\omega_0 t_1 + \theta) + N(t_1)\} \{A \cos(\omega_0 t_2 + \theta) + N(t_2)\} \\ &= A^2 \cos(\omega_0 t_1 + \theta) \cos(\omega_0 t_2 + \theta) + N(t_1) N(t_2) \\ &\quad + A \cos(\omega_0 t_1 + \theta) N(t_2) + A \cos(\omega_0 t_2 + \theta) N(t_1) \\ \therefore R_{YY}(t_1, t_2) &= A^2 E\{\cos(\omega_0 t_1 + \theta) \cos(\omega_0 t_2 + \theta) + R_{NN}(t_1, t_2) \\ &\quad + AE\{\cos(\omega_0 t_1 + \theta)\} E\{N(t_2)\} \\ &\quad + AE\{\cos(\omega_0 t_2 + \theta)\} E\{N(t_1)\}\} \quad (\text{by independence}) \end{aligned}$$

i.e.,  $R_{YY}(\tau) = \frac{A^2}{2} \cos \omega_0 \tau + R_{NN}(\tau)$  [since  $\{N(t)\}$  is stationary]

$$\begin{aligned} \therefore S_{YY}(\omega) &= \frac{A^2}{2} \int_{-\infty}^{\infty} \cos \omega_0 \tau e^{-i\omega\tau} d\tau + S_{NN}(\omega) \\ &= \frac{\pi A^2}{2} \{\delta(\omega - \omega_0) + \delta(\omega + \omega_0)\} + S_{NN}(\omega) \end{aligned}$$

where  $S_{NN}(\omega)$  is given in the problem.

### Exercise 7(A)

#### Part-A (Short-answer Questions)

1. Define a Gaussian process.
2. When is a random process said to be normal?

3. What is the output of (i) a dc meter and (ii) a true rms meter when it measures a voltage  $V(t)$  which is a Gaussian ergodic random process with a mean of zero and variance of 4 V?
4. Give the  $n$ th order density of Gaussian process (or) the  $n$ th order normal density function.
5. What is the importance of Gaussian process in signal analysis?
6. State the properties of a Gaussian process.
7. If a Gaussian process is WSS, prove that it is also SSS.
8. If the member functions of a Gaussian process are uncorrelated, prove that they are independent.
9. Name a few random processes that are defined in terms of stationary Gaussian process.
10. Define square law detector process.
11. Define full-wave linear detector process.
12. Define half-wave linear detector process.
13. Define hard limiter process.
14. Define a band pass process or band pass signal.
15. Define a low pass process or an ideal low pass process.
16. Define a narrow band process (or) when is a random process said to be quasimonochromatic?
17. When is a random process called monochromatic?
18. Give the representation of a narrow band Gaussian process  $\{X(t)\}$  that is frequently used in signal analysis.
19. What is the justification for the representation of a narrow band Gaussian process  $\{X(t)\}$  in the form  
$$X(t) = R_x(t) \cos \{\omega_0 \pm \theta_x(t)\}$$
20. Define the envelope and phase of a narrow-band Gaussian process.
21. What are the distributions followed by the envelope and phase of a narrow-band Gaussian process?
22. Define the inphase and quadrature components of a narrow-band Gaussian process.
23. What is the nature of the inphase and quadrature components of a band pass process.
24. What do you mean by the quadrature representation of a WSS process?
25. What kind of random processes can be represented in the quadrature form?
26. For what kind of random processes is the quadrature representation useful in communication theory?
27. What do you mean by ‘noise’ in signal transmission?
28. How are external and internal noises caused?
29. What is thermal noise? By what type of random process is it represented?
30. If  $\{N(t)\}$  is a thermal noise, what is the nature of the graph of  $S_{NN}(\omega)$ ?
31. Why is thermal noise called white noise?

32. If the input to a linear time-invariant system is white noise  $\{N(t)\}$ , what is PSDF of the output?
33. If the PSD of white noise is  $\frac{N_0}{2}$ , find its ACF.
34. Find the average power (or) the mean square value of the white noise  $\{N(t)\}$ .
35. Why is the spectral density of the white noise  $\{N(t)\}$  not physically reliable? How is this difficulty overcome?
36. Define band-limited white noise.
37. State a few properties of band-limited white noise.
38. Find the ACF of the band-limited white noise.
39. Find the average power of the band-limited white noise.
40. What is meant by a filter in electrical systems.
41. Explain band pass and low pass filters.
42. How is filtering done in electrical systems?

**Part-B**

43. Suppose that  $\{X(t)\}$  is a normal process with  $\eta(t) = 3$  and  $C(t_1, t_2) = 4e^{-0.2|t_1-t_2|}$ .
- Find the probability that  $X(5) \leq 2$ , and
  - Find the probability that  $|X(8) - X(5)| \leq 1$ .
44. Prove that the RVs  $X_1, X_2, \dots, X_n$  are jointly normal, if the sum  $a_1X_1 + a_2X_2 + \dots + a_nX_n$  is a normal RV for any set of constants  $a_1, a_2, \dots, a_n$ .
- [Hint:** From the given condition
- $$Z = \omega_1X_1 + \omega_2X_2 + \dots + \omega_nX_n \text{ is a normal RV}$$
- $$E(Z) = 0 \text{ and } E(Z)^2 = \sigma_z^2 \text{ if } E(X_i) = 0$$
- $$\phi(Z) = e^{-\sigma_z^2/2}]$$
45. Let  $\sum_{k=1}^n X_k$ , where the  $X_k$ 's are a set of independent RVs each normally distributed with mean  $\mu$  and variance  $\sigma^2$ . Show that  $\{Y_n; n = 1, 2, \dots\}$  is a normal (Gaussian) process.
46. Given a normal process  $\{X(t)\}$  with  $\mu_x = 0$  and  $R_x(i) = 4e^{-2|t|}$ , we form the RVs  $Z = X(t+1)$ ,  $W = X(t-1)$ .
- Find  $E(ZW)$  and  $E\{(Z + W)^2\}$  and
  - Find  $f_Z(z)$ ,  $P\{Z < 1\}$  and  $f_{ZW}(z, w)$ .
47. Show that, if the RVs  $X, Y, Z$  are jointly normal and independent in pairs, then they are independent.
- [Hint:** Prove that  $f_{XYZ}(x, y, z) = f_X(x)f_Y(y)f_Z(z)$ .]
48. A voltage  $V(t)$ , which is Gaussian ergodic random process with a mean of zero and a variance of 4 V, is measured (i) by a dc meter, (ii) a true rms meter, and (iii) a meter which first squares  $V(t)$  and then read its dc component. Find the output of each meter and justify your answer.

[**Hint:** If  $\{X(t)\}$  is an ergodic process such as a voltage waveform,  $\mu_x$  can be measured by using a dc voltmeter and  $\sigma_x$  can be measured by using a true rms (ac coupled) voltmeter.]

49.  $\{X(t)\}$  is a Gaussian random process with mean  $\mu_x(t)$  and autocorrelation function  $R_{xx}(t_1, t_2)$ . Find  $E\{X(t_2)/X(t_1)\}$ ,  $t_1 < t_2$ .
50. If  $Z(t) = X \cos \omega_0 t + Y \sin \omega_0 t$ , where  $X$  and  $Y$  are independent Gaussian RVs with zero mean and unit variance and  $\omega_0$  is a constant,
  - (i) Show that  $\{Z(t)\}$  is a Gaussian random process.
  - (ii) Find the joint pdf of  $Z(t_1)$  and  $Z(t_2)$ .
  - (iii) Is the process WSS?
  - (iv) Is the process SSS?
  - (v) Find  $E\{Z(t_2)/Z(t_1)\}$ ,  $t_2 > t_1$ .
51. If  $\{X(t)\}$  is a Gaussian process with  $\mu_x = 0$  and  $R_{xx}(\tau) = 0$  for  $|\tau| > a$ , prove that it is correlation-ergodic.
52. If  $\{X(t)\}$  is a zero mean low pass process with a bandwidth of  $\omega_B$ , prove that
 
$$E\{X(t + \tau) - X(t)\}^2 \leq \omega_B^2 \tau^2 E\{X^2(t)\}$$

[**Hint:** Refer to Worked Example 8.]
53. If  $E\{X(t)\}$  is a band limited low pass process with bandwidth  $\omega_B$ , prove that  $R(\tau) \geq R(0) \cos \omega_B \tau$ , for  $|\tau| < \frac{\pi}{2\omega_B}$ .
54. For a narrow band process  

$$X(t) = X_c(t) \cos \omega_0 t + X_s(t) \sin \omega_0 t$$
, where  $\{X_c(t)\}$  and  $\{X_s(t)\}$  are stationary, uncorrelated, lowpass processes with

$$S_{x_c x_c}(\omega) = S_{x_s x_s}(\omega) = \begin{cases} g(\omega), & |\omega| < \omega_B \\ 0, & |\omega| > \omega_B \end{cases}$$

$$\text{show that } S_{xx}(\omega) = \frac{1}{2}\{g(\omega - \omega_0) + g(\omega + \omega_0)\}$$

55. Determine the autocorrelation of white noise.  
 (See the property of white noise)
56. Show that the narrow band noise  $\{N(t)\}$  can be represented as  $n(t) = n_c(t) \cos \omega_c t - n_s(t) \sin \omega_c t$ , where  $n_c(t)$  and  $n_s(t)$  are inphase and quadrature phase components of  $n(t)$  and  $\omega_c$  is the centre frequency of the band.  

[**Hint:** This is the same as the book-work on ‘quadrature representation’ of a WSS process.  $X(t)$ ,  $I(t)$  and  $\theta(t)$  are in the book-work must be replaced by  $n(t)$ ,  $n_c(t)$  and  $n_s(t)$  respectively].
57. (i) If  $\{N(t)\}$  is a band limited white noise such that

$$S_{NN}(\omega) = \frac{N_0}{2}, \text{ for } |\omega| < \omega_B \\ = 0, \text{ elsewhere}$$

find the autocorrelation of  $\{N(t)\}$ .

- (ii) If  $\{N(t)\}$  is a band limited white noise centered at a carrier frequency  $\omega_0$  such that

$$S_{NN}(\omega) = \frac{N_0}{2}, \text{ for } |\omega - \omega_0| < \omega_B \\ = 0, \text{ elsewhere}$$

Find the autocorrelation of  $\{N(t)\}$ .

58. A source of noise is a Gaussian with a mean of 0.4 V and a standard deviation of 0.15 V. For what percentage of time would you expect the measured noise voltage to exceed 0.7 V?

[Hint: Noise is a normal RV with the given parameters.]

59. If  $\{N(t)\}$  is a band-pass white noise and  $\{N_c(t)\}$  and  $\{N_s(t)\}$  are its quadrature components, prove that rms values of  $\{N(t)\}$ ,  $\{N_c(t)\}$  and  $\{N_s(t)\}$  are equal.

60. If  $X(t) = A \cos \omega_0 t + B \sin \omega_0 t$ , find the Hilbert transform of  $X(t)$ .

61. If the input to a linear time-invariant system is a zero mean, white Gaussian process  $\{N(t)\}$  and  $\{Y(t)\}$  is the output, prove that

(i)  $E\{Y(t)\} = 0$ ,

(ii)  $R_{YY}(\tau) = \frac{N_0}{2} \delta(\tau) * h(\tau) * h(-\tau)$ , and

(iii)  $S_{YY}(\omega) = \frac{N_0}{2} |H(\omega)|^2$ .

62. The impulse response of a low pass filter is  $\alpha e^{-\alpha t} U(t)$ ; where  $\alpha = \frac{1}{RC}$ .

If a zero mean, white Gaussian process  $\{N(t)\}$  is input into this filter, find the mean square value and autocorrelation function of the output.

## Poisson Process

There are many practical situations where the random times of occurrences of some specific events are of primary interest. For example, we may want to study the times at which components fail in a large system or the times at which jobs enter the queue in a computer system or the times of arrival of phone calls at an exchange or the times of emission of electrons from the cathode of a vacuum tube. In these examples, our main interest may not be the event itself but the sequence of random time instants at which the events occur. An ensemble of discrete sets of points from the time domain called a *point process* is used to model and analyse phenomena such as the ones mentioned above. An independent increments points process, i.e., a point process with the property that the number of occurrences in any finite collection of nonoverlapping time intervals are independent RVs, leads to a Poisson process.

**Definition:** If  $X(t)$  represents the number of occurrences of a certain event in  $(0, t)$ , then the discrete random process  $\{X(t)\}$  is called the Poisson process, provided the following postulates are satisfied:

- (i)  $P[1 \text{ occurrence in } (t, t + \Delta t)] = \lambda \Delta t + O(\Delta t)$
- (ii)  $P[0 \text{ occurrence in } (t, t + \Delta t)] = 1 - \lambda \Delta t + O(\Delta t)$
- (iii)  $P[2 \text{ or more occurrences in } (t, t + \Delta t)] = O(\Delta t)$
- (iv)  $X(t)$  is independent of the number of occurrences of the event in any interval prior and after the interval  $(0, t)$ .
- (v) The probability that the event occurs a specified number of times in  $(t_0, t_0 + t)$  depends only on  $t$ , but not on  $t_0$ .

## Probability Law for the Poisson Process $\{X(t)\}$

Let  $\lambda$  be the number of occurrences of the event in unit time.

$$\begin{aligned} \text{Let } P_n(t) &= P\{X(t) = n\} \\ \therefore P_n(t + \Delta t) &= P\{X(t + \Delta t) = n\} \\ &= P\{(n - 1) \text{ calls in } (0, t) \text{ and } 1 \text{ call in } (t, t + \Delta t)\} \\ &\quad + P\{n \text{ calls in } (0, t) \text{ and no call in } (t, t + \Delta t)\} \\ &= P_{n-1}(t) \lambda \Delta t + P_n(t) (1 - \lambda \Delta t) \text{ (by the postulates)} \\ \therefore \frac{P_n(t + \Delta t) - P_n(t)}{\Delta t} &= \lambda \{P_{n-1}(t) - P_n(t)\} \end{aligned}$$

Taking the limits as  $\Delta t \rightarrow 0$

$$\frac{d}{dt} P_n(t) = \lambda \{P_{n-1}(t) - P_n(t)\} \quad (1)$$

Let the solution of the equation (1) be

$$P_n(t) = \frac{(\lambda t)^n}{n!} f(t) \quad (2)$$

Differentiating (2) with respect to  $t$ ,

$$P'_n(t) = \frac{\lambda^n}{n!} \{n t^{n-1} f(t) + t^n f'(t)\} \quad (3)$$

Using (2) and (3) in (1),

$$\frac{\lambda^n}{n!} t^n f'(t) = -\lambda \frac{(\lambda t)^n}{n!} f(t)$$

i.e.,  $f'(t) = -\lambda f(t)$

$$\therefore f(t) = k e^{-\lambda t} \quad (4)$$

From (2),  $f(0) = P_0(0) = P\{X(0) = 0\}$

$$\begin{aligned} &= P\{\text{no event occurs in } (0, 0)\} \\ &= 1 \end{aligned} \quad (5)$$

Using (5) in (4), we get  $k = 1$  and hence,

$$f(t) = e^{-\lambda t} \quad (6)$$

Using (6) in (2),

$$P_n(t) = P\{X(t) = n\} = \frac{e^{-\lambda t} (\lambda t)^n}{n!}, \quad n = 0, 1, 2, \dots, \infty$$

Thus the probability distribution of  $X(t)$  is the Poisson distribution with parameter  $\lambda t$ .

**Note** We have assumed that the rate of occurrence of the event  $\lambda$  is a constant, but it can be function of  $t$  also. When  $\lambda$  is a constant, the process is called a **homogeneous Poisson process**. Unless specified otherwise, the Poisson process will be assumed homogeneous.

## Second-Order Probability Function of a Homogeneous Poisson Process

$$\begin{aligned} P[X(t_1) = n_1, X(t_2) = n_2] &= P[X(t_1) = n_1] P[X(t_2) = n_2 | X(t_1) = n_1], \quad t_2 > t_1 \\ &= P[X(t_1) = n_1] P[\text{the event occurs } (n_2 - n_1) \text{ times in the interval of length } (t_2 - t_1)] \\ &= \frac{e^{-\lambda t_1} (\lambda t_1)^{n_1}}{n_1!} \frac{e^{-\lambda(t_2-t_1)} \{\lambda(t_2-t_1)\}^{n_2-n_1}}{n_2!} \text{ if } n_2 \geq n_1 \\ &= \begin{cases} \frac{e^{-\lambda t_2} \lambda^{n_2} t_1^{n_1} (t_2 - t_1)^{n_2 - n_1}}{n_1! n_2!}, & n_2 \geq n_1 \\ 0, & \text{otherwise} \end{cases} \end{aligned}$$

Proceeding similarly, we can get the third-order probability function as

$$\begin{aligned} P[X(t_1) = n_1, X(t_2) = n_2, X(t_3) = n_3] &= \frac{e^{-\lambda t_3} \lambda^{n_3} t_1^{n_1} t_2^{n_2} (t_2 - t_1)^{n_2 - n_1} (t_3 - t_2)^{n_3 - n_2}}{n_1! n_2! n_3!} \text{ if } n_3 \geq n_2 \geq n_1 \\ &= 0, \quad \text{otherwise} \end{aligned}$$

## Mean and Autocorrelation of the Poisson Process

The probability law of the Poisson process  $\{X(t)\}$  is the same as that of a Poisson distribution with parameter  $\lambda t$ .

$$\begin{aligned} \therefore E\{X(t)\} &= \text{Var}\{X(t)\} = \lambda t \\ \therefore E\{X^2(t)\} &= \lambda t + \lambda^2 t^2 \quad (1) \\ R_{xx}(t_1, t_2) &= E\{X(t_1) X(t_2)\} \\ &= E[X(t_1) \{X(t_2) - X(t_1) + X(t_1)\}] \\ &= E[X(t_1) \{X(t_2) - X(t_1)\}] + E\{X^2(t_1)\} \\ &= E[X(t_1)] E[X(t_2) - X(t_1)] + E\{X^2(t_1)\}. \end{aligned}$$

Since  $\{X(t)\}$  is a process of independent increments,

$$\begin{aligned} &= \lambda t_1 [\lambda (t_2 - t_1) + \lambda t_1 + \lambda^2 t_1^2], \text{ if } t_2 \geq t_1 \\ &= \lambda^2 t_1 t_2 + \lambda t_1, \text{ if } t_2 \geq t_1 \end{aligned} \quad [\text{by (1)}]$$

or  $R_{xx}(t_1, t_2) = \lambda^2 t_1 t_2 + \lambda \min(t_1, t_2)$

$$\begin{aligned} C_{xx}(t_1, t_2) &= R_{xx}(t_1, t_2) - E\{X(t_1)\} E\{X(t_2)\} \\ &= \lambda^2 t_1 t_2 + \lambda t_1 - \lambda^2 t_1 t_2 \\ &= \lambda t_1, \text{ if } t_2 \geq t_1 \end{aligned}$$

or  $= \min(t_1, t_2)$

$$\begin{aligned} r_{xx}(t_1, t_2) &= \frac{C_{xx}(t_1, t_2)}{\sqrt{\text{Var}\{X(t_1)\} \text{Var}\{X(t_2)\}}} \\ &= \frac{\lambda t_1}{\sqrt{\lambda t_1 \lambda t_2}} = \sqrt{\frac{t_1}{t_2}}, \text{ if } t_2 \geq t_1 \end{aligned}$$

**Note** Poisson process is not a stationary process.

## Properties of Poisson Process

1. The Poisson process is a Markov process.

### Proof

Consider  $P[X(t_3) = n_3 | X(t_2) = n_2, X(t_1) = n_1]$

$$\begin{aligned} &= \frac{P[X(t_1) = n_1, X(t_2) = n_2, X(t_3) = n_3]}{P[X(t_1) = n_1, X(t_2) = n_2]} \\ &= \frac{e^{-\lambda(t_3-t_2)} \lambda^{n_3-n_2} (t_3 - t_2)^{n_3-n_2}}{|n_3 - n_2|} \\ &\quad [\text{refer to the second- and third-order probability functions of the Poisson process}] \\ &= P[X(t_3) = n_3 | X(t_2) = n_2] \end{aligned}$$

This means that the conditional probability distribution of  $X(t_3)$  given all the past values  $X(t_1) = n_1, X(t_2) = n_2$  depends only on the most recent value  $X(t_2) = n_2$ .

That is, the Poisson process possesses the Markov property. Hence, the result.

2. **Additive property:** Sum of two independent Poisson processes is a Poisson process.

### Proof

We have already derived in Chapter 4 the characteristic function of a Poisson distribution with parameter  $\lambda$  as  $e^{-\lambda(1-e^{i\omega})}$ .

Therefore, the characteristic functions of  $X_1(t)$  and  $X_2(t)$  are given by

$$\phi_{X_1(t)}(\omega) = e^{-\lambda_1 t(1-e^{i\omega})} \text{ and } \phi_{X_2(t)}(\omega) = e^{-\lambda_2 t(1-e^{i\omega})}$$

Since  $X_1(t)$  and  $X_2(t)$  are independent,

$$\begin{aligned}\phi_{X_1(t)+X_2(t)}(\omega) &= \phi_{X_1(t)}(\omega)\phi_{X_2(t)}(\omega) \\ &= e^{-(\lambda_1+\lambda_2)t(1-e^{i\omega})}\end{aligned}$$

which is characteristic function of Poisson distribution with parameter  $(\lambda_1 + \lambda_2)t$ .

Therefore,  $\{X_1(t) + X_2(t)\}$  is a Poisson process.

### Alternative proof

Let  $X(t) = X_1(t) + X_2(t)$ .

$$\begin{aligned}P\{X(t) = n\} &= \sum_{r=0}^n P\{X_1(t) = r\} P\{X_2(t) = n - r\} \\ &= \sum_{r=0}^n \frac{e^{-\lambda_1 t} (\lambda_1 t)^r}{|r|} \frac{e^{-\lambda_2 t} (\lambda_2 t)^{n-r}}{|n-r|} \\ &= e^{-(\lambda_1+\lambda_2)t} \frac{1}{|n|} \sum_{r=0}^n n C_r (\lambda_1 t)^r (\lambda_2 t)^{n-r} \\ &= e^{-(\lambda_1+\lambda_2)t} [(\lambda_1 + \lambda_2)t]^n / |n|\end{aligned}$$

Therefore,  $X_1(t) + X_2(t)$  is a Poisson process with parameter  $(\lambda_1 + \lambda_2)t$ .

**Note** The additive property holds good for any number of independent Poisson processes.

3. Difference of two independent Poisson processes is not a Poisson process.

### Proof

Let  $X(t) = X_1(t) - X_2(t)$

$$\begin{aligned}E\{X(t)\} &= E\{X_1(t)\} - E\{X_2(t)\} \\ &= (\lambda_1 - \lambda_2)t \\ E\{X^2(t)\} &= E\{X_1^2(t)\} + E\{X_2^2(t)\} - 2 E\{X_1(t)\} E\{X_2(t)\} \\ &\quad \text{(by independence)} \\ &= (\lambda_1^2 t^2 + \lambda_1 t) + (\lambda_2^2 t^2 + \lambda_2 t) - 2(\lambda_1 t)(\lambda_2 t) \\ &= (\lambda_1 + \lambda_2)t + (\lambda_1 - \lambda_2)^2 t^2 \\ &\neq (\lambda_1 - \lambda_2)t + (\lambda_1 - \lambda_2)^2 t^2\end{aligned}$$

Recall that  $E\{X^2(t)\}$  for a Poisson process  $\{X(t)\}$  with parameter  $\lambda$  is given by  $E\{X^2(t)\} = \lambda t + \lambda^2 t^2$ .

Therefore,  $\{X_1(t) - X_2(t)\}$  is not a Poisson process.

4. The interarrival time of a Poisson process, i.e., the interval between two successive occurrences of a Poisson process with parameter  $\lambda$  has an exponential distribution with mean  $1/\lambda$ .

**Proof**

Let two consecutive occurrences of the event be  $E_i$  and  $E_{i+1}$ .

Let  $E_i$  take place at time instant  $t_i$  and  $T$  be the interval between the occurrences of  $E_i$  and  $E_{i+1}$ .  $T$  is a continuous RV.

$$\begin{aligned} P(T > t) &= P\{E_{i+1} \text{ did not occur in } (t_i, t_i + t)\} \\ &= P\{\text{No event occurs in an interval of length } t\} \\ &= P\{X(t) = 0\} \\ &= e^{-\lambda t} \end{aligned}$$

Therefore, the cdf of  $T$  is given by

$$F(t) = P\{T \leq t\} = 1 - e^{-\lambda t}$$

Therefore, the pdf of  $T$  is given by

$$f(t) = \lambda e^{-\lambda t} (t \geq 0)$$

which is an exponential distribution with mean  $1/\lambda$ .

5. If the number of occurrences of an event  $E$  in an interval of length  $t$  is a Poisson process  $\{X(t)\}$  with parameter  $\lambda$  and if each occurrence of  $E$  has a constant probability  $p$  of being recorded and the recordings are independent of each other, then the number  $N(t)$  of recorded occurrences in  $t$  is also a Poisson process with parameter  $\lambda p t$ .

**Proof**

$$\begin{aligned} P\{N(t) = n\} &= \sum_{r=0}^{\infty} P\{\text{$E$ occurs $(n+r)$ times in $t$ and $n$ of them are recorded}\} \\ &= \sum_{r=0}^{\infty} \frac{e^{-\lambda t} (\lambda t)^{n+r}}{|n+r|} (n+r) C_n p^n q^r, \quad q = 1-p \\ &= \sum_{r=0}^{\infty} \frac{e^{-\lambda t} (\lambda t)^{n+r}}{|n+r|} \frac{|n+r|}{|n|r} p^n q^r \\ &= \frac{e^{-\lambda t} (\lambda p t)^n}{|n|} \sum_{r=0}^{\infty} \frac{(\lambda q t)^r}{|r|} \\ &= \frac{e^{-\lambda t} (\lambda p t)^n}{|n|} e^{\lambda q t} \\ &= \frac{e^{-\lambda p t} (\lambda p t)^n}{|n|} \end{aligned}$$

### Worked Example 7(B)

#### Example 1

Suppose that customers arrive at a bank according to a Poisson process with a mean rate of 3 per minute; find the probability that during a time interval of 2 min (a) exactly 4 customers arrive, and (b) more than 4 customers arrive.

Mean of the Poisson process =  $\lambda t$

Mean arrival rate = mean number of arrivals per minute (unit time) =  $\lambda$

Given  $\lambda = 3$

$$P\{X(t) = k\} = \frac{e^{-\lambda t} (\lambda t)^k}{k!}$$

$$(a) \quad P\{X(2) = 4\} = \frac{e^{-6} 6^4}{4!} = 0.133$$

$$\begin{aligned} (b) \quad P\{X(2) > 4\} &= 1 - [P\{X(2) = 0\} + P\{X(2) = 1\} + P\{X(2) = 2\} \\ &\quad + P\{X(2) = 3\} + P\{X(2) = 4\}] \\ &= 1 - \sum_{k=0}^4 e^{-6} 6^k / k! \\ &= 0.715 \end{aligned}$$

#### Example 2

A machine goes out of order, whenever a component fails. The failure of this part follows a Poisson process with a mean rate of 1 per week. Find the probability that 2 weeks have elapsed since last failure. If there are 5 spare parts of this component in an inventory and that the next supply is not due in 10 weeks, find the probability that the machine will not be out of order in the next 10 weeks.

(a) Here, the unit time is 1 week.

Mean failure rate = mean number of failures in a week =  $\lambda = 1$ .

$P\{\text{no failures in the 2 weeks since last failure}\}$

$$\begin{aligned} &= P\{X(2) = 0\} \\ &= \frac{e^{-2\lambda} (2\lambda)^0}{0!} = e^{-2} = 0.135 \end{aligned}$$

(b) There are only 5 spare parts and the machine should not go out of order in the next 10 weeks.

$P\{\text{for this event}\} = P\{X(10) \leq 5\}$

$$\begin{aligned} &= \sum_{k=0}^5 \frac{e^{-10} 10^k}{k!} \\ &= 0.068 \end{aligned}$$

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### Example 3

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If  $\{N_1(t)\}$  and  $\{N_2(t)\}$  are 2 independent Poisson processes with parameters  $\lambda_1$  and  $\lambda_2$  respectively, show that

$$P[N_1(t) = k \mid \{N_1(t) + N_2(t) = n\}] = nC_k p^k q^{n-k}, \text{ where}$$

$$p = \frac{\lambda_1}{\lambda_1 + \lambda_2} \text{ and } q = \frac{\lambda_2}{\lambda_1 + \lambda_2}.$$

Required conditional probability

$$\begin{aligned} &= \frac{P[\{N_1(t) = k\} \cap \{N_1(t) + N_2(t) = n\}]}{P\{N_1(t) + N_2(t) = n\}} \\ &= \frac{P[\{N_1(t) = k\} \cap \{N_2(t) = n - k\}]}{P\{N_1(t) + N_2(t) = n\}} \\ &= \frac{\frac{e^{-\lambda_1 t} (\lambda_1 t)^k}{k!} \times \frac{e^{-\lambda_2 t} (\lambda_2 t)^{n-k}}{(n-k)!}}{\frac{e^{-(\lambda_1 + \lambda_2)t} ((\lambda_1 + \lambda_2)t)^n}{n!}} \\ &\quad \text{(by independence and additive property)} \end{aligned}$$

$$\begin{aligned} &= \frac{n!}{k! (n-k)!} \frac{(\lambda_1 t)^k (\lambda_2 t)^{n-k}}{((\lambda_1 + \lambda_2)t)^n} \\ &= nC_k \left( \frac{\lambda_1}{\lambda_1 + \lambda_2} \right)^k \left( \frac{\lambda_2}{\lambda_1 + \lambda_2} \right)^{n-k} \\ &= nC_k p^k q^{n-k} \end{aligned}$$

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### Example 4

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If customers arrive at a counter in accordance with a Poisson process with a mean rate of 2 per minute, find the probability that the interval between 2 consecutive arrivals is (a) more than 1 min, (b) between 1 min and 2 min, and (c) 4 min or less.

Refer of Property 4 of Poisson processes.

The interval  $T$  between 2 consecutive arrivals follows an exponential distribution with parameter  $\lambda = 2$ .

$$(a) \quad P(T > 1) = \int_1^\infty 2e^{-2t} dt = e^{-2} = 0.135$$

$$(b) \quad P(1 < T < 2) = \int_1^2 2e^{-2t} dt = e^{-2} - e^{-4} = 0.117$$

$$(c) \quad P(T \leq 4) = \int_0^4 2e^{-2t} dt = 1 - e^{-8} = 0.999$$

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### Example 5

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A radioactive source emits particles at a rate of 5 per minute in accordance with Poisson process, each particle emitted has a probability 0.6 of being recorded. Find the probability that 10 particles are recorded in a 4-min period.

Refer to Property 5 of Poisson processes. The number of recorded particles  $N(t)$  follows a Poisson process with parameter  $\lambda p$ .

Here,  $\lambda = 5$  and  $p = 0.6$

$$\therefore P\{N(t) = k\} = \frac{e^{-3t}(3t)^k}{|k|}$$

$$\begin{aligned}\therefore P\{N(4) = 10\} &= \frac{e^{-12}(12)^{10}}{|10|} \\ &= 0.104\end{aligned}$$

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### Example 6

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The number of accidents in a city follows a Poisson process with a mean of 2 per day and the number  $X_i$  of people involved in the  $i$ th accident has the distribution (independent)  $P\{X_i = k\} = \frac{1}{2^k}$  ( $k \geq 1$ ). Find the mean and variance of the number of people involved in accidents per week.

The mean and variance of the distribution  $P\{X_i = k\} = \frac{1}{2^k}$ ,  $k = 1, 2, 3, \dots, \infty$  can be obtained as 2 and 2.

Let the number of accidents on any day be assumed as  $n$ .

The numbers of people involved in these accidents be  $X_1, X_2, \dots, X_n$ .

$X_1, X_2, \dots, X_n$  are independent and identically distributed RVs with mean 2 and variance 2.

Therefore, by central limit theorem,  $(X_1 + X_2 + \dots + X_n)$  follows a normal distribution with mean  $2n$  and variance  $2n$ , i.e., the total number of people involved in all the accidents on a day with  $n$  accidents =  $2n$ .

If  $N$  denotes the number of people involved in accidents on any day, then

$$P\{N = 2n\} = P\{X(t) = n\} \quad [\text{where } X(t) \text{ is the number of accidents}]$$

$$= \frac{e^{-2t}(2t)^n}{|n|} \quad (\text{by data})$$

$$\begin{aligned}\therefore E\{N\} &= \sum_{n=0}^{\infty} \frac{2ne^{-2t}(2t)^n}{|n|} \\ &= 2E\{X(t)\} = 4t\end{aligned}$$

$$\text{Var}\{N\} = E\{N^2\} - E^2(N)$$

$$\begin{aligned}
 &= \sum_{n=0}^{\infty} \frac{4n^2 e^{-2t} (2t)^n}{n!} - 16t^2 \\
 &= 4 E\{X^2(t)\} - 16t^2 \\
 &= 4[\text{Var}(X(t)) + E^2\{X(t)\}] - 16t^2 \\
 &= 4[2t + 4t^2] - 16t^2 = 8t
 \end{aligned}$$

Therefore, mean and variance of the number of people involved in accidents per week are 28 and 56 respectively.

### Example 7

If  $T_n$  is the RV denoting the time of occurrence of the  $n$ th event in a Poisson process with parameter  $\lambda$ , show that the distribution function  $F_n(t)$  of  $T_n$  is given by

$$F_n(t) = \begin{cases} 1 - \sum_{k=0}^{n-1} \frac{(\lambda t)^k}{k!} e^{-\lambda t}, & \text{if } t \geq 0 \\ 0, & \text{if } t < 0 \end{cases}$$

Deduce the density function  $f_n(t)$  of  $T_n$

$$\begin{aligned}
 F_n(t) &= P\{T_n \leq t\} \\
 &= 1 - P\{T_n > t\}
 \end{aligned}$$

When  $T_n > t$ , i.e., the time of occurrence of the  $n$ th event  $> t$ ,  $(n-1)$  or less events must have occurred in  $(0, t)$ .

$$\begin{aligned}
 \therefore F_n(t) &= 1 - P\{X(t) \leq n-1\} \\
 &= 1 - \sum_{k=0}^{n-1} \frac{(\lambda t)^k}{k!} e^{-\lambda t}, \text{ when } t \geq 0
 \end{aligned}$$

Differentiating both sides with respect to  $t$  and noting that  $F'_n(t) = f_n(t)$

$$\begin{aligned}
 f_n(t) &= - \sum_{k=0}^{n-1} \left\{ -\lambda \frac{(\lambda t)^k}{k!} e^{-\lambda t} + \lambda \frac{(\lambda t)^{k-1}}{(k-1)!} e^{-\lambda t} \right\} \\
 &= \lambda e^{-\lambda t} \sum_{k=0}^{n-1} \left\{ \frac{(\lambda t)^k}{k!} - \frac{(\lambda t)^{k-1}}{(k-1)!} \right\} \\
 &= \lambda e^{-\lambda t} \times \left[ 1 + \left\{ \frac{\lambda t}{1!} - 1 \right\} + \left\{ \frac{(\lambda t)^2}{2!} - \frac{\lambda t}{1!} \right\} + \dots + \left\{ \frac{(\lambda t)^{n-1}}{(n-1)!} - \frac{(\lambda t)^{n-2}}{(n-2)!} \right\} \right] \\
 &= \frac{\lambda^n t^{n-1} e^{-\lambda t}}{(n-1)!}, \quad t \geq 0
 \end{aligned}$$

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## Example 8

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If  $\{X(t)\}$  is a Poisson process, prove that

$$\begin{aligned}
 P\{X(s) = r/X(t) = n\} &= nC_r \left(\frac{s}{t}\right)^r \left(1 - \frac{s}{t}\right)^{n-r} \quad \text{where } s < t \\
 P\{X(s) = r/X(t) = n\} &= \frac{P[\{X(s) = r\} \cap \{X(t) = n\}]}{P\{X(t) = n\}} \\
 &= \frac{P\{X(s) = r \cap X(t-s) = n-r\}}{P\{X(t) = n\}} \\
 &= \frac{P\{X(s) = r\} P\{X(t-s) = n-r\}}{P\{X(t) = n\}} \quad (\text{by independence}) \\
 &= \frac{\{e^{-\lambda t} (\lambda s)^r / [r]\} \{e^{-\lambda(t-s)} [\lambda(t-s)]^{n-r} / [n-r]\}}{e^{-\lambda t} (\lambda t)^n / [n]} \\
 &= \frac{\underline{n}}{\underline{r} \underline{n-r}} \frac{s^r (t-s)^{n-r}}{t^n} \\
 &= nC_r \left(\frac{s}{t}\right)^r \left(1 - \frac{s}{t}\right)^{n-r}
 \end{aligned}$$

### Exercise 7(B)

#### Part-A (Short-answer Questions)

1. What is a point process? Given an example.
2. Define a Poisson process.
3. What are the postulates of a Poisson process?
4. State the probability law of a Poisson process.
5. When is a Poisson process said to be homogeneous?
6. If  $\{X(t)\}$  is a homogeneous Poisson process, find  
 $P\{X(t_1) = n_1, X(t_2) = n_2\}, t_2 > t_1$
7. Find the autocorrelation  $R_{xx}(t_1, t_2)$  of the Poisson process  $\{X(t)\}$ .
8. State and prove the additive property of a Poisson process.
9. Prove that the difference of 2 independent Poisson process is not a Poisson process.
10. Prove that the interarrival time of a Poisson process follows an exponential distribution.
11. If the customers arrive at a bank according to a Poisson process with mean rate of 2 per minute, find the probability that, during an 1-min interval, no customer arrives.

**Part-B**

12. Find the first-order characteristic function of a Poisson process.
13. If particles are emitted from a radioactive source at the rate of 20 per hour, find the probability that exactly 5 particles are emitted during a 15-min period.
14. On the average, a submarine on patrol sights 6 enemy ships per hour. Assuming that the number of ships sighted in a given length of time is a Poisson variate, find the probability of sighting
  - (i) 6 ships in the next half-an-hour,
  - (ii) 4 ships in the next 2 h,
  - (iii) at least 1 ship in the next 15 min, and
  - (iv) at least 2 ships in the next 20 min.
15. Patients arrive randomly and independently at a doctor's consulting room from 8 a.m. at an average rate of one in 5 min. The waiting room can hold 12 persons. What is the probability that the room will be full when the doctor arrives at 9 a.m.?
16. Messages arrive at a telegraph office in accordance with the laws of a Poisson process with a mean rate of 3 messages per hour.
  - (i) What is the probability that no message will have arrived during the morning hours, i.e., between 8 a.m. and 12 noon?
  - (ii) What is the distribution of the time at which the first afternoon message arrives?
17. Assume that a circuit has an IC whose time to failure is an exponentially distributed RV with expected lifetime of 3 months. If there are 10 spare IC's and time from failure to replacement is zero, what is the probability that the circuit can be kept operational for at least 1 year?
18. Assume that an office switchboard has 5 telephone lines and that starting at 8 am on Monday, the time that a call arrives on each line is an exponential RV with parameter  $\lambda$ . Also assume that the calls arrive independently on the lines. Show that the time of arrival of the first call (irrespective of which line it arrives on) is exponential with parameter  $5\lambda$ .
19. A radioactive source emits particles at a rate of 6 per minute in accordance with Poisson process. Each particle emitted has a probability of  $1/3$  of being recorded. Find the probability that at least 5 particles are recorded in a 5-min period.
20. Suppose that customers arrive at a counter independently from 2 different sources. Arrivals occur in accordance with a Poisson process with mean rate of 6 per hour from the first source and 4 per hour from the second source. Find the mean interval between any 2 successive arrivals.
21. Assume that a device fails when a cumulative effect of  $k$  shocks occur. If the shocks occur according to a Poisson process with parameter  $\lambda$ , find the density function for the life  $T$  of the device.

[Hint: Refer of Worked Example 7]

22. In the case of a Poisson process, show that the conditional probability that events have occurred at  $\tau_1, \tau_2, \dots, \tau_n$ , given that  $n$  events have occurred in  $(0, t)$ , is given by  $\frac{|n|}{t^n}$ .
23. Passengers arrive at a terminal for boarding the next bus. The times of their arrival are Poisson with an average arrival rate of 1 per minute. The times of departure of each bus are Poisson with an average departure rate of 2 per hour. Assume that the capacity of the bus is large. Find the average number of passengers in (i) each bus, and (ii) the first bus that leaves after 9 a.m.
24. Passengers arrive at a terminal after 9 a.m. The times of their arrival are Poisson with mean density  $\lambda = 1$  per minute. The time interval from 9 am to the departure of the next bus is a RV  $T$ . Find the mean number of passengers in this bus (i) if  $T$  has an exponential density with mean 30 min, and (ii) if  $T$  is uniform between 0 and 60 min.

## Markov Process

Another interesting model of a random process is the one in which the value of the random process depends only upon the most recent previous value and is independent of all values in the more distant past. Such a model is called a Markov model and is often described by saying that a Markov process is one in which the future value is independent of the past values, given the present value. Models in which the future depends only upon the present are common among electrical engineering models.

Consider the experiment of tossing a fair coin a number of times. The possible outcomes at each trial are two—‘head’ with probability  $1/2$  and ‘tail’ with probability  $1/2$ . If we denote the outcome of the  $n$ th toss, which is a RV, by  $X_n$  and the outcomes ‘head’ and ‘tail’ by 1 and 0 respectively, then

$$P\{X_n = 1\} = \frac{1}{2} \text{ and } P\{X_n = 0\} = \frac{1}{2}; n = 1, 2, \dots$$

Thus we have a sequence of independent RVs  $X_1, X_2, \dots$ , since the trials are independent and hence the outcome of the  $n$ th trial does not depend in any way on the previous trials.

Consider now the RV that represents the total number of heads in the first  $n$  trials and is given by  $S_n = X_1 + \dots + X_n$ . The possible values of  $S_n$  are  $0, 1, 2, \dots, n$ . If  $S_n = k$  ( $k = 0, 1, \dots, n$ ), then the RV  $S_{n+1}$  ( $= S_n + X_{n+1}$ ) can assume only 2 possible values, namely  $k + 1$  [if the  $(n + 1)$ th trial results in a head] and  $k$  [if the  $(n + 1)$ th trial results in a tail].

Thus,  $P\{S_{n+1} = k + 1 / S_n = k\} = \frac{1}{2}$

$$P\{S_{n+1} = k / S_n = k\} = \frac{1}{2}$$

These probabilities are not at all affected by the values of the RVs  $S_1, S_2, \dots, S_{n-1}$ . Also the conditional probability of  $S_{n+1}$  given  $S_n$  depends on the value of  $S_n$  and not on the manner in which the value of  $S_n$  was reached. This is a simple example of a Markov chain. Random processes  $\{X(t)\}$  (with Markov property) which take discrete values, whether  $t$  is discrete or continuous, are called *Markov chains*. Poisson process, discussed earlier, is a continuous time Markov chain. In this section, we will discuss discrete time Markov chains.

## Definition of a Markov Chain

If, for all  $n$ ,  $P\{X_n = a_n | X_{n-1} = a_{n-1}, X_{n-2} = a_{n-2}, \dots, X_0 = a_0\} = P\{X_n = a_n | X_{n-1} = a_{n-1}\}$ , then the process  $\{X_n\}$ ,  $n = 0, 1, \dots$ , is called a Markov chain.

$(a_1, a_2, \dots, a_n, \dots)$  are called the states of the Markov chain. The conditional probability  $P\{X_n = a_j | X_{n-1} = a_i\}$  is called the *one-step transition probability* from state  $a_i$  to state  $a_j$  at the  $n^{\text{th}}$  step (trial) and is denoted by  $p_{ij}(n-1, n)$ .

If the one-step transition probability does not depend on the step, i.e.,  $p_{ij}(n-1, n) = p_{ij}(m-1, m)$  the Markov chain is called a *homogeneous Markov chain* or the chain is said to have stationary transition probabilities. The use of the word ‘stationary’ does not imply a stationary random sequence.

When the Markov chain is homogeneous, the one-step transition probability is denoted by  $p_{ij}$ . The matrix  $P = \{p_{ij}\}$  is called (one-step) *transition probability matrix*, shortly, tpm.

**Note** The tpm of a Markov chain is a stochastic matrix, since  $p_{ij} \geq 0$  and  $\sum_j p_{ij} = 1$  for all  $i$ , i.e., the sum of all the elements of any row of the tpm is 1. This is obvious because the transition from state  $a_i$  to any one of the states (including  $a_i$  itself) is a certain event.

The conditional probability that the process is in state  $a_j$  at step  $n$ , given that it was in state  $a_i$  at step 0, i.e.,  $P\{X_n = a_j | X_0 = a_i\}$  is called the *n-step transition probability* and denoted by  $p_{ij}(n)$ .

**Note**  $p_{ij}^{(l)} = p_{ij}$

Let us consider an example in which we explain how the tpm is formed for a Markov chain. Assume that a man is at an integral point of the  $x$ -axis between the origin and the point  $x = 3$ . He takes a unit step either to the right with probability 0.7 or to the left with probability 0.3, unless he is at the origin when he takes a step to the right to reach  $x = 1$  or he is at the point  $x = 3$ , when he takes a step to the left to reach  $x = 2$ . The chain is called ‘Random walk with reflecting barriers’.

The tpm is given below:

$$\text{States of } X_n \begin{array}{cccc} 0 & 1 & 2 & 3 \end{array}$$

$$\text{States of } X_{n-1} = \begin{matrix} 0 \\ 1 \\ 2 \\ 3 \end{matrix} \begin{pmatrix} 0 & 1 & 0 & 0 \\ 0.3 & 0 & 0.7 & 0 \\ 0 & 0.3 & 0 & 0.7 \\ 0 & 0 & 1 & 0 \end{pmatrix}$$

**Note**  $p_{23}$  = the element in the 2nd row, 3rd column of this tpm = 0.7. This means that, if the process is at state 2 at step  $(n-1)$ , the probability that it moves to state 3 at step  $n$  = 0.7, where  $n$  is any positive integer.

**Definition:** If the probability that the process is in state  $a_i$  is  $p_i$  ( $i = 1, 2, \dots, k$ ) at any arbitrary step, then the row vector  $p = (p_1, p_2, \dots, p_k)$  is called *the probability distribution of the process* at that time. In particular,  $p^{(0)} = \{p_1^{(0)}, p_2^{(0)}, \dots, p_k^{(0)}\}$  is the initial probability distribution.

**[Remark:** The transition probability matrix together with the initial probability distribution completely specifies a Markov chain  $\{X_n\}$ . In the example given above, let us assume that the initial probability distribution of the chain is

$$p^{(0)} = \left( \frac{1}{4}, \frac{1}{4}, \frac{1}{4}, \frac{1}{4} \right)$$

i.e.,  $P\{X_0 = i\} = 1/4, i = 0, 1, 2, 3$

Then we have, for the example given above,

$$\begin{aligned} P\{X_1 = 2/X_0 = 1\} &= 0.7; P\{X_2 = 1/X_1 = 2\} = 0.3, \\ P\{X_2 = 1, X_1 = 2/X_0 = 1\} &= P\{X_2 = 1/X_1 = 2\} \times P\{X_1 = 2/X_0 = 1\} \\ &= 0.3 \times 0.7 = 0.21 \end{aligned} \tag{1}$$

$$\begin{aligned} P\{X_2 = 1, X_1 = 2, X_0 = 1\} &= P\{X_0 = 1\} \times P\{X_2 = 1, X_1 = 2/X_0 = 1\} \\ &= 1/4 \times 0.21 = 0.0525 \quad [\text{by (1)}] \end{aligned} \tag{2}$$

$$\begin{aligned} P\{X_3 = 3, X_2 = 1, X_1 = 2, X_0 = 1\} &= P\{X_2 = 1, X_1 = 2, X_0 = 1\} \\ &\quad \times P\{X_3 = 3/X_2 = 1, X_1 = 2, X_0 = 1\} \\ &= 0.0525 P\{X_3 = 3/X_2 = 1\} \quad (\text{Markov property}) [\text{by (2)}] \\ &= 0.0525 \times 0 = 0 \end{aligned}$$

## Chapman–Kolmogorov Theorem

If  $P$  is the tpm of a homogeneous Markov chain, then the  $n$ -step tpm  $P^{(n)}$  is equal to  $P^n$ .

i.e.,  $[p_{ij}^{(n)}] = [p_{ij}]^n$

**Proof**

$$p_{ij}^{(2)} = P\{X_2 = j/X_0 = i\}, \text{ since the chain is homogeneous.}$$

The state  $j$  can be reached from the state  $i$  in 2 steps through some intermediate state  $k$ .

$$\begin{aligned} \text{Now, } p_{ij}^{(2)} &= P\{X_2 = j/X_0 = i\} = P\{X_2 = j, X_1 = k / X_0 = i\} \\ &= P\{X_2 = j/X_1 = k, X_0 = i\} P\{X_1 = k/X_0 = i\} \\ &= p_{kj}^{(1)} p_{ik}^{(1)} \\ &= p_{ik} p_{kj} \end{aligned}$$

Since the transition from state  $i$  to state  $j$  in 2 steps can take place through any one of the intermediate states,  $k$  can assume the values 1, 2, 3, .... The transitions through various intermediate states are mutually exclusive.

$$\text{Hence, } p_{ij}^{(2)} = \sum_k p_{ik} p_{kj}$$

i.e., the  $ij$ -th element of 2 step tpm =

the  $ij$ -th element of the product of the 2 one-step tpm's

$$\text{i.e., } P^{(2)} = P^2$$

$$\begin{aligned} \text{Now, } p_{ij}^{(3)} &= P\{X_3 = j/X_0 = i\} \\ &= \sum_k P\{X_3 = j/X_2 = k\} P\{X_2 = k/X_0 = i\} \\ &= \sum_k p_{kj} p_{ik}^{(2)} \\ &= \sum_k p_{ik}^{(2)} p_{kj} \end{aligned}$$

$$\text{Similarly, } p_{ij}^{(3)} = \sum_k p_{ik} p_{kj}^{(2)}$$

$$\text{i.e., } P^{(3)} = P^2 P = P P^2 = P^3$$

Proceeding further in a similar way, we get

$$P^{(n)} = P^n$$

For example, consider the problem of Random walk with reflecting barriers, discussed above, for which the tpm is

$$\begin{aligned} P &= \begin{pmatrix} 0 & 1 & 0 & 0 \\ 0.3 & 0 & 0.7 & 0 \\ 0 & 0.3 & 0 & 0.7 \\ 0 & 0 & 1 & 0 \end{pmatrix} \\ P^2 &= \begin{pmatrix} 0.3 & 0 & 0.7 & 0 \\ 0 & 0.51 & 0 & 0.49 \\ 0.09 & 0 & 0.91 & 0 \\ 0 & 0.3 & 0 & 0.7 \end{pmatrix} \end{aligned}$$

From this matrix, we see that  $p_{11}^{(2)} = 0.51$ . This is so, because

$$\begin{aligned} p_{11}^{(2)} &= p_{10}p_{01} + p_{11}p_{11} + p_{12}p_{21} + p_{13}p_{31} \\ &= (0.3)(1) + (0)(0) + (0.7)(0.3) + (0)(0) = 0.51 \end{aligned}$$

**Definition:** A stochastic matrix  $P$  is said to be a *regular matrix*, if all the entries of  $P^m$  (for some positive integer  $m$ ) are positive. A homogeneous Markov chain is said to be *regular* if its tpm is regular.

We state below two theorems without proof:

1. If  $p = \{p_i\}$  is the state probability distribution of the process at an arbitrary time, then that after one step is  $pP$ , where  $P$  is the tpm of the chain and that after  $n$  steps in  $pP^n$ .
2. If a homogeneous Markov chain is regular, then every sequence of state probability distributions approaches a unique fixed probability distribution called the *stationary (state) distribution* or *steady-state distribution* of the Markov chain.

That is,  $\lim_{n \rightarrow \infty} \{p^{(n)}\} = \pi$ , where the state probability distribution at step  $n$ ,  $p^{(n)} = (p_1^{(n)}, p_2^{(n)}, \dots, p_k^{(n)})$  and the stationary distribution  $\pi = (\pi_1, \pi_2, \dots, \pi_n)$  are row vectors.

3. Moreover, if  $P$  is the tpm of the regular chain, then  $\pi P = \pi$  ( $\pi$  is a row vector). Using this property of  $\pi$ , it can be found out, as in the worked examples given below:

## Classification of States of a Markov Chain

If  $p_{ij}^{(n)} > 0$  for some  $n$  and for all  $i$  and  $j$ , then every state can be reached from every other state. When this condition is satisfied, the Markov chain is said to be *irreducible*. The tpm of an irreducible chain is an irreducible matrix. Otherwise, the chain is said to be *nonirreducible* or *reducible*.

State  $i$  of a Markov chain is called a *return state*, if  $p_{ij}^{(n)} > 0$  for some  $n > 1$ .

The *period*  $d_i$  of a return state  $i$  is defined as the greatest common divisor of all  $m$  such that  $p_{ii}^{(m)} > 0$ , i.e.,  $d_i = \text{GCD } \{m : p_{ii}^{(m)} > 0\}$ . State  $i$  is said to be *periodic* with period  $d_i$  if  $d_i > 1$  and *aperiodic* if  $d_i = 1$ .

Obviously state  $i$  is aperiodic if  $p_{ii} \neq 0$ . The probability that the chain returns to state  $i$ , having started from state  $i$ , for the first time at the  $n$ th step (or after  $n$  transitions) is denoted by  $f_{ii}^{(n)}$  and called the *first return time probability* or the *recurrence time probability*.  $\{n, f_{ii}^{(n)}\}, n = 1, 2, 3, \dots$ , is the distribution of recurrence times of the state  $i$ .

If  $F_{ii} = \sum_{n=1}^{\infty} f_{ii}^{(n)} = 1$ , the return to state  $i$  is certain.

$\mu_{ii} = \sum_{n=1}^{\infty} n f_{ii}^{(n)}$  is called the *mean recurrence time* of the state  $i$ .

A state  $i$  is said to be *persistent* or *recurrent* if the return to state  $i$  is certain, i.e., if  $F_{ii} = 1$ . The state  $i$  is said to be *transient* if the return to state  $i$  is uncertain,

i.e., if  $F_{ii} < 1$ . The state  $i$  is said to be *nonnull persistent* if its mean recurrence time  $\mu_{ii}$  is finite and *null persistent*, if  $\mu_{ii} = \infty$ .

A non-null persistent and aperiodic state is called *ergodic*.

We give below two theorems, without proof, which will be helpful to classify the states of a Markov chain.

1. If a Markov chain is irreducible, all its states are of the same type. They are all transient, all null persistent or all nonnull persistent. All its states are either aperiodic or periodic with the same period.
2. If a Markov chain finite irreducible, all its states are nonnull persistent.

## Birth and Death Process

Another random process that has wide applications in several fields of natural phenomena such as spread of epidemics, queueing problems, telephone exchange, traffic maintenance and population growth is the birth and death process.

**Definition:** If  $X(t)$  represents the number of individuals present at time  $t$  in a population [or the size of the population at time  $t$ ] in which two types of events occur—one representing birth which contributes to its increase and the other representing death which contributes to its decrease, then the discrete random process  $\{X(t)\}$  is called the *birth and death process*, provided the two events, viz., birth and death are governed by the following postulates:

If  $X(t) = n (n > 0)$ ,

- (i)  $P[1 \text{ birth in } (t, t + \Delta t)] = \lambda_n(t) \Delta t + O(\Delta t)$
- (ii)  $P[0 \text{ birth in } (t, t + \Delta t)] = 1 - \lambda_n(t) \Delta t + O(\Delta t)$
- (iii)  $P[2 \text{ or more births in } (t, t + \Delta t)] = O(\Delta t)$
- (iv) Births occurring in  $(t, t + \Delta t)$  are independent of time since last birth.
- (v)  $P[1 \text{ death in } (t, t + \Delta t)] = \mu_n(t) \Delta t + O(\Delta t)$
- (vi)  $P[0 \text{ death in } (t, t + \Delta t)] = 1 - \mu_n(t) \Delta t + O(\Delta t)$
- (vii)  $P[2 \text{ or more deaths in } (t, t + \Delta t)] = O(\Delta t)$
- (viii) Deaths occurring in  $(t, t + \Delta t)$  are independent of time since last death.
- (ix) The birth and death occur independently of each other at any time.

## Probability Distribution of $X(t)$

Let  $P_n(t) = P\{X(t) = n\}$  = probability that the size of the population is  $n$  at time  $t$ . Then  $P_n(t + \Delta t) = P\{X(t + \Delta t) = n\}$  = probability that the size of the population is  $n$  at time  $(t + \Delta t)$ . Now the event  $X(t + \Delta t) = n$  can happen in any one of the following four mutually exclusive ways:

- (i)  $X(t) = n$  and no birth or death in  $(t, t + \Delta t)$
- (ii)  $X(t) = n - 1$  and 1 birth and no death in  $(t, t + \Delta t)$
- (iii)  $X(t) = n + 1$  and no birth and 1 death in  $(t, t + \Delta t)$
- (iv)  $X(t) = n$  and 1 birth and 1 death in  $(t, t + \Delta t)$

$$\begin{aligned}\therefore P_n(t + \Delta t) &= P(\text{i}) + P(\text{ii}) + P(\text{iii}) + P(\text{iv}) \\ &= P_n(t) (1 - \lambda_n \Delta t) (1 - \mu_n \Delta t) + P_{n-1}(t) \cdot \lambda_n \Delta t\end{aligned}$$

$(1 - \mu_{n-1} \Delta t) + P_{n+1}(t) (1 - \lambda_{n+1} \Delta t) \mu_{n+1} \Delta t + \mu_n(t) \lambda_n \Delta t \cdot \mu_n \Delta t$ , omitting higher powers of  $\Delta t$ .

i.e., 
$$\begin{aligned} P_n(t + \Delta t) &= P_n(t) - (\lambda_n + \mu_n) P_n(t) \Delta t + \lambda_{n-1} \cdot P_{n-1}(t) \cdot \Delta t \\ &\quad + \mu_{n+1} \cdot P_{n+1}(t) \cdot \Delta t, \text{ omitting terms containing } (\Delta t)^2. \end{aligned}$$

$$\therefore \frac{P_n(t + \Delta t) - P_n(t)}{\Delta t} = \lambda_{n-1} \cdot P_{n-1}(t) - (\lambda_n + \mu_n) P_n(t) + \mu_{n+1} P_{n+1}(t) \dots \quad (1)$$

Taking limits on both sides of (1) as  $\Delta t \rightarrow 0$ , we get

$$P'_n(t) = \lambda_{n-1} P_{n-1}(t) - (\lambda_n + \mu_n) P_n(t) + \mu_{n+1} P_{n+1}(t) \quad (2)$$

The difference-differential equation (2) holds good for  $n \geq 1$ . It is not valid when  $n = 0$ , as no death is possible in  $(t, t + \Delta t)$  and  $X(t) = n - 1 = -1$  is meaningless.

$$\therefore P_0(t + \Delta t) = P_0(t) (1 - \lambda_0 \Delta t) + P_1(t) (1 - \lambda_1 \Delta t) \mu_1 \Delta t$$

i.e., 
$$\frac{P_0(t + \Delta t) - P_0(t)}{\Delta t} = -\lambda_0 P_0(t) + \mu_1 P_1(t) \quad (3)$$

Proceeding to limits as  $\Delta t \rightarrow 0$ , we get

$$P'_0(t) = -\lambda_0 P_0(t) + \mu_1 P_1(t) \quad (4)$$

On solving equations (2) and (4), we get  $P_n(t)$  [ $n \geq 0$ ] which gives  $P\{X(t) = n\}$ , the probability distribution of  $X(t)$ .

**Note** The equations (2) and (4) characterize the generalised birth and death process, which assume that the collective birth and death rates in  $(t, t + \Delta t)$  are  $\lambda_n$  and  $\mu_n$  respectively) which depend on the size  $n$  of the population at time  $t$ .

If we further make the simplifying assumptions that the birth rate in  $(t, t + \Delta t)$  is  $\lambda$  for each individual in the population and the death rate in  $(t, t + \Delta t)$  is  $\mu$  for each individual, then  $\lambda_n = n\lambda$  and  $\mu_n = n\mu$ .

In this case, equations (2) and (4) become

$$P'_n(t) = (n - 1) \lambda P_{n-1}(t) - n(\lambda + \mu) P_n(t) + (n + 1)\mu P_{n+1}(t) \quad (5)$$

and  $P'_0(t) = \mu P_1(t) \quad (6)$

Equations (5) and (6) are said to characterize the simple birth and death process or Linear Growth Process.

### Value of $P_n(t)$ for the Simple Birth and Death Process

We define the probability generating function  $G(z, t) = \sum_{n=0}^{\infty} P_n(t) z^n$ .

For the simple birth and death process,  $P_n(t)$  is given by

$$P'_n(t) = (n - 1) \lambda P_{n-1}(t) - n(\lambda + \mu) P_n(t) + (n + 1)\mu P_{n+1}(t) \quad (1)$$

Multiplying both sides of (1) by  $z^n$  and summing over all values of  $n$ , we have

$$\sum_{n=0}^{\infty} P'_n(t) z^n = \lambda \sum_{n=0}^{\infty} (n - 1) P_{n-1}(t) z^n - (\lambda + \mu) \sum_{n=0}^{\infty} n P_n(t) z^n + \mu \sum_{n=0}^{\infty} (n + 1) P_{n+1}(t) z^n$$

$$= \lambda z^2 \sum_{n=0}^{\infty} n P_n(t) z^{n-1} - (\lambda + \mu) z \sum_{n=0}^{\infty} n P_n(t) z^{n-1} + \mu \sum_{n=0}^{\infty} n P_n(t) z^{n-1}$$

i.e.,  $\frac{\partial G}{\partial t} = \lambda z^2 \frac{\partial G}{\partial z} - (\lambda + \mu) z \frac{\partial G}{\partial z} + \mu \frac{\partial G}{\partial z},$

since  $\frac{\partial G}{\partial z} = \sum_{n=0}^{\infty} n P_n(t) z^{n-1}$

i.e.,  $\frac{\partial G}{\partial t} = (\lambda z - \mu)(z - 1) \frac{\partial G}{\partial z} \dots \quad (2)$

This equation (2) is a Lagrange's linear equation. The corresponding subsidiary simultaneous equations are

$$\frac{dt}{-1} = \frac{dz}{(\lambda z - \mu)(z - 1)} = \frac{dG}{0} \quad (3)$$

From the first two ratios, we have

$$\frac{dt}{-1} = \frac{1}{\lambda - \mu} \left( \frac{dz}{z - 1} - \frac{\lambda dz}{\lambda z - \mu} \right), \text{ by partial fractions}$$

$\therefore$  one solution of (3) is

$$\log \left( \frac{z - 1}{\lambda z - \mu} \right) = (\mu - \lambda)t - \log c_1$$

i.e.,  $\frac{z - 1}{\lambda z - \mu} = \frac{1}{c} e^{-(\lambda - \mu)t} \text{ or } \frac{\lambda z - \mu}{z - 1} = c_1 e^{(\lambda - \mu)t} \quad (4)$

The second solution (4) is obviously  $G = c_2$  (5)

$\therefore$  the general solution of (2) is

$$G(z, t) = f \left\{ \left( \frac{\lambda z - \mu}{z - 1} \right) e^{(\mu - \lambda)t} \right\} \quad (6)$$

where  $f$  is an arbitrary function.

Now,  $G(z, 0) = \sum_{n=0}^{\infty} P_n(0) z^n$   
 $= z$ , since  $P_n(0) = P\{X(0) = n\} = 1$ , for  $n = 1$  and  $= 0$ , for all  $n \neq 1$

Using this in (6), we have

$$f \left( \frac{\lambda z - \mu}{\lambda - 1} \right) = z \quad (7)$$

When  $z$  is replaced by  $\phi(z)$ , let the argument of  $f$  reduce to  $z$ .

Then  $\frac{\lambda \phi(z) - \mu}{\phi(z) - 1} = z \quad \therefore \phi(z) = \frac{\mu - z}{\lambda - z}$

But by (7),  $f(z) = \frac{\mu - z}{\lambda - z}$

Using (8) in (6), the required solution is

$$\begin{aligned} G(z, t) &= \frac{\mu - \left( \frac{\lambda z - \mu}{z - 1} \right) e^{(\mu - \lambda)t}}{\lambda - \left( \frac{\lambda z - \mu}{z - 1} \right) e^{(\mu - \lambda)t}} \\ &= \frac{\mu(z-1)(e^{-\lambda-\mu}t) - (\lambda z - \mu)}{\lambda(z-1)e^{(\lambda-\mu)t} - (\lambda z - \mu)} \\ &= \frac{\mu\{1 - e^{(\lambda-\mu)t}\} - \{\lambda - \mu e^{(\lambda-\mu)t}\}z}{\{\mu - \lambda e^{(\lambda-\mu)t}\} - \lambda z\{1 - e^{(\lambda-\mu)t}\}} \end{aligned}$$

Putting  $\alpha(t) = \frac{\mu\{1 - e^{(\lambda-\mu)t}\}}{\mu - \lambda e^{(\lambda-\mu)t}}$

and  $\beta(t) = \frac{\lambda}{\mu} \alpha(t) = \frac{\lambda\{1 - e^{(\lambda-\mu)t}\}}{\mu - \lambda e^{(\lambda-\mu)t}}$ , we get

$$\begin{aligned} G(z, t) &= \frac{\alpha(t) + \{1 - \alpha(t) - \beta(t)\}z}{1 - \beta(t)z} \\ &= [\alpha(t) + \{1 - \alpha(t) - \beta(t)\}z \{1 - \beta(t)z\}^{-1}] \\ &= [\alpha(t) + \{1 - \alpha(t) - \beta(t)\}z \sum_{n=0}^{\infty} \{\beta(t)\}^n z^n] \end{aligned}$$

Now,  $P_n(t)$  = coefficient of  $z^n$  in the expansion of  $G(z, t)$   
 $= \alpha(t) \{\beta(t)\}^n + \{1 - \alpha(t) - \beta(t)\} \{\beta(t)\}^{n-1}$   
 $= [\alpha(t) \beta(t) + 1 - \alpha(t) - \beta(t)] \{\beta(t)\}^{n-1}$   
 $= \{1 - \alpha(t)\} \{1 - \beta(t)\} \{\beta(t)\}^{n-1}; n \geq 1.$

Now,  $P_0(t) = 1 - \sum_{n=1}^{\infty} P_n(t)$   
 $= 1 - \{1 - \alpha(t)\} \{1 - \beta(t)\} \sum_{n=1}^{\infty} \{\beta(t)\}^{n-1}$   
 $= 1 - \{1 - \alpha(t)\} \{1 - \beta(t)\} \{1 - \beta(t)\}^{-1}$   
 $= \alpha(t)$

**Note** Had we assumed that  $X(0) = m$ , instead of  $X(0) = 1$ , we would have got

$$G(z, t) = \left[ \frac{\alpha(t) + \{1 - \alpha(t) - \beta(t)z\}}{1 - \beta(t)z} \right]^m$$

## Mean and Variance of the Population Size in a Linear Birth and Death Process

The probability distribution of the size  $X(t)$  of the population in the linear birth and death process is given by

$$\{n, P_n(t)\}; n = 0, 1, 2, \dots, \infty$$

$$\begin{aligned} \therefore E\{X(t)\} &= \sum_{n=0}^{\infty} n P_n(t) \\ &= (1-\alpha)(1-\beta) \sum_{n=1}^{\infty} n \beta^{n-1}, \text{ where } \alpha = \alpha(t) \text{ and } \beta = \beta(t) \\ &= (1-\alpha)(1-\beta) \{1 + 2\beta + 3\beta^2 + \dots\} \\ &= \frac{1-\alpha}{1-\beta} \\ &= \frac{\mu - \lambda e^{(\lambda-\mu)t} - \mu + \mu e^{(\lambda-\mu)t}}{\mu - \lambda e^{(\lambda-\mu)t} - \lambda + \lambda e^{(\lambda-\mu)t}} \\ &= \frac{(\mu - \lambda) e^{(\lambda-\mu)t}}{\mu - \lambda} = e^{(\lambda-\mu)t} \end{aligned}$$

**Note** When  $\lambda < \mu$  (viz., birth rate is smaller than death rate),  $E\{X(t)\} \rightarrow 0$  as  $t \rightarrow \infty$ .  
 When  $\lambda > \mu$  (viz., birth rate is greater than the death rate),  $E\{X(t)\} \rightarrow \infty$  as  $t \rightarrow \infty$ . Of course, when  $\lambda = \mu$ ,  $E\{X(t)\} = 1 = X(0)$

$$\begin{aligned} \text{Now, } E\{X^2(t)\} &= \sum_{n=0}^{\infty} n^2 P_n(t) \\ &= (1-\alpha)(1-\beta) \left[ \sum_{n=1}^{\infty} \{n(n+1)-n\} \beta^{n-1} \right] \\ &= (1-\alpha)(1-\beta) [\{1.2 + 2.3\beta + 3.4\beta^2 + \dots\} - \{1 + 2\beta + 3\beta^2 + \dots\}] \\ &= (1-\alpha)(1-\beta) \left[ \frac{2}{(1-\beta)^3} - \frac{1}{(1-\beta)^2} \right] \\ &= \frac{(1-\alpha)(1+\beta)}{(1-\beta)^2} \end{aligned}$$

$$\begin{aligned} \text{Var}\{X(t)\} &= E\{X^2(t)\} - E^2\{X(t)\} \\ &= \frac{(1-\alpha)(1+\beta)}{(1-\beta)^2} - \frac{(1-\alpha)^2}{(1-\beta)^2} \end{aligned}$$

$$\begin{aligned}
 &= \frac{(1-\alpha)(\alpha+\beta)}{(1-\beta)^2} \\
 &= \left( \frac{\lambda+\mu}{\lambda-\mu} \right) e^{(\lambda-\mu)t} \{e^{(\lambda-\mu)t} - 1\}
 \end{aligned}$$

## Pure Birth Process

If  $X(t)$  represents the size of a population at time  $t$ , in which only births can take place, then the discrete random process  $\{X(t)\}$  is called *pure birth process*, provided the births are governed by the postulates (i), (ii), (iii) and (iv) of the birth and death process.

The difference-differential equations representing a pure birth process are

$$P'_n(t) = \lambda_{n-1} P_{n-1}(t) - \lambda_n P_n(t) \quad (1)$$

$$\text{and } P'_0(t) = \lambda_0 P_0(t) \quad (2)$$

which are obtained from the corresponding equations representing birth and death process by putting  $\mu_n = 0$ , for  $n \geq 1$ .

If we assume that  $\lambda_n = n\lambda$ , where  $\lambda$  is the birth-rate for all individuals at time  $t$ , the pure birth process is called *Yule-Furry process* or simple birth process.

In this case, equation (2) becomes  $P'_0(t) = 0$  and so  $P_0(t) = \text{constant} = 0$ .

Solution of equation (1) is obtained by putting  $\mu = 0$  in the solution of linear birth and death process.

viz.,  $P_n(t) = e^{-\lambda t} (1 - e^{-\lambda t})^{n-1}; n \geq 1$ .

Also for the simple birth process  $\{X(t)\}$ ,  $E\{X(t)\} = e^{\lambda t}$  and  $\text{Var } \{X(t)\} = e^{\lambda t} \{e^{\lambda t} - 1\}$ .

## Queueing Processes

When the generalised birth and death process is in steady-state, viz., when  $P_n(t)$  and  $P_0(t)$  are independent of time,  $P'_n(t)$  and  $P'_0(t)$  become zero. Hence, the difference-differential equations that characterise the generalised birth and death process reduce to the difference equations.

$$\lambda_{n-1} P_{n-1} - (\lambda_n + \mu_n) P_n + \mu_{n+1} P_{n+1} = 0 \quad (1)$$

$$\text{and } -\lambda_0 P_0 + \mu_1 P_1 = 0 \quad (2)$$

These equations (1) and (2) characterise the Poisson queueing systems, which are discussed in detail in the chapter on "Queueing Theory". In the single server queueing models,  $\lambda_n = \lambda$  and  $\mu_n = \mu$ .

In the multiserver (' $s$ ' servers) queueing models,

$$\lambda_n = \lambda \text{ and } \mu_n = \begin{cases} n\mu, & \text{if } 0 \leq n < s \\ s\mu, & \text{if } n \geq s \end{cases}$$

## Renewal Process

Renewal process is a generalisation of Poisson process. In Poisson process, the time interval between two consecutive occurrences of the event follows an exponential distribution, whereas in renewal process, the inter-arrival times are independent identically distributed continuous Random Variables.

### Definition

If  $N(t)$  represents the number of occurrences of a certain event (the number of renewals of a certain component in a machine) in  $(0, t)$ , then the discrete random process  $\{N(t); t \geq 0\}$  is called a *renewal counting process* or simply *renewal process*, provided the inter-arrival times  $X_1, X_2, X_3, \dots$  are non-negative, independent and identically distributed random variables ( $i \cdot i \cdot d \cdot r \cdot v' \cdot s$ ) with a common distribution function  $F(x)$ .

**Note** If  $X_i$  is assumed to represent the life-time of the components being replaced and the first component is installed at time  $t = 0$ , then it is replaced instantaneously at time  $t = X_1$  (viz., the first renewal has taken place at  $t = X_1$ ). The replaced component is again replaced at time  $t = X_1 + X_2$  and so on. If  $S_n = X_1 + X_2 + \dots + X_n$ , then  $S_n$  represents the time at which the  $n$ th replacement is made.  $N(t)$  is the largest value of  $n$  for which  $S_n \leq t$ .]

## Probability Distribution of the Number of Renewals, $N(t)$ and $E\{N(t)\}$

The distribution of  $N(t)$  is related to that of  $S_n = X_1 + X_2 + \dots + X_n$ , because  $N(t) < n$ , if and only if  $S_n > t$ , as seen from the figure given below.

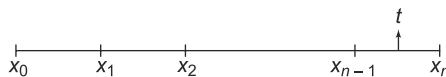


Fig. 7.5

$$\begin{aligned}\therefore P\{N(t) < n\} &= P\{S_n > t\} \\ &= 1 - P\{S_n \leq t\} \\ &= 1 - F_n(t)\end{aligned}$$

where  $F_n(t)$  is the distribution function of  $S_n$  such that  $F_0(t) = 1$

$$\therefore P\{N(t) \geq n\} = F_n(t)$$

$$\text{Now, } P\{N(t) = n\} = P\{N(t) \geq n\} - P\{N(t) \geq n + 1\}$$

$$= F_n(t) - F_{n+1}(t).$$

$$\begin{aligned}
E\{N(t)\} &= \sum_{n=0}^{\infty} nP_n(t), \text{ where } P_n(t) = P\{N(t) = n\} \\
&= \sum_{n=0}^{\infty} n\{F_n(t) - F_{n+1}(t)\} \\
&= \{F_1(t) - F_2(t)\} + 2\{F_2(t) - F_3(t)\} + 3\{F_3(t) - F_4(t)\} + \dots \infty \\
&= F_1(t) + F_2(t) + F_3(t) + \dots \infty \\
&= \sum_{n=1}^{\infty} F_n(t) \text{ or } \sum_{n=1}^{\infty} P\{S_n \leq t\}
\end{aligned}$$

**Note**  $M(t) = E\{N(t)\}$  is called the renewal function of the process  $\{N(t)\}$   
 $M'(t) = m(t)$  is called the renewal density of the process  $\{N(t)\}$

$$\therefore m(t) = \sum_{n=1}^{\infty} F'_n(t) \text{ or } \sum_{n=1}^{\infty} f_n(t), \text{ where } f_n(t) \text{ is the density function of } S_n$$

## Renewal Equation

The integral equation satisfied by the renewal function  $M(t)$ , called the *renewal equation* is given by  $M(t) = F(t) + \int_0^t M(t-x)f(x)dx$ .

**Proof:**

$$\begin{aligned}
M(t) &= E\{N(t)\} \tag{1} \\
&= \int_0^{\infty} E\{N(t)/X_1 = x\} f(x)dx
\end{aligned}$$

where  $f(x)$  is the common pdf of  $X_r$  ( $r = 1, 2, 3, \dots$ )

If  $x > t$  and  $X_1 = x$ , no renewal occurs in  $(0, t)$ , so that  $E\{N(t)\}/X_1 = x\} = 0$ , (2)

If  $0 \leq x \leq t$  and  $X_1 = x$ , one renewal has occurred at time  $x$  and the expected number of renewals in the remaining time interval of length  $t-x$  is  $E\{N(t-x)\}$ .

$$\begin{aligned}
\therefore E\{N(t)\}/X_1 = x\} &= 1 + E\{N(t-x)\} \\
&= 1 + M(t-x)
\end{aligned} \tag{3}$$

Using (2) and (3), we have

$$\begin{aligned}
M(t) &= \int_0^t \{1 + M(t-x)\} f(x)dx \\
&= \int_0^t f(x)dx + \int_0^t M(t-x) f(x)dx
\end{aligned}$$

$$= F(t) + \int_0^t M(t-x) f(x) dx$$

## Poisson Process as a Renewal Process

In the Poisson process, the inter-arrival times  $X_1, X_2, X_3, \dots$  follow identical exponential distributions with pdf  $\lambda e^{-\lambda x}, x \geq 0, \lambda > 0$

$$\therefore S_n = X_1 + X_2 + \dots + X_n$$

follows the Erlang distribution with pdf  $\frac{\lambda^n x^{n-1} e^{-\lambda x}}{(n-1)!}, x \geq 0$ .

$$\therefore F_n(t) = P\{S_n \leq t\}$$

$$= \int_0^t \frac{\lambda^n x^{n-1} e^{-\lambda x}}{(n-1)!} dx \quad (1)$$

$$= 1 - \int_t^\infty \frac{\lambda^n x^{n-1} e^{-\lambda x}}{(n-1)!} dx \quad (\because \text{the integrand is a pdf})$$

$$= 1 - \frac{\lambda^n e^{-\lambda t}}{(n-1)!} \int_0^\infty (y+t)^{n-1} e^{-\lambda y} dy, \text{ on putting } x = y + t$$

$$= 1 - \frac{\lambda^n e^{-\lambda t}}{(n-1)!} \sum_{i=0}^{n-1} (n-1) C_i y^{n-1-i} t^i e^{-\lambda y} dy$$

$$= 1 - \frac{\lambda^n e^{-\lambda t}}{(n-1)!} \sum_{i=0}^{n-1} (n-1) C_i t^i \int_0^\infty y^{n-1-i} e^{-\lambda y} dy$$

$$= 1 - \frac{\lambda^n e^{-\lambda t}}{(n-1)!} \sum_{i=0}^{n-1} \frac{(n-1)! t^i}{i! (n-1-i)!} \frac{(n-1-i)!}{\lambda^{n-i}}$$

$$= 1 - e^{-\lambda t} \sum_{i=0}^{n-1} \frac{(\lambda t)^i}{i!}$$

$$\therefore P\{N(t) = n\} = F_n(t) - F_{n+1}(t)$$

$$= \left[ 1 - e^{-\lambda t} \sum_{i=0}^{n-1} \frac{(\lambda t)^i}{i!} \right] - \left[ 1 - e^{-\lambda t} \sum_{i=0}^n \frac{(\lambda t)^i}{i!} \right]$$

$= e^{-\lambda t} (\lambda t)^n / n!$ , which is the probability law of a Poisson process.

$$E\{N(t)\} = \sum_{n=1}^{\infty} P\{S_n \leq t\}$$

$$= \sum_{n=1}^{\infty} \int_0^t \frac{\lambda^n x^{n-1} e^{-\lambda x}}{(n-1)!} dx, \text{ by (1)}$$

$$\begin{aligned}
 &= \lambda \int_0^t e^{-\lambda x} \left\{ \sum_{n=1}^{\infty} \frac{(\lambda x)^{n-1}}{(n-1)!} \right\} dx \\
 &= \lambda \int_0^t e^{-\lambda x} \cdot e^{\lambda x} dx = \lambda t
 \end{aligned}$$

i.e., renewal function of the Poisson process =  $\lambda t$ .

$\therefore$  Renewal density of the Poisson process =  $\lambda$

**Corollary** If the inter-arrival times  $X_1, X_2, X_3, \dots$  follow identical Erlang distributions with pdf  $\frac{\lambda^k}{(k)} x^{k-1} e^{-\lambda x}; x \geq 0$ , then  $S_n$  will also follow an Erlang

distribution with pdf  $\frac{\lambda^{nk}}{(nk)} x^{nk-1} e^{-\lambda x}; x \geq 0$

(by the reproductive property of Erlang distribution). Proceeding as in the previous case, we can get

$$(i) \quad F_n(t) = 1 - e^{-\lambda x} \sum_{i=0}^{(nk-1)} \frac{(\lambda t)^i}{i!}$$

$$(ii) \quad P\{N(t) = n\} = e^{-\lambda x} \sum_{i=nk}^{(n+1)k-1} \frac{(\lambda t)^i}{i!}; n = 0, 1, 2, \dots$$

For example,

$$P\{N(t) = 0\} = e^{-\lambda x} \left\{ 1 + \frac{\lambda t}{1!} + \frac{(\lambda t)^2}{2!} + \dots + \frac{(\lambda t)^{k-i}}{(k-i)!} \right\}$$

$$\text{and (iii)} \quad E\{N(t)\} = \frac{\lambda t}{k}$$

### Worked Example 7(C)

#### Example 1

The transition probability matrix of a Markov chain  $\{X_n\}, n = 1, 2, 3, \dots$ , having 3 states 1, 2 and 3 is

$$P = \begin{pmatrix} 0.1 & 0.5 & 0.4 \\ 0.6 & 0.2 & 0.2 \\ 0.3 & 0.4 & 0.3 \end{pmatrix}$$

and the initial distribution is  $p^{(0)} = (0.7, 0.2, 0.1)$ .

Find (a)  $P\{X_2 = 3\}$ , and (b)  $P\{X_3 = 2, X_2 = 3, X_1 = 3, X_0 = 2\}$ .

$$\begin{aligned} P^{(2)} = P^2 &= \begin{pmatrix} 0.1 & 0.5 & 0.4 \\ 0.6 & 0.2 & 0.2 \\ 0.3 & 0.4 & 0.3 \end{pmatrix} \begin{pmatrix} 0.1 & 0.5 & 0.4 \\ 0.6 & 0.2 & 0.2 \\ 0.3 & 0.4 & 0.3 \end{pmatrix} \\ &= \begin{pmatrix} 0.43 & 0.31 & 0.26 \\ 0.24 & 0.42 & 0.34 \\ 0.36 & 0.35 & 0.29 \end{pmatrix} \end{aligned}$$

$$\begin{aligned} (a) \quad P\{X_2 = 3\} &= \sum_{i=1}^3 P\{X_2 = 3/X_0 = i\} \times P\{X_0 = i\} \\ &= p_{13}^{(2)} P(X_0 = 1) + p_{23}^{(2)} P(X_0 = 2) + p_{33}^{(2)} P(X_0 = 3) \\ &= 0.26 \times 0.7 + 0.34 \times 0.2 + 0.29 \times 0.1 \\ &= 0.182 + 0.068 + 0.029 \\ &= 0.279 \end{aligned}$$

$$(b) \quad P\{X_1 = 3/X_0 = 2\} = p_{23} = 0.2 \quad (1)$$

$$\begin{aligned} P\{X_1 = 3, X_0 = 2\} &= P\{X_1 = 3/X_0 = 2\} \times P\{X_0 = 2\} \\ &= 0.2 \times 0.2 = 0.04 \text{ [by (1)]} \end{aligned} \quad (2)$$

$$\begin{aligned} P\{X_2 = 3, X_1 = 3, X_0 = 2\} &= P\{X_2 = 3/X_1 = 3, X_0 = 2\} \times P\{X_1 = 3, X_0 = 2\} \\ &= P\{X_2 = 3/X_1 = 3\} \times P\{X_1 = 3, X_0 = 2\} \\ &\quad \text{(by Markov property)} \\ &= 0.3 \times 0.04 \text{ [by (2)]} \\ &= 0.012 \end{aligned} \quad (3)$$

$$\begin{aligned} P\{X_3 = 2, X_2 = 3, X_1 = 3, X_0 = 2\} &= P\{X_3 = 2/X_2 = 3, X_1 = 3, X_0 = 2\} \\ &\quad \times P\{X_2 = 3, X_1 = 3, X_0 = 2\} \\ &= P\{X_3 = 2/X_2 = 3\} \times P\{X_2 = 3, X_1 = 3, X_0 = 2\} \\ &\quad \text{(by Markov property)} \\ &= 0.4 \times 0.012 \text{ [by (3)]} \\ &= 0.0048 \end{aligned}$$

## Example 2

A fair dice is tossed repeatedly. If  $X_n$  denotes the maximum of the numbers occurring in the first  $n$  tosses, find the transition probability matrix  $P$  of the Markov chain  $\{X_n\}$ .

Find also  $P^2$  and  $P(X_2 = 6)$

State space:  $\{1, 2, 3, 4, 5, 6\}$

The tpm is formed using the following analysis.

Let  $X_n$  = the maximum of the numbers occurring in the first  $n$  trials = 3, say

Then  $X_{n+1} = 3$ , if the  $(n+1)$ th trial results in 1, 2, or 3

= 4, if the  $(n+1)$ th trial results in 4

= 5, if the  $(n+1)$ th trial results in 5

= 6, if the  $(n+1)$ th trial results in 6

$$\therefore P\{X_{n+1} = 3 | X_n = 3\} = \frac{1}{6} + \frac{1}{6} + \frac{1}{6} = \frac{3}{6}$$

$$P\{X_{n+1} = i | X_n = 3\} = \frac{1}{6}, \text{ when } i = 4, 5, 6$$

Therefore, the transition probability matrix of the chain is

$$P = \begin{pmatrix} 1/6 & 1/6 & 1/6 & 1/6 & 1/6 & 1/6 \\ 0 & 2/6 & 1/6 & 1/6 & 1/6 & 1/6 \\ 0 & 0 & 3/6 & 1/6 & 1/6 & 1/6 \\ 0 & 0 & 0 & 4/6 & 1/6 & 1/6 \\ 0 & 0 & 0 & 0 & 5/6 & 1/6 \\ 0 & 0 & 0 & 0 & 0 & 1 \end{pmatrix}$$

$$P^2 = \frac{1}{36} \begin{pmatrix} 1 & 3 & 5 & 7 & 9 & 11 \\ 0 & 4 & 5 & 7 & 9 & 11 \\ 0 & 0 & 9 & 7 & 9 & 11 \\ 0 & 0 & 0 & 16 & 9 & 11 \\ 0 & 0 & 0 & 0 & 25 & 11 \\ 0 & 0 & 0 & 0 & 0 & 36 \end{pmatrix}$$

Initial state probability distribution is  $p^{(0)} = \left(\frac{1}{6}, \frac{1}{6}, \frac{1}{6}, \frac{1}{6}, \frac{1}{6}, \frac{1}{6}\right)$  since all the values 1, 2, ..., 6 are equally likely.

$$P\{X_2 = 6\} = \sum_{i=1}^6 P\{X_2 = 6 | X_0 = i\} \times P\{X_0 = i\}$$

$$= \frac{1}{6} \sum_{i=1}^6 p_{i6}^{(2)}$$

$$= \frac{1}{6} \times \frac{1}{36} \times (11 + 11 + 11 + 11 + 11 + 36)$$

$$= \frac{91}{216}$$

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### Example 3

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A man either drives a car or catches a train to go to office each day. He never goes 2 days in a row by train but if he drives one day, then the next day he is just as likely to drive again as he is to travel by train. Now suppose that on the first day of the week, the man tossed a fair die and drove to work if and only if a 6 appeared. Find (a) the probability that he takes a train on the third day, and (b) the probability that he drives to work in the long run.

The travel pattern is a Markov chain, with state space = (train, car)

The tpm of the chain is

$$P = \begin{matrix} T & C \\ \begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix} \end{matrix}$$

The initial state probability distribution is  $p^{(1)} = \left(\frac{5}{6}, \frac{1}{6}\right)$ ,

since  $P(\text{travelling by car}) = P(\text{getting 6 in the toss of the die})$

$$= \frac{1}{6}$$

and  $P(\text{travelling by train}) = \frac{5}{6}$

$$p^{(2)} = p^{(1)}P = \left(\frac{5}{6}, \frac{1}{6}\right) \begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix} = \left(\frac{1}{12}, \frac{11}{12}\right)$$

$$p^{(3)} = p^{(2)}P = \left(\frac{1}{12}, \frac{11}{12}\right) \begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix} = \left(\frac{11}{24}, \frac{13}{24}\right)$$

$$\therefore P(\text{the man travels by train on the third day}) = \frac{11}{24}$$

Let  $\pi = (\pi_1, \pi_2)$  be the limiting form of the state probability distribution or stationary state distribution of the Markov chain.

By the property of  $\pi$ ,  $\pi P = \pi$

$$\text{i.e., } (\pi_1, \pi_2) \begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix} = (\pi_1, \pi_2)$$

$$\text{i.e., } \frac{1}{2}\pi_2 = \pi_1 \quad (1)$$

$$\text{and } \pi_1 + \frac{1}{2}\pi_2 = \pi_2 \quad (2)$$

Equations (1) and (2) are one and the same.

Therefore, consider (1) or (2) with  $\pi_1 + \pi_2 = 1$ , since  $\pi$  is a probability distribution.

Solving,  $\pi_1 = \frac{1}{3}$  and  $\pi_2 = \frac{2}{3}$

$$\therefore P\{\text{the man travels by car in the long run}\} = \frac{2}{3}.$$

### Example 4

Consider a communication system which transmits the 2 digits 0 and 1 through several stages. Let  $X_n (n \geq 1)$  be the digit leaving the  $n$ th stage of the system and  $X_0$  be the digit entering the first stage (or leaving the 0th stage). At each stage there is a constant probability  $q$  that the digit which enters will be transmitted unchanged (i.e., the digit will remain unchanged when it leaves) and the probability  $p$  otherwise (i.e., the digit changes when it leaves), where  $p + q = 1$ . Write down the tpm  $P$  of the homogeneous two-state Markov chain  $\{X_n\}$ . Find  $P^m$ ,  $P^\infty$  and the conditional probability that the digit entering the first stage is 0, given that the digit leaving the  $m$ th stage is 0. Assume that the initial state probability distribution is  $p^{(0)} = (a, 1 - a)$ .

$$\begin{array}{c} \text{State of } X_{n+1} \\ \begin{matrix} 0 & 1 \end{matrix} \\ \text{State space} = (0, 1); P \equiv \text{State of } X_n \begin{pmatrix} q & p \\ p & q \end{pmatrix} \end{array}$$

$$\begin{aligned} \text{Now, } P^2 &= \begin{pmatrix} q & p \\ p & q \end{pmatrix} \begin{pmatrix} q & p \\ p & q \end{pmatrix} \\ &= \begin{pmatrix} p^2 + q^2 & 2pq \\ 2pq & p^2 + q^2 \end{pmatrix} \\ &= \begin{pmatrix} \frac{1}{2}[(q+p)^2 + (q-p)^2] & \frac{1}{2}[(q+p)^2 - (q-p)^2] \\ \frac{1}{2}[(q+p)^2 - (q-p)^2] & \frac{1}{2}[(q+p)^2 + (q-p)^2] \end{pmatrix} \\ &= \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^2 & \frac{1}{2} - \frac{1}{2}r^2 \\ \frac{1}{2} - \frac{1}{2}r^2 & \frac{1}{2} + \frac{1}{2}r^2 \end{pmatrix}, \text{ where } q - p = r \\ P^3 &= \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^3 & \frac{1}{2} - \frac{1}{2}r^3 \\ \frac{1}{2} - \frac{1}{2}r^3 & \frac{1}{2} + \frac{1}{2}r^3 \end{pmatrix} \end{aligned}$$

The values of  $P^2$  and  $P^3$  make us guess that

$$P^m = \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^m & \frac{1}{2} - \frac{1}{2}r^m \\ \frac{1}{2} - \frac{1}{2}r^m & \frac{1}{2} + \frac{1}{2}r^m \end{pmatrix}$$

It is correct as can be proved by induction as follows:

$$\begin{aligned} P^{m+1} &= \begin{pmatrix} q & p \\ p & q \end{pmatrix} \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^m & \frac{1}{2} - \frac{1}{2}r^m \\ \frac{1}{2} - \frac{1}{2}r^m & \frac{1}{2} + \frac{1}{2}r^m \end{pmatrix} \\ &= \begin{pmatrix} \frac{q}{2} + \frac{q}{2}r^m + \frac{p}{2} - \frac{p}{2}r^m & \frac{q}{2} - \frac{q}{2}r^m + \frac{p}{2} + \frac{p}{2}r^m \\ \frac{p}{2} + \frac{p}{2}r^m + \frac{q}{2} - \frac{q}{2}r^m & \frac{p}{2} - \frac{p}{2}r^m + \frac{q}{2} + \frac{q}{2}r^m \end{pmatrix} \\ &= \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^{m+1} & \frac{1}{2} - \frac{1}{2}r^{m+1} \\ \frac{1}{2} - \frac{1}{2}r^{m+1} & \frac{1}{2} + \frac{1}{2}r^{m+1} \end{pmatrix} \\ \therefore P^m &= \begin{pmatrix} \frac{1}{2} + \frac{1}{2}r^m & \frac{1}{2} - \frac{1}{2}r^m \\ \frac{1}{2} - \frac{1}{2}r^m & \frac{1}{2} + \frac{1}{2}r^m \end{pmatrix}, \text{ where } m \text{ is a positive integer } \geq 1 \\ p^\infty &= \lim_{m \rightarrow \infty} (P^m) = \begin{pmatrix} \frac{1}{2} & \frac{1}{2} \\ \frac{1}{2} & \frac{1}{2} \end{pmatrix} \text{ since } |r| < 1 \end{aligned}$$

$$\begin{aligned} \text{Now, } P\{X_m = 0, X_0 = 0\} &= P\{X_m = 0/X_0 = 0\} \times P\{X_0 = 0\} \\ &= ap_{00}^{(m)} \end{aligned}$$

$$\text{and } P\{X_m = 0, X_0 = 1\} = bp_{10}^{(m)} \text{ where } b = 1 - a$$

$$\begin{aligned} \text{Now, } P\{X_0 = 0/X_m = 0\} &= \frac{P\{X_0 = 0\} \times P\{X_m = 0/X_0 = 0\}}{P\{X_0 = 0\} \times p_{00}^{(m)} + P\{X_0 = 1\} \times P_{10}^{(m)}} \\ &\quad \text{(by Bayes' theorem)} \end{aligned}$$

$$= \frac{a \left\{ \frac{1}{2} + \frac{1}{2}r^m \right\}}{a \left\{ \frac{1}{2} + \frac{1}{2}r^m \right\} + b \left\{ \frac{1}{2} - \frac{1}{2}r^m \right\}}$$

$$= \frac{a(1+r^m)}{1+(a-b)r^m}, \text{ where } b = 1-a$$

### Example 5

A gambler has Rs. 2/-. He bets Re. 1 at a time and wins Re. 1 with probability 1/2. He stops playing if he loses Rs. 2 or wins Rs. 4 (a) What is the tpm of the related Markov chain? (b) What is the probability that he has lost his money at the end of 5 plays? (c) What is the probability that the game lasts more than 7 plays?

Let  $X_n$  represent the amount with the player at the end of the  $n^{\text{th}}$  round of the play.

State space of  $\{X_n\} = \{0, 1, 2, 3, 4, 5, 6\}$ , as the game ends, if the player loses all the money ( $X_n = 0$ ) or wins Rs. 4, i.e., has Rs. 6 ( $X_n = 6$ ). The tpm of the Markov chain is

$$P = \begin{pmatrix} & 0 & 1 & 2 & 3 & 4 & 5 & 6 \\ 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\ 1 & \frac{1}{2} & 0 & \frac{1}{2} & 0 & 0 & 0 & 0 \\ 2 & 0 & \frac{1}{2} & 0 & \frac{1}{2} & 0 & 0 & 0 \\ 3 & 0 & 0 & \frac{1}{2} & 0 & \frac{1}{2} & 0 & 0 \\ 4 & 0 & 0 & 0 & \frac{1}{2} & 0 & \frac{1}{2} & 0 \\ 5 & 0 & 0 & 0 & 0 & \frac{1}{2} & 0 & \frac{1}{2} \\ 6 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \end{pmatrix}$$

**Note** This is called a random walk with absorbing barriers at 0 and 6, since the chain cannot come out of the states 0 and 6, once it has entered them.

The initial probability distribution of  $\{X_n\}$  is  $p(0) = (0, 0, 1, 0, 0, 0, 0)$ , as the player has got Rs. 2/- to start with.

$$p^{(1)} = p^{(0)}P = (0, 1/2, 0, 1/2, 0, 0, 0)$$

$$p^{(2)} = p^{(1)}P = (1/4, 0, 1/2, 0, 1/4, 0, 0)$$

$$p^{(3)} = p^{(2)}P = (1/4, 1/4, 0, 3/8, 0, 1/8, 0)$$

$$p^{(4)} = p^{(3)}P = (3/8, 0, 5/16, 0, 1/4, 0, 1/16)$$

$$p^{(5)} = p^{(4)}P = (3/8, 5/32, 0, 9/32, 0, 1/8, 1/16)$$

$P\{\text{the man has lost his money at the end of 5 plays}\}$

$$\begin{aligned} &= P\{X_5 = 0\} = \text{the entry corresponding to state 0 in } p^{(5)} \\ &= 3/8 \end{aligned}$$

$$\text{Again, } p^{(6)} = p^{(5)}P = \left(\frac{29}{64}, 0, \frac{7}{32}, 0, \frac{13}{64}, 0, \frac{1}{8}\right)$$

$$p^{(7)} = p^{(6)} P = \left( \frac{29}{64}, \frac{7}{64}, 0, \frac{27}{128}, 0, \frac{13}{128}, \frac{1}{8} \right)$$

$P\{\text{the game lasts more than 7 rounds}\} = P\{\text{the system is neither in state 0 nor in 6 at the end of the seventh round}\}$

$$\begin{aligned} &= P\{X_7 = 1, 2, 3, 4 \text{ or } 5\} \\ &= \frac{7}{64} + 0 + \frac{27}{128} + 0 + \frac{13}{128} = \frac{27}{64} \end{aligned}$$

### Example 6

There are 2 white marbles in urn A and 3 red marbles in urn B. At each step of the process, a marble is selected from each urn and the 2 marbles selected are interchanged. Let the state  $a_i$  of the system be the number of red marbles in A after  $i$  changes. What is the probability that there are 2 red marbles in A after 3 steps? In the long run, what is the probability that there are 2 red marbles in urn A?

State space of the chain  $\{X_n\} = (0, 1, 2)$ , since the number of balls in the urn A is always 2.

Let the tpm of the chain  $\{X_n\}$  be

$$P = \begin{matrix} & \begin{matrix} 0 & 1 & 2 \end{matrix} \\ \begin{matrix} 0 \\ 1 \\ 2 \end{matrix} & \begin{pmatrix} p_{00} & p_{01} & p_{02} \\ p_{10} & p_{11} & p_{12} \\ p_{20} & p_{21} & p_{22} \end{pmatrix} \end{matrix}$$

$p_{00} = 0$  (since the state cannot remain at 0 after interchange of marbles)

$p_{02} = p_{20} = 0$  (since the number of red marbles in urn cannot increase or decrease by 2 in one interchange)

To start with, A contains 0 red marble. After an interchange, A will contain 1 red marble (and 1 white marble) certainly.

$$\therefore p_{01} = 1$$

Let  $X_n = 1$ , i.e., A contains 1 red marble (and 1 white marble) and B contains 1 white and 2 red marbles.

Then  $X_{n+1} = 0$ , if A contains 0 red marble (and 2 white marbles) and B contains 3 red marbles, i.e., if 1 red marble is chosen from A and 1 white marble is chosen from B and interchanged.

$$\therefore P\{X_{n+1} = 0/X_n = 1\} = p_{10} = \frac{1}{2} \times \frac{1}{3} = \frac{1}{6}$$

$$\text{Similarly, we can find } p_{12} = \frac{1}{2} \times \frac{2}{3} = \frac{1}{3}$$

Since  $P$  is a stochastic matrix,

$$p_{10} + p_{11} + p_{12} = 1$$

$$\therefore p_{11} = \frac{1}{2}$$

Similarly,  $p_{21} = \frac{2}{3}$  and  $p_{22} = 1 - (p_{20} + p_{21}) = \frac{1}{3}$

$$\therefore P = \begin{pmatrix} 0 & 1 & 0 \\ 1/6 & 1/2 & 1/3 \\ 0 & 2/3 & 1/3 \end{pmatrix}$$

Now,  $p^{(0)} = (1, 0, 0)$ , as there is no red marble in A in the beginning.

$$p^{(1)} = p^{(0)} P = (0, 1, 0)$$

$$p^{(2)} = p^{(1)} P = \left( \frac{1}{6}, \frac{1}{2}, \frac{1}{3} \right)$$

$$p^{(3)} = p^{(1)} P = \left( \frac{1}{12}, \frac{23}{26}, \frac{5}{18} \right)$$

$\therefore P\{\text{there are 2 red marbles in } A \text{ after 3 steps}\}$

$$= P\{X_3 = 2\} = p_2^{(3)} = \frac{5}{18}$$

Let the stationary probability distribution of the chain be  $\pi = (\pi_0, \pi_1, \pi_2)$ .

By the property of  $\pi$ ,  $\pi P = \pi$  and  $\pi_0 + \pi_1 + \pi_2 = 1$

$$\text{i.e., } \pi_0, \pi_1, \pi_2 \begin{pmatrix} 0 & 1 & 0 \\ 1/6 & 1/2 & 1/3 \\ 0 & 2/3 & 1/3 \end{pmatrix} = (\pi_0, \pi_1, \pi_2)$$

$$\text{i.e., } \frac{1}{6}\pi_1 = \pi_0$$

$$\pi_0 + \frac{1}{2}\pi_1 + \frac{2}{3}\pi_2 = \pi_1$$

$$\frac{1}{3}\pi_1 + \frac{1}{3}\pi_2 = \pi_2$$

and

$$\pi_1 + \pi_2 + \pi_3 = 1$$

$$\text{Solving, } \pi_0 = \frac{1}{10}, \pi_1 = \frac{6}{10}, \pi_2 = \frac{3}{10}$$

$\therefore P\{\text{there are 2 red marbles in } A \text{ in the long run}\} = 0.3$

### Example 7

Find the nature of the states of the Markov chain with the tpm

$$P = \begin{pmatrix} 0 & 1 & 2 \\ 0 & 1 & 0 \\ 1/2 & 0 & 1/2 \\ 2 & 0 & 1 & 0 \end{pmatrix}$$

$$P^2 = \begin{pmatrix} 1/2 & 0 & 1/2 \\ 0 & 1 & 0 \\ 1/2 & 0 & 1/2 \end{pmatrix}; P^3 = P$$

$$\therefore P^4 = P^2$$

and so on. In general,  $P^{2n} = P^2, P^{2n+1} = P$

We note that  $p_{00}^{(2)} > 0, p_{01}^{(1)} > 0, p_{02}^{(2)} > 0$

$p_{10}^{(1)} > 0, p_{11}^{(2)} > 0, p_{12}^{(1)} > 0$

$p_{20}^{(2)} > 0, p_{21}^{(1)} > 0, p_{22}^{(2)} > 0$

Therefore, the Markov chain is irreducible.

Also  $p_{ii}^{(2)} = p_{ii}^{(4)} = p_{ii}^{(6)} \dots > 0$ , for all  $i$ , all the states of the chain are periodic, with period 2.

Since the chain is finite and irreducible, all its states are nonnull persistent. All states are not ergodic.

### Example 8

Three boys A, B and C are throwing a ball to each other. A always throws the ball to B and B always throws the ball to C, but C is just as likely to throw the ball to B as to A. Show that the process is Markovian. Find the transition matrix and classify the states.

The transition probability matrix of the process  $\{X_n\}$  is given below:

$$\begin{array}{c} \text{State of } X_n \\ \begin{array}{ccc} A & B & C \end{array} \\ \begin{array}{c} \text{State of } X_{n-1} \\ A \\ B \\ C \end{array} \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1/2 & 1/2 & 0 \end{pmatrix} \equiv P, \text{ say} \end{array}$$

States of  $X_n$  depend only on states of  $X_{n-1}$ , but not on states of  $X_{n-2}, X_{n-3}, \dots$ , or earlier states. Therefore,  $\{X_n\}$  is a Markov chain.

$$\text{Now } P^2 = \begin{pmatrix} 0 & 0 & 1 \\ 1/2 & 1/2 & 0 \\ 0 & 1/2 & 1/2 \end{pmatrix}; P^3 = \begin{pmatrix} 1/2 & 1/2 & 0 \\ 0 & 1/2 & 1/2 \\ 1/4 & 1/4 & 1/2 \end{pmatrix}$$

$p_{11}^{(3)} > 0, p_{13}^{(2)} > 0, p_{21}^{(2)} > 0, p_{22}^{(2)} > 0, p_{33}^{(2)} > 0$  and all other  $p_{ij}^{(1)} > 0$ .

Therefore, the chain is irreducible.

$$P^4 = \begin{pmatrix} 0 & 1/2 & 1/2 \\ 1/4 & 1/4 & 1/2 \\ 1/4 & 1/2 & 1/4 \end{pmatrix}; P^5 = \begin{pmatrix} 1/4 & 1/4 & 1/2 \\ 1/4 & 1/2 & 1/4 \\ 1/8 & 3/8 & 1/2 \end{pmatrix}; P^6 = \begin{pmatrix} 1/4 & 1/2 & 1/4 \\ 1/4 & 3/8 & 1/2 \\ 1/8 & 3/8 & 3/8 \end{pmatrix}$$

and so on.

We note that  $p_{ii}^{(2)}, p_{ii}^{(3)}, p_{ii}^{(5)}; p_{ii}^{(6)}$  etc. are  $> 0$  for  $i = 2, 3$ , and GCD of 2, 3, 5, 6, ... = 1.

Therefore, the states 2 and 3 (i.e., B and C) are periodic with period 1, i.e., aperiodic.

We note that  $p_{11}^{(3)}, p_{11}^{(5)}, p_{11}^{(6)}$  etc. are  $> 0$  and GCD of 3, 5, 6, ... = 1

Therefore, the state 1 (i.e., state A) is periodic with period 1, i.e., aperiodic.

Since the chain is finite and irreducible, all its states and nonnull persistent. Moreover all the states are ergodic.

### Exercise 7(C)

#### **Part-A (Short-answer Questions)**

1. Define a Markov process.
2. Define a Markov chain and give an example of a Markov chain.
3. Prove that the Poisson process is a Markov process.
4. When is a Markov chain called homogeneous?
5. When is a homogeneous Markov chain said to be regular?
6. Define transition probability matrix of a Markov chain.
7. What is a stochastic matrix? When is it said to be regular?
8. Prove that the tpm of a Markov chain is a stochastic matrix.
9. Define  $n$ -step transition probability in a Markov chain.
10. State the Chapman-Kolmogorov theorem.
11. What do you mean by probability distribution of a Markov chain?
12. When is a Markov chain completely specified?
13. What is meant by steady-state distribution of a Markov chain?
14. Write down the relation satisfied by the steady-state distribution and the tpm of a regular Markov chain.
15. If the tpm of a Markov chain is  $\begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix}$  find the steady-state distribution of the chain.
16. When is a Markov chain said to be irreducible or ergodic?
17. Prove that the matrix  $\begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1/2 & 1/2 & 0 \end{pmatrix}$  is the tpm of an irreducible Markov chain.
18. What do you mean by an absorbing Markov chain. Give an example.
19. If the initial state probability distribution of a Markov chain is  $p^{(0)} = \left(\frac{5}{6}, \frac{1}{6}\right)$  and the tpm of the chain is  $\begin{pmatrix} 0 & 1 \\ 1/2 & 1/2 \end{pmatrix}$ , find the probability distribution of the chain after 2 steps.

**Part-B**

20. The tpm of a Markov chain with three states 0, 1, 2 is

$$P = \begin{pmatrix} 3/4 & 1/4 & 0 \\ 1/4 & 1/2 & 1/4 \\ 0 & 3/4 & 1/4 \end{pmatrix}$$

and the initial state distribution of the chain is  $P\{X_0 = i\} = \frac{1}{3}, i = 0, 1, 2$ .

- Find (i)  $P\{X_2 = 2\}$ , and (ii)  $P\{X_3 = 1, X_2 = 2, X_1 = 1, X_0 = 2\}$ .
21. A man is at an integral point on the  $x$ -axis between the origin and the point 3. He takes a unit step to the right with probability  $\frac{1}{3}$  or to the left with probability  $2/3$ , unless he is at the origin, where he takes a step to the right to reach the point 1 or is at the point 3, where he takes a step to the left to reach point 2. What is the probability that (i) he is at the point 1 after 3 walks? and (ii) he is at the point 1 in the long run?
22. Suppose that the probability of a dry day (state 0) following a rainy day (state 1) is  $\frac{1}{3}$  and that the probability of a rainy day following a dry day is  $\frac{1}{2}$ . Given that May 1 is a dry day, find the probability that (i) May 3 is also a dry day, and (ii) May 5 is also a dry day.
23. A gambler has Rs. 3/- . At each play of the game, he loses Re. 1 with probability  $\frac{3}{4}$ , but wins Rs. 2/- with probability  $\frac{1}{4}$ . He stops playing if he has lost his initial amount of Rs. 3/- or he has won at least Rs. 3/- . Write down the tpm of the associated Markov chain. Find the probability that there are at least 4 rounds to the game.
24. A communication source can generate 1 of 3 possible messages 1, 2 and 3. Assume that the generation can be described by a homogeneous Markov chain with the following transition probability matrix

<i>Current message</i>	<i>Next message</i>		
	1	2	3
1	0.5	0.3	0.2
2	0.4	0.2	0.4
3	0.3	0.3	0.4

and the initial state probability distribution  $p^{(0)} = (0.3, 0.3, 0.4)$ . Find  $p^{(3)}$ .

25. Assume that the weather in a certain locality can be modeled as the homogeneous Markov chain whose transition probability matrix is given below.

Today's weather	Tomorrow's weather		
	Fair	Cloudy	Rainy
Fair	0.8	0.15	0.05
Cloudy	0.5	0.3	0.2
Rainy	0.6	0.3	0.1

If the initial state distribution is given by  $p^{(0)} = (0.7, 0.2, 0.1)$ , find  $p^{(2)}$  and  $\lim_{n \rightarrow \infty} p^{(n)}$ .

26. A fair coin is tossed until 3 heads occur in a row. Let  $X_n$  be the sequence of heads ending at the  $n$ th trial. What is the probability that there are at least 8 tosses of the coin?
27. There are 2 white marbles in urn  $A$  and 4 red marbles in urn  $B$ . At each step of the process, a marble is selected from each urn and the 2 marbles selected are interchanged. The state of the related Markov chain is the number of red balls in  $A$  after the interchange. What is the probability that there are 2 red balls in urn  $A$  (i) after 3 steps, and (ii) in the long run?
28. A student's study habits are as follows: If he studies one night, he is 70% sure not to study the next night. On the other hand, if he does not study one night, he is 60% sure not to study the next nights as well. In the long run, how often does he study?
29. A salesman's territory consists of 3 cities  $A$ ,  $B$  and  $C$ . He never sells in the same city on successive days. If he sells in city  $A$ , then the next day he sells in  $B$ . However, if he sells either in  $B$  or  $C$ , then the next day he is twice as likely to sell in city  $A$  as in the other city. How often does he sell in each of the cities in the steady state?
30. A housewife buys 3 kinds of cereals,  $A$ ,  $B$  and  $C$ . She never buys the same cereal in successive weeks. If she buys cereal  $A$ , the next week she buys cereal  $B$ . However, if she buys  $B$  or  $C$ , the next week she is 3 times as likely to buy  $A$  as the other cereal. In the long run, how often she buy each of the three cereals?
31. Two boys  $B_1$ ,  $B_2$  and two girls  $G_1$ ,  $G_2$  are throwing a ball from one to another. Each boy throws the ball to the other boy with probability 1/2 and to each girl with probability 1/4. On the other hand, each girl throws the ball to each boy with probability 1/2 and never to the other girl. In the long run, how often does each receive the ball?
32. A gambler's luck follows a pattern. If he wins a game, the probability of his winning the next game is 0.6. However if he loses a game, the probability of his losing the next game is 0.7. There is an even chance that the gambler wins the first game. What is the probability that he wins (i) the second game, (ii) the third game, and (iii) in the long run?

33. The three-state Markov chain is given by the tpm

$$P = \begin{pmatrix} 0 & 2/3 & 1/3 \\ 1/2 & 0 & 1/2 \\ 1/2 & 1/2 & 0 \end{pmatrix}$$

Prove that the chain is irreducible and all the states are aperiodic and non-null persistent. Find also the steady-state distribution of the chain.

34. A man tosses a fair coin until 3 heads occur in a row. Let  $X_n = k$ , if at the  $n^{\text{th}}$  trial, the last tail occurred at the  $(n - k)^{\text{th}}$  trial; i.e.,  $X_n$  denotes the longest string of heads ending at the  $n^{\text{th}}$  trial. Show that the process is Markovian. Find the transition matrix and classify the states.

## ANSWERS

### Exercise 7(A)

3. DC meter measures the mean of the input process;  
 $\therefore$  output of the dc meter = 0 V  
True rms metre measure the SD of the input process.  
 $\therefore$  output of the true rms meter = 2 V
5. Gaussian process is used to model and analyse the effects of thermal noise in electronic circuits used in communication system.
9. Square law detector process, Full-wave linear detector process, Half-wave linear detector process and Hard limiter process.
18.  $X(t) = R_x(t) \cos \{\omega_0 \pm \theta_x(t)\}$
19. In communication systems, information bearing signals are often narrow-band Gaussian processes. When such signals are viewed on an oscilloscope, they appear like a sine wave with slowly varying amplitude and phase. Hence, the representation.
21.  $R_x(t)$  (envelope) follows a Rayleigh distribution and  $\theta_x(t)$  (phase) follows a uniform distribution in  $(0, 2\pi)$ .
23. They are low pass processes.
24. A zero mean WSS process  $\{X(t)\}$  can be represented in the form  $X(t) = I(t) \cos \omega_0 t - Q(t) \sin \omega_0 t$ . This kind of representation is called the quadrature representation.
25. WSS processes with zero mean can be represented in the quadrature form.
26. The quadrature representation is useful in communication theory, only when  $\{X(t)\}$  is a zero mean, WSS, band pass process.
30. The graph of  $S_{NN}(\omega)$  is a straight line parallel to the  $\omega$ -axis.
31.  $S_{NN}(\omega)$  is a constant for all values of  $\omega$ , i.e.,  $S_{NN}(\omega)$  contains all frequencies in equal amount. White noise is called so in analogy to white light which consists of all colours.

32.  $S_{YY}(\omega) = \frac{N_0}{2} |H(\omega)|^2$ , where  $\{Y(t)\}$  is the output process and  $H(\omega)$  is the power transfer function.

$$33. F\left\{\frac{N_0}{2}\delta(\tau)\right\} = \int_{-\infty}^{\infty} \frac{N_0}{2}\delta(\tau)e^{-i\omega t} d\tau = \frac{N_0}{2}$$

$$\therefore R_{NN}(\tau) = F^{-1}\frac{N_0}{2} = \frac{N_0}{2}\delta(\tau)$$

$$34. E\{N^2(t)\} = R_{NN}(0) = \int_{-\infty}^{\infty} S_{NN}(\omega) d\omega = \int_{-\infty}^{\infty} \frac{N_0}{2} d\omega \rightarrow \infty$$

35. Since  $\int_{-\infty}^{\infty} S_{NN}(\omega) d\omega \rightarrow \infty$ , it is not physically realisable. If the frequency band is taken to be finite, say  $(-\omega_B, \omega_B)$ , then  $\int_{-\omega_B}^{\omega_B} S_{NN}(\omega) d\omega = N_0 \omega_B < \infty$ . Thus, the spectral density becomes realisable.

$$38. S_{NN}(\omega) = \begin{cases} \frac{N_0}{2}, & \text{in } |\omega| \leq \omega_B \\ 0, & \text{elsewhere} \end{cases}$$

$$R_{NN}(\tau) = \frac{1}{2\pi} \int_{-\omega_B}^{\omega_B} \frac{N_0}{2} e^{i\tau\omega} d\omega = \frac{N_0}{2\pi} \int_0^{\omega_B} \cos \tau\omega d\omega = \frac{N_0}{2\pi} \frac{\sin \tau\omega_B}{\tau}$$

$$39. E\{N^2(t)\} = R_{NN}(0) = \frac{N_0 \omega_B}{2\pi} \lim_{\tau \rightarrow 0} \left( \frac{\sin \tau\omega_B}{\tau\omega_B} \right) = \frac{N_0 \omega_B}{2\pi}$$

42. Filtering in a stable, linear, time-invariant system is done by selecting carefully the power transfer function  $H(\omega)$  so that certain undesired spectral components of the input signal are removed or filtered out.

43. (i) 0.309; (ii) 0.4

$$46. 4e^{-4}; 8(1+e^{-4}); \frac{1}{2\sqrt{2\pi}} e^{-z^2/8} \quad -\infty < z < \infty; 0.6915;$$

$$\frac{1}{8\pi\sqrt{1-e^{-8}}} \exp\left[-\frac{1}{8(1-e^{-8})}(z^2 - 2e^{-4}zw + w^2)\right] \quad -\infty < z, w < \infty$$

48. 0, 2, 4 volts

$$49. \frac{\{R_{xx}(\tau) - \mu_x^2\}}{R_x(0) - \mu_x^2} X(t_1), \text{ where } \tau = |t_1 - t_2|$$

50. (i)  $Z(t)$  is a Gaussian RV for all  $t$  and hence  $\{Z(t)\}$  is a Gaussian process.

$$(ii) f(z_1, z_2) = \frac{1}{2\pi\sqrt{1-r^2}} \exp\left\{-\frac{1}{2(1-r^2)} z_1^2 - 2rz_1z_2 + z_2^2\right\}$$

$$\text{where } r = \cos \omega(t_1 - t_2)$$

- (iii) Yes      (iv) Yes      (v)  $R_{zz}(\tau) \cdot z(t_1)$  where  $\tau = |t_1 - t_2|$
55.  $\frac{N_0}{2} \delta(\tau)$
57. (i)  $\frac{N_0 \omega_B}{2\pi} \left( \frac{\sin \omega_B \tau}{\omega_B \tau} \right)$     (ii)  $\frac{N_0 \omega_B}{2\pi} \left( \frac{\sin \omega_B \tau}{\omega_B \tau} \right) \cos \omega_0 \tau$
58. 2.3% of the time
60.  $\check{X}(t) = A \sin \omega_0 t - B \cos \omega_0 t$
62.  $\frac{\alpha N_0}{4}; \frac{\alpha N_0}{4} e^{\alpha \tau}, \tau \geq 0$

### Exercise 7(B)

- An ensemble of discrete sets of points from the time domain is called point process, e.g. the times at which components fail in a large system; the times at which phone calls arrive at an exchange.
- $P\{X(t) = k\} = \frac{e^{-\lambda t} (\lambda t)^k}{k!}$ , where  $\lambda = 2$ ;  $P\{X(1) = 0\} = e^{-2}$
- $e^{-\lambda t} (1 - e^{i\omega})$
- 0.1755
- 0.0504      (ii) 0.0054      (iii) 0.7769      (iv) 0.5941
- 0.1144
- (i)  $e^{-12}$     (ii)  $1 - e^{-3(t-12)}$ ,  $t \geq 12$
- 0.9972
- 0.9707
- 6 minutes
- $f_T(t) = \frac{\lambda^k \cdot t^{k-1} e^{-\lambda t}}{|k-1|}, t > 0$
- (i) 30    (ii) 60
- (i) 30    (ii) 30

### Exercise 7(C)

- A square matrix, in which the sum of all the elements of each row is 1, is called a stochastic matrix. A stochastic matrix  $P$  is said to be regular if all the entries of  $P^m$  (for some positive integer  $m$ ) are positive.
- A Markov chain is completely specified when the initial probability distribution and the tpm are given.
- If  $\pi = (\pi_1, \pi_2, \dots, \pi_n)$  is the steady-state distribution of the chain whose tpm is the  $n$ th order square matrix  $P$ , then  $\pi P = \pi$ .

15. If  $(\pi_1, \pi_2)$  is the steady-state distribution of the chain,  $(\pi_1, \pi_2)$ .  $\begin{pmatrix} 0 & 1 \\ 1 & 1 \\ 2 & 2 \end{pmatrix}$

$$= (\pi_1, \pi_2); \frac{1}{2}\pi_2 = \pi_1 \text{ and } \pi_1 + \frac{1}{2}\pi_2 = \pi_2, \text{ i.e., } 2\pi_1 = \pi_2. \text{ Also } \pi_1 + \pi_2 = 1$$

$$\therefore \pi_1 = \frac{1}{3}, \pi_2 = \frac{2}{3}.$$

18. A state  $i$  of a Markov chain is said to be an absorbing state if  $p_{ii} = 1$ , i.e., if it is impossible to leave it.  
A Markov chain is said to be absorbing if it has at least one absorbing state.

**Example:** 
$$\begin{matrix} & 1 & 2 & 3 & 4 \\ 1 & \left( \begin{array}{cccc} 1 & 0 & 0 & 0 \\ 1/2 & 0 & 1/2 & 0 \\ 0 & 1/2 & 0 & 1/2 \\ 0 & 0 & 0 & 1 \end{array} \right) \\ 2 & \\ 3 & \\ 4 & \end{matrix}$$

19.  $p^{(1)} = p^{(0)}P = P^{(2)} = p^{(1)}P = \left( \frac{11}{14}, \frac{13}{24} \right).$

20.  $1/6; 3/64$

21.  $P = \begin{pmatrix} 0 & 1 & 0 & 0 \\ 2/3 & 0 & 1/3 & 0 \\ 0 & 2/3 & 0 & 1/3 \\ 0 & 0 & 1 & 0 \end{pmatrix}$

(i)  $\frac{22}{27}$  (ii)  $\frac{3}{7}$

22. (i)  $\frac{5}{12}$  (ii)  $\frac{173}{432}$

23.  $\frac{27}{64}$

24. ( $\bullet$  4083,  $\bullet$  2727,  $\bullet$  3190)

25. ( $\bullet$  7245,  $\bullet$  1920,  $\bullet$  0835);  $\left( \frac{114}{157}, \frac{30}{157}, \frac{13}{157} \right)$

26.  $81/128$

27.  $3/8; 2/5$

28. 4/11 of the nights

29.  $P = \begin{pmatrix} 0 & 1 & 0 \\ 2/3 & 0 & 1/3 \\ 2/3 & 1/3 & 0 \end{pmatrix}; (40\%, 45\%, 15\%)$

30.  $P = \begin{pmatrix} 0 & 1 & 0 \\ 3/4 & 0 & 1/4 \\ 3/4 & 1/4 & 0 \end{pmatrix}; \left(\frac{15}{35}, \frac{16}{35}, \frac{4}{35}\right)$

31.  $\left(\frac{1}{3}, \frac{1}{3}, \frac{1}{6}, \frac{1}{6}\right)$

32.  $\frac{9}{20}, \frac{87}{200}, \frac{3}{7}$

33.  $\left(\frac{9}{27}, \frac{10}{27}, \frac{8}{27}\right)$

34.  $X_n \quad X_{n+1}$

$$\begin{matrix} & 0 & 1 & 2 & 3 \\ 0 & \left( \begin{matrix} 1/2 & 1/2 & 0 & 0 \end{matrix} \right) \\ 1 & \left( \begin{matrix} 1/2 & 0 & 1/2 & 0 \end{matrix} \right) \\ 2 & \left( \begin{matrix} 1/2 & 0 & 0 & 1/2 \end{matrix} \right) \\ 3 & \left( \begin{matrix} 0 & 0 & 0 & 1 \end{matrix} \right) \end{matrix}$$

The chain is not irreducible. State 3 is absorbing and other states are aperiodic.

# *Chapter* 8

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# Queueing Theory

**T**here are many situations in daily life when a queue is formed. For example, machines waiting to be repaired, patients waiting in a Doctor's room, cars waiting at a traffic signal and passengers waiting to buy tickets in counters form queues. Queue is formed if the service required by the customer (machine, patient, car, etc.) is not immediately available, that is, if the current demand for a particular service exceeds the capacity to provide the service.

Queues may be decreased in size or prevented from forming by providing additional service facilities which results in a drop in the profit. On the other hand, excessively long queues may result in lost sales and lost customers. Hence the problem of interest is how to achieve a balance between the cost associated with long waiting (queues) and the cost associated with the prevention of waiting in order to maximise the profits. As queueing theory provides an answer to this problem, it has become a topic of interest. Before we consider the solutions of queueing problems, we shall consider the general framework of a queueing system.

Although there are many types of queueing systems, all of them can be classified and described according to the following characteristics:

## **1. The Input (or Arrival) Pattern**

The input describes the manner in which the customers arrive and join the queueing system. It is not possible to observe and control the actual moment of arrival of a customer for service. Hence the number of arrivals in one time period or the interval between successive arrivals is not treated as a constant, but a random variable. So the mode of arrival of customers is expressed by means of the probability distribution of the number of arrivals per unit of time or of the inter-arrival time.

We shall mostly deal with only those queueing systems in which the number of arrivals per unit of time has a poisson distribution with mean  $\lambda$ . In this case, the time between consecutive arrivals has an exponential distribution with

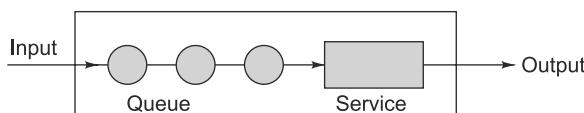
mean  $\frac{1}{\lambda}$  [Refer to Property 4 of the poisson process discussed in the previous Chapter 7].

Further the input process should specify the number of queues that are permitted to form, the maximum queue length and the maximum number of customers requiring service, viz., the nature of the source (finite or infinite) from which the customers emanate.

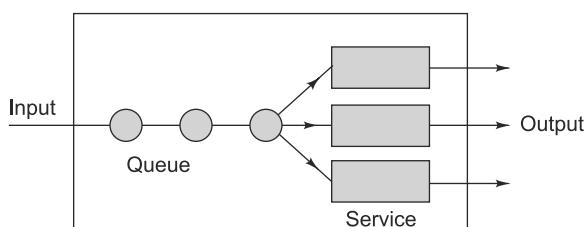
## 2. The Service Mechanism (or Pattern)

The mode of service is represented by means of the probability distribution of the number of customers serviced per unit of time or of the inter-service time. We shall mostly deal with only those queueing systems in which the number of customers serviced per unit of time has a Poisson distribution with mean  $\mu$  or equivalently the inter-service time (viz. the time to complete the service for a customer) has an exponential distribution with mean  $\frac{1}{\mu}$ .

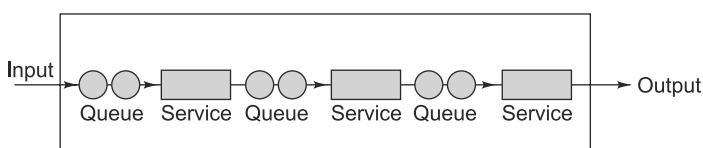
Further the service process should specify the number of servers and the arrangement of servers (in parallel, in series, etc.), as the behaviour of the queueing system depends on them also. The following figures represent the framework of queueing systems in which only one queue is permitted to form:



**Fig. 8.1 Single server queueing system**



**Fig. 8.2 Multiple servers (in parallel) queueing system**



**Fig. 8.3 Multiple servers (in series) queueing system**

### 3. The Queue Discipline

The queue discipline specifies the manner in which the customers form the queue or equivalently the manner in which they are selected for service, when a queue has been formed. The most common discipline is the *FCFS* (First Come First Served) or *FIFO* (First In First Out) as per which the customers are served in the strict order of their arrival. If the last arrival in the system is served first, we have the *LCFS* or *LIFO* (last in First Out) discipline. If the service is given in random order, we have the *SIRO* discipline. In the queueing systems which we deal with, we shall assume that service is provided on the *FCFS* (First Come First Served) basis.

## Symbolic Representation of a Queueing Model

Usually a queueing model is specified and represented symbolically in the form  $(a/b/c):(d/e)$ , where  $a$  denotes the type of distribution of the number of arrivals per unit time,  $b$  the type of distribution of the service time,  $c$  the number of servers,  $d$  the capacity of the system, viz., the maximum queue size and  $e$  the queue discipline.

Accordingly, the first four models which we will deal with will be denoted by the symbols  $(M/M/1):(\infty/FIFO)$ ,  $(M/M/s): (\infty/FIFO)$ ,  $(M/M/1): (k/FIFO)$  and  $(M/M/s): (k/FIFO)$ .

In the above symbols, the letter '*M*' stands for Markov' indicating that the number of arrivals in time  $t$  and the number of completed services in time  $t$  follow Poisson process which is a continuous time Markov chain.

## Difference Equations Related to Poisson Queue Systems

If the characteristics of a queueing system (such as the input and output parameters) are independent of time or equivalently if the behaviour of the system is independent of time, the system is said to be in *steady-state*. Otherwise it is said to be in *transient-state*.

Let  $P_n(t)$  be the probability that there are  $n$  customers in the system at time  $t$  ( $n > 0$ ). Let us first derive the differential equation satisfied by  $P_n(t)$  and then deduce the difference equation satisfied by  $P_n$  (probability of  $n$  customers at any time) in the steady-state.

Let  $\lambda_n$  be the average arrival rate when there are  $n$  customers in the system (both waiting in the queue and being served) and let  $\mu_n$  be the average service rate when there are  $n$  customers in the system.

**Note** The system being in steady-state does not mean that the arrival rate and service rate are independent of the number of customers in the system.

Now,  $P_n(t + \Delta t)$  is the probability of  $n$  customers at time  $t + \Delta t$ .

The presence of  $n$  customers in the system at time  $t + \Delta t$  can happen in any one of the following four mutually exclusive ways:

- (i) Presence of  $n$  customers at  $t$  and no arrival or departure during  $\Delta t$  time.
- (ii) Presence of  $(n - 1)$  customers at  $t$  and one arrival and no departure during  $\Delta t$  time.
- (iii) Presence of  $(n + 1)$  customers at  $t$  and no arrival and one departure during  $\Delta t$  time.
- (iv) Presence of  $n$  customers at  $t$  and one arrival and one departure during  $\Delta t$  time (since more than one arrival/departure during  $\Delta t$  is ruled out).

$$\therefore P_n(t + \Delta t) = P_n(t) (1 - \lambda_n \Delta t) (1 - \mu_n \Delta t) + P_{n-1}(t) \lambda_{n-1} \Delta t (1 - \mu_{n-1} \Delta t) \\ + P_{n+1}(t) (1 - \lambda_{n+1} \Delta t) \mu_{n+1} \Delta t + P_n(t) \cdot \lambda_n \Delta t \cdot \mu_n \Delta t \\ [\text{since } P(\text{an arrival occurs during } \Delta t \text{ time}) = \lambda \Delta t, \text{ etc.}]$$

$$\text{i.e., } P_n(t + \Delta t) = P_n(t) - (\lambda_n + \mu_n) P_n(t) \Delta t + \lambda_{n-1} P_{n-1}(t) \Delta t \\ + \mu_{n+1} P_{n+1}(t) \Delta t, \text{ on omitting terms containing } (\Delta t)^2 \text{ which is negligibly small.}$$

$$\therefore \frac{P_n(t + \Delta t) - P_n(t)}{\Delta t} = \lambda_{n-1} P_{n-1}(t) - (\lambda_n + \mu_n) P_n(t) + \mu_{n+1} P_{n+1}(t) \quad (1)$$

Taking limits on both sides of (1) as  $\Delta t \rightarrow 0$ , we have

$$P'_n(t) = \lambda_{n-1} P_{n-1}(t) - (\lambda_n + \mu_n) P_n(t) + \mu_{n+1} P_{n+1}(t) \quad (2)$$

Equation (2) does not hold good for  $n = 0$ , as  $P_{n-1}(t)$  does not exist. Hence we derive the differential equation satisfied by  $P_0(t)$  independently. Proceeding as before,

$$P_0(t + \Delta t) = P_0(t) (1 - \lambda_0 \Delta t) + P_1(t) (1 - \lambda_1 \Delta t) \mu_1 \Delta t, \\ [\text{by the possibilities (i) and (iii) given above and as no departure is possible when } n = 0]$$

$$\therefore \frac{P_0(t + \Delta t) - P_0(t)}{\Delta t} = -\lambda_0 P_0(t) + \mu_1 P_1(t) \quad (3)$$

Taking limits on both sides of (3) as  $\Delta t \rightarrow 0$ , we have

$$P'_0(t) = -\lambda_0 P_0(t) + \mu_1 P_1(t) \quad (4)$$

Now, in the steady-state,  $P_n(t)$  and  $P_0(t)$  are independent of time and hence  $P'_n(t)$  and  $P'_0(t)$  become zero. Hence, the differential equations (2) and (4) reduce to the difference equations

$$\lambda_{n-1} P_{n-1} - (\lambda_n + \mu_n) P_n + \mu_{n+1} P_{n+1} = 0 \quad (5)$$

$$\text{and } -\lambda_0 P_0 + \mu_1 P_1 = 0 \quad (6)$$

## Values of $P_0$ and $P_n$ for Poisson Queue Systems

From equation (6) derived above, we have

$$P_1 = \frac{\lambda_0}{\mu_1} P_0 \quad (7)$$

Putting  $n = 1$  in (5) and using (7), we have

$$\begin{aligned}\mu_2 P_2 &= (\lambda_1 + \mu_1) P_1 - \lambda_0 P_0 \\ &= (\lambda_1 + \mu_1) \frac{\lambda_0}{\mu_1} P_0 - \lambda_0 P_0 = \frac{\lambda_0 \lambda_1}{\mu_1} P_0 \\ \therefore P_2 &= \frac{\lambda_0 \lambda_1}{\mu_1 \mu_2} P_0\end{aligned}\tag{8}$$

Successively putting  $n = 2, 3, \dots$  in (5) and proceeding similarly, we can get

$$P_3 = \frac{\lambda_0 \lambda_1 \lambda_2}{\mu_1 \mu_2 \mu_3} P_0, \text{ etc.}$$

$$\text{Finally, } P_n = \frac{\lambda_0 \lambda_1 \lambda_2 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \mu_3 \cdots \mu_n} P_0, \text{ for } n = 1, 2, \dots\tag{9}$$

Since the number of customers in the system can be 0 or 1 or 2 or 3, etc.,

which events are mutually exclusive and exhaustive, we have  $\sum_{n=0}^{\infty} P_n = 1$ .

$$\begin{aligned}\text{i.e., } P_0 + \sum_{n=1}^{\infty} \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right) P_0 &= 1 \\ \therefore P_0 &= \frac{1}{1 + \sum_{n=1}^{\infty} \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right)}\end{aligned}\tag{10}$$

Equations (9) and (10) will be used to derive the important characteristics of the four queueing models.

## **Characteristics of Infinite Capacity, Single Server Poisson Queue Model I [M/M/1]: ( $\infty/\text{FIFO}$ ) model], when $\lambda_n = \lambda$ and $\mu_n = \mu$ ( $\lambda < \mu$ )**

1. *Average number  $L_s$  of customers in the system:* Let  $N$  denote the number of customers in the queueing system (i.e., those in the queue and the one who is being served).

$N$  is a discrete random variable, which can take the values  $0, 1, 2, \dots, \infty$

such that  $P(N = n) = P_n = \left(\frac{\lambda}{\mu}\right)^n P_0$ , from equation (9) of the previous discussion.

From equation (10), we have

$$P_0 = \frac{1}{1 + \sum_{n=1}^{\infty} \left( \frac{\lambda}{\mu} \right)^n} = \frac{1}{\sum_{n=0}^{\infty} \left( \frac{\lambda}{\mu} \right)^n} = 1 - \frac{\lambda}{\mu}$$

$$\therefore P_n = \left( \frac{\lambda}{\mu} \right)^n \left( 1 - \frac{\lambda}{\mu} \right)$$

$$\begin{aligned} \text{Now, } L_s &= E(N) = \sum_{n=0}^{\infty} n \times P_n \\ &= \left( \frac{\lambda}{\mu} \right) \left( 1 - \frac{\lambda}{\mu} \right) \sum_{n=1}^{\infty} n \left( \frac{\lambda}{\mu} \right)^{n-1} \\ &= \frac{\lambda}{\mu} \left( 1 - \frac{\lambda}{\mu} \right) \left( 1 - \frac{\lambda}{\mu} \right)^{-2}, \text{ by binomial summation} \\ &= \frac{\frac{\lambda}{\mu}}{1 - \frac{\lambda}{\mu}} = \frac{\lambda}{\mu - \lambda} \end{aligned} \quad (1)$$

2. Average number  $L_q$  of customers in the queue or Average length of the queue:

If  $N$  is the number of customers in the system, then the number of customers in the queue is  $(N - 1)$

$$\begin{aligned} \therefore L_q &= E(N - 1) = \sum_{n=1}^{\infty} (n - 1) P_n \\ &= \left( 1 - \frac{\lambda}{\mu} \right) \sum_{n=1}^{\infty} (n - 1) \left( \frac{\lambda}{\mu} \right)^n \\ &= \left( \frac{\lambda}{\mu} \right)^2 \left( 1 - \frac{\lambda}{\mu} \right) \sum_{n=2}^{\infty} (n - 1) \left( \frac{\lambda}{\mu} \right)^{n-2} \\ &= \left( \frac{\lambda}{\mu} \right)^2 \left( 1 - \frac{\lambda}{\mu} \right) \left( 1 - \frac{\lambda}{\mu} \right)^{-2} \\ &= \frac{\left( \frac{\lambda}{\mu} \right)^2}{1 - \frac{\lambda}{\mu}} = \frac{\lambda^2}{\mu(\mu - \lambda)} \end{aligned} \quad (2)$$

3. Average number  $L_w$  of customers in non-empty queues

$$\begin{aligned} L_w &= E\{(N - 1)/(N - 1) > 0\}, \text{ since the queue is non-empty} \\ &= \frac{E(N - 1)}{P(N - 1 > 0)} = \frac{\lambda^2}{\mu(\mu - \lambda)} \times \frac{1}{\sum_{n=2}^{\infty} P_n} \end{aligned}$$

$$\begin{aligned}
&= \frac{\lambda^2}{\mu(\mu - \lambda)} \times \frac{1}{\sum_{n=2}^{\infty} \left(\frac{\lambda}{\mu}\right)^n \left(1 - \frac{\lambda}{\mu}\right)} \\
&= \frac{\lambda^2}{\mu(\mu - \lambda)} \times \frac{1}{\left(\frac{\lambda}{\mu}\right)^2 \left(1 - \frac{\lambda}{\mu}\right) \sum_{n=0}^{\infty} \left(\frac{\lambda}{\mu}\right)^n} \\
&= \frac{\mu}{\mu - \lambda} \times \frac{1}{\left(1 - \frac{\lambda}{\mu}\right) \left(1 - \frac{\lambda}{\mu}\right)^{-1}} = \frac{\mu}{\mu - \lambda} \tag{3}
\end{aligned}$$

4. Probability that the number of customers in the system exceeds  $k$

$$\begin{aligned}
P(N > k) &= \sum_{n=k+1}^{\infty} P_n = \sum_{n=k+1}^{\infty} \left(\frac{\lambda}{\mu}\right)^n \left(1 - \frac{\lambda}{\mu}\right) \\
&= \left(\frac{\lambda}{\mu}\right)^{k+1} \left(1 - \frac{\lambda}{\mu}\right) \sum_{n=k+1}^{\infty} \left(\frac{\lambda}{\mu}\right)^{n-(k+1)} \\
&= \left(\frac{\lambda}{\mu}\right)^{k+1} \cdot \left(1 - \frac{\lambda}{\mu}\right) \sum_{n=0}^{\infty} \left(\frac{\lambda}{\mu}\right)^n \\
&= \left(\frac{\lambda}{\mu}\right)^{k+1} \cdot \left(1 - \frac{\lambda}{\mu}\right) \left(1 - \frac{\lambda}{\mu}\right)^{-1} = \left(\frac{\lambda}{\mu}\right)^{k+1} \tag{4}
\end{aligned}$$

5. Probability density function of the waiting time in the system

Let  $W_s$  be the continuous random variable that represents the waiting time of a customer in the system, viz, the time between arrival and completion of service.

Let its pdf be  $f(w)$  and let  $f(w/n)$  be the density function of  $W_s$  subject to the condition that there are  $n$  customers in the queueing system when the customer arrives,

$$\text{Then } f(w) = \sum_{n=0}^{\infty} f(w/n) P_n \tag{5}$$

Now,  $f(w/n) = \text{pdf of sum of } (n+1) \text{ service times (one part-service time of the customer being served + } n \text{ complete service times)}$

$= \text{pdf of sum of } (n+1) \text{ independent random variables, each of which is exponentially distributed with parameter } \mu$

$$= \frac{\mu^{n+1}}{n!} e^{-\mu w} w^n; w > 0 \text{ which is the pdf of Erlang distribution.}$$

[ $\because$  the mgf of the exponential distribution ( $\mu$ ) is  $\left(1 - \frac{t}{\mu}\right)^{-1}$  and hence the mgf of the sum of  $(n + 1)$  independent exponential ( $\mu$ ) variables is  $\left(1 - \frac{t}{\mu}\right)^{-n+1}$ , which is the mgf of Erlang distribution with parameters  $\mu$  and  $(n + 1)$ ] (refer to Erlang distribution in Chapter 5).

$$\begin{aligned} \therefore f(w) &= \sum_{n=0}^{\infty} \frac{\mu^{n+1}}{n!} e^{-\mu w} w^n \left(\frac{\lambda}{\mu}\right)^n \left(1 - \frac{\lambda}{\mu}\right), \text{ by (5)} \\ &= \mu e^{-\mu w} \left(1 - \frac{\lambda}{\mu}\right) \sum_{n=0}^{\infty} \frac{1}{n!} (\lambda w)^n \\ &= \mu \left(1 - \frac{\lambda}{\mu}\right) e^{-\mu w} e^{\lambda w}, \text{ by exponential summation} \\ &= (\mu - \lambda) e^{-(\mu - \lambda)w} \quad (6), \text{ which is the pdf of an exponential distribution with parameter } (\mu - \lambda). \end{aligned}$$

6. Average waiting time of a customer in the system:

$W_s$  follows an exponential distribution with parameter  $(\mu - \lambda)$ .

$$\therefore E(W_s) = \frac{1}{\mu - \lambda} \quad (7)$$

( $\because$  the mean of an exponential distribution is the reciprocal of its parameter).

7. Probability that the waiting time of a customer in the system exceeds  $t$

$$\begin{aligned} P(W_s > t) &= \int_t^{\infty} f(w) dw \\ &= \int_t^{\infty} (\mu - \lambda) e^{-(\mu - \lambda)w} dw \\ &= [-e^{-(\mu - \lambda)w}]_t^{\infty} = e^{-(\mu - \lambda)t} \end{aligned} \quad (8)$$

8. Probability density function of the waiting time  $W_q$  in the queue:

$W_q$  represents the time between arrival and reach of service point.

Let the pdf of  $W_q$  be  $g(w)$  and let  $g(w/n)$  be the density function of  $W_q$  subject to the condition that there are  $n$  customers in the system or there are  $(n - 1)$  customers in the queue apart from one customer receiving service. Now  $g(w/n) = \text{pdf of sum of } n \text{ service times} [\text{one residual service time} + (n - 1) \text{ full service times}]$

$$= \frac{\mu^n}{(n-1)!} e^{-\mu w} w^{n-1}; w > 0$$

$$\begin{aligned}
\therefore g(w) &= \sum_{n=1}^{\infty} \frac{\mu^n}{(n-1)!} e^{-\mu w} w^{n-1} \left(\frac{\lambda}{\mu}\right)^n \left(1 - \frac{\lambda}{\mu}\right) \\
&= \left(1 - \frac{\lambda}{\mu}\right) \lambda e^{-\mu w} \sum_{n=1}^{\infty} \frac{1}{(n-1)!} (\lambda w)^{n-1} \\
&= \frac{\lambda}{\mu} (\mu - \lambda) e^{-\mu w} e^{\lambda w} \\
&= \frac{\lambda}{\mu} (\mu - \lambda) e^{-(\mu - \lambda)w}; w > 0
\end{aligned} \tag{9}$$

and  $g(w) = 1 - \frac{\lambda}{\mu}$ , when  $w = 0$

**Note** 1.  $W_q$  is a continuous random variable in  $w > 0$  and it takes the value 0 with a non-zero probability. 2.  $W_q$  does not follow an exponential distribution.

#### 9. Average waiting time of a customer in the queue

$$\begin{aligned}
E(W_q) &= \frac{\lambda}{\mu} (\mu - \lambda) \int_0^{\infty} w e^{-(\mu - \lambda)w} dw \\
&= \frac{\lambda}{\mu} \int_0^{\infty} x e^{-x} \frac{dx}{\mu - \lambda} \\
&= \frac{\lambda}{\mu(\mu - \lambda)} [x(-e^{-x}) - e^{-x}]_0^{\infty} \\
&= \frac{\lambda}{\mu(\mu - \lambda)}
\end{aligned} \tag{10}$$

#### 10. Average waiting time of a customer in the queue, if he has to wait

$$\begin{aligned}
E(W_q/W > 0) &= \frac{E(W_q)}{P(W_q > 0)} \\
&= \frac{E(W_q)}{1 - P(W_q = 0)} \\
&= \frac{E(W_q)}{1 - P(\text{no customer in the queue})}
\end{aligned}$$

$$\begin{aligned}
 &= \frac{E(W_q)}{1 - P_0} \\
 &= \frac{\lambda}{\mu(\mu - \lambda)} \times \frac{\mu}{\lambda} \quad \left( \because P_0 = 1 - \frac{\mu}{\lambda} \right) \\
 &= \frac{1}{\mu - \lambda}
 \end{aligned} \tag{11}$$

## Relations among $E(N_s)$ , $E(N_q)$ , $E(W_s)$ and $E(W_q)$

$$(i) \quad E(N_s) = \frac{\lambda}{\mu - \lambda} = \lambda E(W_s) \quad [\because E(N_s) = L_s]$$

$$(ii) \quad E(N_q) = \frac{\lambda^2}{\mu(\mu - \lambda)} = \lambda E(W_q) \quad [\because E(N_q) = L_q]$$

$$(iii) \quad E(W_s) = E(W_q) + \frac{1}{\mu}$$

$$(iv) \quad E(N_s) = E(N_q) + \frac{\lambda}{\mu}$$

**Note**

- If any one of the quantities  $E(N_s)$ ,  $E(N_q)$ ,  $E(W_s)$  and  $E(W_q)$  is known, the other three can be found out using the relations given above.
- The above relations, called Little's formulas hold good for the models with infinite capacity, but with a slight modification for the models with finite capacity.

## Characteristics of Infinite Capacity, Multiple Server Poisson Queue Model II [M/M/s]: ( $\infty$ /FIFO) model], When $\lambda_n = \lambda$ for all $n$ ( $\lambda < s\mu$ )

- Values of  $P_0$  and  $P_n$

For the Poisson queue system,  $P_n$  is given by

$$P_n = \frac{\lambda_0 \lambda_1 \lambda_2 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \mu_3 \cdots \mu_n} \times P_0, \quad n \geq 1, \tag{1}$$

$$\text{where } P_0 = \left[ 1 + \sum_{n=1}^{\infty} \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right) \right]^{-1} \tag{2}$$

If there is a single server,  $\mu_n = \mu$  for all  $n$ . But there are  $s$  servers working independently of each other.

If there be less than  $s$  customers, i.e., if  $n < s$ , only  $n$  of the  $s$  servers will be busy and the others idle and hence the mean service rate will be  $n\mu$ .

If  $n \geq s$ , all the  $s$  servers will be busy and hence the mean service rate =  $s\mu$ .

$$\text{hence, } \mu_n = \begin{cases} n\mu, & \text{if } 0 \leq n < s \\ s\mu, & \text{if } n \geq s \end{cases} \quad (3)$$

Using (3) in (1) and (2), we have

$$\begin{aligned} P_n &= \frac{\lambda^n}{1\mu \cdot 2\mu \cdot 3\mu \cdots n\mu} P_0, \text{ if } 0 \leq n < s \\ &= \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n P_0, \text{ if } 0 \leq n < s \end{aligned} \quad (4)$$

$$\begin{aligned} \text{and } P_n &= \frac{\lambda^n}{\{1\mu \cdot 2\mu \cdots (s-1)\mu\} \{s\mu \cdot s\mu \cdots (n-s+1) \text{ times}\}} P_0 \\ &= \frac{\lambda^n}{(s-1)! \mu^{s-1} (s\mu)^{n-s+1}} P_0 \\ &= \frac{1}{s! s^{n-s}} \left( \frac{\lambda}{\mu} \right)^n P_0, \text{ if } n \geq s \end{aligned} \quad (5)$$

Now,  $P_0$  is given by  $\sum_{n=0}^{\infty} P_n = 1$

$$\text{i.e. } \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \sum_{n=s}^{\infty} \frac{1}{s! s^{n-s}} \left( \frac{\lambda}{\mu} \right)^n \right] P_0 = 1$$

$$\text{i.e., } \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{s^s}{s!} \sum_{n=s}^{\infty} \left( \frac{\lambda}{\mu s} \right)^n \right] P_0 = 1$$

$$\text{i.e., } \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{s^s}{s!} \left( \frac{\lambda}{\mu s} \right)^s \frac{1}{1 - \frac{\lambda}{\mu s}} \right] P_0 = 1$$

$$\text{i.e., } \left[ \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n \right\} + \left\{ \frac{1}{s! \left( 1 - \frac{\lambda}{\mu s} \right)} \left( \frac{\lambda}{\mu} \right)^s \right\} \right] P_0 = 1$$

$$\text{or} \quad P_0 = \frac{1}{\left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n \right\} + \left\{ \frac{1}{s! \left( 1 - \frac{\lambda}{\mu s} \right)} \left( \frac{\lambda}{\mu} \right)^s \right\}} \quad (6)$$

2. Average number of customers in the queue or average queue length

$$\begin{aligned} L_q &= E(N_q) = E(N - s) = \sum_{n=s}^{\infty} (n - s) P_n \\ &= \sum_{x=0}^{\infty} x P_{x+s} \\ &= \sum_{x=0}^{\infty} x \times \frac{1}{s! s^x} \left( \frac{\lambda}{\mu} \right)^{s+x} P_0 \\ &= \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s P_0 \sum_{x=0}^{\infty} x \left( \frac{\lambda}{\mu s} \right)^x \\ &= \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \frac{\lambda}{\mu s} \cdot P_0 \frac{1}{\left( 1 - \frac{\lambda}{\mu s} \right)^2} \\ &= \frac{1}{s \cdot s!} \cdot \frac{\left( \frac{\lambda}{\mu} \right)^{s+1}}{\left( 1 - \frac{\lambda}{\mu s} \right)^2} P_0 \end{aligned} \quad (7)$$

3. Average number of customers in the system

By Little's formula (iv),

$$\begin{aligned} E(N_s) &= E(N_q) + \frac{\lambda}{\mu} \\ &= \frac{1}{s \cdot s!} \cdot \frac{\left( \frac{\lambda}{\mu} \right)^{s+1}}{\left( 1 - \frac{\lambda}{\mu s} \right)^2} P_0 + \frac{\lambda}{\mu} \end{aligned} \quad (8)$$

Result (8) can also be directly derived by using the definition  $E(N_s) =$

$$\sum_{n=0}^{\infty} n P_n.$$

4. Average time a customer has to spend in the system

By Little's formula (i),

$$\begin{aligned} E(W_s) &= \frac{1}{\lambda} E(N_s) \\ &= \frac{1}{\mu} + \frac{1}{\mu} \cdot \frac{1}{s \cdot s!} \frac{\left(\frac{\lambda}{\mu}\right)^s}{\left(1 - \frac{\lambda}{\mu s}\right)^2} \cdot P_0 \end{aligned} \quad (9)$$

5. Average time a customer has to spend in the queue

By Little's formula (ii),

$$\begin{aligned} E(W_q) &= \frac{1}{\lambda} E(N_q) \\ &= \frac{1}{\mu} \cdot \frac{1}{s \cdot s!} \frac{\left(\frac{\lambda}{\mu}\right)^s}{\left(1 - \frac{\lambda}{\mu s}\right)^2} \cdot P_0 \end{aligned} \quad (10)$$

6. Probability that an arrival has to wait

Required probability = Probability that there are  $s$  or more customers in the system

i.e.,  $P(W_s > 0) = P(N \geq s)$

$$\begin{aligned} &= \sum_{n=s}^{\infty} P_n = \sum_{n=s}^{\infty} \frac{1}{s! s^{n-s}} \left(\frac{\lambda}{\mu}\right)^n P_0 \\ &= \frac{1}{s!} \left(\frac{\lambda}{\mu}\right)^s \cdot P_0 \sum_{n=s}^{\infty} \left(\frac{\lambda}{\mu s}\right)^{n-s} \\ &= \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \end{aligned} \quad (11)$$

7. Probability that an arrival enters the service without waiting

Required probability

$$= 1 - P(\text{an arrival has to wait})$$

$$= 1 - \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \quad (12)$$

8. Mean waiting time in the queue for those who actually wait.

$$\begin{aligned} E(W_q/W_s > 0) &= \frac{E(W_q)}{P(W_s > 0)} \\ &= \frac{1}{\mu} \cdot \frac{1}{s \cdot s!} \frac{\left(\frac{\lambda}{\mu}\right)^s}{\left(1 - \frac{\lambda}{\mu s}\right)^2} P_0 \times \frac{s! \left(1 - \frac{\lambda}{\mu s}\right)}{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0} \quad [\text{using (10) and 11}] \\ &= \frac{1}{\mu s \left(1 - \frac{\lambda}{\mu s}\right)} = \frac{1}{\mu s - \lambda} \end{aligned} \quad (13)$$

9. Probability that there will be someone waiting

Required probability =  $P(N \geq s + 1)$

$$\begin{aligned} &= \sum_{n=s+1}^{\infty} P_n = \sum_{n=s}^{\infty} P_n - P(N=s) \\ &= \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right)} - \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s!} \quad [\text{using (10) and (5)}] \\ &= \frac{\left(\frac{\lambda}{\mu}\right)^s P_0}{s!} \cdot \frac{\left(\frac{\lambda}{\mu s}\right)}{1 - \frac{\lambda}{\mu s}} \end{aligned} \quad (14)$$

10. Average number of customers (in non-empty queues), who have to actually wait.

$$\begin{aligned} L_w &= E(N_q/N_q \geq 1) \\ &= E(N_q)/P(N \geq s) \\ &= \frac{1}{s \cdot s!} \frac{\left(\frac{\lambda}{\mu}\right)^{s+1} \cdot P_0}{\left(1 - \frac{\lambda}{\mu s}\right)^2} \cdot \frac{s! \left(1 - \frac{\lambda}{\mu s}\right)}{\left(\frac{\lambda}{\mu}\right)^s P_0} \end{aligned}$$

$$= \frac{\left(\frac{\lambda}{\mu s}\right)}{1 - \frac{\lambda}{\mu s}} \quad (15)$$

### Characteristics of Finite Capacity, Single-Server Poisson Queue Model III [(M/M/1): (k/FIFO) Model]

#### 1. Values of $P_0$ and $P_n$

For the Poisson queue system,  $P_n = P(N = n)$  in the steady-state is given by the difference equations

$$\begin{aligned} \lambda_{n-1} P_{n-1} - (\lambda_n + \mu_n) P_n + \mu_{n+1} P_{n+1} &= 0; n > 0 \\ \text{and } -\lambda_0 P_0 + \mu_1 P_1 &= 0; n = 0 \end{aligned}$$

This model represents the situation in which the system can accommodate only a finite number  $k$  of arrivals. If a customer arrives and the queue is full, the customer leaves without joining the queue.

Therefore, for this model,

$$\mu_n = \mu, n = 1, 2, 3, \dots$$

and  $\lambda_n = \begin{cases} \lambda, & \text{for } n = 0, 1, 2, \dots (k-1) \\ 0, & \text{for } n = k, k+1, \dots \end{cases}$

Using these values in the difference equations given above, we have

$$\mu P_1 = \lambda P_0 \quad (1)$$

$$\mu P_{n+1} = (\lambda + \mu) P_n - \lambda P_{n-1}, \text{ for } 1 \leq n \leq k-1 \quad (2)$$

and  $\mu P_k = \lambda P_{k-1}, \text{ for } n = k \quad (3) \quad (\because P_{k+1} \text{ has no meaning})$

From (1),  $P_1 = \frac{\lambda}{\mu} P_0$

From (2),  $\mu P_2 = (\lambda + \mu) \frac{\lambda}{\mu} P_0 - \lambda P_0$

$\therefore P_2 = \left(\frac{\lambda}{\mu}\right)^2 P_0$  and so on

In general,  $P_n = \left(\frac{\lambda}{\mu}\right)^n P_0, \text{ true for } 0 \leq n \leq k-1 \quad (3)$

From (3),  $P_k = \frac{\lambda}{\mu} \left(\frac{\lambda}{\mu}\right)^{k-1} P_0 = \left(\frac{\lambda}{\mu}\right)^k P_0$

Now,

$$\sum_{n=0}^k P_n = 1$$

i.e.,  $P_0 \sum_{n=0}^k \left(\frac{\lambda}{\mu}\right)^n = 1$

i.e.,  $P_0 \frac{\left\{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}\right\}}{1 - \frac{\lambda}{\mu}} = 1,$

which is valid even for  $\lambda > \mu$

$$\therefore P_0 = \begin{cases} \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}}, & \text{if } \lambda \neq \mu \\ \frac{1}{k+1}, & \text{if } \lambda = \mu, \text{ since } \lim_{\frac{\lambda}{\mu} \rightarrow 1} \left\{ \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \right\} = \frac{1}{k+1} \end{cases} \quad (4)$$

$$\therefore P_n = \begin{cases} \left(\frac{\lambda}{\mu}\right)^n \cdot \left[ \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \right], & \text{if } \lambda \neq \mu \\ \frac{1}{k+1}, & \text{if } \lambda = \mu \end{cases} \quad (6) \quad (7)$$

2. Average number of customers in the system

$$E(N) = \sum_{n=0}^k n P_n = \frac{\left(1 - \frac{\lambda}{\mu}\right)}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \cdot \sum_{n=0}^k n \left(\frac{\lambda}{\mu}\right)^n$$

$$= \frac{\left(1 - \frac{\lambda}{\mu}\right) \cdot \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \cdot \sum_{n=0}^k \frac{d}{dx} (x^n), \text{ where } x = \frac{\lambda}{\mu}$$

$$\begin{aligned}
&= \frac{\left(1 - \frac{\lambda}{\mu}\right) \cdot \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \frac{d}{dx} \left( \frac{1 - x^{k+1}}{1 - x} \right) \\
&= \frac{(1-x)x}{1-x^{k+1}} \cdot \left[ \frac{(1-x)\{- (k+1)x^k\} + (1-x^{k+1})}{(1-x)^2} \right] \\
&= \frac{x \left[ 1 - (k+1)x^k + kx^{k+1} \right]}{(1-x)(1-x^{k+1})} \\
&= \frac{x(1-x^{k+1}) - (k+1)(1-x)x^{k+1}}{(1-x)(1-x^{k+1})} \\
&= \frac{x}{1-x} - \frac{(k+1)x^{k+1}}{1-x^{k+1}} \\
&= \frac{\lambda}{\mu - \lambda} - \frac{(k+1)\left(\frac{\lambda}{\mu}\right)^{k+1}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}}, \text{ if } \lambda \neq \mu
\end{aligned} \tag{8}$$

and  $E(N) = \sum_{n=0}^k \frac{n}{k+1} = \frac{k}{2}$ , if  $\lambda = \mu$  (9)

3. Average number of customers in the queue.

$$\begin{aligned}
E(N_q) &= E(N-1) = \sum_{n=1}^k (n-1)P_n \\
&= \sum_{n=0}^k nP_n - \sum_{n=1}^k P_n \\
&= E(N) - (1 - P_0)
\end{aligned} \tag{10}$$

As per Little's formula (iv),

$$E(N_q) = E(N) - \frac{\lambda}{\mu},$$

which is true when the average arrival rate is  $\lambda$  throughout. Now, we see that, in step (8),  $1 - P_0 \neq \frac{\lambda}{\mu}$ , because the average arrival rate is  $\lambda$  as long as there is a vacancy in the queue and it is zero when the system is full.

Hence, we define *the overall effective arrival rate*, denoted by  $\lambda'$  or  $\lambda_{\text{eff}}$ , by using step (8) and Little's formula as

$$\frac{\lambda'}{\mu} = 1 - P_0 \text{ or } \lambda' = \mu(1 - P_0) \quad (11)$$

Thus, step (8) can be rewritten as

$$E(N_q) = E(N) - \frac{\lambda'}{\mu}, \quad (12)$$

which is the modified Little's formula for this model.

#### 4. Average waiting times in the system and in the queue:

By the modified Little's formulas,

$$E(W_s) = \frac{1}{\lambda'} E(N) \quad (13)$$

$$\text{and} \quad E(W_q) = \frac{1}{\lambda'} E(N_q) \quad (14)$$

where  $\lambda'$  is the effective arrival rate, given by step (9).

### Characteristics of Finite Queue, Multiple Server Poisson Queue Model IV [(M/M/s): (k/FIFO) Model]

#### 1. Values of $P_0$ and $P_n$

For the Poisson queue system,  $P_n$  is given by

$$P_n = \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \cdot P_0, \quad n \geq 1, \quad (1)$$

$$\text{where} \quad P_0 = \left\{ 1 + \sum_{n=1}^k \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right) \right\}^{-1} \quad (2)$$

For this  $(M/M/s): (k/FIFO)$  model,

$$\lambda_n = \begin{cases} \lambda, & \text{for } n = 0, 1, 2, \dots, k-1 \\ 0, & \text{for } n = k, k+1, \dots \end{cases}$$

$$\mu_n = \begin{cases} n\mu, & \text{for } n = 0, 1, 2, \dots, s-1 \\ s\mu, & \text{for } n = s, s+1, \dots \end{cases}$$

Using these values of  $\lambda_n$  and  $\mu_n$  in (2) and noting that  $1 < s < k$ , we get

$$\begin{aligned} P_0^{-1} &= \left\{ 1 + \frac{\lambda}{1! \mu} + \frac{\lambda^2}{2! \mu^2} + \cdots + \frac{\lambda^{s-1}}{(s-1)! \mu^{s-1}} \right\} + \left\{ \frac{\lambda^s}{(s-1)! \mu^{s-1} \cdot \mu s} \right. \\ &\quad \left. + \frac{\lambda^{s+1}}{(s-1)! \mu^{s-1} \cdot (\mu s)^2} + \cdots + \frac{\lambda^k}{(s-1)! \mu^{s-1} \cdot (\mu s)^{k-s+1}} \right\} \end{aligned}$$

$$= \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{\lambda^s}{s! \mu^s} \left[ 1 + \frac{\lambda}{\mu s} + \left( \frac{\lambda}{\mu s} \right)^2 + \dots + \left( \frac{\lambda}{\mu s} \right)^{k-s} \right]$$

$$= \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \sum_{n=s}^k \left( \frac{\lambda}{\mu s} \right)^{n-s} \quad (3)$$

$$P_n = \begin{cases} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n P_0, & \text{for } n \leq s \\ \frac{1}{s! s^{n-s}} \cdot \left( \frac{\lambda}{\mu} \right)^n \cdot P_0, & \text{for } s < n \leq k \\ 0, & \text{for } n > k \end{cases} \quad (4)$$

2. Average queue length or average number of customers in the queue

$$\begin{aligned} E(N_q) &= E(N - s) = \sum_{n=s}^k (n - s) P_n \\ &= \frac{P_0}{s!} \sum_{n=s}^k (n - s) \left( \frac{\lambda}{\mu} \right)^n / s^{n-s} \quad [\text{using (4)}] \\ &= \frac{\left( \frac{\lambda}{\mu} \right)^s \cdot P_0}{s!} \sum_{x=0}^{k-s} x \cdot \left( \frac{\lambda}{\mu s} \right)^x \\ &= \frac{\left( \frac{\lambda}{\mu} \right)^s \cdot P_0 \rho_{k-s}}{s!} \sum_{x=0}^{k-s} x \rho^{x-1} \text{ where } P = \frac{\lambda}{\mu s} \\ &= \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{P_0 \rho}{s!} \sum_{x=0}^{k-s} \frac{d}{d\rho} (\rho^x) \\ &= \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{P_0 \rho}{s!} \frac{d}{d\rho} \left\{ \frac{1 - \rho^{k-s+1}}{1 - \rho} \right\} \\ &= \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{P_0 \rho}{s!} \left[ \frac{-(1-\rho)(k-s+1)\rho^{k-s} + (1-\rho^{k-s+1})}{(1-\rho)^2} \right] \\ &= \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{P_0 \rho}{s!} \left[ \frac{-(k-s)(1-\rho)\rho^{k-s} - (1-\rho)\rho^{k-s} + 1 - \rho^{k-s+1}}{(1-\rho)^2} \right] \end{aligned}$$

$$\begin{aligned}
 &= P_0 \left( \frac{\lambda}{\mu} \right)^s \frac{\rho}{s!} \left[ \frac{-(k-s)(1-\rho)\rho^{k-s} + 1 - \rho^{k-s}(1-\rho+\rho)}{(1-\rho)^2} \right] \\
 &= P_0 \cdot \left( \frac{\lambda}{\mu} \right)^s \frac{\rho}{s!(1-\rho)^2} [1 - \rho^{k-s} - (k-s)(1-\rho)\rho^{k-s}],
 \end{aligned}$$

where  $\rho = \frac{\lambda}{\mu s}$  (5)

3. Average number of customers in the system

$$\begin{aligned}
 E(N) &= \sum_{n=0}^k nP_n = \sum_{n=0}^{s-1} nP_n + \sum_{n=s}^k nP_n \\
 &= \sum_{n=0}^{s-1} nP_n + \sum_{n=s}^k (n-s)P_n + \sum_{n=s}^k sP_n \\
 &= \sum_{n=0}^{s-1} nP_n + E(N_q) + s \left\{ \sum_{n=0}^k P_n - \sum_{n=0}^{s-1} P_n \right\} \\
 &= E(N_q) + s - \sum_{n=0}^{s-1} (s-n)P_n \quad \left( \because \sum_{n=0}^k P_n = 1 \right) \quad (6)
 \end{aligned}$$

Obviously,  $\left\{ s - \sum_{n=0}^{s-1} (s-n)P_n \right\} \neq \frac{\lambda}{\mu}$ , so that step (6) represents Little's formula.

In order to make (6) to assume the form of Little's formula, we define the *overall effective arrival rate*  $\lambda'$  or  $\lambda_{\text{eff}}$  as follows:

$$\begin{aligned}
 \frac{\lambda'}{\mu} &= s - \sum_{n=0}^{s-1} (s-n)P_n \\
 \text{i.e., } \lambda' &= \mu \left[ s - \sum_{n=0}^{s-1} (s-n)P_n \right] \quad (7)
 \end{aligned}$$

With this definition of  $\lambda'$ , step (6) becomes

$$E(N) = E(N_q) + \frac{\lambda'}{\mu} \quad (8)$$

which is the modified Little's formula for this model.

4. Average waiting time in the system and in the queue:

By the modified Little's formulas,

$$E(W_s) = \frac{1}{\lambda'} E(N) \quad (9)$$

and  $E(W_q) = \frac{1}{\lambda'} E(N_q)$  (10)

where  $\lambda'$  is the effective arrival rate, given by step (7).

## More Poisson Queue Models

In the topic ‘Birth and Death process’ discussed in pages 7.50 and 7.51 of the Chapter 7 on “Special Random Processes”, we have derived that

$$P'_n(t) = \lambda_{n-1} P_{n-1}(t) - (\lambda_n + \mu_n) P_n(t) + \mu_{n+1} P_{n+1}(t) \quad (2)$$

and  $P'_0(t) = -\lambda_0 P_0(t) + \mu_1 P_1(t)$  (4)

where  $P_n(t) = P\{X(t) = n\}$  and  $X(t)$  represents the number of individuals present at time  $t$  in a population.

In the discussion of the topic “Difference equation related to Poisson Queue Systems” given in pages 8.3–8.5, we have noted that in the steady-state,  $P_n(t)$  and  $P_0(t)$  are independent of time and hence  $P'_n(t)$  and  $P'_0(t)$  become zero. Also the differential equations (2) and (4) reduce to the difference equations

$$\lambda_{n-1} P_{n-1} - (\lambda_n + \mu_n) P_n + \mu_{n+1} P_{n+1} = 0 \quad (5)$$

and  $-\lambda_0 P_0 + \mu_1 P_1 = 0$  (6)

Solving (5) and (6), we have derived the values of  $P_0$  and  $P_n$  as given below:

$$P_0^{-1} = 1 + \sum_{n=1}^{\infty} \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right) \quad (10)$$

and  $P_n = \frac{\lambda_0 \lambda_1 \lambda_2 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \mu_3 \cdots \mu_n} P_0$ , for  $n = 1, 2, 3, \text{etc.} \dots$  (9)

### Particular Cases

#### (i) Queues with impatient customers

In this case,  $\lambda_n = \frac{\lambda}{n+1}$ , for  $n \geq 0$  and  $\mu_n = \mu$ , for  $n \geq 1$  viz., the arrival rate decreases with increases in queue length, while the service rate  $\mu_n$  is independent of  $n$  and remains the same as  $\mu$ . This situation arises due to *baulking and reneging* of the customers. The word “baulking” means refusing to go forward, viz. to join the queue. The word “reneging” means going back on one’s determination to join the queue. In this case, the queue itself is called the queue with discouragement.

Using  $\lambda_n = \frac{\lambda}{n+1}$  and  $\mu_n = \mu$  in steps (10) and (9), we get

$$P_0^{-1} = 1 + \sum_{n=1}^{\infty} \frac{\lambda^n}{n! \mu^n} = 1 + \frac{\rho}{1!} + \frac{\rho^2}{2!} + \frac{\rho^3}{3!} + \dots \infty = e^{\rho}$$

and so  $P_0 = e^{-\rho}$ , where  $\rho = \frac{\lambda}{\mu}$ .

and

$$P_n = \frac{\lambda}{1.\mu} \cdot \frac{\lambda}{2\mu} \cdot \frac{\lambda}{3\mu} \cdots \frac{\lambda}{n\mu} P_0 = \left( \frac{1}{n!} \rho^n \right) e^{-\rho}, \text{ for } n \geq 0$$

Since  $P_n = P(X = n)$ ,  $X$ , the number of customers in the system at any time follows a poison distribution with parameter  $\rho = \frac{\lambda}{\mu}$ .

**Note** For this model,  $E(N_S) = E(X) = \rho$ . By Little's formulas,

$$E(N_Q) = E(N_S) - \rho = 0, \quad E(W_Q) = \frac{1}{\lambda} E(N_Q) = 0$$

and

$$E(W_S) = E(W_Q) + \frac{1}{\mu} = \frac{1}{\mu}$$

### (ii) Self-service Mode

In this case,  $\lambda_n = \lambda$ , for  $n \geq 0$  and  $\mu_n = n\mu$ , for  $n \geq 1$ . viz., the arrival rate  $\lambda_n (= \lambda)$  is independent of  $n$ , while the service rate  $\mu_n (= n\mu)$  increases with increases of  $n$ . viz., if there are  $n$  customers in the system,  $n$  servers will be available for serving for  $n = 1, 2, 3, \dots \infty$ . Thus, each arrival will enter the service facility without waiting and act as a server himself/herself. Hence, the queue model is called *self-service model*.

As the number of servers is unlimited, since each customer is also the server, the queue itself is called *Queue with infinite number of channels* and is denoted as  $(M/M/\infty)$ : ( $\infty$ /FIFO).

**Note** For this model also,  $P_0 = e^{-\rho}$  and  $P_n = \frac{1}{n!} \rho^n e^{-\rho}$ , for  $n \geq 0$  where  $\rho = \frac{\lambda}{\mu}$ . Moreover

$$E(N_Q) = 0, \quad E(W_Q) = 0, \quad \text{as no queue will be formed and } E(N_S) = \rho \text{ and } E(W_S) = \frac{1}{\mu}.$$

## Worked Example 8

### Example 1

Arrivals at a telephone booth are considered to be Poisson with an average time of 12 min. between one arrival and the next. The length of a phone call is assumed to be distributed exponentially with mean 4 min.

- (a) Find the average number of persons waiting in the system.
- (b) What is the probability that a person arriving at the booth will have to wait in the queue?
- (c) What is the probability that it will take him more than 10 min. altogether to wait for the phone and complete his call?

- (d) Estimate the fraction of the day when the phone will be in use.  
 (e) The telephone department will install a second booth, when convinced that an arrival has to wait on the average for at least 3 min. for phone. By how much the flow of arrivals should increase in order to justify a second booth?  
 (f) What is the average length of the queue that forms from time to time?

$$\text{Mean inter-arrival time} = \frac{1}{\lambda} = 12 \text{ min.}$$

Therefore, mean arrival rate  $= \lambda = \frac{1}{12}$  per minute.

$$\text{Mean service time} = \frac{1}{\mu} = 4 \text{ min.}$$

Therefore, mean service rate  $= \mu = \frac{1}{4}$  per minute.

(a)  $E(N) = \frac{\lambda}{\mu - \lambda}$ , (by formula (1) of model I)

$$= \frac{\frac{1}{12}}{\frac{1}{4} - \frac{1}{12}} = 0.5 \text{ customer}$$

(b)  $P(W > 0) = 1 - P(W = 0)$   
 $= 1 - P(\text{no customer in the system})$   
 $= 1 - P_0$   
 $= 1 - \left(1 - \frac{\lambda}{\mu}\right)$  (by the formula for  $P_0$  of model I)

$$= \frac{\lambda}{\mu} = \frac{1/12}{1/4} = \frac{1}{3}$$

(c)  $P(W > 10) = e^{-(\mu - \lambda) \times 10}$  [by formula (8) of model I]  
 $= e^{-\left(\frac{1}{4} - \frac{1}{12}\right) \times 10}$   
 $= e^{-\frac{5}{3}} = 0.1889$

(d)  $P(\text{the phone will be idle}) = P(N = 0) = P_0$   
 $= 1 - \frac{\lambda}{\mu} = \frac{2}{3}$

$$\therefore P(\text{the phone will be in use}) = 1 - \frac{2}{3} = \frac{1}{3}$$

or the fraction of the day when the phone will be in use =  $\frac{1}{3}$ .

- (e) The second phone will be installed, if  $E(W_q) > 3$ .

$$\text{i.e., if } \frac{\lambda}{\mu(\mu - \lambda)} > 3 \text{ [by formula (10) of model I]}$$

$$\text{i.e., if } \frac{\lambda_R}{\frac{1}{4}\left(\frac{1}{4} - \lambda_R\right)} > 3,$$

where  $\lambda_R$  is the required arrival rate.

$$\text{i.e., if } \lambda_R > \frac{3}{4} \left( \frac{1}{4} - \lambda_R \right)$$

$$\text{i.e., if } \lambda_R > \frac{3}{28}$$

Hence, the arrival rate should increase by  $\frac{3}{28} - \frac{1}{12} = \frac{1}{42}$  per minute, to justify a second phone.

- (f)  $E(N_q/\text{the queue is always available})$

$$= E(N_q/N_q > 0)$$

$$= E(N_q/N > 1)$$

$$= \frac{E(N_q)}{P(N > 1)} = \frac{E(N_q)}{1 - P_0 - P_1}$$

$$= \frac{\lambda^2}{\mu(\mu - \lambda)} \times \frac{1}{1 - \left(1 + \frac{\lambda}{\mu}\right)P_0} \quad [\text{by formula (2) of model 1}]$$

$$= \frac{\lambda^2}{\mu(\mu - \lambda)} \cdot \frac{1}{1 - \left(1 + \frac{\lambda}{\mu}\right)\left(1 - \frac{\lambda}{\mu}\right)}$$

$$= \frac{\lambda^2}{\mu(\mu - \lambda)} \cdot \frac{\mu^2}{\lambda^2} = \frac{\mu}{\mu - \lambda} = \frac{1/4}{1/4 - 1/12} = 1.5 \text{ persons.}$$

---

**Example 2**

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Customers arrive at a one-man barber shop according to a Poisson process with a mean interarrival time of 12 min. Customers spend an average of 10 min in the barber's chair.

- What is the expected number of customers in the barber shop and in the queue?
- Calculate the percentage of time an arrival can walk straight into the barber's chair without having to wait.
- How much time can a customer expect to spend in the barber's shop?
- Management will provide another chair and hire another barber, when a customer's waiting time in the shop exceeds 1.25 h. How much must the average rate of arrivals increase to warrant a second barber?
- What is the average time customers spend in the queue?
- What is the probability that the waiting time in the system is greater than 30 min?
- Calculate the percentage of customers who have to wait prior to getting into the barber's chair.
- What is the probability that more than 3 customers are in the system?

$$\frac{1}{\lambda} = 12 \quad \therefore \lambda = \frac{1}{12} \text{ per minute}$$

$$\frac{1}{\mu} = 10 \quad \therefore \mu = \frac{1}{10} \text{ per minute}$$

$$(a) E(N_s) = \frac{\lambda}{\mu - \lambda} = \frac{1/12}{1/10 - 1/12} = 5 \text{ customers [by formula (1) of model I]}$$

$$E(N_q) = \frac{\lambda^2}{\mu(\mu - \lambda)} \text{ [by formula (2) of model I]}$$

$$= \frac{\frac{1}{144}}{\frac{1}{10} \left( \frac{1}{10} - \frac{1}{12} \right)} = 4.17 \text{ customers}$$

$$(b) P(\text{a customer straight goes to the barber's chair})$$

$$= P(\text{No customer in the system})$$

$$= P_0 = 1 - \frac{\lambda}{\mu} = 1 - \frac{\frac{1}{12}}{\frac{1}{10}} = \frac{1}{6}$$

Therefore, percentage of time an arrival need not wait = 16.7.

$$(c) \quad E(W) = \frac{1}{\mu - \lambda}, [\text{by formula (7) of model I}]$$

$$= \frac{1}{\frac{1}{10} - \frac{1}{12}} = 60 \text{ min or } 1 \text{ h}$$

$$(d) \quad E(W) > 75, \text{ if } \frac{1}{\mu - \lambda_r} > 75$$

$$\text{i.e.,} \quad \text{if } \lambda_r > \mu - \frac{1}{75}$$

$$\text{i.e.,} \quad \text{if } \lambda_r > \frac{1}{10} - \frac{1}{75}$$

$$\text{i.e.,} \quad \text{if } \lambda_r > \frac{13}{150}$$

Hence, to warrant a second barber, the average arrival rate must increase by  $\frac{13}{150} - \frac{1}{12} = \frac{1}{300}$  per minute.

$$(e) \quad E(W_q) = \frac{\lambda}{\mu(\mu - \lambda)}, [\text{by formula (10) of model I}]$$

$$= \frac{\frac{1}{12}}{\frac{1}{10} \left( \frac{1}{10} - \frac{1}{12} \right)} = 50 \text{ min}$$

$$(f) \quad P(W > t) = e^{-(\mu - \lambda)t}, [\text{by formula (8) of model I}]$$

$$\therefore P(W > 30) = e^{-\left(\frac{1}{10} - \frac{1}{12}\right) \times 30}$$

$$= e^{-0.5} = 0.6065$$

$$(g) \quad P(\text{a customer has to wait}) = P(W > 0)$$

$$= 1 - P(W = 0)$$

$$= 1 - P(N = 0) = 1 - P_0$$

$$= \frac{\lambda}{\mu} = \frac{1/12}{1/10} = \frac{5}{6}$$

$\therefore$  Percentage of customers who have to wait

$$= \frac{5}{6} \times 100 = 83.33$$

$$\begin{aligned}
 (h) \quad P(N > 3) &= P_4 + P_5 + P_6 + \dots \\
 &= 1 - \{P_0 + P_1 + P_2 + P_3\} \\
 &= 1 - \left(1 - \frac{\lambda}{\mu}\right) \left\{1 + \frac{\lambda}{\mu} + \left(\frac{\lambda}{\mu}\right)^2 + \left(\frac{\lambda}{\mu}\right)^3\right\} \\
 [\text{since } P_n &= \left(\frac{\lambda}{\mu}\right)^n \left(1 - \frac{\lambda}{\mu}\right), \text{ for } n \geq 0, \text{ for model I}] \\
 &= \left(\frac{\lambda}{\mu}\right)^4 = \left(\frac{5}{6}\right)^4 = 0.4823
 \end{aligned}$$

---

### Example 3

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At what average rate must a clerk in a supermarket work in order to ensure a probability of 0.90 that the customer will not wait longer than 12 min? It is assumed that there is only one counter at which customers arrive in a Poisson fashion at an average rate of 15 per hour and that the length of the service by the clerk has an exponential distribution.

$$\lambda = 15/\text{hour}; \mu = \mu_R/\text{hour}$$

$$P\left(W_q \leq \frac{1}{5}\right) = 0.90$$

$$\text{i.e.,} \quad P\left(W_q > \frac{1}{5}\right) = 0.10$$

$$\text{i.e.,} \quad \int_{0.2}^{\infty} g(w) dw = 0.10$$

$$\text{i.e.,} \quad \int_{0.2}^{\infty} \frac{\lambda}{\mu} (\mu - \lambda) e^{-(\mu - \lambda)w} dw = 0.10 \quad [\text{by formula (9) of model I}]$$

$$\text{i.e.,} \quad \left[ -\frac{\lambda}{\mu} e^{-(\mu - \lambda)w} \right]_{0.2}^{\infty} = 0.1$$

$$\text{i.e.,} \quad \frac{15}{\mu_R} e^{-(\mu_R - 15) \times 0.2} = 0.1$$

$$\text{i.e.,} \quad (15 - \mu_R) \times 0.2 = \log(0.1) - \log 15 + \log \mu_R$$

$$\text{i.e.,} \quad 0.2 \mu_R + \log \mu_R = 3 + \log 150 \approx 8 \quad (1)$$

Solving (1), we get  $\mu_R = 24$  approximately.

That is, the clerk must serve at the rate of 24 customer per hour.

---

### Example 4

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If people arrive to purchase cinema tickets at the average rate of 6 per minute, it takes an average of 7.5 seconds to purchase a ticket. If a person arrives 2 min before the picture starts and if it takes exactly 1.5 min to reach the correct seat after purchasing the ticket,

- (a) Can he expect to be seated for the start of the picture?
- (b) What is the probability that he will be seated for the start of the picture?
- (c) How early must he arrive in order to be 99% sure of being seated for the start of the picture?

$$\lambda = 6/\text{minute}; \mu = 8/\text{minute}$$

$$(a) E(W) = \frac{1}{\mu - \lambda} \quad [\text{by formula (7) of model I}]$$

$$= \frac{1}{8 - 6} = \frac{1}{2} \text{ min}$$

$\therefore E(\text{total time required to purchase the ticket and to reach the seat})$

$$= \frac{1}{2} + 1 \frac{1}{2} = 2 \text{ min}$$

Hence, he can just be seated for the start of the picture.

$$(b) P(\text{total time} < 2 \text{ min})$$

$$\begin{aligned} &= P\left(W < \frac{1}{2}\right) = 1 - P\left(W > \frac{1}{2}\right) \\ &= 1 - e^{-\mu} \left(1 - \frac{\lambda}{\mu}\right) \times \frac{1}{2} \quad [\text{by formula (8) of model I}] \\ &= 1 - e^{-1} = 0.63 \end{aligned}$$

$$(c) P(W < t) = 0.99$$

$$\text{i.e., } P(W > t) = 0.01$$

$$\text{i.e., } e^{-(\mu - \lambda)t} = 0.1$$

$$\text{i.e., } -2t = \log(0.1) = -2.3$$

$$\therefore t = 1.15 \text{ min}$$

$$\text{i.e., } P(\text{ticket purchasing time} < 1.15) = 0.99$$

$$\therefore P[\text{total time to get the ticket and to go to the seat} < (1.15 + 1.5)] = 0.99$$

Therefore, the person must arrive at least 2.65 min early so as to be 99% sure of seeing the start of the picture.

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### Example 5

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A duplicating machine maintained for office use is operated by an office assistant who earns Rs. 5 per hour. The time to complete each job varies according to an exponential distribution with mean 6 min. Assume a Poisson input with an average arrival rate of 5 jobs per hour. If an 8-h day is used as a base, determine

- the percentage idle time of the machine,
- the average time a job is in the system, and
- the average earning per day of the assistant.

$$\lambda = 5/\text{hour}; \mu = \frac{60}{6} = 10/\text{hour}$$

(a)  $P(\text{the machine is idle}) = P(N = 0) = P_0$

$$= 1 - \frac{\lambda}{\mu} \quad (\text{by the formula for } P_0 \text{ in model I})$$

$$= 1 - \frac{5}{10} = \frac{1}{2}$$

$\therefore$  Percentage of idle time of the machine = 50

(b)  $E(W) = \frac{1}{\mu - \lambda}$  [by formula (7) of model I]

$$= \frac{1}{10 - 5} = \frac{1}{5} \text{ h or } 12 \text{ min}$$

(c)  $E(\text{earning per day})$

$$= E(\text{number of jobs done/day}) \times \text{earning per job}$$

$$= E(\text{number of jobs done/day}) \times E(\text{time in hour/job}) \times \text{earning/hour}$$

$$= (8 \times 5) \times \frac{1}{5} \times 5 = \text{Rs. } 40.$$

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### Example 6

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The mean rate of arrival of planes at an airport during the peak period is 20 per hour, but the actual number of arrivals in any hour follows a Poisson distribution. The airport can land 60 planes per hour on an average in good weather or 30 planes per hour in bad weather, but the actual number landed in any hour follows a Poisson distribution with respective averages. When there is congestion, the planes are forced to fly over the field in the stack awaiting the landing of other planes that arrived earlier.

- How many planes would be flying over the field in the stack on an average in good weather and in bad weathers?

- (b) How long a plane would be in the stack and in the process of landing in good and bad weathers?
- (c) How much stack and landing time to allow so that priority to land out of order will have to be requested only 1 in 20 times.

$$\lambda = 20 \text{ per hour}$$

$$\mu = \begin{cases} 60 \text{ per hour in good weather} \\ 30 \text{ per hour in bad weather} \end{cases}$$

**Note** Landing time is service time; the planes flying over the field in the stack are assumed to form the queue.

- (a)  $E(N_q) = \text{Average number of planes flying over the field} = \frac{\lambda^2}{\mu(\mu - \lambda)}$  [by formula (2) of model I]

$$= \begin{cases} \frac{20^2}{60(60 - 20)}, & \text{in good weather} \\ \frac{20^2}{30(30 - 20)}, & \text{in bad weather} \end{cases}$$

$$= \begin{cases} \frac{1}{6}, & \text{in good weather} \\ \frac{4}{3}, & \text{in bad weather} \end{cases}$$

- (b)  $E(W) = \text{Average time for flying in the stack and for landing}$

$$= \frac{1}{\mu - \lambda} \text{ [by formula (7) of model I]}$$

$$= \begin{cases} \frac{1}{40} \text{ h or } 1.5 \text{ min in good weather} \\ \frac{1}{10} \text{ h or } 6 \text{ min in bad weather} \end{cases}$$

- (c) Let  $t_R$  be the maximum stack and landing time to be allowed, beyond which priority out of order is to be requested.

$$\text{Then } P(W > t_R) = \frac{1}{20}$$

$$\text{i.e., } e^{-\mu \left(1 - \frac{\lambda}{\mu}\right) t_R} = 0.05 \text{ [by formula (8) of model I]}$$

$$\text{i.e., } \begin{cases} e^{-40t_R} = 0.05, & \text{for good weather} \\ e^{-10t_R} = 0.05, & \text{for bad weather} \end{cases}$$

$$\text{i.e., } \begin{cases} t_R = 0.075 \text{ h or } 4.5 \text{ min for good weather} \\ = 0.299 \text{ h or } 18 \text{ min for bad weather} \end{cases}$$

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### Example 7

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There are three typists in an office. Each typist can type an average of 6 letters per hour. If letters arrive for being typed at the rate of 15 letters per hour,

- (a) What fraction of the time all the typists will be busy?
- (b) What is the average number of letters waiting to be typed?
- (c) What is the average time a letter has to spend for waiting and for being typed?
- (d) What is the probability that a letter will take longer than 20 min waiting to be typed and being typed?

$$\lambda = 15/\text{hour}; \mu = 6/\text{hour}; s = 3.$$

Hence, this is a problem in multiple server  $[(M/M/s): (\infty/\text{FIFO})]$  model, i.e., model II.

- (a)  $P(\text{all the typists are busy}) = P(N \geq 3)$

$$\begin{aligned} &= \frac{\left(\frac{\lambda}{\mu}\right)^3 \cdot P_0}{3! \left(1 - \frac{\lambda}{3\mu}\right)} && [\text{by formula (11) of model II}] \\ &= \frac{(2.5)^3 P_0}{6 \times \left(1 - \frac{2.5}{3}\right)} \end{aligned} \tag{1}$$

Now,  $P_0 = \frac{1}{\left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left(\frac{\lambda}{\mu}\right)^n \right\} + \left\{ \frac{1}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \cdot \left(\frac{\lambda}{\mu}\right)^s \right\}}$

[by formula (6) of model II]

$$= \frac{1}{\left\{ 1 + 2.5 + \frac{1}{2} \times (2.5)^2 \right\} + \left\{ \frac{1}{6 \times \left(1 - \frac{5}{6}\right)} \times (2.5)^3 \right\}}$$

$$= \frac{1}{22.25} = 0.0449 \quad (2)$$

Using (2) in (1), we have  $P(N \geq 3) = 0.7016$ .

Hence, the fraction of the time all the typists will be busy = 0.7016.

$$(b) \quad E(N_q) = \frac{1}{s.s!} \cdot \frac{\left(\frac{\lambda}{\mu}\right)^{s+1} \cdot P_0}{\left(1 - \frac{\lambda}{\mu s}\right)^2} \quad [\text{by formula (7) of model II}]$$

$$= \frac{1}{3 \times 6} \times \frac{(2.5)^4}{\left(1 - \frac{2.5}{3}\right)^2} \times 0.0449 = 3.5078$$

$$(c) \quad E(W) = \frac{1}{\lambda} E(N) \quad [\text{by Little's formula (i)}]$$

$$= \frac{1}{\lambda} \left\{ E(N_q) + \frac{\lambda}{\mu} \right\}, \quad [\text{by Little's formula (iv)}]$$

$$= \frac{1}{15} \{ 3.5078 + 2.5 \} = 0.4005 \text{ h}$$

or  $= 24 \text{ min, nearly}$

$$(d) \quad P(W > t) = e^{-\mu t} \left\{ 1 + \frac{\left(\frac{\lambda}{\mu}\right)^s \left[ 1 - e^{-\mu t \left(s-1-\frac{\lambda}{\mu}\right)} \right] P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right) \left(s-1-\frac{\lambda}{\mu}\right)} \right\}$$

(This formula has not been derived; it may be assumed.)

$$\begin{aligned} \therefore P\left(W > \frac{1}{3}\right) &= e^{-6 \times \frac{1}{3}} = \left[ 1 + \frac{(2.5)^3 \{1 - e^{(-2 \times -0.5)}\} \times 0.0449}{6 \left(1 - \frac{2.5}{3}\right) (-0.5)} \right] \\ &= e^{-2} \left[ 1 + \frac{0.7016 (1 - e)}{(-0.5)} \right] \\ &= 0.4616 \end{aligned}$$

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## Example 8

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Given an average arrival rate of 20 per hour, is it better for a customer to get service at a single channel with mean service rate of 22 customers per hour or at one of two channels in parallel with mean service rate of 11 customers per hour for each of the two channels. Assume both queues to be of Poisson type.

For the single-channel service,

$$\lambda = 20/\text{hour} \text{ and } \mu = 22/\text{hour}.$$

$$\begin{aligned} E(W) &= \frac{1}{\mu - \lambda} && [\text{by formula (7) of model I}] \\ &= \frac{1}{2} \text{ h} \end{aligned}$$

For the two-channel service,

$$\lambda = 20/\text{hour} \text{ and } \mu = 11/\text{hour}.$$

$$\begin{aligned} E(W) &= \frac{1}{\mu} + \frac{1}{\mu} \cdot \frac{1}{s \cdot s!} \cdot \frac{\left(\frac{\lambda}{\mu}\right)^s}{\left(1 - \frac{\lambda}{\mu s}\right)^2} \times P_0 \\ &\quad [\text{by formula (9) of model II}] \end{aligned}$$

$$\begin{aligned} &= \frac{1}{11} + \frac{1}{11 \times 2 \times 2} \times \frac{\left(\frac{20}{11}\right)^2}{\left(1 - \frac{20}{22}\right)^2} \times P_0 \\ &= 0.0909 + 9.0909 \times P_0 \end{aligned} \tag{1}$$

$$\begin{aligned} \text{Now, } P_0^{-1} &= \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \cdot \left(\frac{\lambda}{\mu}\right)^n \right\} + \left\{ \frac{1}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \cdot \left(\frac{\lambda}{\mu}\right)^s \right\} \\ &\quad [\text{by formula (6) of model II}] \end{aligned}$$

$$\begin{aligned} &= 1 + \frac{20}{11} + \frac{1}{2 \times \frac{1}{11}} \times \left(\frac{20}{11}\right)^2 \\ &= 21 \end{aligned}$$

$$\therefore P_0 = 0.0476 \tag{2}$$

Using (2) in (1), we have

$$E(W) = 0.5236 \text{ h}$$

As the average waiting time in single channel service is less than that in two channel service, the customer has to prefer the former.

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### Example 9

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A telephone company is planning to install telephone booths in a new airport. It has established the policy that a person should not have to wait more than 10% of the times he tries to use a phone. The demand for use is estimated to be Poisson with an average of 30 per hour. The average phone call has an exponential distribution with a mean time of 5 min. How many phone booths should be installed?

$$\lambda = 30/\text{hour} \text{ and } \mu = 12/\text{hour}$$

In order that infinite queue may not build up, the traffic intensity  $\frac{\lambda}{\mu s} < 1$ , for multiserver model.

$$\text{i.e., } s > \frac{\lambda}{\mu}$$

$$\text{i.e., } s > \frac{30}{12} (= 2.5)$$

Therefore, the telephone company must install at least 3 booths.

Now, we have to find the number  $s$  of telephone booths such that

$$P(W > 0) \leq 0.10 \text{ or equivalently}$$

$$P(N \geq s) \leq 0.10$$

i.e., we have to find  $s$  such that

$$\frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \leq 0.10 \quad [\text{by formula (11) model II}]$$

This inequation is not easily solvable. Hence, we proceed by trials and find out the least value of  $s$  that satisfies this inequation.

Let  $s = 3$ :

$$\text{Then } P(W > 0) = \frac{(2.5)^3 \cdot P_0}{6 \left(1 - \frac{2.5}{3}\right)} = 15.625 P_0,$$

$$\text{where } P_0 = \left[ \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left(\frac{\lambda}{\mu}\right)^n \right\} + \frac{\left(\frac{\lambda}{\mu}\right)^s}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \right]^{-1}$$

[by formula (6) of model II]

i.e., 
$$P_0 = \left[ \left\{ 1 + 2.5 + \frac{1}{2} \times (2.5)^2 \right\} + \frac{(2.5)^3}{6 \times \left( 1 - \frac{2.5}{3} \right)} \right]^{-1}$$
  

$$= (22.25)^{-1} = 0.0449$$

$\therefore P(W > 0) = 15.625 \times 0.0449 = 0.7022 < 0.10$

Let  $s = 4$ :

Then 
$$P(W > 0) = \frac{(2.5)^4 \cdot P_0}{24 \left( 1 - \frac{2.5}{4} \right)} = 4.3403 P_0$$

where 
$$P_0 = \left[ \left\{ 1 + 2.5 + \frac{1}{2} \times (2.5)^2 + \frac{1}{6} \times (2.5)^3 \right\} + \frac{(2.5)^4}{24 \left( 1 - \frac{2.5}{4} \right)} \right]^{-1}$$
  

$$= 0.0737$$

$\therefore P(W > 0) = 4.3403 \times 0.0737 = 0.3199 < 0.10$

Similarly, when  $s = 5$ ,  $P(W > 0) = 0.1304 < 0.10$ .

When  $s = 6$ ,  $P(W > 0) = 0.047 < 0.10$ .

Hence the number of booths to be installed = 6.

### Example 10

A bank has two tellers working on savings accounts. The first teller handles withdrawals only. The second teller handles deposits only. It has been found that the service time distributions for both deposits and withdrawals are exponential with a mean service time of 3 min per customer. Depositors are found to arrive in a Poisson fashion throughout the day with mean arrival rate of 16 per hour. Withdrawers also arrive in a Poisson fashion with a mean arrival rate of 14 per hour. What would be the effect on the average waiting time for the customers if each teller could handle both withdrawals and deposits. What would be the effect, if this could only be accomplished by increasing the service time to 3.5 min?

When there is a separate channel for the depositors,  $\lambda_1 = 16/\text{hour}$ ,  $\mu = 20/\text{hour}$

$$\begin{aligned} \therefore E(W_q \text{ for depositors}) &= \frac{\lambda_1}{\mu(\mu - \lambda_1)} && [\text{by formula (10) of model I}] \\ &= \frac{16}{20(20-16)} = \frac{1}{5} \text{ h or } 12 \text{ min} \end{aligned}$$

When there is a separate channel for the withdrawers,  $\lambda_2 = 14/\text{hour}$ ,  $\mu = 20/\text{hour}$ .

$$\begin{aligned}\therefore E(W_q \text{ for withdrawers}) &= \frac{\lambda_2}{\mu(\mu - \lambda_2)} \\ &= \frac{14}{20(20-14)} = \frac{7}{60} \text{ h or } 7 \text{ min}\end{aligned}$$

If both tellers do both services,

$$s = 2, \mu = 20/\text{hour}, \lambda = \lambda_1 + \lambda_2 = 30/\text{hour}$$

$$\begin{aligned}\therefore E(W_q \text{ for any customer}) &= \frac{1}{\mu} \cdot \frac{1}{s.s!} \cdot \frac{\left(\frac{\lambda}{\mu}\right)^s \times P_0}{\left(1 - \frac{\lambda}{\mu s}\right)^2}, \\ &\quad [\text{by formula (10) of model II}] \\ &= \frac{1}{20} \times \frac{1}{2 \times 2} \times \frac{(1.5)^2}{(1 - 0.75)^2} \times P_0 \\ &= 0.45 \times P_0\end{aligned}\tag{1}$$

$$\begin{aligned}\text{Now } P_0 &= \left[ \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \cdot \left( \frac{\lambda}{\mu} \right)^n \right\} + \frac{1}{s! \left( 1 - \frac{\lambda}{\mu s} \right)} \left( \frac{\lambda}{\mu} \right)^s \right]^{-1} \\ &\quad [\text{by formula (6) of, model II}] \\ &= \left[ 1 + 1.5 + \frac{(1.5)^2}{2 \times 0.25} \right]^{-1} \\ &= \frac{1}{7}\end{aligned}\tag{2}$$

Using (2) in (1),

$$E(W_q \text{ for any customer}) = 0.45 \times \frac{1}{7} \text{ h or } 3.86 \text{ min}$$

Hence, if both tellers do both types of service, the customers get benefited as their waiting time is considerably reduced.

Now, if both tellers do both types of service but with increased service time,

$$s = 2, \lambda = 30, \mu = \frac{60}{3.5} = \frac{120}{7} \text{ per hour.}$$

$E(W_q$  of any customer)

$$= \frac{7}{120} \times \frac{1}{2 \times 2} \times \frac{(1.75)^2}{\left(1 - \frac{7}{8}\right)^2} \times P_0 = 2.86 P_0, \text{ where}$$

$$P_0 = \left[ 1 + 1.75 + \frac{(1.75)^2}{2 \times \frac{1}{8}} \right]^{-1} = \frac{1}{15}$$

$$\therefore E(W_q \text{ of any customer}) = \frac{2.86}{15} \text{ h or } 11.44 \text{ min}$$

If this arrangement is adopted, withdrawers stand to lose as their waiting time is increased considerably and depositors get slightly benefited.

### Example 11

A supermarket has two girls attending to sales at the counters. If the service time for each customer is exponential with mean 4 min and if people arrive in Poisson fashion at the rate of 10 per hour,

- (a) what is the probability that a customer has to wait for service?
- (b) what is the expected percentage of idle time for each girl?
- (c) if the customer has to wait in the queue, what is the expected length of his waiting time?

$$s = 2, \lambda = \frac{1}{6} \text{ per minute}, \mu = \frac{1}{4} \text{ per minute}$$

- (a)  $P(\text{a customer has to wait for service})$

$$= P(N \geq 2) = 1 - P_0 - P_1 \quad (1)$$

$$P_0 = \left[ \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n \right\} + \left\{ \frac{\left( \frac{\lambda}{\mu} \right)^s}{s! \left( 1 - \frac{\lambda}{\mu s} \right)} \right\} \right]^{-1}$$

[by formula (6) of model II]

$$= \left[ 1 + \frac{2}{3} + \frac{\left( \frac{2}{3} \right)^2}{2 \times \left( 1 - \frac{1}{3} \right)} \right]^{-1} = \frac{1}{2} \quad (2)$$

$$P_1 = \frac{\lambda}{\mu} \cdot P_0, \quad \text{[by formula (4) of model II]}$$

$$= \frac{2}{3} \times \frac{1}{2} = \frac{1}{3} \quad (3)$$

Using (2) and (3) in (1), we have

$$P(N \geq 2) = 1 - \frac{1}{2} - \frac{1}{3} = \frac{1}{6}$$

(b) Fraction of time when the girls are busy =  $\frac{\lambda}{\mu s} = \frac{1}{3}$

$$\therefore \text{Fraction of time when the girls are idle} = \frac{2}{3}$$

$$\therefore \text{Expected percentage of idle time for each girl} = \frac{2}{3} \times 100 = 67$$

(c)  $E(W_q/W_s > 0) = \frac{1}{\mu s - \lambda}$  [by formula (13) of model II]

$$= \frac{1}{\frac{1}{4} \times 2 - \frac{1}{6}} = 3 \text{ min}$$

### Example 12

A petrol pump station has 4 pumps. The service times follow the exponential distribution with a mean of 6 min and cars arrive for service in a Poisson process at the rate of 30 cars per hour.

- (a) What is the probability that an arrival would have to wait in line?
- (b) Find the average waiting time, average time spent in the system and the average number of cars in the system.
- (c) For what percentage of time would a pump be idle on an average?

$$s = 4, \lambda = 30/\text{hour}, \mu = 10/\text{hour}$$

(a)  $P(\text{an arrival has to wait}) = P(W > 0)$

$$= \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \quad [\text{by formula (11) of model II}]$$

$$= \frac{3^4 \times P_0}{24 \times \left(1 - \frac{3}{4}\right)} = 13.5 \times P_0 \quad (1)$$

$$P_0 = \left[ \left\{ \sum_{n=0}^{s-1} \frac{1}{n!} \left(\frac{\lambda}{\mu}\right)^n \right\} + \frac{\left(\frac{\lambda}{\mu}\right)^s}{s! \left(1 - \frac{\lambda}{\mu s}\right)} \right]^{-1}$$

[by formula (6) of model II]

$$\begin{aligned}
 &= \left[ \left( 1 + 3 + \frac{1}{2} \times 9 + \frac{1}{6} \times 27 \right) + \frac{3^4}{24 \times \left( 1 - \frac{3}{4} \right)} \right]^{-1} \\
 &= 0.0377
 \end{aligned} \tag{2}$$

Using (2) in (1),  $P(W > 0) = 0.5090$

$$\begin{aligned}
 (b) \quad E(W_q) &= \frac{1}{\mu} \cdot \frac{1}{s \times s!} \times \frac{\left( \frac{\lambda}{\mu} \right)^s}{\left( 1 - \frac{\lambda}{\mu s} \right)^2} \times P_0 \text{ [by formula (10) of model II]} \\
 &= \frac{1}{10 \times 4 \times 24} \times \frac{3^4}{\left( 1 - \frac{3}{4} \right)^2} \times 0.0377 = 0.0509 \text{ h}
 \end{aligned}$$

or 3.05 min

$$\begin{aligned}
 E(W_s) &= \frac{1}{\mu} + E(W_q) \quad \text{[by formulas (9) and (10) of model II]} \\
 &= 6 + 3.05 = 9.05 \text{ min}
 \end{aligned}$$

$$\begin{aligned}
 E(N) &= \frac{1}{s \times s!} \times \frac{\left( \frac{\lambda}{\mu} \right)^{s+1}}{\left( 1 - \frac{\lambda}{\mu s} \right)^2} \times P_0 + \frac{\lambda}{\mu} \text{ [by formula (8) of model II]} \\
 &= \frac{1}{4 \times 24} \times \frac{3^5}{\left( 1 - \frac{3}{4} \right)^2} \times 0.0377 + 3 \\
 &= 4.53 \text{ cars}
 \end{aligned}$$

$$\begin{aligned}
 (c) \quad \text{The fraction of time when the pumps are busy} &= \text{traffic intensity} = \frac{\lambda}{\mu s} \\
 &= \frac{3}{4} \\
 \therefore \text{The fraction of time when the pumps are idle} &= \frac{1}{4}
 \end{aligned}$$

Therefore, required percentage = 25%

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### Example 13

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In a single server queueing system with Poisson input and exponential service times, if the mean arrival rate is 3 calling units per hour, the expected service time

is 0.25 h and the maximum possible number of calling units in the system is 2, find  $P_n$  ( $n \geq 0$ ), average number of calling units in the system and in the queue and average waiting time in the system and in the queue.

The situation in this problem is one of finite capacity, single server Poisson queue models.

$$\lambda = 3, \mu = 4 \text{ and } k = 2$$

$$\text{As } \lambda \neq \mu, \quad P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \quad [\text{by formula (4) of model III}]$$

$$= \frac{1 - \frac{3}{4}}{1 - \left(\frac{3}{4}\right)^3} = \frac{16}{37} = 0.4324$$

$$\text{Since } \lambda \neq \mu, \quad P_n = \left(\frac{\lambda}{\mu}\right)^n \left[ \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \right] \quad [\text{by formula (6) of model III}]$$

$$= (0.4324) (0.75)^n$$

$$E(N) = \frac{\lambda}{\mu - \lambda} - \frac{(k+1)\left(\frac{\lambda}{\mu}\right)^{k+1}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \quad [\text{by formula (8) of model IV}]$$

$$= \frac{3}{4-3} - \frac{3 \times \left(\frac{3}{4}\right)^3}{1 - \left(\frac{3}{4}\right)^3} = 3 - \frac{81}{37} = \frac{30}{37} \approx 0.8 \text{ calling unit}$$

$$E(N_q) = E(N) - (1 - P_0) \quad [\text{by formula (10) of model III}]$$

$$= \frac{30}{37} - \left(1 - \frac{16}{37}\right) = \frac{9}{37} = 0.24 \text{ calling unit}$$

$$E(W_s) = \frac{1}{\lambda'} E(N) \quad [\text{by formula (13) of model III}]$$

where  $\lambda' = \mu (1 - P_0)$ ,  $[\text{by formula (11) of model III}]$

$$\begin{aligned}
 &= 4 \left( 1 - \frac{16}{37} \right) = \frac{84}{37} \\
 \therefore E(W_s) &= \frac{37}{84} \times \frac{30}{37} = \frac{5}{14} \text{ h or } 21.4 \text{ min} \\
 E(W_q) &= \frac{1}{\lambda'} E(N_q) \quad [\text{by formula (14) of model III}] \\
 &= \frac{37}{84} \times \frac{9}{37} \\
 &= \frac{3}{28} \text{ h or } 6.4 \text{ min}
 \end{aligned}$$

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### Example 14

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The local one-person barber shop can accommodate a maximum of 5 people at a time (4 waiting and 1 getting hair-cut). Customers arrive according to a Poisson distribution with mean 5 per hour. The barber cuts hair at an average rate of 4 per hour (Exponential service time).

- (a) What percentage of time is the barber idle?
- (b) What fraction of the potential customers are turned away?
- (c) What is the expected number of customers waiting for a hair-cut?
- (d) How much time can a customer expect to spend in the barber shop?

$$\lambda = 5, \mu = 4, k = 5$$

- (a)  $P(\text{the barber is idle}) = P(N = 0)$

$$\begin{aligned}
 &= P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left( \frac{\lambda}{\mu} \right)^{k+1}} \quad [\text{by formula (4) of model III}] \\
 &= \frac{1 - \frac{5}{4}}{1 - \left( \frac{5}{4} \right)^6} = 0.0888
 \end{aligned}$$

$\therefore$  Percentage of time when the barber is idle  $\approx 9\%$ .

- (b)  $P(\text{a customer is turned away}) = P(N > 5)$

$$= \left( \frac{\lambda}{\mu} \right)^5 \cdot \left[ \frac{1 - \frac{\lambda}{\mu}}{1 - \left( \frac{\lambda}{\mu} \right)^{k+1}} \right] \quad [\text{by formula (6) of model III}]$$

$$= \left(\frac{5}{4}\right)^5 \cdot \left[ \frac{1 - \frac{5}{4}}{1 - \left(\frac{5}{4}\right)^6} \right]$$
$$= \frac{3125}{11529} = 0.2711$$

Therefore,  $0.2711 \times$  potential customers are turned away.

(c)  $E(N_q) = E(N) - (1 - P_0)$

$$= \frac{\lambda}{\mu - \lambda} - \frac{(k+1)\left(\frac{\lambda}{\mu}\right)^{k+1}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} - (1 - P_0),$$

[by formulas (6) and (10) of model III]

$$= \left\{ -5 - \frac{6 \times \left(\frac{5}{4}\right)^6}{1 - \left(\frac{5}{4}\right)^6} \right\} - (1 - 0.0888)$$

$$= \frac{6 \times \frac{15625}{4096}}{\frac{11529}{4096}} - 5.9112 \approx 2.2 \text{ customers}$$

(d)  $E(W) = \frac{1}{\lambda}, E(N)$  [by formula (13) of model III]

$$= \frac{1}{\mu(1 - P_0)} \times E(N)$$

$$= \frac{3.1317}{3.6448} \approx 0.8592 \text{ h}$$

or  $51.5 \text{ min}$

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### Example 15

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At a railway station, only one train is handled at a time. The railway yard is sufficient only for 2 trains to wait, while the other is given signal to leave the station. Trains arrive at the station at an average rate of 6 per hour and the railway station can handle them on an average of 6 per hour. Assuming Poisson arrivals and exponential service distribution, find the probabilities for the numbers of trains in the system. Also find the average waiting time of a new train coming

into the yard. If the handling rate is doubled, how will the above results get modified?

$$(a) \quad \lambda = 6 \text{ per hour}, \mu = 6 \text{ per hour}, k = 2 + 1 = 3$$

$$\text{Since } \lambda = \mu, P_0 = \frac{1}{k+1}$$

$$= \frac{1}{4} \quad [\text{by formula (5) of model III}]$$

$$P_n = \frac{1}{k+1} = \frac{1}{4} \text{ for } n = 1, 2, 3 \quad [\text{by formula (7) of model III}]$$

$$E(N) = \frac{k}{2} \quad [\text{by formula (9) of model III}] \\ = 1.5 \text{ trains}$$

$$E(W) = \frac{1}{\lambda'} E(N) \quad [\text{by formula (13) of model III}]$$

$$= \frac{1.5}{\mu(1 - P_0)}$$

$$= \frac{1.5}{6 \times \frac{3}{4}} = \frac{1}{3} \text{ h or } 20 \text{ min}$$

$$(b) \quad \lambda = 6; \mu = 12, k = 3$$

$$\text{Since } \lambda \neq \mu, P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \quad [\text{by formula (4) of model III}]$$

$$= \frac{1 - \frac{1}{2}}{1 - \left(\frac{1}{2}\right)^4} = \frac{8}{15}$$

$$P_n = \left(\frac{\lambda}{\mu}\right)^n \left\{ \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \right\} \quad [\text{by formula (6) of model III}]$$

$$= \frac{8}{15} \cdot \left(\frac{1}{2}\right)^n, \text{ for } n = 1, 2, 3.$$

$$E(N) = \frac{\lambda}{\mu - \lambda} - \frac{(k+1) \left( \frac{\lambda}{\mu} \right)^{k+1}}{1 - \left( \frac{\lambda}{\mu} \right)^{k+1}} \quad [\text{by formula (8) of model III}]$$

$$= 1 - \frac{4 \times \left( \frac{1}{2} \right)^4}{1 - \left( \frac{1}{2} \right)^4} = 1 - \frac{4}{15} = \frac{11}{15} \approx 0.73 \text{ train}$$

$$E(W) = \frac{1}{\lambda'} E(N) \quad [\text{by formula (13) of model III}]$$

$$= \frac{1}{\mu(1 - P_0)} \times E(N)$$

$$= \frac{\frac{11}{15}}{12 \left( 1 - \frac{8}{15} \right)} = \frac{11}{84} \text{ h or } 7.9 \text{ min}$$

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### Example 16

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Patients arrive at a clinic according to Poisson distribution at a rate of 30 patients per hour. The waiting room does not accommodate more than 14 patients. Examination time per patient is exponential with mean rate of 20 per hour.

- (a) Find the effective arrival rate at the clinic.
- (b) What is the probability that an arriving patient will not wait?
- (c) What is the expected waiting time until a patient is discharged from the clinic?
- (a)  $\lambda = 30 \text{ per hour}$ ,  $\mu = 20 \text{ per hour}$ ,  $k = 14 + 1 = 15$

$$\text{Since } \lambda \neq \mu, P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left( \frac{\lambda}{\mu} \right)^{k+1}} \quad [\text{by formula (4) of model III}]$$

$$= \frac{1 - \frac{3}{2}}{1 - \left( \frac{3}{2} \right)^{16}} = 0.00076$$

Effective arrival rate  $\lambda' = \mu(1 - P_0)$  [by formula (11) of model III]

$$= 20 \times (1 - 0.00076)$$

$$= 19.98 \text{ per hour}$$

(b)  $P(\text{a patient will not wait})$

$$= P_0 = 0.00076$$

$$(c) E(N) = \frac{\lambda}{\mu - \lambda} - \frac{(k+1) \left( \frac{\lambda}{\mu} \right)^{k+1}}{1 - \left( \frac{\lambda}{\mu} \right)^{k+1}}$$

$$= -3 - \frac{16 \times \left( \frac{3}{2} \right)^{16}}{1 - \left( \frac{3}{2} \right)^{16}} = 13 \text{ patients nearly}$$

$$E(W) = \frac{E(N)}{\lambda'} = \frac{13}{19.98} = 0.65 \text{ h or } 39 \text{ min}$$

### Example 17

A 2-person barber shop has 5 chairs to accommodate waiting customers. Potential customers, who arrive when all 5 chairs are full, leave without entering barber shop. Customers arrive at the average rate of 4 per hour and spend an average of 12 min in the barber's chair. Compute  $P_0$ ,  $P_1$ ,  $P_7$ ,  $E(N_q)$  and  $E(W)$ .

The situation in this problem is one of finite capacity, multiserver Poisson queue models.

$$\lambda = 4 \text{ per hour}, \mu = 5 \text{ per hour}, s = 2, k = 2 + 5 = 7$$

$$(a) P_0 = \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \sum_{n=s}^k \left( \frac{\lambda}{\mu s} \right)^{n-s} \right]^{-1}$$

[by formula (3) of model IV]

$$= \left[ \sum_{n=0}^1 \frac{1}{n!} \left( \frac{4}{5} \right)^n + \frac{1}{2} \cdot \left( \frac{4}{5} \right)^2 \sum_{n=2}^7 \left( \frac{2}{5} \right)^{n-2} \right]^{-1}$$

$$= \left[ 1 + \frac{4}{5} + \frac{8}{25} \left\{ 1 + \frac{2}{5} + \left( \frac{2}{5} \right)^2 + \left( \frac{2}{5} \right)^3 + \left( \frac{2}{5} \right)^4 + \left( \frac{2}{5} \right)^5 \right\} \right]^{-1}$$

$$= \left[ \frac{9}{5} + \frac{8}{25} \left\{ \frac{1 - (0.4)^7}{1 - 0.4} \right\} \right]^{-1} = 0.4287$$

$$(b) \quad P_n = \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n P_0, \text{ for } n \leq s \quad [\text{by formula (4) of model IV}]$$

$$\therefore \quad P_1 = \left( \frac{4}{5} \right) \times 0.4287 = 0.3430$$

$$(c) \quad P_n = \frac{1}{s! \cdot s^{n-s}} \left( \frac{\lambda}{\mu} \right)^n \cdot P_0, \text{ for } s < n \leq k \quad [\text{by formula (4) of model IV}]$$

$$\begin{aligned} \therefore \quad P_7 &= \frac{1}{2 \times 2^{7-2}} \times \left( \frac{4}{5} \right)^7 \times 0.4287 \\ &= 0.0014 \end{aligned}$$

$$(d) \quad E(N_q) = P_0 \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{\rho}{s!(1-\rho)^2} [1 - p^{k-s} - (k-s)(1-\rho)\rho^{k-s}],$$

$$\text{where } \rho = \frac{\lambda}{\mu s} \quad [\text{by formula (5) of model IV}]$$

$$\begin{aligned} &= (0.4287) \cdot (0.8)^2 \cdot \frac{(0.4)}{2 \times (0.6)^2} [1 - (0.4)^5 - 5 \times 0.6 \times (0.4)^5] \\ &= 0.15 \text{ customer} \end{aligned}$$

$$(e) \quad E(N) = E(N_q) + s - \sum_{n=0}^{s-1} (s-n) P_n \quad [\text{by formula (6) of model IV}]$$

$$= 0.1462 + 2 - \sum_{n=0}^1 (2-n) P_n$$

$$= 2.1462 - (2 \times P_0 + 1 \times P_1)$$

$$= 2.1462 - (2 \times 0.4287 + 1 \times 0.3430)$$

$$= 0.95 \text{ customer}$$

$$E(W) = \frac{1}{\lambda'} \cdot E(N) \quad [\text{by formula (9) of model IV}]$$

$$\text{where } \lambda' = \mu \left[ s - \sum_{n=0}^{s-1} (s-n) P_n \right] \quad [\text{by formula (7) of model IV}]$$

$$= 4[2 - (2 \times 0.4287 + 1 \times 0.3430)]$$

$$= 3.1984$$

$$\therefore E(W) = \frac{0.9458}{3.1984} = 0.2957 \text{ h or } 17.7 \text{ min}$$

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**Example 18**

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At a port there are 6 unloading berths and 4 unloading crews. When all the berths are full, arriving ships are diverted to an overflow facility 20 kms down the river. Tankers arrive according to a Poisson process with a mean of 1 every 2 h. It takes for an unloading crew, on the average, 10 h to unload a tanker, the unloading time following an exponential distribution. Find

- how many tankers are at the port on the average?
- how long does a tanker spend at the port on the average?
- what is the average arrival rate at the overflow facility?

$$\lambda = \frac{1}{2} \text{ per hour}, \mu = \frac{1}{10} \text{ per hour}, s = 4, k = 6$$

$$(a) P_0 = \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \sum_{n=s}^k \left( \frac{\lambda}{\mu s} \right)^{n-s} \right]^{-1}$$

[by formula (3) of model IV]

$$= \left[ \left( 1 + 5 + \frac{1}{2} \times 5^2 + \frac{1}{6} \times 5^3 \right) + \frac{1}{24} \times 5^4 \times \left\{ \left( \frac{5}{4} \right)^0 + \left( \frac{5}{4} \right)^1 + \left( \frac{5}{4} \right)^2 \right\} \right]^{-1}$$

$$= 0.0072$$

$$E(N_q) = P_0 \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{\rho}{s!(1-\rho)^2} [1 - \rho^{k-s} - (k-s)(1-\rho)\rho^{k-s}],$$

$$\text{where } \rho = \frac{\lambda}{\mu s} \text{ [by formula (5) of model IV]}$$

$$= 0.0072 \times 5^4 \times \frac{1.25}{24 \times (.25)^2} [1 - (1.25)^2 - 2 \times (-.25)(1.25)^2]$$

$$= 0.8203 \text{ tanker}$$

$$E(N) = E(N_q) + s - \sum_{n=0}^{s-1} (s-n) P_n \quad \text{[by formula (6) of model IV]}$$

$$= 4.8203 - (4P_0 + 3P_1 + 2P_2 + P_3)$$

$$= 4.8203 - \{4 \times 0.0072 + 3 \times 0.0360 + 2 \times 0.09 + 0.15\}$$

$$= 4.3535 \text{ tankers}$$

(b)  $E(W) = \frac{1}{\lambda'} E(N)$  [by formula (9) of model IV]

where  $\lambda' = \mu \left[ s - \sum_{n=0}^{s-1} (s-n) P_n \right]$  [by formula (7) of model IV]

$$= \frac{1}{10} [4 - \{4P_0 + 3P_1 + 2P_2 + P_3\}]$$

$$= \frac{1}{10} [4 - 0.4668] = 0.3533$$

$$\therefore E(W) = \frac{4.3535}{0.3533} = 12.32 \text{ h}$$

- (c) When  $N = 6$ , i.e., when the number of tankers in the port is 6, overflow occurs.

$$P(N=6) = \frac{1}{s!s^{n-s}} \left( \frac{\lambda}{\mu} \right)^n P_0, \text{ for } n = k \text{ [by formula (4) of model IV]}$$

$$= \frac{1}{24 \times 4^2} \times 5^6 \times 0.0072$$

$$= 0.2930$$

Average arrival rate at the overflow facility = (Average arrival rate at the port)  $\times$  (Probability that overflow occurs)

$$= \frac{1}{2} \times 0.2930 = 0.586 \text{ per hour}$$

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### Example 19

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A car servicing station has 2 bays where service can be offered simultaneously. Because of space limitation, only 4 cars are accepted for servicing. The arrival pattern is Poisson with 12 cars per day. The service time in both the bays is exponentially distributed with  $\mu = 8$  cars per day per bay. Find the average number of cars in the service station, the average number of cars waiting for service and the average time a car spends in the system.

$$\lambda = 12 \text{ per day}, \mu = 8 \text{ per day}, s = 2, k = 4$$

(a)  $P_0 = \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \sum_{n=s}^k \left( \frac{\lambda}{\mu s} \right)^{n-s} \right]^{-1}$

[by formula (3) of model IV]

$$= \left[ 1 + \frac{1.5}{1} + \frac{1}{2} \times (1.5)^2 \{ 1 + (.75) + (.75)^2 \} \right]^{-1}$$

$$= 0.1960$$

$$E(N_q) = P_0 \left( \frac{\lambda}{\mu} \right)^s \cdot \frac{\rho}{s!(1-\rho)^2} [1 - \rho^{k-s} - (k-s)(1-\rho)\rho^{k-s}],$$

where  $\rho = \frac{\lambda}{\mu s}$  [by formula (5) of model IV]

$$\text{i.e., } E(N_q) = 0.1960 \times (1.5)^2 \times \frac{0.75}{2 \times (0.25)^2} \times [1 - (0.75)^2 - 2 \times 0.25 \times (0.75)^2]$$

$$= 0.4134 \text{ car}$$

(b)  $E(N)$  = Average number of cars in the service station

$$= E(N_q) + s - \sum_{n=0}^{s-1} (s-n) P_n \text{ [by formula (6) of model IV]}$$

$$= 0.4134 + 2 - \sum_{n=0}^1 (2-n) P_n$$

$$= 2.4134 - (2P_0 + P_1)$$

$$= 2.4134 - (2 \times 0.1960 + 1.5 \times 0.1960)$$

$$= 1.73 \text{ cars}$$

(c)  $E(W) = \frac{1}{\lambda'} E(N)$  [by formula (9) of model IV]

where  $\lambda' = \mu \left[ s - \sum_{n=0}^{s-1} (s-n) P_n \right]$  [by formula (7) of model IV]

$$= 8[2 - (2P_0 + P_1)]$$

$$= 10.512$$

$$\therefore E(W) = \frac{1.73}{10.512} = 0.1646 \text{ day}$$

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### Example 20

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A group of engineers has 2 terminals available to aid in their calculations. The average computing job requires 20 min of terminal time and each engineer requires some computation about once every half an hour. Assume that these are distributed according to an exponential distribution. If there are 6 engineers in the group, find

- (a) the expected number of engineers waiting to use one of the terminals and in the computing centre, and  
 (b) the total time lost per day.

$$\lambda = 2 \text{ per hour}, \mu = 3 \text{ per hour}, s = 2, k = 6$$

$$(a) P_0 = \left[ \sum_{n=0}^{s-1} \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n + \frac{1}{s!} \left( \frac{\lambda}{\mu} \right)^s \sum_{n=s}^k \left( \frac{\lambda}{\mu s} \right)^{n-s} \right]^{-1}$$

[by formula (3) of model IV]

$$= \left[ 1 + \frac{2}{3} + \frac{1}{2} \times \left( \frac{2}{3} \right)^2 \left\{ 1 + \frac{1}{3} + \left( \frac{1}{3} \right)^2 + \left( \frac{1}{3} \right)^3 + \left( \frac{1}{3} \right)^4 \right\} \right]^{-1}$$

$$= 0.5003$$

$$E(N_q) = P_0 \times \left( \frac{\lambda}{\mu} \right)^s \times \frac{\rho}{s!(1-\rho)^2} [1 - \rho^{k-s} - (k-s)(1-\rho)\rho^{k-s}],$$

$$\text{where } \rho = \frac{\lambda}{\mu s} \quad \quad \quad \text{[by formula (5) of model IV]}$$

$$\text{i.e., } E(N_q) = 0.5003 \times \left( \frac{2}{3} \right)^2 \times \frac{\left( \frac{1}{3} \right)}{2 \times \left( \frac{2}{3} \right)^2} \left[ 1 - \left( \frac{1}{3} \right)^4 - 4 \times \frac{2}{3} \times \left( \frac{1}{3} \right)^4 \right]$$

$$= 0.0796$$

$$E(N) = E(N_q) + s - \sum_{n=0}^{s-1} (s-n)P_n \quad \quad \quad \text{[by formula (6) of model IV]}$$

$$= 0.0796 + 2 - \sum_{n=0}^1 (2-n)P_n$$

$$= 2.0796 - (2P_0 + P_1)$$

$$= 2.0796 - \left( 2 \times 0.5003 + \frac{2}{3} \times 0.5003 \right)$$

$$= 0.75$$

$$(b) E(W_q) = \frac{1}{\lambda'} E(N_q) \quad \quad \quad \text{[by formula (10) of model IV]}$$

$$\text{where } \lambda' = \mu \left[ s - \sum_{n=0}^{s-1} (s-n) P_n \right] \quad \quad \quad \text{[by formula (7) of model IV]}$$

$$= 3 \left[ 2 - \sum_{n=0}^1 (2-n) P_n \right]$$

$$\begin{aligned}
 &= 3 [2 - (2P_0 + P_1)] \\
 &= 3 \left[ 2 - \left( 2 \times 0.5003 + \frac{2}{3} \times 0.5003 \right) \right] \\
 &= 1.9976 \\
 \therefore E(W_q) &= \frac{0.0796}{1.9976} = 0.0398 \text{ h}
 \end{aligned}$$

Every time an engineer approaches the computer centre, he has to lose 0.0398 h by way of waiting.

If the day consists of 8 working hours, he has to approach the centre 16 times.

$$\begin{aligned}
 \therefore \text{Time lost in waiting in a day per engineer} \\
 &= 16 \times 0.0398 = 0.6368 \text{ h}
 \end{aligned}$$

$\therefore$  Total time lost in waiting in a day by all the 6 engineers =  $6 \times 0.6368 = 3.82$  h.

### Example 21

A shipping company has a single unloading berth with ships arriving in a Poisson fashion at an average rate of 3 per day. The unloading time distribution for a ship is found to be exponential with average unloading time of  $\frac{1}{2}$  day. The company has a large labour supply without regular working hours and to avoid waiting of ships for getting unloaded, the company will use as many unloading labourers as there are arriving ships for getting unloaded. Under these conditions, (a) what will be the average number of unloading labourers working at any time? (b) what is the probability that more than 4 labourers will be needed?

Mean arrival rate =  $\lambda = 3$  ships/day

Mean unloading (service) time =  $\frac{1}{2}$  day

$\therefore$  Mean service rate =  $\mu = 2$  ships/day

Given that  $\lambda_n$  is a constant ( $=\lambda$ ) and  $\mu_n$  increases with  $n$  i.e.,  $\mu_n = n \mu$ .

$\therefore$  this is a problem under case (ii) of the birth and death process/self-service

model with  $\rho = \frac{\lambda}{\mu} = \frac{3}{2}$

(a)  $\therefore E(N_S) =$  The average number of workers unloading at any time

$$= \rho = \frac{3}{2} = 1.5 \text{ workers}$$

(b)  $P$  (more than 4 labourers will be needed)

$$\begin{aligned}
 &= P(N_S \geq 5) \text{ or } P(X \geq 5) \text{ or } P_n \geq 5 \\
 &= 1 - (P_0 + P_1 + P_2 + P_3 + P_4)
 \end{aligned}$$

$$= 1 - \left\{ 1 + \frac{1.5}{1!} + \frac{(1.5)^2}{2!} + \frac{(1.5)^3}{3!} + \frac{(1.5)^4}{4!} \right\} e^{-1.5}$$

$$= 0.019$$

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**Example 22**


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In a general Poisson queue system, find the steady-state probability  $P_n$  and the average number of customers in the system, given that

$$\lambda_n = \begin{cases} \lambda, & \text{for } n = 0, 1, 2, \dots, (k-1) \\ 0, & \text{for } n = k, k+1, \dots \end{cases}$$

and

$$\mu_n = n \mu, \text{ for } n = 1, 2, \dots$$

$$P_0^{-1} = 1 + \sum_{n=1}^k \left( \frac{\lambda_0 \lambda_1 \lambda_2 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \mu_3 \cdots \mu_n} \right)$$

$$= 1 + \sum_{n=1}^k \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n \text{ or } \sum_{n=0}^k \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n$$

$$P_k = \frac{\lambda_0 \lambda_1 \cdots \lambda_{k-1}}{\mu_1 \mu_2 \cdots \mu_k}, P_0 = \frac{1}{k!} \left( \frac{\lambda}{\mu} \right)^k \div \sum_{n=0}^k \frac{1}{n!} \left( \frac{\lambda}{\mu} \right)^n$$

and

$$E(N_s) = \sum_{k=0}^{\infty} k P_k$$

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**Example 23**


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In a general Poisson queue system, find the steady-state probability  $P_r$  and the average number of customers in the system, given that

$$\lambda_r = \begin{cases} (k-r)\lambda, & \text{for } r = 0, 1, 2, \dots, (k-1) \\ 0, & \text{for } r = k, k+1, \dots \end{cases}$$

and  $\mu_r = r\mu$ , for  $r \geq 1$

$$P_0^{-1} = 1 + \sum_{r=1}^k \left( \frac{\lambda_0 \lambda_1 \lambda_2 \cdots \lambda_{r-1}}{\mu_1 \mu_2 \mu_3 \cdots \mu_r} \right)$$

$$= 1 + \frac{k\lambda}{1!\mu} + \frac{k\lambda(k-1)\lambda}{2!\mu} + \frac{k\lambda(k-1)\lambda(k-2)\lambda}{3!\mu} + \cdots + \frac{k\lambda \cdot (k-1)\lambda \cdots 1\lambda}{k!\mu}$$

$$= \left( 1 + \frac{\lambda}{\mu} \right)^k \therefore P_0 = \left( 1 + \frac{\lambda}{\mu} \right)^{-k}$$

$$P_r = \frac{\lambda_0 \cdot \lambda_1 \cdots \lambda_{r-1}}{\mu_1 \cdot \mu_2 \cdots \mu_r} P_0 = \frac{k(k-1)\cdots(k-r+1)}{r!} \left( \frac{\lambda}{\mu} \right)^r \cdot \left( 1 + \frac{\lambda}{\mu} \right)^{-k}$$

$$\begin{aligned}
 &= kC_r \cdot \left( \frac{\lambda}{\mu} \right)^r \cdot \left( \frac{\mu}{\lambda + \mu} \right)^k \\
 &= kC_r \cdot \left( \frac{\mu}{\lambda + \mu} \right)^{k-r} \left( \frac{\mu}{\lambda + \mu} \right)^r \cdot \left( \frac{\lambda}{\mu} \right)^r \\
 &= kC_r \left( \frac{\mu}{\lambda + \mu} \right)^{k-r} \left( \frac{\lambda}{\lambda + \mu} \right)^r
 \end{aligned}$$

[ $r = 0, 1, 2, \dots, k$ ] which is the probability law of a binomial distribution, with  $p = \frac{\lambda}{\lambda + \mu}$  and  $q = \frac{\mu}{\lambda + \mu}$

$$\therefore E(N_s) = E(X) = kp = \frac{k\lambda}{\lambda + \mu}$$

## Exercise 8

### Part-A (Short-answer Questions)

1. What are the characteristics of a queueing system?
2. What do the letters in the symbolic representation ( $a/b/c$ ): ( $d/e$ ) of a queueing model represent?
3. What do you mean by transient state and steady-state queueing systems?
4. Write down the difference equations that give the probability that there are  $n$  customers ( $n \geq 0$ ) in a Poisson queueing system in steady-state.
5. Write down the formulas for  $P_0$  and  $P_n$  in a Poisson queue system in the steady-state.
6. Give the formulas for the average number of customers (i) in the system, (ii) in the queue, and (iii) in the non-empty queues for the ( $M/M/1$ ): ( $\infty/FIFO$ ) model.
7. Obtain the variance of queue length for the ( $M/M/1$ ): ( $\infty/FIFO$ ) model.
8. In the usual notation of a ( $M/M/1$ ): ( $\infty/FIFO$ ) queue system, find  $P(N > 2)$ , if  $\lambda = 12$  per hour and  $\mu = 30$  per hour.
9. In the usual notation of a ( $M/M/1$ ): ( $\infty/FIFO$ ) queue system if  $\lambda = 12$  per hour and  $\mu = 24$  per hour, find the average number of customers in the system and in the queue.
10. Give the formulas for the waiting time of a customer in the queue and in the system for the ( $M/M/1$ ): ( $\infty/FIFO$ ) model.
11. If a customer has to wait in a ( $M/M/1$ ): ( $\infty/FIFO$ ) queue system, what is his average waiting time in the queue, if  $\lambda = 8$  per hour and  $\mu = 12$  per hour?

12. What is the probability that a customer has to wait more than 15 min to get his service completed in a  $(M/M/1)$ :  $(\infty/FIFO)$  queue system, if  $\lambda = 6$  per hour and  $\mu = 10$  per hour?
13. Write down the probability density function of the waiting time of a customer in the  $(M/M/1)$ :  $(\infty/FIFO)$  queue system.
14. Write down the Little's formulas that hold good for all the Poisson queue models.
15. Write down the Little's formulas that hold good for the infinite capacity Poisson queue models.
16. What is the probability that there are no customers in the  $(M/M/s)$ :  $(\infty/FIFO)$  queueing system?
17. Write down the formula for  $P_n$  in terms of  $P_0$  for the  $(M/M/s)$ :  $(\infty/FIFO)$  queueing system.
18. Give the formulas for the average number of customers in the system and in the queue for the  $(M/M/s)$ :  $(\infty/FIFO)$  queueing model.
19. If there are 2 servers in an infinite capacity Poisson queue system with  $\lambda = 10$  per hour and  $\mu = 15$  per hour, what is the percentage of idle time for each server?
20. In a 3-server infinite capacity Poisson queue model if  $\lambda/s\mu = \frac{2}{3}$ , find  $P_0$ .
21. In a 3-server infinite capacity Poisson queue model if  $\lambda/s\mu = 2/3$  and  $P_0 = 1/9$  find the average number of customers in the queue and in the system.
22. If  $\lambda/s\mu = \frac{2}{3}$  in a  $(M/M/s)$ :  $(\infty/FIFO)$  queue system find the average number of customers in the non-empty queue.
23. What is the probability that an arrival to an infinite-capacity 3-server Poisson queue system with  $\lambda/s\mu = 2/3$  and  $P_0 = 1/9$  will have to wait?
24. What is the probability that an arrival to an infinite-capacity 3-server Poisson queue system with  $\lambda/s\mu = 2/3$  and  $P_0 = 1/9$  enters the service without waiting?
25. Give the formulas for the average waiting time of a customer in the system and in the queue for the  $(M/M/s)$ :  $(\infty/FIFO)$  queueing model.
26. What is the average waiting time of a customer in the 3-server infinite-capacity Poisson queue if he happens to wait, given that  $\lambda = 6$  per hour and  $\mu = 4$  per hour.
27. Give the probability that there is no customer in an  $(M/M/1)$ :  $(k/FIFO)$  queueing system.
28. Write down the probability that there are  $n$  customers in an  $(M/M/1)$ :  $(k/FIFO)$  queueing system.
29. If  $\lambda = 4$  per hour and  $\mu = 12$  per hour in an  $(M/M/1)$ :  $(4/FIFO)$  queueing system, find the probability that there is no customer in the system. If  $\lambda = \mu$ , what is the value of this probability?

30. Write the formulas for the average number of customers in the  $(M/M/1):(k/FIFO)$  queueing system and also in the queue.
31. Define effective arrival rate with respect to an  $(M/M/1):(k/FIFO)$  queueing model.
32. How are  $N_s$  and  $N_q$  related in an  $(M/M/1):(k/FIFO)$  queueing model?
33. Write down the Little's formulas for the average waiting time in the system and in the queue for an  $(M/M/1):(k/FIFO)$  queueing model.
34. If  $\lambda = 3$  per hour,  $\mu = 4$  per hour and maximum capacity  $k = 7$  in a  $(M/M/1):(k/FIFO)$  system, find the average number of customers in the system.
35. Write the formula for the probability that there is no customer in an  $(M/M/s):(k/FIFO)$  queue system.
36. Write the formula for the probability that there are  $n$  customers in an  $(M/M/s):(k/FIFO)$  queueing system.
37. Write down the formula for the average queue length in an  $(M/M/s):(k/FIFO)$  queueing model.
38. Define effective arrival rate with respect to an  $(M/M/s):(k/FIFO)$  queueing model.
39. How are  $N_s$  and  $N_q$  related in an  $(M/M/s):(k/FIFO)$  queueing model?
40. Write down Little's formulas for the average waiting time in the system and in the queue for an  $(M/M/s):(k/FIFO)$  queueing model.

### **Part B**

41. Arrivals at a telephone booth are considered to be Poisson with an average time of 10 min between one arrival and the next. The length of a phone call is assumed to be distributed exponentially with mean 3 min.
  - (i) Find the average number of persons waiting in the system.
  - (ii) What is the probability that a person arriving at the booth will have to wait in the queue?
  - (iii) What is the probability that it will take him more than 10 min altogether to wait for phone and complete his call?
  - (iv) Estimate the fraction of the day when the phone will be in use.
  - (v) The telephone department will install a second booth when convinced that an arrival has to wait on the average for at least 3 min for phone. By how much the flow of arrivals should increase in order to justify a second booth?
42. Customers arrive at a one man barber shop according to a Poisson process with a mean interarrival time of 20 min. Customers spend an average of 15 min in the barber's chair.
  - (i) What is the expected number of customers in the barber shop? in the queue?
  - (ii) What is the probability that a customer will not have to wait for a hair cut?
  - (iii) How much can a customer expect to spend in the barber shop?

- (iv) Management will put another chair and hire another barber when a customer's average waiting time in the shop exceeds 1.25 h. How much must the average rate of arrivals increase to warrant a second barber?
  - (v) What is average time customers spend in the queue?
  - (vi) What is the probability that the waiting time in the system is greater than 30 min?
  - (vii) What is the probability that there are more than 3 customers in the system?
43. If customers arrive for service according to a Poisson distribution at the average rate of 5 per day, how fast must they be serviced on the average (assume exponential service time) in order to keep the average number of customers in the system less than 4?
44. Patients arrive at an hospital for emergency service at the rate of one every hour. Currently only one emergency can be handled at a time. Patients spend an average of 20 min for receiving emergency service. How much the average service time need to be decreased to keep the average time to wait and receive the service less than 25 min?
45. A departmental secretary receives an average of 8 jobs per hour. Many are short jobs, while others are quite long. Assume, however, that the time to perform a job has an exponential distribution with a mean of 6 min.
- (i) What is the average elapsed time from the time the secretary receives a job until it is completed?
  - (ii) Calculate  $E(N)$ ,  $E(W_q)$ ,  $P(W > 2h)$ ,  $P(N > 5)$  and the fraction of time the secretary is busy.
46. A service station expects a customer every 4 min on the average. Service takes, on the average, 3 min. Assume Poisson input and exponential service.
- (i) What is the average number of customers waiting for service?
  - (ii) How long can a customer expect to wait for service?
  - (iii) What is the probability that a customer will spend less than 15 min waiting for and getting service?
  - (iv) What is the probability that a customer will spend longer than 10 min waiting for and getting service?
47. A dress shop has 3 sales persons. Assume that arrivals follow Poisson pattern with an average of 10 min between arrivals. Also assume that any salesperson can provide the desired service for any customer. If the time to provide service for a customer is exponentially distributed with a mean of 20 min per customer, calculate  $E(N)$ ,  $E(N_q)$ ,  $E(W)$ ,  $E(W_q)$  and  $P_n$  for  $n = 0, 1$  and  $2$ .
48. If the mean arrival rate is 24 per hour, find from the customer's point of view of the time spent in the system, whether 3 channels in parallel with a mean service rate of 10 per hour is better or worse than a single channel with mean service rate of 30 per hour.

49. Four counters are being run on the frontier of a country to check the passports and necessary papers of the tourists. The tourists choose a counter at random. If the arrivals at the frontier is Poisson at the rate of  $\lambda$  and the service time is exponential with parameter  $\frac{\lambda}{2}$ , what is the steady-state average queue at each counter?
50. An insurance company has 3 claim registers in its branch office. People with claims against the company are found to arrive in a Poisson fashion at an average rate of 20/8-h day. The amount of time that an adjuster spends with a claimant is found to have exponential distribution with mean service time 40 min. Claimants are processed in the order of their appearance.
- (i) How many hours a week can an adjuster expect to spend with claimants?
  - (ii) How much time, on the average, does a claimant spend in the branch office?
51. A telephone exchange has 2 long-distance operators. The telephone company finds that during the peak load, long distance calls arrive in a Poisson fashion at an average rate of 15 per hour. The length of service on these calls is approximately exponentially distributed with mean length 5 min.
- (a) What is the probability that a subscriber will have to wait for his long distance call during the peak hours of the day?
  - (b) If the subscribers will wait and are serviced in turn, what is the expected waiting time?
52. A petrol pump station has 4 pumps. The service times follow the exponential distribution with a mean of 6 min and cars arrive for service in a Poisson process at the rate of 30 cars per hour.
- (i) What is the probability that an arrival would have to wait in line?
  - (ii) Find the average waiting time in the queue, average time spent in the system and the average number of cars in the system.
  - (iii) For what percentage of time would a pump be idle on an average?
53. A one-person barber shop has 6 chairs to accommodate people waiting for a hair cut. Assume that customers who arrive when all the 6 chairs are full leave without entering the barber shop. Customers arrive at the average rate of 3 per hour and spend an average of 15 min in the barber's chair.
- (i) What is the probability that a customer can get directly into the barber's chair upon arrival?
  - (ii) What is the expected number of customers waiting for a hair cut?
  - (iii) How much time can a customer expect to spend in the barber shop?
  - (iv) What fraction of potential customers are turned away?
54. Assume that the goods trains are coming in a yard at the rate of 30 trains per day and suppose that the inter-arrival times follow an exponential

- distribution. The service time for each train is assumed to be exponential with an average of 36 min. If the yard can admit 9 trains at a time, calculate the probability that the yard is empty and the average queue length.
55. A car park contains 5 cars. The arrival of cars is Poisson at a mean rate of 10 per hour. The length of time each car spends in the car park has negative exponential distribution with mean of 2 min. How many cars are in the car park on an average and what is the probability of a newly arriving customer finding the car park full and leaving to park his car elsewhere.
56. A stenographer is attached to 5 officers for whom she performs stenographic work. She gets calls from the officers at the rate of 4 per hour and takes on the average 10 min to attend to each call. If arrival rate is Poisson and service time is exponential find (i) the average waiting time for an arriving call (ii) the average number of waiting calls, and (iii) the average time an arriving call spends in the system.
57. A 2-person barber shop has 5 chairs to accommodate waiting customers. Potential customers, who arrive when all the 5 chairs are full, leave without entering the barber shop. Customers arrive at the average rate of 3.7634 per hour and spend an average of 15 min in the barber's chair. Compute  $P_0$ ,  $P_1$ ,  $P_7$ ,  $E(N_q)$  and  $E(W)$ .
58. A barber shop has 2 barbers and 3 chairs for waiting customers. Assume that customers arrive in Poisson fashion at a rate of 5 per hour and that each barber services customers according to exponential distribution with mean of 15 min. Further, if a customer arrives and there are no empty chairs in the shop, he will leave. Find the steady-state probabilities. What is the probability that the shop is empty? What is the expected number of customers in the shop?
59. An automobile inspection station has 3 inspection stalls. Assume that cars wait in such a way that when a stall becomes vacant, the car at the head of the line pulls up to it. The station can accommodate at most 4 cars waiting (7 in the station) at one time. The arrival pattern is Poisson with a mean of 1 car every minute during the peak hours. The service time is exponential with mean 6 min. Find the average number of customers in the system during peak hours, the average waiting time and the average number of cars per hour that cannot enter the house because of full capacity.
60. A mechanic repairs 4 machines. The mean time between service requirements is 5 h for each machine and forms an exponential distribution. The mean repair time is 1 h and also follows the same distribution pattern. Machine downtime costs Rs. 25 per hour and the mechanic costs Rs. 55 per day.
- Find the expected number of machines under service and in waiting.
  - Determine the expected downtime cost per 8-h day?

- (iii) Would it be economical to engage 2 mechanics, each repairing only 2 machines?
61. For the general queue system, find the steady-state probability, when  $\lambda_n = \lambda$ , for  $n \geq 0$  and  $\mu_n = n^k \mu$ , for  $n \geq 1$ .
62. For the general queue system, find the steady-state probability and the average number of customers in the system, when

$$\lambda_n = \begin{cases} \lambda, & \text{for } n = 0, 1, 2, \dots, (k-1) \\ 0, & \text{for } n \geq k \end{cases}$$

- and  $\mu_n = \mu$ , for  $n \geq 1$ .
63. Research students visit a computer centre to solve their problems (one problem each) in a Poisson fashion at an average rate of 5 per day. The time to solve a problem by a student has an exponential distribution with a mean time of  $\frac{1}{3}$  day. Assuming that the centre can provide any number of computers to the students and that the average solving time is inversely proportional to the number of students working on the problem, find the average time spent by a student in the centre.
64. In a big departmental store, self-service pattern is allowed. Customers enter the store in a Poisson fashion at an average rate of 15 per hour. The time taken by a customer to select his/her requirements follows an exponential distribution with a mean time of 10 minutes. Assuming that billing and payment take negligible time, find the average number of customers in the store and the probability that there are no customers in the store.
65. There are two sections—cosmetics and provisions—in a big departmental store, where self-service of customers is permitted. Customers enter the cosmetics section and provisions section in Poisson fashion at an average rate of 15 per hour and 20 per hour respectively. The time taken by a customer to select his/her requirements follows an exponential distribution with a mean time of 6 minutes in the cosmetics section and 12 minutes in the provisions section. Assuming that billing and payment take negligible time in both sections, find the average number of customers in the stores and the probability that there are less than 5 customers in both the sections put together.

## **ANSWERS**

### **Exercise 8**

- |                                 |            |                   |                   |
|---------------------------------|------------|-------------------|-------------------|
| 11. 5 min                       | 12. 0.3679 | 19. 50%           | 20. $\frac{1}{9}$ |
| 21. $\frac{8}{9}, \frac{26}{9}$ | 22. 2      | 23. $\frac{4}{9}$ | 24. $\frac{5}{9}$ |

26. 10 min      29.  $\frac{81}{121}, \frac{1}{5}$       34. 2.11
41. (i) 0.43      (ii) 0.3      (iii) 0.097      (iv) 0.3  
 (v) 0.16 person per minute
42. (i) 3; 2.25      (ii) 0.25      (iii) 1 h      (iv) 3.2 per hour  
 (v)  $\frac{3}{4}$  h      (vi) 0.61      (vi) 0.32
43. 3.84 h/service
44. 17.65 min per patient
45. (i) 30 min      (ii) 4; 24 min; 0.0183; 0.2621; 0.8
46. (i) 2.25      (ii) 9 min      (iii) 0.7135      (iv) 0.8465
47. 2.1739; 0.1739; 21.739 min; 1.739 min;  $P_0 = 0.1304$ ,  $P_1 = 0.2608$ ;  
 $P_2 = 0.2602$
48. Single channel is better
49.  $\frac{4}{23}$
50. 22.2 h; 49 min
51. (i) 0.48      (ii) 3.2 min
52. (i) 0.3826      (ii) 3.05 min; 9.054 min; 4.53 cars; (iii) 24.98%
53. (i) 0.2778      (ii) 1.3878      (iii) 43.8 min      (iv) 3.7%
54. 0.28; 1.55
55. 0.49, 0.0027
56. (i) 12.45 min      (ii) 0.79 customer (iii) 22.42 min
57. 0.36133; 0.33996; 0.00368; 0.2457; 19 min
58.  $P_n = (0.56)(0.625)^n$ , for  $2 \leq n \leq 5$ ;  $P_1 = 0.35$ ;  $P_0 = 0.28$ ; 2.956
59. 6.06 cars; 12.3 min; 30.4 cars
60. (i) 1      (ii) Total cost = Rs. 255  
 (iii) Total cost with 2 machines = Rs. 270; Use of 2 machines is not economical.

61.  $P_n = \frac{\left(\frac{\lambda}{\mu}\right)^n}{(n!)^\alpha} \cdot P_0$ , where  $P_0^{-1} = \sum_{n=0}^{\infty} \frac{\left(\frac{\lambda}{\mu}\right)^n}{(n!)^\alpha}$

62.  $P_n = \left(\frac{1-\rho}{1-\rho^{k+1}}\right) \cdot \rho^n, 0 \leq n < k$ , where  $\rho = \frac{\lambda}{\mu}$

$$E(N_S) = [1 - (1 + k)\rho^k + k\rho^{k+1}] \rho / (1 - \rho)(1 - \rho^{k+1})$$

63. 8 hours

64. 2.5;  $e^{-2.5}$

65. 5.5;  $\sum_{r=0}^4 \frac{1}{r!} (5.5)^r \cdot r^{-5.5}$

# Chapter 9

## Advanced Queue Models and Queue Networks

### Finite Source Model

So far in all the queue models we have discussed, we assumed that the number of customers who require service is very large. In other words, the customers requiring service emanate from an infinite population (source). However, there are situations, where maximum limit is imposed on the incoming customers. In other words, the flow of customers is from a finite source.

A practical situation corresponding to this model is that of machine servicing where the calling population is that of machines and the arrival corresponds to the broken-down machine. The repairmen/mechanics are of course the servers. Symbolically, this machine-servicing model is represented as  $(M/M/1)$ :  $(k/k/\text{FIFO})$  or  $(M/M/s)$ :  $(k/k/\text{FIFO})$ , depending on the number of repairmen.

### Characteristics of Machine-Service Model $(M/M/s)$ : $(k/k/\text{FIFO})$

For the general Poisson queue systems, the steady-state probabilities  $P_0$  and  $P_n$  are given by

$$P_0^{-1} = 1 + \sum_{n=0}^{\infty} \left( \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} \right) \quad (1)$$

and 
$$P_n = \frac{\lambda_0 \lambda_1 \cdots \lambda_{n-1}}{\mu_1 \mu_2 \cdots \mu_n} P_0, \quad \text{for } n \geq 1 \quad (2)$$

For the machine-service model, we assume that there are  $k$  machines, one or more of which are likely to break down and, hence, become customers for the queue system in which there are ‘ $s$ ’ repairmen or servers.

Thus, 
$$\lambda_n = \begin{cases} (k-n)\lambda, & \text{for } 0 \leq n < k \\ 0, & \text{for } n \geq k \end{cases} \quad (3)$$

and  $\mu_n = \begin{cases} n\mu, & \text{for } 0 \leq n < s \\ s\mu, & \text{for } s \leq n \leq k \end{cases}$  (4)

Using these values in (1) and (2), we get

$$P_0^{-1} = \sum_{n=0}^{s-1} kC_n \rho^n + \sum_{n=s}^k kC_n \frac{n!}{s! s^{n-s}} \rho^n, \text{ where } \rho = \frac{\lambda}{\mu} \quad (5)$$

and  $P_n = \begin{cases} kC_n \rho^n \cdot P_0, & \text{for } 0 \leq n < s \\ kC_n \frac{n!}{s! s^{n-s}} \rho^n \cdot P_0, & \text{for } s \leq n \leq k \end{cases}$  (6)

$$E(N_Q) = \sum_{n=s+1}^k (n-s) P_n \quad (7)$$

$$E(N_S) = E(N_Q) + \frac{\lambda'}{\mu} \quad (8)$$

where  $\lambda' = E(\lambda_n) = \sum_{n=0}^{k-1} (k-n) \lambda P_n = \lambda \{k - E(N_S)\}$  (9)

$$E(W_Q) = \frac{1}{\lambda'} E(N_Q) \quad (10)$$

$$E(W_S) = \frac{1}{\lambda'} E(N_S) \quad (11)$$

**Note**

If there is only one repairman (server), viz.,  $s = 1$ , the results (5) to (11) become

$$P_0^{-1} = \sum_{n=0}^1 kC_n \cdot \rho^n \quad (5');$$

$$P_n = kC_n \rho^n \cdot P_0; 0 \leq n \leq 1 \quad (6');$$

$$E(N_Q) = k - \left(1 + \frac{1}{\rho}\right)(1 - P_0) \quad (7');$$

$$E(N_S) = k - \frac{1}{\rho}(1 - P_0) \quad (8');$$

$$\lambda' = \mu(1 - P_0) \quad \text{or} \quad \frac{\lambda}{\rho}(1 - P_0) \quad (9');$$

$$E(W_Q) = \frac{1}{\lambda'} E(N_Q) \quad (10');$$

and  $E(W_S) = \frac{1}{\lambda'} E(N_S) \quad (11)'$

## Non-Markovian Queueing Model [(M/G/1): ( $\infty$ /GD) Model]

So far we have discussed Markovian queue models in which the inter-arrival and inter-service times were assumed to follow exponential distributions with parameters  $\lambda$  and  $\mu$ . When the arrivals and departures do not follow Poisson distributions, the discussion of the queueing models is tedious. However we can derive the characteristics of a particular non-Markovian model ( $M/G/1$ ) : ( $\infty/GD$ ), where  $M$  indicates that the number of arrivals in time  $t$  follows a Poisson process,  $G$  indicates that the service time follows a general (arbitrary) distribution and  $GD$  indicates general queue discipline (viz., any kind of queue discipline).

The average number  $L$  of customers in the  $M/G/1$  queueing system is given by a formula, known as Pollaczek–Khinchine formula, which is derived below:

### **Pollaczek–Khinchine Formula**

Let  $N$  and  $N'$  be the numbers of customers in the system at times  $t$  and  $t + T$ , when two consecutive customers have just left the system after getting service.

Thus,  $T$  is the random service time, which is a continuous random variable. Let  $f(t)$ ,  $E(T)$ ,  $\text{Var}(T)$  be the pdf, mean and variance of  $T$ . Also let  $M$  be the number of customers arriving in the system during the service time  $T$ .

$$\text{Hence } N' = \begin{cases} M, & \text{if } N = 0 \\ N - 1 + M, & \text{if } N > 0 \end{cases}$$

where  $M$  is a discrete random variable, taking the values  $0, 1, 2, \dots$

$$\text{Equivalently, } N' = N - 1 + M + \delta \quad (1)$$

$$\text{where } \delta = \begin{cases} 1, & \text{if } N = 0 \\ 0, & \text{if } N > 0 \end{cases}$$

$$\therefore E(N') = E(N) - 1 + E(M) + E(\delta) \quad (2)$$

When the system has reached the steady-state, the probability of the number of customers in the system will be a constant.

$$\text{Hence, } E(N) = E(N') \text{ and } E(N^2) = E(N'^2) \quad (3)$$

$$\text{Using this in (2), we get } E(\delta) = 1 - E(M) \quad (4)$$

Squaring both sides of (1), we have

$$N^2 = N^2 + (M - 1)^2 + \delta^2 + 2N(M - 1) + 2(M - 1)\delta + 2N\delta \quad (5)$$

$$\text{Now, } \delta^2 = \delta \quad (\because \delta^2 = 0 \text{ or } 1, \text{ according as } \delta = 0 \text{ or } 1)$$

$$\text{and } N\delta = \begin{cases} 0 \times 1, & \text{if } N = 0 \\ N \times 0, & \text{if } N > 0 \end{cases}$$

$$= 0$$

Using these values in (5), we have

$$N'^2 = N^2 + M^2 + 2N(M - 1) + (2M - 1)\delta - 2M + 1$$

$$\text{i.e., } 2N(1 - M) = N^2 - N'^2 + M^2 + (2M - 1)\delta - 2M + 1$$

$$\therefore 2E\{N(1 - M)\} = E(N^2) - E(N'^2) + E(M^2) + E\{(2M - 1)\delta\} - 2E(M) + 1$$

$$\text{i.e., } 2E(N)\{1 - E(M)\} = E(M^2) + \{2E(M) - 1\}E(\delta) - 2E(M) + 1$$

[by independence and by (3)]

$$\therefore E(N) = \frac{E(M^2) + \{2E(M) - 1\}\{1 - E(M)\} - 2E(M) + 1}{2\{1 - E(M)\}}, \text{ by (4)}$$

$$= \frac{E(M^2) - 2E(M) + E(M)}{2\{1 - E(M)\}} \quad (6)$$

Since the number  $M$  of arrivals in time  $T$  follows a Poisson process with parameter  $\lambda$ , say, then  $E(M) = \lambda T$  and  $\text{Var}(M) = \lambda T$  or  $E(M^2) = (\lambda T)^2 + \lambda T$

$$\text{Now, } E(M) = E\{E(M/T)\}$$

$$= E(\lambda T) = \lambda E(T) \quad (7)$$

$$\begin{aligned} E(M^2) &= E\{E(M^2/T)\} = E\{\lambda^2 T^2 + \lambda T\} \\ &= \lambda^2 \{\text{Var}(T) + E^2(T)\} + \lambda E(T) \end{aligned} \quad (8)$$

Using (6) and (7) in (5), we have

$$\begin{aligned} L_s &= E(N) = \frac{\lambda^2 V(T) + \lambda^2 E^2(T) + \lambda E(T) - 2\lambda^2 E^2(T) + \lambda E(T)}{2\{1 - \lambda E(T)\}} \\ &= \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}} \end{aligned}$$

### Note

1. The other characteristics  $L_q = E(N_q)$ ,  $E(W_q)$  and  $E(W_{q'})$  of this model can be obtained by using Little's formulas.
2.  $\lambda E(T)$  must be less than 1, otherwise  $L_s$  becomes negative, which is meaningless.
3. In this  $M/G/1$  model, if  $G \equiv M$ , viz., the service time  $T$  follows an exponential distribution with parameter  $\mu$ , then

$$E(T) = \frac{1}{\mu} \text{ and } V(T) = \frac{1}{\mu^2}$$

$$L_s = \frac{\lambda}{\mu} + \frac{\lambda^2 \left\{ \frac{1}{\mu^2} + \frac{1}{\mu^2} \right\}}{1 \left( 1 - \frac{\lambda}{\mu} \right)} = \frac{\lambda}{\mu - \lambda}$$

which has already been derived for  $M/M/1$  model.

## Particular Cases

### 1. $(M/D/1);(\infty/\text{FIFO})$ Model

This is a particular case of the  $(M/G/1); (\infty/GD)$  Model, where  $D$  means that the service time  $T$  is not random, but it is fixed. In this case,  $E(T)$  is assumed as  $\frac{1}{\mu}$  and  $V(T) = 0$ .

Using these values in the Pollaczek–Khintchine formula for  $E(N_S)$ , we get

$$E(N_S) = \frac{\lambda}{\mu} + \frac{\lambda^2 \left\{ 0 + \frac{1}{\mu_2} \right\}}{2 \left\{ 1 - \frac{\lambda}{\mu} \right\}} = \rho + \frac{\rho^2}{2(1-\rho)}$$

where  $\rho = \frac{\lambda}{\mu}$  or  $\lambda E(T)$

### 2. $(M/E_k/1);(\infty/\text{FIFO})$ Model

Here,  $E_k$  means that the service time  $T$  follows an Erlang's distribution with  $k$  exponential phases, each with parameter  $\mu$ .

$$\text{Hence, } E(T) = \frac{k}{\mu} \text{ and } V(T) = \frac{k}{\mu^2}$$

Using these values in the Pollaczek–Khintchine formula for  $E(N_S)$ , we get

$$\begin{aligned} E(N_S) &= \frac{k\lambda}{\mu} + \frac{\lambda^2 \left\{ \frac{k}{\mu^2} + \frac{k^2}{\mu^2} \right\}}{2 \left\{ 1 - \frac{k\lambda}{\mu} \right\}} \\ &= k\rho + \frac{k(k+1)\rho^2}{2(1-k\rho)}, \text{ where } \rho = \frac{\lambda}{\mu} \end{aligned}$$

### Worked Example 9(A)

#### Example 1

A machine shop has 4 machines in working condition initially and one mechanic is attached to the shop. The mean time between service requirements is 5 hours and forms an exponential distribution. The mean service time is 1 hour and follows an exponential distribution. Machine downtime costs Rs. 55/- per day. Determine (a) the probability that the service facility will be idle, (b) the average number of machines waiting to be repaired and being repaired, and (c) the expected downtime cost per day.

This is a problem in the machine-service model  $(M/M/1):(k/k/\text{FIFO})$ .

Here,  $k =$  the number of machines in the calling population  $= 4$ ;  $\lambda = \frac{1}{5}$  per hour;  $\mu = 1$  per hour;  $\rho = \frac{\lambda}{\mu} = 0.2$

$$\begin{aligned} \text{(a)} \quad P_0^{-1} &= 4C_0\rho^0 + 4C_1\rho^1 + 4C_2\cdot 2!\rho^2 + 4C_3\cdot 3!\rho^3 + 4C_4\cdot 4!\rho^4 \\ &= 1 + 4 \times 0.2 + 6 \times 2 \times 0.04 + 4 \times 6 \times 0.008 + 24 \times 0.0016, \text{ since } s = 1 \\ &= 2.5104 \end{aligned}$$

$$\therefore P_0 = P(\text{the service facility is idle}) = (2.5104)^{-1} = 0.3983$$

$$\text{(b)} \quad E(N_S) = k - \frac{1}{\rho}(1 - P_0) = 4 - 5 \times (1 - 0.3983) = 0.9915$$

$$\text{(c)} \quad \lambda' = \mu(1 - P_0) = 1 \times (1 - 0.3983) = 0.6017$$

$$E(W_S) = \frac{1}{\lambda'} E(N_S) = \frac{0.9915}{0.6017} = 1.6478 \text{ hours}$$

$$E(\text{Downtime cost}) = 55 + \frac{8}{1.6478} \times 25, \text{ assuming that the machine shop works}$$

for 8 hours a day  $= 55 + 121.37 = \text{Rs. } 176.37$

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### Example 2

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In the machine shop of a small-scale industry, machines break down with a mean rate of 2 per hour. The maintenance shop of the industry has 2 mechanics who can attend to the breakdown machines individually. The service rate of each of the mechanics is 1.5 machines per hour. Initially there are 5 working machines in the workshop. Find  $P_0$ ,  $E(N_Q)$ ,  $E(N_S)$ ,  $E(W_Q)$  and  $E(W_S)$ .

This is a problem in the machine-service model  $(M/M/s):(k/k/\text{FIFO})$ , when  $s = 2$  and  $k = 5$

$$P_0^{-1} = \sum_{n=0}^{s-1} kC_n \rho^n + \sum_{n=s}^k kC_n \frac{n!}{s! s^{n-s}} \rho^n,$$

$$\text{where } \rho = \frac{\lambda}{\mu} = \frac{2}{1.5} = \frac{4}{3}$$

$$\begin{aligned} &= 1 + 5C_1 \cdot \rho + 5C_2 \cdot \frac{2!}{2! 2^0} \cdot \rho^2 + 5C_3 \cdot \frac{3!}{2! 2^1} \rho^3 \\ &\quad + 5C_4 \cdot \frac{4!}{2! 2^2} \cdot \rho^4 + 5C_5 \cdot \frac{5!}{2! 2^3} \cdot \rho^5 = 140.0122 \end{aligned}$$

$$\therefore P_0 = 0.0072$$

$$\begin{aligned} E(N_Q) &= \sum_{n=s+1}^k (n-s)P_n = \sum_{n=3}^5 (n-2)P_n = \sum_{n=3}^5 (n-2)5C_n \cdot \frac{n!}{2 \times 2^{n-2}} \cdot P_0 \\ &= 1.604 \text{ machines} \end{aligned}$$

$$\begin{aligned}
\lambda' &= \sum_{n=0}^{k-1} (k-n) \lambda \cdot P_n \\
&= 2[5 \times P_0 + 4 \times P_1 + 3 \times P_2 + 2 \times P_3 + 1 \times P_4] \\
&= 2[5 + 26.6667 + 53.3333 + 35.5556 + 47.4074] \times 0.0072 \\
&= 2.937/\text{hour} \\
E(N_S) &= E(N_Q) + \frac{\lambda'}{\mu} = 1.604 + \frac{2.937}{1.5} = 3.562 \text{ machines} \\
E(W_Q) &= \frac{1}{\lambda'} E(N_Q) = \frac{1.604}{2.937} = 0.546 \text{ hour} \\
E(W_S) &= \frac{1}{\lambda'} E(N_S) = \frac{3.562}{2.937} = 1.213 \text{ hour}
\end{aligned}$$

### Example 3

The town police department has 5 patrol cars. A patrol car breaks down and requires service once in every 30 days. The police department has 2 repair workers, each of whom takes an average of 3 days to repair a car. Break times and service times are exponential. (a) Determine the average number of patrol cars in good condition, (b) Find the average downtime for a patrol car that needs service, and (c) Find the fraction of time a service mean is idle.

This is a problem in the machine-service model  $(M/M/s):(k/k/\text{FIFO})$ , where  $s = 2$  and  $k = 5$ .

$$\lambda = \frac{1}{30} \text{ per day and } \mu = \frac{1}{3} \text{ per day}$$

$$\begin{aligned}
P_0^{-1} &= \sum_{n=0}^2 5C_n \cdot \rho^n + \sum_{n=3}^5 5C_n \cdot \frac{n!}{2!2^{n-2}} \cdot \rho^n, \text{ where } \rho = \frac{\lambda}{\mu} = 0.1 \\
&= \{1 + 5 \times 0.1 + 10 \times (0.1)^2\} + \\
&\quad \left\{ 10 \times \frac{6}{2 \times 2} (0.1)^3 + 5 \times \frac{24}{2 \times 4} \times (0.1)^4 + 1 \times \frac{120}{2 \times 8} \times (0.1)^5 \right\} \\
&= 1.6166
\end{aligned}$$

$$\therefore P_0 = 0.6186$$

$$\begin{aligned}
E(N_Q) &= \sum_{n=s+1}^k (n-s) P_n = \sum_{n=s+1}^k (n-s) \cdot kC_n \frac{n!}{s! s^{n-s}} \rho^n \cdot P_0 \\
&= \left[ 1 \times 10 \times \frac{3!}{2 \times 2^1} \times (0.1)^3 + 2 \times 5 \times \right. \\
&\quad \left. \frac{4!}{2 \times 2^2} \times (0.1)^4 + 3 \times 1 \times \frac{5!}{2 \times 2^3} \times (0.1)^5 \right] \times 0.6186 \\
&= 0.0112
\end{aligned}$$

$$\begin{aligned}
 \lambda' &= \sum_{n=0}^4 (5-n) \lambda P_n \\
 \lambda' &= \frac{1}{30} [5P_0 + 4P_1 + 3P_2 + 2P_3 + P_4] \\
 &= \frac{0.6186}{30} = [5 + 2 + 0.3 + 0.0075 + 0.003] = 0.1507 \\
 E(N_S) &= E(N_Q) + \frac{\lambda'}{\mu} = 0.0112 + \frac{0.1507}{(1/3)} = 0.4633; \\
 \overline{E(N_S)} &= 5 - 0.4633 \\
 &= 4.5367 \\
 E(W_S) &= \frac{1}{\lambda'} E(N_S) = \frac{0.4633}{0.1507} = 3.0743 \text{ days}
 \end{aligned}$$

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### Example 4

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In a college canteen, it was observed that there is only one waiter who takes exactly 4 minutes to serve a cup of coffee once the order has been placed with him. If students arrive in the canteen at an average rate of 10 per hour, how much time one is expected to spend waiting for his turn to place the order.

This is a problem in  $(M/D/1):(\infty/\text{FIFO})$  model, where  $D$  indicates that the service  $T$  is a constant.

Here,  $E(T) = 4$  minutes and  $\lambda = 10$  per hour or  $\frac{1}{6}$  per minute. By the Pollackzek–Khinchine formula,

$$\begin{aligned}
 E(N_S) &= \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}} \\
 &= \frac{1}{6} \times 4 + \frac{\frac{1}{36} \{0 + 16\}}{2 \left\{1 - \frac{4}{6}\right\}}, \text{ since } V(T) = 0, \text{ as } T \text{ is a constant} \\
 &= \frac{2}{3} + \frac{4}{9} \times \frac{3}{2} = \frac{4}{3}
 \end{aligned}$$

By Little's formula,  $E(W_S) = \frac{1}{\lambda} E(N_S) = 6 \times \frac{4}{3} = 8$  minutes

$$\begin{aligned}
 E(W_Q) &= E(W_S) - \frac{1}{\mu} = 8 - 4 \left\{ \text{since } \mu = \frac{1}{E(T)} \right\} \\
 &= 4 \text{ minutes}
 \end{aligned}$$

---

### Example 5

---

A one-man barber shop takes exactly 25 minutes to complete one haircut. If customers arrive at the barber shop in a Poisson fashion at an average rate of one every 40 minutes, how long on the average does a customer spend in the shop? Also, find the average time a customer must wait for service.

The service time  $T$  is a constant = 25 min viz.,  $T$  follows a distribution with  $E(T) = 25$  and  $V(T) = 0$ . Also  $\lambda = \frac{1}{40}$ .

$\therefore$  by the Pollaczek–Khinchine formula,

$$\begin{aligned} E(N_s) &= \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}} \\ &= \frac{25}{40} + \frac{\frac{1}{40^2} \{0 + 25^2\}}{2\left\{1 - \frac{1}{40} \times 25\right\}} \\ &= \frac{5}{8} + \frac{25/64}{2 \times (3/8)} = \frac{55}{48} \end{aligned}$$

By Little's formula,

$$E(W_s) = \frac{1}{\lambda} E(N_s) = 40 \times \frac{55}{48} = 45.8 \text{ minutes}$$

$$E(W_q) = E(W_s) - \frac{1}{\mu} = E(W_s) - E(T) = 20.8 \text{ min.}$$

i.e., a customer has to spend 45.8 minutes in the shop and has to wait for service for 20.8 minutes on the average.

---

### Example 6

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For an  $(M/E_2/1):(\infty/\text{FIFO})$  model with  $\lambda = \frac{6}{5}$  per hour and  $\mu = \frac{3}{2}$  per hour, find the average waiting time of a customer. Also find the average time he spends in the system.

The service time  $T$  follows on Erlang distribution with  $k = 2$  exponential phases

$$\text{Hence, } E(T) = \frac{k}{\mu} = \frac{4}{3} \text{ hours and } V(T) = \frac{k}{\mu^2} = \frac{8}{9}$$

By the Pollackzek–Khinchine and Little's formulas,

$$\begin{aligned} E(W_Q) &= \text{Average waiting time of a customer} \\ &= \left( \frac{k+1}{2k} \right) \cdot \frac{\lambda}{\mu(\mu - \lambda)} \end{aligned}$$

$$= \frac{3}{4} \cdot \frac{6/5}{\frac{3}{2} \left( \frac{3}{2} - \frac{6}{5} \right)} = \frac{3}{4} \times \frac{6}{5} \times \frac{2}{3} \times \frac{10}{3} = 2 \text{ hours}$$

$E(W_s)$  = Average time spent in the system

$$= E(W_Q) + \frac{1}{\mu} = \left( 2 + \frac{2}{3} \right) \text{ hours} = 2 \text{ hours, } 40 \text{ minutes}$$

### Example 7

A patient who goes to a single-doctor clinic for a general check-up has to go through 4 phases. The doctor takes on the average 4 minutes for each phase of the check-up and the time taken for each phase is exponentially distributed. If the arrivals of the patients at the clinic are approximately Poisson at the average rate of 3 per hour, what is the average time spent by a patient (a) in the examination? (b) waiting in the clinic?

Let  $X_1, X_2, X_3, X_4$  denote the times required for the 4 phases of the check-up.

Each  $X_r$  is exponential with a mean of 4 min or with parameter  $\frac{1}{4}$ .

Since the  $X_r$ 's independent,  $(X_1 + X_2 + X_3 + X_4)$  follows an Erlang's distribution with parameters ' $\lambda$ ' =  $\frac{1}{4}$  and ' $k$ ' = 4 [Refer to Problem (43) in Exercise 5(B)]

The mean and variance of Erlang's distribution with parameters ' $\lambda$ ' and ' $k$ ' are  $\frac{k}{\lambda}$  and  $\frac{k}{\lambda^2}$ .

Thus, if  $T$  represents the service time for a patient,

$$E(T) = \frac{k}{\lambda} = \frac{4}{1/4} = 16$$

and  $V(T) = \frac{k}{\lambda^2} = \frac{4}{1/16} = 64$

$\therefore$  Average time for examination of each patient = 16 min. If  $\lambda_c$  represents the arrival rate in the clinic, then by  $P - K$  formula,

$$\begin{aligned} E(N_s) &= \lambda_c E(T) + \frac{\lambda_c^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda_c E(T)\}} \\ &= \frac{1}{20} \times 16 + \frac{\frac{1}{400} \{64 + 256\}}{2 \left\{ 1 - \frac{1}{20} \times 16 \right\}} \quad \left( \because \lambda_c = 3/\text{hour} \text{ or } \frac{1}{20}/\text{min} \right) \end{aligned}$$

$$= \frac{4}{5} + \frac{\frac{4}{5}}{2 \times \frac{1}{5}} = \frac{14}{5}$$

By Little's formula,

$$E(W_s) = \frac{1}{\lambda_c} E(N_s) = 20 \times \frac{14}{5} = 56 \text{ minutes}$$

and  $E(W_q) = E(W_s) - \frac{1}{\mu} = 56 - \frac{1}{1/E(T)} = 56 - 16 = 40 \text{ min.}$

i.e., a patient has to wait 40 minutes for check-up in the clinic.

---

### Example 8

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A car-wash facility operates with only one bay. Cars arrive according to a Poisson distribution with a mean of 4 cars per hour and may wait in the facility's parking lot if the bay is busy. The parking lot is large enough to accommodate any number of cars. Find the average number of cars waiting in the parking lot, if the time for washing and cleaning a car follows

- (a) uniform distribution between 8 and 12 minutes,
- (b) a normal distribution with mean 12 minutes and S.D. 3 minutes, and
- (c) a discrete distribution with values equal to 4, 8 and 15 minutes and corresponding probabilities 0.2, 0.6 and 0.2.

(a)  $\lambda = 4/\text{hour}$  or  $\frac{1}{15}$  per minute.

$$\begin{aligned} E(T) &= \text{mean of the uniform distribution in } (8, 12) \\ &= \frac{1}{2} (8 + 12) = 10 \text{ minutes} \end{aligned}$$

$$V(T) = \frac{1}{12} (12 - 8)^2 = \frac{4}{3}$$

By P-K formula,

$$\begin{aligned} E(N_s) &= \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda \cdot E(T)\}} \\ &= \frac{1}{15} \times 10 + \frac{\frac{1}{225} \left\{ \frac{4}{3} + 100 \right\}}{2 \left\{ 1 - \frac{1}{15} \times 10 \right\}} \\ &= \frac{2}{3} + \frac{152}{225} = 1.342 \text{ cars.} \end{aligned}$$

By Little's formula,

$$E(N_q) = E(N_s) - \frac{\lambda}{\mu} = 1.342 - \frac{\left(\frac{1}{15}\right)}{\left(\frac{1}{10}\right)} \quad \left(\because \mu = \frac{1}{E(T)}\right)$$

$$= 0.675 \text{ car}$$

(b)  $\lambda = \frac{1}{15}$ ;  $E(T) = 12$  min. and  $V(T) = 9$

$$\therefore \mu = \frac{1}{E(T)} = \frac{1}{12}$$

By P-K formula,

$$E(N_s) = \frac{1}{15} \times 12 + \frac{\frac{1}{225} \{9 + 144\}}{2 \left\{1 - \frac{1}{15} \times 12\right\}}$$

$$= \frac{4}{5} + \frac{153}{90} = 2.5 \text{ cars}$$

By Little's formula,

$$E(N_q) = E(N_s) - \frac{\lambda}{\mu} = 2.5 - \frac{(1/15)}{(1/12)} = 1.7 \text{ cars}$$

- (c) The service time  $T$  follows the discrete distribution given below:

$T:$	4	8	15
$p(T):$	0.2	0.6	0.2

$$E(T) = \sum T \cdot p(T) = 0.8 + 4.8 + 3 = 8.6 \text{ minutes}$$

$$E(T^2) = \sum T^2 \cdot p(T) = 3.2 + 38.4 + 45 = 86.6$$

$$V(T) = E(T^2) - E^2(T) = 86.6 - (8.6)^2 = 12.64$$

By P-K formula,

$$E(N_s) = \frac{1}{15} \times 8.6 + \frac{\frac{1}{225} \{12.64 + 73.96\}}{2 \left\{1 - \frac{1}{15} \times 8.6\right\}}$$

$$= 0.573 + \frac{86.6}{225} \times \frac{15}{12.8} = 0.573 + 0.451 = 1.024 \text{ cars}$$

By Little's formula,

$$E(N_q) = E(N_s) - \frac{\lambda}{\mu} = 1.024 - \frac{(1/5)}{(1/8.6)} = 0.451 \text{ car}$$

### Exercise 9(A)

**Part-A (Short-answer Questions)**

1. Explain the symbolic representation  $(M/M/s):(k/k/\text{FIFO})$  of the machine-service model.
2. Write down the formulas for  $P_0$  and  $P_n$  for the  $(M/M/s):(k/k/\text{FIFO})$  model.
3. How do the relations between (i)  $E(N_Q)$  and  $E(W_Q)$ , and (ii)  $E(N_S)$  and  $E(W_S)$  get modified from the usual Little's formulas in the case of  $(M/M/s):(k/k/\text{FIFO})$  model.
4. Write down the formulas for  $E(N_Q)$  and  $E(N_S)$  for the  $(M/M/1):(k/k/\text{FIFO})$  model.
5. Write down the general Pollaczek–Khinchine formula for  $E(N_S)$ .
6. Write down the simplified formulas for  $E(N_S)$  for the  $(M/D/1):(\infty/\text{FIFO})$  and  $(M/E_k/1):(\infty/\text{FIFO})$  models.

**Part-B**

7. In the machine shop of a medium-scale industry, machines break down with a mean rate of 2.5 per hour. The maintenance shop of the industry has only one mechanic who can attend to the broken-down machines. The service rate of the mechanic is 1.5 machines per hour. Initially, there are 6 working machines in the machine shop. Compute  $P_0$ ,  $E(N_Q)$ ,  $E(N_S)$ ,  $E(W_Q)$  and  $E(W_S)$ .
8. A mechanic services 4 machines. For each machine, the mean time between service requirements is 10 hours and is assumed to form an exponential distribution. The repair time tends to follow the same distribution and has a mean time of 2 hours. When a machine is down for repair, the time last has a value of Rs. 20 per hour. The mechanic costs Rs. 50 day. Find (i) the expected number of machines in operation and (ii) the expected cost of downtime per day. (iii) Would it be desirable to engage two mechanics, each to service only two machines?
9. There are 5 machines in a workshop, each of which suffers breakdown at an average rate of 2 per hour. There are 2 servicemen and only one man can work on a machine at a time. If  $n$  machines are out of order when  $n > 2$ , then  $(n - 2)$  of them wait until a serviceman is free. Once a serviceman starts work on a machine, the time to complete the repair has an exponential distribution with a mean of 5 minutes. Find the distribution of the number of machines out of action at a given time. Find also the average time an out of action machine has to spend waiting for the repairs to start.
10. 2 repairmen are attending to 5 machines in a workshop. Each machine breaks down according to exponential distribution with a mean of 3 per hour. The repair time per machine is exponential with a mean of

- 15 minutes. (i) Find the probability that the two repairmen are idle, (ii) that one repairman is idle, and (iii) the expected number of idle machine not being serviced.
11. The 5 machines in a machine shop break down periodically and the company has 2 repairmen to service the machines when breakdown occurs. When a machine is fixed, the time until the next breakdown is exponentially distributed with a mean of 30 hours. The repair time of each repairman is exponentially distributed with a mean of 3 hours. Find the average number of machines in operation at any given time, the expected down-time of a machine that require repair and the expected percent of idle time of each repairman.
  12. In a car-manufacturing plant, a loading crane takes exactly 10 minutes to load a car into a wagon and again come back to position to load another car. If the arrival of cars is a Poisson stream at an average of 1 every 20 minutes, calculate the average waiting time of a car.
  13. A barber runs his own saloon. It takes him exactly 25 minutes to complete one hair cut. Customers arrive in a Poisson fashion at an average rate of one every 35 minutes. (i) For what per cent of time would the barber be idle? (ii) What is the average time a customer spent in the shop?
  14. Repairing a certain type of faulty machine consists of 5 basic steps that must be performed sequentially. The time taken to perform each of the 5 steps is exponentially distributed with a mean of 5 minutes and independent of the other steps. If the machines become faulty in a Poisson fashion at an average rate of 2 per hour and if there is only one repairman, what is the average idle time for each faulty machine?
  15. In a certain bank, the customers arrive according to an exponential distribution with a mean of 4 per hour. From observations on the teller's performance, the mean service time is estimated to be 10 minutes, with a variance of 25. It is felt that the Erlang would be a reasonable assumption for the distribution of the teller's service times. Also it is assumed that there is no limit on the number of customers. Bank officials wish to know, on the average, how long a customer must wait until he gets the service and how many customers are waiting for service.

[**Hint:** Find the value of  $k$ , using  $\mu = \frac{1}{10}$  and  $\sigma^2 = \frac{1}{k\mu^2} = 25$  ]

## QUEUE NETWORKS

### Introduction

Queue network can be regarded as a group of ' $k$ ' inter-connected nodes, where each node represents a service facility of some kind with  $s_i$  servers at node  $i$ . ( $s_i \geq 1$ )

## Series Queues

A *series queue* model or a tandem queue model is, in general, one in which (i) customers may arrive from outside the system at any node and may leave the system from any node, (ii) customers may enter the system at some node, traverse from node to node in the system and leave the system from some node, not necessarily following the same order of nodes, and (iii) customers may return to the nodes previously visited, skip some nodes and even choose to remain in the system for ever.

In the following sections, we shall discuss a few kinds of series queues.

### ***Series Queues with Blocking***

This is a sequential queue model consisting of two service points  $S_1$  and  $S_2$ , at each of which there is only one server and where no queue is allowed to form at either point.

**Note**

*It is a misnomer by which we called the model as a series queue model!*

An entering customer will first go to  $S_1$ , after he gets the service completed at  $S_1$ , he will go to  $S_2$  if it is empty or will wait in  $S_1$  until  $S_2$  becomes empty [viz., the system is blocked for a new customer.] This means that a potential customer will enter the system only when  $S_1$  is empty, irrespective of whether  $S_2$  is empty or not, since the model is a sequential model, viz., all the customers require service at  $S_1$  and then at  $S_2$ .

Let us now proceed to find the steady-state probabilities  $P(m, n)$  that there is  $m$  customer ( $m = 0$  or  $1$ ) in  $S_1$  and  $n$  customer ( $n = 0$  or  $1$ ) in  $S_2$ . Any state of the model will be denoted by  $(m, n)$ . The possible states of the system are given below with their interpretation:

State	Interpretation
$(0, 0)$	No customer in either service point
$(1, 0)$	Only one customer in $S_1$
$(0, 1)$	Only one customer in $S_2$
$(1, 1)$	One customer each in $S_1$ and $S_2$
$(b, 1)$	One customer each in $S_1$ and $S_2$ but the customer in $S_1$ having finished his work at $S_1$ is waiting for $S_2$ to become free, while the customer in $S_2$ is being served.

We assume that potential customers arrive in accordance with a Poisson process with parameter  $\lambda$  and the service times at  $S_1$  and  $S_2$  follow exponential distributions with parameters  $\mu_1$  and  $\mu_2$  respectively. To get the values of  $P(m, n)$ , we shall first write down the steady-state balance equations using a state transition diagram, as given below, which consists of a small circle for each state and directed lines labeled by the rates at which the process goes from one state to another.

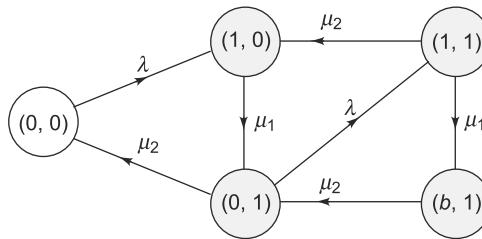


Fig. 9.1 State Transition Diagram

**Note** The arrow from state  $(0, 0)$  to state  $(1, 0)$  is labeled  $\lambda$ . This means that when the system (consisting of  $S_1$  and  $S_2$ ) is empty, it goes to state  $(1, 0)$  through an arrival which enters the system at a rate  $\lambda$ . Similar explanation holds good for the arrow from  $(0, 1)$  to  $(1, 1)$ . The transition from  $(1, 0)$  to  $(0, 1)$  takes place, when the customer finishes his job at  $S_1$  at the rate  $\mu_1$ . Similar explanation holds good for the arrow from  $(1, 1)$  to  $(b, 1)$ . When in state  $(b, 1)$ , the process will go to state  $(0, 1)$ , when the customer at  $S_2$  completes his service, that occurs at the rate  $\mu_2$ . Similar explanations hold good for the transitions from  $(1, 1)$  to  $(1, 0)$  and from  $(0, 1)$  to  $(0, 0)$ .

Now, the balance equation corresponding to any state  $(m, n)$  is obtained by equating the rate at which the process leaves that state and the rate at which the process enters that state. The rate at which the process leaves the state  $(m, n) = P(m, n) \times$  sum of the labels of the arrows that leave the state  $(m, n)$ . The rate at which the process enters the state  $(m, n) =$  sum of the labels of the arrows that enter  $(m, n)$  multiplied by the relevant probabilities of the states from which they emanate. Accordingly, we get the following balance equations for the 5 states of the system:

State	Balance equation	
$(0, 0)$	$\lambda P(0, 0) = \mu_2 P(0, 1)$	(1)
$(1, 0)$	$\mu_1 P(1, 0) = \lambda P(0, 0) + \mu_2 P(1, 1)$	(2)
$(0, 1)$	$(\lambda + \mu_2) P(0, 1) = \mu_1 P(1, 0) + \mu_2 P(b, 1)$	(3)
$(1, 1)$	$(\mu_1 + \mu_2) P(1, 1) = \lambda P(0, 1)$	(4)
$(b, 1)$	$\mu_2 P(b, 1) = \mu_1 P(1, 1)$	(5)

Since the process has to be in any one of the 5 mutually exclusive and exhaustive states, we have

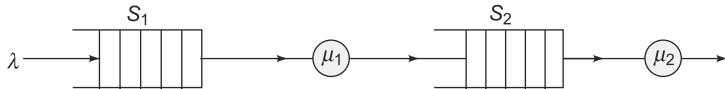
$$P(0, 0) + P(1, 0) + P(0, 1) + P(1, 1) + P(b, 1) = 1 \quad (6)$$

Solving the above six equations, we can get the five steady-state probabilities.

## Two-Stage (Service Point) Series Queues (Two-Stage Tandem Queues)

Let us consider a two (service) stage queuing system in which customers arrive from outside at a Poisson rate  $\lambda$  to  $S_1$ . After being served at  $S_1$ , they then join the queue in front of  $S_2$ . After receiving service at  $S_2$ , they leave the system.

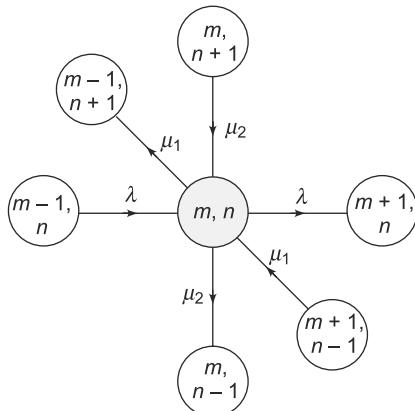
It is assumed that there is infinite waiting (queueing space) at each service point. Each server serves one customer at a time and the service times at  $S_1$  and  $S_2$  follow exponential distributions with parameters  $\mu_1$  and  $\mu_2$  respectively.



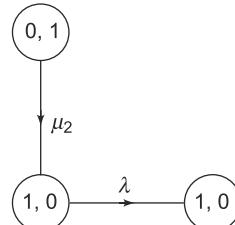
**Fig. 9.2**

To find the steady-state joint probability  $P(m, n)$  of  $m$  customers in  $S_1$  and  $n$  customers in  $S_2$ , where  $m \geq 0$  and  $n \geq 0$ , we shall write down the balance equations using the state transition diagrams given below. By taking  $(m, n)$  as the central state, we have shown only the arrows entering into and leaving from it.

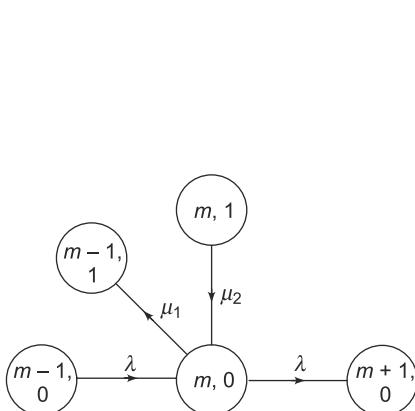
**Note** The state transition diagram [Fig. 9.4] for the central state  $(0, 0)$  is extracted from Fig. 9.3 by retaining only those states for which  $m \geq 0$  and  $n \geq 0$  and the connecting arrows. Similarly the state transition diagrams [Fig. 9.5 and Fig. 9.6] for the central states  $(m, 0)$  and  $(n, 0)$  are obtained.



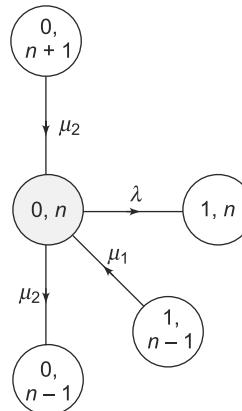
**Fig. 9.3**



**Fig. 9.4**



**Fig. 9.5**



**Fig. 9.6**

The balance equations for this model are obtained as given below:

$$\lambda P(0, 0) = \mu_2 P(0, 1) \quad (1)$$

$$(\lambda + \mu_1) P(m, 0) = \lambda P(m - 1, 0) + \mu_2 P(m, 1) [m > 0] \quad (2)$$

$$(\lambda + \mu_2) P(0, n) = \mu_1 P(1, n - 1) + \mu_2 P(0, n + 1) [n > 0] \quad (3)$$

$$(\lambda + \mu_1 + \mu_2) P(m, n) = \lambda P(m - 1, n) + \mu_1 P(m + 1, n - 1) \\ + \mu_2 P(m, n + 1) \quad (4)$$

Also  $\sum_m \sum_n P(m, n) = 1$  (5)

We shall not attempt to solve the above balance equations directly, but shall guess the solution by using *Burke's theorem* which is stated below without proof:

For an  $M/M/s$  queueing system, the output (departure) process is also Poisson with the same rate  $\lambda$  as the input (arrival) process in the steady-state, where  $s \geq 1$ .

This means that the arrival process to  $S_2$ , which is the same as the departure process from  $S_1$ , is also Poisson with parameter  $\lambda$ . Thus, the queuing systems at both  $S_1$  and  $S_2$  are  $M/M/1$  models (if we assume that there is only one server at  $S_1$  and at  $S_2$ ).

From the discussion on the characteristics of  $M/M/1$  model

$$P(m \text{ customers at } S_1 \text{ system}) = \left( \frac{\lambda}{\mu_1} \right)^m \left( 1 - \frac{\lambda}{\mu_1} \right) \text{ and}$$

$$P(n \text{ customers at } S_2 \text{ system}) = \left( \frac{\lambda}{\mu_2} \right)^n \left( 1 - \frac{\lambda}{\mu_2} \right)$$

Since the number of customers at  $S_1$  and  $S_2$  are independent random variables, the joint probability of  $m$  customers at  $S_1$  and  $n$  customers at  $S_2$  is given by

$$P(m, n) = \left( \frac{\lambda}{\mu_1} \right)^m \cdot \left( 1 - \frac{\lambda}{\mu_1} \right)^m \cdot \left( \frac{\lambda}{\mu_2} \right)^n \cdot \left( 1 - \frac{\lambda}{\mu_2} \right)^n, \text{ for } m \geq 0 \text{ and } n \geq 0 \quad (6)$$

**Note** It is easily verified that (6) satisfies the balance equations (1) to (5). For example, on using (6) in (4),

$$\begin{aligned} \text{R.S. of (4)} &= \left( 1 - \frac{\lambda}{\mu_1} \right) \left( 1 - \frac{\lambda}{\mu_2} \right) \\ &\quad \left[ \lambda \cdot \left( \frac{\lambda}{\mu_1} \right)^{m-1} \cdot \left( \frac{\lambda}{\mu_2} \right)^n + \mu_1 \left( \frac{\lambda}{\mu_1} \right)^{m+1} \cdot \left( \frac{\lambda}{\mu_2} \right)^{n-1} + \mu_2 \left( \frac{\lambda}{\mu_1} \right)^m \cdot \left( \frac{\lambda}{\mu_2} \right)^{n+1} \right] \\ &= \left( 1 - \frac{\lambda}{\mu_1} \right) \left( 1 - \frac{\lambda}{\mu_2} \right) \left( \frac{\lambda}{\mu_1} \right)^{m-1} \left( \frac{\lambda}{\mu_2} \right)^{n-1} \left[ \frac{\lambda^2}{\mu_2} + \frac{\lambda^2}{\mu_1} + \frac{\lambda^3}{\mu_1 \mu_2} \right] \end{aligned}$$

$$\begin{aligned}
&= \left(1 - \frac{\lambda}{\mu_1}\right) \left(1 - \frac{\lambda}{\mu_2}\right) \left(\frac{\lambda}{\mu_1}\right)^{m-1} \left(\frac{\lambda}{\mu_2}\right)^{n-1} \frac{\lambda^2}{\mu_1 \mu_2} (\lambda + \mu_1 + \mu_2) \\
&= (\lambda + \mu_1 + \mu_2) \left[ \left(\frac{\lambda}{\mu_1}\right)^m \left(1 - \frac{\lambda}{\mu_1}\right) \right] \left[ \left(\frac{\lambda}{\mu_2}\right)^n \left(1 - \frac{\lambda}{\mu_2}\right) \right] \\
&= (\lambda + \mu_1 + \mu_2) P(m, n) \\
&= L.S. \text{ of (4)]}
\end{aligned}$$

The average number of customers in the system is given by

$$\begin{aligned}
L_S &= \sum_m \sum_n (m+n) P(m, n) \\
&= \sum_{m=0}^{\infty} m \left(\frac{\lambda}{\mu_1}\right)^m \left(1 - \frac{\lambda}{\mu_1}\right) + \sum_{n=0}^{\infty} n \left(\frac{\lambda}{\mu_2}\right)^n \left(1 - \frac{\lambda}{\mu_2}\right) \\
&= \left(1 - \frac{\lambda}{\mu_1}\right) \cdot \frac{\lambda}{\mu_1} \cdot \left(1 - \frac{\lambda}{\mu_1}\right)^{-2} + \left(1 - \frac{\lambda}{\mu_2}\right) \cdot \frac{\lambda}{\mu_2} \cdot \left(1 - \frac{\lambda}{\mu_2}\right)^{-2} \\
&= \frac{\lambda}{\mu_1 - \lambda} + \frac{\lambda}{\mu_2 - \lambda} = L_{S1} + L_{S2}
\end{aligned}$$

The average waiting time of a customer in the system is given by

$$\begin{aligned}
E(W_S) &= \frac{1}{\lambda} L_S, \text{ by Little's formula.} \\
&= \frac{1}{\mu_1 - \lambda} + \frac{1}{\mu_2 - \lambda} = E(W_{S1}) + E(W_{S2})
\end{aligned}$$

### Note

1. As the arrival rate  $\lambda$  in a 2-state tandem queue model increases, the node with the larger value of  $\rho_i = \frac{\lambda}{\mu_i}$  will introduce instability. Hence, the node with the larger value of  $\rho_i$  is called the **bottleneck** of the system.
2. The discussion given above holds good, even when there are many servers at each stage, provided the service time of each server at stage  $i$  ( $i = 1, 2$ ) is exponential with mean  $\frac{1}{\mu_i}$ .

In this case, it is to be noted that the queueing system at each stage is an M/M/s model and hence the relevant results corresponding to this model should be used.

## Open Jackson Networks

A network of  $k$  service facilities or nodes is called an *open Jackson network*, if it satisfies the following characteristics:

1. Arrivals, from outside, to the node  $i$  follow a Poisson process with mean rate  $r_i$  and join the queue at  $i$  and wait for his turn for service.

2. Service times at the channels at node  $i$  are independent and each is exponentially distributed with parameter  $\mu_i$ .
3. Once a customer gets the service completed at node  $i$ , he joins the queue at node  $j$  with probability  $P_{ij}$  (whatever be the number of customers waiting at  $j$  for service), where  $i = 1, 2, \dots, k$  and  $j = 0, 1, 2, \dots, k$ .  $P_{i0}$  represents the probability that a customer leaves the system from node  $i$  after getting the service at  $i$ .

**Note** Derivations and the methods of solving the balance equations for this model are beyond the scope of the book. However we just give below Jackson's solution for the balance equations of this model.

If we denote the total arrival rate of customers to server  $j$  [viz., the sum of the arrival rate  $r_j$  (Note: It is not  $\lambda_j$ ) to  $j$  coming from outside and the rates of departure  $\lambda_i$  from the servers  $i$ ] by  $\hat{\lambda}_j$ , then

$$\hat{\lambda}_j = r_j + \sum_{i=1}^k \lambda_i P_{ij}; j = 1, 2, \dots, k \quad (1)$$

$P_{ij}$  is the probability that a departure from server  $i$  joins the queue at server  $j$  and hence  $\lambda_i P_{ij}$  is the rate of arrival to server  $j$  from among those coming out from server  $i$ .

Equations (1) are called *Traffic equations or Flow balance equations*.

Jackson has proved that the steady-state solutions of these traffic equations with single server at each node is

$$P(n_1, n_2, \dots, n_k) = \rho_1^{n_1} (1 - \rho_1) \cdot \rho_2^{n_2} (1 - \rho_2) \dots \rho_k^{n_k} (1 - \rho_k) \quad (2)$$

where  $\rho_j = \frac{\lambda_j}{\mu_j}$ , provided  $\rho_j < 1$  for all  $j$ .

Since  $P(n_1, n_2, \dots, n_k) = [\rho_1^{n_1} (1 - \rho_1)][\rho_2^{n_2} (1 - \rho_2)] \dots [\rho_k^{n_k} (1 - \rho_k)]$

viz., the joint probability is equal to the product of the marginal probabilities, we can interpret that the network acts as if the queue at each node  $i$  is an independent  $M/M/1$  queue with rates  $\lambda_i$  and  $\mu_i$ .

**Note** 1. Even though (2) may be misinterpreted as the arrival process at node  $i$  is Poisson in this network model, it need not be so, as the customer may visit a server more than once (known as feedback situation).

Thus, whether the arrival processes at various nodes are Poisson or not, step (2) alone holds good and hence the network behaves as if its nodes were independent  $M/M/1$  queue models.

$$2. \sum_{i=1}^k E\{N_{s_i}\}, \text{ but } E(W_s) = \frac{1}{\sum_{i=1}^k r_i} E(N_s)$$

3. Jackson's open-network concept can be extended when the nodes are multi server nodes. In this case the network behaves as if each node is an independent  $M/M/s$  model.

## Closed Jackson Networks

A queueing network of  $k$  nodes is called a closed Jackson network, if new customers never enter into and the existing customers never depart from the system, viz., if  $r_i = 0$  and  $P_{i0} = 0$  for all  $i$ . In other words, it is equivalent to a finite source queueing system of  $N$  customers who traverse continuously inside the network where the service time of server  $i$  is exponentially distributed with rate  $\mu_i$ ;  $i = 1, 2, \dots, k$ .

When a customer completes service at  $S_i$ , he then joins the queue at  $S_j$ ,  $j = 1, 2, \dots, k$  with probability  $P_{ij}$  where it is assumed that  $\sum_{j=1}^k P_{ij} = 1$  for all  $i = 1, 2, \dots, k$ . We note that the matrix  $P = [P_{ij}]$  is similar to one-step transition probability matrix of a Markov chain, that is stochastic and irreducible.

The flow balance equations of this model become

$$\lambda_j = \sum_{i=1}^k \lambda_i P_{ij}; j = 1, 2, \dots, k \quad [\because r_j = 0] \dots \quad (1)$$

The matrix  $[P_{ij}]$  is called the *routing probability matrix* in this context. Jackson has proved that the steady-state solution of equation (1) is

$$P(n_1, n_2, \dots, n_k) = C_N \rho_1^{n_1} \rho_2^{n_2} \dots \rho_k^{n_k}, \text{ where}$$

$$C_N^{-1} = \sum_{n_1+n_2+\dots+n_k=N} \rho_1^{n_1} \rho_2^{n_2} \dots \rho_k^{n_k}, \text{ where } \rho_j = \frac{\lambda_j}{\mu_j}$$

**Note**  $N$  customers must be allocated among  $k$  nodes, such that  $n_1 + n_2 + \dots + n_k = N$ . This allocation can be done in  $\binom{N+k-1}{N}$  ways.

If there are  $s_i$  servers at node  $i$ , the solution is given by

$$P(n_1, n_2, \dots, n_k) = C_N \frac{\rho_1^{n_1}}{a_1(n_1)} \cdot \frac{\rho_2^{n_2}}{a_2(n_2)} \dots \cdot \frac{\rho_k^{n_k}}{a_k(n_k)},$$

where  $C_N^{-1} = \sum_{n_1+n_2+\dots+n_k=N} \frac{\rho_1^{n_1}}{a_1(n_1)} \frac{\rho_2^{n_2}}{a_2(n_2)} \dots \frac{\rho_k^{n_k}}{a_k(n_k)}$ ,

where  $a_i(n_i) = \begin{cases} n_i, & \text{if } n_i < s_i \\ s_i^{n_i-s_i} | s_i, & \text{if } n_i \geq s_i \end{cases}$

## Mean-Value Analysis (MVA)

This is another method which does not require evaluating  $C_N$  to find  $P(n_1, n_2, \dots, n_k)$ , that is cumbersome and time consuming, for large values of  $N$  and  $k$ . It is

based on a concept called the *Arrival Theorem* which we simply state below without proof:

The queue length observed by an arrival at node  $i$  is the same as the general queue length in a closed network with one less customer.

Using this, we can find the average waiting time at a node  $i$  in terms of the mean service time in the node and the average number of customers in the system as observed by an arrival at that node.

For the  $M/M/1$  model, we know that  $L = E(N_s) = \frac{\lambda}{\mu - \lambda}$  and  $W = E(W_s) = \frac{1}{\mu - \lambda}$ . From these, it follows that  $W = \frac{1+L}{\mu}$  (1). Assuming that all the nodes in the closed network have a single server, (1) takes the form

$$W_i(N) = \frac{1 + L_i(N-1)}{\mu_i} \quad (2),$$

$W_i(N)$  = mean waiting time at node  $i$  for the closed network containing  $N$  customers

$\mu_i$  = mean service rate for the single server at node  $i$

$L_i(N-1)$  = mean number of customers at node  $i$  in a closed network containing  $(N-1)$  customers.

**Note** Formula for  $W_i(N)$  in (2) does not hold good for a multiple server closed network.  
It is modified later.

Assuming that Little's formula  $E(N_s) = \lambda E(W_s)$  holds good for all the nodes of the closed network, we have  $L_i(N) = \lambda_i W_i(N)$  (3), where  $\lambda_i$  is the throughput (arrival) rate for node  $i$  in an  $N$ -customer closed network.

Equations (2) and (3) give us a method for recursively calculating  $L_i(N)$  and  $W_i(N)$  for any given  $N$  starting with an empty network for which  $L_i(0) = 0$  and  $W_i(0) = \frac{1}{\mu_i}$ . The method is explained below in the form of an algorithm (MVA algorithm for a  $k$ -node, single server-per-node network).

(1) Solve the modified traffic equations  $v_j = \sum_{i=1}^k v_i P_{ij}$  ( $j = 1, 2, \dots, k$ ) obtained

from  $\lambda_j = \sum_{i=1}^k \lambda_i P_{ij}$  ( $j = 1, 2, \dots, k$ ), where  $v_i = \frac{\lambda_i}{\lambda_l}$  so that  $v_l = 1$ , for some  $l$ .

(2) Initialise,  $L_i(0) = 0$ , for  $i = 1, 2, \dots, k$ .

(3) For  $r = 1$  to  $N$ , calculate

$$(a) \quad W_i(r) = \frac{1 + L_i(r-1)}{\mu_i}, \text{ for } i = 1, 2, \dots, k$$

$$(b) \quad \lambda_l(r) = \frac{r}{\sum_{i=1}^k v_i W_i(r)}, \text{ where } v_l = 1$$

- (c)  $\lambda_l(r) = v_i \lambda_l(r)$ , for  $i = 1, 2, \dots, k; i \neq l$   
 (d)  $L_i(r) = \lambda_l(r) W_l(r)$  for  $i = 1, 2, \dots, k$ .

MVA algorithm for finding  $L_i(N)$  and  $W_i(N)$  in a  $k$ -node, multiple-server ( $s_i$  servers at node  $i$ ) closed network with routing matrix  $P = \{P_{ij}\}$

**Note** Before we give the algorithm, we denote the marginal probability that there are  $m$  customers at node  $i$  in an  $N$ -customer, multiple server, closed network by  $P_i(m, N)$ .

Also  $P_i(m, N) = \frac{\lambda_i(N)}{\alpha_i(m)\mu_i} P_i(m-1, N-1)$ , for  $i \leq m \leq N-1$ ,

where  $\alpha_i(m) = \frac{\alpha_i(m)}{\alpha_i(m-1)} = \begin{cases} m, & \text{for } m \leq s_i \\ s_i, & \text{for } m \geq s_i \end{cases}$

- Solve the modified traffic equations  $v_j = \sum_{m=0}^k v_i P_{ij}$  ( $j = 1, 2, \dots, k$ ), where

$v_l = 1$ , for some  $l$ . [These equations hold good, even for a multiple server closed network]

- Initialise  $L_i(0) = 0; P_i(0, 0) = 1; P_i(m, 0) = 0, (m \neq 0)$
- For  $r = 1$  to  $N$ , calculate

(a)  $W_i(r) = \frac{1}{s_i \mu_i} \left[ 1 + L_i(r-1) + \sum_{m=0}^{s_i-2} (s_i - 1 - m) P_i(m, r-1) \right]$ , for  $i = 1, 2, \dots, k$ .

(b)  $\lambda_l(r) = r \sqrt[k]{\sum_{i=1}^k v_i W_i(r)}$ , where  $v_l = 1$

(c)  $\lambda_i(r) = v_l \lambda_l(r)$ , for  $i = 1, 2, \dots, k; i \neq l$

(d)  $L_i(r) = \lambda_i(r) W_i(r)$ , for  $i = 1, 2, \dots, k$

(e)  $P_i(m, r) = \frac{\lambda_i(r)}{\alpha_i(m)\mu_i} P_i(m-1, r-1)$ , for  $m = 1, 2, \dots, r, i = 1, 2, \dots, k$ .

**Note** Even if we are interested only in the values of  $W_i$  and  $L_i$ , we still require the marginal probabilities  $P_i(m, r-1)$ , for  $m = 0, 1, 2, \dots, (s_i - 2)$ .

### Worked Example 9(B)

#### Example 1

There are two salesmen in a ration shop, one in charge of billing and receiving payment and the other in charge of weighing and delivering the items. Due to limited availability of space, only one customer is allowed to enter the shop, that too when the billing clerk is free. The customer who has finished his billing job

has to wait there until the delivery section becomes free. If customers arrive in accordance with a Poisson process at rate 1 and the service times of two clerks are independent and have exponential rates of 3 and 2, find

- (a) the proportion of customers who enter the ration shop,
- (b) the average number of customers in the shop, and
- (c) the average amount of time that an entering customer spends in the shop.

The system in this problem is a sequential queue model with blocking;  $\lambda = 1$ ,  $\mu_1 = 3$  and  $\mu_2 = 2$

The balance equations for this model with the given value of  $\lambda, \mu_1, \mu_2$  are

$$P(0, 0) = 2P(0, 1), \dots \quad (1); \quad 3P(1, 0) = P(0, 0) + 2P(1, 1) \quad \dots \quad (2);$$

$$3P(0, 1) = 3P(1, 0) + 2P(b, 1) \dots \quad (3); \quad 5P(1, 1) = P(0, 1) \quad \dots \quad (4);$$

$$\begin{aligned} 2P(b, 1) &= 3P(1, 1), \dots \quad (5) \quad \text{and} \quad P(0, 0) + P(1, 0) + P(0, 1) \\ &\quad + P(1, 1) + P(b, 1) = 1 \dots \quad (6) \end{aligned}$$

$$\text{From (1),} \quad P(0, 1) = \frac{1}{2}P(0, 0) \quad (7)$$

$$\text{Using (7) in (4),} \quad P_{11} = \frac{1}{10}P(0, 0) \quad (8)$$

$$\text{Using (8) in (5),} \quad P(b, 1) = \frac{3}{20}P(0, 0) \quad (9)$$

$$\text{Using (8) in (2),} \quad P(1, 0) = \left(\frac{1}{3} + \frac{1}{15}\right)P(0, 0) \text{ or } \frac{2}{5}P(0, 0) \quad (10)$$

Using (7), (8) (9) and (10) in (6), we get

$$\left(1 + \frac{2}{5} + \frac{1}{2} + \frac{1}{10} + \frac{3}{20}\right)P(0, 0) = 1$$

$$\text{i.e.,} \quad \frac{43}{20}P(0, 0) = 1 \text{ or } P(0, 0) = \frac{20}{43}$$

$$\therefore P(0, 1) = \frac{10}{43}, P(1, 0) = \frac{8}{43}, P(1, 1) = \frac{2}{43}, P(b, 1) = \frac{3}{43}$$

- (a) Proportion of customers entering the shop

$$= P(0, 0) + P(0, 1) = \frac{30}{43}$$

- (b) Average number of customers in the shop

$$= 1 \times \{P(1, 0) + P(0, 1)\} + 2\{P(1, 1) + P(b, 1)\}$$

$$= 1 \times \frac{18}{43} + 2 \times \frac{5}{43} = \frac{28}{43}$$

- (c) Average amount of time that an entering customer spends in the shop is given by

$$\begin{aligned} E(W_s) &= \frac{1}{\lambda_a} \cdot E(N_s), \text{ where } \lambda_a = \text{average rate of entering the shop} \\ &\quad = \lambda \cdot \{P(0, 0) + P(0, 1)\} \\ &= \frac{28}{30} \text{ or } \frac{14}{15} &= 1 \times \frac{30}{43} \end{aligned}$$

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### Example 2

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If, in Example (1), customers reach the ration shop (not enter) according to Poisson process at the rate of 5 per hour and both the salesmen take 6 minutes each to serve a customer on the average and the service times follow exponential distribution, find the average number of customers in the shop and the average time spent by a customer who has entered the shop.

Here,  $\lambda = 5/\text{hour}$  and  $\mu_1 = \mu_2 = 10/\text{hour}$

The balance equations are

$$5 P(0, 0) = 10 P(0, 1) \dots (1); \quad 10 P(1, 0) = 5 P(0, 0) + 10 P(1, 1) \dots (2)$$

$$15 P(0, 1) = 10 \cdot P(1, 0) + 10 P(b, 1) \dots (3); \quad 20 P(1, 1) = 5 P(0, 1) \dots (4)$$

$$\begin{aligned} 10 \cdot P(b, 1) &= 10 \cdot P(1, 1) \dots (5); \quad P(0, 0) + P(1, 0) + P(0, 1) \\ &\quad + P(1, 1) + P(b, 1) = 1 \dots (6) \end{aligned}$$

Solving the above questions, we get

$$P(0, 1) = \frac{1}{2} P(0, 0); \quad P(1, 1) = \frac{1}{8} P(0, 0) = P(b, 1); \quad P(1, 0) = \frac{5}{8}$$

$$P(0, 0) \text{ and } \frac{19}{8} P(0, 0) = 1$$

$$\therefore P(0, 0) = \frac{8}{19}, \quad P(0, 1) = \frac{4}{19}, \quad P(1, 0) = \frac{5}{19}, \quad P(1, 1) = P(b, 1) = \frac{1}{19}$$

Average number of customers in the shop

$$= 1 \times \left( \frac{5}{19} + \frac{4}{19} \right) + 2 \times \left( \frac{1}{19} + \frac{1}{19} \right) = \frac{13}{19}$$

Average time spent by the customers inside the shop

$$= \frac{\frac{13}{19}}{5 \left( \frac{8}{19} + \frac{4}{19} \right)} = \frac{13}{60} \text{ hour or } 13 \text{ minutes}$$

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### Example 3

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In an ophthalmic clinic, there are two sections-one section for assessing the power approximately and the other for final assessment and prescription of glasses. Patients arrive at the clinic in a Poisson fashion at the rate of 3 per

hour. The assistant in the first section takes nearly 15 minutes per patient for approximate assessment of power and the doctor in the second section takes nearly 6 minutes per patient for final prescription. If the service times in the two sections are approximately exponential, find the probability that there are 3 patients in the first section and 2 patients in the second section. Find also the average number of patients in the clinic and the average waiting time of a patient in the clinic. Assume that enough space is available for the patients to wait in front of both sections.

The situation in this problem is comparable with 2-stage Tandem queue with single server at each state.

$$\lambda = 3/\text{hour}; \mu_1 = 4/\text{hour} \text{ and } \mu_2 = 10/\text{hour}$$

- (a)  $P(m, n) = P(m \text{ customers in the first section and } n \text{ customers in the second section})$

$$\begin{aligned} &= \left( \frac{\lambda}{\mu_1} \right)^m \left( 1 - \frac{\lambda}{\mu_1} \right) \cdot \left( \frac{\lambda}{\mu_2} \right)^n \left( 1 - \frac{\lambda}{\mu_2} \right) \\ \therefore P(3, 2) &= \left( \frac{3}{4} \right)^3 \left( 1 - \frac{3}{4} \right) \cdot \left( \frac{3}{10} \right)^2 \cdot \left( 1 - \frac{3}{10} \right) = \frac{1701}{2,56,000} \end{aligned}$$

- (b) Average number of patients in the clinic

$$\begin{aligned} &= \frac{1}{\mu_1 - \lambda} + \frac{1}{\mu_2 - \lambda} \\ &= \frac{3}{4 - 3} + \frac{3}{10 - 3} = 3 \frac{3}{7} \end{aligned}$$

- (c) Average waiting time of a patient in the clinic

$$\begin{aligned} &= \frac{1}{\mu_1 - \lambda} + \frac{1}{\mu_2 - \lambda} \\ &= \frac{1}{4 - 3} + \frac{1}{10 - 3} = 1 \frac{1}{7} \text{ hour} \end{aligned}$$

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### Example 4

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A repair facility shared by a large number of machines has 2 sequential stations with respective service rates of 2 per hour and 3 per hour. The cumulative failure rate of all the machines is 1 per hour. Assuming that the system behaviour may be approximated by the 2-stage tandem queue, find

- (a) the average repair time including the waiting time  
(b) the probability that both the service stations are idle, and (c) the bottleneck of the repair facility.

Here,  $\lambda = 1$ ;  $\mu_1 = 2$  and  $\mu_2 = 3$

- (a) The average number of machines in service at both the stations

$$\begin{aligned} &= \frac{\lambda}{\mu_1 - \lambda} + \frac{\lambda}{\mu_2 - \lambda} \\ &= \frac{1}{2-1} + \frac{1}{3-1} = 1\frac{1}{2}. \end{aligned}$$

The average repair time including the waiting time

$$\begin{aligned} &= \frac{1}{\mu_1 - \lambda} + \frac{1}{\mu_2 - \lambda} \\ &= \frac{1}{2-1} + \frac{1}{3-1} = 1\frac{1}{2} \text{ hours} \end{aligned}$$

- (b)  $P(\text{both the service stations are idle})$

$$\begin{aligned} &= P(0, 0) \\ &= \left( \frac{\lambda}{\mu_1} \right)^0 \cdot \left( 1 - \frac{\lambda}{\mu_1} \right) \cdot \left( \frac{\lambda}{\mu_2} \right)^0 \cdot \left( 1 - \frac{\lambda}{\mu_2} \right) \\ &= \left( 1 - \frac{1}{2} \right) \left( 1 - \frac{1}{3} \right) = \frac{1}{3} \end{aligned}$$

- (c) Since  $\frac{\lambda}{\mu_1} > \frac{\lambda}{\mu_2}$ , the service station 1 is the bottleneck of the repair facility.

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### Example 5

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In a busy medical shop located in the heart of a city, there are 3 salesmen who receive the customers, supply the drugs and prepare the bills. On finishing the job with any one of the salesmen, the customer goes to the payment counter, manned by the owner himself and leaves the shop after paying the bill. If customers enter the medical shop in a Poisson fashion at the rate of 30 per hour, each salesman takes on the average 5 minutes to serve a customer in an exponential fashion and the owner takes on the average 1 minute per customer to check the bill and receive the payment in an exponential fashion, find the average number of customers in the shop and the average time each customer spends in the shop. If the owner wishes that no customer should stay in the shop for more than 8 minutes and the number of customers in the shop should not exceed 5 and at the same time he wishes to take on the average 1.5 minutes per customer for his job, find if these can be accomplished by appointing one more salesman of the same calibre as others.

*Case (i)* For the multiple server queueing system in the supply section,  $\lambda = 30/\text{hour}$ ,  $\mu_1 = 12/\text{hour}$  and  $s = 3$

$$E(N_s) = \frac{1}{s \cdot \underline{s}} \frac{\left(\frac{\lambda}{\mu_1}\right)^{s+1} \cdot P_0}{\left(1 - \frac{\lambda}{\mu_1 s}\right)^2} + \frac{\lambda}{\mu_1} \quad (1)$$

where  $P_0^{-1} = \sum_{r=0}^{s-1} \frac{1}{\underline{r}} \left(\frac{\lambda}{\mu_1}\right)^r + \left\{ \frac{1}{\underline{s} \left(1 - \frac{\lambda}{\mu_1 s}\right)} \left(\frac{\lambda}{\mu_1}\right)^s \right\}$

$$= \sum_{r=0}^2 \frac{1}{\underline{r}} \left(\frac{5}{2}\right)^r + \left\{ \frac{1}{\underline{3} \left(1 - \frac{5}{6}\right)} \cdot \left(\frac{5}{2}\right)^3 \right\}$$

$$= 1 + \frac{5}{2} + \frac{25}{8} + \frac{125}{8} = \frac{89}{4} \quad (2)$$

Using (2) in (1),

$$E(N_s) = \frac{1}{3 \cdot \underline{3}} \cdot \frac{\left(\frac{5}{2}\right)^4}{\left(\frac{1}{6}\right)^2} \times \frac{4}{89} + \frac{5}{2} = 6.01$$

$$E(W_s) = \frac{1}{\lambda} \cdot E(N_s) = \frac{6.01}{30} \times 60 \text{ minutes or } 12.02 \text{ minutes.}$$

For the payment section,  $\lambda = 30/\text{hour}$ ;  $\mu_2 = 60/\text{hour}$  and the system is a single server system.

$$E(N_s) = \frac{\lambda}{\mu_2 - \lambda} = \frac{30}{60 - 30} = 1$$

and  $E(W_s) = \frac{1}{\lambda} E(N_s) = \frac{1}{30} \text{ hour or } 2 \text{ minutes}$

$\therefore$  Number of customer in the shop = 7.01 and total waiting time in the shop for each customer = 14.02 minutes

*Case (ii):* For the supply section,  $\lambda = 30$ ,  $\mu_2 = 12$  and  $s = 4$

$$P_0^{-1} = \sum_{r=0}^3 \frac{1}{\underline{r}} \cdot \left(\frac{5}{2}\right)^r + \left\{ \frac{\left(\frac{5}{2}\right)^4}{\underline{4} \cdot \left(1 - \frac{5}{8}\right)} \right\}$$

$$= 1 + \frac{5}{2} + \frac{25}{8} + \frac{125}{48} + \frac{625}{144} = \frac{1954}{144} \text{ or } \frac{977}{72}$$

$$E(N_s) = \frac{1}{4 \cdot [4]} \cdot \frac{\left(\frac{5}{2}\right)^5}{\left(\frac{3}{8}\right)^2} \times \frac{72}{977} + \frac{5}{2} = 3.03$$

$$E(W_s) = \frac{1}{30} \times 3.03 \text{ hour or } 6.06 \text{ minutes}$$

For the payment section,  $\lambda = 30$ ,  $\mu_2 = 40$

$$E(N_s) = \frac{\lambda}{\mu_2 - \lambda} = \frac{30}{40 - 30} = 3$$

$$E(W_s) = \frac{1}{30} \times 3 \text{ hour or } 6 \text{ minutes.}$$

$\therefore$  Number of customers in the shop = 6.03 and the total waiting time in the shop for each customer = 12.06 minutes. Though there is some advantage due to the appointment of one more customer, the owner's wish could not be accomplished fully.

### Example 6

In the railway reservation section of a city junction, there is enough space for the customers to assemble, form a queue and fill up the reservation forms. There are 5 reservation counters in front of which also there is enough space for the customers to wait. Customers arrive at the reservation section at the rate of 50/hour according to Poisson process, take 1 minute each on the average to fill up the forms and then move to the reservation counter section. Each reservation clerk takes 5 minutes on the average to complete the business of a customer in an exponential manner. (a) Find the probability that a customer has to wait to get the service in the reservation counter section, (b) find the total waiting time for a customer in the entire reservation section, and (c) find the total number of customers in the entire section. Assume that only those who have the filled up reservation forms will be allowed into the counter section.

The queueing system in the form-filling portion is only a  $M/M/1$  model, since each customer is served by himself or herself (viz., one server)

For this system,  $\lambda = 50/\text{hour}$  and  $\mu = 60/\text{hour}$ .

$$E(N_s) = \frac{\lambda}{\mu - \lambda} = \frac{50}{60 - 50} = 5;$$

$$E(W_s) = \frac{1}{\mu - \lambda} = \frac{1}{10} \text{ hour or } 6 \text{ min}$$

The queuing system in the reservation counter section in an  $M/M/s$  model with  $\lambda = 50/\text{hour}$ ,  $\mu = 12/\text{hour}$  and  $s = 5$  [Since the output of the  $M/M/1$  system is the same as the input of that system, by Burke's theorem and the output of this system (namely, 50/hour) becomes the input of the  $M/M/s$  system.]

- (a)  $P(\text{a customer has to wait in the counter section})$

$$= \frac{\left(\frac{\lambda}{\mu}\right)^s \cdot P_0}{\frac{s}{s} \left(1 - \frac{\lambda}{\mu s}\right)},$$

$$\text{where } P_0^{-1} = \sum_{r=0}^{s-1} \frac{1}{r!} \cdot \left(\frac{\lambda}{\mu}\right)^r + \left\{ \frac{(\lambda/\mu)^s}{\frac{s}{s} \left(1 - \frac{\lambda}{\mu s}\right)} \right\}$$

$$= \sum_{r=0}^{s-1} \frac{1}{r!} \left(\frac{25}{6}\right)^r + \frac{\left(\frac{25}{6}\right)^s}{\frac{5}{5} \cdot \left(1 - \frac{5}{6}\right)}$$

$$= 1 + \frac{25}{6} + \frac{1}{2} \cdot \left(\frac{25}{6}\right)^2 + \frac{1}{6} \left(\frac{25}{6}\right)^3 + \frac{1}{24} \left(\frac{25}{6}\right)^4 + \frac{\left(\frac{25}{6}\right)^5}{20}$$

$$= 164.0490$$

$$\therefore P_0 = 0.0061$$

$$\therefore \text{Required probability} = \frac{\left(\frac{25}{6}\right)^s \times 0.0061}{\frac{5}{5} \cdot \frac{1}{6}} = 0.3824.$$

- (b) Average number of customers in the counter section is given by

$$E(N_s) = \frac{1}{s!} \cdot \frac{(\lambda/\mu)^{s+1}}{(1 - \lambda/\mu s)^2} P_0 + \frac{\lambda}{\mu}$$

$$= \frac{1}{5 \cdot 5!} \cdot \frac{\left(\frac{25}{6}\right)^6}{\left(1 - \frac{5}{6}\right)^2} \times 0.0061 + \frac{25}{6} = 6.0819$$

$\therefore$  Average waiting time of a customer in the counter section is  $E(W_s) = \frac{1}{\lambda} E(N_s) = 0.1216$  hour or 7.298 minutes

$\therefore$  Total waiting time of a customer in the entire reservation room =  $6 + 7.298 = 13.298$  minutes.

- (c) Total number of customers in the entire reservation room =  $5 + 6.0819 = 11.0819$ .

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## Example 7

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In a book shop, there are two sections, one for textbooks and the other for notebooks. Customers from outside arrive at the textbook section at a Poisson rate of 4 per hour and at the notebook section at a Poisson rate of 3 per hour. The service rates of the T.B. section and N.B. section are respectively 8 and 10 per hour. A customer upon completion of service at T.B. section is equally likely to go to the N.B. section or to leave the book shop, whereas a customer upon completion of service at N.B. section will go to the T.B. section with probability  $\frac{1}{3}$  and will leave the book shop otherwise. Find the joint steady-state probability that there are 4 customers in the T.B. section and 2 customers in the N.B. section. Find also the average number of customers in the book shop and the average waiting time of a customer in the shop. Assume that there is only one salesman in each section.

The system given in this problem is a Jackson's open queueing system.

Let  $\lambda_1$  and  $\lambda_2$  be the total (resultant) arrival rates of  $S_1$  (T.B. section) and  $S_2$  (N.B. section).

Jackson's flow balance equations for this open model are  $\lambda_j = r_j + \sum_{i=1}^k \lambda_i P_{ij}$ ,  
 $j = 1, 2, \dots, k$ .

$$\text{Thus } \lambda_1 = 4 + \sum_{i=1}^2 \lambda_i P_{i1} = 4 + \lambda_1 P_{11} + \lambda_2 P_{21}$$

$$\text{i.e., } \lambda_1 = 4 + \frac{1}{3} \lambda_2 \quad (\text{Since } P_{11} = 0) \quad (1)$$

$$\lambda_2 = 3 + \sum_{i=1}^2 \lambda_i P_{i2} = 3 + \lambda_1 P_{12} + \lambda_2 P_{22}$$

$$\text{i.e., } \lambda_2 = 3 + \frac{1}{2} \lambda_1 \quad (\text{Since } P_{22} = 0) \quad (2)$$

Solving (1) and (2), we get  $\lambda_1 = 6$  and  $\lambda_2 = 6$

$P(n_1 \text{ customers in } S_1 \text{ and } n_2 \text{ customers in } S_2)$

$$\begin{aligned} &= P(n_1, n_2) = \left( \frac{\lambda_1}{\mu_1} \right)^{n_1} \left( 1 - \frac{\lambda_1}{\mu_1} \right) \left( \frac{\lambda_2}{\mu_2} \right)^{n_2} \left( 1 - \frac{\lambda_2}{\mu_2} \right) \\ \therefore &P(4, 2) = \left( \frac{6}{8} \right)^4 \cdot \left( 1 - \frac{6}{8} \right) \cdot \left( \frac{6}{10} \right)^2 \cdot \left( 1 - \frac{6}{10} \right) \\ &= \left( \frac{3}{4} \right)^4 \cdot \frac{1}{4} \cdot \left( \frac{3}{5} \right)^2 \cdot \frac{2}{5} = 0.0114 \end{aligned}$$

$$\begin{aligned}
 E(N_s) &= [E(N_s)]_{S_1} + [E(N_s)]_{S_2} \\
 &= \frac{\lambda_1}{\mu_1 - \lambda_1} + \frac{\lambda_2}{\mu_2 - \lambda_2} = \frac{6}{8-6} + \frac{6}{10-6} = 4.5 \\
 E(W_s) &= \frac{1}{\lambda} E(N_s), \text{ where } \lambda = r_1 + r_2 = 4 + 3 = 7 \\
 \therefore E(W_s) &= \frac{9}{2 \times 7} = \frac{9}{14} \text{ hour or } 38.6 \text{ minutes.}
 \end{aligned}$$

### Example 8

In a network of 3 service stations 1, 2, 3, customers arrive at 1, 2, 3, from outside, in accordance with Poisson process having rates 5, 10, 15 respectively. The service times at the 3 stations are exponential with respective rates 10, 50, 100. A customer completing service at station 1 is equally likely to (a) go to station 2, (b) go to station 3, or (c) leave the system. A customer departing from service at station 2 always goes to station 3. A departure from service at station 3 is equally likely to go to station 2 or leave the system.

- (i) What is the average number of customers in the system, consisting of all the three stations?
- (ii) What is the average time a customer spends in the system?

The system given is a Jackson's open queue system.

Let  $\lambda_1, \lambda_2, \lambda_3$ , be the resultant arrival rates at  $S_1, S_2$  and  $S_3$  respectively.

Jackson's flow balance equations are

$$\lambda_j = r_j + \sum_{i=1}^3 \lambda_i p_{ij}, j = 1, 2, 3 \quad (1)$$

We note that  $P_{12} = \frac{1}{3}, P_{13} = \frac{1}{3}, P_{23} = 1, P_{32} = \frac{1}{2}$

Putting  $j = 1$  in (1), we get

$$\lambda_1 = 5 + \lambda_1 P_{11} + \lambda_2 P_{21} + \lambda_3 P_{31} \quad (2)$$

i.e.,

$$\lambda_1 = 5 \quad (2)$$

(since  $P_{11} = P_{21} = P_{31} = 0$ )

Putting  $j = 2$  in (1), we get

$$\lambda_2 = r_2 + \lambda_1 P_{12} + \lambda_2 P_{22} + \lambda_3 P_{32}$$

i.e.,

$$\lambda_2 = 10 + \frac{1}{3} \lambda_1 + \frac{1}{2} \lambda_3 \quad (\text{since } P_{22} = 0)$$

i.e.,

$$\lambda_2 = \frac{35}{3} + \frac{1}{2} \lambda_3 \quad (3)$$

by using (2)

Putting  $j = 3$  in (1), we get

$$\lambda_3 = r_3 + \lambda_1 P_{13} + \lambda_2 P_{23} + \lambda_3 P_{33}$$

i.e.,  $\lambda_3 = 15 + \frac{1}{3} \lambda_1 + \lambda_2$  (since  $P_{33} = 0$ )

i.e.,  $\lambda_3 = \frac{50}{3} + \lambda_2$  (4)

by using (2)

Solving (3) and (4), we get  $\lambda_2 = 40$  and  $\lambda_3 = \frac{170}{3}$

$$\begin{aligned} \text{(i)} \quad E(N_s) &= \frac{\lambda_1}{\mu_1 - \lambda_1} + \frac{\lambda_2}{\mu_2 - \lambda_2} + \frac{\lambda_3}{\mu_3 - \lambda_3} \\ &= \frac{5}{10 - 5} + \frac{40}{50 - 40} + \frac{170/3}{100 - \frac{170}{3}} \\ &= 1 + 4 + \frac{170}{130} \text{ or } 6 \frac{4}{13} \end{aligned}$$

$$\begin{aligned} \text{(ii)} \quad E(W_s) &= \frac{E(N_s)}{\lambda}, \text{ where } \lambda = r_1 + r_2 + r_3 = 30 \\ &= \frac{82}{13 \times 30} = \frac{41}{195} \end{aligned}$$

### Example 9

Customers arrive at a service centre consisting of 2 the service points  $S_1$  and  $S_2$  at a Poisson rate of 35/hour and form a queue at the entrance. On studying the situation at the centre, they decide to go to either  $S_1$  or  $S_2$ . The decision making takes on the average 30 seconds in an exponential fashion. Nearly 55% of the customers go to  $S_1$ , that consists of 3 parallel servers and the rest go to  $S_2$ , that consist of 7 parallel servers. The service times at  $S_1$  are exponential with a mean of 6 minutes and those at  $S_2$  with a mean of 20 minutes. About 2% of customers, on finishing service at  $S_1$  go to  $S_2$  and about 1% of customers, on finishing service at  $S_2$  go to  $S_1$ . Find the average queue sizes in front of each node and the total average time a customer spends in the service centre.

The system given is a Jackson's open queue system with 3 nodes. At the first node, we have a single server queue and at the second and third nodes, we have multiple server queues.

Let  $\lambda_1, \lambda_2, \lambda_3$  be the resultant arrival rates at the 3 nodes.  $r_1 = 35, r_2 = 0, r_3 = 0$  are the arrival rates at the 3 nodes from outside.

Jackson's flow balance equations are given by

$$\lambda_j = r_j + \sum_{i=1}^3 \lambda_i P_{ij}; j = 1, 2, 3 \quad (1)$$

Noting that  $P_{12} = 0.55$ ,  $P_{13} = 0.45$ ,  $P_{23} = 0.02$  and  $P_{32} = 0.01$  and using them in (1), we get

$$\lambda_1 = 35 \quad (2)$$

$$\lambda_2 = 0.55 \lambda_1 + 0.01 \lambda_3, \text{ i.e., } \lambda_1 - 0.01 \lambda_3 = 19.25 \quad (3)$$

$$\lambda_3 = 0.45 \lambda_1 + 0.02 \lambda_2, \text{ i.e., } 0.02 \lambda_2 - \lambda_3 = -15.75 \quad (4)$$

Solving (3) and (4), we get  $\lambda_2 = 19.4114$  and  $\lambda_3 = 16.1382$

Also  $\mu_1 = 120$ ,  $\mu_2 = 10$  and  $\lambda_3 = 3$

At node 1, we have a single server system

$$\therefore (L_q)_1 = \frac{\lambda_1^2}{\mu_1(\mu_1 - \lambda_1)} = \frac{35 \times 35}{120 \times 85} = 0.1201$$

At node 2, have a 3-server queue system.

$$\text{For this system, } P_0^{-1} = \sum_{r=0}^2 \frac{1}{r!} \left( \frac{\lambda_2}{\mu_2} \right)^r + \left\{ \frac{\left( \frac{\lambda_2}{\mu_2} \right)^3}{[3 \left( 1 - \frac{\lambda_2}{3\mu_2} \right)]} \right\}$$

$$= 1 + 1.9411 + \frac{1}{2} \times (1.9411)^2 + \frac{(1.9411)^3}{6 \times (1 - 0.6470)}$$

$$= 8.2782$$

$$\therefore P_0 = 0.1208$$

$$\text{Now } (L_q)_2 = \frac{1}{s_2 \cdot \left[ s_2 \left( 1 - \frac{\lambda_2}{\mu_2 s_2} \right)^2 \right]} \cdot P_0$$

$$= \frac{1}{3 \times 6} \times \frac{(1.9411)^4}{(0.3530)^2} \times 0.1208 = 0.7646$$

At node 3, we have a 7 server queue system.

$$\text{For this system, } P_0^{-1} = \sum_{r=0}^6 \frac{1}{r!} \left( \frac{\lambda_3}{\mu_3} \right)^r + \frac{(\lambda_3 / \mu_3)^7}{[s_3 \left( 1 - \frac{\lambda_3}{s_3 \mu_3} \right)]}$$

$$= 1 + 5.3794 + \frac{1}{2} \times (5.3794)^2 + \frac{1}{6} \times (5.3794)^3$$

$$+ \frac{1}{24} \times (5.3794)^4 + \frac{1}{120} \times (5.3794)^5$$

$$+ \frac{1}{720} \times (5.3794)^6 + \frac{(5.3794)^7}{[7 \times 0.2315]}$$

$$\begin{aligned}
&= 1 + 5.3794 + 14.4690 + 25.9448 + 34.8919 \\
&\quad + 37.5395 + 33.6566 + 111.7263 \\
&= 2640.6075 \\
\therefore P_0 &= 0.0038
\end{aligned}$$

Now,

$$\begin{aligned}
(L_q)_3 &= \frac{1}{s_3 \underline{s_3}} \cdot \frac{\left(\frac{\lambda_3}{\mu_3}\right)^{s_3+1}}{\left(1 - \frac{\lambda_3}{s_3 \mu_3}\right)^2} \cdot P_0 \\
&= \frac{1}{7 \times \underline{7}} \times \frac{(5.3794)^8}{(1 - 0.7685)^2} \times 0.0038 \\
&= 1.4094
\end{aligned}$$

$$\begin{aligned}
(L_s)_1 &= (L_q)_1 + \frac{\lambda_1}{\mu_1} \text{ or } \frac{\lambda_1}{\mu_1 - \lambda_1} = 0.4118 \\
(L_s)_2 &= (L_q)_2 + \frac{\lambda_2}{\mu_2} = 0.7646 + 1.9411 = 2.7057 \\
(L_s)_3 &= (L_q)_3 + \frac{\lambda_3}{\mu_3} = 1.4094 + 5.3794 = 6.7888
\end{aligned}$$

Total waiting time for a customer in the system is given by

$$\begin{aligned}
E(W_s) &= \frac{(L_s)_1 + (L_s)_2 + (L_s)_3}{r_1 + r_2 + r_3} \\
&= 0.2830 \text{ hour or 17 minutes nearly}
\end{aligned}$$

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### Example 10

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There are 2 clerks in a bank, one processing housing loan applications and the other processing agricultural loan applications. While processing, they get doubts according to an exponential distribution each with a mean of  $\frac{1}{2}$ . To get clarifications, a clerk goes to the Deputy manager with probability  $\frac{3}{4}$  and to the senior manager with probability  $\frac{1}{4}$ . After completing the job with D.M. a clerk goes to S.M. with probability  $\frac{1}{3}$  and returns to his seat otherwise. Completing the job with S.M., a clerk always returns to his seat. If the D.M. clarifies the doubts

and advises a clerk according to an exponential distribution with parameter 1 and the S.M. with parameter 3, find

- (a) the steady-state probabilities  $P(n_1, n_2, n_3)$  for all possible values of  $n_1, n_2, n_3$
- (b) the probability that both the managers are idle
- (c) the probability that at least one manager is idle.

The situation in this problem is a Jackson's closed network.

We assume that the processing point by the clerks is node 1 and the service points by D.M. and S.M. are node 2 and node 3 respectively.

$$\text{Now } \lambda = \mu_1 = 2, \mu_2 = 1, \mu_3 = 3, \text{ and } P_{12} = \frac{3}{4}, P_{13} = \frac{1}{4}, P_{23} = \frac{1}{3}$$

$P(n_1, n_2, n_3)$  means the probability that there are  $n_1$  clerks in their seats,  $n_2$  customers in D.M.'s room and  $n_3$  customers in S.M.'s room.

There are 6 possible combinations for the values of  $n_1, n_2, n_3$ , namely  $(2, 0, 0), (0, 2, 0), (0, 0, 2), (1, 1, 0), (1, 0, 1)$  and  $(0, 1, 1)$ , since  $n_1 + n_2 + n_3 = 2$ .

The flow balance equations  $\lambda_j = \sum_{i=1}^3 \lambda_i P_{ij}, (j = 1, 2, 3)$  can be put in the matrix form as

$$(\lambda_1, \lambda_2, \lambda_3) = (\lambda_1, \lambda_2, \lambda_3) \begin{pmatrix} 0 & P_{12} & P_{13} \\ P_{21} & 0 & P_{23} \\ P_{31} & P_{32} & 0 \end{pmatrix} \quad (1)$$

Putting  $\rho_i = \frac{\lambda_i}{\mu_i}$  and noting that  $\rho_i = \frac{\lambda_1}{\lambda}$  in (1), we get  $(\lambda\rho_1, \rho_2 \mu_2, \rho_3 \mu_3)$

$$= (\lambda\rho_1, \rho_2 \mu_2, \rho_3 \mu_3) \begin{pmatrix} 0 & \frac{3}{4} & \frac{1}{4} \\ \frac{2}{3} & 0 & \frac{1}{3} \\ 1 & 0 & 0 \end{pmatrix}$$

$$\text{i.e., } \frac{2}{3}\rho_2 + 3\rho_3 = 2\rho_1 \quad (2)$$

$$\frac{3}{2}\rho_1 = \rho_2 \quad (3)$$

$$\frac{1}{2}\rho_1 + \frac{1}{3}\rho_2 = 3\rho_3 \quad (4)$$

Equations (2), (3) and (4) are homogeneous linear equations in  $\rho_1, \rho_2$  and  $\rho_3$ . To get one solution we put  $\rho_2 = 1$  and solve the equations.

We then get

$$\rho_1 = \frac{2}{3}, \rho_2 = 1 \text{ and } \rho_3 = \frac{2}{9}$$

Since there are 2 clerks (customers) at node 1, the system there is a 2-server queueing model.

$$\therefore P(n_1, n_2, n_3) = C_N \frac{\rho_1^{n_1}}{a_{1(n_1)}}, \rho_2^{n_2} \cdot \rho_3^{n_3},$$

where

$$N = n_1 + n_2 + n_3 = 2$$

$$\text{Now, } C_N^{-1} = \sum_{n_1+n_2+n_3=2} \frac{\rho_1^{n_1}}{|n_1|} \cdot \rho_2^{n_2} \cdot \rho_3^{n_3}, \text{ since}$$

$$a_1(n_1) = \begin{cases} |n_1|, & \text{when } n_1 < s_1 \\ |s_1| \cdot s_1^{n_1-s_1} & \text{when } n_1 \geq s_1 \end{cases}$$

for all possible combinations of  $n_1, n_2, n_3$  such that  $n_1, n_2, n_3 = 2$

$$\begin{aligned} \text{Thus, } C_N^{-1} &= \frac{\rho_1^2}{|2|} \cdot \rho_2^0 \cdot \rho_3^0 + \frac{\rho_1^0}{|0|} \cdot \rho_2^2 \cdot \rho_3^0 + \frac{\rho_1^0}{|0|} \cdot \rho_2^0 \cdot \rho_3^2 \\ &\quad + \frac{\rho_1^1}{|1|} \rho_2^1 \cdot \rho_3^0 + \frac{\rho_1^1}{|1|} \cdot \rho_2^0 \cdot \rho_3^1 + \frac{\rho_1^0}{|0|} \cdot \rho_2^1 \cdot \rho_3^1 \end{aligned}$$

$$= \frac{2}{9} + 1 + \frac{4}{81} + \frac{2}{3} + \frac{4}{27} + \frac{2}{9} = \frac{187}{81} \text{ or } 2.3086$$

$$\therefore C_N = \frac{1}{2.3086} = 0.4332$$

$$\text{Now, } P(2, 0, 0) = 0.4332 \times \frac{(2/3)^2}{2} = 0.0962; P(0, 2, 0) = 0.4332$$

$$P(0, 0, 2) = 0.0214; P(1, 1, 0) = 0.2888; P(1, 0, 1) = 0.0642;$$

$$P(0, 1, 1) = 0.0962$$

$$P(\text{both the managers are idle}) = P(2, 0, 0) = 0.0962$$

$$\begin{aligned} P(\text{at least one manager is idle}) &= P(1, 1, 0) + P(1, 0, 1) + P(2, 0, 0) \\ &= 0.2888 + 0.0642 + 0.0962 \\ &= 0.4492 \end{aligned}$$

### Example 11

Assuming that there is only one clerk (i.e.,  $N = 1$ ) in the loan section in the previous example, compute the probabilities that he is in nodes 1, 2, and 3. Verify that these are the same as the average lengths of the queues at the nodes. Check the correctness of these values by comparing with the values obtained by using mean-value Aanalysis.

Let  $P_i(1)$  = probability that the single clerk is in node 1.

When  $N = 1 = n_1 + n_2 + n_3$ , we have single server queue systems at all the nodes 1, 2, 3.

$$\begin{aligned}\therefore P(n_1, n_2, n_3) &= C_N \cdot \rho_1^{n_1} \rho_2^{n_2} \rho_3^{n_3}, \text{ where } C_N \equiv C_1 \\ &= C_1 \left(\frac{2}{3}\right)^{n_1} \cdot \left(\frac{2}{9}\right)^{n_3}\end{aligned}$$

$C_1$  is found out more easily as follows:

$$P(1, 0, 0) = C_1 \cdot \frac{2}{3}; P(0, 1, 0) = C_1; P(0, 0, 1) = C_1 \cdot \frac{2}{9}$$

Since the presence of the lonely clerk in the nodes is mutually exclusive and exhaustive,

$$\begin{aligned}P(1, 0, 0) + P(0, 1, 0) + P(0, 0, 1) &= 1 \\ \text{i.e., } \left(\frac{2}{3} + 1 + \frac{2}{9}\right) C_1 &= 1 \quad \therefore C_1 = \frac{9}{17}\end{aligned}$$

$$\therefore P_1(1) = P(1, 0, 0) = \frac{6}{17}; P_2(1) = P(0, 1, 0) = \frac{9}{17};$$

$$P_3(1) = P(0, 0, 1) = \frac{2}{17}$$

Now,  $L_1(1) = \text{Average length of the queueing system at node 1}$

$$\begin{aligned}&= \sum_{r=0}^1 r P \text{ (r clerk in node 1)} \\ &= 0 \times \left(\frac{9}{17} + \frac{2}{17}\right) + 1 \times \frac{6}{17} = \frac{6}{17}\end{aligned}$$

$$\text{Similarly, } L_2(1) = 0 \times \left(\frac{6}{17} + \frac{2}{17}\right) + 1 \times \frac{9}{17} = \frac{9}{17}$$

$$\text{and } L_3(1) = 0 \times \left(\frac{6}{17} + \frac{9}{17}\right) + 1 \times \frac{2}{17} = \frac{2}{17}$$

Thus,  $P_r(1) = L_r(1)$ , for  $r = 1, 2, 3$ .

Now, let us use the MVA algorithm and compute  $L_r(1)$

The flow balance equations of Example (10) are

$$(\lambda_1, \lambda_2, \lambda_3) = (\lambda_1, \lambda_2, \lambda_3) \begin{pmatrix} 0 & \frac{3}{4} & \frac{1}{4} \\ \frac{2}{3} & 0 & \frac{1}{3} \\ 1 & 0 & 0 \end{pmatrix} \quad (1)$$

Putting  $v_i = \frac{\lambda_i}{\mu_2}$  and taking  $v_1 = v_2 = 1$  in (1), we have

$$(v_1, 1, v_3) = (v_1, 1, v_3) \begin{pmatrix} 0 & \frac{3}{4} & \frac{1}{4} \\ \frac{2}{3} & 0 & \frac{1}{3} \\ 1 & 0 & 0 \end{pmatrix}$$

i.e.,

$$v_1 = \frac{2}{3} + v_3 \quad (2)$$

$$1 = \frac{3}{4} v_1 \quad (3)$$

and

$$v_3 = \frac{1}{4} v_1 + \frac{1}{3} \quad (4)$$

Solving (2), (3) and (4), we get  $v_1 = \frac{4}{3}$ ,  $v_2 = 1$ ,  $v_3 = \frac{2}{3}$

Assuming  $L_1(0) = L_2(0) = L_3(0) = 0$  and taking  $r = 1$ , from step 3 (a) of M.V.A. algorithm, we have

$$W_1(1) = \frac{1}{\mu_1} = \frac{1}{\lambda} = \frac{1}{2}; W_2(1) = \frac{1}{\mu_2} = 1;$$

$$W_3(1) = \frac{1}{\mu_3} = \frac{1}{3}$$

For  $r = 1$ , from step 3 (b) of M.V.A.A.,

$$\lambda_2(1) = \frac{1}{\sum_{i=1}^3 v_i W_i(1)} = \frac{1}{\frac{4}{3} \times \frac{1}{2} + 1 \times 1 + \frac{2}{3} \times \frac{1}{3}} = \frac{9}{17}$$

For  $r = 1$ , from step 3 (c) of M.V.A.A.,

$$\lambda_1(1) = v_1 \lambda_2(1) = \frac{4}{3} \times \frac{9}{17} = \frac{12}{17}$$

$$\lambda_3(1) = v_3 \lambda_2(1) = \frac{2}{3} \times \frac{9}{17} = \frac{6}{17}$$

From  $r = 1$ , from step 3 (d) of M.V.A.A.,

$$L_1(1) = \lambda_1(1) \cdot W_1(1) = \frac{12}{17} \times \frac{1}{2} = \frac{6}{17}$$

$$L_2(1) = \lambda_2(1) \cdot W_2(1) = \frac{9}{17} \times 1 = \frac{9}{17}$$

$$L_3(1) = \lambda_3(1) \cdot W_3(1) = \frac{6}{17} \times \frac{1}{3} = \frac{2}{17}$$

These values of  $L_i(1)$  agree with the corresponding values obtained otherwise.

### Example 12

Use M.V.A. algorithms for (a) single server, (b) multiple servers to compute  $L_1(2)$ ,  $L_2(2)$  and  $L_3(2)$ , where  $L_i(N)$  denotes the average length of the queueing system at node  $i$ , assuming that there are  $N = 2$  clerks in the system of Example (10).

*Case (i):* For  $r = 2$ , from step 3(a) of M.V.A.A.,

$$W_1(2) = \frac{1 + L_1(1)}{\mu_1} = \frac{1 + \frac{6}{17}}{2} = \frac{23}{34}; W_2(2) = \frac{1 + L_2(1)}{\mu_2} = \frac{1 + \frac{9}{17}}{1} = \frac{26}{17};$$

$$W_3(2) = \frac{1 + L_3(1)}{\mu_3} = \frac{1 + \frac{2}{17}}{3} = \frac{19}{51}$$

For  $r = 2$ , from step 3(b) of M.V.A.A.,

$$\lambda_2(2) = \frac{2}{\frac{4}{3} \times \frac{23}{24} + 1 \times \frac{26}{17} + \frac{2}{3} \times \frac{19}{27}} = \frac{153}{205}$$

For  $r = 2$ , from step 3(c) of M.V.A.A.,

$$\lambda_1(2) = v_1 \lambda_2(2) = \frac{4}{3} \times \frac{153}{205} = \frac{204}{205}$$

$$\lambda_3(2) = v_3 \lambda_2(2) = \frac{2}{3} \times \frac{153}{205} = \frac{102}{205}$$

For  $r = 2$ , from step 3(d) of M.V.A.A.,

$$L_1(2) = \lambda_1(2) \cdot W_1(2) = \frac{204}{205} \times \frac{23}{34} = \frac{138}{205} = 0.6732$$

$$L_2(2) = \lambda_2(2) \cdot W_2(2) = \frac{153}{205} \times \frac{26}{17} = \frac{234}{205} = 1.1415$$

$$L_3(2) = \lambda_3(2) \cdot W_3(2) = \frac{102}{205} \times \frac{19}{51} = \frac{38}{205} = 0.1853$$

It can be verified that  $L_1(2) + L_2(2) + L_3(2) = 2$ .

*Case (ii):* When  $r = 1 = N$ , step 3 (a) of the multiple server M.V.A. algorithm reduces to the step 3 (a) of the single server M.V.A. algorithm

Also steps 3(b), 3(c) and 3(d) hold good for  $r = 1$

∴ proceeding as in Example 11, we get

$$L_1(1) = \frac{6}{17}, L_2(1) = \frac{9}{17} \text{ and } L_3(1) = \frac{2}{17}$$

Step 3(e) of the algorithm gives

$$P_1(1, 1) = \frac{\lambda_1(1)}{\alpha_1(1) \cdot \mu_1} \cdot P_1(0, 0) = \frac{12/17}{1 \times 2} = \frac{6}{17}$$

$$P_2(1, 1) = \frac{\lambda_2(1)}{\alpha_2(1) \cdot \mu_2} \cdot P_2(0, 0) = \frac{9/17}{1 \times 1} = \frac{9}{17}$$

$$P_3(1, 1) = \frac{\lambda_3(1)}{\alpha_3(1) \cdot \mu_3} = \frac{6/17}{1 \times 3} = \frac{2}{17}$$

Now  $P_1(0, 1) = P(0, 1, 0) + P(0, 0, 1) = \frac{9}{17} + \frac{2}{17} = \frac{11}{17}$

$$P_2(0, 1) = P(0, 0, 1) + P(1, 0, 0) = \frac{2}{17} + \frac{6}{17} = \frac{8}{17}$$

$$P_3(0, 1) = P(0, 1, 0) + P(1, 0, 0) = \frac{9}{17} + \frac{6}{17} = \frac{15}{17}$$

We now proceed to the second iteration with  $r = 2$ . Step 3(a) gives

$$\begin{aligned} W_1(2) &= \frac{1}{2\mu_1} \{1 + L_1(1) + (2-1) P_1(0, 1)\} \\ &= \frac{1}{4} \left( 1 + \frac{6}{17} + \frac{11}{17} \right) = \frac{1}{2} = 0.5 \end{aligned}$$

[Why  $W_1(2)$  values got by the two algorithms differ?]

$$W_2(2) = \frac{1}{\mu_2} \{1 + L_2(1) + 0\} = 1 + \frac{9}{17} = \frac{26}{17} = 1.5294$$

$$W_3(2) = \frac{1}{\mu_3} \{1 + L_3(1) + 0\} = \frac{1}{3} \left( 1 + \frac{2}{17} \right) = \frac{19}{51} = 0.3725$$

Step 3(b) gives

$$\begin{aligned} \lambda_2(2) &= \frac{2}{v_1 W_1(2) + v_2 W_2(2) + v_3 W_3(2)} \\ &= \frac{2}{\frac{4}{3} \times 0.5 + 1 \times 1.5294 + \frac{2}{3} \times 0.3725} = 0.8182 \end{aligned}$$

Step 3(c) gives

$$\lambda_1(2) = v_1 \cdot \lambda_2(2) = \frac{4}{3} \times 0.8182 = 1.0909$$

$$\lambda_3(2) = v_3 \lambda_2(2) = \frac{2}{3} \times 0.8182 = 0.5455$$

Step 3(d) gives

$$L_1(2) = \lambda_1(2) \cdot W_1(2) = 1.0909 \times 0.5 = 0.5455$$

$$L_2(2) = \lambda_2(2) \cdot W_2(2) = 0.8182 \times 1.5294 = 1.2513$$

$$L_3(2) = \lambda_3(2) \cdot W_3(2) = 0.5455 \times 0.3725 = 0.2032$$

We note that  $\sum_{i=1}^3 L_i(2) = 2$ .

### Exercise 9(B)

#### Part-A (Short-answer Questions)

1. What is meant by queue network?
2. What are the features of a general tandem queue?
3. What do you mean by series queue with blocking?
4. Draw the state transition diagram of a two-stage sequential queue model with blocking for the stage 2.
5. Write down the balance equations for the model given in (4) above.
6. Define a two stage series queue.
7. Draw the state transition diagrams for a 2-stage series queue by taking  $(n_1, n_2)$  as the central state, when  $n_1 \geq 0$  and  $n_2 \geq 0$ .
8. Write down the balance equations for a 2-stages series queue model.
9. Stage Burke's theorem used in queueing theory.
10. Write the formula for the steady-state joint probability for  $m$  and  $n$  customers in the nodes  $S_1$  and  $S_2$  respectively for a 2-stage series queue model.
11. Define an open Jackson network.
12. Write down the flow balance equations for an open Jackson network.
13. Write down the formula for  $P(n_1, n_2 \dots, n_k)$  for a single server Jackson's open network.
14. Define a closed Jackson network.
15. Write down the flow balance equations for a closed Jackson network.
16. Write down the formula for the steady-state probability  $P(n_1, n_2 \dots, n_k)$  for a single server Jackson's closed network.
17. Write down the formula for the steady-state probability  $P(n_1, n_2 \dots, n_k)$  for multiple server Jackson's closed network.
18. State 'Arrival theorem' used in the study of Jackson's network.
19. Write down the M.V.A. algorithm used to find the characteristics of a  $k$ -node, single server per node, Jackson's closed network.
20. Write down the M.V.A. algorithm used to study a  $k$ -node multiple server Jackson's closed network.

**Part-B**

21. For a 2-stage (service point) sequential queue model with blockage, compute the average number of customers in the system and the average time that a customer has to spend in the system, if  $\lambda = 1$ ,  $\mu_1 = 2$  and  $\mu_2 = 1$ .
22. For a 2-stage (service point) sequential queue model with blockage, compute  $L_s$  and  $W_s$ , if  $\lambda = 1$ ,  $\mu_1 = 1$  and  $\mu_2 = 2$ .
23. If, in a 2-stage (service point) sequential queue model with blocking, the arrivals follow Poisson process with parameter  $\lambda$  and the service times are equal with the same parameter  $\mu$ , find the steady-state probability  $P(0, 0)$  in terms of  $\lambda$  and  $\mu$ .
24. There are 2 chairs in a barber shop, each of which is managed by a barber. B1 is a specialist in hair-cutting and B2 is a specialist in shaving and washing. B1 and B2 do their jobs according to exponential distributions with parameters 1 and 3 respectively. Customers who require both hair-cutting and shaving enter the shop only if B1 is free. If a customer has finished his job with B1 goes to B2, if he is free; otherwise he waits in B1 chair until B2 becomes free. Find the average number of customers in the shop and the average amount that an entering customer has to spend in the shop.
25. For a walk-in interview conducted by a company, candidates arrive at the interview hall at a Poisson rate of 6/hour. In the interview hall, verification of certificates is done by a clerk and the personal interview is conducted by an officer in different rooms. Both of them do their jobs in an exponential manner, each taking 6 minutes on the average. If no queue is allowed to form in front of the rooms and a candidate has to wait in the clerk's room when the officer is busy, find the average number of candidates present in the interview hall and the average time spent by a candidate in the hall.
26. On the first day of admission of freshers in an engineering college, freshers, after surrendering the original certificates, join the queue in front of the Principal's chamber for getting actual admission at the rate of 18 per hour in a Poisson manner. After getting the nod from the Principal, they join the queue in front of the cash collection counter to make payment of fees. If the Principal and the cashier do their service in an exponential fashion taking 2.5 and 3 minutes respectively, find the average number of freshers inside the college office and the average time spent by a fresher inside the office. Find also the probability that there are 3 students in front of the Principal's chamber and 4 students in front of the cash counter.
27. In a big factory, there are a large number of operating machines and two sequential repair shops, which do the service of the damaged machines exponentially with respective rates of 1/hour and 2/hour.

- If the cumulative failure rate of all the machines in the factory is  $0.5/\text{hour}$ , find (i) the probability that both repair shops are idle, (ii) the average number of machines in the service section of the factory, and (iii) the average repair time of a machine.
28. In a festival season cracker shop, there are 4 salesmen and 1 manager. Customers arrive at the cracker shop according to Poisson process at the rate of  $15/\text{hour}$ . On entry, the customer stands in the queue, goes to any of the salesmen and gets his requirements and also the bill. After finishing the job with the salesman the customer goes to the manager's counter and joins the queue there. The manager checks the bill and receives the payment. If the service time of each salesman is exponentially distributed with a mean 6 minutes and the same for the manager is 3 minutes, find the average number of customers in the cracker shop and the average waiting time for a customer in it.
29. During peak hours, customers arrive at a super market according to a Poisson process at a mean rate of  $40/\text{hour}$ . In the self-service shopping section it takes  $\frac{3}{4}$  hour for a customer to fill his shopping cart, the filling times being approximately exponentially distributed. Further the check-out times are also approximately exponentially distributed with a mean of 4 minutes, regardless of the particular check-out counter chosen by a customer. Assume that any number of customers (with their carts) can be accommodated both in the shopping section and the check-out lounge.
- (i) What is the minimum number of check-out counters required for operation during peak periods?
  - (ii) If it is decided to add one more than the minimum number of counters required for operation what is the average waiting time in the check-out lounge?
  - (iii) How many customers will be in (a) the lounge and (b) the entire super market on the average?
30. In the outpatient section of a private nursing home, there are 3 doctors and 2 salesmen in the attached medical shop. Patients arrive at the nursing home according to Poisson process with a mean of 30 per hour and join the queue in front of the O.P. section. After getting diagnosed and getting the prescription of medicines by any one of the doctors, the patients go out and join the queue in front of the medical shop. After purchasing the medicines, the patients leave the nursing home. If the service time of each doctor is exponential with a mean of 3 minutes, and that of a salesman is exponential with a mean of 3 minutes find the average number of out patients in the nursing home and the average time spent by a patient in the nursing home.
31. There are 2 counters in a snack stall, one dealing with sweets and savouries and the other with cakes and biscuits. Customers arrive at the stall according to a Poisson process with a mean rate of  $20/\text{hour}$ . 60% of

- the arrivals go to the sweets counter and the rest to the bakery counter. After finishing the purchase in the sweets counter, 20% go to the bakery section and the rest leave the stall. 10% of those who have completed their purchase in the bakery then go to the sweets section and the rest leave the stall. It takes on the average 4 minutes for a customer to complete his purchase in the sweets section and 5 minutes in the bakery section, the service times being exponential. (i) how many, on the average, are in the snack stall? (ii) What is the average waiting time at each counter? (iii) If a customer wants both sweet items and bakery items, how long, on the average, does he spend in the stall?
32. In a departmental store, there are 2 sections, namely grocery section and perishable (vegetables and fruits) section. Customers from outside arrive at the G-section according to a Poisson process at a mean rate of 10/hour and they reach the P-section at a mean rate of 2/hour. The service times at both the sections are exponentially distributed with parameters 15 and 12 respectively. On finishing the job in the G-section, a customer is equally likely to go to the P-section or to leave the store, whereas a customer on finishing his job in the P-section will go to the G-section with probability 0.25 and leave the store otherwise. Assuming that there is only one salesman in each section, find the probability that there are 3 customers in the G-section and 2 customers in the P-section. Find also the average number of customers in the store and the average waiting time of a customer in the store.
33. In a textile shop, there are 3 sections, namely, Gents' (GS), Ladies' (LS) and Kids' (KS) sections, in each of which there is only one salesman. Customers, from outside arrive at the 3 sections in a Poisson manner at the rates of 3, 5 and 6 per hour respectively. The service rates in the 3 sections are exponential with rates 40, 60 and 50 respectively. A customer completing purchase in G.S. will go to L.S. with probability 0.5 and to the K.S. with probability 0.3; otherwise leaves the shop. A customer completing purchase in this L.S. is equally likely to go to the G.S. and to the K.S. A customer completing purchase in the K.S. will go to the G.S. with probability 0.2 and to the L.S. with probability 0.6; otherwise he leaves the shop. Find the average number of customers in the entire textile shop. Find also the average time that a customer has to spend in the shop.
34. In a factory there are 2 machines which are expected to be operational at all times and 2 service men, one of whom will rectify ordinary defects and the other will do the service in respect of serious defects. The machines break down according to an exponential distribution with parameter 2. When a machine breaks down, it has a probability 0.8 of being serviced for ordinary defect and a probability 0.2 of being serviced for serious defect. After service for ordinary defect, a machine

will require service for serious defect with probability 0.4 and return to operation otherwise. After service for serious defect, the machine will always become operational. Treating the operational status of the machines as node 1, find the probability that (i) both the servicemen are idle, and (ii) at least one machine is operational. Assume that the service times of the repairmen are exponentially distributed with parameters 3 and 4 respectively.

35. In problem (34), if there is only one machine in operation, find the average lengths of the machines in the 3 nodes of the system (i) directly, and (ii) using M.V.A. algorithm.
36. In problem (34), assuming that there are 2 machines in operation, find the average number of machines in the 3 nodes of the system, using single server M.V.A. algorithm.
37. Solve problem (36), using multiple server M.V.A. algorithm.

## **ANSWERS**

### **Exercise 9(A)**

7.  $P_0 = 0.00003556$ ;  $E(N_Q) = 4.4$ ;  $E(N_S) = 5.4$ ;  $E(W_Q) = 2.9335$  hours;  $E(W_S) = 3.6$  hours.
8. (a) 3; (b) Rs. 160 per day; (c) Engaging only one mechanic is desirable.
9. 
$$P_0 = \frac{648}{1993}; p_n = \begin{cases} 5C_n \left(\frac{1}{6}\right)^n \cdot P_0, & \text{when } 0 \leq n \leq 2 \\ 5C_n \cdot 2 \cdot n! \left(\frac{1}{12}\right)^n \cdot P_0, & \text{when } 2 \leq n \leq 5 \end{cases}$$
10. (a) 0.04305; (b) 0.16144; (c) 0.911
11. 4.53 machines; 3.11 hours and 78%
12. 5 minutes
13. (a) 29%; (b) 30 minutes
14. 100 minutes
15. 12.5 minutes; 0.832 customer

### **Exercise 9(B)**

21.  $L_S = 1, W_S = \frac{11}{6}$
22.  $L_S = \frac{28}{37}, W_S = \frac{28}{18}$
23.  $P(0, 0) = 2\mu^3/(3\lambda^2 + 4\mu\lambda + 2\mu^2)$

24.  $L = \frac{41}{46}$ ,  $W = \frac{41}{30}$
25.  $L = 0.5883$ ,  $W = 10$  min
26.  $L = 12$ ,  $W = 40$  min,  $P = 0.0069$
27.  $\frac{3}{8}, \frac{4}{3}, \frac{8}{3}$
28.  $L_S = 4.54$ ,  $W_S = 18.18$  min
29. 3, 1.136 min, 0.7573, 3.424
30. 8.9, 17.8 min
31. 14.3885,  $(W_S)_1 = 30.947$  min,  $(W_S)_2 = 43.2339$  min, 43.1505 min
32. 0.0152, 6, 30 min
33. 10, 42.8 min
34. 0.2815, 0.7263
35. (i) 0.5587, 0.2961, 0.1452; (ii) 0.5576, 0.2974, 0.1450
36. 1.2229, 0.5433, 0.2338
37. 1.0048, 0.6958, 0.2294.



# *Chapter* 10

## Tests of Hypotheses

### **Introduction**

Every statistical investigation aims at collecting information about some aggregate or collection of individuals or of their attributes, rather than the individuals themselves. In statistical language, such a collection is called *a population* or *universe*. For example, we have the population of products turned out by a machine, of lives of electric bulbs manufactured by a company etc. A population is finite or infinite, according as the number of elements is finite or infinite. In most situations, the population may be considered infinitely large. A finite subset of a population is called *a sample* and the process of selection of such samples is called *sampling*. The basic objective of the theory of sampling is to draw inference about the population using the information of the sample.

### **Parameters and Statistics**

Generally, in statistical investigations, our ultimate interest will lie in one or more characteristics possessed by the members of the population. If there is only one characteristic of importance, it can be assumed to be a variable  $x$ . If  $x_i$  be the value of  $x$  for the  $i$ th member of the sample, then  $(x_1, x_2, \dots, x_n)$  are referred to as sample observations. Our primary interest will be to know the values of different statistical measures such as mean and variance of the population distribution of  $x$ . Statistical measures, calculated on the basis of population values of  $x$  are called *parameters*. Corresponding measures computed on the basis of sample observations are called *statistics*.

### **Sampling Distribution**

If a number of samples, each of size  $n$ , (i.e. each containing  $n$  elements) are drawn from the same population and if for each sample the value of some statistic, say, mean is calculated, a set of values of the statistic will be obtained.

**Note** The values of the statistic will normally vary from one sample to another, as the values of the population members included in different samples, though drawn from the same population, may be different. These differences in the values of a statistic are said to be due to sampling fluctuations.

If the number of samples is large, the values of the statistic may be classified in the form of a frequency table. The probability distribution of the statistic that would be obtained if the number of samples, each of same size were infinitely large is called the *sampling distribution* of the statistic. If we adopt random sampling technique that is the most popular and frequently used method of sampling [the discussion of which is beyond the scope of this book], the nature of the sampling distribution of a statistic can be obtained theoretically, using the theory of probability, provided the nature of the population distribution is known.

Like any other distribution, a sampling distribution will have its mean, standard deviation and moments of higher order. The standard deviation of the sampling distribution of a statistic is of particular importance in tests of hypotheses and is called the *standard error* of the statistic.

## Estimation and Testing of Hypotheses

In sampling theory, we are primarily concerned with two types of problems which are given below:

- (i) Some characteristic or feature of the population in which we are interested may be completely unknown to us and we may like to make a guess about this characteristic entirely on the basis of a random sample drawn from the population. This type of problem is known as the problem of *estimation*.
- (ii) Some information regarding the characteristic or feature of the population may be available to us and we may like to know whether the information is tenable (or can be accepted) in the light of the random sample drawn from the population and if it can be accepted, with what degree of confidence it can be accepted. This type of problem is known as the problem of *testing of hypotheses*.

## Tests of Hypotheses and Tests of Significance

When we attempt to make decisions about the population on the basis of sample information, we have to make assumptions or guesses about the nature of the population involved or about the value of some parameter of the population. Such assumptions, which may or may not be true, are called *statistical hypotheses*. Very often, we set up a hypothesis which assumes that there is no significant difference between the sample statistic and the corresponding population parameter or between two sample statistics. Such a hypothesis of no difference is

called a *null hypothesis* and is denoted by  $H_0$ . A hypothesis that is different from (or complementary to) the null hypothesis is called an *alternative hypothesis* and is denoted by  $H_1$ . A procedure for deciding whether to accept or to reject a null hypothesis (and hence to reject or to accept the alternative hypothesis respectively) is called *the test of hypothesis*.

If  $\theta_0$  is a parameter of the population and  $\theta$  is the corresponding sample statistic, usually there will be some difference between  $\theta_0$  and  $\theta$  since  $\theta$  is based on sample observations and is different for different samples. Such a difference which is caused due to sampling fluctuations is called *insignificant difference*. The difference that arises due to the reason that either the sampling procedure is not purely random or that the sample has not been drawn from the given population is known as *significant difference*. This procedure of testing whether the difference between  $\theta_0$  and  $\theta$  is significant or not is called *the test of significance*.

## Critical Region and Level of Significance

If we are prepared to reject a null hypothesis when it is true or if we are prepared to accept that the difference between a sample statistic and the corresponding parameter is significant, when the sample statistic lies in a certain region or interval, then that region is called the *critical region* or *region of rejection*. The region complementary to the critical region is called *the region of acceptance*.

In the case of large samples, the sampling distributions of many statistics tend to become normal distributions. If ' $t$ ' is a statistic in large samples, then  $t$  follows a normal distribution with mean  $E(t)$ , which is the corresponding population parameter, and S.D. equal to S.E. ( $t$ ). Hence,  $Z = \frac{t - E(t)}{\text{S.E.}(t)}$  is a standard normal variate, i.e.,  $Z$  (called *the test statistic*) follows a normal distribution with mean zero and S.D. unity.

It is known from the study of normal distribution, that the area under the standard normal curve between  $t = -1.96$  and  $t = +1.96$  is 0.95. Equivalently the area under the general normal curve of ' $t$ ' between  $[E(t) - 1.96 \text{ S.E.}(t)]$  and  $[E(t) + 1.96 \text{ S.E.}(t)]$  is 0.95. In other words, 95 per cent of the values of  $t$  will lie between  $[E(t) \mp 1.96 \text{ S.E.}(t)]$  or only 5 per cent of values of  $t$  will lie outside this interval.

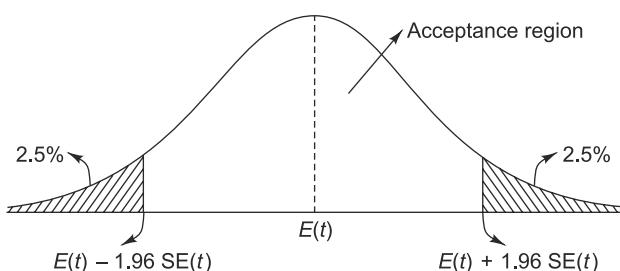


Fig. 10.1

If we are prepared to accept that the difference between  $t$  and  $E(t)$  is significant when  $t$  lies in either of the regions  $[-\infty, E(t) - 1.96 \text{ S.E.}(t)]$  and  $[E(t) + 1.96 \text{ S.E.}(t), \infty]$  then these two regions constitute the critical region of ' $t$ '.

The probability ' $\alpha$ ' that a random value of the statistic lies in the critical region is called *the level of significance* and is usually expressed as a percentage.

From the study of normal distributions, it is known that

$$P\{E(t) - 1.96 \text{ S.E.}(t) < t < E(t) + 1.96 \text{ S.E.}(t)\} = 0.95$$

$$\text{i.e. } P\left\{\left|\frac{t - E(t)}{\text{S.E.}(t)}\right| < 1.96\right\} = 0.95$$

$$\text{i.e. } P\{|Z| > 1.96\} = 0.05 \quad \text{or} \quad 5\%$$

Thus, when  $t$  lies in either of the two regions constituting the critical region given above, the level of significance is 5 per cent.

**Note** *The level of significance can also be defined as the maximum probability with which we are prepared to reject  $H_0$  when it is true. In other words, the total area of the region of rejection expressed as a percentage is called the level of significance.*

(The specification of critical region and the choice of level of significance will depend upon the nature of the problem and is a matter of judgement for those who carry out the investigation. Usually, the levels of significance are taken as 5%, 2% or 1%).

## Errors in Hypotheses Testing

The level of significance is fixed by the investigator and as such it may be fixed at a higher level by his wrong judgement. Due to this, the region of rejection becomes larger and the probability of rejecting a null hypothesis, when it is true, becomes greater. The error committed in rejecting  $H_0$ , when it is really true, is called *Type I error*. This is similar to a good product being rejected by the consumer and hence Type I error is also known as *producer's risk*. The error committed in accepting  $H_0$ , when it is false, is called *Type II error*. As this error is similar to that of accepting a product of inferior quality, it is also known as *consumer's risk*.

The probabilities of committing Type I and II errors are denoted by  $\alpha$  and  $\beta$  respectively. It is to be noted that the probability  $\alpha$  of committing Type I error is the level of significance.

## One-Tailed and Two-Tailed Tests

If  $\theta_0$  is a population parameter and  $\theta$  is the corresponding sample statistic and if we set up the null hypothesis  $H_0 : \theta = \theta_0$ , then the alternative hypothesis which is complementary to  $H_0$  can be any one of the following:

- (i)  $H_1 : \theta \neq \theta_0$ , i.e.  $\theta > \theta_0$  or  $\theta < \theta_0$
- (ii)  $H_1 : \theta > \theta_0$

(iii)  $H_1 : \theta < \theta_0$ .

$H_1$  given in (i) is called a two tailed alternative hypothesis, whereas  $H_1$  given in (ii) is called a right-tailed alternative hypothesis and  $H_1$  given in (iii) is called a left-tailed alternative hypothesis.

When  $H_0$  is tested while  $H_1$  is a one-tailed alternative (right or left), the test of hypothesis is called a *one-tailed test*.

When  $H_0$  is tested while  $H_1$  is two-tailed alternative, the test of hypothesis is called a *two-tailed test*.

The application of one-tailed or two-tailed test depends upon the nature of the alternative hypothesis. The choice of the appropriate alternative hypothesis depends on the situation and the nature of the problem concerned.

## Critical Values or Significant Values

The value of the test statistic  $z$  for which the critical region and acceptance region are separated is called *the critical value* or *the significant value* of  $z$  and denoted by  $z_\alpha$ , when  $\alpha$  is the level of significance. It is clear that the value of  $z_\alpha$  depends not only on  $\alpha$  but also on the nature of alternative hypothesis.

When 
$$z = \frac{t - E(t)}{S \cdot E'(t)}$$
, we have seen that

$$P(|z| < 1.96) = 95 \text{ per cent} \quad P(|z| > 1.96) = 5 \text{ per cent}.$$

Thus,  $z = \pm 1.96$  separate the critical region and the acceptance region at 5% level of significance for a two-tailed test. That is the critical values of  $z$  in this case are  $\pm 1.96$ .

In general, the critical value  $z_\alpha$  for the level of significance  $\alpha$  is given by the equation  $P(|z| > z_\alpha) = \alpha$  for a two-tailed test, by the equation  $P(z > z_\alpha) = \alpha$  for the right-tailed test and by the equation

$P(z < -z_\alpha) = \alpha$  for the left-tailed test. These equations are solved by using the normal tables.

**Note** If  $z_\alpha$  is the critical value of  $z$  corresponding to the level of significance  $\alpha$  in the right-tailed test, then  $P(z > z_\alpha) = \alpha$ .

By symmetry of the standard normal distribution followed by  $z$ ,  $P(z < -z_\alpha) = \alpha$ .

$$\begin{aligned}\therefore P(|z| > z_\alpha) &= P\{(z > z_\alpha) + (z < -z_\alpha)\} \\ &= P\{z > z_\alpha\} + P(z < -z_\alpha) \\ &= 2\alpha\end{aligned}$$

That is,  $z_\alpha$  is the critical value of  $z$  corresponding to the LOS (level of significance)  $2\alpha$ .

Thus, the critical value of  $z$  for a single tailed test (right or left) at LOS ' $\alpha$ ' is the same as that for a two-tailed test of LOS ' $2\alpha$ '.

The critical values for some standard LOS's are given in the following table both for two-tailed and one-tailed tests.

**Table 10.1**

Nature of test \ LOS	1% (.01)	2% (.02)	5% (.05)	10% (.1)
Two-tailed	$ z_\alpha  = 2.58$	$ z_\alpha  = 2.33$	$ z_\alpha  = 1.96$	$ z_\alpha  = 1.645$
Right-tailed	$z_\alpha = 2.33$	$z_\alpha = 2.055$	$z_\alpha = 1.645$	$z_\alpha = 1.28$
Left-tailed	$z_\alpha = -2.33$	$z_\alpha = -2.055$	$z_\alpha = -1.645$	$z_\alpha = -1.28$

## Procedure for Testing of Hypothesis

1. Null hypothesis  $H_0$  is defined.
2. Alternative hypothesis  $H_1$  is also defined after a careful study of the problem and also the nature of the test (whether one-tailed or two tailed) is decided.
3. LOS ' $\alpha$ ' is fixed or taken from the problem if specified and  $z_\alpha$  is noted.
4. The test-statistic  $z = \frac{t - E(t)}{S.E.(t)}$  is computed.
5. Comparison is made between  $|z|$  and  $z_\alpha$ . If  $|z| < z_\alpha$ ,  $H_0$  is accepted or  $H_1$  is rejected, i.e., it is concluded that the difference between  $t$  and  $E(t)$  is not significant at  $\alpha\%$  L.O.S.

On the other hand, if  $|z| > z_\alpha$ ,  $H_0$  is rejected or  $H_1$  is accepted, i.e. it is concluded that the difference between  $t$  and  $E(t)$  is significant at  $\alpha\%$  L.O.S.

## Interval Estimation of Population Parameters

It was pointed out that the objective of the theory of sampling is to estimate population parameters with the help of the corresponding sample statistics. Estimation of a parameter by single value is referred to as *point estimation*, the study of which is beyond the scope of this book. However, an alternative procedure is to give an interval within which the parameter may be supposed to lie. This is called *interval estimation*. The interval within which the parameter is expected to lie is called *the confidence interval* for that parameter. The end points of the confidence interval are called *confidence limits* or *fiducial limits*.

We have already seen that

$$P\{|z| \leq 1.96\} = 0.95$$

$$\text{i.e. } P\left\{\left|\frac{t - E(t)}{S.E.(t)}\right| \leq 1.96\right\} = 0.95$$

$$\text{i.e. } P\{t - 1.96 S.E.(t) \leq E(t) \leq t + 1.96 S.E.(t)\} = 0.95$$

This means that we can assert, with 95% confidence, that the parameter  $E(t)$  will lie between  $t - 1.96 S.E.(t)$  and  $t + 1.96 S.E.(t)$ . Thus,  $\{t - 1.96 S.E.(t), t + 1.96 S.E.(t)\}$  are the 95% confidence limits for  $E(t)$ .

Similarly,  $\{t - 2.58 S.E.(t), t + 2.58 S.E.(t)\}$  is the 99% confidence interval for  $E(t)$ .

## Tests of Significance for Large Samples

It is generally agreed that, if the size of the sample exceeds 30, it should be regarded as a large sample. The tests of significance used for large samples are different from the ones used for small samples for the reason that the following assumptions made for large samples do not hold for small samples:

1. The sampling distribution of a statistic is approximately normal, irrespective of whether the distribution of the population is normal or not.
2. Sample statistics are sufficiently close to the corresponding population parameters and hence may be used to calculate the standard error of the sampling distribution.

### Test 1 Test of significance of the difference between sample proportion and population proportion.

Let  $X$  be the number of successes in  $n$  independent Bernoulli trials in which the probability of success for each trial is a constant  $= P$  (say). Then it is known that  $X$  follows a binomial distribution with mean  $E(X) = nP$  and variance  $V(X) = nPQ$ .

When  $n$  is large,  $X$  follows  $N(nP, \sqrt{nPQ})$ , i.e. a normal distribution with mean  $nP$  and S.D.  $\sqrt{nPQ}$ , where  $Q = 1 - P$ .

$$\therefore \frac{X}{n} \text{ follows } N\left\{\frac{nP}{n}, \sqrt{\frac{nPQ}{n^2}}\right\}$$

Now,  $\frac{X}{n}$  is the proportion of successes in the sample consisting of  $n$  trials, that is denoted by  $p$ . Thus, the sample proportion  $p$  follows  $N\left(P, \sqrt{\frac{PQ}{n}}\right)$ . Therefore, test statistic  $z = \frac{p - P}{\sqrt{\frac{PQ}{n}}}$ .

If  $|z| \leq z_{\alpha/2}$ , the difference between the sample proportion  $p$  and the population proportion  $P$  is not significant at  $\alpha\%$  L.O.S.

**Note** 1. If  $P$  is not known, we assume that  $p$  is nearly equal to  $P$  and hence S.E. ( $p$ ) is taken as  $\sqrt{\frac{pq}{n}}$ . Thus  $z = \frac{p - P}{\sqrt{\frac{pq}{n}}}$ .

2. 95 per cent confidence limits for  $P$  are then given by  $\frac{|P - p|}{\sqrt{\frac{pq}{n}}} \leq 1.96$ , i.e. they are  $\left(p - 1.96 \sqrt{\frac{pq}{n}}, p + 1.96 \sqrt{\frac{pq}{n}}\right)$ .

### Test 2 Test of significance of the difference between two sample proportions.

Let  $p_1$  and  $p_2$  be the proportions of successes in two large samples of size  $n_1$  and  $n_2$  respectively drawn from the same population or from two population with the same proportion  $P$ .

Then,  $p_1$  follows  $N\left(P, \sqrt{\frac{PQ}{n_1}}\right)$  and  $p_2$  follows  $N\left(P, \sqrt{\frac{PQ}{n_2}}\right)$ .

Therefore,  $p_1 - p_2$ , which is a linear combination of two normal variables, also follows a normal distribution.

$$\text{Now } E(p_1 - p_2) = E(p_1) - E(p_2) = P - P = 0$$

$$V(p_1 - p_2) = V(p_1) + V(p_2) \quad (\because \text{the two samples are independent})$$

$$= PQ \left( \frac{1}{n_1} + \frac{1}{n_2} \right)$$

$$\therefore (p_1 - p_2) \text{ follows } N \left\{ 0, \sqrt{PQ \left( \frac{1}{n_1} + \frac{1}{n_2} \right)} \right\}$$

$$\therefore \text{the test statistic } z = \frac{p_1 - p_2}{\sqrt{PQ \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}}.$$

If  $P$  is not known, an unbiased estimate of  $P$  based on both samples, given by

$$\hat{P} = \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2}, \text{ is used in the place of } P.$$

As before, if  $|z| \leq z_{\alpha}$  the difference between the two sample proportions  $p_1$  and  $p_2$  is not significant at  $\alpha$  per cent L.O.S.

**Note** A sample statistic  $\theta$  is said to be an unbiased estimate of the parameter  $\theta_0$ , if  $E(\theta) = \theta_0$ . In the present case,

$$\begin{aligned} E \left\{ \frac{n_1 \hat{p}_1 + n_2 \hat{p}_2}{n_1 + n_2} \right\} &= \frac{1}{n_1 + n_2} \{n_1 E(p_1) + n_2 E(p_2)\} \\ &= \frac{1}{n_1 + n_2} (n_1 P + n_2 P) = P. \end{aligned}$$

$$\therefore \text{an unbiased estimate of } P \text{ is } \left( \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2} \right).$$

**Test 3** Test of significance of the difference between sample mean and population mean.

Let  $X_1, X_2, \dots, X_n$  be the sample observations in a sample of size  $n$ , drawn from a population that is  $N(\mu, \sigma^2)$ .

Then each  $X_i$  follows  $N(\mu, \sigma^2)$ .

It is known that if  $X_i$  ( $i = 1, 2, \dots, n$ ) are independent normal variates with mean  $\mu_i$  and variance  $\sigma_i^2$ , then  $\sum c_i X_i$  is a normal variate with mean  $\mu = \sum c_i \mu_i$  and variance  $\sigma^2 = \sum c_i^2 \sigma_i^2$ .

Now, putting  $c_i = \frac{1}{n}$ ,  $\mu_i = \mu$  and  $\sigma_i = \sigma$ , we get

$$\sum c_i X_i = \frac{1}{n} \sum X_i = \bar{X}, \sum c_i \mu_i = \frac{1}{n} \mu + \frac{1}{n} \mu + \dots + \frac{1}{n} \mu \text{ (n terms)} = \mu$$

$$\begin{aligned} \text{and } \sum c_i^2 \sigma_i^2 &= \frac{1}{n^2} \sigma^2 + \frac{1}{n^2} \sigma^2 + \dots + \frac{1}{n^2} \sigma^2 \text{ (n terms)} \\ &= \frac{\sigma^2}{n}. \end{aligned}$$

Thus, if  $X_i$  are  $n$  independent normal variates with the same mean  $\mu$  and same variance  $\sigma^2$ , then their mean  $\bar{X}$  follows a  $N\left(\mu, \frac{\sigma}{\sqrt{n}}\right)$ . Even if the population, from which the sample is drawn, is non-normal, it is known (from central limit theorem) that the above result holds good, provided  $n$  is large.

$$\therefore \text{the test statistic } z = \frac{\bar{X} - \mu}{\sigma/\sqrt{n}}.$$

As usual, if  $|z| \leq z_{\alpha}$ , the difference between the sample mean  $\bar{X}$  and the population mean  $\mu$  is not significant at  $\alpha\%$  L.O.S.

### Note

1. If  $\sigma$  is not known, the sample S.D. 's' can be used in its place, as s is nearly equal to  $\sigma$  when  $n$  is large.

2. 95% confidence limits for  $\mu$  are given by  $\frac{|\mu - \bar{X}|}{\sigma/\sqrt{n}} \leq 1.96$ , i.e.  $\left( \bar{X} - 1.96 \frac{\sigma}{\sqrt{n}}, \bar{X} + 1.96 \frac{\sigma}{\sqrt{n}} \right)$ , if  $\sigma$  is known. If  $\sigma$  is not known, then the 95% confidence interval is  $\left( \bar{X} - \frac{1.96 s}{\sqrt{n}}, \bar{X} + \frac{1.96 s}{\sqrt{n}} \right)$ .

### Test 4 Test of significance of the difference between the means of two samples.

Let  $\bar{X}_1$  and  $\bar{X}_2$  be the means of two large samples of sizes  $n_1$  and  $n_2$  drawn from two populations (normal or non-normal) with the same mean  $\mu$  and variances  $\sigma_1^2$  and  $\sigma_2^2$  respectively.

Then  $\bar{X}_1$  follows a  $N\left(\mu, \frac{\sigma_1^2}{\sqrt{n_1}}\right)$  and  $\bar{X}_2$  follows a  $N\left(\mu, \frac{\sigma_2^2}{\sqrt{n_1}}\right)$  either exactly or approximately.

$$\therefore \bar{X}_1 - \bar{X}_2 \text{ also follows a normal distribution.}$$

$$E(\bar{X}_1 - \bar{X}_2) = E(\bar{X}_1) - E(\bar{X}_2) = \mu - \mu = 0.$$

$$V(\bar{X}_1 - \bar{X}_2) = V(\bar{X}_1) + V(\bar{X}_2)$$

( $\because \bar{X}_1$  and  $\bar{X}_2$  are independent, as the samples are independent)

$$= \frac{\sigma_1^2}{n_1} + \frac{\sigma_2^2}{n_2}$$

Thus,  $(\bar{X}_1 - \bar{X}_2)$  follows a  $N \left( 0, \sqrt{\frac{\sigma_1^2}{n_1} + \frac{\sigma_2^2}{n_2}} \right)$

$$\therefore \text{the test statistic } z = \frac{\bar{X}_1 - \bar{X}_2}{\sqrt{\frac{\sigma_1^2}{n_1} + \frac{\sigma_2^2}{n_2}}} \quad (1)$$

If  $|z| \leq z_{\alpha/2}$ , the difference between  $(\bar{X}_1 - \bar{X}_2)$  and 0 or the difference between  $\bar{X}_1$  and  $\bar{X}_2$  is not significant at  $\alpha$  per cent L.O.S.

**Note** If the samples are drawn from the same population, i.e., if  $\sigma_1 = \sigma_2 = \sigma$  then

$$z = \frac{\bar{X}_1 - \bar{X}_2}{\sigma \sqrt{\frac{1}{n_1} + \frac{1}{n_2}}} \quad (2)$$

2. If  $\sigma_1$  and  $\sigma_2$  are not known and  $\sigma_1 \neq \sigma_2$ ,  $\sigma_1$  and  $\sigma_2$  can be approximated by the sample S.D.'s  $s_1$  and  $s_2$ . Hence, in such a situation,

$$z = \frac{\bar{X}_1 - \bar{X}_2}{\sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}}} \quad (\beta) [from (1)]$$

3. If  $\sigma_1$  and  $\sigma_2$  are equal and not known, then  $\sigma_1 = \sigma_2 = \sigma$  is approximated by  $\sigma^2 = \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2}$ . Hence, in such a situation,

$$z = \frac{\bar{X}_1 - \bar{X}_2}{\sqrt{\left( \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2} \right) \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}}, \quad from (2)$$

$$i.e. \quad z = \frac{\bar{X}_1 - \bar{X}_2}{\sqrt{\frac{s_1^2}{n_2} + \frac{s_2^2}{n_1}}} \quad (4)$$

4. The difference in the denominators of the values of  $z$  given in (3) and (4) may be noted.

**Test 5** Test of significance of the difference between sample S.D. and population S.D.

Let 's' be the S.D. of a large sample of size  $n$  drawn from a normal population with S.D.  $\sigma$ . Then it is known that  $s$  follows a  $N\left(\sigma, \frac{\sigma}{\sqrt{2n}}\right)$  approximately.

$$\therefore \text{the test statistic } z = \frac{s - \sigma}{\sigma/\sqrt{2n}}$$

As before, the significance of the difference between  $s$  and  $\sigma$  is tested.

**Test 6** Test of significance of the difference between S.D.'s of two large samples.

Let  $s_1$  and  $s_2$  be the S.D.'s of two large samples of sizes  $n_1$  and  $n_2$  drawn from a normal population with S.D.  $\sigma$ .

$$s_1 \text{ follows a } N\left(\sigma, \frac{\sigma}{\sqrt{2n_1}}\right) \text{ and } s_2 \text{ follows a } N\left(\sigma, \frac{\sigma}{\sqrt{2n_2}}\right).$$

$$\therefore (s_1 - s_2) \text{ follows a } N\left(0, \sigma \sqrt{\frac{1}{2n_1} + \frac{1}{2n_2}}\right).$$

$$\therefore \text{the test statistic } z = \frac{s_1 - s_2}{\sigma \sqrt{\frac{1}{2n_1} + \frac{1}{2n_2}}}.$$

As usual, the significance of the difference between  $s_1$  and  $s_2$  is tested.

**Note**

If  $\sigma$  is not known, it is approximated by  $\sqrt{\left(\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2}\right)}$ , when  $n_1$  and  $n_2$  are large. In this situation,

$$\chi^2 = \frac{s_1 - s_2}{\sqrt{\left(\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2}\right)} \left( \frac{1}{2n_1} + \frac{1}{2n_2} \right)}$$

i.e.

$$\chi^2 = \frac{s_1 - s_2}{\sqrt{\frac{s_1^2}{2n_2} + \frac{s_2^2}{2n_1}}}$$

## Worked Example 10(A)

### Example 1

Experience has shown that 20 per cent of a manufactured product is of top quality. In one day's production of 400 articles, only 50 are of top quality. Show that either the production of the day chosen was not a representative sample or the hypothesis of 20 per cent was wrong. Based on the particular day's production, find also the 95 per cent confidence limits for the percentage of top quality product.

$$H_0 : P = \frac{1}{5}, \text{ i.e., 20 per cent of the products manufactured is of top quality.}$$

$$H_1 : P \neq \frac{1}{5}.$$

$p$  = proportion of top quality products in the sample

$$= \frac{50}{400} = \frac{1}{8}$$

From the alternative hypothesis  $H_1$ , we note that two-tailed test is to be used. Let us assume that LOS (level of significance)

$$= 5\%. \quad \therefore z_\alpha = 1.96$$

$$z = \frac{p - P}{\sqrt{\frac{PQ}{n}}} = \frac{\frac{1}{8} - \frac{1}{5}}{\sqrt{\frac{1}{5} \times \frac{4}{5} \times \frac{1}{400}}}, \text{ since the size of the sample} = 400.$$

$$= -\frac{3}{40} \times 50 = -3.75$$

Now,  $|z| = 3.75 > 1.96$ .

The difference between  $p$  and  $P$  is significant at 5 per cent level.

Also,  $H_0$  is rejected. Hence  $H_0$  is wrong or the production of the particular day chosen is not a representative sample.

95 per cent confidence limits for  $P$  are given by

$$\frac{|p - P|}{\sqrt{\frac{pq}{n}}} \leq 1.96$$

**Note**

We have taken  $\sqrt{\frac{pq}{n}}$  in the denominator, because  $P$  is assumed to be unknown, for which we are trying to find the confidence limits and  $P$  is nearly equal to  $p$ .

$$\text{i.e. } p - \sqrt{\frac{pq}{n}} \times 1.96 \leq P \leq p + \sqrt{\frac{pq}{n}} \times 1.96$$

$$\text{i.e. } 0.125 - \sqrt{\frac{1}{8} \times \frac{7}{8} \times \frac{1}{400}} \times 1.96 \leq P \leq 0.125 + \sqrt{\frac{1}{8} \times \frac{7}{8} \times \frac{1}{400}} \times 1.96$$

$$\text{i.e. } 0.093 \leq P \leq 0.157$$

$\therefore$  95 per cent confidence limits for the percentage of top quality product are 9.3 and 15.7.

### Example 2

The fatality rate of typhoid patients is believed to be 17.26 per cent. In a certain year 640 patients suffering from typhoid were treated in a metropolitan hospital and only 63 patients died. Can you consider the hospital efficient?

$H_0 : p = P$ , i.e. the hospital is not efficient.  $H_1 : p < P$ .

One-tailed (left-tailed) test is to be used

Let us assume that LOS = 1%.  $\therefore z_\alpha = -2.33$

$$z = \frac{p - P}{\sqrt{\frac{PQ}{n}}}, \text{ where } p = \frac{63}{640} = 0.0984 \text{ and}$$

$$P = 0.1726 \text{ and hence } Q = 0.8274.$$

$$z = \frac{0.0984 - 0.1726}{\sqrt{\frac{0.1726 \times 0.8274}{640}}} = -4.96$$

$$\therefore |z| > |z_\alpha|$$

$\therefore$  the difference between  $p$  and  $P$  is significant, i.e.,  $H_0$  is rejected and  $H_1$  is accepted.

i.e. the hospital is efficient in bringing down the fatality rate of typhoid patients.

### Example 3

A salesman in a departmental store claims that at most 60 percent of the shoppers entering the store leaves without making a purchase. A random sample of 50 shoppers showed that 35 of them left without making a purchase. Are these sample results consistent with the claim of the salesman? Use a level of significance of 0.05.

Let  $P$  and  $p$  denote the population and sample proportions of shoppers not making a purchase.

$$H_0 : p = P$$

$$H_1 : p > P, \text{ since } p = 0.7 \text{ and } P = 0.6$$

One-tailed (right-tailed) test is to be used.

$$\text{LOS} = 5\% \quad \therefore z_\alpha = 1.645$$

$$z = \frac{p - P}{\sqrt{\frac{PQ}{n}}} = \frac{0.7 - 0.6}{\sqrt{\frac{0.6 \times 0.4}{50}}} = 1.443$$

$$\therefore |z| < z_\alpha$$

$\therefore$  the difference between  $p$  and  $P$  is not significant at 5 percent level.

i.e.  $H_0$  is accepted and  $H_1$  is rejected.

i.e. the sample results are consistent with the claim of the salesman.

#### Example 4

Show that for a random sample of size 100, drawn with replacement, the standard error of sample proportion cannot exceed 0.05.

The items of the sample are drawn one after another with replacement.

$\therefore$  the proportion (probability) of success in the population, i.e.  $P$  remains a constant.

We know that the sample proportion  $p$  follows a  $N\left(P, \sqrt{\frac{PQ}{n}}\right)$

$$\text{i.e. standard error of } p = \sqrt{\frac{PQ}{n}} = \frac{1}{10} \sqrt{PQ} \quad (\because n = 100) \quad (1)$$

$$\text{Now } (\sqrt{P} - \sqrt{Q})^2 \geq 0$$

$$\text{i.e. } P + Q - 2\sqrt{PQ} \geq 0$$

$$\text{i.e. } 1 - 2\sqrt{PQ} \geq 0 \quad \text{or} \quad \sqrt{PQ} \leq \frac{1}{2} \quad (2)$$

Using (2) in (1), we get,

$$\text{S.E. of } p \leq \frac{1}{20}. \quad \text{i.e. S.E. of } p \text{ cannot exceed 0.05.}$$

#### Example 5

A cubical die is thrown 9000 times and a throw of three or four is observed 3240 times. Show that the die cannot be regarded as an unbiased one and find the extreme limits between which the probability of a throw of three or four lies.

$H_0$  : the die is unbiased, i.e.  $P = \frac{1}{3}$  (= the probability of getting 3 or 4)

$H_1$  :  $P \neq \frac{1}{3}$

Two-tailed test is to be used.

$$\text{Let LOS} = 5\% \quad \therefore z_\alpha = 1.96$$

Though we may test the significance of the difference between the sample and population proportions, we shall test the significance of the difference between the number  $X$  of successes in the sample and that in the population.

When  $n$  is large,  $X$  follows a  $N(nP, \sqrt{n}PQ)$  [Refer to Test 1].

$$\therefore z = \frac{X - nP}{\sqrt{n}PQ} = \frac{3240 - \left(9000 \times \frac{1}{3}\right)}{\sqrt{9000 \times \frac{1}{3} \times \frac{2}{3}}} = 5.37$$

$$\therefore |z| > z_\alpha$$

$\therefore$  The difference between  $X$  and  $nP$  is significant. i.e.,  $H_0$  is rejected.  
i.e., the die cannot be regarded as unbiased.

If  $X$  follows a  $N(\mu, \sigma)$ , then the reader can easily verify that  $P(\mu - 3\sigma \leq X \leq \mu + 3\sigma) = .9974$ .

The limits  $\mu \pm 3\sigma$  are considered as the extreme (confidence) limits within which  $X$  lies.

Accordingly, the extreme limits for  $P$  are given by

$$\frac{|P - p|}{\sqrt{\frac{pq}{n}}} \leq 3 \quad [\text{Refer to Example (1)}]$$

$$\text{i.e., } p - 3\sqrt{\frac{pq}{n}} \leq P \leq p + 3\sqrt{\frac{pq}{n}}$$

$$\text{i.e., } 0.36 - 3\sqrt{\frac{0.36 \times 0.64}{9000}} \leq P \leq 0.36 + 3\sqrt{\frac{0.36 \times 0.64}{9000}}$$

$$\text{i.e., } 0.345 \leq P \leq 0.375.$$

### Example 6

In a large city  $A$ , 20 per cent of a random sample of 900 school boys had a slight physical defect. In another large city  $B$ , 18.5 percent of a random sample of 1600 school boys had the same defect. Is the difference between the proportions significant?

$$p_1 = 0.2, \quad p_2 = 0.185, \quad n_1 = 900 \quad \text{and} \quad n_2 = 1600$$

$$H_0 : p_1 = p_2$$

$$H_1 : p_1 \neq p_2$$

Two-tailed test is to be used.

Let L.O.S. be 5%  $\therefore z_\alpha = 1.96$

$$z = \frac{p_1 - p_2}{\sqrt{PQ \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}} \quad (1)$$

Since  $P$ , the population proportion, is not given, we estimate it as  $\hat{P} = \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2} = \frac{180 + 296}{900 + 1600} = 0.1904$ .

Using in (1), we have

$$z = \frac{0.2 - 0.185}{\sqrt{0.1904 \times 0.8096 \times \left( \frac{1}{900} + \frac{1}{1600} \right)}} = 0.92$$

$|z| \leq z_\alpha$ . Therefore, the difference between  $p_1$  and  $p_2$  is not significant at 5 per cent level.

### Example 7

Before an increase in excise duty on tea, 800 people out of a sample of 1000 were consumers of tea. After the increase in duty, 800 people were consumers of tea in a sample of 1200 persons. Find whether there is significant decrease in the consumption of tea after the increase in duty.

Let  $p_1$  and  $p_2$  be the proportions of the consumers before and after the increase in duty respectively.

$$\text{Then } p_1 = \frac{800}{1000} = \frac{4}{5} \quad \text{and} \quad p_2 = \frac{800}{1200} = \frac{2}{3}.$$

$$H_0 : p_1 = p_2$$

$$H_1 : p_1 > p_2$$

One-tailed (right-tailed) test is to be used. Let LOS be 1%.  $\therefore z_\alpha = 2.33$ .

$$\begin{aligned} z &= \frac{p_1 - p_2}{\sqrt{PQ \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}}, \text{ where } p = \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2} \\ &= \frac{\frac{4}{5} - \frac{2}{3}}{\sqrt{0.7273 \times 0.2727 \times \left( \frac{1}{1000} + \frac{1}{1200} \right)}} \\ &= \frac{0.13 \times \sqrt{1000 \times 1200}}{\sqrt{0.7273 \times 0.2727 \times 2200}} = 6.82 \end{aligned}$$

Now,  $|z| > z_\alpha$

$\therefore$  the difference between  $p_1$  and  $p_2$  is significant at 1% level.

i.e.,  $H_0$  is rejected and  $H_1$  is accepted.

i.e., there is significant decrease in the consumption of tea after the increase in duty.

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### Example 8

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15.5 per cent of a random sample of 1600 undergraduates were smokers, whereas 20% of a random sample of 900 postgraduates were smokers in a state. Can we conclude that less number of undergraduates are smokers than the postgraduates?

$$p_1 = 0.155 \quad \text{and} \quad p_2 = 0.2; \quad n_1 = 1600 \quad \text{and} \quad n_2 = 900$$

$$H_0 : p_1 = p_2$$

$$H_1 : p_1 < p_2$$

One-tailed (left-tailed) test is to be used. Let LOS be 5%.  $\therefore z_\alpha = -1.645$ .

$$z = \frac{p_1 - p_2}{\sqrt{PQ\left(\frac{1}{n_1} + \frac{1}{n_2}\right)}}, \text{ where } P \approx \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2} = 0.1712$$

$$= \frac{0.155 - 0.2}{\sqrt{0.1712 \times 0.8288 \times \left(\frac{1}{1600} + \frac{1}{900}\right)}} = \frac{-0.045 \times 1200}{\sqrt{0.1712 \times 0.8288 \times 2500}}$$

$$= -2.87$$

Now,  $|z| > |z_\alpha|$

$\therefore$  the difference between  $p_1$  and  $p_2$  is significant.

i.e.  $H_0$  is rejected and  $H_1$  is accepted.

i.e. The habit of smoking is less among the undergraduates than among the postgraduates.

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### Example 9

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A sample of 100 students is taken from a large population. The mean height of the students in this sample is 160 cm. Can it be reasonably regarded that, in the population, the mean height is 165 cm, and the S.D. is 10 cm?

$$\bar{x} = 160, \quad n = 100, \quad \mu = 165 \quad \text{and} \quad \sigma = 10.$$

$H_0 : \bar{x} = \mu$  (i.e. the difference between  $\bar{x}$  and  $\mu$  is not significant)

$H_1 : \bar{x} \neq \mu$ .

Two-tailed test is to be used.

Let LOS be 1%  $\therefore z_\alpha = 2.58$

$$z = \frac{\bar{x} - \mu}{\sigma / \sqrt{n}} = \frac{160 - 165}{10 / \sqrt{100}} = -5$$

Now,  $|z| > z_\alpha$

$\therefore$  the difference between  $\bar{x}$  and  $\mu$  is significant at 1% level.

i.e.  $H_0$  is rejected.

i.e. it is not statistically correct to assume that  $\mu = 165$ .

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### Example 10

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The mean breaking strength of the cables supplied by a manufacturer is 1800 with a S.D. of 100. By a new technique in the manufacturing process, it is claimed that the breaking strength of the cable has increased. In order to test this claim, a sample of 50 cables is tested and it is found that the mean breaking strength is 1850. Can we support the claim at 1 per cent level of significance?

$$\bar{x} = 1850, \quad n = 50, \quad \mu = 1800 \quad \text{and} \quad \sigma = 100$$

$$H_0: \bar{x} = \mu$$

$$H_1: \bar{x} > \mu$$

One-tailed (right-tailed) test is to be used.

$$\text{LOS} = 1\% \quad \therefore z_\alpha = 2.33$$

$$z = \frac{\bar{x} - \mu}{\sigma / \sqrt{n}} = \frac{1850 - 1800}{100 / \sqrt{50}} = 3.54$$

Now,  $|z| > z_\alpha$

$\therefore$  the difference between  $\bar{x}$  and  $\mu$  is significant at 1 per cent level.

i.e.  $H_0$  is rejected and  $H_1$  is accepted.

i.e. based on the sample data, we may support the claim of increase in breaking strength.

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### Example 11

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The mean value of a random sample of 60 items was found to be 145 with a S.D. of 40. Find the 95% confidence limits for the population mean. What size of the sample is required to estimate the population mean within five of its actual value with 95% or more confidence, using the sample mean?

95% confidence limits for  $\mu$  are given by

$$\frac{|\mu - \bar{x}|}{\sigma / \sqrt{n}} \leq 1.96$$

Since the population S.D.  $\sigma$  too is not given, we can approximate it by the sample S.D.s. therefore 95% confidence limits for  $\mu$  are given by  $\frac{|\mu - \bar{x}|}{s / \sqrt{n}} \leq 1.96$

$$\text{i.e., } \bar{x} - 1.96 \frac{s}{\sqrt{n}} \leq \mu \leq \bar{x} + 1.96 \frac{s}{\sqrt{n}}.$$

$$\text{i.e., } 145 - \frac{1.96 \times 40}{\sqrt{60}} \leq \mu \leq 145 + \frac{1.96 \times 40}{\sqrt{60}}$$

$$\text{i.e., } 134.9 \leq \mu \leq 155.1$$

We have to find the value of  $n$  such that

$$P\{\bar{x} - 5 \leq \mu \leq \bar{x} + 5\} \geq .95$$

i.e.,  $P\{-5 \leq \mu - \bar{x} \leq 5\} \geq .95$

i.e.,  $P\{|\mu - \bar{x}| \leq 5\} \geq .95 \quad \text{or}$

$$P\{|\bar{x} - \mu| \leq 5\} \geq .95$$

$$\therefore P\left\{\frac{|\bar{x} - \mu|}{\sigma/\sqrt{n}} \leq \frac{5}{\sigma/\sqrt{n}}\right\} \geq .95$$

i.e.,  $P\left\{|z| \leq \frac{5\sqrt{n}}{\sigma}\right\} \geq .95$ , where  $z$  is the standard normal variate (1)

We know that  $P\{|z| \leq 1.96\} = .95$

$$\therefore \text{the least value of } n = n_L \text{ that will satisfy (1) is given by } \frac{5\sqrt{n_L}}{\sigma} = 1.96$$

i.e.,  $\sqrt{n_L} = \frac{1.96 s}{5} \quad (\because \sigma \approx s)$

i.e.,  $n_L = \left(\frac{1.96 \times 40}{5}\right)^2$

i.e.,  $n_L = 245.86$

$\therefore$  the least size of the sample = 246.

### Example 12

A normal population has a mean of 0.1 and S.D. of 2.1. Find the probability that the mean of a sample of size 900 drawn from this population will be negative.

Since  $\bar{x}$  follows a  $N\left(\mu, \frac{\sigma}{\sqrt{n}}\right)$ ,  $z = \frac{\bar{x} - \mu}{\sigma/\sqrt{n}}$  is the standard normal variate.

Now,  $P(\bar{x} < 0) = P\{\bar{x} - 0.1 < -0.1\}$

$$= P\left\{\frac{\bar{x} - 0.1}{(2.1)/\sqrt{900}} < \frac{-0.1}{(2.1)/\sqrt{900}}\right\}$$

$$= P\{z < -1.43\}$$

$$= P\{z > 1.43\},$$

by symmetry of the standard normal distribution.

$$= 0.5 - P\{0 < z < 1.43\}$$

$$= 0.5 - 0.4236 \text{ (from the normal tables)}$$

$$= 0.0764.$$

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### Example 13

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In a random sample of size 500, the mean is found to be 20. In another independent sample of size 400, the mean is 15. Could the samples have been drawn from the same population with S.D. 4?

$$\bar{x}_1 = 20, \quad n_1 = 500; \quad \bar{x}_2 = 15, \quad n_2 = 400; \quad \sigma = 4$$

$H_0 : \bar{x}_1 = \bar{x}_2$ , i.e. the samples have been drawn from the same population.

$$H_1 : \bar{x}_1 \neq \bar{x}_2.$$

Two-tailed test is to be used.

Let LOS be 1%  $\therefore z_\alpha = 2.58$

$$\begin{aligned} z &= \frac{\bar{x}_1 - \bar{x}_2}{\sigma \sqrt{\frac{1}{n_1} + \frac{1}{n_2}}} && \text{(Refer to Note 1 under Test 4)} \\ &= \frac{20 - 15}{4 \sqrt{\frac{1}{500} + \frac{1}{400}}} = 18.6 \end{aligned}$$

Now,  $|z| > z_\alpha$

$\therefore$  the difference between  $\bar{x}_1$  and  $\bar{x}_2$  is significant at 1% level.

i.e.  $H_0$  is rejected

i.e. the samples could not have been drawn from the same population.

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### Example 14

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A simple sample of heights of 6400 Englishmen has a mean of 170 cm and a S.D. of 6.4 cm, while a simple sample of heights of 1600 Americans has a mean of 172 cm and a S.D. of 6.3 cm. Do the data indicate that Americans are, on the average, taller than Englishmen?

$$n_1 = 6400, \bar{x}_1 = 170 \quad \text{and} \quad s_1 = 6.4$$

$$n_2 = 1600, \bar{x}_2 = 172 \quad \text{and} \quad s_2 = 6.3$$

$$H_0 : \mu_1 = \mu_2 \quad \text{or} \quad \bar{x}_1 = \bar{x}_2,$$

i.e., the samples have been drawn from two different populations with the same mean.

$$H_1 : \bar{x}_1 < \bar{x}_2 \quad \text{or} \quad \mu_1 < \mu_2.$$

Left-tailed test is to be used.

Let LOS be 1%.  $\therefore z_\alpha = -2.33$

$$z = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{\sigma_1^2}{n_1} + \frac{\sigma_2^2}{n_2}}} = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}}}$$

[ $\because \sigma_1 \approx s_1$  and  $\sigma_2 \approx s_2$ . Refer to Note 2 under Test 4]

$$= \frac{170 - 172}{\sqrt{\frac{(6.4)^2}{6400} + \frac{(6.3)^2}{1600}}} = -11.32$$

Now,  $|z| > |z_\alpha|$

$\therefore$  the difference between  $\bar{x}_1$  and  $\bar{x}_2$  (or  $\mu_1$  and  $\mu_2$ ) is significant at 1% level.

i.e.  $H_0$  is rejected and  $H_1$  is accepted.

i.e. Americans are, on the average, taller than Englishmen.

### Example 15

Test the significance of the difference between the means of the samples, drawn from two normal populations with the same S.D. from the following data:

Table 10.2

	Size	Mean	S.D.
Sample 1	100	61	4
Sample 2	200	63	6

$$H_0 : \bar{x}_1 = \bar{x}_2 \quad \text{or} \quad \mu_1 = \mu_2$$

$$H_1 : \bar{x}_1 \neq \bar{x}_2 \quad \text{or} \quad \mu_1 \neq \mu_2$$

Two-tailed test is to be used.

Let LOS be 5%  $\therefore z_\alpha = 1.96$

$$z = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{s_1^2}{n_2} + \frac{s_2^2}{n_1}}}$$

[Refer to Note 3 under Test 4; The populations have the same S.D.]

$$= \frac{61 - 63}{\sqrt{\frac{4^2}{200} + \frac{6^2}{100}}} = -3.02$$

Now,  $|z| > z_\alpha$

$\therefore$  the difference between  $\bar{x}_1$  and  $\bar{x}_2$  (or  $\mu_1$  and  $\mu_2$ ) is significant at 5% level.

i.e.  $H_0$  is rejected and  $H_1$  is accepted.

i.e. The two normal populations, from which the samples are drawn, may not have the same mean, though they may have the same S.D.

### Example 16

The average marks scored by 32 boys is 72 with a S.D. of 8, while that for 36 girls is 70 with a S.D. of 6. Test at 1% level of significance whether the boys perform better than girls.

$$H_0: \bar{x}_1 = \bar{x}_2 \quad (\text{or } \mu_1 = \mu_2)$$

$$H_1: \bar{x}_1 > \bar{x}_2$$

Right-tailed test is to be used.

$$\text{LOS} = 1\% \quad \therefore z_\alpha = 2.33$$

$$z = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}}}$$

(The two populations are assumed to have S.D.'s  $\sigma_1 \approx s_1$  and  $\sigma_2 \approx s_2$ )

$$= \frac{72 - 70}{\sqrt{\frac{8^2}{32} + \frac{6^2}{36}}} = 1.15$$

$$|z| < z_\alpha$$

$\therefore$  the difference between  $\bar{x}_1$  and  $\bar{x}_2$  ( $\mu_1$  and  $\mu_2$ ) is not significant at 1% level.

i.e.,  $H_0$  is accepted and  $H_1$  is rejected.

i.e., statistically, we cannot conclude that boys perform better than girls.

### Example 17

The heights of men in a city are normally distributed with a mean of 171 cm and S.D. of 7 cm, while the corresponding values for women in the same city are 165 cm and 6 cm respectively. If a man and a woman are chosen at random from this city, find the probability that the woman is taller than the man.

Let  $\bar{x}_1$  and  $\bar{x}_2$  denote the mean heights of men and women respectively.

Then  $\bar{x}_1$  follows a  $N(171, 7)$  and  $\bar{x}_2$  follows a  $N(165, 6)$

$\therefore \bar{x}_1 - \bar{x}_2$  also follows a normal distribution.

$$E(\bar{x}_1 - \bar{x}_2) = E(x_1) - E(\bar{x}_2) = 171 - 165 = 6$$

$$V(\bar{x}_1 - \bar{x}_2) = V(\bar{x}_1) + V(\bar{x}_2) = 49 + 36 = 85 \quad [\text{Refer to Test 4}]$$

$$\therefore \text{S.D. of } (\bar{x}_1 - \bar{x}_2) = \sqrt{85} = 9.22.$$

$$\therefore \bar{x}_1 - \bar{x}_2 \text{ follows a } N(6, 9.22)$$

$$\text{Now, } P(\bar{x}_2 > \bar{x}_1) = P(\bar{x}_1 - \bar{x}_2 < 0)$$

$$= P\left\{ \frac{(\bar{x}_1 - \bar{x}_2) - 6}{9.22} < \frac{-6}{9.22} \right\}$$

$$= P\{z < -0.65\}, \text{ where } z \text{ is the standard normal variate.}$$

$$= P\{z > 0.65\}, \text{ by symmetry.}$$

$$= 0.5 - P(0 < z < 0.65)$$

$$= 0.5 - 0.2422 = 0.2578.$$

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**Example 18**

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Two populations have the same mean, but the S.D. of one is twice that of the other. Show that in samples, each of size 500, drawn under simple random conditions, the difference of the means will, in all probability, not exceed  $0.3 \sigma$ , where  $\sigma$  is the smaller S.D.

Let  $\bar{x}_1$  and  $\bar{x}_2$  be the means of the samples of size 500 each. Let their S.D.'s be  $\sigma$  and  $2\sigma$  respectively.

$$\bar{x}_1 \text{ follows a } N\left(\mu, \frac{\sigma}{\sqrt{500}}\right) \text{ and}$$

$$\bar{x}_2 \text{ follows a } N\left(\mu, \frac{2\sigma}{\sqrt{500}}\right)$$

$\therefore \bar{x}_1 - \bar{x}_2$  also follows a normal distribution

$$E(\bar{x}_1 - \bar{x}_2) = E(\bar{x}_1) - E(\bar{x}_2) = \mu - \mu = 0$$

$$V(\bar{x}_1 - \bar{x}_2) = V(\bar{x}_1) + V(\bar{x}_2)$$

$$= \frac{\sigma^2}{500} + \frac{4\sigma^2}{500} = \frac{\sigma^2}{100}$$

$$\therefore \text{S.D. of } (\bar{x}_1 - \bar{x}_2) = \frac{\sigma}{10}$$

Thus,  $(\bar{x}_1 - \bar{x}_2)$  follows a  $N\left(0, \frac{\sigma}{10}\right)$ .

$$\therefore P\{|\bar{x}_1 - \bar{x}_2| \leq 0.3 \sigma\}$$

$$= P\left\{ \frac{|(\bar{x}_1 - \bar{x}_2) - 0|}{\sigma/10} \leq \frac{0.3 \sigma}{\sigma/10} \right\}$$

$$= P\{|z| \leq 3\}, \text{ where } z \text{ is the standard normal variate}$$

$$= 0.9974 \approx 1.$$

$\therefore |\bar{x}_1 - \bar{x}_2|$  will not exceed  $0.3 \sigma$  almost certainly.

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**Example 19**

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A manufacturer of electric bulbs, according to a certain process, finds the S.D. of the life of lamps to be 100 hours. He wants to change the process, if the new process results in a smaller variation in the life of lamps. In adopting a new process, a sample of 150 bulbs gave an S.D. of 95 hours. Is the manufacturer justified in changing the process?

$$\sigma = 100, \quad n = 150 \quad \text{and} \quad s = 95$$

$$H_0 : s = \sigma$$

$$H_1 : s < \sigma$$

Left-tailed test is to be used.

Let LOS be 5%.  $\therefore z_\alpha = -1.645$

$$z = \frac{s - \sigma}{\sigma / \sqrt{2n}} = \frac{95 - 100}{100 / \sqrt{300}} = -0.866$$

Now,  $|z| < |z_\alpha|$

$\therefore$  the difference between  $s$  and  $\sigma$  is not significant at 5% level.

i.e.  $H_0$  is accepted and  $H_1$  is rejected.

i.e. The manufacturer is not justified in changing the process.

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### Example 20

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The S.D. of a random sample of 1000 is found to be 2.6 and the S.D. of another random sample of 500 is 2.7. Assuming the samples to be independent, find whether the two samples could have come from populations with the same S.D.

$$n_1 = 1000, \quad s_1 = 2.6; \quad n_2 = 500, \quad s_2 = 2.7$$

$$H_0 : s_1 = s_2 \quad (\text{or } \sigma_1 = \sigma_2)$$

$$H_1 : s_1 \neq s_2$$

Two-tailed test is to be used.

Let LOS be 5%.  $\therefore z_\alpha = 1.96$

$$z = \frac{s_1 - s_2}{\sqrt{\frac{s_1^2}{2n_1} + \frac{s_2^2}{2n_2}}}, \text{ since } \sigma \text{ is not known.}$$

$$= \frac{2.6 - 2.7}{\sqrt{\frac{(2.6)^2}{1000} + \frac{(2.7)^2}{2000}}} = -0.98$$

Now,  $|z| < z_\alpha$

$\therefore$  the difference between  $s_1$  and  $s_2$  (and hence between  $\sigma_1$  and  $\sigma_2$ ) is not significant at 5% level,

i.e.,  $H_0$  is accepted.

i.e., the two samples could have come from populations with the same S.D.

### Exercise 10(A)

#### Part-A (Short-answer Questions)

1. What is the difference between population and sample?
2. Distinguish between parameter and statistic.

3. What do you mean by sampling distribution?
4. What is meant by standard error?
5. What do you mean by estimation?
6. What is meant by hypothesis testing?
7. Define null hypothesis and alternative hypothesis.
8. What is meant by test of significance?
9. What do you mean by critical region and acceptance region?
10. Define level of significance.
11. Give the general form of a test statistic.
12. Define type I and type II errors.
13. Define producer's risk and consumer's risk.
14. What is the relation between type I error and level of significance?
15. Define one-tailed and two-tailed tests.
16. Define critical value of a test statistic.
17. What is the relation between the critical value and level of significance?
18. What is the relation between the critical values for a single tailed test and a two-tailed test?
19. Write down the 1% and 5% critical values for right-tailed and two-tailed tests.
20. What do you mean by interval estimation and confidence limits?
21. Write down the general form of 95% confidence limits of a population parameter in terms of the corresponding sample statistic.
22. What is the standard error of the sample proportion, when the population proportion is (i) known, and (ii) not known?
23. What is the standard error of the difference between two sample proportions when the population proportion is (i) known, and (ii) not known?
24. What do you mean by unbiased estimate? Give an example.
25. Write down the form of the 98% confidence interval for the population mean in terms of (i) population S.D., and (ii) Sample S.D.
26. What is the standard error of the difference between the means of two large samples drawn from different populations with (i) known S.D.'s, and (ii) unknown S.D.'s?
27. What is the standard error of the difference between the means of two large samples drawn from the same population with (i) known S.D., and (ii) unknown S.D.?
28. What is the standard error of the difference between the S.D.'s of two large samples drawn from the same population with (i) known S.D., and (ii) unknown S.D.?

**Part-B**

29. Out of 200 individuals, 40 per cent show a certain trait and the number expected on a certain theory is 50 per cent. Find whether the number observed differs significantly from expectation.

30. A coin is thrown 400 times and is found to result in ‘Head’ 245 times. Test whether the coin is a fair one.
31. A manufacturer of light bulbs claims that on the average 2 per cent of the bulbs manufactured by his firm are defective. A random sample of 400 bulbs contained 13 defective bulbs. On the basis of this sample, can you support the manufacturer’s claim at 5% level of significance?
32. 100 people were affected by cholera and out of them only 90 survived. Would you reject the hypothesis that the survival rate, if affected by cholera, is 85 per cent in favour of the hypothesis that it is more at 5 per cent level of significance?
33. A random sample of 400 mangoes was taken from a big consignment and 40 were found to be bad. Prove that the percentage of bad mangoes in the consignment will, in all probability, lie between 5.5 and 14.5.
34. A random sample of 64 articles produced by a machine contained 14 defectives. Is it reasonable to assume that only 10 per cent of the articles produced by the machine are defective? If not, find the 99 per cent confidence limits for the percentage of defective articles produced by the machine.
35. Certain crosses of the pea gave 5321 yellow and 1804 green seeds. The expectation is 25 per cent of green seeds based on a certain theory. Is this divergence significant or due to sampling fluctuations?
36. During a countrywide investigation, the incidence of T.B. was found to be 1 per cent. In a college with 400 students, 5 are reported to be affected whereas in another with 1200 students, 10 are found to be affected. Does this indicate any significant difference?
37. A random sample of 600 men chosen from a certain city contained 400 smokers. In another sample of 900 men chosen from another city, there were 450 smokers. Do the data indicate that (i) the cities are significantly different with respect to smoking habit among men? (ii) the first city contains more smokers than the second?
38. A sample of 300 spare parts produced by a machine contained 48 defectives. Another sample of 100 spare parts produced by another machine contained 24 defectives. Can you conclude that the first machine is better than the second?
39. In two large populations, there are 30 per cent and 25 per cent respectively of fair haired people. Is this difference likely to be hidden in samples of sizes 1200 and 900 respectively drawn from the two populations?

**Hint:**  $H_0 : P_1 - P_2 = 0$ , or  $P_1 = P_2$  and

$$H_1 : P_1 \neq P_2 \text{ and } z = \frac{P_1 - P_2}{\sqrt{\frac{P_1 Q_1}{n_1} + \frac{P_2 Q_2}{n_2}}}$$

40. A machine produces 16 defective bolts in a batch of 500 bolts. After the machine is overhauled, it produces three defective bolts in a batch of 100 bolts. Has the machine improved?
41. There were 956 births in a year in town A, of which 52.5 per cent were males, while in towns A and B combined together this proportion in a total of 1406 births was 0.496. Is there any significant difference in the proportion of male births in the two towns?
42. A cigarette-manufacturing company claims that its brand A cigarettes outsells its brand B by 8 per cent. It is found that 42 out of a sample of 200 smokers prefer brand A and 18 out of another sample of 100 smokers prefer brand B. Test at 5 per cent L.O.S. whether the 8 per cent difference is a valid claim.

**Hint:**  $H_0 : P_1 - P_2 = .08$ ;  $H_1 : P_1 - P_2 \neq .08$  and

$$Z = \frac{(p_1 - p_2) - (P_1 - P_2)}{\sqrt{PQ\left(\frac{1}{n_1} + \frac{1}{n_2}\right)}}, \text{ where } P = \frac{n_1 p_1 + n_2 p_2}{n_1 + n_2}$$

43. A sample of 900 items is found to have a mean of 3.47 cm. Can it be reasonably regarded as a simple sample from a population with mean 3.23 cm and S.D. 2.31 cm?
44. A sample of 400 observations has mean 95 and S.D. 12. Could it be a random sample from a population with mean 98? What should be the maximum value of the population mean so that the sample can be regarded as one drawn from it almost certainly?
45. A manufacturer claims that, the mean breaking strength of safety belts for air passengers produced in his factory is 1275 kgs. A sample of 100 belts was tested and the mean breaking strength and S.D. were found to be 1258 kg and 90 kg respectively. Test the manufacturer's claim at 5 per cent level of significance.
46. An I.Q. test was given to a large group of boys in the age group of 18 to 20 years, who scored an average of 62.5 marks. The same test was given to a fresh group of 100 boys of the same age group. They scored an average of 64.5 marks with a S.D. 12.5 marks. Can we conclude that the fresh group of boys have better I.Q.?
47. The guaranteed average life of a certain brand of electric bulb is 1000 hours with a S.D. of 125 hours. It is decided to sample the output so as to ensure that 90 per cent of the bulbs do not fall short of the guaranteed average by more than 2.5 per cent. What should be the minimum sample size?

48. A random sample of 100 students gave a mean weight of 58 kg with a S.D. of 4 kg. Find the 95 per cent and 99 per cent confidence limits of the mean of the population.
49. The means of two simple samples of 1000 and 2000 items are 170 cm and 169 cm. Can the samples be regarded as drawn from the same population with S.D. 10, at 5 per cent level of significance?
50. The mean and S.D. of sample sizes of 400 are 250 and 40 respectively. Those of another sample of size 400 are 220 and 55. Test at 1% level of significance whether the means of the two populations from which the samples have been drawn are equal.
51. Intelligence tests were given to two groups of boys and girls of the same age group chosen from the same college and the following results were got:

**Table 10.3**

	<i>Size</i>	<i>Mean</i>	<i>S.D.</i>
Boys	100	73	10
Girls	60	75	8

Examine if the difference between the means is significant.

52. A sample of 100 bulbs of brand A gave a mean lifetime of 1200 hours with a S.D. of 70 hours, while another sample of 120 bulbs of brand B gave a mean lifetime of 1150 hours with a S.D. of 85 hours. Can we conclude that brand A bulbs are superior to brand B bulbs?
53. In a college, 60 junior students are found to have a mean height of 171.5 cm and 50 senior students are found to have a mean height of 173.8 cm. Can we conclude, based on this data, that the juniors are shorter than seniors at (i) 5% level of significance, and (ii) 1% level of significance, assuming that the S.D. of students of that college is 6.2 cm?
54. Two samples drawn from two different populations gave the following results:

**Table 10.4**

	<i>Size</i>	<i>Mean</i>	<i>S.D.</i>
Sample I	400	124	14
Sample II	250	120	12

Find the 95% confidence limits for the difference of the population means.

$$\text{Hint: } (\bar{x}_1 - \bar{x}_2) - 1.96 \sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}} \leq \mu_1 - \mu_2 \leq (\bar{x}_1 - \bar{x}_2) + 1.96 \sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}}$$

55. Two samples drawn from two different populations gave the following results:

**Table 10.5**

	Size	Mean	S.D.
Sample I	100	582	24
Sample II	100	540	28

Test the hypothesis, at 5% level of significance, that the difference of the means of the populations is 35.

$$\text{Hint: } z = \frac{(\bar{x}_1 - \bar{x}_2) - 35}{\sqrt{\frac{s_1^2}{n_1} + \frac{s_2^2}{n_2}}}$$

56. Two populations have their means equal, but the S.D. of one is twice the other. Show that, in the samples of size 2000 drawn one from each, the difference of the means will in all probability, not exceed 0.15  $\sigma$ , where  $\sigma$  is the smaller S.D.
57. In a certain random sample of 72 items, the S.D. is found to be 8. Is it reasonable to suppose that it has been drawn from a population with S.D. 7?
58. In a random sample of 200 items, drawn from a population with S.D. 0.8, the sample S.D. is 0.7. Can we conclude that the sample S.D. is less than the population S.D. at 1% level of significance?
59. The S.D. of a random sample of 900 members is 4.6 and that of another independent sample of 1600 members is 4.8. Examine if the two samples could have been drawn from a population with S.D. 4.0?
60. Examine whether the two samples for which the data are given in Table 10.6 could have been drawn from populations with the same S.D.:

**Table 10.6**

	Size	S.D.
Sample I	100	5
Sample II	200	7

## Tests of Significance for Small Samples

The tests of significance discussed in the previous section hold good only for large samples, i.e., only when the size of the sample  $n \geq 30$ . When the sample is small, i.e.,  $n < 30$ , the sampling distributions of many statistics are not normal, even though the parent populations may be normal. Moreover, the assumption of near equality of population parameters and the corresponding sample statistics will not be justified for small samples. Consequently, we have to develop entirely different tests of significance that are applicable to small samples.

### Student's $t$ -Distribution

A random variable  $T$  is said to follow student's  $t$ -distribution or simply  $t$ -distribution, if its probability density function is given by

$$f(t) = \frac{1}{\sqrt{v} \beta \left(\frac{v}{2}, \frac{1}{2}\right)} \cdot \left(1 + \frac{t^2}{v}\right)^{-(v+1)/2}, \quad -\infty < t < \infty.$$

$v$  is called the number of degrees of freedom of the  $t$ -distribution.

**Note**

$t$ -distribution was defined by the mathematician W.S.D. Gosset whose pen name is Student.

### Properties of $t$ -Distribution

1. The probability curve of the  $t$ -distribution is similar to the standard normal curve and is symmetric about  $t = 0$ , bell-shaped and asymptotic to the  $t$ -axis as shown in Fig. 10.2.
2. For sufficiently large value of  $n$ , the  $t$ -distribution tends to the standard normal distribution.
3. The mean of the  $t$ -distribution is zero.
4. The variance of the  $t$ -distribution is  $\frac{v}{v-2}$ , if  $n > 2$  and is greater than 1, but it tends to 1 as  $v \rightarrow \infty$ .

### Uses of $t$ -Distribution

The  $t$ -distribution is used to test the significance of the difference between

1. The mean of a small sample and the mean of the population.
2. The means of two small samples and
3. The coefficient of correlation in the small sample and that in the population, assumed zero.

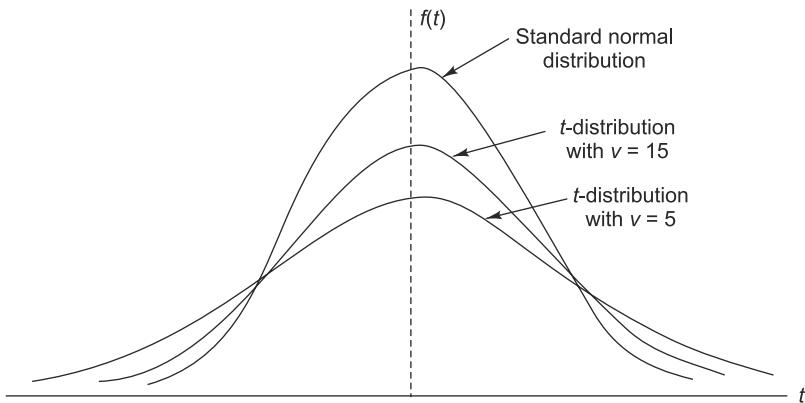


Fig. 10.2

### Note on Degree of Freedom

The number of degrees of freedom, usually denoted by the Greek alphabet  $v$ , can be interpreted as the number of useful bits of information generated by a sample of given size for estimating a population parameter. Suppose we wish to find the mean of a sample with observations  $x_1, x_2, \dots, x_n$ . We have to use all the ' $n$ ' values taken by the variable with full freedom (i.e., without putting any constraint or restriction on them) for computing  $\bar{x}$ . Hence,  $\bar{x}$  is said to have  $n$  degrees of freedom.

Suppose we wish to further compute the S.D. 's' of this sample using the formula  $s^2 = \frac{1}{n} \sum (x_i - \bar{x})^2$ . Though we use the  $n$  values  $x_1 - \bar{x}, x_2 - \bar{x}, \dots, x_n - \bar{x}$  for this computation, they do not have ' $n$ ' degrees of freedom, as they depend on  $\bar{x}$  which has been already calculated and fixed. Since there is one restriction regarding the value of  $\bar{x}$ , 's' is said to have  $(n - 1)$  degrees of freedom.

If we compute another statistic of the sample based on  $\bar{x}$  and  $s$ , then that statistic will be assumed to have  $(n - 2)$  degrees of freedom, and so on.

Thus, the number of independent variates used to compute the test statistic is known as the number of degrees of freedom of that statistic. In general, the number of degrees of freedom is given by  $v = n - k$ , where  $n$  is the number of observations in the sample and  $k$  is the number of constraints imposed on them or  $k$  is the number of values that have been found out and specified by prior calculations.

### Critical Values of $t$ and the $t$ -Table

The critical value of  $t$  at level of significance  $\alpha$  and degrees of freedom  $v$  is given by  $P \{|t| > t_v(\alpha)\} = \alpha$  for two-tailed test, as in the case of normal distribution and large samples and by  $P \{t > t_v(\alpha)\} = \alpha$  for the right-tailed test also, as in the case

of normal distribution. The critical value of  $t$  for a single (right or left) tailed test at LOS ' $\alpha$ ' corresponding to  $v$  degrees of freedom is the same as that for a two-tailed test at LOS ' $2\alpha$ ' corresponding to the same degrees of freedom.

Critical values  $t_v(\alpha)$  of the  $t$ -distribution for two-tailed tests corresponding to a few important levels of significance and a range of values of  $v$  have been published by Prof. R.A. Fisher in the form of a table, called the  $t$ -table, which is given in the Appendix.

**Test 1** *Test of significance of the difference between sample mean and population mean.*

If  $\bar{x}$  is the mean of a sample of size  $n$ , drawn from a population  $N(\mu, \sigma)$ , we have seen that  $z = \frac{\bar{x} - \mu}{\sigma / \sqrt{n}}$  follows a  $N(0, 1)$ .

If  $\sigma$ , the S.D. of the population is not known, we have to estimate it using the sample S.D.'s'. From the theory of estimation, it is known that  $s \sqrt{\frac{n}{n-1}}$  is an unbiased estimate of  $\sigma$  with  $(n-1)$  degrees of freedom. When  $n$  is large,  $\frac{n}{n-1} \approx 1$  and, hence,  $s$  was taken as a satisfactory estimate of  $\sigma$  and hence  $z = \frac{\bar{x} - \mu}{s / \sqrt{n}}$  was assumed to follow a  $N(0, 1)$ . But when  $n$  is small, we cannot use  $s$  as an estimate of  $\sigma$ , since

$$\frac{\bar{x} - \mu}{\sigma / \sqrt{n}} = \frac{\bar{x} - \mu}{s \sqrt{\frac{n}{n-1}} \cdot \frac{1}{\sqrt{n}}} = \frac{\bar{x} - \mu}{s / \sqrt{n-1}}$$

Now,  $\frac{\bar{x} - \mu}{s / \sqrt{n-1}}$  does not follow a normal distribution, but follows a  $t$ -distribution with number of degrees of freedom  $v = n-1$ . Hence,  $\frac{\bar{x} - \mu}{s / \sqrt{n-1}}$  is denoted by  $t$  and is taken as the test-statistic.

Sometimes  $t = \frac{\bar{x} - \mu}{s / \sqrt{n-1}}$  is also taken as  $t = \frac{\bar{x} - \mu}{S / \sqrt{n}}$ ,

where  $S^2 = \frac{1}{n-1} \sum_{r=1}^n (x_r - \bar{x})^2$  and is called students 't'.

We shall use only  $t = \frac{\bar{x} - \mu}{s / \sqrt{n-1}}$ , where  $s$  is the sample S.D.

We get the value of  $t_v(\alpha)$  for the L.O.S.  $\alpha$  and  $v = n - 1$  from the  $t$ -table.

If the calculated value of  $t$  satisfies  $|t| < t_v(\alpha)$ , the null hypothesis  $H_0$  is accepted at L.O.S. ' $\alpha$ ' otherwise,  $H_0$  is rejected at L.O.S. ' $\alpha$ '.

**Note**

95% confidence interval of  $m$  is given by

$$\left| \frac{\bar{x} - \mu}{s / \sqrt{n-1}} \right| \leq t_{0.05}, \quad \text{since } P \left\{ \left| \frac{\bar{x} - \mu}{s / \sqrt{n-1}} \right| \leq t_{0.05} \right\} = 0.95$$

i.e. by  $\bar{x} - t_{0.05} \frac{s}{\sqrt{n-1}} \leq m \leq \bar{x} + t_{0.05} \times \frac{s}{\sqrt{n-1}}$ , where  $t_{0.05}$  is the 5 per cent

critical value of  $t$  for  $n (= n - 1)$  degrees of freedom for a two-tailed test.

**Test 2** Test of significance of the difference between means of two small samples drawn from the same normal population.

In Test (4) for large samples, the test statistic used to test the significance of the difference between the means of two samples from the same normal population was taken as

$$z = \frac{\bar{x}_1 - \bar{x}_2}{\sigma \sqrt{\frac{1}{n_1} + \frac{1}{n_2}}}, \text{ which follows a } N(0, 1) \quad (1)$$

If  $\sigma$  is not known, we may assume that  $\sigma = \sqrt{\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2}}$ , when  $n_1$  and  $n_2$

are large, where  $s_1$  and  $s_2$  are the sample S.D.'s. This assumption no longer holds good when  $n_1$  and  $n_2$  are small.

In fact, it is known from the theory of estimation, that an estimate of  $\sigma$  is  $\sqrt{\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2}}$  with  $(n_1 + n_2 - 2)$  degrees of freedom, when  $n_1$  and  $n_2$  are small.

Using this value of  $\sigma$  in (1), the test statistic becomes

$$\frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left( \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2} \right) \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}},$$

which does not follow a  $N(0, 1)$ , but follows a  $t$ -distribution with  $v = (n_1 + n_2 - 2)$  degrees of freedom. Hence, the  $t$ -test is applied in this case.

**Note**

1. If  $n_1 = n_2 = n$  and if the samples are independent i.e., the observations in the two samples are not at all related, then the test statistic is given by

$$t = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{s_1^2 + s_2^2}{n-1}}} \quad \text{with} \quad v = 2n - 2 \quad (2)$$

2. If  $n_1 = n_2 = n$  and if the pairs of values of  $x_1$  and  $x_2$  are associated in some way (or correlated), the formula (2) for  $t$  in Note (1) should not be used. In this case, we shall assume that  $H_0 : \bar{d} (= \bar{x} - \bar{y}) = 0$  and test the significance of the difference between  $\bar{d}$  and 0, using the test statistic  $t = \frac{\bar{d}}{s / \sqrt{n-1}}$  with  $v = n - 1$ , where  $d_i = x_i - y_i$  ( $i = 1, 2, \dots, n$ ),

$$\bar{d} = \bar{x} - \bar{y}; \text{ and } s = \text{S.D. of } d\text{'s} = \frac{1}{n} \sum_{i=1}^n (d_i - \bar{d})^2.$$

## Snedecor's F-Distribution

A random variable  $F$  is said to follow Snedecor's  $F$ -distribution or simply  $F$ -distribution, if its probability density function is given by

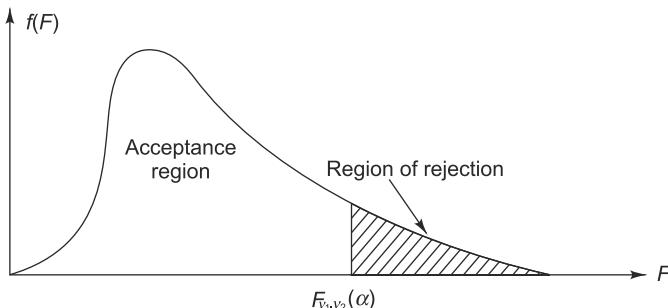
$$f(F) = \frac{(v_1/v_2)^{v_1/2}}{\beta\left(\frac{v_1}{2}, \frac{v_2}{2}\right)} \cdot \frac{F^{\frac{v_1}{2}-1}}{\left(1 + \frac{v_1 F}{v_2}\right)^{(v_1+v_2)/2}}, \quad F > 0.$$

**Note**

(The mathematical variable corresponding to the random variable  $F$  is also taken as  $F$ .)  $v_1$  and  $v_2$  used in  $f(F)$  are the degrees of freedom associated with the  $F$ -distribution.

## Properties of the F-Distribution

- The probability curve of the  $F$ -distribution is roughly sketched in Fig. 10.3.



**Fig. 10.3**

2. The square of the  $t$ -variate with  $n$  degrees of freedom follows a  $F$ -distribution with 1 and  $n$  degrees of freedom.
3. The mean of the  $F$ -distribution is  $\frac{v_2}{v_2 - 2}$  ( $v_2 > 2$ ).
4. The variance of the  $F$ -distribution is  $\frac{2v_2^2(v_1 + v_2 - 2)}{v_1(v_2 - 2)^2(v_2 - 4)}$  ( $v_2 > 4$ ).

### **Use of $F$ -Distribution**

$F$ -distribution is used to test the equality of the variance of the populations from which two small samples have been drawn.

*F-test of significance of the difference between population variances and F-table.*

To test the significance of the difference between population variances, we shall first find their estimates,  $\hat{\sigma}_1^2$  and  $\hat{\sigma}_2^2$  based on the sample variances  $s_1^2$  and  $s_2^2$  and then test their equality. It is known that  $\hat{\sigma}_1^2 = \frac{n_1 s_1^2}{n_1 - 1}$  with the number of degree of freedom  $v_1 = n_1 - 1$  and  $\hat{\sigma}_2^2 = \frac{n_2 s_2^2}{n_2 - 1}$  with the number of degrees of freedom  $v_2 = n_2 - 1$ .

It is also known that  $F = \frac{\hat{\sigma}_1^2}{\hat{\sigma}_2^2}$  follows a  $F$ -distribution with  $v_1$  and  $v_2$  degrees of freedom. If  $\hat{\sigma}_1^2 = \hat{\sigma}_2^2$ , then  $F = 1$ . Hence, our aim is to find how far any observed value of  $F$  can differ from unity due to fluctuations of sampling.

Snedecor has prepared tables that give, for different values of  $v_1$  and  $v_2$ , the 5 per cent and 1 per cent critical values of  $F$ . An extract from these tables is given in the Appendix. If  $F$  denotes the observed (calculated) value and  $F_{v_1, v_2}(\alpha)$  denotes the critical (tabulated) value of  $F$  at LOS  $\alpha$ , then  $P\{F > F_{v_1, v_2}(\alpha)\} = \alpha$ .

#### **Note**

*F-test is not a two-tailed test and is always a right-tailed test, since  $F$  cannot be negative. Thus, if  $F > F_{v_1, v_2}(\alpha)$ , then the difference between  $F$  and 1, i.e., the difference between  $\hat{\sigma}_1^2$  and  $\hat{\sigma}_2^2$  is significant at LOS ' $\alpha$ '. In other words, the samples may not be regarded as drawn from the same population or from populations with the same variance. If  $F < F_{v_1, v_2}(\alpha)$ , the difference is not significant at LOS  $\alpha$ .*

1. We should always make  $F > 1$ . This is done by taking the larger of the two estimates of  $\sigma^2$  as  $\hat{\sigma}_1^2$  and by assuming that the corresponding degree of freedom as  $v_1$ .
2. To test if two small samples have been drawn from the same normal population, it is not enough to test if their means differ significantly or not, because in this test we assumed that the two samples came from the same population or from populations with equal variance. So, before applying the  $t$ -test for the significance of the difference of two sample means, we should satisfy ourselves about the equality of the population variances by  $F$ -test.

### Worked Example 10(B)

#### Example 1

Tests made on the breaking strength of 10 pieces of a metal wire gave the results: 578, 572, 570, 568, 572, 570, 570, 572, 596 and 584 kg. Test if the mean breaking strength of the wire can be assumed as 577 kg.

Let us first compute sample mean  $\bar{x}$  and sample S.D.'s' and then test if  $\bar{x}$  differs significantly from the population mean  $\mu = 577$ .

$$\text{We take the assumed mean } A = \frac{568 + 596}{2} = 582$$

$$d_i = x_i - A$$

$$\therefore x_i = d_i + A$$

$$\begin{aligned}\therefore \bar{x} &= \frac{1}{n} \sum x_i = \frac{1}{n} \sum d_i + A \\ &= \frac{1}{10} \times (-68) + 582 = 575.2 \text{ (see Table 10.7 given below)}\end{aligned}$$

**Table 10.7**

$x_i$	$d_i = x_i - A$	$d_i^2$
578	-4	16
572	-10	100
570	-12	144
568	-14	196
572	-10	100
570	-12	144
570	-12	144
572	-10	100
596	14	196
584	2	4
Total	-68	1144

$$\begin{aligned}s^2 &= \frac{1}{n} \sum d_i^2 - \left( \frac{1}{n} \sum d_i \right)^2 \\ &= \frac{1}{10} \times 1144 - \left( \frac{1}{10} \times -68 \right)^2 = 68.16\end{aligned}$$

$$\therefore s = 8.26$$

Now,  $t = \frac{\bar{x} - \mu}{s/\sqrt{n-1}} = \frac{575.2 - 577}{8.26/\sqrt{9}} = -0.65$

and  $v = n - 1 = 9$ .

$$H_0: \bar{x} = \mu \text{ and } H_1: \bar{x} \neq \mu.$$

Let LOS be 5%. Two-tailed test is to be used.

From the  $t$ -table, for  $v = 9$ ,  $t_{5\%} = 2.26$ . Since  $|t| < t_{5\%}$ , the difference between  $\bar{x}$  and  $\mu$  is not significant or  $H_0$  is accepted.  $\therefore$  the mean breaking strength of the wire can be assumed as 577 kg at 5% LOS

### Example 2

A machinist is expected to make engine parts with axle diameter of 1.75 cm. A random sample of 10 parts shows a mean diameter 1.85 cm. with a S.D. of 0.1 cm. On the basis of this sample, would you say that the work of the machinist is inferior?

$$\bar{x} = 1.85, s = 0.1, n = 10 \text{ and } \mu = 1.75.$$

$$H_0: \bar{x} = \mu; H_1: \bar{x} \neq \mu$$

Two-tailed test is to be used. Let L.O.S. be 5%

$$t = \frac{\bar{x} - \mu}{s/\sqrt{n-1}} = \frac{0.10}{0.1/\sqrt{9}} = 3 \text{ and } v = n - 1 = 9.$$

From the  $t$ -table, for  $v = 9$ ,  $t_{0.05} = 2.26$  and  $t_{0.01} = 3.25$ .

$$\therefore |t| > t_{0.05} \text{ and } |t| < t_{0.01}$$

$\therefore H_0$  is rejected and  $H_1$  is accepted at 5% level, but  $H_0$  is accepted and  $H_1$  is rejected at 1% level. That is, at 5% LOS, the work of the machinist can be assumed to be inferior, but at 1% LOS, the work cannot be assumed to be inferior.

### Example 3

A certain injection administered to each of 12 patients resulted in the following increases of blood pressure:

$$5, 2, 8, -1, 3, 0, 6, -2, 1, 5, 0, 4.$$

Can it be concluded that the injection will be, in general, accompanied by an increase in B.P.?

$$\text{The mean of the sample is given by } \bar{x} = \frac{1}{n} \sum x = \frac{31}{12} = 2.58$$

The S.D. 's' of the sample is given by

$$s^2 = \frac{1}{n} \sum x^2 - \left( \frac{1}{n} \sum x \right)^2 = \frac{1}{12} \times 185 - (2.58)^2 = 8.76$$

$$\therefore s = 2.96$$

$$H_0: \bar{x} = \mu,$$

where  $\mu = 0$ , i.e. the injection will not result in increase in B.P.

$$H_1: \bar{x} > \mu$$

Right-tailed test is to be used. Let L.O.S. be 5%. Now,  $t_{5\%}$  for one-tailed test for ( $v = 11$ ) =  $t_{10\%}$  for two-tailed test for ( $v = 11$ ) = 1.80 (from  $t$ -table)

$$\text{Now } t = \frac{\bar{x} - \mu}{s / \sqrt{n-1}} = \frac{2.58 - 0}{2.96 / \sqrt{11}} = 2.89$$

We see that  $|t| > t_{10\%}$  ( $v = 11$ )

$\therefore H_0$  is rejected and  $H_1$  is accepted.

i.e. we may conclude that the injection is accompanied by an increase in B.P.

### Example 4

The mean lifetime of a sample of 25 bulbs is found as 1550 hours with an S.D. of 120 hours. The company manufacturing the bulbs claims that the average life of their bulbs is 1600 hours. Is the claim acceptable at 5% level of significance?

$$\bar{x} = 1550, s = 120, n = 25 \text{ and } \mu = 1600.$$

$$H_0: \bar{x} = \mu \text{ and } H_1: \bar{x} < \mu.$$

Left-tailed test is to be used. LOS = 5%

$$\text{Now, } t = \frac{\bar{x} - \mu}{s / \sqrt{n-1}} = \frac{-50 \sqrt{24}}{120} = -2.04 \text{ and } v = 24$$

$t_{5\%}$  for one-tailed test for ( $v = 24$ ) =  $t_{10\%}$  for two-tailed test for ( $v = 24$ ) = 1.71.

We see that  $|t| > |t_{10\%}|$

$\therefore H_0$  is rejected and  $H_1$  is accepted at 5% LOS,

i.e., the claim of the company cannot be accepted at 5% LOS

### Example 5

The heights of ten males of a given locality are found to be 175, 168, 155, 170, 152, 170, 175, 160, 160 and 165 cms. Based on this sample, find the 95% confidence limits for the height of males in that locality.

We shall first find the mean  $\bar{x}$  and S.D. 's' of the sample, by taking the assumed mean  $A = 165$  (Table 10.8).

$$d_i = x_i - A$$

$$\therefore \bar{x} = A + \bar{d}$$

$$= 165 + \frac{1}{10} \times 0 = 165.$$

$$s^2 = \frac{1}{n} \sum d_i^2 - \left( \frac{1}{n} \sum d_i \right)^2$$

$$= \frac{1}{10} \times 578 = 57.8$$

$$\therefore s = 7.6$$

From the  $t$ -table,

$$t_{5\%} (v = 9) = 2.26.$$

The 95% confidence limits for  $\mu$  are

$$\left( \bar{x} - 2.26 \frac{s}{\sqrt{n-1}}, \bar{x} + 2.26 \frac{s}{\sqrt{n-1}} \right)$$

i.e.

$$\left( 165 - \frac{2.26 \times 7.6}{\sqrt{9}}, 165 + \frac{2.26 \times 7.6}{\sqrt{9}} \right)$$

i.e. (159.3, 170.7)

i.e. the heights of males in the locality are likely to lie within 159.3 cm and 170.7 cm.

**Table 10.8**

$x_i$	$d_i = x_i - A$	$d_i^2$
175	10	100
168	3	9
155	-10	100
170	5	25
152	-13	169
170	5	25
175	10	100
160	-5	25
160	-5	25
165	0	0
Total	0	578

---

### Example 6

---

Two independent samples of sizes 8 and 7 contained the following values:

Sample I : 19, 17, 15, 21, 16, 18, 16, 14

Sample II : 15, 14, 15, 19, 15, 18, 16

Is the difference between the sample means significant?

**Table 10.9**

<i>Sample I</i>			<i>Sample II</i>		
$x_1$	$d_1 = x_1 - 18$	$d_1^2$	$x_2$	$d_2 = x_2 - 16$	$d_2^2$
19	1	1	15	-1	1
17	-1	1	14	-2	4
15	-3	9	15	-1	1
21	3	9	19	3	9
16	-2	4	15	-1	1
18	0	0	18	2	4
16	-2	4	16	0	0
14	-4	16			
Total	-8	44	Total	0	20

$$\text{For sample I, } \bar{x}_1 = 18 + \bar{d}_1 = 18 + \frac{1}{8} \sum d_1 \\ = 18 + \frac{1}{8} \times (-8) = 17.$$

$$s_1^2 = \frac{1}{n_1} \sum d_1^2 - \left( \frac{1}{n_1} \sum d_1 \right)^2 \\ = \frac{1}{8} \times 44 - \left( \frac{1}{8} \times -8 \right)^2 = 4.5 \\ \therefore s_1 = 2.12.$$

$$\text{For sample II, } \bar{x}_2 = 16 + \bar{d}_2 = 16 + \frac{1}{7} \sum d_2 = 16.$$

$$s_2^2 = \frac{1}{n_2} \sum d_2^2 - \left( \frac{1}{n_2} \sum d_2 \right)^2 \\ = \frac{1}{7} \times 20 - \left( \frac{1}{7} \times 0 \right)^2 = 2.857$$

$$\therefore s_2 = 1.69$$

$$H_0: \bar{x}_1 = \bar{x}_2 \quad \text{and} \quad H_1: \bar{x}_1 \neq \bar{x}_2$$

Two-tailed test is to be used. Let LOS be 5 %

$$t = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left( \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2} \right) \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}} = \frac{17 - 16}{\sqrt{\left( \frac{8 \times 4.5 + 7 \times 2.857}{13} \right) \left( \frac{1}{8} + \frac{1}{7} \right)}} \\ = 0.93$$

$$\text{Also, } v = n_1 + n_2 - 2 = 13.$$

From the  $t$ -table,  $t_{5\%} (v=13) = 2.16$

Since  $|t| < t_{5\%}$ ,  $H_0$  is accepted and  $H_1$  is rejected.

i.e. the two sample means do not differ significantly at 5% LOS

### Example 7

Table 10.10 gives the biological values of protein from cow's milk and buffalo's milk at a certain level. Examine if the average values of protein in the two samples significantly differ.

**Table 10.10**

Cow's milk ( $x_1$ ):	1.82	2.02	1.88	1.61	1.81	1.54
Buffalo's milk ( $x_2$ ):	2.00	1.83	1.86	2.03	2.19	1.88

$$n = 6$$

$$\bar{x}_1 = \frac{1}{6} \times 10.68 = 1.78$$

$$s_1^2 = \frac{1}{6} \times \sum x_1^2 - (\bar{x}_1)^2 = \frac{1}{6} \times 19.167 - (1.78)^2 = 0.0261$$

$$\bar{x}_2 = \frac{1}{6} \times 11.79 = 1.965$$

$$s_2^2 = \frac{1}{6} \times \sum x_2^2 - (\bar{x}_2)^2 = \frac{1}{6} \times 23.2599 - (1.965)^2 = 0.0154$$

As the two samples are independent, the test statistic is given by

$$t = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\frac{s_1^2 + s_2^2}{n-1}}}$$

with  $v = 2n - 2$  [Refer to Note (2) under Test (2)]

$$t = \frac{1.78 - 1.965}{\sqrt{\frac{0.0261 + 0.0154}{5}}} = \frac{-0.185}{\sqrt{0.0083}} = -2.03 \text{ and } v = 10.$$

$$H_0: \bar{x}_1 = \bar{x}_2 \quad \text{and} \quad H_1: \bar{x}_1 \neq \bar{x}_2.$$

Two-tailed test is to be used. Let LOS be 5%

From  $t$ -table,  $t_{5\%}(v=10) = 2.23$ .

Since  $|t| < t_{5\%}(v=10)$ ,  $H_0$  is accepted.

i.e., the difference between the mean protein values of the two varieties of milk is not significant at 5% level.

### Example 8

Samples of two types of electric bulbs were tested for length of life and the following data were obtained.

	Size	Mean	S.D.
Sample I	8	1234 hours	36 hours
Sample II	7	1036 hours	40 hours

Is the difference in the means sufficient to warrant that type I bulbs are superior to type II bulbs?

$$\bar{x}_1 = 1234, \quad s_1 = 36, \quad n_1 = 8; \quad \bar{x}_2 = 1036, \quad s_2 = 40, \quad n_2 = 7$$

$$H_0: \bar{x}_1 = \bar{x}_2; \quad H_1: \bar{x}_1 > \bar{x}_2.$$

Right-tailed test is to be used. Let LOS be 5%.

$$t = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left(\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2}\right)\left(\frac{1}{n_1} + \frac{1}{n_2}\right)}} = \frac{198}{\sqrt{\left(\frac{21568}{13}\right)\left(\frac{1}{8} + \frac{1}{7}\right)}} = \frac{198}{21.0807}$$

$$= 9.39$$

$$v = n_1 + n_2 - 2 = 13$$

$t_{5\%}$  ( $v = 13$ ) for one-tailed test =  $t_{10\%}$  ( $v = 13$ ) for two tailed test = 1.77 (from  $t$ -table)

Now,  $t > t_{10\%}$  ( $v = 13$ )

$\therefore H_0$  is rejected and  $H_1$  is accepted.

i.e. Type I bulbs may be regarded superior to type II bulbs at 5% LOS.

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### Example 9

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The mean height and the S.D. height of eight randomly chosen soliders are 166.9 cm. and 8.29 cm. respectively. The corresponding values of six randomly chosen sailors are 170.3 cm and 8.50 cm. respectively. Based on this data, can we conclude that soliders are, in general, shorter than sailors?

$$\bar{x}_1 = 166.9, \quad s_1 = 8.29, \quad n_1 = 8; \quad \bar{x}_2 = 170.3, \quad s_2 = 8.50, \quad n_2 = 6.$$

$$H_0: \bar{x}_1 = \bar{x}_2; \quad H_1: \bar{x}_1 < \bar{x}_2.$$

Left-tailed test is to be used. Let LOS be 5%.

$$t = \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left(\frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2}\right)\left(\frac{1}{n_1} + \frac{1}{n_2}\right)}} = \frac{-3.4}{\sqrt{\left(\frac{983.29}{12}\right)\left(\frac{1}{8} + \frac{1}{6}\right)}} \\ = -0.695$$

$$v = n_1 + n_2 - 2 = 12$$

$t_{5\%}$  ( $v = 12$ ) for one-tailed test =  $t_{10\%}$  ( $v = 12$ ) for two tailed test = 1.78 (from  $t$ -table)

Now,  $|t| < t_{10\%}$  ( $v = 12$ )

$\therefore H_0$  is accepted and  $H_1$  is rejected.

i.e. based on the given data, we cannot conclude that soliders are in general, shorter than sailors.

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### Example 10

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The following data relate to the marks obtained by 11 students in two tests, one held at the beginning of a year and the other at the end of the year after intensive coaching. Do the data indicate that the students have benefited by coaching?

Test 1: 19, 23, 16, 24, 17, 18, 20, 18, 21, 19, 20

Test 2: 17, 24, 20, 24, 20, 22, 20, 20, 18, 22, 19

The given data relate to the marks obtained in two tests by the same set of students. Hence, the marks in the two tests can be regarded as correlated and so the  $t$ -test for paired values should be used.

Let  $d = x_1 - x_2$ ,

where  $x_1, x_2$  denote the marks in the two tests.

Thus, the values of  $d$  are  $2, -1, -4, 0, -3, -4, 0, -2, 3, -3, 1$ .

$$\Sigma d = -11 \quad \text{and} \quad \Sigma d^2 = 69$$

$$\therefore \bar{d} = \frac{1}{n} \Sigma d = \frac{1}{11} \times -11 = -1$$

$$s^2 = s_d^2 = \frac{1}{n} \Sigma d^2 - (\bar{d})^2 = \frac{1}{11} \times 69 - (-1)^2 = 5.27$$

$$\therefore s = 2.296$$

$H_0 : \bar{d} = 0$ , i.e. the students have not benefited by coaching;  $H_1 : \bar{d} < 0$  (i.e.  $\bar{x}_1 < \bar{x}_2$ ).

One-tailed test is to be used. Let LOS be 5%

$$t = \frac{\bar{d}}{s / \sqrt{n-1}} = \frac{-1}{2.296 / \sqrt{10}} = -1.38 \quad \text{and} \quad v = 10$$

$t_{5\%} (v=10)$  for one-tailed test =  $t_{10\%} (v=10)$  for two-tailed test = 1.81 (from  $t$ -table).

Now,  $|t| < t_{10\%} (v=10)$

$\therefore H_0$  is accepted and  $H_1$  is rejected.

i.e. there is no significant difference between the two sets of marks.

i.e. the students have not benefitted by coaching.

### Example 11

A sample of size 13 gave an estimated population variance of 3.0, while another sample of size 15 gave an estimate of 2.5. Could both samples be from populations with the same variance?

$$n_1 = 13, \quad \hat{\sigma}_1^2 = 3.0 \quad \text{and} \quad v_1 = 12$$

$$n_2 = 15, \quad \hat{\sigma}_2^2 = 2.5 \quad \text{and} \quad v_2 = 14.$$

$H_0 : \hat{\sigma}_1^2 = \hat{\sigma}_2^2$ , i.e. The two samples have been drawn from populations with the same variance.

$$H_1 : \hat{\sigma}_1^2 \neq \hat{\sigma}_2^2. \text{ Let L.O.S. be } 5\%$$

$$F = \frac{\hat{\sigma}_1^2}{\hat{\sigma}_2^2} = \frac{3.0}{2.5} = 1.2$$

$$v_1 = 12 \quad \text{and} \quad v_2 = 14.$$

$$F_{5\%} (v_1 = 12, v_2 = 14) = 2.53, \text{ from the } F\text{-table.}$$

$$F < F_{5\%}. \therefore H_0 \text{ is accepted}$$

i.e. the two samples could have come from two normal populations with the same variance.

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**Example 12**

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Two samples of sizes nine and eight gave the sums of squares of deviations from their respective means equal to 160 and 91 respectively. Can they be regarded as drawn from the same normal population?

$$n_1 = 9, \quad \Sigma(x_i - \bar{x})^2 = 160, \quad \text{i.e. } n_1 s_1^2 = 160$$

$$n_2 = 8, \quad \Sigma(y_i - \bar{y})^2 = 91, \quad \text{i.e. } n_2 s_2^2 = 91$$

$$\hat{\sigma}_1^2 = \frac{n_1 s_1^2}{n_1 - 1} = \frac{1}{8} \times 160 = 20; \quad \hat{\sigma}_2^2 = \frac{n_2 s_2^2}{n_2 - 1} = \frac{1}{7} \times 91 = 13$$

Since  $\hat{\sigma}_1^2 > \hat{\sigma}_2^2$ ,  $v_1 = n_1 - 1 = 8$  and  $v_2 = n_2 - 1 = 7$

$$H_0: \hat{\sigma}_1^2 = \hat{\sigma}_2^2 \quad \text{and} \quad H_1: \hat{\sigma}_1^2 \neq \hat{\sigma}_2^2.$$

Let the LOS be 5%.

$$F = \frac{\hat{\sigma}_1^2}{\hat{\sigma}_2^2} = \frac{20}{13} = 1.54$$

$F_{5\%}(v_1 = 8, v_2 = 7) = 3.73$ , from the  $F$ -table.

Since  $F < F_{5\%}$ ,  $H_0$  is accepted.

i.e. the two samples could have come from two normal populations with the same variance.

We cannot say that the samples have come from the same population, as we are unable to test if the means of the samples differ significantly or not.

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**Example 13**

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Two independent samples of eight and seven items respectively had the following values of the variable.

Sample 1 : 9, 11, 13, 11, 15, 9, 12, 14

Sample 2 : 10, 12, 10, 14, 9, 8, 10

Do the two estimates of population variance differ significantly at 5% level of significance?

For the first sample,  $\Sigma x_1 = 94$  and  $\Sigma x_1^2 = 1138$

$$\begin{aligned} \therefore s_1^2 &= \frac{1}{n_1} \sum x_1^2 - \left( \frac{1}{n_1} \sum x_1 \right)^2 \\ &= \frac{1}{8} \times 1138 - \left( \frac{1}{8} \times 94 \right)^2 = 4.19 \end{aligned}$$

For the second sample,  $\Sigma x_2 = 73$  and  $\Sigma x_2^2 = 785$

$$\therefore s_2^2 = \frac{1}{n_2} \sum x_2^2 - \left( \frac{1}{n_2} \sum x_2 \right)^2$$

$$= \frac{1}{7} \times 785 - \left( \frac{1}{7} \times 73 \right)^2 = 3.39$$

$$\hat{\sigma}_1^2 = \frac{n_1}{n_1 - 1} s_1^2 = 4.79 \text{ and } \hat{\sigma}_2^2 = \frac{n_2}{n_2 - 1} s_2^2 = 3.96$$

Since  $\hat{\sigma}_1^2 > \hat{\sigma}_2^2$ ,  $v_1 = 7$  and  $v_2 = 6$

$$H_0: \hat{\sigma}_1^2 = \hat{\sigma}_2^2 \text{ and } H_1: \hat{\sigma}_1^2 \neq \hat{\sigma}_2^2$$

$$F = \frac{\hat{\sigma}_1^2}{\hat{\sigma}_2^2} = \frac{4.79}{3.96} = 1.21$$

$F_{5\%}(v_1 = 7, v_2 = 6) = 4.21$ , from the  $F$ -table. Since  $F < F_{5\%}$ ,  $H_0$  is accepted.

i.e.  $\hat{\sigma}_1^2$  and  $\hat{\sigma}_2^2$  do not differ significantly at 5% level of significance.

### Example 14

Two random samples gave the following data:

	Size	Mean	Variance
Sample I	8	9.6	1.2
Sample II	11	16.5	2.5

Can we conclude that the two samples have been drawn from the same normal population?

Refer to Note (2) under  $F$ -test. To conclude that the two samples have been drawn from the same population, we have to check first that the variances of the populations do not differ significantly and then check that the sample means (and hence the population means) do not differ significantly.

$$\hat{\sigma}_1^2 = \frac{8 \times 1.2}{7} = 1.37; \quad \hat{\sigma}_2^2 = \frac{11 \times 2.5}{10} = 2.75$$

$$F = \frac{\hat{\sigma}_2^2}{\hat{\sigma}_1^2} = \frac{2.75}{1.37} = 2.007 \text{ with degrees of freedom 10 and 7.}$$

From the  $F$ -table,  $F_{5\%}(10, 7) = 3.64$

If  $H_0: \hat{\sigma}_1^2 = \hat{\sigma}_2^2$  and  $H_1: \hat{\sigma}_1^2 \neq \hat{\sigma}_2^2$ ,

$H_0$  is accepted, since  $F < F_{5\%}$

i.e. the variances of the populations from which samples are drawn may be regarded as equal.

$$\begin{aligned} t &= \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left( \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2} \right) \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}} = \frac{-6.9}{\sqrt{\left( \frac{9.6 + 27.5}{17} \right) \left( \frac{1}{8} + \frac{1}{11} \right)}} \\ &= -\frac{6.9}{0.6864} = -10.05 \end{aligned}$$

and  $v = n_1 + n_2 - 2 = 17$ .

$t_{5\%} (v=17) = 2.11$ , from the  $t$ -table.

If  $H_0: \bar{x}_1 = \bar{x}_2$  and  $H_1: \bar{x}_1 \neq \bar{x}_2$ ,  $H_0$  is rejected, since  $|t| > t_{5\%}$ .

i.e. the means of two samples (and so the populations) differ significantly.

$\therefore$  the two samples could not have been drawn from the same normal population.

### Example 15

The nicotine contents in two random samples of tobacco are given below.

Sample I : 21    24    25    26    27

Sample II : 22    27    28    30    31    36.

Can you say that the two samples came from the same population?

$$\bar{x}_1 = \text{Mean of sample I} = \frac{123}{5} = 24.6$$

$$\bar{x}_2 = \text{Mean of sample II} = \frac{174}{6} = 29.0$$

$$s_1^2 = \text{Variance of sample I} = \frac{1}{5} \sum (x_i - 24.6)^2 = 4.24$$

$$s_2^2 = \text{Variance of sample II} = \frac{1}{6} \sum (x_i - 29.0)^2 = 18.0$$

$$\hat{\sigma}_1^2 = \frac{5}{4} \times 4.24 = 5.30 \text{ and } v = 4; \quad \hat{\sigma}_2^2 = \frac{6}{5} \times 18.0 = 21.60 \text{ and } v = 5$$

$$H_0: \hat{\sigma}_1^2 = \hat{\sigma}_2^2; \quad H_1: \hat{\sigma}_1^2 \neq \hat{\sigma}_2^2$$

$$F = \frac{\hat{\sigma}_2^2}{\hat{\sigma}_1^2} = \frac{21.60}{5.30} = 4.07$$

$$F_{5\%}(5, 4) = 6.26.$$

Since  $F < F_{5\%}$ ,  $H_0$  is accepted.

$\therefore$  the variances of the two populations can be regarded as equal.

$$\begin{aligned} t &= \frac{\bar{x}_1 - \bar{x}_2}{\sqrt{\left( \frac{n_1 s_1^2 + n_2 s_2^2}{n_1 + n_2 - 2} \right) \left( \frac{1}{n_1} + \frac{1}{n_2} \right)}} = \frac{-4.4}{\sqrt{\left( \frac{21.2 + 108.0}{9} \right) \left( \frac{1}{5} + \frac{1}{6} \right)}} \\ &= \frac{-4.4}{2.2943} = -1.92 \end{aligned}$$

and  $v = 9$ .

From  $t$ -table,  $F_{5\%}(v=9) = 2.26$ .

If  $H_0: \bar{x}_1 = \bar{x}_2$  and  $H_1: \bar{x}_1 \neq \bar{x}_2$ ,  $H_0$  is accepted

since  $|t| < F_{5\%}$ .

That is, the means of two samples (and hence the populations) do not differ significantly. Therefore, the two samples could have been drawn from the same normal population.

### Exercise 10(B)

#### **Part-A (Short-answer Questions)**

1. Write down the probability density of student's  $t$ -distribution.
2. State the important properties of the  $t$ -distribution.
3. Give any two uses of  $t$ -distribution.
4. What do you mean by degrees of freedom?
5. How will you get the critical value of  $t$  for a single-tailed test at level of significance  $\alpha$ ?
6. What is the test statistic used to test the significance of the difference between small sample mean and population mean?
7. Give the 95% confidence interval of the population mean in terms of the mean and S.D. of a small sample.
8. What is the test statistic used to test the significance of the difference between the means of two small samples?
9. Give an estimate of the population variance in terms of variances of two small samples. What is the associated number of degrees of freedom?
10. What is the test statistic used to test the significance of the difference between the means of two small samples of the same size? What is the associated number of degree of freedom?
11. What is the test statistic used to test the significance of the difference between the means of two small samples of the same size, when the sample items are correlated?
12. Write down the probability density function of the  $F$ -distribution.
13. State the important properties of the  $F$ -distribution.
14. What is the use of  $F$ -distribution?
15. Why is the  $F$ -distribution associated with two numbers of degrees of freedom?

#### **Part-B**

16. A random sample of ten boys had the following I.Q.'s: 70, 120, 110, 101, 88, 83, 95, 98, 107, 100. Does the data support the assumption of a population mean I.Q. of 100? Find a reasonable range in which most of the mean I.Q. values of samples of ten boys lie.
17. A random sample of 16 values from a normal population showed a mean of 103.75 cm. and the sum of the squares of deviations from this mean is equal to 843.75 square cms. Show that the assumption of a mean of 108.75 cm for the population is not reasonable. Obtain 95% and 99% fiducial limits for the same.

18. The mean weekly sales of soap bars in departmental stores is 145 bars per store. After an advertising campaign, the mean weekly sales in 17 stores for a typical week increased to 155 and showed a S.D. of 16. Was the advertising campaign successful?
19. The annual rainfall at a certain place is normally distributed with mean 30. If the rainfalls during the past eight years are 31.1, 30.7, 24.3, 28.1, 27.9, 32.2, 25.4 and 29.1, can we conclude that average rainfall during the past eight years is less than the normal rainfall?
20. A machine is expected to produce nails of 7 cm length. A random sample of 10 nails were found to measure : 7.2, 7.3, 7.1, 6.9, 6.8, 6.5, 6.9, 6.8, 7.1 and 7.2 cm. On the basis of this sample, what can be said about the reliability of the machine?
21. A random sample of eight envelopes is taken from the letterbox of a post office and their weights in grams are found to be: 12.2, 11.9, 12.5, 12.3, 11.6, 11.7, 12.2 and 12.4. Find the 95% and 99% confidence limits for the mean weight of the envelopes in the letter box.
22. The average production of 16 workers in a factory was 107 with a S.D. of 9, while 12 workers in another comparable factory had an average production of 111 with a S.D. of 10. Can we say that the production rate of workers in the latter factory is more than that in the former factory?
23. Two different types of drugs *A* and *B* were tried on certain patients for increasing weight. 5 persons were given drug *A* and 7 persons were given drug *B*. The increase in weight (in kg.) is given below
- |          |     |     |     |     |     |     |     |
|----------|-----|-----|-----|-----|-----|-----|-----|
| Drug A : | 3.6 | 5.5 | 5.9 | 4.1 | 1.4 |     |     |
| Drug B : | 4.5 | 3.6 | 5.5 | 6.8 | 2.7 | 3.6 | 5.0 |
- Do the two drugs differ significantly with regard to their effect in increasing weight?
24. Samples of 12 foremen in one division and 10 foremen in another division of a factory were selected at random and the following data were obtained (Table 10.11).

**Table 10.11**

	<i>Division I</i>	<i>Division II</i>
Sample size	12	10
Average monthly salary of foremen (Rs.)	5250	4900
S.D. of salary (Rs.)	152	165

Can you conclude that foremen in Division I get more salary than foremen in Division II?

25. Two independent groups of 10 children were tested to find how many digits they could repeat from memory after hearing them. The results are as follows.

Group A: 8    6    5    7    6    8    7    4    5    6  
 Group B: 10    6    7    8    6    9    7    6    7    7

Is the difference between mean scores of the two groups significant?

26. Table 10.12 gives the marks obtained by 12 students in two tests, one held before coaching and the other after coaching. Does the data indicate that the coaching was effective in improving the performance of students?

**Table 10.12**

Test I:	55,	60,	65,	75,	49,	25,	18,	30,	35,	54,	61,	72
Test II:	63,	70,	70,	81,	54,	29,	21,	38,	32,	50,	70,	80

27. In one sample of 8 items, the sum of the squares of deviations of the sample values from the sample mean was 84.4 and in another sample of 10 observations it was 102.6. Test whether this difference is significant at 5% level.
28. Two random samples drawn from two normal populations gave the following observations.  
 Sample I: 20, 16, 26, 27, 23, 22, 18, 24, 25, 19  
 Sample II: 17, 23, 32, 25, 22, 24, 28, 18, 31, 33, 20, 27  
 Test whether the two populations have the same variance.
29. Two random samples gave the following results. (Table 10.13)

**Table 10.13**

Sample No.	Size	Mean	Variance
I	16	440	40
II	25	460	42

Test whether the samples have been drawn from the same normal population.

30. Two random samples gave the following results. (Table 10.14)

**Table 10.14**

Sample No.	Size	Sum of the values	Sum of the squares of values
I	10	150	2340
II	12	168	2460

Test whether the samples have been drawn from the same normal population.

## Chi-Square Distribution

If  $X_1, X_2, \dots, X_n$  are normally distributed independent random variables, then it is known that  $(X_1^2 + X_2^2 + \dots + X_n^2)$  follows a probability distribution, called chi-square ( $\chi^2$ -distribution) distribution with  $n$  degrees of freedom.

The probability density function of the  $\chi^2$ -distribution is given by

$$f(\chi^2) = \frac{1}{2^{v/2} \sqrt{\left(\frac{v}{2}\right)}} \cdot (\chi^2)^{v/2-1} e^{-\chi^2/2}$$

$0 < \chi^2 < \infty$ , where  $v$  is the number of degrees of freedom.

### Properties of $\chi^2$ -Distribution

1. A rough sketch of the probability curve of the  $\chi^2$ -distribution for  $v = 3$  and  $v = 6$  is given in Fig. 10.4.

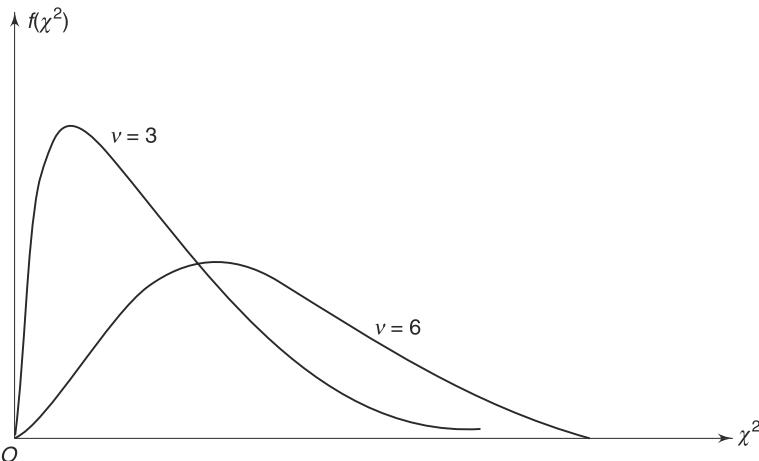


Fig. 10.4

2. As  $v$  becomes smaller and smaller, the curve is skewed more and more to the right. As  $v$  increases, the curve becomes more and more symmetrical.
3. The mean and variance of the  $\chi^2$ -distribution are  $v$  and  $2v$  respectively.
4. As  $n$  tends to  $\infty$ , the  $\chi^2$ -distribution becomes a normal distribution.

### Uses of $\chi^2$ -Distribution

1.  $\chi^2$ -distribution is used to test the goodness of fit. i.e., it is used to judge whether a given sample may be reasonably regarded as a simple sample from a certain hypothetical population.
2. It is used to test the independence of attributes. i.e., if a population is known to have two attributes (or traits), then  $\chi^2$ -distribution is used to test whether the two attributes are associated or independent, based on a sample drawn from the population.

## $\chi^2$ -Test of Goodness of Fit

On the basis of the hypothesis assumed about the population, we find the expected frequencies  $E_i (i = 1, 2, \dots, n)$ , corresponding to the observed frequencies

$O_i (i = 1, 2, \dots, n)$  such that  $\sum E_i = \sum O_i$ . It is known that  $\chi^2 = \sum_{i=1}^n \frac{(O_i - E_i)^2}{E_i}$

follows approximately a  $\chi^2$ -distribution with degrees of freedom equal to the number of independent frequencies. In order to test the goodness of fit, we have to determine how far the differences between  $O_i$  and  $E_i$  can be attributed to fluctuations of sampling and when we can assert that the differences are large enough to conclude that the sample is not a simple sample from the hypothetical population. In other words, we have to determine how large a value of  $\chi^2$  we can get so as to assume that the sample is a simple sample from the hypothetical population.

The critical value of  $\chi^2$  for  $v$  degrees of freedom at  $\alpha$  level of significance, denoted by  $\chi_v^2(\alpha)$  is given by

$$P[\chi^2 > \chi_v^2(\alpha)] = \alpha.$$

Critical values of the  $\chi^2$ -distribution corresponding to a few important levels of significance and a range of values of  $v$  are available in the form of a table called  $\chi^2$ -table, which is given in the Appendix.

If the calculated  $\chi^2 < \chi_v^2(\alpha)$ , we will accept the null hypothesis  $H_0$  which assumes that the given sample is one drawn from the hypothetical population, i.e. we will conclude that the difference between the observed and expected frequencies is not significant at  $\alpha$  % LOS If  $\chi^2 > \chi_v^2(\alpha)$ , we will reject  $H_0$  and conclude that the difference is significant.

### **Conditions for the Validity of $\chi^2$ -Test**

1. The number of observations  $N$  in the sample must be reasonably large, say  $\geq 50$ .
2. Individual frequencies must not be too small, i.e.  $O_i \geq 10$ . In case  $O_i < 10$ , it is combined with the neighbouring frequencies, so that the combined frequency is  $\geq 10$ .
3. The number of classes  $n$  must be neither too small nor too large i.e.,  $4 \leq n \leq 16$ .

### **$\chi^2$ -Test of Independence of Attributes**

If the population is known to have two major attributes  $A$  and  $B$ , then  $A$  can be divided into  $m$  categories  $A_1, A_2, \dots, A_m$  and  $B$  can be divided into  $n$  categories  $B_1, B_2, \dots, B_n$ . Accordingly the members of the population and hence those of the sample can be divided into  $mn$  classes. In this case, the sample data may be presented in the form of a matrix containing  $m$  rows and  $n$  columns and hence  $mn$  cells and showing the observed frequencies  $O_{ij}$ , in the various cells, where  $i = 1, 2, \dots, m$  and  $j = 1, 2, \dots, n$ .  $O_{ij}$  means the number of observed frequencies possessing the attributes  $A_i$  and  $B_j$ . The matrix or tabular form of the sample data, called an  $(m \times n)$  contingency table is given below:

**Table 10.15**

$A \setminus B$	$B_1$	$B_2$	-	$B_j$	-	$B_n$	Row Total
$A_1$	$O_{11}$	$O_{12}$	-	$O_{1j}$	-	$O_{1n}$	$O_{1*}$
$A_2$	$O_{21}$	$O_{22}$	-	$O_{2j}$	-	$O_{2n}$	$O_{2*}$
:	-	-	-	-	-	-	-
$A_i$	$O_{i1}$	$O_{i2}$	-	$O_{ij}$	-	$O_{in}$	$O_{i*}$
:	-	-	-	-	-	-	-
$A_m$	$O_{m1}$	$O_{m2}$	$O_{mj}$	-	-	$O_{mn}$	$O_{m*}$
Column Total	$O_{*1}$	$O_{*2}$	-	$O_{*j}$	-	$O_{*n}$	$N$

Now, based on the null hypothesis  $H_0$ , i.e., the assumption that the two attributes  $A$  and  $B$  are independent, we compute the expected frequencies  $E_{ij}$  for various cells, using the following formula  $E_{ij} = \frac{O_{i*} \cdot O_{*j}}{N}$ ,  $i = 1, 2, \dots, m$ ; and  $j = 1, 2, \dots, n$

i.e. 
$$E_{ij} = \left\{ \begin{array}{l} (\text{Total of observed frequencies in the } i^{\text{th}} \text{ row}) \times \\ (\text{total of observed frequencies in the } j^{\text{th}} \text{ column}) \\ \hline \text{Total of all cell frequencies} \end{array} \right\}$$

Then we compute  $\chi^2 = \sum_{i=1}^m \sum_{j=1}^n \left\{ \frac{(O_{ij} - E_{ij})^2}{E_{ij}} \right\}$

The number of degrees of freedom for this  $\chi^2$  computed from the  $(m \times n)$  contingency table is  $v = (m-1)(n-1)$ .

If  $\chi^2 < \chi^2_v(\alpha)$ ,  $H_0$  is accepted at  $\alpha$  % LOS, i.e., the attributes  $A$  and  $B$  are independent.

If  $\chi^2 > \chi^2_v(\alpha)$ ,  $H_0$  is rejected at  $\alpha$  % LOS, i.e.,  $A$  and  $B$  are not independent.

### Worked Example 10(C)

#### Example 1

The following table shows the distribution of digits in the numbers chosen at random from a telephone directory:

**Table 10.16**

Digit:	0	1	2	3	4	5	6	7	8	9	Total
Frequency:	1026,	1107,	997,	966,	1075,	933,	1107,	972,	964,	853	10,000

Test whether the digits may be taken to occur equally frequently in the directory.

$H_0$  : The digits occur equally frequently, i.e., they follow a uniform distribution.

Based on  $H_0$ , we compute the expected frequencies.

The total number of digits = 10,000.

If the digits occur uniformly, then each digit will occur

$$\frac{10,000}{10} = 1000 \text{ times}$$

$$O_i : 1026, 1107, \dots, 853$$

$$E_i : 1000, 1000, \dots, 1000$$

$$\begin{aligned}\chi^2 &= \sum \frac{(O_i - E_i)^2}{E_i} \\ &= \frac{1}{1000} \left\{ (26)^2 + (107)^2 + (-3)^2 + (-34)^2 + (75)^2 \right. \\ &\quad \left. + (-67)^2 + (107)^2 + (-28)^2 + (-36)^2 + (-147)^2 \right\} \\ &= 58.542\end{aligned}$$

Since  $\Sigma E_i$  was taken equal to  $\Sigma O_i$  (i.e., an information from the sample),  $v = n - 1 = 10 - 1 = 9$ . From the  $\chi^2$ -table,

$$\chi^2_{5\%} (n = 9) = 16.919$$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected, i.e., the digits do not occur uniformly in the directory.

## Example 2

Table 10.17 gives the number of air-craft accidents that occurred during the various days of a week. Test whether the accidents are uniformly distributed over the week.

Table 10.17

Day:	Mon	Tues	Wed	Thu	Fri	Sat
No. of accidents:	15	19	13	12	16	15

$H_0$  : Accidents occur uniformly over the week.

Total number of accidents = 90

Based on  $H_0$ , the expected number of accidents on any day =  $\frac{90}{6} = 15$ .

$$\begin{array}{lllllll} O_i & : & 15 & 19 & 13 & 12 & 16 & 15 \\ E_i & : & 15 & 15 & 15 & 15 & 15 & 15 \end{array}$$

$$\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i} = \frac{1}{15} (0 + 16 + 4 + 9 + 1 + 0) = 2.$$

Since  $\sum E_i = \sum O_i, v = 6 - 1 = 5$

From the  $\chi^2$ -table,  $\chi^2_{5\%} (v=5) = 11.07$ .

Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$  is accepted.

i.e. accidents may be regarded to occur uniformly over the week.

### Example 3

Table 10.18 shows defective articles produced by four machines:

**Table 10.18**

Machine:	A	B	C	D
Production time:	1 hour	1 hour	2 hours	3 hours
No. of defectives:	12	30	63	98

Do the figures indicate a significant difference in the performance of the machines?

$H_0$  : Production rates of the machines are the same.

Total number of defectives = 203.

Based on  $H_0$ , the expected numbers of defectives produced by the machines are

$$E_i : \quad \frac{1}{7} \times 203, \quad \frac{1}{7} \times 203, \quad \frac{2}{7} \times 203, \quad \frac{3}{7} \times 203$$

$$\text{i.e. } E_i : \quad 29, \quad 29, \quad 58, \quad 87 \\ O_i : \quad 12, \quad 30, \quad 63, \quad 98$$

$$\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i} = \frac{17^2}{29} + \frac{1^2}{29} + \frac{5^2}{58} + \frac{11^2}{87} = 11.82$$

Since  $\sum E_i = \sum O_i, v = 4 - 1 = 3$

From the  $\chi^2$ -table,  $\chi^2_{5\%} (v=3) = 7.815$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected.

i.e., there is significant difference in the performance of machines.

### Example 4

The following data represents the monthly sales (in Rs) of a certain retail stores in a leap year. Examine if there is any seasonality in the sales.

6100, 5600, 6350, 6050, 6250, 6200, 6300, 6250, 5800, 6000, 6150 and 6150.

$H_0$  : There is no seasonability in the sales, i.e. the daily sales are uniform throughout the year or the daily sales follow a uniform distribution.

Based on  $H_0$ , we compute the expected frequencies.

The total sales in the year = Rs. 73,200.

If the daily sales are uniform, then the sales on each day

$$= \frac{73,200}{366} = \text{Rs } 200$$

$O_i$  : 6100, 5600, 6350, 6050, 6250, 6200, 6300, 6250, 5800, 6000, 6150, 6150.

Assuming that the months are taken in the usual calendar order, namely, January, February etc. the expected monthly sales are

$E_i$  : 6200, 5800, 6200, 6000, 6200, 6000, 6200, 6000, 6200, 6000, 6200

Then  $\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i}$

$$= \frac{(-100)^2}{6200} + \frac{(-200)^2}{5800} + \dots + \frac{(-50)^2}{6200} = 38.913$$

Since  $\sum E_i$  was found as  $\sum O_i$  from the sample,  $v = n - 1 = 12 - 1 = 11$ .

From the  $\chi^2$ -table,  $\chi^2_{5\%}(v = 11) = 19.675$ .

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected, i.e., the daily sales are not uniform throughout the year.

### Example 5

Theory predicts that the proportion of beans in four groups  $A, B, C, D$  should be  $9 : 3 : 3 : 1$ . In an experiment among 1600 beans, the numbers in the four groups were 882, 313, 287 and 118. Does the experiment support the theory?

$H_0$  : The experiment supports the theory, i.e., the numbers of beans in the four groups are in the ratio  $9 : 3 : 3 : 1$

Based on  $H_0$ , the expected numbers of beans in the four groups are as follows:

$$E_i : \frac{9}{16} \times 1600, \quad \frac{3}{16} \times 1600, \quad \frac{3}{16} \times 1600, \quad \frac{1}{16} \times 1600$$

i.e.  $E_i$  : 900, 300, 300, 100

$$O_i : 882, \quad 313, \quad 287, \quad 118$$

$$\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i} = \frac{18^2}{900} + \frac{13^2}{300} + \frac{13^2}{300} + \frac{18^2}{100} = 4.73$$

Since  $\sum E_i = \sum O_i$ ,  $v = 4 - 1 = 3$

From the  $\chi^2$ -table,  $\chi^2_{5\%}(v = 3) = 7.82$

Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$  is accepted.

i.e., the experimental data support the theory.

### Example 6

A survey of 320 families with five children each revealed the following distribution:

**Table 10.19**

No. of boys:	0	1	2	3	4	5
No. of girls:	5	4	3	2	1	0
No. of families:	12	40	88	110	56	14

Is this result consistent with the hypothesis that male and female births are equally probable?

$H_0$  : Male and female births are equally probable, i.e.,  $P(\text{male birth}) = p = 1/2$  and  $P(\text{female birth}) = q = 1/2$ .

Based on  $H_0$ , the probability that a family of 5 children has  $r$  male children  
 $= 5C_r \left(\frac{1}{2}\right)^5$  (by binomial law)

$$\therefore \text{Expected number of families having } r \text{ male children} = 320 \times 5 C_r \times \frac{1}{2^5}$$

$$= 10 \times 5 C_r.$$

$$\begin{array}{llllll} \text{Thus } E_i : & 10 & 50 & 100 & 100 & 50 & 10 \\ \text{and } O_i : & 12 & 40 & 88 & 110 & 56 & 14 \end{array}$$

$$\begin{aligned} \chi^2 &= \sum \frac{(O_i - E_i)^2}{E_i} = \frac{2^2}{10} + \frac{10^2}{50} + \frac{12^2}{100} + \frac{10^2}{100} + \frac{6^2}{50} + \frac{4^2}{10} \\ &= 7.16 \end{aligned}$$

We have used the sample data to get  $\sum E_i$  only. The values of  $p$  and  $q$  have not been found by using the sample data.

$$\therefore v = n - 1 = 6 - 1 = 5 \quad \text{and} \quad \chi^2_{5\%} (v=5) = 11.07$$

Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$  is accepted.

i.e. male and female births are equally probable.

### Example 7

Twelve dice were thrown 4096 times and a throw of six was considered a success. The observed frequencies were as given below.

No. of successes: 0    1    2    3    4    5    6    7 and over

Frequency:    447    1145    1180    796    380    115    25    8

Test whether the dice were unbiased.

$H_0$  : All the dice were unbiased, i.e.,  $P(\text{getting 6}) = p = \frac{1}{6}$ .  $\therefore q = \frac{5}{6}$ .

Based on  $H_0$ , the probability of getting exactly ' $r$ ' successes =  $12 C_r p^r q^{12-r}$   
 $(r = 0, 1, 2, \dots, 12)$

$\therefore$  Expected number of times in which ' $r$ ' successes are obtained

$$= 4096 \times 12 C_r \left(\frac{1}{6}\right)^r \cdot \left(\frac{5}{6}\right)^{12-r}$$

$$= 4096 \times 12 C_r \times \frac{5^{12-r}}{6^{12}} \quad (r = 0, 1, 2, \dots, 12)$$

i.e.

$$E_0 = N(0 \text{ success}) = N(r=0) = 459.39$$

$$E_1 = N(r=1) = 1102.54$$

$$E_2 = N(r=2) = 1212.80$$

$$E_3 = N(r=3) = 808.53$$

$$E_4 = N(r=4) = 363.84$$

$$E_5 = N(r=5) = 116.43$$

$$E_6 = N(r=6) = 27.17$$

$$E_7 = N(r \geq 7) = 5.30$$

Converting  $E_i$ 's into whole numbers subject to the condition that  $\sum E_i = 4096$ , we get

$$\begin{aligned} E_i : & 459, \quad 1103, \quad 1213, \quad 809, \quad 364, \quad 116, \quad 27, \quad 5 \\ O_i : & 447, \quad 1145, \quad 1180, \quad 796, \quad 380, \quad 115, \quad 25, \quad 8, \end{aligned}$$

Since  $E$  and  $O$  corresponding to the last class, i.e., 5 and 8 are less than 10, we combine the last two classes and consider as a single class.

$$\begin{aligned} \chi^2 &= \sum \frac{(O_i - E_i)^2}{E_i} = \frac{12^2}{459} + \frac{42^2}{1103} + \frac{33^2}{1213} + \frac{13^2}{809} + \frac{16^2}{364} + \frac{1^2}{116} + \frac{1^2}{32} \\ &= 3.76 \end{aligned}$$

$v = n - 1$ , since only  $\sum E_i$  has been found using the sample data.

$$\begin{aligned} &= 7 - 1 \quad [n \text{ must be taken as the number of classes after combination} \\ &\quad \text{of end classes, if any}] \\ &= 6 \end{aligned}$$

and  $\chi^2_{5\%}(v=6) = 12.59$ , from the  $\chi^2$ -table. Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$  is accepted, i.e. the dice were unbiased.

### Example 8

Fit a binomial distribution for the following data and also test the goodness of fit.

x :	0	1	2	3	4	5	6	Total
f :	5	18	28	12	7	6	4	80

To find the binomial distribution  $N(q+p)^n$ , which fits the given data, we require  $p$ .

We know that the mean of the binomial distribution is  $np$ , from which we can find  $p$ . Now the mean of the given distribution is found out and is equated to  $np$ .

x :	0	1	2	3	4	5	6	Total
f :	5	18	28	12	7	6	4	80
fx :	0	18	56	36	28	30	24	192

$$\bar{x} = \frac{\sum f x}{\sum f} = \frac{192}{80} = 2.4$$

i.e.  $np = 2.4$  or  $6p = 2.4$ , since the maximum value taken by  $x$  is  $n$ .

$\therefore p = 0.4$  and hence  $q = 0.6$

$\therefore$  the expected frequencies are given by the successive terms in the expansion of 80  $(0.6 + 0.4)^6$ .

Thus,  $E_i$ : 3.73, 14.93, 24.88, 22.12, 11.06, 2.95, 0.33

Converting the  $E_i$ 's into whole number such that  $\sum E_i = \sum O_i = 80$ , we get

$E_i$ :	4	15	25	22	11	3	0
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Let us now proceed to test the goodness of binomial fit.

$O_i$ :	5	18	28	12	7	6	4
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The first class is combined with the second and the last two classes are combined with the last but second class in order to make the expected frequency in each class greater than or equal to 10. Thus, after regrouping, we have,

$E_i$ :	19	25	22	14
---------	----	----	----	----

$O_i$ :	23	28	12	17
---------	----	----	----	----

$$\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i} = \frac{4^2}{19} + \frac{3^2}{25} + \frac{10^2}{22} + \frac{3^2}{14} = 6.39$$

We have used the given sample to find

$$\sum E_i (= \sum O_i) \text{ and } p \text{ through its mean.}$$

Hence,  $v = n - k$

$$= 4 - 2 = 2$$

$$\chi^2_{5\%} (v = 2) = 5.99, \text{ from the } \chi^2\text{-table.}$$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$ , which assumes that the given distribution is approximately a binomial distribution, is rejected, i.e., the binomial fit for the given distribution is not satisfactory.

### Example 9

Fit a Poisson distribution for the following distribution and also test the goodness of fit.

$x$ :	0	1	2	3	4	5	Total
$f$ :	142	156	69	27	5	1	400

To find the Poisson distribution whose probability law is

$$P(X = r) = \frac{e^{-\lambda} \lambda^r}{r!}, r = 0, 1, 2, \dots,$$

we require  $\lambda$ , which is the mean of the Poisson distribution.

We find the mean of the given distribution and assume it as  $\lambda$ .

$x :$	0	1	2	3	4	5	Total
$f :$	142	156	69	27	5	1	400
$fx :$	0	156	138	81	20	5	400

$$\bar{x} = \frac{\sum fx}{\sum f} = \frac{400}{400} = 1 = \lambda$$

The expected frequencies are given by

$$\frac{N \cdot e^{-\lambda} \lambda^r}{r!} \text{ or } \frac{400 \times e^{-1}}{r!}, r = 0, 1, 2, \dots, \infty$$

Thus,

$$E_i : 147.15, 147.15, 73.58, 24.53, 6.13, 1.23$$

The values of  $E_i$  are very small for  $i = 6, 7, \dots$  and hence neglected.

Converting the values of  $E_i$ 's into whole numbers such that  $\sum E_i = 400$ , we get

$$E_i : 147, 147, 74, 25, 6, 1$$

Let us now proceed to test the goodness of Poisson fit.

$$O_i : 142, 156, 69, 27, 5, 1$$

The last three classes are combined into one, so that the expected frequency in that class may be  $\geq 10$ . Thus, after regrouping, we have

$O_i :$	142	156	69	33
$E_i :$	147	147	74	32

$$\chi^2 = \sum \frac{(O_i - E_i)^2}{E_i} = \frac{5^2}{147} + \frac{9^2}{147} + \frac{5^2}{74} + \frac{1^2}{32} \\ = 1.09$$

We have used the sample data to find  $\sum E_i$  and  $\lambda$ . Hence,

$$n = n - k = 4 - 2 = 2$$

From the  $\chi^2$ -table,  $\chi^2_{5\%} (v = 2) = 5.99$ .

Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$ , which assumes that the given distribution is nearly Poisson, is accepted,

i.e., the Poisson fit for the given distribution is satisfactory.

### Example 10

Test the normality of the following distribution by using  $\chi^2$ -test of goodness of fit

$x :$	125, 135, 145, 155, 165, 175, 185, 195, 205	Total
$f :$	1, 1, 14, 22, 25, 19, 13, 3, 2	100

Let us first fit a normal distribution to the given data and then test the goodness of fit.

To fit a normal distribution and hence find the expected frequencies, we require the density function of the normal distribution which involves the mean and S.D. Let us now compute the mean  $\bar{x}$  and S.D. 's' of the sample distribution and assume them as  $\mu$  and  $\sigma$ .

**Table 10.20**

$x$	$f$	$d = \frac{x - 165}{10}$	$fd$	$fd^2$
125	1	-4	-4	16
135	1	-3	-3	9
145	14	-2	-28	56
155	22	-1	-22	22
165	25	0	0	0
175	19	1	19	19
185	13	2	26	52
195	3	3	9	27
205	2	4	8	32
Total	100	-	5	233

$$\bar{x} = A + \frac{c}{N} \sum f d = 165 + \frac{10}{100} \times 5 = 165.5$$

$$\begin{aligned}s^2 &= c^2 \left\{ \frac{1}{N} \sum f d^2 - \left( \frac{1}{N} \sum f d \right)^2 \right\} \\ &= 10^2 (2.33 - 0.0025) = 232.75\end{aligned}$$

$$\therefore s = 15.26$$

$\therefore$  the density function of the normal distribution which fits the given distribution is  $f(x) = \frac{1}{15.26\sqrt{2\pi}} e^{-(x-165.5)^2/465.5}$ .

To find the expected frequency corresponding to a given  $x$ , we find  $y = f(x)$  and multiply  $y$  by the class-width and then by the total frequency  $N$ .

**Note**  $y = \frac{1}{\sigma\sqrt{2\pi}} e^{-(x-\mu)^2/2\sigma^2}$ . If we put  $\frac{x-\mu}{\sigma} = z$ , then  $y = \frac{1}{\sigma} \left\{ \frac{1}{\sqrt{2\pi}} e^{-z^2/2} \right\} = \frac{\phi(z)}{\sigma}$ , where  $\phi(z)$  is density function of the standard normal distribution. Values of  $\phi(z)$  are got from the normal table given in the Appendix.

**Table 10.21**

$x$	$y = \frac{x - 165.5}{15.26}$	$\phi(z)$	$\frac{c \cdot \phi(z)}{\sigma} = \frac{10\phi(z)}{15.26}$	Expected frequency $= N \phi(z)/\sigma$
125	-2.65	.0119	.0078	0.78
135	-2.00	.0540	.0354	3.54
145	-1.34	.1626	.1066	10.66
155	-0.69	.3144	.2060	20.60
165	-0.03	.3988	.2613	26.13
175	0.62	.3292	.2157	21.57
185	1.28	.1758	.1152	11.52
195	1.93	.0620	.0406	4.06
205	2.59	.0139	.0091	0.91

Converting the expected frequencies as whole numbers such that  $\sum E_i = 100$ , we get

$$E_i : 1, 3, 11, 21, 26, 22, 11, 4, 1$$

Let us now proceed to test the goodness of normal fit.

Combining the end classes so as to make the individual frequencies greater than 10,

$$\begin{aligned} E_i &: 15, 21, 26, 22, 16 \\ O_i &: 16, 22, 25, 19, 18 \end{aligned}$$

$$\begin{aligned} \chi^2 &= \sum \frac{(O_i - E_i)^2}{E_i} = \frac{1^2}{15} + \frac{1^2}{21} + \frac{1^2}{26} + \frac{3^2}{22} + \frac{2^2}{16} \\ &= 0.82 \end{aligned}$$

We have used the sample data to find  $\sum E_i$ ,  $\mu$  and  $\sigma$ . Hence  $v = n - k = 5 - 3 = 2$ .

From the  $\chi^2$ -table,  $\chi^2_{5\%} (v = 2) = 5.99$ .

Since  $\chi^2 < \chi^2_{5\%}$ ,  $H_0$ , which assumes that the given distribution is nearly normal, is accepted.

i.e. the normal fit for the given distribution is satisfactory.

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### Example 11

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The following data are collected on two characters (Table 10.22).

**Table 10.22**

	Smokers	Non-smokers
Literates :	83	57
Illiterates :	45	68

Based on this, can you say that there is no relation between smoking and literacy?

$H_0$  : Literacy and smoking habit are independent.

**Table 10.23**

	Smokers	Non-smokers	Total
Literates	83	57	140
Illiterates	45	68	113
Total	128	125	253

**Table 10.24**

$O$	$E$	$E$ (rounded)	$(O - E)^2/E$
83	$\frac{128 \times 140}{253} = 70.83$	71	$122/71 = 2.03$
57	$\frac{125 \times 140}{253} = 69.17$	69	$122/69 = 2.09$
45	$\frac{128 \times 113}{253} = 57.17$	57	$122/57 = 2.53$
68	$\frac{125 \times 113}{253} = 55.83$	56	$122/56 = 2.57$
			$\chi^2 = 9.22$

$$\begin{aligned} v &= (m-1)(n-1) \\ &= (2-1)(2-1) = 1. \end{aligned}$$

From the  $\chi^2$ -table,  $\chi^2_{5\%} (v=1) = 3.84$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected.

i.e. there is some association between literacy and smoking.

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### Example 12

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Prove that the value of  $\chi^2$  for the  $2 \times 2$  contingency table is given by

$a$	$b$
$c$	$d$

$$\chi^2 = \frac{N(ad - bc)^2}{(a+b)(c+d)(a+c)(b+d)}, \text{ where } N = a + b + c + d.$$

Hence, compute  $\chi^2$  for the  $2 \times 2$  contingency table given in Example 11.

The value of  $E$  corresponding to the cell in which  $O = a$  is given by

$$E = \frac{(a+b)(a+c)}{(a+b+c+d)}.$$

$\therefore$  the value of  $\chi^2$  corresponding to this cell is given by

$$\begin{aligned}\chi^2 &= \left\{ a - \frac{(a+b)(a+c)}{a+b+c+d} \right\}^2 \div \frac{(a+b)(a+c)}{(a+b+c+d)} \\ &= \frac{\{a(a+b+c+d) - (a+b)(a+c)\}^2}{N(a+b)(a+c)} \\ &= \frac{(ad-bc)^2}{N(a+b)(a+c)}\end{aligned}$$

Similarly, the value of  $\chi^2$  are found out for the other three cells.

$\therefore \chi^2$  for the table

$$\begin{aligned}&= \frac{(ad-bc)^2}{N} \left\{ \frac{1}{(a+b)(a+c)} + \frac{1}{(a+b)(b+d)} + \frac{1}{(a+c)(c+d)} + \frac{1}{(b+d)(c+d)} \right\} \\ &= \frac{(ad-bc)^2}{N(a+b)(c+d)(a+c)(b+d)} \{ (b+d)(c+d) + (a+c)(c+d) \\ &\quad + (a+b)(b+d) + (a+b)(a+c) \} \\ &= \frac{(ad-bc)^2 \left[ \sum a^2 + 2 \sum ab \right]}{N(a+b)(c+d)(a+c)(b+d)} = \frac{(ad-bc)^2 (a+b+c+d)^2}{N(a+b)(c+d)(a+c)(b+d)} \\ &= \frac{N(ad-bc)^2}{(a+b)(c+d)(a+c)(b+d)}\end{aligned}\tag{1}$$

Using (1) for the contingency table in Example 11,

we get  $\chi^2 = \frac{253(83 \times 68 - 45 \times 57)^2}{140 \times 113 \times 128 \times 125} = 9.48$ .

### Example 13

Two batches each of 12 animals are taken for test of inoculation. One batch was inoculated and the other batch was not inoculated. The numbers of dead and surviving animals are given in Table 10.25 in both cases. Can the inoculation be regarded as effective against the disease? Make Yate's correction for continuity of  $\chi^2$ .

Table 10.25

	Dead	Survived	Total
Inoculated	2	10	12
Not inoculated	8	4	12
Total	10	14	24

Note on Yate's correction for continuity of  $\chi^2$ .

The  $\chi^2$ -table was prepared using the theoretical  $\chi^2$ -distribution which is continuous, whereas the approximate values of  $\chi^2$  that we are using are discrete. To rectify this defect, Yates has shown that, when

$$\chi^2 = \sum \left[ \frac{\left\{ |O_i - E_i| - \frac{1}{2} \right\}^2}{E_i} \right]$$

is used, the  $\chi^2$ -approximation is improved. Yate's correction is used only when  $v = 1$  and, hence, for a  $2 \times 2$  contingency table. It is used only when some cell frequency is small, i.e., less than 5.

In the present problem, two cell frequencies are less than 5 each. Hence, we apply Yate's correction (Table 10.26).

**Table 10.26**

$O$	$E$	$ O - E  - 0.5$	$\{ O - E  - 0.5\}^2/E$
2	$\frac{12 \times 10}{24} = 5$	2.5	$6.25/5 = 1.25$
10	$\frac{12 \times 14}{24} = 7$	2.5	$6.25/7 = 0.89$
8	$\frac{12 \times 10}{24} = 5$	2.5	$6.25/5 = 1.25$
4	$\frac{12 \times 14}{24} = 7$	2.5	$6.25/7 = 0.89$
	$v = (2 - 1)(2 - 1) = 1$		$\chi^2 = 4.28$

From the  $\chi^2$ -table,  $\chi^2_{5\%}(v = 1) = 3.84$

If  $H_0$  : Inoculation and effect on the diseases are independent, then  $H_0$  is rejected as  $\chi^2 > \chi^2_{5\%}$  i.e. Inoculation can be regarded as effective against the disease.

**Note**

Even if Yate's correction is not made, we would have arrived at the same conclusion.

**Example 14**

A total number of 3759 individuals were interviewed in a public opinion survey on a political proposal. Of them, 1872 were men and the rest women. 2257 individuals were in favour of the proposal and 917 were opposed to it. 243 men were undecided and 442 women were opposed to the proposal. Do you justify or contradict the hypothesis that there is no association between sex and attitude?

A careful analysis of the problem results in the following contingency (Table 10.27).

**Table 10.27**

	<i>favoured</i>	<i>Opposed</i>	<i>Undecided</i>	<i>Total</i>
Men	1154	475	243	1872
Women	1103	442	342	1887
Total	2257	917	585	3759

$H_0$ : Sex and attitude are independent, i.e., there is no association between sex and attitude.

**Table 10.28**

<i>O</i>	<i>E (rounded E)</i>	$(O - E)^2/E$
1154	$\frac{1872 \times 2257}{3759} \approx 1124$	$302/1124 = 0.80$
475	$\frac{1872 \times 917}{3759} \approx 457$	$182/457 = 0.71$
243	$\frac{1872 \times 585}{3759} \approx 291$	$482/291 = 7.92$
1103	$\frac{1887 \times 2257}{3759} \approx 1133$	$302/1133 = 0.79$
442	$\frac{1887 \times 917}{3759} \approx 460$	$182/460 = 0.70$
342	$\frac{1887 \times 585}{3759} \approx 294$	$482/294 = 7.84$
	$v = (3 - 1)(2 - 1) = 2$	$\chi^2 = 18.76$

From the  $\chi^2$ -table,  $\chi^2_{5\%}(v = 2) = 5.99$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected.

That is, sex and attitude are not independent i.e. there is some association between sex and attitude.

### Example 15

The following table gives for a sample of married women, the level of education and the marriage adjustment score:

**Table 10.29**

Level of education	Marriage adjustment				
	<i>Very low</i>	<i>low</i>	<i>high</i>	<i>very high</i>	<i>Total</i>
College	24	97	62	58	241
High school	22	28	30	41	121
Middle school	32	10	11	20	73
Total	78	135	103	119	435

Can you conclude from the above data that the higher the level of education, the greater is the degree of adjustment in marriage?

$H_0$  : There is no relation between the level of education and adjustment in marriage.

$$v = (4 - 1)(3 - 1) = 6$$

$$\chi^2_{5\%} (v = 6) = 12.59$$

**Table 10.30**

$O$	$E$ (rounded)	$(O-E)^2/E$
24	43	$19^2/43 = 8.40$
97	75	$22^2/75 = 6.45$
62	57	$5^2/57 = 0.44$
58	66	$8^2/66 = 0.97$
22	22	$0^2/22 = 0.00$
28	37	$9^2/37 = 2.19$
30	29	$1^2/29 = 0.03$
41	33	$8^2/33 = 1.94$
32	13	$19^2/13 = 27.77$
10	23	$13^2/23 = 7.35$
11	17	$6^2/17 = 2.12$
20	20	$0^2/20 = 0.00$
$\chi^2_{5\%} (v = 6) = 12.59$		$\chi^2 = 57.66$

Since  $\chi^2 > \chi^2_{5\%}$ ,  $H_0$  is rejected.

That is, the level of education and adjustment in marriage are associated.

Thus, we may conclude that the higher the level of education, the greater is the degree of adjustment in marriage.

### Exercise 10(C)

#### Part-A (Short-answer Questions)

1. Define Chi-square distribution.
2. Write down the probability density function of the  $\chi^2$ -distribution.
3. State the important properties of  $\chi^2$ -distribution.
4. Give two uses of  $\chi^2$ -distribution.
5. What is  $\chi^2$ -test of goodness of fit?
6. State the conditions under which  $\chi^2$ -test of goodness of fit is valid.
7. What is  $\chi^2$ -test of independence of attributes?
8. What is contingency table?
9. Write down the value of  $\chi^2$  for a  $2 \times 2$  contingency table with cell frequencies  $a, b, c$  and  $d$ .

10. What is Yate's correction for continuity of  $\chi^2$ ?

**Part-B**

11. In 250 digits from the lottery numbers, the frequencies of the digits were as follows:

Digit:	0	1	2	3	4	5	6	7	8	9
Frequency:	23	25	20	23	23	22	29	25	33	27

Test the hypothesis that the digits were randomly drawn.

12. The following table gives the number of fatal road accidents that occurred during the seven days of the week. Find whether the accidents are uniformly distributed over the week.

Day :	Sun	Mon	Tues	Wed	Thu	Fri	Sat
Number :	8	14	16	12	11	14	9

13. In 120 throws of a single dice, the following distribution of faces are obtained:

Face :	1	2	3	4	5	6
Frequency :	30	25	18	10	22	15

Do these results support the equal probability hypothesis?

14. The number of demands for a particular spare part in a shop was found to vary from day to day. In a sample study, the following information was obtained:

Day :	Mon	Tues	Wed	Thu	Fri	Sat
No. of demands :	124	125	110	120	126	115

Test the hypothesis that the number of parts demanded does not depend on the day of the week.

15. According to genetic theory, children having one parent of blood type  $M$  and the other of blood type  $N$  will always be one of the three types- $M$ ,  $MN$  and  $N$  and the average proportions of these types will be  $1 : 2 : 1$ . Out of 300 children, having one  $M$  parent and one  $N$  parent, 30 per cent were found to be of type  $M$ , 45 per cent of type  $MN$  and the remaining of type  $N$ . Test the genetic theory by  $\chi^2$ -test.

16. 5 coins are tossed 256 times. The number of heads observed is given below. Examine if the coins are true.

No. of heads :	0	1	2	3	4	5
Frequency :	5	35	75	84	45	12

17. 5 dice were thrown 243 times and the numbers of times 1 or 2 was thrown ( $x$ ) are given below:

$x :$	0	1	2	3	4	5
Frequency :	30	75	76	47	13	2

Examine if the dice were unbiased.

18. Fit a binomial distribution for the following data and also test the goodness of fit.

$x :$	0	1	2	3	4
$f :$	5	29	36	25	5

19. Fit a binomial distribution for the following data and also test the goodness of fit

$x:$  0 1 2 3 4 5 6 7 8 9  
 $f:$  3 8 11 15 16 14 12 11 9 1

20. Fit a Poisson distribution for the following distribution and also test the goodness of fit.

$x:$  0 1 2 3 4 5 6 7  
 $f:$  314 335 204 86 29 9 3 0

21. Fit a Poisson distribution for the following distribution and also test the goodness of fit.

$x:$  0 1 2 3 4  
 $f:$  123 59 14 3 1

22. The figures given below are (i) the observed frequencies of a distribution, and (ii) the expected frequencies of the normal distribution having the same mean, S.D. and total frequency as in (i).

(i) 1, 12, 66, 220, 495, 792, 924, 792, 495, 220, 66, 12, 1  
(ii) 2, 15, 66, 210, 484, 799, 943, 799, 484, 210, 66, 15, 2

Do you think that the normal distribution provides a good fit to the data?

23. Fit a normal distribution to the following data and find also the goodness of fit.

$x:$  4 6 8 10 12 14 16 18 20 22 24  
 $f:$  1 7 15 22 35 43 38 20 13 5 1

24. In an epidemic of certain disease, 92 children contacted the disease. Of these 41 received no treatment and of these 10 showed after effects. Of the remainder who did receive the treatment, 17 showed after effects. Test the hypothesis that the treatment was not effective.

25. Out of 1660 candidates who appeared for a competitive examination, 422 were successful. Out of these, 256 had attended a coaching class and 150 of them came out successful. Examine whether coaching was effective as regards the success in the examination.

26. In a pre-poll survey, out of 1000 rural voters, 620 favoured A and the rest B. Out of 1000 urban voters, 450 favoured B and the rest A. Examine if the nature of the area is related to voting preference.

27. The following information was obtained in a sample of 40 small general shops:

Table 10.31

	<i>Shops in urban areas</i>	<i>Shops in rural areas</i>
Owned by men	17	18
Owned by women	3	12

Can it be said that there are more women owners in rural areas than in urban areas? Use Yate's correction for continuity.

28. A certain drug is claimed to be effective in curing cold. In an experiment on 500 persons with cold, half of them were given the drug and half of them were given the sugar pills. The patients' reaction to the treatment are recorded in the following table.

**Table 10.32**

	<i>Helped</i>	<i>Harmed</i>	<i>No effect</i>
Drug	150	30	70
Sugar pills	130	40	80

On the basis of this data, can it be concluded that the drug and sugar pills differ significantly in curing cold?

29. A survey of radio listeners' preference for two types of music under various age groups gave the following information.

**Table 10.33**

<i>Type of music</i>	<i>Age group</i>		
	<i>19-25</i>	<i>26-35</i>	<i>Above 36</i>
Carnatic music :	80	60	90
Film music :	210	325	44
Indifferent :	16	45	132

Is preference for type of music influenced by age?

30. The table given below shows the results of a survey in which 250 respondents were classified according to levels of education and attitude towards students' agitation in a certain town. Test whether the two criteria of classification are independent.

**Table 10.34**

<i>Education</i>	<i>Attitude</i>		
	<i>Against</i>	<i>Neutral</i>	<i>For</i>
Middle school:	40	25	5
High school:	40	20	5
College:	30	15	30
Postgraduate:	15	15	10

## **ANSWERS**

### **Exercise 10(A)**

29.  $z = 2.83$ ; significant      30.  $z = 4.5$ ; the coin is not fair  
 31.  $z = 1.79$ ; claim cannot be supported.

32. No, since  $z (= 1.40) < z_\alpha (= 1.645)$ .
34. No ; (16.7, 27.0)
35. difference due to sampling fluctuations
36.  $z = 0.725$  ; not significant
37.  $z = 6.49$  ; Yes; Yes.
38. Yes, since  $z (= 1.80) > z_{.05} (= 1.645)$
39.  $z = 2.56$  ; the difference cannot be hidden
40.  $z (= .104) < z_{.05} (= 1.645)$  ; the machine has improved
41.  $z (= 3.17) > z_{.05} (= 1.96)$  ; difference significant.
42.  $|z| = 1.02$  ; the claim is valid
43.  $z = 3.12$  ; No
44. No, since  $z = 5$  ; 96.8
45. Claim cannot be true as  $z = 1.89$  and  $z_{5\%} = 1.645$
46. Yes, since  $z (= 1.6) < z_{5\%} (= 1.645)$
47. 41
48. (57.2, 58.8) and (57.0, 59.0)
49. No, since  $z = 2.58$
50. No, since  $z = 8.82$
51. No, since  $z = 1.32$
52. Yes, since  $z (= 4.78) > z_{1\%} (= 2.33)$
53. Yes, at 5% level, since  $|z| (= 1.937) > z_{5\%} (= 1.645)$  and No, at 1% level, since  $|z| (= 1.937) < |z_\alpha| (= 2.33)$
54. (1.98, 6.02)
55.  $H_0 : \mu_1 - \mu_2 = 35$  accepted, as  $z = 1.90$
57. Yes, as  $z = 1.71$
58. Yes, as  $|z| (= 2.5) > z_{1\%} (= 2.33)$
59. Yes, as  $z = 1.70$
60. No, as  $z = 3.61$ .

## Exercise 10(B)

16.  $|t| = 0.62$  ; Yes ;  $83.66 < \mu < 110.74$
17.  $|t| = 2.67$ ; (99.76, 107.74) and (98.22, 109.29)
18.  $t = 2.5$  ; the campaign was successful
19.  $|t| = 1.44$ ;  $t_{10\%} = 1.90$  ;  $\bar{x}$  is not less than  $\mu$ .
20.  $|t| = 0.26$ ; machine is reliable
21. (11.82, 12.38); (11.69, 12.51)
22.  $t = -1.067$  ; No
23.  $|t| = 0.424$  ; No
24.  $t = 4.33$  ; Yes
25.  $|t| = 1.85$  ; Not significant
26.  $t = 4.0$  ; coaching was effective
27.  $F = 1.057$  ; Not significant
28. The populations have the same variance

29. No, though the difference between variances is not significant, the difference between the mean is significant.
30. Yes, as the differences between the means and between the variance are not significant.

### **Exercise 10(C)**

11.  $\chi^2 = 5.2$ ;  $v = 9$ ; digits randomly drawn.
12.  $\chi^2 = 4.17$ ;  $v = 6$ ; accidents occur uniformly.
13.  $\chi^2 = 12.9$ ;  $v = 5$ ; equal probability hypothesis is refuted.
14.  $\chi^2 = 1.68$ ;  $v = 5$ ; the demand does not depend on the day of the week
15.  $\chi^2 = 4.5$ ;  $v = 2$ ; genetic theory may be correct
16.  $\chi^2 = 3.54$ ;  $v = 3$ ; coins are true.
17.  $\chi^2 = 2.76$ ;  $v = 4$ ; dice are unbiased
18.  $E_i : 7, 26, 37, 24, 6$ ;  $\chi^2 = 0.06$ ;  $v = 1$ ; binomial fit is good. ( $E_i$ : 1, 5, 11, 18, 21, 19, 13, 8, 3, 1)
19.  $\chi^2 = 11.30$ ;  $v = 4$ ; binomial fit is not satisfactory;
20.  $E_i : 301, 362, 217, 87, 26, 6, 1$ ;  $\chi^2 = 5.40$ ;  $v = 4$ ; Poisson fit is good.
21.  $E_i : 121, 61, 15, 3, 0$ ;  $\chi^2 = 0.99$ ;  $v = 1$ ; Poisson fit is good.
22.  $x^2 = 3.84$ ;  $v = 8$ ; Normal fit is good.
23.  $E_i : 2, 5, 13, 25, 37, 42, 36, 23, 12, 4, 1$ ;  $\chi^2 = 1.68$ ;  $v = 4$ ; Normal fit is good.
24.  $\chi^2 = 0.85$ ;  $v = 1$ ; No association between treatment and after-effect.
25.  $\chi^2 = 176.12$ ;  $v = 1$ ; coaching was effective.
26.  $\chi^2 = 10.09$ ;  $v = 1$ ; some relation between area and voting preference
27.  $\chi^2 = 2.48$ ;  $v = 1$ ; No, as there is no relation between area and sex of ownership
28.  $\chi^2 = 3.52$ ;  $v = 2$ ; do not differ significantly
29.  $\chi^2 = 373.40$ ;  $v = 4$ ; preference for type of music influenced by age.
30.  $\chi^2 = 35.42$ ;  $v = 6$ ; the two criteria are not independent.



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# Design of Experiments

**B**y ‘experiment’, we mean collection of data (which usually consist of a series of measurement of some feature of an object) for a scientific investigation, according to certain specified sampling procedures. Statistics provides not only the principles and the basis for the proper planning of the experiments but also the methods for proper interpretation of the results of the experiment.

In the beginning, the study of the design of experiments was associated only with agricultural experimentation. The need to save time and money has led to a study of ways to obtain maximum information with the minimum cost and labour. Such motivations resulted in the subsequent acceptance and wide use of the design of experiments and the related analysis of variance techniques in all fields of scientific experimentation. In this chapter we consider some aspects of experimental design briefly and analysis of data from such experiments using analysis of variance techniques.

## Aim of the Design of Experiments

A statistical experiment in any field is performed to verify a particular hypothesis. For example, an agricultural experiment may be performed to verify the claim that a particular manure has got the effect of increasing the yield of paddy. Here the quantity of the manure used and the amount of yield are the two variables involved directly. They are called *experimental variables*. Apart from these two, there are other variables such as the fertility of the soil, the quality of the seed used and the amount of rainfall, which also affect the yield of paddy. Such variables are called *extraneous variables*. The main aim of the design of experiments is to control the extraneous variables and hence to minimise the experimental error so that the results of the experiments could be attributed only to the experimental variables.

## Basic Principles of Experimental Design

In order to achieve the objective mentioned above, the following three principles are adopted while designing the experiments— (1) randomisation, (2) replication, and (3) local control.

### 1. Randomisation

As it is not possible to eliminate completely the contribution of extraneous variables to the value of the response variable (the amount of yield of paddy), we try to control it by randomisation. The group of experimental units (plots of the same size) in which the manure is used is called the *experimental group* and the other group of plots in which the manure is not used and which will provide a basis for comparison is called the *control group*. If any information regarding the extraneous variables and the nature and magnitude of their effect on the response variable in question is not available, we resort to randomisation. That is, we select the plots for the experimental and control groups in a random manner, which provides the most effective way of eliminating any unknown bias in the experiment.

### 2. Replication

In a comparative experiment, in which the effects of different manures on the yield are studied, each manure is used in more than one plot. In other words, we resort to replication which means repetition. It is essential to carry out more than one test on each manure in order to estimate the amount of the experimental error and hence to get some idea of the precision of the estimates of the manure effects.

### 3. Local Control

To provide adequate control of extraneous variables, another essential principle used in the experimental design is the local control. This includes techniques such as grouping, blocking and balancing of the experimental units used in the experimental design. By *grouping*, we mean combining sets of homogeneous plots into groups, so that different manures may be used in different groups. The number of plots in different groups need not necessarily be the same. By *blocking*, we mean assigning the same number of plots in different blocks. The plots in the same block may be assumed to be relatively homogeneous. We use as many manures as the number of plots in a block in a random manner. By *balancing*, we mean adjusting the procedures of grouping, blocking and assigning the manures in such a manner that a balanced configuration is obtained.

## SOME BASIC DESIGNS OF EXPERIMENT

### 1. Completely Randomised Design (CRD)

Let us suppose that we wish to compare ' $h$ ' treatments (use of ' $h$ ' different manures) and there are  $n$  plots available for the experiment.

Let the  $i$ th treatment be replicated (repeated)  $n_i$  times, so that  $n_1 + n_2 + \dots + n_h = n$ .

The plots to which the different treatments are to be given are found by the following randomisation principle. The plots are numbered from 1 to  $n$  serially.  $n$  identical cards are taken, numbered from 1 to  $n$  and shuffled thoroughly. The numbers on the first  $n_1$  cards drawn randomly give the numbers of the plots to which the first treatment is to be given. The numbers on the next  $n_2$  cards drawn at random give the numbers of the plots to which the second treatment is to be given and so on. This design is called a completely randomised design, which is used when the plots are homogeneous or the pattern of heterogeneity of the plots is unknown.

### ***Analysis of Variance (ANOVA)***

The analysis of variance is a widely used technique developed by R.A. Fisher. It enables us to divide the total variation (represented by variance) in a group into parts which are ascribable to different factors and a residual random variation which could not be accounted for by any of these factors. The variation due to any specific factor is compared with the residual variation for significance by applying the F-test, with which the reader is familiar. The details of the procedure will be explained in the sequel.

### ***Analysis of Variance for One Factor of Classification***

Let a sample of  $N$  values of a given random variable  $X$  (representing the yield of paddy) be subdivided into ' $h$ ' classes according to some factor of classification (different manures).

We wish to test the null hypothesis that the factor of classification has no effect on the variable, viz., there is no difference between various classes, viz., the classes are homogeneous. Let  $x_{ij}$  be the value of the  $j^{\text{th}}$  member of the  $i^{\text{th}}$  class, which contains  $n_i$  members. Let the general mean of all the  $N$  values be  $\bar{x}$  and the mean of  $n_i$  values in the  $i^{\text{th}}$  class be  $\bar{x}_i$ .

$$\begin{aligned} \text{Now, } \sum_i \sum_j (x_{ij} - \bar{x})^2 &= \sum_i \sum_j \{(x_{ij} - \bar{x}_i) + (\bar{x}_i - \bar{x})\}^2 \\ &= \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 + \sum_i \sum_j (\bar{x}_i - \bar{x})^2 \\ &\quad + 2 \sum_i \sum_j (x_{ij} - \bar{x}_i)(\bar{x}_i - \bar{x}) \\ &= \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 + \sum_i \sum_{j=1}^{n_i} (x_i - \bar{x})^2 \\ &\quad + 2 \sum_i (\bar{x}_i - \bar{x}) \sum_{j=1}^{n_i} (x_{ij} - \bar{x}_i) \end{aligned}$$

$$= \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 + \sum_i n_i (x_i - \bar{x})^2$$

$\left[ \because \sum_{j=1}^{n_i} (x_{ij} - \bar{x}_i) = \text{sum of the deviations of } n_i \text{ values of the } x_{ij} \text{ in the } i^{\text{th}} \text{ class from their mean} = \bar{x}_i = 0 \right]$

i.e.,  $Q = Q_2 + Q_1$ , say, where

$Q_1 = \sum_i n_i (\bar{x}_i - \bar{x})^2 = \text{sum of the squared deviations of class means from the general mean (variation between classes)}$

$Q_2 = \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 = \text{sum of the squared deviations of variates from the corresponding class means (variation within classes)}$  and  $Q = \text{total variation}$ .

Since  $Q_2 = Q - Q_1$ , viz., the variation  $Q_2$  within classes is got after removing the variation  $Q_1$  between classes from the total variation  $Q$ ,  $Q_2$  is the residual variation.

If  $s^2$  is the variance of a sample of size  $n$  drawn from a population with variance  $\sigma^2$ , then it is known from the theory of estimation that  $\left( \frac{ns^2}{n-1} \right)$  is an unbiased estimate of  $\sigma^2$ .

$$\text{i.e., } E \left( \frac{ns^2}{n-1} \right) = \sigma^2.$$

Since the items in the  $i^{\text{th}}$  class with variance  $\frac{1}{n_i} \sum_{j=1}^{n_i} (x_{ij} - \bar{x}_i)^2$  may be considered as a sample of size  $n_i$  drawn from a population with variance  $\sigma^2$ ,

$$E \left\{ \frac{n_i}{n_i - 1} \cdot \frac{1}{n_i} \sum_{j=1}^{n_i} (x_{ij} - \bar{x}_i)^2 \right\} = \sigma^2$$

$$\text{i.e., } E \left[ \sum_j (x_{ij} - \bar{x}_i)^2 \right] = (n_i - 1) \sigma^2$$

$$\therefore E \left[ \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 \right] = \sum_{i=1}^h (n_i - 1) \sigma^2$$

$$\text{i.e., } E(Q_2) = (N - h) \sigma^2 \text{ or } E \left\{ \frac{Q_2}{N - h} \right\} = \sigma^2$$

i.e.,  $\frac{Q_2}{N-h}$  is an unbiased estimate of  $\sigma^2$  with  $(N-h)$  degrees of freedom.

Now, if we consider the entire group of  $N$  items with variance  $\frac{1}{N} \sum_i \sum_j (x_{ij} - \bar{x})^2$

as a sample of size  $N$  drawn from the same population,

$$E \left\{ \frac{N}{N-1} \frac{1}{N} \sum_i \sum_j (x_{ij} - \bar{x})^2 \right\} = \sigma^2$$

i.e.,  $E \left( \frac{Q}{N-1} \right) = \sigma^2$

i.e.,  $\frac{Q}{N-1}$  is an unbiased estimate of  $\sigma^2$  with  $(N-1)$  degrees of freedom.

Now,  $Q_1 = Q - Q_2$

$$\begin{aligned} \therefore E(Q_1) &= E(Q) - E(Q_2) \\ &= (N-1)\sigma^2 - (N-h)\sigma^2 \\ &= (h-1)\sigma^2 \text{ or } E \left( \frac{Q_1}{h-1} \right) = \sigma^2 \end{aligned}$$

i.e.,  $\frac{Q_1}{h-1}$  is also an unbiased estimate of  $\sigma^2$  with  $(h-1)$  degrees of freedom.

If we assume that the sampled population is normal, then the estimates  $\frac{Q_1}{h-1}$  and  $\frac{Q_2}{N-h}$  are independent and, hence, the ratio  $\frac{Q_1/(h-1)}{Q_2/(N-h)}$  follows a F-distribution with  $(h-1, N-h)$  degrees of freedom or the ratio  $\frac{Q_2/(N-h)}{Q_1/(h-1)}$

follows a F-distribution with  $(N-h, h-1)$  degrees of freedom. Choosing the ratio which is greater than one, we employ the F-test.

If the calculated value of  $F < F_{5\%}$ , the null hypothesis is accepted, viz., different treatments do not contribute significantly different yields.

These results are displayed in the form of a table, called the ANOVA table, as given below:

**Table 11.1 ANOVA table for one factor of classification**

<i>Source of variation (S.V.)</i>	<i>Sum of squares (S.S.)</i>	<i>Degree of freedom (d.f.)</i>	<i>Mean square (M.S.)</i>	<i>Variance ratio (F)</i>
Between classes	$Q_1$	$h - 1$	$Q_1/(h - 1)$	$\frac{Q_1 / (h - 1)}{Q_2 / (N - h)}$ (OR)
Within classes	$Q_2$	$N - h$	$Q_2/(N - h)$	$\frac{Q_2 / (N - h)}{Q_1 / (h - 1)}$
Total	$Q$	$N - 1$	—	—

**Note** For calculating  $Q$ ,  $Q_1$ ,  $Q_2$ , the following computational formulas may be used:

$$\begin{aligned} Q &= N \left\{ \frac{1}{N} \sum \sum x_{ij}^2 - \bar{x}^2 \right\} \\ &= N \left\{ \frac{1}{N} \sum \sum x_{ij}^2 - \left( \frac{1}{N} \sum \sum x_{ij} \right)^2 \right\} \\ &= \sum \sum x_{ij}^2 - \frac{T^2}{N}, \text{ where } T = \sum \sum x_{ij} \end{aligned}$$

Similarly, for the  $i$ th class,

$$\sum_j (x_{ij} - \bar{x}_i)^2 = \sum_j x_{ij}^2 - \frac{T_i^2}{n_i}, \text{ where } T_i = \sum_j x_{ij}.$$

$$\therefore Q_2 = \sum_i \sum_j (x_{ij} - \bar{x}_i)^2 = \sum_i \sum_j x_{ij}^2 - \sum_i \frac{T_i^2}{n_i}$$

Hence,

$$\begin{aligned} Q_1 &= Q - Q_2 \\ &= \sum_i \frac{T_i^2}{n_i} - \frac{T^2}{N} \end{aligned}$$

## 2. Randomised Block Design (RBD)

Let us consider an agricultural experiment using which we wish to test the effect of ' $k$ ' fertilizing treatments on the yield of a crop. We assume that we know some information about the soil fertility of the plots. Then we divide the plots into ' $h$ ' blocks, according to the soil fertility, each block containing ' $k$ ' plots. Thus, the plots in each block will be of homogeneous fertility as far as possible.

Within each block, the ' $k$ ' treatments are given to the ' $k$ ' plots in a perfectly random manner, such that each treatment occurs only once in any block. But the same ' $k$ ' treatments are repeated from block to block. This design is called Randomised Block Design.

### **Analysis of Variance for Two Factors of Classification**

Let the  $N$  variate values  $\{x_{ij}\}$  (representing the yield of paddy) be classified according to two factors. Let there be ' $h$ ' rows (blocks) representing one factor of classification (soil fertility) and ' $k$ ' columns representing the other factor (treatment), so that  $N = hk$ .

We wish to test the null hypothesis that the rows and columns are homogeneous viz., there is no difference in the yields of paddy between the various rows and between the various columns.

Let  $x_{ij}$  be the variate value in the  $i^{\text{th}}$  row and  $j^{\text{th}}$  column.

Let  $\bar{x}$  be the general mean of all the  $N$  values,  $\bar{x}_{i^*}$  be the mean of the  $k$  values in the  $i^{\text{th}}$  row and  $\bar{x}_{*j}$  be the mean of the  $h$  values in the  $j^{\text{th}}$  column.

$$\text{Now, } x_{ij} - \bar{x} = (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x}) + (\bar{x}_{i^*} - \bar{x}) + (\bar{x}_{*j} - \bar{x})$$

$$\begin{aligned} \therefore \sum \sum (x_{ij} - \bar{x})^2 &= \sum \sum (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x})^2 + \sum \sum (\bar{x}_{i^*} - \bar{x})^2 \\ &\quad + \sum \sum (x_{*j} - \bar{x})^2 + 2 \sum \sum (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x})(\bar{x}_{i^*} - \bar{x}) \\ &\quad + 2 \sum \sum (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x})(\bar{x}_{*j} - \bar{x}) \\ &\quad + 2 \sum \sum (\bar{x}_{i^*} - \bar{x})(\bar{x}_{*j} - \bar{x}) \end{aligned} \quad (1)$$

Now, the fourth member in the R.H.S. of (1)

$$\begin{aligned} &= 2 \sum_i (\bar{x}_{i^*} - \bar{x}) \sum_{j=1}^k (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x}) \\ &= 2 \sum_i (\bar{x}_{i^*} - \bar{x})(k \bar{x}_{i^*} - k \bar{x}_{i^*} - k \bar{x} + k \bar{x}) \\ &= 0 \end{aligned}$$

Similarly, the last two members in the R.H.S. of (1) also become zero each.

$$\text{Also, } \sum_i \sum_j (\bar{x}_{i^*} - \bar{x})^2 = k \sum_i (\bar{x}_{i^*} - \bar{x})^2 = Q_1, \text{ say}$$

$$\sum_i \sum_j (\bar{x}_{*j} - \bar{x})^2 = h \sum_j (\bar{x}_{*j} - \bar{x})^2 = Q_2, \text{ say}$$

$$\text{Let } Q = \sum \sum (x_{ij} - \bar{x})^2 \text{ and}$$

$$Q_3 = \sum \sum (x_{ij} - \bar{x}_{i^*} - \bar{x}_{*j} + \bar{x})^2$$

Using all these in (1), we get

$$Q = Q_1 + Q_2 + Q_3, \text{ where}$$

$Q$  = total variation.

$Q_1$  = sum of the squares due to the variations in the rows,

$Q_2$  = that in the columns and

$Q_3$  = that due to the residual variations.

Proceeding as in one factor of classification, we can prove that  $\frac{Q_1}{h-1}$ ,  $\frac{Q_2}{k-1}$ ,

$\frac{Q_3}{(h-1)(k-1)}$  and  $\frac{Q}{hk-1}$  are unbiased estimates of the population variance  $\sigma^2$

with degrees of freedom  $h-1$ ,  $k-1$ ,  $(h-1)(k-1)$  and  $(hk-1)$  respectively. If the sampled population is assumed normal, all these estimates are independent.

$\therefore \frac{Q_1 / (h-1)}{Q_3 / (h-1)(k-1)}$  follows a F-distribution with  $\{h-1, (h-1)(k-1)\}$

degrees of freedom and  $\frac{Q_2 / (k-1)}{Q_3 / (h-1)(k-1)}$  follows a F-distribution with  $\{k-1, (h-1)(k-1)\}$  degrees of freedom. Then the F-tests are applied as usual and the significance of difference between rows and between columns is analysed.

**Table 11.2 The ANOVA table for two factors of classifications**

S.V.	S.S.	d.f.	M.S.	F
Between rows	$Q_1$	$h-1$	$Q_1 / (h-1)$	$\left[ \frac{Q_1 / (h-1)}{Q_3 / (h-1)(k-1)} \right]^{\pm 1}$
Between columns	$Q_2$	$k-1$	$Q_2 / (k-1)$	$\left[ \frac{Q_2 / (k-1)}{Q_3 / (h-1)(k-1)} \right]^{\pm 1}$
Residual	$Q_3$	$(h-1)(k-1)$	$Q_3 / (h-1)(k-1)$	—
Total	$Q$	$hk-1$	—	—

**Note** The following working formulas that can be easily derived may be used to compute  $Q$ ,  $Q_1$ ,  $Q_2$  and  $Q_3$ :

$$1. Q = \sum \sum x_{ij}^2 - \frac{T^2}{N}, \text{ where } T = \sum \sum x_{ij}$$

$$2. Q_1 = \frac{1}{k} \sum T_i^2 - \frac{T^2}{N}, \text{ where } T_i = \sum_{j=1}^k x_{ij}$$

$$3. Q_2 = \frac{1}{b} \sum T_j^2 - \frac{T^2}{N}, \text{ where } T_j = \sum_{i=1}^b x_{ij}$$

$$4. Q_3 = Q - Q_1 - Q_2$$

$$\text{It may be verified that } \sum_i T_i = \sum_j T_j = T.$$

### 3. Latin Square Design (LSD)

We consider an agricultural experiment, in which  $n^2$  plots are taken and arranged in the form of an  $n \times n$  square, such that the plots in each row will be homogeneous as far as possible with respect to one factor of classification, say, soil fertility and plots in each column will be homogeneous as far as possible with respect to another factor of classification, say, seed quality.

Then  $n$  treatments are given to these plots such that each treatment occurs only once in each row and only once in each column. The various possible arrangements obtained in this manner are known as Latin squares of order  $n$ . This design of experiment is called the Latin Square Design.

#### **Analysis of Variance for Three Factors of Classifications**

Let the  $N (= n^2)$  variate values  $\{x_{ij}\}$ , representing the yield of paddy, be classified according to three factors. Let the rows, columns and letters stand for the three factors, say soil fertility, seed quality and treatment respectively.

We wish to test the null hypothesis that the rows, columns and letters are homogeneous, viz., there is no difference in the yield of paddy between the rows (due to soil fertility), between the columns (due to seed quality) and between the letters (due to the treatments).

Let  $x_{ij}$  be the variate value corresponding to the  $i^{\text{th}}$  row,  $j^{\text{th}}$  column and  $k^{\text{th}}$  latter.

Let  $\bar{x} = \frac{1}{n^2} \sum \sum x_{ij}$ ,  $\bar{x}_{i*} = \frac{1}{n} \sum_j x_{ij}$ ,  $\bar{x}_{*j} = \frac{1}{n} \sum_i x_{ij}$  and  $\bar{x}_k$  be the mean of the values of  $x_{ij}$  corresponding to the  $k^{\text{th}}$  treatment.

$$\text{Now, } x_{ij} - \bar{x} = (\bar{x}_{i*} - \bar{x}) + (\bar{x}_{*j} - \bar{x}) + (\bar{x}_k - \bar{x}) + (x_{ij} - \bar{x}_{i*} - \bar{x}_{*j} - \bar{x}_k + 2\bar{x})$$

$$\begin{aligned} \therefore \sum \sum (x_{ij} - \bar{x})^2 &= n \sum_i (\bar{x}_{i*} - \bar{x})^2 + n \sum_j (\bar{x}_{*j} - \bar{x})^2 \\ &\quad + n \sum_k (\bar{x}_k - \bar{x})^2 + \sum_i \sum_j (x_{ij} - \bar{x}_{i*} - \bar{x}_{*j} - \bar{x}_k + 2\bar{x})^2 \end{aligned}$$

( $\because$  all the product terms vanish as in two factor classification)

$$\text{i.e., } Q = Q_1 + Q_2 + Q_3 + Q_4$$

As before, we can prove that  $\frac{Q_1}{n-1}, \frac{Q_2}{n-1}, \frac{Q_3}{n-1}, \frac{Q_4}{(n-1)(n-2)}$  and  $\frac{Q}{n^2-1}$  are unbiased estimates of the population variance  $\sigma^2$  with degrees of freedom  $n-1, n-1, n-1, (n-1)(n-2)$  and  $(n^2-1)$  respectively.

If the sampled population is assumed normal, all these estimates are independent.

$\therefore$  each of  $\frac{Q_1/(n-1)}{Q_4/(n-1)(n-2)}$ ,  $\frac{Q_2/(n-1)}{Q_4/(n-1)(n-2)}$  and  $\frac{Q_3/(n-1)}{Q_4/(n-1)(n-2)}$

follows a F-distribution with  $\{(n-1), (n-1)(n-2)\}$  degrees of freedom.

Then the F-tests are applied as usual and the significance of differences between rows, columns and treatments is analysed.

**Table 11.3 The ANOVA table for three factors of classification**

S.V.	S.S.	d.f.	M.S.	F
Between rows	$Q_1$	$n - 1$	$Q_1 / (n - 1) = M_1$	$\left(\frac{M_1}{M_4}\right)^{\pm 1}$
Between columns	$Q_2$	$n - 1$	$Q_2 / (n - 1) = M_2$	$\left(\frac{M_2}{M_4}\right)^{\pm 1}$
Between letters	$Q_3$	$n - 1$	$Q_3 / (n - 1) = M_3$	$\left(\frac{M_3}{M_4}\right)^{\pm 1}$
Residual	$Q_4$	$(n - 1)(n - 2)$	$Q_4 / (n - 1)(n - 2) = M_4$	—
Total	$Q$	$n^2 - 1$	—	—

**Note** The following working formulas may be used to compute the  $Q$ 's:

$$1. Q = \sum \sum x_{ij}^2 - \frac{T^2}{n^2}, \text{ where } T = \sum \sum x_{ij}$$

$$2. Q_i = \frac{1}{n} \sum T_i^2 - \frac{T^2}{n^2}, \text{ where } T_i = \sum_{j=1}^n x_{ij}$$

$$3. Q_2 = \frac{1}{n} \sum T_j^2 - \frac{T^2}{n^2}, \text{ where } T_j = \sum_{i=1}^n x_{ij}$$

$$4. Q_3 = \frac{1}{n} \sum T_k^2 - \frac{T^2}{n^2}, \text{ where } T_k \text{ is the sum of all } x_{ij} \text{ receiving the } k^{\text{th}} \text{ treatment.}$$

$$5. Q_4 = Q - Q_1 - Q_2 - Q_3.$$

$$\text{Also } T = \sum_i T_i = \sum_j T_j = \sum_k T_k$$

## Comparison of RBD and LSD

1. The number of replications of each treatment is equal to the number of treatments in LSD, whereas there is no such restrictions on treatments and replication in RBD.
2. LSD can be performed on a square field, while RBD can be performed either on a square field or a rectangular field.
3. LSD is known to be suitable for the case when the number of treatments is between 5 and 12, whereas RBD can be used for any number of treatments.
4. The main advantage of LSD is that it controls the effect of two extraneous variables, whereas RBD controls the effect of only one extraneous variable. Hence the experimental error is reduced to a larger extent in LSD than in RBD.

## Note on Simplification of Computational Work

The variance of a set of values is independent of the origin and so a shift of origin does not affect the variance calculations. Hence, in analysis of variance problems, we can subtract a convenient number from the original values and work out the problems with the new values obtained. Also, since we are concerned with the variance ratios, change of scale also may be introduced without affecting the values of F.

### Worked Example 11

#### Example 1

A completely randomised design experiment with 10 plots and 3 treatments gave the following results:

Plot No.	:	1	2	3	4	5	6	7	8	9	10
Treatment	:	A	B	C	A	C	C	A	B	A	B
Yield	:	5	4	3	7	5	1	3	4	1	7

Analyse the results for treatment effects.

**Rearranging the data according to the treatments, we have the following table:**

Treatment	Yield from plots ( $x_{ij}$ )				$T_i$	$T_i^2$	$n_i$	$\frac{T_i^2}{n_i}$
A	5	7	3	1	16	256	4	64
B	4	4	7	—	15	225	3	75
C	3	5	1	—	9	81	3	27
Total				$T = 40$	—	$N = 10$	166	

$$\begin{aligned}\sum \sum x_{ij}^2 &= (25 + 49 + 9 + 1) + (16 + 16 + 49) + (9 + 25 + 1) \\ &= 84 + 81 + 35 = 200\end{aligned}$$

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 200 - \frac{40^2}{10} = 200 - 160 = 40$$

$$Q_1 = \sum \frac{T_i^2}{n_i} - \frac{T^2}{N} = 166 - 160 = 6$$

$$\therefore Q_2 = Q - Q_1 = 40 - 6 = 34$$

### ANOVA table

S.V.	S.S.	d.f.	M.S.	$F_0$
Between classes (treatments)	$Q_1 = 6$	$h - 1 = 2$	3.0	$\frac{4.86}{3.0}$
Within classes	$Q_2 = 34$	$N - h = 7$	4.86	= 1.62
Total	$Q = 40$	$N - 1 = 9$	-	-

From the  $F$ -table,  $F_{5\%}(v_1 = 7, v_2 = 2) = 19.35$

We note that  $F_0 < F_{5\%}$

Let  $H_0$ : The treatments do not differ significantly.

$\therefore$  the null hypothesis is accepted.

i.e., the treatments are not significantly different.

### Example 2

The following table shows the lives in hours of four brands of electric lamps:

*Brand*

A : 1610, 1610, 1650, 1680, 1700, 1720, 1800

B : 1580, 1640, 1640, 1700, 1750

C : 1460, 1550, 1600, 1620, 1640, 1660, 1740, 1820

D : 1510, 1520, 1530, 1570, 1600, 1680

Perform an analysis of variance and test the homogeneity of the mean lives of the four brands of lamps.

We subtract 1640 (= the average of the extreme values) from the given values and work out with the new values of  $x_{ij}$

<i>Brand</i>	<i>Lives of lamps (<math>x_{ij}</math>)</i>								$Ti$	$n_i$	$\frac{T_i^2}{n_i}$
A	-30	-30	10	40	60	80	160	-	290	7	12014
B	-60	0	0	60	110	-	-	-	110	5	2420
C	-180	-90	-40	-20	0	20	100	180	-30	8	113
D	-130	-120	-110	-70	-40	40	-	-	-430	6	30817
	Total								-60	26	45364

$$\begin{aligned}\sum \sum x_{ij}^2 &= (900 + 900 + 100 + 1600 + 3600 + 6400 + 25600) \\&\quad + (3600 + 0 + 0 + 3600 + 12100) \\&\quad + (32400 + 8100 + 1600 + 400 + 0 + 400 + 10000 + 32400) \\&\quad + (16900 + 14400 + 12100 + 4900 + 1600 + 1600) \\&= 39100 + 19300 + 85300 + 51500 = 195200\end{aligned}$$

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 1,95,200 - 138 = 1,95,062$$

$$Q_1 = \sum \frac{T_i^2}{n_i} - \frac{T^2}{n} = 45,364 - 138 = 45,226$$

$$Q_2 = Q - Q_1 = 1,95,062 - 45,226 = 1,49,836$$

#### ANOVA table

S.V.	S.S.	d.f.	M.S.	$F_0$
Between brands	$Q_1 = 45,226$	$h - 1 = 3$	15,075	$\frac{15,075}{6,811}$
Within brands	$Q_2 = 1,49,836$	$N - h = 22$	6,811	= 2.21
Total	$Q = 1,95,062$	$N - 1 = 25$	-	-

From the  $F$ -tables,  $F_{5\%}(v_1 = 3, v_2 = 22) = 3.06$

$$F_0 < F_{5\%}$$

Hence, the null hypothesis  $H_0$ , namely, the means of the lives of the four brands are homogeneous, is accepted viz., the lives of the four brands of lamps do not differ significantly.

**Note**

We could have used a change of scale also. viz., we could have made the change  
old  $x_{ij} - 1640$   
New  $x_{ij} = \frac{\text{old } x_{ij} - 1640}{10}$  and simplified the numerical work still further

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#### Example 3

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A car rental agency, which uses 5 different brands of tyres in the process of deciding the brand of tyre to purchase as standard equipment for its fleet, finds that each of 5 tyres of each brand last the following number of kilometres (in thousands):

**Tyre brands**

A	B	C	D	E
36	46	35	45	41
37	39	42	36	39
42	35	37	39	37
38	37	43	35	35
47	43	38	32	38

Test the hypothesis that the five tyre brands have almost the same average life.

We shift the origin to 40 and work out with the new values of  $x_{ij}$ .

Tyre brand	$x_{ij}$					$T_i$	$n_i$	$\frac{T_i^2}{n_i}$	$\sum_{j=1}^5 x_{ij}^2$
A	-4	-3	2	-2	7	0	5	0	82
B	6	-1	-5	-3	3	0	5	0	80
C	-5	2	-3	3	-2	-5	5	5	51
D	5	-4	-1	-5	-8	-13	5	33.8	131
E	1	-1	-3	-5	-2	-10	5	20	40
Total					-28	25	58.8	384	

$$T = \sum_i T_i = -28; \sum \sum x_{ij}^2 = \sum_i \left( \sum_j x_{ij}^2 \right) = 384$$

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 384 - \frac{(-28)^2}{25} = 352.64$$

$$Q_1 = \sum \frac{T_i^2}{n_i} - \frac{T^2}{N} = 58.8 - 31.36 = 27.44$$

$$Q_2 = Q - Q_1 = 352.64 - 27.44 = 325.20$$

**ANOVA table**

S.V.	S.S.	d.f.	M.S.	F <sub>0</sub>
Between tyre brands	$Q_1 = 27.44$	$h - 1 = 4$	6.86	$\frac{16.26}{6.86}$
Within tyre brands	$Q_2 = 325.20$	$N - h = 20$	16.26	$= 2.37$
Total	$Q = 352.64$	$N - 1 = 24$	-	-

From the F-tables,  $F_{5\%} (v_1 = 20, v_2 = 4) = 5.80$

$$F_0 < F_{5\%}$$

Hence, ' $H_0$ ' : the five tyre brands have almost the same average life is accepted viz., the five tyre brands do not differ significantly in their lives.

### Example 4

In order to determine whether there is significant difference in the durability of 3 makes of computers, samples of size 5 are selected from each make and the frequency of repair during the first year of purchase is observed. The results are as follows:

Makes		
A	B	C
5	8	7
6	10	3
8	11	5
9	12	4
7	4	1

In view of the above data, what conclusion can you draw?

Make	$x_{ij}$					$T_i$	$n_i$	$\frac{T_i^2}{n_i}$	$\sum_j x_{ij}^2$
A	5	6	8	9	7	35	5	245	255
B	8	10	11	12	4	45	5	405	445
C	7	3	5	4	1	20	5	80	100
	Total					100	15	730	800

$$T = \sum T_i = 100; \sum \sum x_{ij}^2 = 800; N = \sum n_i = 15$$

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 800 - \frac{100^2}{15} = 133.33$$

$$Q_1 = \sum \frac{T_i^2}{n_i} - \frac{T^2}{N} = 730 - 666.67 = 63.33$$

$$Q_2 = Q - Q_1 = 70$$

**ANOVA table**

S.V.	S.S.	d.f.	M.S.	F <sub>0</sub>
Between makes	$Q_1 = 63.33$	$h - 1 = 2$	31.67	$\frac{31.67}{5.83}$
Within makes	$Q_2 = 70$	$N - h = 12$	5.83	5.43
Total	$Q = 133.33$	$N - 1 = 14$	—	—

From the  $F$ -tables,  $F_{5\%} (v_1 = 2, v_2 = 12) = 3.88$

$$F_0 > F_{5\%}$$

Hence, the null hypothesis ( $H_0$ : the 3 makes of computers do not differ in the durability) is rejected.

viz., there is significant difference in the durability of the 3 makes of computers.

### Example 5

Three varieties of a crop are tested in a randomised block design with four replications, the layout being as given below: The yields are given in kilograms. Analyse for significance.

C48	A51	B52	A49
A47	B49	C52	C51
B49	C53	A49	B50

Rewriting the data such that the rows represent the blocks and the columns represent the varieties of the crop (as assumed in the discussion of analysis of variance for two factors of classification), we have the following table:

Blocks	Crops		
	A	B	C
1	47	49	48
2	51	49	53
3	49	52	52
4	49	50	51

We shift the origin to 50 and work out with the new values of  $x_{ij}$ .

### Crops

Blocks	A	B	C	$T_i$	$T_i^2/k$	$\sum_j x_{ij}^2$
1	-3	-1	-2	-6	$36/3 = 12$	14
2	1	-1	3	3	$9/3 = 3$	11
3	-1	2	2	3	$9/3 = 3$	9
4	-1	0	1	0	$0/3 = 0$	2
$T_j$	-4	0	4	$T = 0$	$\sum \frac{T_i^2}{k} = 18$	36
$T_j^2/h$	$\frac{16}{4} = 4$	$\frac{0}{4} = 0$	$\frac{16}{4} = 4$	$\sum \frac{T_i^2}{h} = 8$		
$\sum_i x_{ij}^2$	12	6	18	36		

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 36 - \frac{0^2}{12} = 36$$

$$Q_1 = \frac{1}{k} \sum T_i^2 - \frac{T^2}{N} = 18 - 0 = 18$$

$$Q_2 = \frac{1}{h} \sum T_j^2 - \frac{T^2}{N} = 8 - 0 = 8$$

$$Q_3 = Q - Q_1 - Q_2 = 36 - 18 - 8 = 10$$

### ANOVA table

S.V.	S.S.	d.f.	M.S.	$F_0$
Between rows (blocks)	$Q_1 = 18$	$h - 1 = 3$	6	$\frac{6}{1.67} = 3.6$
Between columns (crops)	$Q_2 = 8$	$k - 1 = 2$	4	$\frac{4}{1.67} = 2.4$
Residual	$Q_3 = 10$	$(h - 1)(k - 1) = 6$	1.67	—
Total	$Q = 36$	$hk - 1 = 11$	—	—

From  $F$ -tables,  $F_{5\%}(v_1 = 3, v_2 = 6) = 4.76$  and  $F_{5\%}(v_1 = 2, v_2 = 6) = 5.14$   
Considering the difference between rows, we see that  $F_0 (= 3.6) < F_{5\%} (= 4.76)$

Hence, the difference between the rows is not significant. ( $H_0$  is accepted)  
viz., the blocks do not differ significantly with respect to the yield.

Considering the difference between columns, we see that  $F_0 (= 2.4) < F_{5\%} (= 5.14)$

Hence, the difference between the columns is not significant. ( $H_0$  is accepted)  
viz., the varieties of crop do not differ significantly with respect to the yield.

### Example 6

Five breeds of cattle  $B_1, B_2, B_3, B_4, B_5$  were fed on four different rations  $R_1, R_2, R_3, R_4$ . Gains in weight in kg. over a given period were recorded and given below:

	$B_1$	$B_2$	$B_3$	$B_4$	$B_5$
$R_1$	1.9	2.2	2.6	1.8	2.1
$R_2$	2.5	1.9	2.3	2.6	2.2
$R_3$	1.7	1.9	2.2	2.0	2.1
$R_4$	2.1	1.8	2.5	2.3	2.4

Is there a significant difference between (a) breeds, and (b) rations?

We effect the change of origin and scale using  $y_{ij} = \frac{x_{ij} - 2}{1/10} = 10(x_{ij} - 2)$  and work out with  $y_{ij}$  values.

	$B_1$	$B_2$	$B_3$	$B_4$	$B_5$	$T_i$	$\frac{T_i^2}{k}$	$\sum_j y_{ij}^2$
$R_1$	-1	2	6	-2	1	6	7.2	46
$R_2$	5	-1	3	6	2	15	45.0	75
$R_3$	-3	-1	2	0	1	-1	0.2	15
$R_4$	1	-2	5	3	4	11	24.2	55
$T_j$	2	-2	16	7	8	$T = 31$	$\sum \frac{T_i^2}{k} = 76.6$	191
$T_j^2 / h$	1	1	64	12.25	16	$\sum T_j^2 / h = 94.5$		
$\sum_i y_{ij}^2$	36	10	74	49	22	191		

$$Q = \sum \sum y_{ij}^2 - \frac{T^2}{N} = 191 - \frac{(31)^2}{20} = 142.95$$

$$Q_1 = \frac{1}{k} \sum T_i^2 - \frac{T^2}{N} = 76.6 - 48.05 = 28.55$$

$$Q_2 = \frac{1}{h} \sum T_j^2 - \frac{T^2}{N} = 94.25 - 48.05 = 46.20$$

$$Q_3 = Q - Q_1 - Q_2 = 142.95 - (28.55 + 46.20) = 68.20$$

#### ANOVA table

S.V.	S.S.	d.f.	M.S.	$F_0$
Between rows (rations)	$Q_1 = 28.55$	$h - 1 = 3$	9.52	$9.52/5.68 = 1.68$
Between Cols. (breeds)	$Q_2 = 46.20$	$K - 1 = 4$	11.55	$11.55/5.68 = 2.03$
Residual	$Q_3 = 68.20$	$(h - 1)(K - 1) = 12$	5.68	-
Total	$Q = 142.95$	$hk - 1 = 19$	-	-

From the F-tables,  $F_{5\%} (v_1 = 3, v_2 = 12) = 3.49$  and  $F_{5\%} (v_1 = 4, v_2 = 12) = 3.26$

With respect to the rows,  $F_0 (= 1.68) < F_{5\%} (= 3.49)$

With respect to the columns,  $F_0 (= 2.03) < F_{5\%} (= 3.26)$

Hence, the difference between the rations and that between the breeds are not significant.

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### Example 7

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The following data represent the number of units of production per day turned out by 5 different workers using 4 different types of machines:

	<b>Machine Type</b>			
	<i>A</i>	<i>B</i>	<i>C</i>	<i>D</i>
1	44	38	47	36
2	46	40	52	43
Workers:	3	34	36	44
	4	43	38	46
	5	38	42	49
				39

- (a) Test whether the five men differ with respect to mean productivity.
- (b) Test whether the mean productivity is the same for the four different machine types.

We subtract 40 from the given values and work out with new values of  $x_{ij}$ .

<i>Worker</i>	<i>Machine Type</i>				$T_i$	$T^2_i/k$	$\sum_j x_{ij}^2$
	<i>A</i>	<i>B</i>	<i>C</i>	<i>D</i>			
1	4	−2	7	−4	5	6.25	85
2	6	0	12	3	21	110.25	189
3	−6	−4	4	−8	−14	49.00	132
4	3	−2	6	−7	0	0	98
5	−2	2	9	−1	8	16.00	90
$T_j$	5	−6	38	−17	$T = 20$	$\sum \frac{T_j^2}{k} = 181.5$	594
$T^2 j/h$	5	7.2	288.8	57.8	$\sum T_j^2 / h = 358.8$		
$\sum_i x_{ij}^2$	101	28	326	139	594		

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 594 - \frac{400}{20} = 574$$

$$Q_1 = \sum \frac{T_i^2}{k} - \frac{T^2}{N} = 181.5 - 20 = 161.5$$

$$Q_2 = \sum \frac{T_j^2}{h} - \frac{T^2}{N} = 358.8 - 20 = 338.8$$

$$Q_3 = Q - Q_1 - Q_2 = 574 - (161.5 + 338.8) = 73.7$$

**ANOVA table**

S.V.	S.S.	d.f.	M.S.	$F_0$
Between rows (workers)	$Q_1 = 161.5$	$h - 1 = 4$	40.375	$\frac{40.375}{6.142} = 6.57$
Between Cols. (machines)	$Q_2 = 338.8$	$k - 1 = 3$	112.933	$\frac{112.933}{6.142} = 18.39$
Residual	$Q_3 = 73.7$	$(h - 1)(k - 1) = 12$	6.142	—
Total	$Q = 574$	$hk - 1 = 19$	—	—

From the  $F$ -tables,  $F_{5\%}(v_1 = 4, v_2 = 12) = 3.26$

and  $F_{5\%}(v_1 = 3, v_2 = 12) = 3.49$

With respect to the rows,  $F_0(= 6.57) > F_{5\%}(= 3.26)$

With respect to the columns,  $F_0(= 18.39) > F_{5\%}(= 3.49)$

Hence, the 5 workers differ significantly and the 4 machine types also differ significantly with respect to mean productivity.

### Example 8

Four doctors each test four treatments for a certain disease and observe the number of days each patient takes to recover. The results are as follows (recovery time in days)

Doctor	Treatment			
	1	2	3	4
A	10	14	19	20
B	11	15	17	21
C	9	12	16	19
D	8	13	17	20

Discuss the difference between (a) doctors, and (b) treatments.

We subtracted 15 from the given values and work out with the new values of  $x_{ij}$ .

Doctor	Treatment				$T_i$	$\frac{T_i^2}{k}$	$\sum_j x_{ij}^2$
	1	2	3	4			
A	-5	-1	4	5	3	2.25	67
B	-4	0	2	6	4	4.00	56
C	-6	-3	1	4	-4	4.00	62
D	-7	-2	2	5	-2	1.00	82
$T_j$	-22	-6	9	20	$T = 1$	$\sum \frac{T_j^2}{k} = 11.25$	267
$T_j^2/h$	121	9	20.25	100	$\sum T_j^2/h = 250.25$		
$\sum_i x_{ij}^2$	126	14	25	102	267		

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 267 - \frac{1}{16} = 266.94$$

$$Q_1 = \sum \frac{T_i^2}{k} - \frac{T^2}{N} = 11.25 - 0.0625 = 11.19$$

$$Q_2 = \sum \frac{T_j^2}{h} - \frac{T^2}{N} = 250.25 - 0.0625 = 250.19$$

$$Q_3 = Q - Q_1 - Q_2 = 266.94 - 261.38 = 5.56$$

**ANOVA table**

S.V.	S.S.	d.f.	M.S.	$F_0$
Between rows (doctors)	$Q_1 = 11.19$	$h - 1 = 3$	3.73	$\frac{3.73}{0.62} = 6.02$
Between cols. (treatments)	$Q_2 = 250.19$	$k - 1 = 3$	83.40	$\frac{83.40}{0.62} = 134.52$
Residual	$Q_3 = 5.56$	$(h - 1)(k - 1) = 9$	0.62	—
Total	$Q = 266.94$	$hk - 1 = 15$	—	—

From the  $F$ -tables,  $F_{5\%} (v_1 = 3, v_2 = 9) = 3.86$

Since  $F_0 > F_{5\%}$  with respect to rows and columns, the difference between the doctors is significant and that between the treatments is highly significant.

### Example 9

The following data resulted from an experiment to compare three burners  $B_1$ ,  $B_2$  and  $B_3$ . A Latin square design was used as the tests were made on 3 engines and were spread over 3 days.

	Engine 1	Engine 2	Engine 3
Day 1	$B_1 - 16$	$B_2 - 17$	$B_3 - 20$
Day 2	$B_2 - 16$	$B_3 - 21$	$B_1 - 15$
Day 3	$B_3 - 15$	$B_1 - 12$	$B_2 - 13$

Test the hypothesis that there is no difference between the burners.

We subtract 16 from the given values and work out with new values of  $x_{ij}$ .

	$E_1$	$E_2$	$E_3$	$T_i$	$\frac{T_i^2}{n}$	$\sum_j x_{ij}^2$
$D_1$	$0(B_1)$	$1(B_2)$	$4(B_3)$	5	8.33	17
$D_2$	$0(B_2)$	$5(B_3)$	$-1(B_1)$	4	5.33	26
$D_3$	$-1(B_3)$	$-4(B_1)$	$-3(B_2)$	-8	21.33	26
$T_j$	-1	2	0	$T = 1$	$\sum T_i^2/n = 35$	69
$T_j^2/n$	0.33	1.33	0	$\sum T_i^2/n = 1.66$		
$\sum_i x_{ij}^2$	1	42	26	69		

Rearranging the data values according to the burners, we have

<i>Burner</i>	$x_k$			$T_k$	$T_k^2/n$
$B_1$	0	-1	-4	-5	8.33
$B_2$	1	0	-3	-2	1.33
$B_3$	4	5	-1	8	21.33
Total			$T = 1$	$\sum \frac{T_k^2}{n} = 31$	

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 69 - \frac{1}{9} = 68.89$$

$$Q_1 = \frac{1}{n} \sum T_i^2 - \frac{T^2}{N} = 35 - \frac{1}{9} = 34.89$$

$$Q_2 = \frac{1}{n} \sum T_j^2 - \frac{T^2}{N} = 1.67 - \frac{1}{9} = 1.56$$

$$Q_3 = \frac{1}{n} \sum T_k^2 - \frac{T^2}{N} = 31 - \frac{1}{9} = 30.89$$

$$Q_4 = Q - Q_1 - Q_2 - Q_3 = 1.55$$

***ANOVA table***

S.V.	S.S.	d.f.	M.S.	$F_0$
Between rows (days)	$Q_1 = 34.89$	$n - 1 = 2$	17.445	$\frac{17.445}{0.775} = 22.51$
Between Cols. (engines)	$Q_2 = 1.56$	$n - 1 = 2$	0.780	$\frac{0.780}{0.775} = 1.01$
Between letters (burners)	$Q_3 = 30.89$	$n - 1 = 2$	15.445	$\frac{15.445}{0.775} = 19.93$
Residual	$Q_4 = 1.55$	$(n - 1)(n - 2) = 2$	0.775	—
Total	$Q = 68.89$	$n^2 - 1 = 8$	—	—

From the  $F$ -tables,  $F_{5\%} (v_1 = 2, v_2 = 2) = 19.00$

Since  $F_0 (= 19.93) > F_{5\%} (= 19.00)$  for the burners, there is significant difference between the burners.

Incidentally, since  $F_0 > F_{5\%}$  for the rows, the difference between the days is significant and since  $F_0 < F_{5\%}$  for the columns, the difference between the engine is not significant.

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### **Example 10**

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Analyse the variance in the following Latin square of yields (in kgs) of paddy where  $A, B, C, D$  denote the different methods of cultivation:

D122	A121	C123	B122
B124	C123	A122	D125
A120	B119	D120	C121
C122	D123	B121	A122

Examine whether the different methods of cultivation have given significantly different yields.

We subtract 120 from the given values and work out with the new values of  $x_{ij}$ .

$i \backslash j$	$I$	2	3	4	$T_i$	$T_i^2/n$	$\sum_j x_{ij}^2$
1	D2	A1	C3	B2	8	16	18
2	B4	C3	A2	D5	14	49	54
3	A0	B-1	D0	C1	0	0	2
4	C2	D3	B1	A2	8	16	18
$T_j$	8	6	6	10	$T = 30$	$\sum T_i^2/n = 81$	92
$T_j^2/n$	16	9	9	25	$\sum T_j^2/n = 59$		
$\sum_i x_{ij}^2$	24	20	14	34	92		

Rearranging the data according to the letters, we have

Letter	$x_k$				$T_k$	$T_k^2/n$
A	1	2	0	2	5	6.25
B	2	4	-1	1	6	9.00
C	3	3	1	2	9	20.25
D	2	5	0	3	10	25.00
Total					30	60.50

$$Q = \sum \sum x_{ij}^2 - \frac{T^2}{N} = 92 - \frac{30^2}{16} = 35.75$$

$$Q_1 = \frac{1}{n} \sum T_i^2 - \frac{T^2}{N} = 81 - 56.25 = 24.75$$

$$Q_2 = \frac{1}{n} \sum T_j^2 - \frac{T^2}{N} = 59 - 56.25 = 2.75$$

$$Q_3 = \frac{1}{n} \sum T_k^2 - \frac{T^2}{N} = 60.50 - 56.25 = 4.25$$

$$\begin{aligned} Q_4 &= Q - Q_1 - Q_2 - Q_3 = 35.75 - (24.75 + 2.75 + 4.25) \\ &= 4.0 \end{aligned}$$

***ANOVA table***

S.V.	S.S.	d.f.	M.S.	F <sub>0</sub>
Between rows	$Q_1 = 24.75$	$n - 1 = 3$	8.25	
Between columns	$Q_2 = 2.75$	$n - 1 = 3$	0.92	
Between letters	$Q_3 = 4.25$	$n - 1 = 3$	1.42	$\frac{1.42}{0.67} = 2.12$
Residual	$Q_4 = 4.0$	$(n - 1)(n - 2) = 6$	0.67	—
Total	$Q = 35.75$	$n^2 - 1 = 15$	—	—

From the F-tables,  $F_{5\%} (v_1 = 3, v_2 = 6) = 4.76$ .

Since  $F_0 (= 2.12) < F_{5\%} (= 4.76)$  with respect to the letters, the difference between the methods of cultivation is not significant.

### Exercise 11

#### **Part-A (Short-answer Questions)**

1. What do you mean by the term ‘experiment’ in Design of experiments?
2. What motivated the adoption of design of experiments technique in scientific problems?
3. What is the aim of the design of experiments?
4. Distinguish between experimental and extraneous variables.
5. Name the basic principles of experimental design.
6. What do you mean by experimental group and control group?
7. What are the techniques frequently used in the local control of extraneous variables?
8. Name three basic designs of experiment.
9. What do you mean by analysis of variance.
10. Explain completely randomised design briefly.
11. Write down the format of the ANOVA table for one factor of classification.
12. Explain randomised block design briefly.
13. Write down the format of the ANOVA table for two factors of classification.
14. Explain Latin square design briefly.
15. Is a  $2 \times 2$  Latin square design possible? Why?  
[Hint: No, as the degree of freedom for the residual variation is zero]
16. Write down the format of ANOVA table for three factors of classification.
17. Compare RBD and LSD.
18. What is the main advantage of LSD over RBD?

19. What is the total number of all possible Latin squares of order 3?
20. What is the total number of all possible Latin squares of order 4?

**Part B**

21. The following tables gives the yields of wheat from 16 plots, all of approximately equal fertility, when 4 varieties of wheat were cultivated in a completely randomised fashion. Test the hypothesis that the varieties are not significantly different.

Plot No. :	1	2	3	4	5	6	7	8	9	10
Variety :	A	B	D	C	B	C	A	D	B	D
Yield :	32	34	29	31	33	34	34	26	36	30
Plot No. :	11	12	13	14	15	16				
Variety :	A	C	B	A	B	C				
Yield :	33	35	37	35	35	32				

22. A random sample is selected from each of 3 makes of ropes and their breaking strength (in certain units) are measured with the following results:

I : 70, 72, 75, 80, 83

II : 60, 65, 57, 84, 87, 73

III : 100, 110, 108, 112, 113, 120, 107

Test whether the breaking strength of the ropes differ significantly.

23. The weights in gm of a number of copper wires, each of length 1 metre, were obtained. These are shown classified according to the dye from which they come:

$D_1$ :	1.30,	1.32,	1.36,	1.35,	1.32,	1.37
$D_2$ :	1.28,	1.35,	1.33,	1.34		
$D_3$ :	1.32,	1.29,	1.31,	1.28,	1.33,	1.30
$D_4$ :	1.31,	1.29,	1.33,	1.31,	1.32	
$D_5$ :	1.30,	1.32,	1.30,	1.33		

Test the hypothesis that there is no difference between the mean weights of wires coming from different dyes.

24. It is suspected that four machines used in a canning operation fills cans to different levels on the average. Random samples of cans produced by each machine were taken and the fill (in ounces) was measured. The results are tabulated below:

**Machine**

<i>A</i>	<i>B</i>	<i>C</i>	<i>D</i>
10.20	10.22	10.17	10.15
10.18	10.27	10.22	10.27
10.36	10.26	10.34	10.28
10.21	10.25	10.27	10.40
10.25	—	—	10.30

Do the machines appear to be filling the cans at different average levels?

25. Different numbers of leaves were taken from each of 6 trees and their lengths measured. The following are the lengths in millimetres:

Tree	Lengths						
	82	87	86	90	81	84	
1	82	87	86	90	81	84	
2	85	84	91	92	88		
3	92	90	84	86	88	93	89
4	80	86	87	81	82	82	
5	87	86	88	90	85	86	87
6	90	86	84	85	85	86	87
						84	87

Can all these leaves be regarded as having come from the same species of trees?

26. There are 3 typists working in an office. The times (in minutes) they spend for tea-break in addition to the allowed lunch tea break are observed and noted below:

A :	25	18	30	32	35	37	19
B :	24	22	26	28	30	32	28
C :	28	20	27	19	29	35	30
						23	27
							32

Can the difference in average times that the 3 typists spend for tea break be attributed to chance variation?

27. Four machines A, B, C, D are used to produce a certain kind of cotton fabric. 4 Samples with each unit of size 100 square metres are selected from the outputs of the machines at random and the number of flaws in each 100 square metres are counted, with the following results:

A	B	C	D
8	6	14	20
9	8	12	22
11	10	18	25
12	4	9	23

Do you think that there is a significant difference in the performance of the four machines?

28. The following table shows the yield (in certain units) of lima beans on 20 plots of land subject to 4 different treatments, 5 plots per treatment. Set up an analysis of variance table to test the significance of the differences between the yields due to different treatments.

$T_1$ :	26.3	30.0	54.2	25.7	52.4
$T_2$ :	18.5	21.1	29.3	17.2	12.4
$T_3$ :	36.9	21.8	24.0	18.5	10.2
$T_4$ :	39.8	28.7	21.2	39.4	29.0

29. To test the significance of the variation of the retail prices of a certain commodity in the 4 principal cities Mumbai, Kolkata, Delhi and Chennai, 7 shops were chosen at random in each city and the prices (in Rs.) observed were as follows:

Mumbai : 100, 97, 91, 87, 87, 81, 79  
 Kolkata : 102, 100, 98, 97, 94, 86, 80  
 Delhi : 106, 102, 98, 86, 86, 84, 84  
 Chennai : 97, 95, 94, 92, 90, 86, 82

Do the data indicate that the prices in the 4 cities are significantly different?

30. Steel wire was made by 4 manufacturers  $A, B, C$  and  $D$ . In order to compare their products, 10 samples were randomly drawn from a batch of wires made by each manufacturer and the strength of each piece of wire was measured. The (coded) values are given below:

$A$  : 55, 50, 80, 60, 70, 75, 40, 45, 80, 70  
 $B$  : 70, 80, 85, 105, 65, 100, 90, 95, 100, 70  
 $C$  : 70, 60, 65, 75, 90, 40, 95, 70, 65, 75  
 $D$  : 90, 115, 80, 70, 95, 100, 105, 90, 100, 60

Carry out an analysis of variance and give your conclusions.

31. A randomised block experiment was laid out (with 4 blocks, each block containing 4 plots) to test 4 varieties of manure  $A, B, C, D$  and the yields per acre are given as below. Test for the significance of the difference among the 4 varieties of manure.

Block I	A155	B152	C157	D156
Block II	B152	C150	D156	A154
Block III	C156	D153	A161	B162
Block IV	D153	A154	B156	C155

32. The following table gives the gains in weights of 4 different types of pigs fed on 3 different rations over a period. Test whether  
 (i) the difference in the rations significant  
 (ii) the 4 types of pigs differ significantly in gaining weight

Ration	Types of pig			
	I	II	III	IV
A	13.8	15.7	16.0	20.2
B	8.7	11.8	9.0	12.9
C	12.0	16.5	13.3	12.5

33. Four experiments determine the moisture content of samples of a powder, each observer taking a sample from each of six consignments. The assessments are given below:

Observer	Consignment					
	1	2	3	4	5	6
1	9	10	9	10	11	11
2	12	11	9	11	10	10
3	11	10	10	12	11	10
4	12	13	11	14	12	10

Perform an analysis of variance on these data and discuss whether there is any significant difference between consignments or between observers.

34. In order to compare three burners  $B_1$ ,  $B_2$  and  $B_3$ , one observation is made on each burner on each of four successive days. The data are tabulated below:

	$B_1$	$B_2$	$B_3$
Day 1	21	23	24
Day 2	18	17	23
Day 3	18	21	20
Day 4	17	20	22

Perform an analysis of variance on these data and find whether the difference between (i) the days, and (ii) the burners significant at 5% LOS.

35. A company appoints 4 salesmen  $A$ ,  $B$ ,  $C$  and  $D$  and observes their sales in 3 seasons summer, winter and monsoon. The figures (in lakhs of Rs) are given in the following table:

<i>Season</i>	<i>Salesmen</i>			
	<i>A</i>	<i>B</i>	<i>C</i>	<i>D</i>
Summer	36	36	21	35
Winter	28	29	31	32
Monsoon	26	28	29	29

Carry out an analysis of variance.

36. The following data represent the numbers of units of production per day turned out by 4 different workers using 5 different types of machines:

<i>Worker</i>	<i>Machine type</i>				
	<i>A</i>	<i>B</i>	<i>C</i>	<i>D</i>	<i>E</i>
1	4	5	3	7	6
2	6	8	6	5	4
3	7	6	7	8	8
4	3	5	4	8	2

On the basis of this information, can it be concluded that (i) the mean productivity is the same for different machines (ii) the workers do not differ with regard to productivity?

37. The number of automobiles arriving at 4 toll gates were recorded for a 2 hours time period (10 am to 12 noon) for each of six working days. The data are as follows:

<i>Day</i>	<i>Gate 1</i>	<i>Gate 2</i>	<i>Gate 3</i>	<i>Gate 4</i>
Mon	200	228	212	301
Tues	208	230	215	305
Wed	225	240	228	288
Thur	223	242	224	212
Fri	228	210	235	215
Sat	220	208	245	200

Determine whether the rate of arrival (i) is the same at each toll gate, and (ii) differs significantly during the six days or not.

38. The following table gives the number of refrigerators sold by 4 salesmen in 3 months:

Months	Salesman			
	I	II	III	IV
May	50	40	48	39
June	46	48	50	45
July	39	44	40	39

Determine whether (i) there is any difference in average sales made by the four salesmen, and (ii) the sales differ with respect to different months.

39. Four different drugs have been developed for a certain disease. These drugs are used in 3 different hospitals and the results, given below, show the number of cases of recovery from the disease per 100 people who have taken the drugs:

	$D_1$	$D_2$	$D_3$	$D_4$
$H_1$	19	8	23	8
$H_2$	10	9	12	6
$H_3$	11	13	13	10

What conclusions can you draw based on an analysis of variance ?

40. The following table gives the additional hours of sleep due to 3 soporofic drugs  $A, B, C$  tried on one patient each from 4 different age groups. Examine whether age has got any significant effect on the gain in sleep. Also examine whether the 3 drugs are similar in their effects or not

Drug	Age group			
	30–40	40–50	50–60	60–70
A	2.0	1.2	1.0	0.3
B	1.1	0.8	0.0	- 0.1
C	1.5	1.3	0.9	0.1

41. The following table gives the results of experiments on 4 varieties of a crop in 5 blocks of plots. Prepare the ANOVA table to test the significance of the difference between the yields of the 4 varieties:

Variety	$B_1$	$B_2$	$B_3$	$B_4$	$B_5$
A	32	34	33	35	37
B	34	33	36	37	35
C	31	34	35	32	36
D	29	26	30	28	29

42. In the table given below are the yields of 6 varieties of a crop in a 4 replicate RBD experiment. Analyse the data:

Replicates	Varieties					
	1	2	3	4	5	6
1	18.5	15.7	16.2	14.1	13.0	13.6
2	11.7	14.25	12.9	14.4	14.9	12.5
3	15.4	14.6	15.5	20.3	18.4	21.5
4	16.5	18.6	12.7	15.7	16.5	18.0

43. Analyse the variance in the following Latin square:

A8	C18	B9
C9	B18	A16
B11	A10	C20

44. Analyse the variance in the following Latin square:

20 B	17 C	25 D	34 A
23 A	21 D	15 C	24 B
24 D	26 A	21 B	19 C
26 C	23 B	27 A	22 D

45. A varietal trial was conducted on wheat with 4 varieties *A, B, C, D* in a Latin square design. The plan of the experiment and the per plot yield are given below.

C25	B23	A20	D20
A19	D19	C21	B18
B19	A14	D17	C20
D17	C20	B21	A15

Analyse the data and interpret the result.

46. The following is the Latin square layout of a design when 4 varieties of seeds are tested. Set up the analysis of variance table and state your conclusions.

A105	B95	C125	D115
C115	D125	A105	B105
D115	C95	B105	A115
B95	A135	D95	C115

47. The table given below shows the yield of a certain crop in kg per plot. The letters *A, B, C, D* refer to 4 different manurial treatments. Carry out an analysis of variance.

A260	B300	C335	D371
B280	A300	D300	C410
D320	C345	B340	A254
C372	D395	A290	B328

48. The following results were obtained in a textile experiment to compare the effects of sizing treatments A, B, C, D on the number of warps breaking per hour. Is the difference between the treatments significant?

		Loom			
		1	2	3	4
Period	1	A 54	B 31	C 70	D 45
	2	B 59	A 23	D 100	C 22
	3	C 40	D 41	B 74	A 33
	4	D 83	C 29	A 100	B 28

49. An agricultural experiment on the Latin square plan gave the following results for the yield of wheat per acre, letters corresponding to varieties.

A16	B10	C11	D9	E9
E10	C9	A14	B12	D11
B15	D8	E8	C10	A18
D12	E6	B13	A13	C12
C13	A11	D10	E7	B14

Discuss the variation of yield with each of the factors corresponding to the rows and columns.

50. The following is a Latin square design of five treatments:

A13	B9	C21	D7	E6
D9	E8	A15	B7	C16
B11	C17	D8	E10	A17
E8	A15	B7	C10	D7
C11	D9	E8	A11	B15

Analyse the data and interpret the results.

## **ANSWERS**

### **Exercise 11**

19. 12      20. 576
21.  $Q_1 = 46.08$ ,  $Q_2 = 73.67$ ,  $F_0 = 2.50$ ,  $F_{5\%} = 3.49$ ;  
Difference between varieties not significant.
22.  $Q_1 = 5838.4$ ,  $Q_2 = 1126$ ,  $F_0 = 38.89$ ,  $F_{5\%} = 3.68$ ; Breaking strengths of ropes differ significantly.
23.  $Q_1 = 35.98$ ,  $Q_2 = 99.38$ ,  $F_0 = 1.81$ ,  $F_{5\%} = 2.87$ ; Mean weights of wires do not differ significantly.

24.  $Q_1 = 44.44, Q_2 = 696, F_0 = 2.98, F_{5\%} = 3.35$ ; No, the machines appear to fill at same level.
25.  $Q_1 = 151.95, Q_2 = 255, F_0 = 4.17, F_{5\%} = 2.50$ ; Leaves have not come from the same species.
26.  $Q_1 = 2.52, Q_2 = 29.27, F_0 = 11.62, F_{5\%} = 19.45$ ; Difference may be attributed to chance variation.
27.  $Q_1 = 540.65, Q_2 = 85.75, F_0 = 25.21, F_{5\%} = 3.49$ ; Performances of the machines differ significantly.
28.  $Q_1 = 34845.93, Q_2 = 10032.78, F_0 = 3.47, F_{5\%} = 3.24$ ; Treatments give significantly different yields.
29.  $Q_1 = 94.97, Q_2 = 1446.03, F_0 = 1.9, F_{5\%} = 8.64$ ; Prices do not differ significantly.
30.  $Q_1 = 5151, Q_2 = 8348, F_0 = 7.41, F_{5\%} = 8.60$ ; Strengths of wire do not differ significantly.
31.  $Q_1 = 42.75, Q_2 = 6.75, Q_3 = 96.25, F_0 = 4.75, F_{5\%} = 8.82$ ; Difference between manures is not significant.
32.  $Q_1 = 3393.59, Q_2 = 878.44, Q_3 = 344.36$ ,  
 $F_0$  (rows) = 9.85 and  $F_{5\%} = 5.14$ ,  
 $F_0$  (columns) = 2.55 and  $F_{5\%} = 4.76$ ;  
 Difference between rations significant. Difference between pigs is not significant.
33.  $Q_1 = 13.13, Q_2 = 9.71, Q_3 = 13.12$ ,  
 $F_0$  (rows) = 5.03 and  $F_{5\%} = 3.29$ ,  
 $F_0$  (columns) = 2.23 and  $F_{5\%} = 5.05$ ;  
 Difference between observers is significant, Difference between consignments is not significant.
34.  $Q_1 = 22.00, Q_2 = 28.17, Q_3 = 14.50$ ,  
 $F_0$  (rows) = 3.03 and  $F_{5\%} = 4.76$ ,  
 $F_0$  (columns) = 5.83 and  $F_{5\%} = 5.14\%$ ;  
 Difference between days is not significant; Difference between burners is significant.
35.  $Q_1 = 32, Q_2 = 42, Q_3 = 136$ ,  
 $F_0$  (rows) = 1.42 and  $F_{5\%} = 19.33$ ,  
 $F_0$  (columns) = 1.62, and  $F_{5\%} = 8.94$ ;  
 Differences between seasons and between salesmen are not significant.
36.  $Q_1 = 22.0, Q_2 = 12.8, Q_3 = 30.0$ ;  
 $F_0$  (rows) = 2.93 and  $F_{5\%} = 3.49$ ,  
 $F_0$  (columns) = 1.28 and  $F_{5\%} = 3.26$ ;  
 Differences between the workers and between machine types are not significant.
37.  $Q_1 = 2279.83, Q_2 = 1470.05, Q_3 = 820.12$ ,  
 $F_0$  (rows) = 1.80 and  $F_{5\%} = 4.64$ ,  
 $F_0$  (columns) = 1.79 and  $F_{5\%} = 3.29$ ;  
 Differences between the days and between the gates are not significant.

38.  $Q_1 = 109.5, Q_2 = 42.0, Q_3 = 64.5,$   
 $F_0$  (rows) = 5.09 and  $F_{5\%} = 5.14,$   
 $F_0$  (columns) = 1.30 and  $F_{5\%} = 4.76;$   
Differences between the months and between salesmen are not significant.
39.  $Q_1 = 55.17, Q_2 = 113.0, Q_3 = 89.5,$   
 $F_0$  (rows) = 1.85 and  $F_{5\%} = 5.14;$   
 $F_0$  (columns) = 2.52 and  $F_{5\%} = 4.76;$   
Differences between the hospitals and between the drugs are not significant.
40.  $Q_1 = 98.17, Q_2 = 341.58, Q_3 = 25.17,$   
 $F_0$  (rows) = 11.69 and  $F_{5\%} = 5.14,$   
 $F_0$  (columns) = 27.11 and  $F_{5\%} = 4.76;$   
Age has significant effect on the gain in sleep; Drugs differ significantly in their effect.
41.  $Q_1 = 134.0, Q_2 = 21.7, Q_3 = 29.5,$   
 $F_0$  (row) = 18.16 and  $F_{5\%} = 3.49,$   
 $F_0$  (columns) = 2.21 and  $F_{5\%} = 3.26;$   
Difference between the yields of 4 varieties is significant.
42.  $Q_1 = 56.76, Q_2 = 12.58, Q_3 = 80.18,$   
 $F_0$  (rows) = 3.30 and  $F_{5\%} = 3.34,$   
 $F_0$  (columns) = 2.27 and  $F_{5\%} = 4.65;$   
Difference between the varieties is significant.
43.  $Q_1 = 11.56, Q_2 = 68.23, Q_3 = 29.56, Q_4 = 68.21,$   
 $F_0$  (rows) = 5.90,  $F_0$  (cols.) = 1,  $F_0$  (letters) = 2.31,  
 $F_{5\%}$  (for all) = 19.0;  
The differences between rows, between columns and between letters are not significant.
44.  $Q_1 = 34.19, Q_2 = 22.69, Q_3 = 141.19, Q_4 = 96.87, F_0$  (rows) = 1.42 and  $F_{5\%} = 8.94;$   $F_0$  (columns) = 2.14 and  $F_{5\%} = 8.94;$   $F_0$  (letters) = 2.91 and  $F_{5\%} = 4.76;$  Differences between rows, between columns and between letters are not significant.
45.  $Q_1 = 46.5, Q_2 = 7.5, Q_3 = 48.5, Q_4 = 10.5,$   
 $F_0$  (rows) = 8.86,  $F_0$  (columns) = 1.43,  $F_0$  (letters) = 9.24,  $F_{5\%} = 4.76.$   
Difference between varieties is significant.
46.  $Q_1 = 2, Q_2 = 4, Q_3 = 22, Q_4 = 60, F_0$  (rows) = 15,  $F_0$  (columns) = 7.5,  
 $F_0$  (letters) = 1.36,  $F_{5\%} = 8.94;$  Difference between rows is significant, but differences between columns and between letters are not significant.
47.  $Q_1 = 2540.5, Q_2 = 2853.75, Q_3 = 18690, Q_4 = 7515.75, F_0$  (rows) = 1.48, and  $F_{5\%} = 8.94;$   $F_0$  (columns) = 1.32 and  $F_{5\%} = 8.94;$   $F_0$  (letters) = 4.97 and  $F_{5\%} = 4.76;$  Differences between rows and between columns are not significant, but difference between treatments is significant.

48.  $Q_1 = 376, Q_2 = 8184, Q_3 = 1547.5, Q_4 = 284.5$   
 $F_0$  (rows) = 2.64,  $F_0$  (columns) = 57.53,  $F_0$  (letters) = 10.88;  $F_{5\%}$  (for all) = 4.76; Difference between periods is not significant; Differences between looms and between treatments are significant.
49.  $Q_1 = 2.16, Q_2 = 66.56, Q_3 = 122.56, Q_4 = 5.28, F_0$  (rows) = 1.2,  $F_0$  (columns) = 37.8,  $F_0$  (letters) = 69.6,  $F_{5\%}$  (for all) = 3.26; Difference between rows is not significant, but differences between columns and between varieties are significant.
50.  $Q_1 = 26, Q_2 = 34, Q_3 = 224.4, Q_4 = 103.6, F_0$  (rows) = 1.33 and  $F_{5\%} = 5.91, F_0$  (columns) = 1.02 and  $F_{5\%} = 5.91, F_0$  (letters) = 6.50 and  $F_{5\%} = 3.26$ ; Differences between rows and between columns are not significant, but difference between treatments is significant.



# *Chapter* 12

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# Statistical Quality Control

## **Introduction**

In these days of tough business competition, it has become essential to maintain the quality of the goods manufactured and market them at reasonable price. If the consumers feel satisfied with regard to the quality, price, etc. of the product manufactured by a certain company, it will result in goodwill for the product and in increase in sales. If not and if proper attention is not given to the complaints of the consumers regarding quality, the manufacturer cannot push through his product in the market and ultimately he has to quit the market. Hence it is important to maintain and improve the quality of the manufactured products for the manufacturer to remain and flourish in his business.

*Quality* does not mean the highest standard of the manufactured product, but conforming to the prescribed standard of the product so as to satisfy the consumers, even if it may be below the highest absolute standard.

Though the quality standard might have been specified, it is not possible to avoid some variation in the quality of the product. For example, a machine is set and hence expected to produce per day 10,000 bolts, each of length 2 cm. It is very unlikely that all the bolts are of length 2 cm exactly. Some of them may be slightly less than 2 cm and some slightly more than 2 cm in length. Such variation in quality of the product can be divided into two kinds, namely, *chance variation* and *assignable variation*.

*Chance variation* is the variation in the quality of the product which occurs due to many minor, but random causes, such as slight changes in temperature, pressure and metal hardness. *Assignable variation* is the variation that occurs due to non-random causes like poor quality of input raw material, mechanical faults, handling of machines by inexperienced operators, etc. Though no method is available by which the chance variation can be controlled or eliminated, assignable variation can be eliminated, if detected early during the production process.

*Statistical Quality Control (SQC)* is a statistical method for finding whether the variation in the quality of the product is due to random causes or assignable causes. SQC does not involve inspecting each and every item produced for quality standards, but involves inspection of samples of items produced and application of tests of significance.

Statistical Quality Control methods are applied to two distinct phases of manufacturing operation. Process control and Acceptance sampling or product control.

*Process Control* means control of the quality of the goods while they are in the process of production. To achieve process control, repeated random samples are taken from the population of items, as and when they are being produced, the sample results are subjected to statistical analysis by means of simple graphical device, known as control charts and the faults in the production process are rectified then and there.

*Control Chart* is a graphical device mainly used for the study and control of the manufacturing process. It is simple to construct and easy to interpret. The manufacturer can find out, at a glance, whether or not the process is under control so that the proportion of defective items is not excessive. Control chart is also called *shewhart chart*.

There are two types of control charts, namely,

1. Control charts of variables (Mean and range charts).
2. Control charts of attributes (p-chart and c-chart).

*Variables* are the quality characteristics of a product that are measurable e.g., diameter of a hole bored by a drilling machine. For the construction of control charts of a variable, a record of the actual measurements of that variable for the sampled items must be known.

*Attributes* are the quality characteristics of a product that are not amenable for measurement, but are identified by their presence or absence, e.g. the presence (and hence the number) of defective items in a sample.

### ***Control Charts of Variables***

Normally, a production process is expected to turn out quality products or to be under control. So the production process is allowed to operate and to produce items. To ensure that the process continues to be in a state of control, a sample of the produced items is drawn periodically and tested for their quality. If the quality of the sampled items is satisfactory in the statistical sense, the production process is allowed to continue. Otherwise, corrective measures are taken to restore the quality of the items.

Instead of inspecting each item in a sample for its quality, it will be advantageous and prudent to arrive at a conclusion regarding the quality of the items on the basis of the sample mean. Hence, the control chart for the sample mean is constructed and used for taking decisions regarding the quality of the items. The sample mean alone may not reflect the quality variations of the items,

but the range of the sampled values is known to be a simple measure of the quality variations. Hence the control chart for the sample range is also constructed and used for taking decisions regarding the quality of the items.

### **Control Limits for the Sample Mean $\bar{X}$ and Sample Range $R$**

Let  $X$  be the random variable that represents the measurable quality characteristic of the population of the items produced. We assume that  $X$  follows a normal distribution with mean  $\mu$  and standard deviation  $\sigma$ .

Then  $\bar{X}$ , the mean of a sample of size  $n$ , is also a random variable that is normally distributed with mean  $\mu$  and S.D.  $\frac{\sigma}{\sqrt{n}}$ .

We know, from the property of normal distribution that  $P\left\{\mu - \frac{3\sigma}{\sqrt{n}} \leq \bar{X} \leq \mu + \frac{3\sigma}{\sqrt{n}}\right\} = 0.9973$ . This means that 99.73% of the sample means (viz., almost all the sample mean values) will be within  $\mu - \frac{3\sigma}{\sqrt{n}}$  and  $\mu + \frac{3\sigma}{\sqrt{n}}$ . The variation in the sample mean values within these limits is due to random causes. If an observed value of  $\bar{X}$  lies outside  $\mu \pm \frac{3\sigma}{\sqrt{n}}$ , it indicates the presence of some assignable cause.  $\mu \pm \frac{3\sigma}{\sqrt{n}}$  are called the *control limits* for the sample mean  $\bar{X}$ ;  $\mu - \frac{3\sigma}{\sqrt{n}}$  is called the *Lower Control Limit (LCL)* and  $\mu + \frac{3\sigma}{\sqrt{n}}$  is called the *Upper Control Limit (UCL)*.

Similarly  $R$ , the sample range follows a normal distribution with mean  $\bar{R}$  and S.D.  $\sigma_R$ . Since  $P\{\bar{R} - 3\sigma_R \leq R \leq \bar{R} + 3\sigma_R\} = 0.9973$ , the lower and upper control limits for the sample range  $R$  are  $\bar{R} \pm 3\sigma_R$ , where  $\bar{R} = \frac{1}{N} \sum R_i$ , where  $N$  is the number of samples each of size  $n$ .

Now,  $\mu$  and  $\sigma$  are not known and computation of  $\sigma_R$  involves some numerical work. Hence, they are estimated approximately by using  $\mu = \bar{\bar{X}} = \frac{1}{N} \sum \bar{X}_i$ ,  $\frac{3\sigma}{\sqrt{n}} = A_2 \bar{R}$ ,  $\sigma_R$  (in LCL) =  $\frac{1}{3}(1 - D_3) \bar{R}$  and  $\sigma_R$  (in UCL) =  $\frac{1}{3}(D_4 - 1) \bar{R}$ , where  $A_2$ ,  $D_3$  and  $D_4$  are *control chart constants*, whose values depend on the sample size  $n$  and which are readily available in the table of control chart constants (Table 12.1), given at the end of the chapter. Using these approximations, the control limits are obtained as follows:

$$\text{LCL for } \bar{X} = \bar{\bar{X}} - A_2 \bar{R} \text{ and UCL for } \bar{X} = \bar{\bar{X}} + A_2 \bar{R}$$

$$\text{LCL for } R = D_3 \bar{R} \text{ and UCL for } R = D_4 \bar{R}$$

### Procedure to Draw the $\bar{X}$ -chart and $R$ -chart

1. The sample values in each of the  $N$  samples each of size  $n$  will be given. Let  $\bar{X}_1, \bar{X}_2, \dots, \bar{X}_N$  be the means of the  $N$  samples and  $R_1, R_2, \dots, R_N$  be the ranges of the  $N$  samples. By range of a sample, we mean the maximum sample value minus the minimum sample value in that sample.
2. We then compute  $\bar{\bar{x}} = \frac{1}{N}(\bar{x}_1 + \bar{x}_2 + \dots + \bar{x}_N)$  and  $\bar{R} = \frac{1}{N}(R_1 + R_2 + \dots + R_N)$ .
3. The values of  $A_2, D_3, D_4$  for the given sample size  $n$  are taken from the table of control chart constants.
4. Then the values of the control limits  $\bar{X} \pm A_2 \bar{R}$  (for the mean chart) and the control limits  $D_3 \bar{R}$  and  $D_4 \bar{R}$  (for the range chart) are computed.
5. On the ordinary graph sheet, the sample numbers are represented on the  $x$ -axis and the sample means on the  $y$ -axis (for the mean chart) and the sample ranges on the  $y$ -axis (for the range chart).
6. For drawing the mean chart, we draw the three lines  $y = \bar{X}$ ,  $y = \bar{X} - A_2 \bar{R}$  and  $y = \bar{X} + A_2 \bar{R}$  which represent respectively the central line, the LCL line and UCL line. Also we plot the points whose coordinates are  $(1, \bar{X}_1), (2, \bar{X}_2), \dots, (N, \bar{X}_N)$  and join adjacent points by line segments. The graph thus obtained is the  $\bar{X}$ -chart.
7. For drawing the range chart, we draw the three line  $y = \bar{R}$ ,  $y = D_3 \bar{R}$  and  $y = D_4 \bar{R}$  which represent respectively the central line, the LCL line and UCL line. Also we plot the points whose co-ordinates are  $(1, R_1), (2, R_2), \dots, (N, R_N)$  and join adjacent points by line segments. The graph thus obtained is the  $R$ -chart.

### Comments on State of Control of the Process

If the plotted points fall within the LCL and UCL lines, there is nothing to worry, as in such a case the variation between the samples is attributed to chance causes and the process is under control.

But when one or more plotted points lie outside the control lines, it is to be considered as a danger signal, indicating that the variations between samples are caused by assignable causes and the process is out of control and that necessary corrective action should be taken at once.

Sometimes even though the plotted points may lie between the control lines, a sizeable number of successive points may show a tendency to lie on definite curves going towards the LCL or UCL lines or may lie on the same side of the central line. Such a pattern of sample points should also be considered as a danger signal, warranting a change in the production process.

### **Control Chart for Sample Standard Deviation or *s*-chart**

Since the standard deviation is an ideal measure of dispersion, a combination of control charts for the sample mean  $\bar{X}$  and the sample S.D.  $s$  is more appropriate than the combination of  $\bar{X}$  and  $R$  charts for controlling process average and process variability. We know that  $s$ , the S.D. of a sample of size  $n$ , is a random variable that is normally distributed with mean  $\sigma$  and S.D.  $\frac{s}{\sqrt{2n}}$ , where  $\sigma$  is the S.D. of the population from which the sample is drawn.

$$\text{Hence, } P \left\{ \sigma - \frac{3\sigma}{\sqrt{2n}} \leq s \leq \sigma + \frac{3\sigma}{\sqrt{2n}} \right\} = 0.9973$$

$$\therefore \text{the lower and upper control limits for } s \text{ are } \sigma - \frac{3\sigma}{\sqrt{2n}} \text{ and } \sigma + \frac{3\sigma}{\sqrt{2n}}.$$

Since  $\sigma$  is not known, it is estimated approximately by  $\bar{s} = \frac{1}{N} (s_1 + s_2 + \dots + s_N)$ , where  $s_i$  is the S.D. of the  $i$ th sample and  $N$  is the number of samples considered. Hence, LCL for  $s = \left(1 - \frac{3}{\sqrt{2n}}\right) \bar{s} \approx B_3 \bar{s}$  and UCL for  $s = \left(1 + \frac{3\sigma}{\sqrt{2n}}\right) \bar{s} \approx B_4 \bar{s}$ .

The values of  $B_3$  and  $B_4$  can be read for various values of sample size  $n$  from the table of control chart constants.

The procedure for drawing the  $s$ -chart is similar to that for  $\bar{x}$ -chart and  $R$ -chart.

If  $\bar{x}$  values and  $s$  values only are given, then CL for  $\bar{x} = \bar{\bar{x}}$ , LCL for  $\bar{x} = \bar{\bar{x}} - A_1 \sqrt{\frac{n-1}{n} \bar{s}}$  and UCL for  $\bar{x} = \bar{\bar{x}} + A_1 \sqrt{\frac{n-1}{n} \bar{s}}$ , when  $n \leq 25$ .

#### **Note**

*The difficulty of computation of  $s$  makes the use of  $s$ -chart almost impractical in most industrial situations. R-chart is preferred to  $s$ -chart because of its computational ease.*

## **Control Charts for Attributes**

To control the quality of certain products whose attributes are available, the following control charts are used:

- (i)  $p$ -chart for proportion of defectives
- (ii)  $np$ -chart for number of defectives
- (iii)  $c$ -chart for the number of defects in a unit

Of the above,  $p$ -chart is used when all the samples drawn from the produced items are of the same size or of different size;  $np$ -chart is used only when all the samples are of the same size  $n$ ;  $c$ -chart is used only when each sample consists of only one item.

### **np-chart**

If the proportion of defectives (successes) in the population of items produced is  $p$  and the number of defectives in a sample of size  $n$  (viz.,  $n$  trials) is  $X$ , then  $X$  follows a binomial distribution with mean  $np$  and S.D.  $\sqrt{npq}$ . When  $n$  is sufficiently large and when neither  $p$  nor  $q$  is very small,  $X$  follows a normal distribution with mean  $np$  and S.D.  $\sqrt{npq}$ .

Hence,  $P\{np - 3\sqrt{npq} \leq X \leq np + 3\sqrt{npq}\} = 0.9973$ .

$\therefore$  the control limits for  $X$ , the number of defectives, are  $np \mp 3\sqrt{npq}$ .

As the population proportion of defectives  $p$  will not be known, it is estimated approximately as  $\bar{p} = \frac{1}{N} (p_1 + p_2 + \dots + p_N)$  or equivalently  $n\bar{p} = \frac{1}{N} (np_1 + np_2 + \dots + np_N)$ , where all the samples are of the same size  $n$ , the number of samples is  $N$  and the number of defectives in the  $i$ th sample is  $np_i$ .

Hence, the control limits for the number of sample defectives are  $n\bar{p} \pm 3\sqrt{n\bar{p}(1 - \bar{p})}$ , where each sample is of size  $n$ . Here, the target value or the central (line) value of  $np$  is  $n\bar{p}$ . To draw the  $np$ -chart, the sample number is represented on the  $x$ -axis and the number of defectives is represented on the  $y$ -axis.

The lines  $y = n\bar{p}$ ,  $y = n\bar{p} - 3\sqrt{n\bar{p}(1 - \bar{p})}$  and  $y = n\bar{p} + 3\sqrt{n\bar{p}(1 - \bar{p})}$ , which represent respectively the central line, LCL line and UCL line are drawn. We plot the points  $(1, np_1)$ ,  $(2, np_2)$ , ...,  $(N, np_N)$  on the graph sheet and adjacent points are joined by line segments. The state of control of the process is decided as before.

### **p-chart**

Under the same assumptions as those of  $np$ -chart, since  $X$  follows a normal distribution with mean  $np$  and S.D.  $\sqrt{npq}$ , the proportion of defectives, viz;  $\frac{X}{n}$  follows a normal distribution with mean  $p$  and S.D.  $\sqrt{\frac{pq}{n}}$ .

Hence,  $P\left\{p - 3\sqrt{\frac{pq}{n}} \leq \frac{X}{n} \leq p + 3\sqrt{\frac{pq}{n}}\right\} = 0.9973$ .

$\therefore$  the control limits for  $\frac{X}{n}$ , the proportion of defectives, are  $p \mp 3\sqrt{\frac{pq}{n}}$ .

As in the previous case,  $p$  is estimated as  $\bar{p} = \frac{1}{N} (p_1 + p_2 + \dots + p_N)$ , where the proportion of defectives in the  $i$ th sample is  $p_i$ .

Hence, the control limits for the fraction or proportion of sample defective are

$$\bar{p} \mp 3\sqrt{\frac{\bar{p}(1 - \bar{p})}{n}} \quad (1)$$

The formula (1) holds good when all the samples are of the same size  $n$ .

If the size of the sample differs from sample to sample, then the formula for the control limits will be  $\bar{p} \pm 3\sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}}$ , where  $\bar{n}$  is the average sample size given by  $\bar{n} = \frac{1}{N} (n_1 + n_2 + \dots + n_N)$ , where  $n_i$  is the size of the  $i$ th sample. This holds good, when  $n_i$  values do not differ very much from  $\bar{n}$ . This method is applied if  $0.75 \bar{n} < n_i < 1.25 \bar{n}$  for all  $i$ .

To draw the  $p$ -chart, the sample number is represented on the  $x$ -axis and the proportion of defectives is represented on the  $y$ -axis. The lines

$$y = \bar{p}, y = \bar{p} - 3\sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}} \quad \text{or} \quad y = \bar{p} + 3\sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}}$$

and  $y = \bar{p} + 3\sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}} \quad \text{or} \quad y = \bar{p} + 3\sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}}$

which represent respectively the central line, LCL line and UCL line are drawn. We plot the points  $(1, p_1), (2, p_2), \dots, (N, p_N)$  on the graph paper and adjacent points are joined by line segments. The state of control of the process is decided as before.

### Note

$np$ -chart and  $p$ -chart are used when  $\bar{p} \geq 0.05$  or  $n\bar{p} \geq 4$ .

### c-chart

When it is required to control (minimise) the number of defects per unit,  $c$ -chart is used. ' $c$ ' represents the number of defects in a unit. For construction of  $c$ -chart, a record of the number of defects in each of the  $N$  articles inspected should be known. Since the probability of occurrence of a defect in a unit is very small, the number  $X$  of defects in a unit follows a Poisson distribution with parameter  $\lambda$ , viz., with mean  $\lambda$  and S.D.  $\sqrt{\lambda}$ . In the limit,  $X$  follows a normal distribution with mean  $\lambda$  and S.D.  $\sqrt{\lambda}$ .

Hence,  $P \{ \lambda - 3\sqrt{\lambda} \leq X \leq \lambda + 3\sqrt{\lambda} \} = 0.9973$

$\therefore$  the control limits for  $X$ , the number of defects in a unit are  $\lambda \pm 3\sqrt{\lambda}$ .

As the value of  $\lambda$  will not be known, it is estimated approximately by  $\bar{c}$ , where  $\bar{c} = \frac{1}{N} (c_1 + c_2 + \dots + c_N)$ , where  $c_i$  is the number of defects in the  $i$ th unit.

Hence, the control limits for the number of defects  $c$  in a unit are  $\bar{c} \pm 3\sqrt{\bar{c}}$ .

To draw the  $c$ -chart, the item number is represented on the  $x$ -axis and the number of defects in a unit is represented on the  $y$ -axis. The lines  $y = \bar{c}$ ,  $y = \bar{c} - 3\sqrt{\bar{c}}$  and  $y = \bar{c} + 3\sqrt{\bar{c}}$ , which represent respectively the central line, LCL line and UCL line are drawn. We plot the points  $(1, c_1), (2, c_2), \dots, (N, c_N)$  on

the graph sheet and the adjacent points are joined by line segments. The state of control of the process is decided as before.

**Note** *c-chart is used when  $\bar{c} \geq 4$  or when  $\bar{c}$  is small compared with the maximum number of defects given in the data.*

### Specification Limits and Tolerance Limits

Though it is desirable to specify a single target value as the acceptable quality standard of a manufactured product, it is not possible to avoid some variation from the target value due to chance and assignable causes. Hence, the quality of the product is taken as acceptable, if the measurement of the quality characteristic lies within an interval that encloses the target value. The end values of such an interval are called *specification limits*. For example, when the target length of a bolt is 2.5 cm, the specification limits may be assumed as 2.4 cm and 2.6 cm, viz.,  $2.5 \pm 0.1$  cm.

*Tolerance Limits* of a quality characteristic are defined as those values between which nearly all the manufactured items will lie.

If the measurable quality characteristics  $X$  is assumed to be normally distributed with mean  $\mu$  and S.D.  $\sigma$ , then the tolerance limits are usually taken as  $\mu \pm 3\sigma$ , since only 0.27% of all the items produced can be expected to fall outside these limits.

As  $\mu$  and  $\sigma$  will not be known, we get the tolerance limits approximately using the control charts for  $\bar{X}$  and  $R$  as explained below:  $N$  samples, each of size  $n$ , are taken from the population of items produced. Let  $\bar{X}_1, \bar{X}_2, \dots, \bar{X}_N$  be the means of these samples and  $R_1, R_2, \dots, R_N$  be the ranges of these samples. The  $\bar{X}$ -chart and  $R$ -chart are constructed using these values. If the variations of the sample mean and range values are due to chance causes only, viz., if the process is under control with respect to both  $\bar{X}$  and  $R$ , then the tolerance limits are computed as

$\bar{\bar{X}} \pm 3 \frac{\bar{R}}{d_2}$ , since the estimates of the mean and S.D. of the population are given

by  $\hat{\mu} = \bar{\bar{X}}$  and  $\hat{\sigma} = \frac{\bar{R}}{d_2}$ , where  $d_2$  is a control chart constant to be read from the table of control chart constants.

If the process is not under control with respect to  $\bar{X}$  or  $R$  or both, then the samples whose means or ranges go out of control are removed and a new set of  $\bar{\bar{X}}$  and  $\bar{R}$  values are computed using the remaining samples. Using these values, a new set of control limits are computed and the control of the process is checked. This procedure is repeated until the process comes under control. After ascertaining that the process is under control with respect to both the sample mean and range, the tolerance limits are computed as  $\hat{\mu} = \bar{\bar{X}} \pm \frac{3\bar{R}}{d_2}$ , where  $\bar{\bar{X}}$  and  $\bar{R}$  are computed using the samples that remain under control ultimately.

If these tolerance limits are within the specification limits, then the process is assumed to operate at an acceptable level. If they do not fall within the specification limits, the process is bound to produce some defective (unacceptable) items, even though the process may be under control.

### Worked Example 12

#### Example 1

Given below are the values of sample mean  $\bar{X}$  and sample range  $R$  for 10 samples, each of size 5. Draw the appropriate mean and range charts and comment on the state of control of the process.

Sample No. :	1	2	3	4	5	6	7	8	9	10
Mean :	43	49	37	44	45	37	51	46	43	47
Range :	5	6	5	7	7	4	8	6	4	6

$$\begin{aligned}\bar{\bar{X}} &= \frac{1}{N} \sum \bar{x}_i \\ &= \frac{1}{10} (43 + 49 + 37 + \dots + 47) \\ &= 44.2\end{aligned}$$

$$\begin{aligned}\bar{R} &= \frac{1}{N} \sum R_i \\ &= \frac{1}{10} (5 + 6 + 5 + \dots + 6) \\ &= 5.8\end{aligned}$$

From the table of control chart constants, for sample size  $n = 5$ , we have

$$A_2 = 0.577, D_3 = 0 \text{ and } D_4 = 2.115$$

Control limit for  $\bar{X}$  = chart:

$$\text{CL (central line)} = \bar{\bar{X}} = 44.2$$

$$\text{LCL} = \bar{\bar{X}} - A_2 \bar{R} = 44.2 - 0.577 \times 5.8 = 40.85$$

$$\text{UCL} = \bar{\bar{X}} + A_2 \bar{R} = 44.2 + 0.577 \times 5.8 = 47.55$$

#### Control Limit for R-chart

$$\text{CL} = \bar{R} = 5.8; \text{LCL} = D_3 \bar{R} = 0; \text{UCL} = D_4 \bar{R} = 2.115 \times 5.8 = 12.27$$

The mean chart and range chart relevant to this problem are given in Fig. 12.1 (a) and (b):

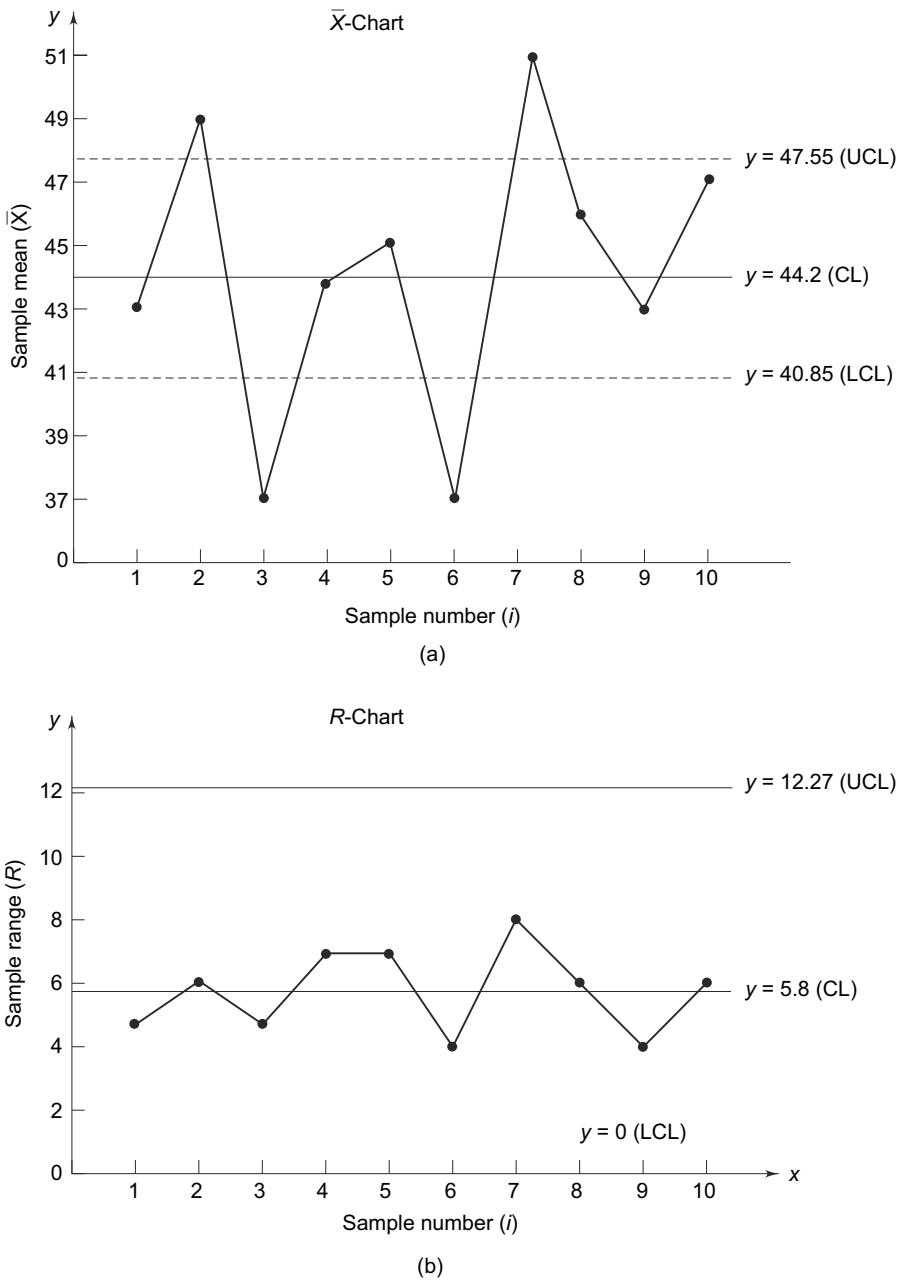


Fig. 12.1

### State of Control

All the sample points in the range chart lie within the control lines. Hence, as far as the variability of the sample values is concerned, the process is under control.

But in the mean chart, two points lie above the upper control line and two points lie below the lower control line. Hence, as far as the average of the sample values is concerned, the process is not under control. On the whole, we conclude that the process is out of control.

**Note**

*Even though we could have arrived at this conclusion regarding the state of control without drawing the control charts, it is necessary to draw the control charts, as it is a part of the solution to the given problem.*

**Example 2**

A machine fills boxes with dry cereal. 15 samples of 4 boxes are drawn randomly. The weights of the sampled boxes are shown as follows. Draw the control charts for the sample mean and sample range and determine whether the process is in a state of control.

Sample Number	1	2	3	4	5	6	7	8
Weights of boxes (X)	10.0	10.3	11.5	11.0	11.3	10.7	11.3	12.3
	10.2	10.9	10.7	11.1	11.6	11.4	11.4	12.1
	11.3	10.7	11.4	10.7	11.9	10.7	11.1	12.7
	12.4	11.7	12.4	11.4	12.1	11.0	10.3	10.7

9	10	11	12	13	14	15
11.0	11.3	12.5	11.9	12.1	11.9	10.6
13.1	12.1	11.9	12.1	11.1	12.1	11.9
13.1	10.7	11.8	11.6	12.1	13.1	11.7
12.4	11.5	11.3	11.4	11.7	12.0	12.1

As the  $\bar{X}$ -chart and  $R$ -chart are to be drawn, we first compute the means and ranges of the given samples.

Sample No. (i)	1	2	3	4	5	6	7	8	9	10
$\Sigma X$	43.9	43.6	46.0	44.2	46.9	43.8	44.1	47.8	49.6	45.6
$\bar{X}_i$	11.0	10.9	11.5	11.1	11.7	11.0	11.0	12.0	12.4	11.4
$R_i$	2.4	1.4	1.7	0.7	0.8	0.7	1.1	2.0	2.1	1.4

11	12	13	14	15
47.5	47.0	47.0	49.1	46.3
11.9	11.8	11.8	12.3	11.6
1.2	0.7	1.0	1.2	1.5

$$\text{Now, } \bar{\bar{X}} = \frac{1}{N} \sum \bar{X}_i = \frac{1}{15} (11.0 + 10.9 + 11.5 + \dots + 11.6)$$

$$= \frac{173.4}{15} = 11.56$$

$$\bar{R} = \frac{1}{N} \sum R_i = \frac{1}{15} (2.4 + 1.4 + 1.7 + \dots + 1.5)$$

$$= \frac{19.9}{15} = 1.33$$

From the table of control chart constants, for the sample size  $n = 4$ , we have

$$A_2 = 0.729, D_3 = 0 \text{ and } D_4 = 2.282$$

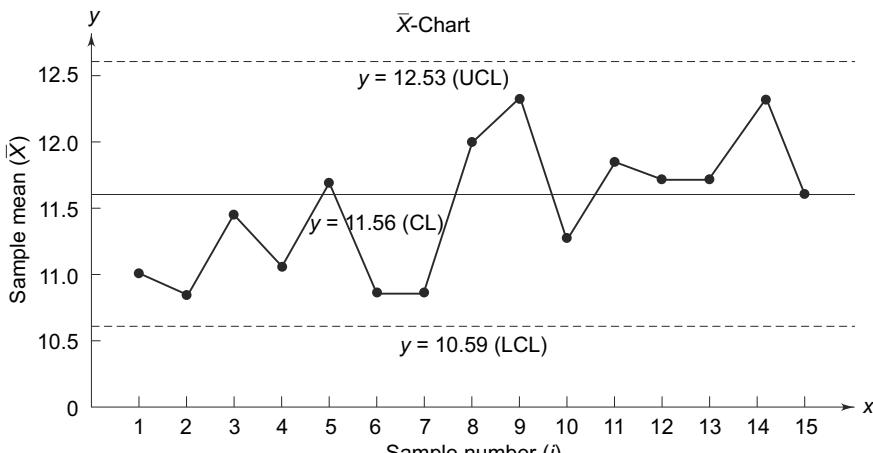
### Control Limits for $\bar{X}$ -chart

$$\begin{aligned} \text{CL} &= \bar{\bar{X}} = 11.56; \text{LCL} = \bar{\bar{X}} - A_2 \bar{R} = 11.56 - 0.729 \times 1.33 \\ &= 10.59 \end{aligned}$$

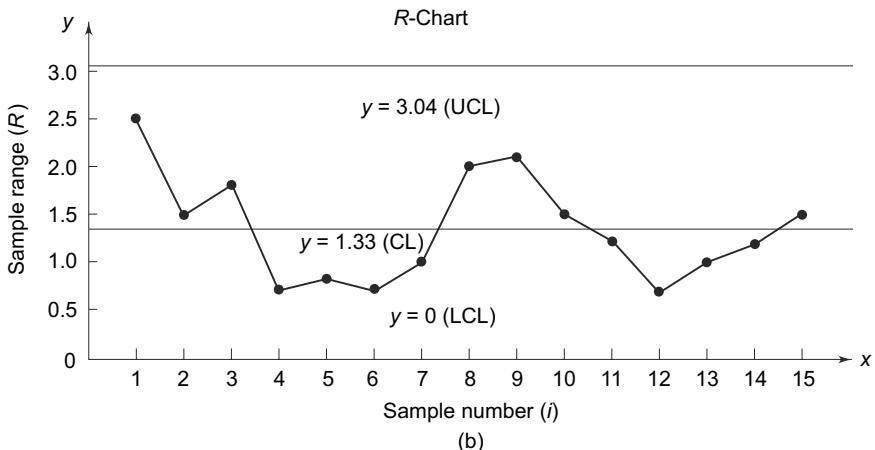
$$\text{UCL} = \bar{\bar{X}} + A_2 \bar{R} = 11.56 + 0.729 \times 1.33 = 12.53$$

### Control Limits for $R$ -chart

$$\begin{aligned} \text{CL} &= \bar{R} = 1.33; \text{LCL} = D_3 \bar{R} = 0 \text{ and } \text{UCL} = D_4 \bar{R} = 2.282 \times 1.33 \\ &= 3.04 \end{aligned}$$



(a)

**Fig. 12.2****State of Control**

Since all the sample points lie within upper and lower control lines both in the  $\bar{X}$ -chart and in the  $R$ -chart, the process is under control.

**Example 3**

The values of sample mean  $\bar{X}$  and sample standard deviation  $s$  for 15 samples, each of size 4, drawn from a production process are given below. Draw the appropriate control charts for the process average and process variability. Comment on the state of control.

Sample No.	1	2	3	4	5	6	7	8	9	10
Mean	15.0	10.0	12.5	13.0	12.5	13.0	13.5	11.5	13.5	13.0
S.D.	3.1	2.4	3.6	2.3	5.2	5.4	6.2	4.3	3.4	4.1

11	12	13	14	15
14.5	9.5	12.0	10.5	11.5
3.9	5.1	4.7	3.3	3.3

$$\bar{\bar{X}} = \frac{1}{N} \sum \bar{X}_i = \frac{1}{15} \times 185.5 = 12.36$$

$$\bar{s} = \frac{1}{N} \sum s_i = \frac{1}{15} \times 60.3 = 4.02$$

From the table of control chart constants, for sample size  $n = 4$ , we have

$$A_1 = 1.880, B_3 = 0 \text{ and } B_4 = 2.266$$

### Control Limits for $\bar{X}$ -chart

$$CL = \bar{\bar{X}} = 12.36$$

$$LCL = \bar{\bar{X}} - A_1 \sqrt{\frac{n}{n-1} \bar{s}} = 12.36 - 1.880 \sqrt{\frac{4}{3}} \times 4.02 = 3.63$$

$$UCL = \bar{\bar{X}} + A_1 \sqrt{\frac{n}{n-1} \bar{s}} = 12.36 + 1.880 \sqrt{\frac{4}{3}} \times 4.02 = 21.09$$

### Control Limits for $s$ -chart

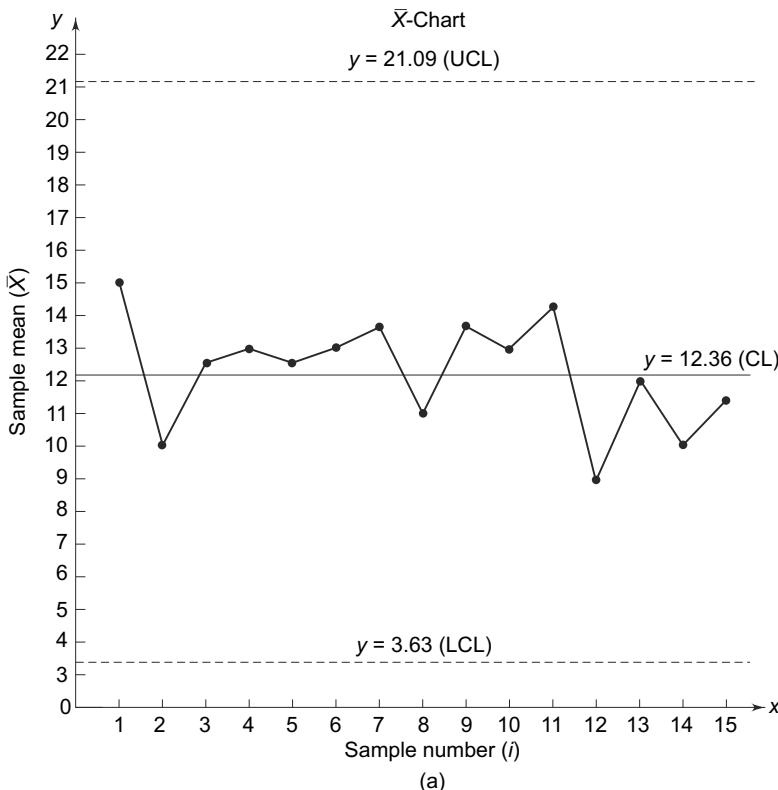
$$CL = \bar{s} = 4.02; LCL = B_3 \bar{s} = 0;$$

$$UCL = B_4 \bar{s} = 2.266 \times 4.02 = 9.11$$

The mean chart and S.D. chart relevant to this problem are given in Fig. 12.3.

### State of Control

Even before drawing the control charts, we observe that the given sample mean values lie between 3.63 and 21.09 and that the given S.D. values fall within 0 and 9.11. Hence, the process is under control with respect to average and variability.



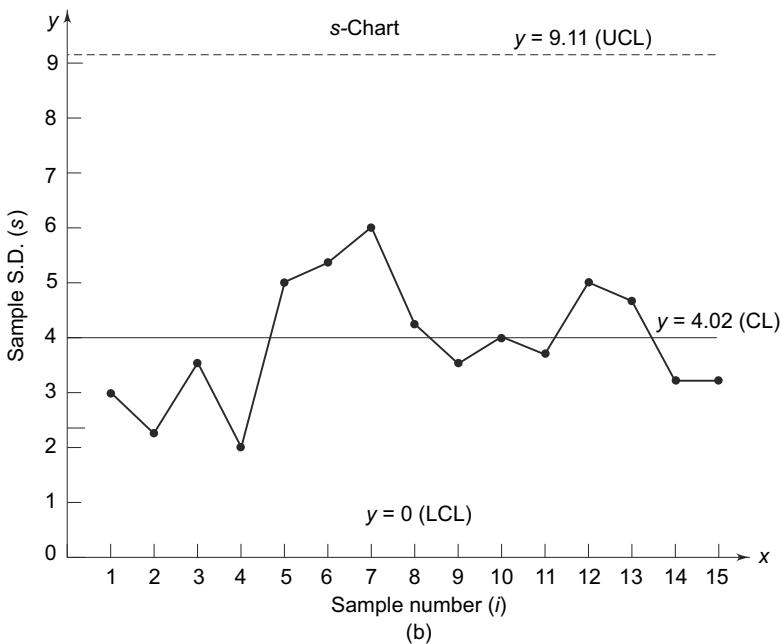


Fig. 12.3

**Example 4**

The following data give the coded measurements of 10 samples each of size 5, drawn from a production process at intervals of 1 hour. Calculate the sample means and S.D.'s and draw the control charts for  $\bar{X}$  and  $s$ .

Sample Number	1	2	3	4	5	6	7	8	9	10
Coded measurements ( $X$ )	9	10	10	8	7	12	9	15	10	16
	15	11	13	13	9	15	9	15	13	14
	14	13	8	11	10	7	9	10	14	12
	9	6	12	10	4	16	13	13	7	14
	13	10	7	13	5	10	5	17	11	14

We first compute mean and S.D. for each sample.

Sample Number	1	2	3	4	5	6	7	8	9	10
$\sum X$	60	50	50	55	35	60	45	70	55	70
$\bar{X}$	12	10	10	11	7	12	9	14	11	14
$\sum (X - \bar{X})^2$	32	26	26	18	26	54	32	28	30	8
$s$	2.5	2.3	2.3	1.9	2.3	3.3	2.5	2.4	2.4	1.3

$$\bar{\bar{X}} = \frac{1}{N} \sum \bar{X}_i$$

$$= \frac{1}{10} \times (12 + 10 + 10 + \dots + 14) = \frac{110}{10} = 11$$

$$\bar{s} = \frac{1}{N} \sum s_i$$

$$= \frac{1}{10} \times (2.5 + 2.3 + \dots + 1.3) = \frac{23.2}{10} = 2.32$$

From the table of control chart constants, for sample size  $n = 5$ , we have

$$A_1 = 1.596; B_3 = 0 \text{ and } B_4 = 2.089$$

### Control Limits for $\bar{X}$ -chart [Fig. 12.4 (a)]

$$CL = \bar{\bar{X}} = 11;$$

$$\begin{aligned} LCL &= \bar{\bar{X}} - A_1 \cdot \sqrt{\frac{n}{n-1} \bar{s}} \\ &= 11 - 1.596 \sqrt{\frac{5}{4}} \times 2.32 = 6.86 \end{aligned}$$

$$\begin{aligned} UCL &= \bar{\bar{X}} + A_1 \cdot \sqrt{\frac{n}{n-1} \bar{s}} \\ &= 11 + 1.596 \sqrt{\frac{5}{4}} \times 2.32 = 15.14 \end{aligned}$$

### Control Limits for $s$ -chart [Fig. 12.4 (b)]

$$CL = \bar{s} = 2.32$$

$$LCL = B_3 \bar{s} = 0$$

$$UCL = B_4 \bar{s} = 2.089 \times 2.32 = 4.85$$

( $s$ -chart is given on page 12.17)

### State of Control

The given sample mean ( $\bar{X}$ ) values lie between 6.86 and 15.14 and the given S.D. ( $s$ ) values between 0 and 4.85. Hence the process is under control with respect to average and variability.

#### Note

*Had we computed LCL and UCL for  $\bar{X}$  chart using ( $\bar{\bar{X}} \mp A_2 \bar{R}$ ), the process would be out of control with respect to the mean.*

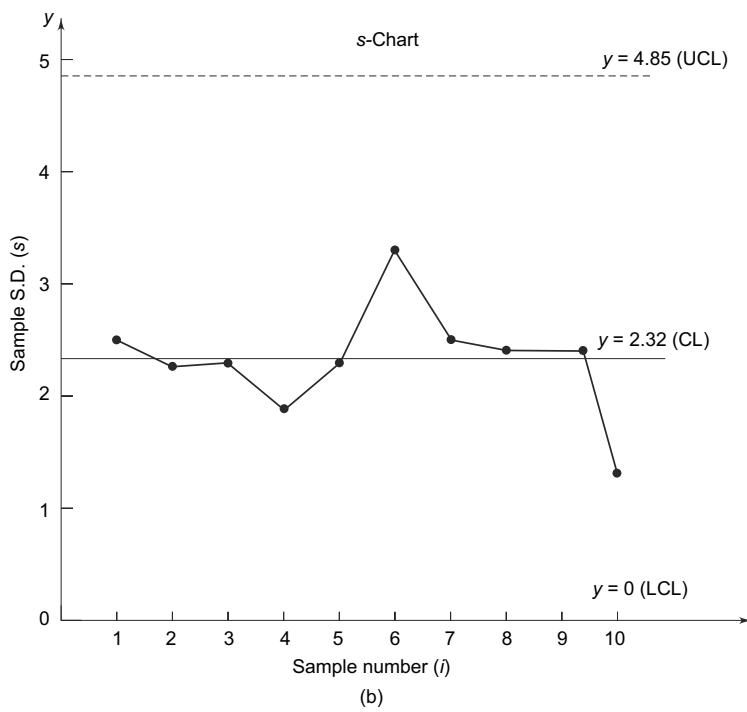
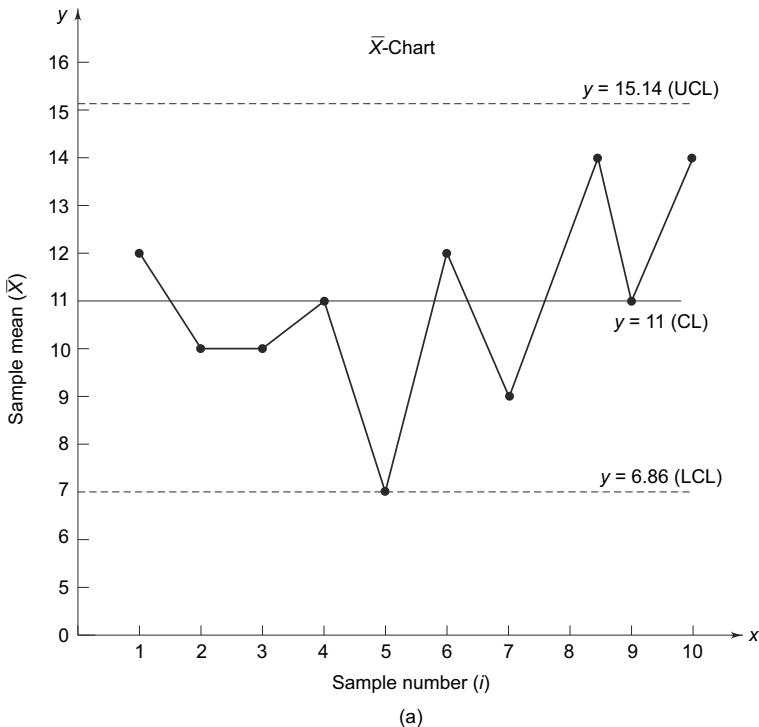


Fig. 12.4

**Example 5**

In a factory producing spark plugs, the number of defectives found in the inspection of 15 lots of 100 each is given below: Draw the control chart for the number of defectives and comment on the state of control.

Sample number ( $i$ ) :	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
Number of defective ( $np$ )	5	10	12	8	6	4	6	3	4	5	4	7	9	3	4

$$\sum np = 5 + 10 + 12 + \dots + 3 + 4 = 90$$

$$\therefore n\bar{p} = \frac{1}{N} \sum np = \frac{90}{15} = 6$$

$$\text{and } \bar{p} = \frac{1}{n} \times 6 = \frac{1}{100} \times 6 = 0.06 \quad (\because \text{each sample contains 100 items})$$

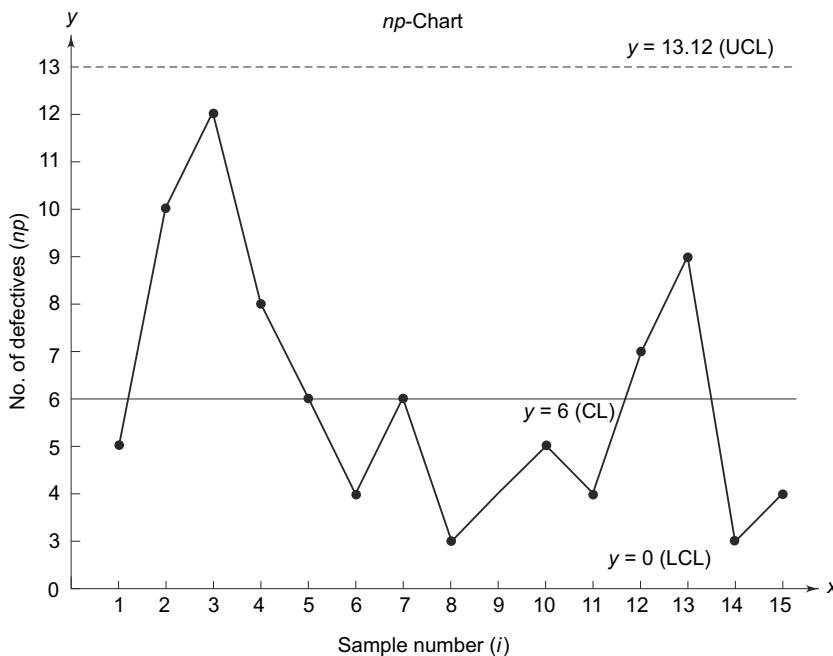
**For the  $np$ -chart (Fig. 12.5)**

$$CL = n\bar{p} = 6$$

$$LCL = n\bar{p} - 3\sqrt{n\bar{p}(1-\bar{p})} = 6 - 3\sqrt{6 \times 0.94} = -1.12$$

Since LCL cannot be negative, LCL = 0

$$UCL = n\bar{p} + 3\sqrt{n\bar{p}(1-\bar{p})} = 6 + 3\sqrt{6 \times 0.94} = 13.12$$

**Fig. 12.5**

Since all the sample points lie between the upper and lower control lines, the process is under control.

### Example 6

15 samples of 200 items each were drawn from the output of a process. The number of defective items in the samples are given below. Prepare a control chart for the fraction defective and comment on the state of control.

Sample No. : (i)	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
No. of defective ( $np$ )	12	15	10	8	19	15	17	11	13	20	10	8	9	5	8

$$\sum np = 12 + 15 + 10 + \dots + 5 + 8 = 180$$

$$\therefore n\bar{p} = \frac{1}{N} \sum np = \frac{180}{15} = 12$$

$$\therefore \bar{p} = \frac{12}{200} = 0.06 \quad (\because \text{each sample contains 200 items})$$

### For the $p$ -chart

$$CL = \bar{p} = 0.06$$

$$LCL = \bar{p} - 3\sqrt{\frac{\bar{p}(1-\bar{p})}{n}} = 0.06 - 3\sqrt{\frac{0.06 \times 0.94}{200}} = 0.01$$

$$UCL = \bar{p} + 3\sqrt{\frac{\bar{p}(1-\bar{p})}{n}} = 0.11$$

The fraction defectives (values of  $p$ ) for the given samples are 0.06, 0.075, 0.05, 0.04, 0.095, 0.075, 0.085, 0.055, 0.065, 0.1, 0.05, 0.04, 0.045, 0.025, 0.04 (1)  $p$ -chart is given on page 12.20.

Since all the sample points lie between the LCL and UCL lines, the process is under control.

### Example 7

10 samples each of size 50 were inspected and the number of defectives in the inspection were: 2, 1, 1, 2, 3, 5, 5, 1, 2, 3. Draw the appropriate control chart for defectives.

Since the number of defectives in 10 samples, each of size 50, are given, we may construct either number of defectives ( $np$ ) chart or proportion of defectives ( $p$ ) chart. We shall construct both and compare the charts.

$$\sum np = 2 + 1 + 1 + \dots + 2 + 3 = 25$$

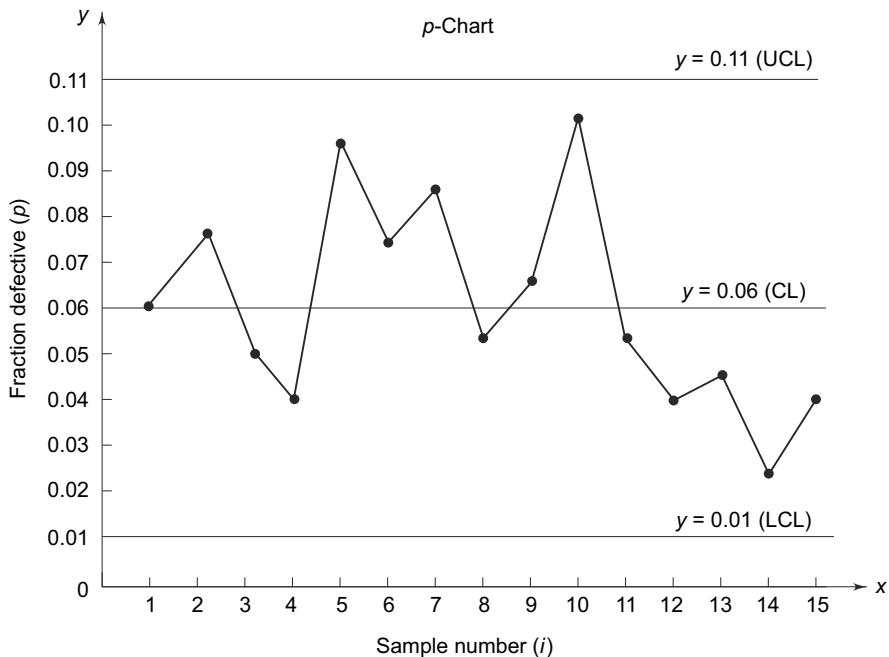


Fig. 12.6

$$\therefore n\bar{p} = \frac{1}{N} \sum np = \frac{25}{10} = 2.5$$

$$\text{and } \bar{p} = \frac{1}{n} \times n\bar{p} = \frac{1}{50} \times 2.5 = 0.05$$

### For the $np$ -chart [Fig. 12.7 (a)]

$$CL = n\bar{p} = 2.5;$$

$$LCL = n\bar{p} - 3\sqrt{n\bar{p}(1-\bar{p})} = 2.5 - 3\sqrt{2.5 \times 0.95} = -2.12$$

Since LCL cannot be negative, we take LCL = 0.

$$UCL = n\bar{p} + 3\sqrt{n\bar{p}(1-\bar{p})} = 2.5 + 3\sqrt{2.5 \times 0.95} = 7.12$$

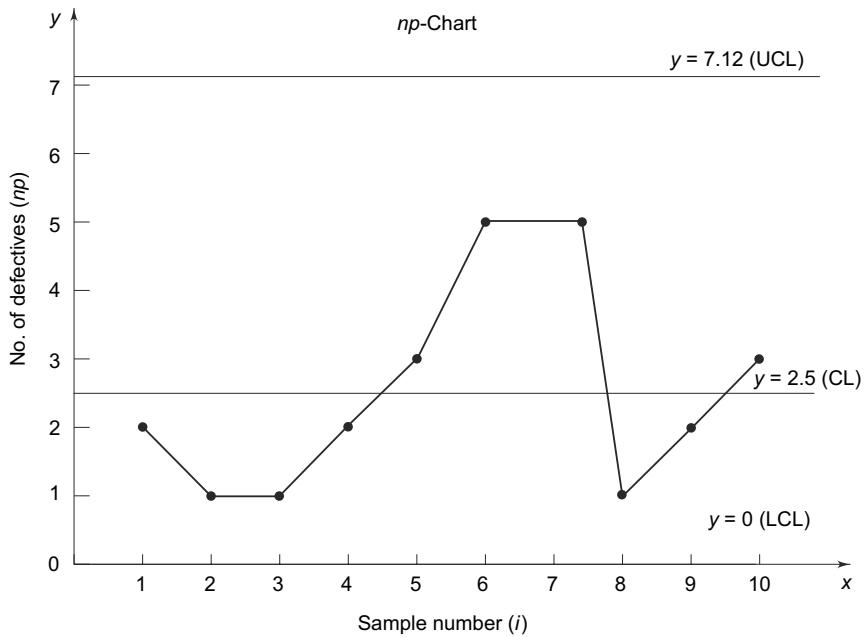
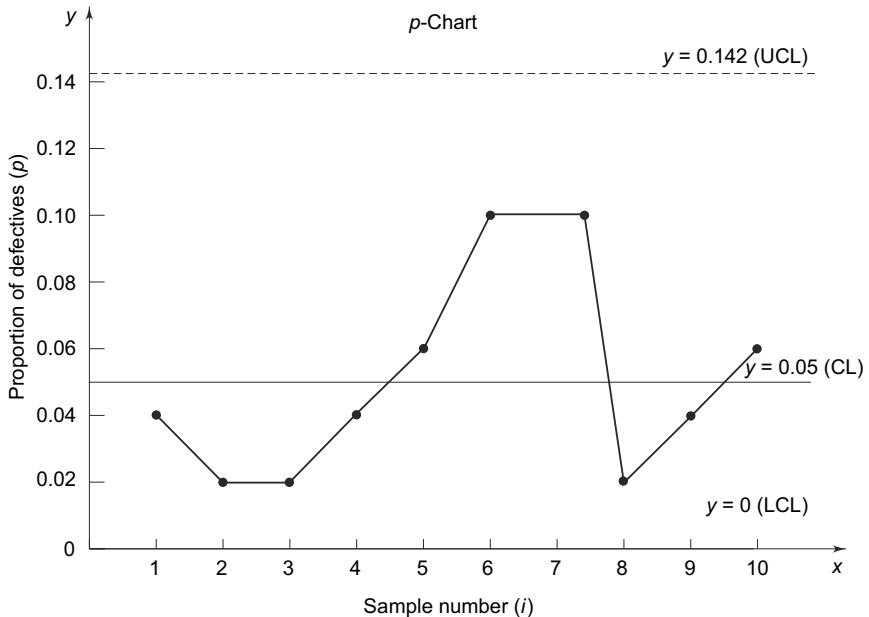
### For the $p$ -chart [Fig. 12.7 (b)]

$$CL = \bar{p} = 0.05;$$

$$LCL = \bar{p} - 3\sqrt{\frac{\bar{p}(1-\bar{p})}{n}} = 0.05 - 3\sqrt{\frac{0.05 \times 0.95}{50}} = -0.042$$

As LCL cannot be negative, we take LCL = 0.

$$UCL = \bar{p} + 3\sqrt{\frac{\bar{p}(1-\bar{p})}{n}} = 0.05 + 3\sqrt{\frac{0.05 \times 0.95}{50}} = 0.142$$

**Fig. 12.7 (a)****Fig. 12.7 (b)**

The proportion of defectives ( $p$ ) for the given samples are 0.04, 0.02, 0.02, 0.04, 0.06, 0.10, 0.10, 0.02, 0.04, 0.06.

Since all the sample points lie within the LCL and UCL lines in both  $np$  and  $p$ -charts, the process is under control. If we use suitable scales, we may see that the  $np$ -chart and  $p$ -chart are identical.

### Example 8

Construct a control chart for defectives for the following data:

Sample No.	:	1	2	3	4	5	6	7	8	9	10
No. inspected	:	90	65	85	70	80	80	70	95	90	75
No. of defectives	:	9	7	3	2	9	5	3	9	6	7

We note that the size of the sample varies from sample to sample. Hence, we cannot construct the  $np$ -chart. We can construct  $p$ -chart, provided  $0.75 \bar{n} < n_i < 1.25 \bar{n}$ , for all  $i$ .

$$\begin{aligned}\text{Here, } \bar{n} &= \frac{1}{N} \sum n_i = \frac{1}{10} \times (90 + 65 + \dots + 90 + 75) \\ &= \frac{1}{10} \times 800 = 80\end{aligned}$$

Hence,  $0.75 \bar{n} = 60$  and  $1.25 \bar{n} = 100$

The values of  $n_i$  be between 60 and 100. Hence,  $p$ -chart, given in Fig. 12.8, can be drawn by the method given below.

**Note** If the condition ( $0.75 \bar{n} < n_i < 1.25 \bar{n}$ ) is not satisfied, other available methods may be used. They are beyond the scope of this book.

$$\begin{aligned}\text{Now, } \bar{p} &= \frac{\text{Total no. of defectives}}{\text{Total no. of items inspected}} \\ &= \frac{60}{800} = 0.075\end{aligned}$$

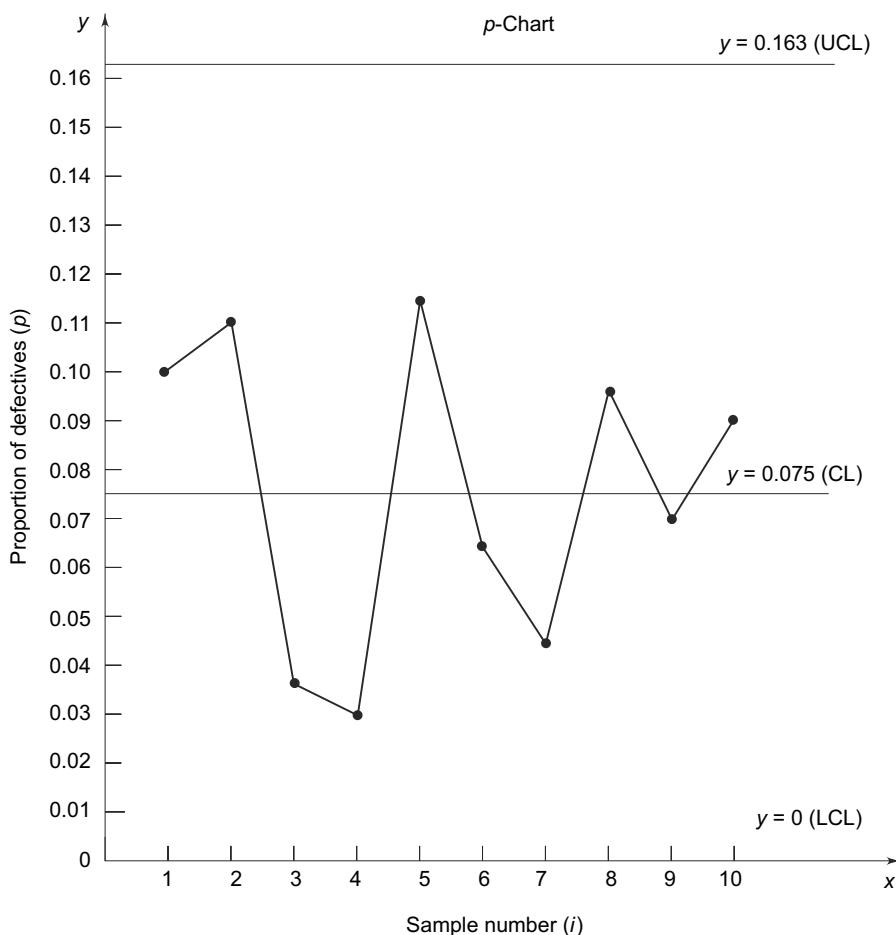
Hence, for the  $p$ -chart to be constructed,

$$\text{CL} = \bar{p} = 0.075$$

$$\begin{aligned}\text{LCL} &= \bar{p} - 3 \sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}} = 0.075 - 3 \sqrt{\frac{0.075 \times 0.925}{80}} \\ &= -0.013\end{aligned}$$

Since LCL cannot be negative, it is taken as 0.

$$\begin{aligned}\text{UCL} &= \bar{p} + 3 \sqrt{\frac{\bar{p}(1-\bar{p})}{\bar{n}}} = 0.075 + 3 \sqrt{\frac{0.075 \times 0.925}{80}} \\ &= 0.163\end{aligned}$$

**Fig. 12.8**

The values of  $p_i$  for the various samples are 0.100, 0.108, 0.035, 0.029, 0.113, 0.063, 0.043, 0.095, 0.067, 0.093.

Since all the sample points lie within the control lines, the process is under control.

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### Example 9

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15 tape-recorders were examined for quality control test. The number of defects in each tape-recorder is recorded below. Draw the appropriate control chart and comment on the state of control.

Unit no. ( $i$ ) :	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
No. of defects ( $c$ ) :	2	4	3	1	1	2	5	3	6	7	3	1	4	2	1

Since the number of defects per sample containing only one item is given, we can draw the  $c$ -chart (Fig. 12.9).

$$\begin{aligned}\bar{c} &= \frac{1}{N} \sum c_i = \frac{1}{15} (2 + 4 + 3 + \dots + 2 + 1) \\ &= \frac{45}{15} = 3\end{aligned}$$

**Note**

Even though  $\bar{c} < 4$ , we draw the  $c$ -chart, as it is the only possible chart.

$$\text{CL} = \bar{c} = 3; \text{LCL} = \bar{c} - 3\sqrt{\bar{c}} = 3 - 3\sqrt{3} = 2.20$$

Since LCL cannot be negative, we take LCL = 0

$$\text{UCL} = \bar{c} + 3\sqrt{\bar{c}} = 3 + 3\sqrt{3} = 8.20$$

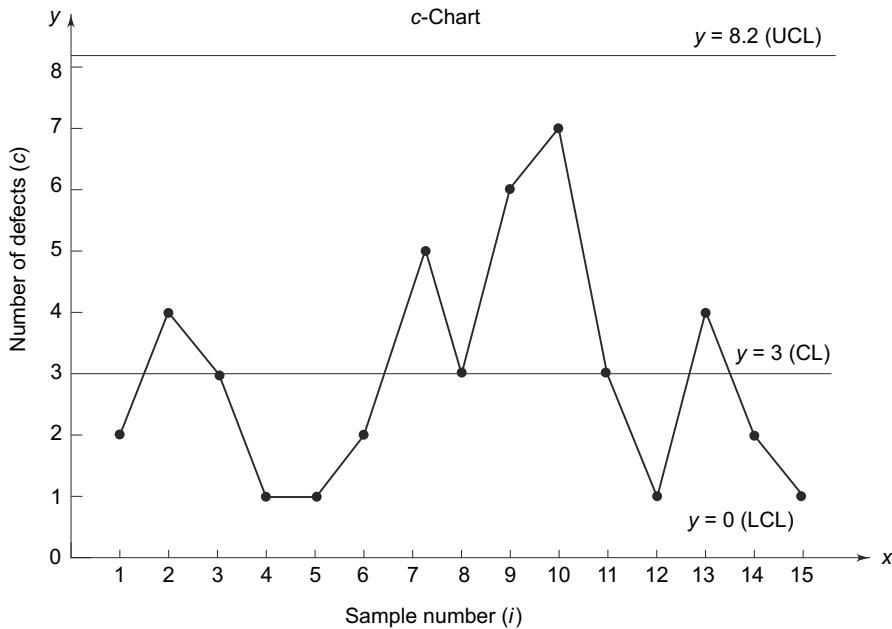


Fig. 12.9

Since all the sample points lie within the LCL and UCL lines, the process is under control.

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### Example 10

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A plant produces paper for newsprint and rolls of paper are inspected for defects. The results of inspection of 20 rolls of papers are given below: Draw the  $c$ -chart and comment on the state of control.

Roll No. ( $i$ ):	1	2	3	4	5	6	7	8	9	10
No. of defects ( $c$ ):	19	10	8	12	15	22	7	13	18	13
( $i$ )	11	12	13	14	15	16	17	18	19	20
( $c$ )	16	14	8	7	6	4	5	6	8	9

$$\bar{c} = \frac{1}{N} \sum c_i = \frac{1}{20} \times 220 = 11$$

For the  $c$ -chart (Fig. 12.10),

$$CL = \bar{c} = 11; LCL = \bar{c} - 3\sqrt{\bar{c}} = 11 - 3\sqrt{11} = 1.05$$

$$UCL = \bar{c} + 3\sqrt{\bar{c}} = 11 + 3\sqrt{11} = 20.95$$

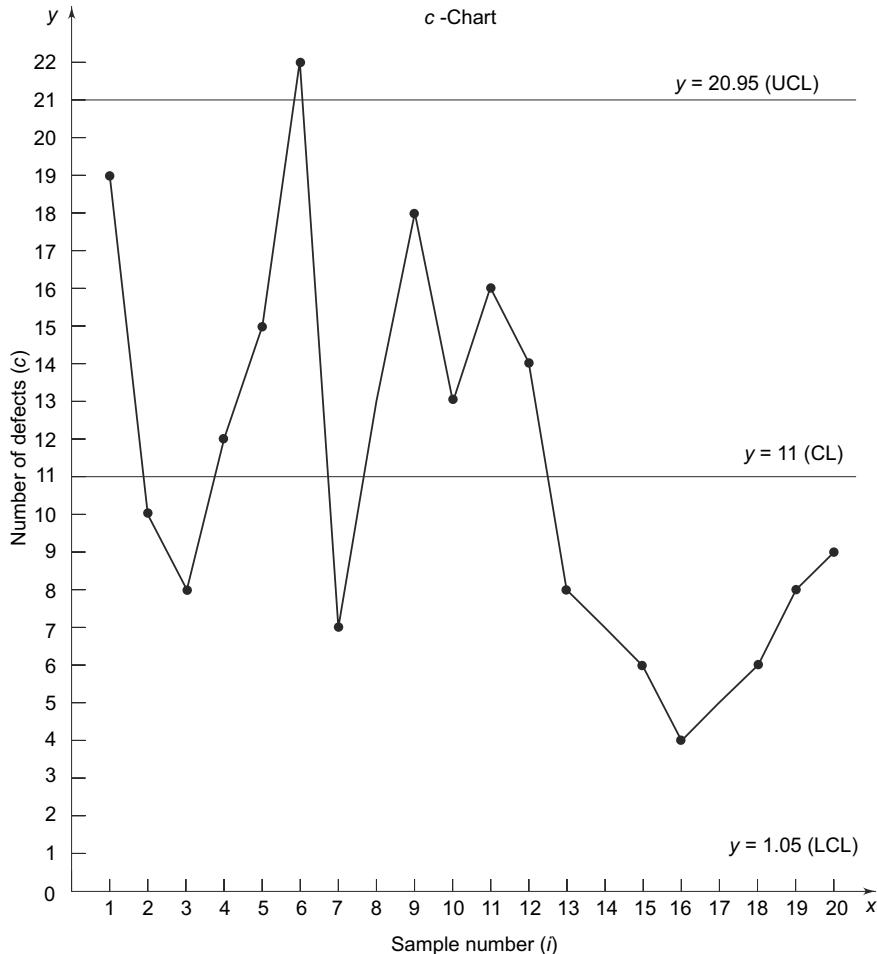


Fig. 12.10

Since one point falls outside the control lines, the process is out of control.

### Example 11

The specifications for a certain quality characteristic area  $15.0 \pm 6.0$  (in coded values). 15 samples of 4 readings each gave the following values for  $\bar{X}$  and  $R$ .

Sample No. ( $i$ ):	1	2	3	4	5	6	7
$\bar{X}$ :	16.1	15.2	14.2	13.9	15.4	15.7	15.2
$R$ :	3.0	2.1	5.6	2.4	4.1	2.7	2.3

$i$ :	8	9	10	11	12	13	14	15
$\bar{X}$ :	15.0	16.5	14.9	15.3	17.8	15.9	14.6	15.2
$R$ :	3.8	5.0	2.9	13.8	14.2	4.8	5.0	2.2

Compute the control limits for  $\bar{X}$  and  $R$ -charts using the above data for all the samples. Hence, examine if the process is in control. If not, remove the doubtful samples and recompute the values of  $\bar{\bar{X}}$  and  $\bar{\bar{R}}$ . After testing the state of control, estimate the tolerance limits and find if the process will meet the required specifications.

Let us consider all the 15 samples given.

$$\bar{\bar{X}} = \frac{1}{N} \sum \bar{X} = \frac{1}{15} (16.1 + 15.2 + \dots + 15.2)$$

$$= \frac{1}{15} \times 230.9 = 15.39$$

$$\bar{\bar{R}} = \frac{1}{N} \sum R = \frac{1}{15} (3.0 + 2.1 + \dots + 2.2)$$

$$= \frac{1}{15} \times 73.9 = 4.93$$

#### For the $\bar{X}$ -chart

$$CL = \bar{\bar{X}} = 15.39;$$

$$LCL = \bar{\bar{X}} - A_2 \bar{R} = 15.39 - 0.729 \times 4.93 = 11.80$$

$$UCL = \bar{\bar{X}} + A_2 \bar{R} = 15.39 + 0.729 \times 4.93 = 18.98$$

#### For the $R$ -chart

$$CL = \bar{\bar{R}} = 4.93; LCL = D_3 \bar{R} = 0$$

$$UCL = D_4 \bar{R} = 2.282 \times 4.93 = 11.25$$

The process is under control with respect to the average ( $\bar{X}$ -chart), but it is not under control with respect to variability ( $R$ -chart), since  $R_{11}$  ( $R$  value for the sample No. 11) = 13.8 and  $R_{12}$  = 14.2 exceed UCL = 11.25.

Hence, the process is not under control. So we remove the samples numbered 11 and 12 from the given data.

Let us now recompute  $\bar{\bar{X}}$  and  $\bar{R}$  based on the remaining 13 samples.

$$\bar{\bar{X}} = \frac{1}{13} (16.1 + 15.2 + \dots + 14.9 + 15.9 + 14.6 + 15.2)$$

$$= \frac{1}{13} \times 197.8 = 15.22$$

$$\bar{R} = \frac{1}{13} (3.0 + 2.1 + \dots + 2.9 + 4.8 + 5.0 + 2.2)$$

$$= \frac{1}{13} \times 45.9 = 3.53$$

Let us now recompute the revised control limits for  $\bar{X}$  and  $R$  charts.

### For the $\bar{X}$ -chart,

$$CL = \bar{\bar{X}} = 15.22;$$

$$\begin{aligned} LCL &= \bar{\bar{X}} - A_2 \bar{R} = 15.22 - 0.729 \times 3.53 \\ &= 12.65 \end{aligned}$$

$$UCL = \bar{\bar{X}} + A_2 \bar{R} = 15.22 + 0.729 \times 3.53 = 17.79$$

### For the $R$ -chart

$$CL = \bar{R} = 3.53; LCL = D_3 \bar{R} = 0;$$

$$UCL = D_4 \bar{R} = 2.282 \times 3.53 = 8.06$$

We see that the process is under control with respect to the 13 samples considered.

Now we can compute the tolerance limits using the revised values of  $\bar{\bar{X}}$  and  $\bar{R}$ .

The tolerance limits are given by

$$\bar{\bar{X}} \mp \frac{3\bar{R}}{d_2}$$

$$= 15.22 \mp \frac{3 \times 3.53}{2.059}$$

(The value of  $d_2$  is read from the table of control chart constants for  $n = 4$ )  
 $= 15.22 \pm 5.14$

Thus, the required tolerance limits are  $(10.08, 20.36)$

Since these tolerance limits lie within the specification limits  $(9.0, 21.0)$ , the process meets the required specifications.

### Example 12

The specifications for a certain quality characteristic are  $(60 \pm 24)$  in coded values. The table given below gives the measurements obtained in 10 samples. Find the tolerance limits for the process and test if the process meets the specifications.

Sample No. ( $i$ )	1	2	3	4	5	6	7	8	9	10
Measurements ( $X$ )	75	48	57	61	55	49	74	67	66	62
	66	79	55	71	68	98	63	70	65	68
	50	53	53	66	58	65	62	68	58	66
	62	61	61	69	62	64	57	56	52	68
	52	49	72	77	75	66	62	61	58	73
	70	56	63	53	63	64	64	66	50	68

The values of  $\bar{X}$  and  $R$  computed for all the samples are tabulated below:

$i :$	1	2	3	4	5	6	7	8	9	10
$\sum X_i :$	375	346	361	397	381	406	382	388	349	405
$\bar{X}_i :$	62.5	57.7	60.2	66.2	63.5	67.7	63.7	64.7	58.2	67.5
$R_i :$	25	31	19	24	20	49	17	14	16	11

$$\text{Now, } \bar{\bar{X}} = \frac{1}{N} \sum \bar{X}_i = \frac{1}{10} \times 631.9 = 63.19$$

$$\bar{R} = \frac{1}{N} \sum R_i = \frac{1}{10} \times 226 = 22.6$$

#### For the $\bar{X}$ -chart

$$CL = \bar{\bar{X}} = 63.19;$$

$$LCL = \bar{\bar{X}} - A_2 \bar{R} = 63.19 - 0.483 \times 22.6 = 52.27$$

$$UCL = \bar{\bar{X}} + A_2 \bar{R} = 63.19 + 0.483 \times 22.6 = 74.11$$

### For the R-chart

$$\text{CL} = \bar{R} = 22.6; \text{LCL} = D_3 \bar{R} = 0;$$

$$\text{UCL} = D_4 \bar{R} = 2.004 \times 22.6 = 45.29$$

We note that

$$\text{LCL} (= 52.27) < \bar{X}_i < \text{UCL} (= 74.11),$$

but  $R_6$  (sample 6) > 45.29.

Hence, the process is not under control.

Now, we remove sample No. 6 from the data and recompute the values of  $\bar{\bar{X}}$  and  $\bar{R}$  based on the remaining 9 samples.

$$\bar{\bar{X}} = \frac{1}{9} \times (62.5 + 57.7 + \dots + 63.5 + 63.7 + \dots + 67.5)$$

$$= \frac{1}{9} \times 564.2 = 62.69$$

$$\bar{R} = \frac{1}{9} \times (25 + 31 + \dots + 20 + 17 + \dots + 11)$$

$$= \frac{1}{9} \times 177 = 19.67$$

Let us now recompute the revised control limits for  $\bar{X}$ - and  $R$ -charts.

### For the $\bar{X}$ -chart

$$\text{CL} = \bar{\bar{X}} = 62.69;$$

$$\text{LCL} = \bar{\bar{X}} - A_2 \bar{R} = 62.69 - 0.483 \times 19.67 = 53.19$$

$$\text{UCL} = \bar{\bar{X}} + A_2 \bar{R} = 62.69 + 0.483 \times 19.67 = 72.19$$

### For the R-chart

$$\text{CL} = \bar{R} = 19.67, \text{LCL} = D_3 \bar{R} = 0;$$

$$\text{UCL} = D_4 \bar{R} = 2.004 \times 19.67$$

$$= 39.42$$

Now,  $53.19 < \bar{X}_i < 72.19$  and  $0 < R_i < 39.42$ , for all  $i$  ( $\neq 6$ )

Hence, the process is under control.

The tolerance limits are given by

$$\bar{\bar{X}} \mp \frac{3\bar{R}}{d_2}$$

$$\begin{aligned}
 &= 62.69 \pm \frac{3 \times 19.67}{2.534} \\
 &= 62.69 \pm 23.29
 \end{aligned}
 \quad (\because d_2 = 2.534, \text{ for } n = 6)$$

i.e., the tolerance limits are (39.40, 85.98)

The specification limits are (36, 84).

Since the upper tolerance limit exceeds the upper specification limit, the process does not meet the specifications.

## Exercise 12

### Part-A (Short-answer Questions)

1. What is meant by 'quality' in the term Statistical Quality Control?
2. What do you mean by Statistical Quality Control?
3. What is the difference between chance variation and assignable variation?
4. What do you understand by process control?
5. What is control chart? Name the types of control charts.
6. Distinguish between variables and attributes in connection with the quality characteristics of a product.
7. Name any two control charts, for each, of variables and attributes.
8. Find the lower and upper control limits for  $\bar{X}$ -chart and  $R$ -chart, when each sample is of size 4 and  $\bar{\bar{X}} = 10.80$  and  $\bar{R} = 0.46$ .
9. When do you say that a process is out of control?
10. Find the lower and upper control limits for  $\bar{X}$ -chart and  $s$ -chart, if  $n = 5$ ,  $\bar{\bar{X}} = 15$  and  $\bar{s} = 2.5$ .
11. Under what situations  $p$ -chart is drawn instead of  $np$ -chart?
12. When  $n$  is constant, will the  $p$ -chart and  $np$ -chart lead to the same conclusions regarding the process of control?
13. Find the lower and upper control limits for  $p$ -chart and  $np$ -chart, when  $n = 100$  and  $\bar{p} = 0.085$ .
14. Distinguish between  $p$ -chart and  $c$ -chart.
15. Find the lower and upper control limits for the  $c$ -chart, when  $\bar{c} = 6$ .
16. What do you mean by specification limits in a manufacturing process?
17. What is meant by tolerance limits?
18. Distinguish between control limits and tolerance limits.
19. When the process is under control and if  $n = 5$ ,  $\bar{\bar{X}} = 1.1126$  and  $\bar{R} = 0.0054$ , find the tolerance limits.
20. How will you decide whether a process is operating at an acceptable level?

**Part-B**

21. The following data give the average life in hours and range in hours of 12 samples each of 5 lamps. Construct the control charts for  $\bar{X}$  and  $R$  and comment on the state of control.

$\bar{X}$ :	120	127	152	157	160	134	137	123	140	144	120	127
$R$ :	30	44	60	34	38	35	45	62	39	50	35	41

22. Draw the mean chart and range chart using the following data relating to 15 samples each of size 5 and comment on the state of control.

$\bar{X}$ :	65.0	64.6	64.1	68.5	68.4	67.9	65.0	64.6	
$\bar{X}$ :	64.1	63.2	62.9	62.4	67.0	66.6	66.1		
$R$ :	9.8	9.8	8.4	3.9	7.6	8.7	0.1	9.7	
$R$ :	7.7	7.5	1.2	9.8	6.4	0.6	6.3		

23. A food company puts mango juice in cans, each of which is advertised to contain 10 ounces of the juice. The weights of the juice drained from cans immediately after filling 20 samples each of 4 cans are taken by random sampling method (at an interval of 30 minutes) and given in the following table in units of 0.01 ounce in excess of 10 ounces. To control the excess weights of mango juice drained while filling, draw the  $\bar{X}$ -chart and  $R$ -chart and comment on the nature of control.

Sample No.	1	2	3	4	5	6	7	8	9	10
Weights drained	15	10	8	12	18	20	15	13	9	6
	12	8	15	17	13	16	19	23	8	10
	13	8	17	11	15	14	23	14	18	24
	20	14	10	12	4	20	17	16	5	20

Sample No.	11	12	13	14	15	16	17	18	19	20
Weights drained	5	3	6	12	15	18	13	10	5	6
	12	15	18	9	15	17	16	20	15	14
	20	18	12	15	6	8	5	8	10	12
	15	18	10	18	16	15	4	10	12	14

24. The table below gives the (coded) measurements obtained in 20 samples. Construct control charts based on mean and range and comment on the state of control:

Sample No.	1	2	3	4	5	6	7	8	9	10	11	12	13
Values of X	2	0	1	1	-1	-1	-1	1	1	1	1	1	1
	1	1	0	0	0	2	0	2	-1	-2	-3	-1	-3
	0	0	0	-1	0	0	-2	-1	0	2	2	0	2
	1	1	1	0	-1	-2	1	0	0	1	1	0	1

Sample No.	14	15	16	17	18	19	20
Values of X	0	-1	1	2	2	0	3
	0	2	-1	1	0	2	-3
	-1	1	2	-1	1	1	-1
	0	1	0	0	0	-1	1
	1	2	2	0	1	1	2

25. The following data give the coded values of the crushing strengths of concrete blocks obtained from 20 samples each of size 5. Draw the  $\bar{X}$  and  $R$ -charts and comment on the state of control.

Sample No.	1	2	3	4	5	6	7	8	9	10
Values of X	11.1	9.6	9.7	10.1	12.4	10.1	11.0	11.2	10.6	8.3
	9.4	10.8	10.0	8.4	10.0	10.2	11.5	10.0	10.4	10.2
	11.2	10.1	10.0	10.2	10.7	10.2	11.8	10.9	10.5	9.8
	10.4	10.8	9.8	9.4	10.1	11.2	11.0	11.2	10.5	9.5
	10.1	11.0	10.4	11.0	11.3	10.1	11.3	11.0	10.9	9.8

Sample No.	11	12	13	14	15	16	17	18	19	20
Values of X	10.6	10.8	10.7	11.3	11.4	10.1	10.7	11.9	10.8	12.4
	9.9	10.2	10.7	11.4	11.2	10.1	12.8	11.9	12.1	11.1
	10.7	10.5	10.8	10.4	11.4	9.7	11.2	11.6	11.8	10.8
	10.2	8.4	8.6	10.6	10.1	9.8	11.2	12.4	9.4	11.0
	11.4	9.9	11.4	11.1	11.6	10.5	11.3	11.4	11.6	11.9

26. The following data gives the measurements of 10 samples each of size 5, in a production process taken at intervals of 2 hours. Draw the control charts for the mean and range and comment on the state of control.

Sample No.	1	2	3	4	5	6	7	8	9	10
Measurements (X)	47	52	48	49	50	55	50	54	49	53
	49	55	53	49	53	55	51	54	55	50
	50	47	51	49	48	50	53	52	54	54
	44	56	50	53	52	53	46	54	49	47
	45	50	53	45	47	57	50	56	53	51

27. The following table gives the sample means and S.D.'s for 15 samples, each size 4, in the production of a certain component. Draw the  $\bar{X}$  and  $s$ -charts and comment on the state of control.

Sample No. (i):	1	2	3	4	5	6	7	8	9
$\bar{X}$ :	1.75	1.32	1.18	0.48	2.30	1.25	1.52	1.78	1.90
$s$ :	0.36	0.53	0.14	0.18	0.55	0.74	0.38	0.45	0.87

Sample No. (i):	10	11	12	13	14	15			
$\bar{X}$ :	1.72	2.40	3.20	2.52	2.05	1.68			
$s$ :	0.83	0.76	0.99	0.65	0.22	0.14			

28. The following data give the mean and S.D. values of 10 samples, each of size 5 drawn from a production process taken at intervals of one hour. Construct the mean and S.D. charts and comment on the state of control.

Sample No.:	1	2	3	4	5	6	7	8	9	10
$\bar{X}$ :	54	51	54	49	52	47	51	50	50	52
$s$ :	3.3	2.4	3.8	3.3	3.4	4.6	1.9	2.5	2.5	2.9

29. The mean and S.D. values of 15 samples each of size 5 are given in the following table. Draw the mean and S.D. charts and comment on the state of control.

Sample No. (i):	1	2	3	4	5	6	7	8	9
$\bar{X}$ :	11.0	12.4	8.6	12.2	9.2	10.6	10.0	9.4	10.6
$s$ :	3.63	3.14	2.87	3.49	1.47	3.01	2.45	3.98	2.15

Sample No. (i):	10	11	12	13	14	15			
$\bar{X}$ :	11.4	14.8	10.6	11.6	9.0	12.8			
$s$ :	2.15	2.14	3.01	3.14	1.67	2.56			

30. Draw the mean and S.D. charts using the data values for the first 10 samples given in Exercise 24 and comment on the state of control.

31. 15 samples of 4 machined items were taken periodically from production and inside diameter of whole in each item was measured. The measurements (coded) are given in the following table. Draw the  $\bar{X}$  and  $s$ -charts and comment on the state of control:

Sample No. ( $i$ ):	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
Measurements ( $X$ )	15	11	8	15	6	13	16	12	15	18	11	14	11	12	5
	14	16	11	14	7	7	10	9	11	16	13	12	9	6	10
	13	6	9	5	10	11	7	16	10	14	10	15	12	4	6
	15	15	9	15	7	11	7	10	10	7	8	12	14	9	10

32. Fifteen samples each of size 50 were inspected and the number of defectives in the inspection were  
2, 3, 4, 2, 3, 0, 1, 2, 2, 3, 5, 5, 1, 2, 3  
Draw the control chart for the number of defectives and comment on the state of control.
33. In a manufacturing company where spark-plugs are produced, the number rejected by inspection of 20 lots of 100 plugs each is given below: Construct the  $np$ -chart and comment on the state of control.

Lot No. ( $i$ ):	1	2	3	4	5	6	7	8	9	10
No. rejected ( $np$ ):	10	6	5	10	8	5	12	8	4	3

Lot No. ( $i$ ):	11	12	13	14	15	16	17	18	19	20
No. rejected ( $np$ ):	6	5	2	8	7	6	3	3	5	4

34. Using the following data, construct the  $np$ -chart and comment on the state of control. Assume that 200 items are inspected each day.

Day ( $i$ ):	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
No. of defective ( $np$ ):	6	6	6	5	0	0	6	14	4	0	1	8	2	4	7

35. On inspection of 20 lots each of 50 items, the following numbers of defectives were found: Construct the control chart for the fraction defectives and comment on the state of control.  
1, 4, 3, 4, 4, 3, 2, 4, 2, 3, 5, 1, 2, 4, 2, 4, 8, 3, 2, 2
36. In an integrated circuit production line, 15 samples of 100 units are checked for electrical specifications on alternate days of a month and the number of defectives found are tabulated below: Draw the  $p$ -chart and comment on the nature of control.  
24, 62, 26, 38, 33, 44, 45, 34, 30, 52, 44, 52, 36, 34 and 38.

37. On inspection of 10 samples, each of size 400, the numbers of defective articles were

19, 4, 9, 12, 9, 15, 26, 14, 15, 17.

Draw the  $np$ -chart and  $p$ -chart and comment on the state of control.

38. Draw the appropriate control chart for the following data and comment on the state of control:

Day:	1	2	3	4	5	6	7	8	9	10
No. inspected:	150	184	181	196	180	174	210	210	195	210
No. of defectives:	25	10	3	14	6	15	43	28	39	25

39. Ten samples of varying size are taken from a production line and the number of defectives is found in each sample. The results are given below. Draw the appropriate control chart and comment on the state of control.

Sample No.	1	2	3	4	5	6	7	8	9	10
No. inspected:	155	160	156	156	164	160	161	173	148	167
No. of defectives:	8	8	8	7	8	6	5	10	7	9

40. Twenty half-litre milk bottles are selected at random from a process and the numbers of air bubbles (defects) observed from the bottles are given in the following table. Draw the appropriate control chart and comment on the nature of control.

Bottle no. ( $i$ ):	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15	16	17	18	19	20
No. of defects ( $c$ ):	4	5	7	3	3	5	6	2	4	8	3	5	4	3	4	5	7	3	6	13

41. The following data relate to the number of defects in each of 15 units drawn randomly from a production process. Draw the control chart for the number of defects and comment on the state of control.

6, 4, 9, 10, 11, 12, 20, 10, 9, 10, 15, 10, 20, 15, 10.

42. Construct a  $c$ -chart for the number of defects from the following data which represent the number of imperfections in 20 pieces of cloth of equal length in a certain production of a mill. Is the process under control?

No. of imperfections: 3, 3, 4, 10, 10, 3, 3, 3, 6, 5, 6, 10, 4, 7, 4, 7, 4, 8, 4 and 7.

43. The specifications for the length of a certain product are  $(30 \pm 10)$  mm. 15 samples each of size 5 gave the following values for  $\bar{X}$  and  $R$  in mm.

Compute the tolerance limits of the production process and also find if the process will meet the specifications:

Sample no. ( $i$ ) :	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
$\bar{X}$ :	25	30	23	38	15	28	38	18	25	40	29	39	37	29	36
$R$ :	20	10	20	21	11	22	22	11	20	19	14	9	16	18	34

44. The specification limits for a quality characteristic are 1.100 and 1.120 (in certain units). 15 samples of 5 measurements each gave the following values for  $\bar{X}$  and  $R$  (Range values are multiplied by  $10^3$ ). Compute the tolerance limits of the process and check if the process will meet the specifications:

Sample no. ( $i$ ):	1	2	3	4	5	6	7	8
$\bar{X}$ :	1.115	1.116	1.114	1.112	1.114	1.112	1.114	1.112
$R$ :	18	17	8	6	7	5	5	7

Sample no. ( $i$ ):	9	10	11	12	13	14	15
$\bar{X}$ :	1.113	1.111	1.113	1.114	1.111	1.113	1.111
$R$ :	3	4	6	4	3	5	7

45. The following data give the measurements of 10 samples, each of size 5, drawn from a process at regular intervals. Find the tolerance limits for the process and test if the specifications ( $50 \pm 10$ ) are met by the process.

Sample no.	1	2	3	4	5	6	7	8	9	10
Measurement ( $X$ )	49	50	50	48	47	52	49	55	53	54
	55	51	53	53	49	55	49	55	50	54
	54	53	48	51	50	47	49	50	54	52
	49	46	52	50	44	56	53	53	47	54
	53	50	47	53	45	50	45	57	51	56

**Table 12.1 Table of Control Chart Constants**

Sample size $n$	Factors for $p$ -chart						Factors for s-chart						Factors for R-chart					
	$A$	$A_1$	$A_2$	$c_2$	$B_1$	$B_2$	$B_3$	$B_4$	$d_2$	$D_1$	$D_2$	$D_3$	$D_4$	$D_5$	$D_6$	$D_7$	$D_8$	
2	2.121	3.760	1.880	0.5642	0	1.843	0	3.267	1.128	0	3.686	0	3.262					
3	1.732	2.394	1.023	0.7236	0	1.858	0	2.568	1.693	0	4.358	0	2.575					
4	1.500	1.880	0.729	0.7979	0	1.808	0	2.266	2.059	0	4.698	0	2.282					
5	1.342	1.596	0.577	0.8407	0	1.756	0	2.089	2.326	0	4.918	0	2.115					
6	1.225	1.410	0.483	0.8686	0.026	1.711	0.030	1.970	2.534	0	5.078	0	2.004					
7	1.134	1.277	0.419	0.88882	0.105	1.672	0.118	1.882	2.704	0.205	5.203	0.076	1.924					
8	1.061	1.175	0.373	0.9027	1.167	1.638	0.185	1.815	2.847	0.387	5.307	5.136	0.864					
9	1.000	1.094	0.337	0.9139	0.219	0.609	0.239	1.760	2.970	0.546	5.394	0.184	1.816					
10	0.949	1.028	0.308	0.9227	0.262	1.584	0.284	1.716	3.078	0.687	5.469	0.223	1.777					
11	0.905	0.973	0.285	0.9300	0.299	1.561	0.321	1.679	3.173	0.812	5.534	0.256	0.744					
12	0.866	0.925	0.266	0.9359	0.331	1.541	0.354	1.646	3.258	0.924	5.592	0.284	1.716					
13	0.832	0.884	0.249	0.9410	0.359	1.523	0.382	1.618	3.336	1.026	5.646	0.308	1.692					
14	0.802	0.848	0.235	0.9453	0.384	1.507	0.406	1.594	3.407	1.121	5.693	0.329	1.671					
15	0.775	0.816	0.223	0.9490	0.406	1.492	0.428	1.572	3.472	1.207	5.737	0.348	1.652					
16	0.750	0.788	0.212	0.9523	0.427	1.478	0.448	1.552	3.532	1.285	5.779	0.364	1.636					
17	0.728	0.762	0.203	0.9551	0.445	1.465	0.466	1.534	3.588	1.359	5.817	0.379	1.621					
18	0.707	0.738	0.194	0.9576	0.461	0.459	0.482	1.518	3.640	1.426	5.854	0.392	1.608					
19	0.688	0.717	0.184	0.9599	0.477	0.443	0.497	1.503	3.689	1.490	5.888	0.404	1.596					
20	0.671	0.697	0.110	0.9619	0.491	1.433	0.510	1.490	3.735	1.544	5.922	0.418	1.586					

***Guide for using appropriate control chart constants***

<i>Quality Controlled</i>	<i>Central Line</i>	<i>LCL size (n)</i>	<i>UCL</i>	<i>Sample</i>	<i>Remarks</i>
$\bar{X}$ (Mean)	$\mu$	$\mu - A\sigma$	$\mu + A\sigma$	Any value of $n$	When $\mu, \sigma$ are given
$\bar{X}$	$\bar{\bar{X}}$	$\bar{\bar{X}} - A_2 \bar{R}$	$\bar{\bar{X}} + A_2 \bar{R}$	$n \leq 10$	When $\mu, s$ are not known,
$\bar{X}$	$\bar{\bar{X}}$	$\bar{\bar{X}} - A_1 \sqrt{\frac{n-1}{n}}$	$\bar{\bar{X}} + A_1 \sqrt{\frac{n-1}{n}}$	$n \leq 25$ ( $n$ constant)	When $\mu, s$ are not known, but $\bar{\bar{X}} \bar{s}$ , are known
$\bar{X}$	$\bar{\bar{X}}$	$\bar{\bar{x}} - \frac{3\bar{s}}{\sqrt{n}}$	$\bar{\bar{X}} + \frac{3\bar{s}}{\sqrt{n}}$	$n > 25$ ( $n$ varying slightly)	-Do-
$R$ (Range)	$d_2 \sigma$	$D_1 \sigma$	$D_2 \sigma$	$n \leq 10$	When $\sigma$ is given
$R$	$\bar{R}$	$D_3 \bar{R}$	$D_4 \bar{R}$	$n \leq 10$	When $\sigma$ is not known
$\sigma$ (S.D.)	$c_2 \sqrt{\frac{n}{n-1}} \sigma$	$B_1 \sqrt{\frac{n}{n-1}} \sigma$	$B_2 \sqrt{\frac{n}{n-1}} \sigma$	Any value of $n$	When $\sigma$ is given
$s$	$\bar{s}$	$B_3 \bar{s}$	$B_4 \bar{s}$	$n \leq 25$ ( $n$ constant)	When $\sigma$ is not known, but $\bar{s}$ is known
$s$	$\bar{s}$	$\bar{s} - \frac{3\bar{s}}{\sqrt{2n}}$	$\bar{s} + \frac{3\bar{s}}{\sqrt{2n}}$	$n > 25$ ( $n$ varying slightly)	-Do-

**ANSWERS****Exercise 12**

8. For  $\bar{X}$ -chart, LCL = 10.46, UCL = 11.14;  
For  $R$ -chart, LCL = 0, UCL = 1.05.
10. For  $\bar{X}$ -chart, LCL = 11.43, UCL = 18.57;  
For  $s$ -chart, LCL = 0, UCL = 5.22.
13. For  $p$ -chart, LCL = 0.0013, UCL = 0.1687;  
For  $np$ -chart, LCL = 0.134, UCL = 16.867.

19. 1.1056, 1.1196
21. For  $\bar{X}$ -chart, LCL = 112.08, UCL = 161.42;  
For  $R$ -chart, LCL = 0, UCL = 90.42  
Process under control.
22. For  $\bar{X}$ -chart, LCL = 61.61, UCL = 69.11;  
For  $R$ -chart, LCL = 0, UCL = 13.75.  
Process under control.
23. LCL for  $\bar{X} = 5.46$ , UCL for  $\bar{X} = 20.84$ ;  
LCL for  $R = 0$ , UCL for  $R = 24.84$ ; under control
24. LCL for  $\bar{X} = 1.488$ , UCL for  $\bar{X} = 2.148$ ;  
LCL for  $R = 0$ , UCL for  $R = 6.662$ ; under control
25. LCL for  $\bar{X} = 9.74$ , UCL for  $\bar{X} = 11.58$ ;  
LCL for  $R = 0$ , UCL for  $R = 3.35$ ; out of control
26. LCL for  $\bar{X} = 47.25$ , UCL for  $\bar{X} = 54.75$ ;  
LCL for  $R = 0$ , UCL for  $R = 13.75$ ; out of control
27. LCL for  $\bar{X} = 0.67$ , UCL for  $\bar{X} = 2.93$ ;  
LCL for  $s = 0$ , UCL for  $s = 1.18$ ; out of control
28. LCL for  $\bar{X} = 45.5$ , UCL for  $\bar{X} = 56.5$ ;  
LCL for  $s = 0$ , UCL for  $s = 6.39$ ; under control
29. LCL for  $\bar{X} = 6.09$ , UCL for  $\bar{X} = 15.81$ ;  
LCL for  $s = 0$ , UCL for  $s = 5.69$ ; under control
30. LCL for  $\bar{X} = -1.55$ , UCL for  $\bar{X} = 2.07$ ;  
LCL for  $s = 0$ , UCL for  $s = 2.12$ ; under control
31. LCL for  $\bar{X} = 5.62$ , UCL for  $\bar{X} = 16.18$ ;  
LCL for  $s = 0$ , UCL for  $s = 5.51$ ; under control
32. LCL = 0, UCL = 7.18; under control
33. LCL = 0, UCL = 13.12; under control
34. LCL = 0, UCL = 10.96; out of control
35. LCL = 0, UCL = 0.166; under control
36. LCL = 0.2481, UCL = 0.5413; out of control
37. LCL for  $np = 2.97$ , UCL for  $np = 25.03$ ;  
LCL for  $p = 0.028$ , UCL for  $p = 0.063$ ; out of control.
38. LCL = 0.042, UCL = 0.178; out of control.
39. LCL = 0, UCL = 0.0979; under control.
40. LCL = 0, UCL = 11.708; out of control.
41. LCL = 1.27, UCL = 21.13; under control.
42. LCL = 0, UCL = 12.62; under control.
43. (7.77, 56.39); does not meet the specifications.
44. (1.1056, 1.1196); meets the specifications.
45. (42.98, 59.90); meets the specifications.



## Appendix

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# Statistical Tables

*(Normal Table, t-table,  $\chi^2$ -table and F-table)*

$z$	Ordinate	$Area = \int_0^z$	$z$	Ordinate	$Area = \int_0^z$
0.00	0.3989	0.0000	1.55	0.1200	0.4394
0.05	0.3984	0.0199	1.60	0.1109	0.4452
0.10	0.3970	0.0398	1.65	0.1023	0.4505
0.15	0.3945	0.0596	1.70	0.941	0.4554
0.20	0.3910	0.0793	1.75	0.0863	0.4099
0.25	0.3867	0.0987	1.80	0.0790	0.4641
0.30	0.3814	0.1179	1.85	0.0721	0.4678
0.35	0.3752	0.1368	1.90	0.0656	0.4719
0.40	0.3683	0.1554	1.95	0.0596	0.4744
0.45	0.3605	0.1736	2.00	0.0540	0.4773
0.50	0.3521	0.1910	2.05	0.0488	0.4798
0.55	0.3429	0.2088	2.10	0.0440	0.4821
0.60	0.3332	0.2258	2.15	0.0396	0.4842
0.65	0.3230	0.2422	2.20	0.0355	0.4861
0.70	0.3123	0.2080	2.25	0.0317	0.4878
0.75	0.3011	0.2734	2.30	0.0283	0.4893
0.80	0.2897	0.2882	2.35	0.0252	0.4906
0.85	0.2780	0.3023	2.40	0.0224	0.4918
0.90	0.2661	0.3159	2.45	0.0224	0.4929
0.95	0.2541	0.3289	2.50	0.0198	0.4938
1.00	0.2420	0.3413	2.55	0.0175	0.4946
1.05	0.2299	0.3531	2.60	0.0155	0.4953
1.10	0.2179	0.3643	2.65	0.0136	0.4960
1.15	0.2059	0.3749	2.70	0.0119	0.4965
1.20	0.1942	0.3849	2.75	0.0104	0.4970

*(Contd.)*

$z$	<i>Ordinate</i>	$Area = \int_0^z$	$z$	<i>Ordinate</i>	$Area = \int_0^z$
1.25	0.1827	0.3944	2.80	0.0091	0.4974
1.30	0.1714	0.4032	2.85	0.0079	0.4978
1.35	0.1604	0.4115	2.90	0.0060	0.4981
1.40	0.1497	0.4192	2.95	0.0051	0.4984
1.45	0.1394	0.4265	3.00	0.0044	0.4987
1.50	0.1295	0.4332	3.05	0.0038	0.4989

**t-Table**

$n$	<i>Probability</i>				
	0.9	0.1	0.05	0.02	0.01
1	0.158	6.314	12.706	31.821	63.657
2	0.142	2.920	4.303	6.965	9.925
3	0.137	2.353	3.182	4.541	5.841
4	0.134	2.132	2.776	3.747	4.604
5	0.132	2.015	2.571	3.365	4.032
6	0.131	1.943	2.447	3.143	3.707
7	0.130	1.895	2.365	2.998	3.496
8	0.130	1.860	2.306	2.896	3.355
9	0.129	1.833	2.262	2.821	3.250
10	0.129	1.812	2.228	2.764	3.169
11	0.129	1.796	2.201	2.718	3.106
12	0.128	1.782	2.179	2.681	3.055
13	0.128	1.771	2.160	2.650	3.012
14	0.128	1.761	2.145	2.624	2.977
15	0.128	1.753	2.131	2.602	2.947
16	0.128	1.746	2.120	2.583	2.921
17	0.128	1.740	2.110	2.567	2.898
18	0.127	1.734	2.101	2.552	2.878
19	0.127	1.729	2.093	2.539	2.861
20	0.127	1.725	2.086	2.528	2.845
21	0.127	1.721	2.080	2.518	2.831
22	0.127	1.717	2.074	2.508	2.819
23	0.127	1.714	2.069	2.500	2.807
24	0.127	1.711	2.064	2.492	2.797
25	0.127	1.708	2.060	2.485	2.787
30	0.127	1.697	2.042	2.457	2.750
40	0.126	1.684	2.021	2.423	2.704
60	0.126	1.671	2.000	2.390	2.660
120	0.126	1.658	1.980	2.358	2.617
$\infty$	0.126	1.645	1.960	2.326	2.576

$\chi^2$ -Table

n	Probability					
	0.99	0.95	0.10	0.05	0.02	0.01
1	0.000157	0.00393	2.706	3.841	5.412	6.635
2	0.0201	0.103	4.605	5.991	7.824	9.210
3	0.115	0.352	6.251	7.815	9.837	11.345
4	0.297	0.711	7.779	9.488	11.668	13.277
5	0.554	1.145	9.236	11.070	13.388	15.086
6	0.872	1.635	10.645	12.592	15.033	16.812
7	1.238	2.167	12.017	14.067	16.622	18.475
8	1.646	2.733	13.362	15.507	18.168	20.090
9	2.088	3.325	14.684	16.919	19.670	21.666
10	2.558	3.940	15.987	18.307	21.161	23.209
11	3.053	4.575	17.275	19.675	22.618	24.725
12	3.571	5.226	18.549	21.026	24.054	26.217
13	4.107	5.982	19.812	22.362	25.472	27.688
14	4.660	6.571	21.064	23.685	26.873	29.141
15	5.229	7.261	22.307	24.996	28.259	30.578
16	5.812	7.962	23.542	26.296	29.633	32.000
17	6.408	8.672	24.768	27.587	30.995	33.409
18	7.015	9.390	25.989	28.869	32.346	34.805
19	7.633	10.117	27.204	30.114	33.687	36.191
20	8.260	10.851	28.412	31.410	35.020	37.566
21	8.897	11.581	29.615	32.671	36.343	38.932
22	9.542	12.338	30.813	33.924	37.659	40.289
23	10.196	13.091	32.007	35.172	38.968	41.638
24	10.856	13.848	33.196	36.415	40.270	42.980
25	11.524	14.611	34.382	37.652	41.566	44.314
26	12.198	15.379	35.563	38.885	42.856	45.642
27	12.879	16.151	36.741	40.113	44.140	46.963
28	13.565	16.928	37.916	41.337	45.419	48.278
29	14.256	17.708	39.087	42.557	46.693	49.588
30	14.953	18.493	40.256	43.773	47.962	50.892

For larger values of  $n$ , the expression  $\sqrt{2\chi^2} - \sqrt{2n-1}$  may be used as a normal variate with unit variance.

**Table of F (Variance ratio)—1 Per Cent Points**

$v_2$	$v_1$									
	1	2	3	4	5	6	8	12	24	$\infty$
1	4052	4999	5403	5625	5764	5859	5981	6106	6234	6366
2	98.49	99.01	99.17	99.25	99.30	99.33	99.36	99.42	99.46	99.50
3	34.12	30.81	29.46	28.71	28.24	27.91	27.49	27.05	26.60	26.12
4	21.20	18.00	16.69	15.98	15.52	15.21	14.80	14.37	13.93	13.46
5	16.26	13.27	12.06	11.39	10.97	10.67	10.27	9.89	9.47	9.02
6	13.74	10.92	9.78	9.15	8.75	8.47	8.10	7.72	7.31	6.88
7	12.25	9.55	8.45	7.85	7.46	7.19	6.84	6.47	6.07	5.65
8	11.26	8.65	7.59	7.01	6.63	6.37	6.03	5.67	5.28	4.86
9	10.56	8.02	6.99	6.42	6.06	5.80	5.47	5.11	4.73	4.31
10	10.04	7.56	6.55	5.99	5.64	5.39	5.06	4.71	4.33	3.91
11	9.65	7.20	6.22	5.67	5.32	5.07	4.74	4.40	4.02	3.60
12	9.33	6.93	5.95	5.41	5.06	4.82	4.50	4.16	3.78	3.36
13	9.07	6.70	5.74	5.20	4.86	4.62	4.30	3.96	3.59	3.16
14	8.86	6.51	5.56	5.03	4.69	4.46	4.14	3.80	3.43	3.00
15	8.68	6.36	5.42	4.89	4.56	4.32	4.00	3.67	3.29	2.87
16	8.53	6.23	5.29	4.77	4.44	4.20	3.89	3.55	3.18	2.75
17	8.40	6.11	5.18	4.67	4.34	4.10	3.79	3.45	3.08	2.65
18	8.28	6.01	5.09	4.58	4.25	4.01	3.71	3.37	3.00	2.57
19	8.18	5.93	5.01	4.50	4.17	3.94	3.63	3.30	2.92	2.49
20	8.10	5.85	4.94	4.43	4.10	3.87	3.56	3.23	2.86	2.42
21	8.02	5.78	4.87	4.37	4.04	3.81	3.51	3.17	2.80	2.36
22	7.94	5.72	4.82	4.31	3.99	3.76	3.45	3.12	2.75	2.31
23	7.88	5.66	4.76	4.26	3.94	3.71	3.41	3.07	2.70	2.26
24	7.82	5.61	4.72	4.22	3.90	3.67	3.36	3.03	2.66	2.21
25	7.77	5.57	4.68	4.18	3.86	3.63	3.32	2.99	2.62	2.17
26	7.72	5.53	4.64	4.14	3.82	3.59	3.29	2.96	2.58	2.13
27	7.68	5.49	4.60	4.11	3.79	3.56	3.26	2.93	2.55	2.10
28	7.64	5.45	4.57	4.07	3.76	3.53	3.23	2.90	2.52	2.06
29	7.60	5.42	4.54	4.04	3.73	3.50	3.20	2.87	2.49	2.03
30	7.56	5.39	4.51	4.02	3.70	3.47	3.17	2.84	2.47	2.01
40	7.31	5.18	4.31	3.83	3.51	3.29	2.99	2.66	2.29	1.81
60	7.08	4.98	4.13	3.65	3.34	3.12	2.82	2.50	2.12	1.60
120	6.85	4.79	3.95	3.48	3.17	2.96	2.66	2.34	1.95	1.38
$\infty$	6.64	4.60	3.78	3.32	3.02	2.80	2.51	2.18	1.79	1.00

**Table of F—5 Per Cent Points**

$v_2$	$v_1$									
	1	2	3	4	5	6	8	12	24	$\infty$
1	161.4	199.5	215.7	224.6	230.2	234.0	238.9	243.9	249.0	253.4
2	18.51	19.00	19.16	19.25	19.30	19.33	19.37	19.41	19.45	19.50
3	10.13	9.55	9.28	9.12	9.01	8.94	8.84	8.74	8.64	8.53
4	7.71	6.94	6.59	6.39	6.26	6.16	6.04	5.91	5.77	5.63
5	6.61	5.79	5.41	5.19	5.05	4.95	4.82	4.68	4.53	4.36
6	5.99	5.14	4.76	4.53	4.39	4.28	4.15	4.00	3.84	3.67
7	5.59	4.74	4.35	4.12	3.97	3.87	3.73	3.57	3.41	3.23
8	5.32	4.46	4.07	3.84	3.69	3.58	3.44	3.28	3.12	2.93
9	5.12	4.26	3.86	3.63	3.48	3.37	3.23	3.07	2.90	2.71
10	4.96	4.10	3.71	3.48	3.33	3.22	3.07	2.91	2.74	2.54
11	4.84	3.98	3.59	3.36	3.20	3.09	2.95	2.79	2.61	2.40
12	4.75	3.88	3.49	3.26	3.11	3.00	2.85	2.69	2.50	2.30
13	4.67	3.80	3.41	3.18	3.02	2.92	2.77	2.60	2.42	2.21
14	4.60	3.74	3.34	3.11	2.96	2.85	2.70	2.53	2.35	2.13
15	4.54	3.68	3.29	3.06	2.90	2.79	2.64	2.48	2.29	2.07
16	4.49	3.63	3.24	3.01	2.85	2.74	2.59	2.42	2.24	2.01
17	4.45	3.59	3.20	2.96	2.81	2.70	2.55	2.38	2.19	1.96
18	4.41	3.55	3.16	2.93	2.77	2.66	2.51	2.34	2.15	1.92
19	4.38	3.52	3.13	2.90	2.74	2.63	2.48	2.31	2.11	1.88
20	4.35	3.49	3.10	2.87	2.71	2.60	2.45	2.28	2.08	1.84
21	4.32	3.47	3.07	2.84	2.68	2.57	2.42	2.25	2.05	1.81
22	4.30	3.44	3.05	2.82	2.66	2.55	2.40	2.23	2.03	1.78
23	4.28	3.42	3.03	2.80	2.64	2.53	2.38	2.20	2.00	1.76
24	4.26	3.40	3.10	2.78	2.62	2.51	2.36	2.18	1.98	1.73
25	4.24	3.38	2.99	2.76	2.60	2.49	2.34	2.16	1.96	1.71
26	4.22	3.37	2.98	2.74	2.59	2.47	2.32	2.15	1.95	1.69
27	4.21	3.35	2.96	2.73	2.57	2.46	2.30	2.13	1.93	1.67
28	4.20	3.34	2.95	2.71	2.56	2.44	2.29	2.12	1.91	1.65
29	4.18	3.33	2.93	2.70	2.54	2.43	2.28	2.10	1.90	1.64
30	4.17	3.32	2.92	2.69	2.53	2.42	2.27	2.09	1.89	1.62
40	4.08	3.23	2.84	2.61	2.45	2.34	2.18	2.00	1.79	1.51
60	4.00	3.15	2.76	2.52	2.37	2.25	2.10	1.92	1.70	1.39
120	3.92	3.07	2.68	2.45	2.29	2.17	2.02	1.83	1.61	1.25
$\infty$	3.84	2.99	2.60	2.37	2.21	2.09	1.94	1.75	1.52	1.00



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# Solved Question Papers

**BE/BTech Degree Examination, April/May 2011**

**Fourth Semester**

**Electronics and Communication Engineering**

**MA2261—Probability and Random Process**

**(Common to Bio-Medical Engineering)**

**(Regulation 2008)**

**Time: Three hours**

**Maximum: 100 marks**

**Answer ALL Questions**

**Part-A (10 × 2 = 20 marks)**

1. The CDF of a continuous random variable is given by  $F(x) = \begin{cases} 0, & x < 0 \\ 1 - e^{-x/5}, & 0 \leq x < \infty \end{cases}$ . Find the PDF and mean of  $X$ .
2. Establish the memoryless property of the exponential density function.
3. Let  $X$  and  $Y$  be continuous random variables with joint probability density function  $f_{XY}(x, y) = \frac{x(x-y)}{8}$ ,  $0 < x < 2$ ,  $-x < y < x$  and  $f_{XY}(x, y) = 0$  elsewhere. Find  $f_{Y|X}(y/x)$ .
4. Find the acute angle between the two lines of regression, assuming the two lines of regression.
5. Prove that a first-order stationary process has a constant mean.
6. State the postulates of a Poisson process.
7. The autocorrelation function of a stationary random process is  $R(\tau) = 16 + \frac{9}{1+16\tau^2}$ . Find the mean and variance of the process.
8. Prove that for a WSS process  $\{X(t)\}$ ,  $R_{XX}(t, t + \tau)$  is an even function of  $\tau$ .
9. Find the system transfer function, if a linear time invariant system has an impulse function
$$H(t) = \begin{cases} \frac{1}{2c}; & |\tau| \leq c \\ 0; & |\tau| \geq c \end{cases}$$
10. Define white noise.

**Part-B (5 × 16 = 80 marks)**

11. (a) The probability density function of a random variable  $X$  is given by

$$f_X(x) = \begin{cases} x, & 0 < x < 1 \\ k(2-x), & 1 \leq x \leq 2 \\ 0, & \text{otherwise} \end{cases}$$

- (1) Find the value of ' $k$ '. (4)
- (2) Find  $P(0.2 < x < 1.2)$  (3)
- (3) What is  $P[0.5 < x < 1.5/x \geq 1]$  (4)
- (4) Find the distribution function of  $f(x)$ . (5)

Or

- (b) (i) Derive the mgf of a Poisson distribution and, hence, or otherwise deduce its mean and variance.
- (ii) The marks obtained by a number of students in a certain subject are assumed to be normally distributed with a mean of 65 and a standard deviation of 5. If 3 students are selected at random from this set, what is the probability that exactly 2 of them will have marks over 70?
12. (a) (i) If  $X$  and  $Y$  are independent Poisson random variables with respective parameters  $\lambda_1$  and  $\lambda_2$ , calculate the conditional distribution of  $X$ , given that  $X + Y = n$ .
- (ii) The regression equation of  $X$  on  $Y$  is  $3Y - 5X + 108 = 0$ . The mean value of  $Y$  is 44 and the variance of  $X$  is  $9/16^{\text{th}}$  of the variance of  $Y$ . Find the mean value of  $X$  and the correlation coefficient.

Or

- (b) (i) If  $X$  and  $Y$  are independent random variables with density function
- $$f_X(x) = \begin{cases} 1, & 1 \leq x \leq 2 \\ 0, & \text{otherwise} \end{cases} \quad \text{and} \quad f_Y(y) = \begin{cases} \frac{y}{6}, & 2 \leq y \leq 4 \\ 0, & \text{otherwise} \end{cases}, \quad \text{find the density function of } Z = XY.$$
- (ii) The lifetime of a particular variety of electric bulb may be considered as a random variable with a mean of 1200 hours and a standard deviation of 250 hours. Using the central limit theorem, find the probability that the average life of 60 bulbs exceeds 1250 hours.
13. (a) (i) A random process  $X(t)$  defined by  $X(t) = A \cos t + B \sin t$ ,  $-\infty < t < \infty$ , where  $A$  and  $B$  are independent random variables each of which takes a value of  $-2$  with a probability of  $1/3$  and a value of  $1$  with a probability of  $2/3$ . Show that  $X(t)$  is wide-sense stationary.
- (ii) A random process has sample functions of the form  $X(t) = A \cos(\omega t + \theta)$ , where  $\omega$  is constant,  $A$  is a random variable with mean zero and variance one and  $\theta$  is a random variable that is uniformly distributed between  $0$  and  $2\pi$ . Assume that the random variables  $A$  and  $\theta$  are independent. Is  $X(t)$  a mean-ergodic process?

Or

- (b) (i) If  $\{X(t)\}$  is a Gaussian process with  $\mu(t) = 10$  and  $C(t_1, t_2) = 16e^{-|t_1-t_2|}$ , find the probability that
- $X(10) \leq 8$
  - $|X(10) - X(6)| \leq 4$
- (ii) Prove that the interval between two successive occurrences of a Poisson process with parameter  $\lambda$  has an exponential distribution with mean  $\frac{1}{\lambda}$ .
14. (a) (i) The power spectral density function of a zero mean WSS process  $X(t)$  is given by  $S(\omega) = \begin{cases} 1, & |\omega| < \omega_0 \\ 0, & \text{otherwise} \end{cases}$ . Find  $R(\tau)$  and show that  $X(t)$  and  $X\left(t + \frac{\pi}{\omega_0}\right)$  are uncorrelated.
- (ii) The autocorrelation of a WSS process is given by  $R(\tau) = \alpha^2 e^{-2\lambda|\tau|}$  determine the power spectral density of the process.
- Or
- (b) (i) State and prove the Weiner–Khintchine theorem.
- (ii) The cross-power spectrum of a real random processes  $\{X(t)\}$  and  $\{Y(t)\}$  is given by  $S_{xy}(\omega) = \begin{cases} a + bj\omega, & \text{for } |\omega| < 1 \\ 0, & \text{elsewhere} \end{cases}$ . Find the cross correlation function.
15. (a) (i) Consider a system with transfer function  $\frac{1}{1+j\omega}$ . An input signal with autocorrelation function  $m\delta(\tau) + m^2$  is fed as input to the system. Find the mean and mean-square value of the output.
- (ii) A stationary random process  $X(t)$  having the autocorrelation function  $R_{XX}(\tau) = A\delta(\tau)$  is applied to a linear system at time  $t = 0$  where  $f(\tau)$  represents the impulse function. The linear system has the impulse response of  $h(t) = e^{-bt} u(t)$  where  $u(t)$  represents the unit step function. Find  $R_{YY}(\tau)$ . Also find the mean and variance of  $Y(t)$ .
- Or
- (b) (i) A linear system is described by the impulse response  $h(t) = \frac{1}{RC} e^{\frac{-t}{RC}} u(t)$ . Assume an input process whose autocorrelation function is  $B\delta(\tau)$ . Find the mean and autocorrelation functions of the output process.

(ii) If  $\{N(t)\}$  is a bandlimited white noise centred at a carrier frequency

$$\omega_0 \text{ such that } S_{NN}(\omega) = \begin{cases} \frac{N_0}{2}, & \text{for } |\omega - \omega_0| < \omega_B \\ 0, & \text{elsewhere} \end{cases}$$

Find the autocorrelation of  $\{N(t)\}$ .

## Solutions

### Part-A

$$1. \quad F(x) = \begin{cases} 0, & \text{in } x < 0 \\ 1 - e^{-x/5}, & \text{in } 0 \leq x < \infty \end{cases}$$

$$\therefore f(x) = F'(x) = \begin{cases} 0, & \text{in } x < 0 \\ \frac{1}{5}e^{-x/5}, & \text{in } 0 \leq x < \infty \end{cases}$$

$$E(x) = \int_0^{\infty} x f(x) dx = \frac{1}{5} \left[ x \frac{e^{-x/5}}{-1/5} - \frac{e^{-x/5}}{1/25} \right]_0^{\infty} = \frac{1}{5} \times 25 = 5$$

2. Standard property, the derivation is available in the book.
3. This is a worked example in the book. ‘c’ given in the book must be taken as  $\frac{1}{8}$ .
4. This is a worked example in the book.
5. This property is proved in the note under ‘stationarity’.
6. Postulates of a Poisson process are available in the book.

$$7. \quad R_{XX}(\tau) = 16 + \frac{9}{1 + 16\tau^2}$$

$$\mu_X^2 = \lim_{\tau \rightarrow \infty} [R_{XX}(\tau)] = 16 \quad \therefore \mu_X = 4$$

$$E[X^2(t)] = R_{XX}(0) = 25$$

$$\therefore V\{X(t)\} = E\{X^2(t)\} - E^2\{X(t)\} = 25 - 16 = 9$$

8.  $R_{XX}(\tau) = E\{X(t) \cdot X(t + \tau)\} \quad \therefore R_{XX}(-\tau) = E\{X(t) \cdot X(t - \tau)\} = R_{XX}(\tau)$   
 $\therefore R_{XX}(\tau)$  is an even function of  $\tau$ .

$$9. \quad \bar{H}(\omega) = F\{H(t)\} = \int_{-\infty}^{\infty} H(t) e^{-i\omega t} dt = \int_{-c}^c \frac{1}{2c} e^{-i\omega t} dt$$

$$= \frac{1}{c} \int_0^c \cos \omega t dt = \frac{1}{c\omega} [\sin \omega t]_0^c = \frac{\sin c\omega}{c\omega}$$

10. Definition of ‘White Noise’ is available in the book.

**Part-B**

11. (a) (1)  $\int_{R_X} f(x) dx = 1$ ; viz.,  $\int_0^1 x dx + k \int_1^2 (2-x) dx = 1$

$$\text{i.e., } \left[ \frac{x^2}{2} \right]_0^1 + k \left[ 2x - \frac{x^2}{2} \right]_1^2 = 1$$

$$\text{i.e., } \frac{1}{2} + k \left( 2 - \frac{3}{2} \right) = 1 \quad \therefore \quad k = 1$$

$$\begin{aligned} (2) \quad P(0.2 < x < 1.2) &= \int_{0.2}^1 f(x) dx + \int_1^{1.2} f(x) dx = \int_{0.2}^1 x dx + \int_1^{1.2} (2-x) dx \\ &= \left( \frac{x^2}{2} \right)_{0.2}^1 + \left( 2x - \frac{x^2}{2} \right)_1^{1.2} = \frac{1}{2}(1 - 0.04) + \frac{1}{2}(3.36 - 3) \\ &= \frac{1}{2} \times 1.32 = 0.66 \end{aligned}$$

$$(3) \quad P\{0.5 < x < 1.5/x \geq 1\} = \frac{P\{(0.5 < x < 1.5) \cap (x \geq 1)\}}{P(x \geq 1)}$$

$$\begin{aligned} &= \frac{P(1 \leq X \leq 1.5)}{P(1 \leq X \leq 2)} = \int_1^{1.5} (2-x) dx \div \int_1^2 (2-x) dx \\ &= \left( 2x - \frac{x^2}{2} \right)_1^{1.5} \div \left( 2x - \frac{x^2}{2} \right)_1^2 = \frac{1.875 - 1.5}{2 - 1.5} = \frac{0.375}{0.5} = 0.75 \end{aligned}$$

(4) Distribution function of  $X$  [not  $f(x)$  as given in the question]

$$F(x) = P(X \leq x) = 0, \text{ when } x < 0$$

$$= \int_0^x f(x) dx = \int_0^x x dx = \frac{x^2}{2}, \text{ when } 0 \leq x \leq 1$$

$$= \int_0^1 x dx + \int_1^x (2-x) dx$$

$$= \frac{1}{2} + \left( 2x - \frac{x^2}{2} \right)_1^x = 2x - \frac{x^2}{2} - 1, \text{ when } 1 \leq x \leq 2$$

$$= \int_0^1 x dx + \int_1^2 (2-x) dx + \int_2^x 0 \cdot dx, \text{ when } x \geq 2$$

$$= 1, \text{ when } x \geq 2$$

$$\begin{aligned}
 11. \text{ (b) (i)} \quad M_X(t) &= \sum_{r=0}^{\infty} e^{tr} e^{-\lambda} \lambda^r / r! = e^{-\lambda} \sum_{r=0}^{\infty} (\lambda e^t)^r / r! = e^{-\lambda} e^{\lambda e^t} \\
 &= e^{-\lambda(1-e^t)} \\
 M_X(t) &= e^{\lambda \left( \frac{t}{1} + \frac{t^2}{2} + \dots \right)} \\
 &= 1 + \frac{\lambda}{1} \left( t + \frac{t^2}{2} + \dots \right) + \frac{\lambda^2 t^2}{2} \left( 1 + \frac{t}{2} + \dots \right)^2 + \dots \quad (1)
 \end{aligned}$$

$$E(X) = \text{coefficient of } \frac{t}{1} \text{ in (1)} = \lambda$$

$$E(X^2) = \text{coefficient of } \frac{t^2}{2} \text{ in (1)} = \lambda + \lambda^2;$$

$$V(X) = E(X^2) - E^2(X) = \lambda$$

- (ii) Let  $X$  represent the marks obtained by the students in the subject. Then  $X$  follows a  $N(65, 5)$

$$\begin{aligned}
 P(\text{a student scores over 70}) &= P(X > 70) = P\left\{\frac{X - 65}{5} > \frac{70 - 65}{5}\right\} \\
 &= P(z > 1) = P(1 < z < \infty) = 0.5 - P(0 < z < 1) \\
 &= 0.5 - 0.3413 = 0.1587
 \end{aligned}$$

Let  $p = P(\text{a student scores above 70}) = 0.1587$  and  $q = 0.8413$ ;  $n = 3$

Since  $p$  is the same for all the students, the number  $Y$  of (successes) students scoring more than 70, follows a binomial distribution.

$$\begin{aligned}
 P(\text{exactly 2 students score more than 70}) &= P(2 \text{ success}) \\
 &= "nC_2 p^r q^{n-r}" = 3C_2 \times (0.1587)^2 \times (0.8413)^1 = 0.0636
 \end{aligned}$$

12. (a) (i) This is a worked example in the book.

$$\text{(ii) Regression line of } X \text{ on } Y \text{ is } 3y - 5x + 108 = 0, \text{ i.e., } x = \frac{3}{5}y + \frac{108}{5} \quad (1)$$

Standard formula for the regression line of  $X$  on  $Y$  is

$$x - \bar{x} = b_{XY}(y - \bar{y})$$

$$\text{i.e., } x = b_{XY}y + (\bar{x} - b_{XY}\bar{y}) \quad (2)$$

$$\text{Comparing (1) and (2); } b_{XY} = \frac{3}{5} \text{ and } \bar{x} - \frac{3}{5} \times 44 = \frac{108}{5}$$

$$\therefore \bar{x} = E(X) = \frac{132}{5} + \frac{108}{5} = 48$$

$$\frac{b_{XY}}{b_{YX}} = \frac{\sigma_X^2}{\sigma_Y^2} = \frac{9}{16} \quad \therefore b_{YX} = \frac{16}{9} \times b_{XY} = \frac{16}{9} \times \frac{3}{5} = \frac{16}{15}$$

$$r^2 = b_{XY}, b_{YX} = \frac{3}{5} \times \frac{16}{15} = \frac{16}{25} \quad \therefore r_{XY} = \frac{4}{5}$$

12. (b) (i)  $Z = XY$ . Let the auxiliary RV be  $W = Y$

$$\therefore x = \frac{z}{w} \text{ and } y = w; J = \begin{vmatrix} 1 & -\frac{z}{w^2} \\ \frac{1}{w} & 1 \end{vmatrix} = \frac{1}{w}$$

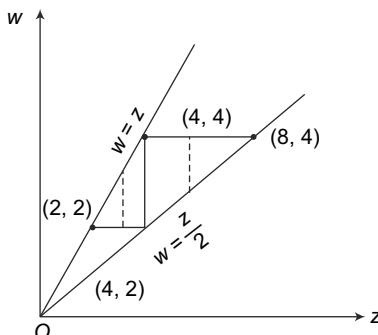
$\therefore$  The joint pdf of  $(Z, W)$  is given by  $f_{ZW}(z, w) = \frac{1}{|w|} f_{XY}(x, y)$

$$\text{i.e., } f_{ZW}(z, w) = \frac{1}{w} \cdot \frac{y}{6} \text{ or } \frac{1}{w} \cdot \frac{w}{6} = \frac{1}{6}, \text{ since } f_X(x, y) = 1 \times \frac{y}{6}, \text{ as } X$$

and  $Y$  are independent.

Range space of  $(X, Y)$  is given by  $1 \leq x \leq 2$  and  $2 \leq y \leq 4$

$\therefore$  Range space of  $(Z, W)$  is given by  $w \leq z \leq 2w$  and  $2 \leq w \leq 4$



**Fig. 1**

The pdf of  $Z$  is given by

$$f_Z(z) = \begin{cases} \int_2^z \frac{1}{6} dw, & \text{if } 2 \leq z \leq 4 \\ \int_{z/2}^4 \frac{1}{6} dw, & \text{if } 4 \leq z \leq 8 \end{cases}$$

$$\text{i.e., } f_Z(z) = \begin{cases} \frac{1}{6}(z - 2), & \text{if } 2 \leq z \leq 4 \\ \frac{1}{12}(8 - z) & \text{if } 4 \leq z \leq 8 \end{cases}$$

(ii) This is a worked example in the book.

13. (a) (i)  $E\{X(t)\} = \cos t \cdot E(A) + \sin t \cdot E(B)$ , where the probability distribution of  $A$  or  $B$  is given below:

$$\therefore E(A) = E(B) = -\frac{2}{3} + \frac{2}{3} = 0$$

$$\text{and } E(A^2) = E(B^2) = \frac{4}{3} + \frac{2}{3} = 2$$

$E\{X(t)\} = 0$ and	<table border="1" style="width: 100%; border-collapse: collapse;"> <tr> <td style="padding: 5px; text-align: center;"><math>A</math> or <math>B</math></td><td style="padding: 5px; text-align: center;">-2</td><td style="padding: 5px; text-align: center;">1</td></tr> <tr> <td style="padding: 5px; text-align: center;"><math>P(A)</math> or <math>P(B)</math></td><td style="padding: 5px; text-align: center;"><math>\frac{1}{3}</math></td><td style="padding: 5px; text-align: center;"><math>\frac{2}{3}</math></td></tr> </table>	$A$ or $B$	-2	1	$P(A)$ or $P(B)$	$\frac{1}{3}$	$\frac{2}{3}$	$E(AB) = E(A) \cdot E(B) =$
$A$ or $B$	-2	1						
$P(A)$ or $P(B)$	$\frac{1}{3}$	$\frac{2}{3}$						

0, since  $A$  and  $B$  are independent.

$$E\{X^2(t)\} = \cos^2 t \cdot E(A^2) + \sin^2 t \cdot E(B^2) + 2 \sin t \cos t E(AB) = 2$$

$$\therefore V\{X(t)\} = 2$$

Since  $E\{X(t)\}$  and  $V\{X(t)\}$  are constants  $\{X(t)\}$  is a WSS process.

$$\begin{aligned}
 \text{(ii)} \quad & \text{Let } \bar{X}_T = \frac{1}{2T} \int_{-T}^T X(t) dt = \frac{1}{2T} \int_{-T}^T A \cos(\omega t + \theta) dt \\
 & E(\bar{X}_T) = \frac{1}{2\pi} \int_{-T}^T E(A) \cdot E\{\cos(\omega t + \theta)\} dt = 0, \text{ since } E(A) = 0 \\
 R(t_1, t_2) &= E\{X(t_1) \times (t_2)\} = E\{A^2 \cos(\omega t_1 + \theta) \cdot \cos(\omega t_2 + \theta)\} \\
 &= E(A^2) E\left[\frac{1}{2} \{\cos \omega(t_1 + t_2) + 2\theta\} + \cos \omega(t_1 - t_2)\right] \\
 &= \frac{1}{2} \int_0^{2\pi} \frac{1}{2\pi} \{\cos \omega(t_1 + t_2) + 2\theta\} + \cos \omega(t_1 - t_2) d\theta, \\
 & \qquad \qquad \qquad \text{since } E(A^2) = 1 \\
 &= \frac{1}{2} \cos \omega(t_1 - t_2)
 \end{aligned}$$

$$\bar{X}_T^2 = \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T X(t_1) X(t_2) dt_1 dt_2$$

$$\begin{aligned}
 \therefore E\{\bar{X}_T^2\} &= \frac{1}{4T^2} \int_{-T}^T \int_{-T}^T R(t_1, t_2) dt_1 dt_2 = \frac{1}{4T^2} \int_{-T}^T \int_0^T \cos \omega(t_1 - t_2) dt_1 dt_2 \\
 &= \frac{1}{4T^2} \int_{-T}^T \frac{1}{\omega} [\sin \omega(t_1 - t_2)]_{t_1=0}^{t_1=T} dt_2 \\
 &= \frac{1}{4\omega T^2} \int_{-T}^T \{\sin \omega(T - t_2) + \sin t_2\} dt_2 \\
 &= 0 \quad \therefore \text{Var}\{\bar{X}_T\} = 0 \quad \therefore \lim_{T \rightarrow \infty} \text{Var}\{\bar{X}_T\} = 0
 \end{aligned}$$

- ∴ by Mean-Ergodic theorem,  $\{X(t)\}$  is a mean-ergodic process.
13. (b) (i) This is a worked example in the book.  
 (ii) This is a standard property, the proof of which is available in the book.
14. (a) (i) This is a worked example in the book.  
 (ii) This is a worked example in the book.
14. (b) (i) Statement and proof of Weiner–Khinchine theorem is available in the book.

$$(ii) R_{XY}(\tau) = \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{XY}(\omega) e^{i\tau\omega} d\omega = \frac{1}{2\pi} \int_{-1}^1 (a + ib\omega) e^{i\tau\omega} d\omega$$

$$= \frac{1}{\pi} \left[ \int_0^1 a \cos \tau\omega d\omega - \int_0^1 b\omega \sin \tau\omega d\omega \right]$$

$$= \frac{1}{\pi} \left[ a \left( \frac{\sin \tau\omega}{\tau} \right)_0^1 - b \left\{ \omega \left( \frac{-\cos \tau\omega}{\tau} \right) + \frac{\sin \tau\omega}{\tau^2} \right\}_0^1 \right]$$

$$= \frac{1}{\pi} \left[ a \frac{\sin \tau}{\tau} - b \left\{ -\frac{\cos \tau}{\tau} + \frac{\sin \tau}{\tau^2} \right\} \right]$$

$$= \frac{1}{\pi} \left[ \frac{a}{\tau} \sin \tau + \frac{b}{\tau} \cos \tau - \frac{b}{\tau^2} \sin \tau \omega \right]$$

$$15. (a) (i) R_{XX}(\tau) = m\delta(\tau) + m^2; \quad \therefore S_{XX}(\omega) = \int_{-\infty}^{\infty} m\delta(\tau) e^{-i\omega\tau} d\tau + F(m^2)$$

$$= m + 2\pi m^2 \delta(\omega)$$

$$S_{YY}(\omega) = S_{XX}(\omega) \cdot |H(\omega)|^2 = [m + 2\pi m^2 \delta(\omega)] \cdot \frac{1}{1 + \omega^2}$$

$$\therefore R_{YY}(\tau) = F^{-1}\{S_{YY}(\omega)\}$$

$$= \frac{1}{2\pi} \int_{-\infty}^{\infty} \frac{m}{1 + \omega^2} e^{i\tau\omega} d\omega + m^2 \int_{-\infty}^{\infty} \frac{1}{1 + \omega^2} \delta(\omega) e^{i\tau\omega} d\omega$$

$$= \frac{m}{2\pi} \pi e^{-\tau} + m^2 \cdot 1 = \frac{m}{2} e^{-\tau} + m^2 \quad \left[ \because \int_{-\infty}^{-\infty} \frac{e^{iaz}}{z^2 + 1} dz = \pi e^{-a} \right]$$

$$E^2\{Y(t)\} = \lim_{\tau \rightarrow \infty} [R_{YY}(\tau)] = m^2; \quad \therefore E\{Y(t)\} = m$$

$$E\{Y^2(t)\} = R_{YY}(0) = \frac{m}{2} + m^2$$

$$(ii) \quad S_{XX}(\omega) = F\{R_{XX}(\tau)\} = \int_{-\infty}^{\infty} A\delta(\tau)e^{-i\omega\tau} d\tau = A$$

$$\begin{aligned} H(\omega) &= F\{h(\tau)\} = \int_{-\infty}^{\infty} e^{-bt} U(t) e^{-i\omega\tau} dt = \int_0^{\infty} e^{-(b+i\omega)t} dt \\ &= \left[ \frac{e^{-(b+i\omega)t}}{-(b+i\omega)} \right]_0^{\infty} = \frac{1}{b+i\omega} \end{aligned}$$

$$S_{YY}(\omega) = S_{XX}(\omega) \cdot |H(\omega)|^2 = \frac{A}{\omega^2 + b^2}$$

$$\begin{aligned} \therefore R_{YY}(\tau) &= F^{-1}\{S_{YY}(\omega)\} = \frac{1}{2\pi} \int_{-\infty}^{\infty} \frac{A}{\omega^2 + b^2} e^{i\tau\omega} d\omega \\ &= \frac{A}{2\pi} \cdot \frac{\pi}{b} e^{-b\tau} = \frac{A}{2b} e^{-b\tau} \end{aligned}$$

$$\mu_Y^2 = \lim_{\tau \rightarrow \infty} R_{YY}(\tau) = 0$$

$$E\{Y^2(t)\} = R_{YY}(0) = \frac{A}{2b}$$

$$\therefore V\{Y(t)\} = E\{Y^2(t)\} - E^2\{Y(t)\} = \frac{A}{2b}$$

$$15. \quad (b) \quad (i) \quad S_{XX}(\omega) = F\{R_{XX}(\omega)\} = \int_{-\infty}^{\infty} B\delta(\tau) e^{-i\omega\tau} d\tau = B$$

$$\begin{aligned} H(\omega) &= F\{h(t)\} = \int_{-\infty}^{\infty} \frac{1}{RC} e^{-t/RC} U(t) e^{-i\omega t} dt = \frac{1}{RC} \int_0^{\infty} e^{-t/RC} e^{-i\omega t} dt \\ &= \frac{1}{RC} \left[ \frac{e^{-\left(\frac{1}{RC}+i\omega\right)t}}{-\left(\frac{1}{RC}+i\omega\right)} \right]_0^{\infty} = \frac{1}{1+iRC\omega} \end{aligned}$$

$$S_{YY}(\omega) = S_{XX}(\omega) \cdot |H(\omega)|^2 = \frac{B}{1+R^2C^2\omega^2}$$

$$\therefore R_{YY}(\tau) = F^{-1}\left\{ \frac{B}{1+R^2C^2\omega^2} \right\} = \frac{B}{2\pi R^2 C^2} \int_{-\infty}^{\infty} \frac{e^{i\tau\omega}}{\left(\frac{1}{RC}\right)^2 + \omega^2} d\omega$$

$$= \frac{B}{2\pi R^2 C^2}, \frac{\pi}{(1/RC)} e^{-\tau/RC} = \frac{B}{2RC} e^{-\tau/RC}$$

$$\mu_Y^2 = \lim_{\tau \rightarrow \infty} [R_{YY}(\tau)] = 0 \quad \therefore \quad \mu_Y = 0$$

(ii)  $R_{NN}(\tau) = F^{-1}\{S_{NN}(\omega)\}$

$$\begin{aligned} &= \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{NN}(\omega) e^{i\tau\omega} d\omega = \frac{N_0}{4\pi} \int_{\omega_0 - \omega_B}^{\omega_0 + \omega_B} e^{i\tau\omega} d\omega \\ &= \frac{N_0}{4\pi} \left[ \frac{e^{i\tau\omega}}{i\tau} \right]_{\omega_0 - \omega_B}^{\omega_0 + \omega_B} = \frac{N_0}{4\pi i\tau} \{e^{-i\tau(\omega_0 + \omega_B)} - e^{i\tau(\omega_0 - \omega_B)}\} \\ &= \frac{N_0 e^{i\tau\omega_0}}{2\pi\tau} \left\{ \frac{e^{i\tau\omega_B} - e^{-i\tau\omega_B}}{2i} \right\} \\ &= \frac{N_0}{2\pi\tau} e^{i\tau\omega_0} \sin \tau\omega_B \end{aligned}$$

**BE/BTech Degree Examination, April/May 2011**

**Fourth Semester**

**Computer Science and Engineering**

**MA2262—Probability and Queueing Theory**

**(Common to BTech Information Technology)**

**(Regulation 2008)**

**Time: Three hours**

**Maximum: 100 marks**

**Answer ALL Questions**

**Part-A (10 × 2 = 20 marks)**

1. The cumulative distribution function of the random variable X is given by

$$F_X(x) = \begin{cases} 0, & x < 0 \\ x + \frac{1}{2}, & 0 \leq x \leq \frac{1}{2} \\ 1, & x > \frac{1}{2} \end{cases} \text{. Compute } P[X > 1/4].$$

2. Let the random variable X denote the sum obtained in rolling a pair of fair dice. Determine the probability mass function of X.
3. Given the two regression lines  $3X + 12Y = 19$ ,  $3Y + 9X = 46$ , find the coefficient of correlation between X and Y.
4. State the Central Limit theorem.
5. Define:
  - (a) Continuous-time random process
  - (b) Discrete state random process
6. Find the transition probability matrix of the process represented by the state transition diagram.
7. Arrival rate of telephone calls at a telephone booth is according to Poisson distribution with an average time of 9 minutes between two consecutive arrivals. The length of a telephone call is assumed to be exponentially distributed with a mean of 3 minutes. Determine the probability that a person arriving at the booth will have to wait.
8. Trains arrive at the yard every 15 minutes and the service time is 33 minutes. If the line capacity of the yard is limited to 4 trains, find the probability that the yard is empty.

9. Given that the service time is Erlang with parameters  $m$  and  $\mu$ . Show that the Pollaczek–Khintchine formula reduces to  $L_s = mp + \frac{m(1+m)\rho^2}{2(1-m\rho)}$ .
10. Give any two examples for series queueing situations.

**Part-B (5 × 16 = 80 marks)**

11. (a) (i) Find the moment-generating function of the binomial random variable with parameters  $m$  and  $p$  and, hence, find its means and variance. (10)  
(ii) Define Weibull distribution and write its mean and variance. (6)
- Or
- (b) (i) Derive mean and variance of a geometric distribution. Also establish the forgetfulness property of the geometric distribution.  
(ii) Suppose that telephones calls arriving at a particular switchboard follow a Poisson process with an average of 5 calls coming per minute. What is the probability that up to a minute will elapse until 2 calls have come in to the switchboard?

12. (a) Given the joint density function  $f(x, y) = \begin{cases} x\frac{(1+3y^2)}{4}, & 0 < x < 2, 0 < y < 1 \\ 0, & \text{elsewhere} \end{cases}$ .

Find the marginal densities  $g(x)$ ,  $h(y)$  and the conditional density  $f(x|y)$  and evaluate  $P\left[\frac{1}{4} < x < \frac{1}{2} / Y = \frac{1}{3}\right]$ .

Or

- (b) (i) Determine whether the random variables  $X$  and  $Y$  are independent, given their joint probability density function as  

$$f(x, y) = \begin{cases} x^2 + \frac{xy}{3}, & 0 \leq x \leq 1, 0 \leq y \leq 2 \\ 0, & \text{otherwise} \end{cases}$$
.
- (ii) If  $X$  and  $Y$  are independent random variables having density functions  $f(x) = \begin{cases} 2e^{-2x}, & x \geq 0 \\ 0, & x < 0 \end{cases}$  and  $f(y) = \begin{cases} 3e^{-3y}, & y \geq 0 \\ 0, & y < 0 \end{cases}$ , respectively, find the density functions of  $z = X - Y$ .
13. (a) (i) Show that random process  $\{X(t)\} = A \cos t + B \sin t$ ,  $-\infty < t < \infty$  is a wide sense stationary process where  $A$  and  $B$  are independent random variables, each of which has a value of  $-2$  with a probability of  $\frac{1}{3}$  and a value of  $1$  with a probability of  $\frac{2}{3}$ .  
(ii) Derive probability distribution of Poisson process and, hence, find its autocorrelation function.

Or

- (b) (i) Find the limiting-state probabilities associated with the following transition probability matrix
- $$\begin{bmatrix} 0.4 & 0.5 & 0.1 \\ 0.3 & 0.3 & 0.4 \\ 0.3 & 0.2 & 0.5 \end{bmatrix}.$$
- (ii) Show that the difference of two independent Poisson processes is not a Poisson process.
14. (a) (i) Customers arrive at a one window drive-in bank according to Poisson distribution with a mean of 10 per hour. Service time per customer is exponential with a mean of 5 minutes. The space in front of the window, including that for the serviced car can accommodate a maximum of three cars. Other cars can wait outside this space.
- (1) What is the probability that an arriving customer can drive directly to the space in front of the window?
- (2) What is the probability that an arriving customer will have to wait outside the indicated space?
- (3) How long is an arriving customer expected to wait before being served? (10)
- (ii) Show that for the (M/M/1): (FCFS/ $\infty/\infty$ ), the distribution of waiting time in the system is  $w(t) = (\mu - \lambda)e^{-(\mu - \lambda)t}$ ,  $t > 0$  (6)
- Or
- (b) Find the steady-state solution for the multiserver M/M/C model and hence find  $L_9$ ,  $W_9$ ,  $W_s$  and  $L_s$  by using Little formula.
15. (a) Derive the expected steady-state size for the single server queues with Poisson input and General service. (16)
- Or
- (b) Write short notes on
- (i) Series queues (8)
- (ii) Open and closed queue networks (8)

## Solutions

### Part-A

1.  $P(X > \frac{1}{4}) = P\left(\frac{1}{4} < X < \infty\right) = F(\infty) - F\left(\frac{1}{4}\right) = 1 - \left(\frac{1}{4} + \frac{1}{2}\right) = \frac{1}{4}$

2.  $X$  takes the values 2, 3, 4, ...12.

$$P(X = 2) = P(1 \text{ from } D_1 \text{ and } 1 \text{ from } D_2) = \frac{1}{6} \cdot \frac{1}{6} = \frac{1}{36}$$

$$P(X = 3) = P\{1 \text{ from } D_1 \text{ and } 2 \text{ from } D_2 \text{ or } 2 \text{ from } D_1 \text{ and } 1 \text{ from } D_2\} = \frac{2}{36}$$

and so on.

$X:$	2	3	4	5	5	7	8	9	10	11	12
$P(X):$	$\frac{1}{36}$	$\frac{2}{36}$	$\frac{3}{36}$	$\frac{4}{36}$	$\frac{5}{36}$	$\frac{6}{36}$	$\frac{5}{36}$	$\frac{4}{36}$	$\frac{3}{36}$	$\frac{2}{36}$	$\frac{1}{36}$

is the required pmf of  $X$ .

3. Rewriting  $3x + 12y = 19$ ,  $y = -\frac{3}{12}x + \frac{19}{12}$  is the regression line of  $Y$  on  $X$ .

Rewriting  $3y + 9x = 46$ ,  $x = -\frac{3}{9}y + \frac{46}{9}$  is the regression line of  $X$  on  $Y$ .

$$\therefore b_{YX} = -\frac{3}{12} \text{ and } b_{XY} = -\frac{3}{9}; r_{XY}^2 = b_{YX} \cdot b_{XY} = \frac{9}{108} = \frac{1}{12}$$

$$\therefore r_{XY} = -\frac{1}{2\sqrt{3}} \text{ (-sign is taken, since } b_{YX} \text{ and } b_{XY} \text{ are both -ve)}$$

4. Statement of CLT is available in the book.  
 5. Definitions is available in the book.  
 6. The question is incomplete, as the state transition diagram is not given.

7.  $\frac{1}{\lambda} = 9 \text{ min } \therefore \lambda = \frac{1}{9} / \text{min}; \frac{1}{\mu} = 3 \text{ min } \therefore \mu = \frac{1}{3} / \text{min}$

$$P(W > 0) = 1 - P(W = 0) = 1 - P_0 = 1 - \left(1 - \frac{\lambda}{\mu}\right) = \frac{1}{\frac{1}{\lambda}} = \frac{1}{\frac{1}{3}} = 3$$

8.  $\lambda = \frac{1}{15} / \text{min}; \mu = \frac{1}{33} / \text{min}; k = 4$

$$P_0 = \frac{\frac{1-\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} = \frac{1 - \frac{11}{5}}{1 - \left(\frac{11}{5}\right)^5} = 0.0237$$

9. P.H. formula is  $L_S = \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}}$

$T$  follows  $Er(m, \mu)$   $\therefore E(T) = \frac{m}{\mu}$  and  $V(T) = \frac{m}{\mu^2}$

$$\therefore L_S = \lambda \cdot \frac{m}{\mu} + \left[ \lambda^2 \left\{ \frac{m}{\mu^2} + \frac{m^2}{\mu^2} \right\} \middle/ 2 \left\{ 1 - \lambda \cdot \frac{m}{\mu} \right\} \right]$$

$$= m\rho + \frac{m(1+m)\rho^2}{2(1-m\rho)}, \text{ where } \frac{\lambda}{\mu} = \rho .$$

10. (i) A ration shop in which one queue is to be formed in the billing section and the other in the ration issuing section.
- (ii) A medical shop in which one queue is formed to receiving drugs and the other in the billing and payment counter.

**Part-B**

11. (a) (i) Standard Book-work (available in the book).  
(ii) Definition, mean and variance of Weibull distribution (available in the book)
11. (b) (i) Standard book-work and property (available in the book).  
(ii)  $\lambda = 5/\text{min}$ ;

$$P\{X(t) = k\} = \frac{e^{-\lambda t} (\lambda t)^k}{k!} \quad \therefore P\{X(1) = 2\} = \frac{e^{-5} \cdot 5^2}{2!} = \frac{25}{2} e^{-5}$$

$$12. (a) g(x) = \int_0^1 f(x, y) dy = \int_0^1 \frac{x}{4} (1 + 3y^2) dy$$

$$= \frac{x}{4} (y + y^3) \Big|_0^1 = \frac{x}{2}; 0 < x < 2$$

$$h(y) = \int_0^2 f(x, y) dx = \int_0^2 \left( \frac{1 + 3y^2}{4} \right) x dx$$

$$= \left( \frac{1 + 3y^2}{4} \right) \left( \frac{x^2}{2} \right) \Big|_0^2 = \frac{1 + 3y^2}{2}; 0 < y < 1$$

$$f(x/y) = \frac{f(x, y)}{h(y)} = \frac{x(1 + 3y^2)}{4} \div \frac{(1 + 3y^2)}{2} = \frac{x}{2}; 0 < x < 2.$$

$$P\left[\left(\frac{1}{4} < x < \frac{1}{2}\right) \middle| y = \frac{1}{3}\right] = \frac{P\left[\left(\frac{1}{4} < x < \frac{1}{2}\right) \cap y = \frac{1}{3}\right]}{P\left(y = \frac{1}{3}\right)}$$

$$\begin{aligned} & \int_{\frac{1}{4}}^{\frac{1}{2}} x \left( \frac{1 + 3y^2}{4} \right) dy \\ &= \frac{\int_{\frac{1}{4}}^{\frac{1}{2}} x \left( \frac{1 + 3y^2}{4} \right) dy}{\int_0^{\frac{1}{2}} x \left( \frac{1 + 3y^2}{4} \right) dy} \end{aligned}$$

$$= \left( \frac{x^2}{2} \right) \Big|_{\frac{1}{4}}^{\frac{1}{2}} \div \left( \frac{x^2}{2} \right) \Big|_0^{\frac{1}{2}}$$

$$= \frac{1}{4} - \frac{1}{16} = \frac{3}{64}$$

[Note:  $y = \frac{1}{3}$  is meant as  $\frac{1}{3} - \frac{dy}{2} \leq y \leq \frac{1}{3} + \frac{dy}{2}$ ]

$$12. \text{ (b) (i)} \quad g(x) = \int_0^2 f(x, y) dy = \int_0^2 \left( x^2 + \frac{xy}{3} \right) dy = \left( x^2 y + \frac{xy^2}{6} \right)_{y=0}^{y=2} = 2x^2 + \frac{2}{3}x; \\ 0 \leq x \leq 1$$

$$h(y) = \int_0^1 f(x, y) dx = \int_0^1 \left( x^2 + \frac{xy}{3} \right) dx = \left( \frac{x^3}{3} + \frac{yx^2}{6} \right)_{x=0}^{x=1} = \frac{1}{3} + \frac{y}{6}; \\ 0 \leq y \leq 2$$

$f(x, y) \neq g(x) \cdot h(y)$   $\therefore X$  and  $Y$  are not independent.

$$\text{(ii)} \quad f_{XY}(x, y) = 6e^{-(2x+3y)}; x, y \geq 0 \quad (\because X \text{ and } Y \text{ are independent})$$

Let  $V = X - Y$  and  $W = Y$   $\therefore x = u + v$  and  $y = v$

$$\therefore J = \begin{vmatrix} x_u & x_v \\ y_u & y_v \end{vmatrix} = 1$$

$$f_{UV}(u, v) = |J| \cdot f_{XY}(x, y) = 6e^{-(2u+5v)} = 6e^{-(2u+5v)}$$

$$f_U(u) = \int_{-u}^{\infty} 6e^{-(2u+5v)} dv = \frac{6}{5} e^{3u}, \text{ when } u < 0$$

$$f_U(u) = \int_0^{\infty} 6e^{-(2u+5v)} dv = \frac{6}{5} e^{-2u}, \text{ when } u > 0$$

13. (a)  $E\{X(b)\} = \cos t \cdot E(A) + \sin t \cdot E(B)$ , where the probability distribution of  $A$  or

$$\text{(i)} \quad B \text{ is } \begin{cases} A: & -2 & 1 \\ p(A): & \frac{1}{3} & \frac{2}{3} \end{cases}$$

$$\therefore E(A) = E(B) = -\frac{2}{3} + \frac{2}{3} = 0 \text{ and } E(A^2) = E(B^2) = \frac{4}{3} + \frac{2}{3} = 2$$

$$E(AB) = E(A) \cdot E(B) = 0; E\{X(t)\} = 0$$

$$E\{X^2(t)\} = \cos^2 t \cdot E(A^2) + \sin^2 t \cdot E(B^2) + 2 \sin t \cos t E(A) \cdot E(B) \\ (\because A \text{ and } B \text{ are independent}).$$

$$= 2(\cos^2 t + \sin^2 t) = 2$$

$$\therefore \text{Var}\{X(t)\} = 2$$

Since  $E\{X(t)\}$  and  $\text{Var}\{X(t)\}$  are constants,  $\{X(t)\}$  is a WSS process.

(ii) Standard Book-work (available in the book)

13. (b) (i) Let  $\pi = (\pi_1, \pi_2, \pi_3)$  be the limiting state probability distribution of the Markov chain.

$$\text{Then } \pi P = \pi; \text{i.e., } (\pi_1, \pi_2, \pi_3) \begin{bmatrix} 0.4 & 0.5 & 0.1 \\ 0.3 & 0.3 & 0.4 \\ 0.3 & 0.2 & 0.5 \end{bmatrix} = (\pi_1, \pi_2, \pi_3)$$

$$\text{viz., } -6\pi_1 + 3\pi_2 + 3\pi_3 = 0, 5\pi_1 - 7\pi_2 + 2\pi_3 = 0, \pi_1 + 4\pi_2 - 5\pi_3 = 0$$

$$\text{Solving, } \pi_1 = \pi_2 = \pi_3 = k. \text{ But } \pi_1 + \pi_2 + \pi_3 = 1 \quad \therefore k = \frac{1}{3}$$

- (ii) Standard property (proof available in the book)

14. (a) (i)  $\lambda = 10 \text{ hour}; \frac{1}{\mu} = 5 \text{ min.} \quad \therefore \mu = \frac{1}{5} / \text{min or } 12/\text{hour}; k = 3$

$$(1) P(\text{an arriving car gets service immediately on arrival}) =$$

$$P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} = \frac{1 - \frac{5}{6}}{1 - \left(\frac{5}{6}\right)^4} = 0.3219$$

$$(2) P(\text{an arriving car has to wait outside the indicated space})$$

$$P_3 = \left(\frac{\lambda}{\mu}\right)^3 \left\{ \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^4} \right\} = 0.5787 \times 0.3219 = 0.1863$$

$$(3) E(W_Q) = \frac{1}{\lambda'} E(N_Q), \text{ where } \lambda' = \mu(1 - P_0) = 12 \times (1 - 0.1863)$$

$$= 9.7644$$

$$\begin{aligned} E(N_S) &= \frac{\lambda}{\mu - \lambda} - \frac{(k+1)\left(\frac{\lambda}{\mu}\right)^{k+1}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} \\ &= \frac{10}{12 - 10} - \frac{4 \times \left(\frac{5}{6}\right)^4}{1 - \left(\frac{5}{6}\right)^4} = 5 - \frac{4 \times 0.48225}{0.51775} \\ &= 1.2743 \end{aligned}$$

$$E(N_Q) = E(N_S) - \frac{\lambda}{\mu} = 1.2743 - 0.8333 = 0.4110$$

$$\therefore E(W_Q) = \frac{0.4110}{9.7644} \text{ hour or } \frac{0.4110}{9.7644} \times 60 \text{ min} = 2.5255 \text{ min}$$

- (ii) A standard characteristic of (M/M/1): ( $\infty$ /FCFS) model; derivation available in the book.
14. (b) Steady-state solution for the (M/M/C): ( $\infty$ /FIFO) model and the derivations of the formulas for  $L_Q$  [ $E(N_Q)$ ],  $W_Q$ ,  $W_S$  and  $L_S$  [ $E(N_S)$ ] are standard results (available in the book).
15. (a) What is required is the formula for  $E(N_S)$  in a (M/G/1) model. This is the same as the derivation of Pollaczek–Khinchine formula (available in the book).
15. (b) (i) and (ii) Answers are to be lifted from the relevant pages in the book.



**BE/BTech Degree Examination, May/June 2012**

**Fourth Semester**

**Computer Science and Engineering**

**MA2262—Probability and Queueing Theory**

**(Common to Information Technology)**

**(Regulation 2008)**

**Time: Three hours**

**Maximum: 100 marks**

**Statistical Tables may be permitted.**

**Answer ALL Questions**

**Part-A (10 × 2 = 20 marks)**

1. Check whether the following is a probability density function or not:  
$$\begin{cases} \lambda e^{-\lambda x}, & x \leq 0, \lambda > 0 \\ 0, & \text{elsewhere} \end{cases}$$
2. If a random variable has the moment-generating function given by  $M_X(t) = \frac{2}{2-t}$ , determine the variance of  $X$ .
3. The regression equations of  $X$  on  $Y$  and  $Y$  on  $X$  are respectively  $5x - y = 22$  and  $64x - 45y = 24$ . Find the means of  $X$  and  $Y$ .
4. State the Central Limit theorem.
5. Define wide sense stationary process.
6. If the initial state probability distribution of a Markov chain is  $p^{(0)} = \begin{pmatrix} 5 & 1 \\ 6 & 6 \end{pmatrix}$  and the transition probability matrix of the chain is  $\begin{pmatrix} 0 & 1 \\ 1 & 1 \\ 2 & 2 \end{pmatrix}$ , find the probability distribution of the chain after 2 steps.
7. State Little's formula for a (M/M/1): (GD/N/∞) queueing model.
8. Define steady-state and transient state in Queueing theory.
9. When will an M/G/1 queueing model become a classic M/M/1 queueing model?
10. State Pollaczek–Khinchine formula for the average number of customers in an M/G/1 queueing model.

**Part-B (5 × 16 = 80 Marks)**

11. (a) (i) A random variable  $X$  has the following probability function:

$X$	0	1	2	3	4	5	6	7
$P(x)$	0	$k$	$2k$	$2k$	$3k$	$k^2$	$2k^2$	$7k^2 + k$

- (1) Find the value of  $k$ .  
 (2) Evaluate  $p(X < 6)$ ,  $P(X \geq 6)$   
 (3) If  $p(X \leq c) > \frac{1}{2}$ , find the minimum value of  $c$ .

- (ii) Find the moment-generating function of an exponential random variable and, hence, find its mean and variance.

Or

- (b) (i)  $X$  is a Poisson variate such that

$$p(X = 2) = 9p(X = 4) + 90 p(X = 6). \text{ Find}$$

- (1) Mean and  $E(X^2)$   
 (2)  $p(X \geq 2)$ .

- (ii) In a certain city, the daily consumption of electric power in millions of kilowatt-hours can be treated as a random variable having Gamma distribution with parameters  $\lambda = \frac{1}{2}$  and  $v = 3$ . If the power plant of this city has a daily capacity of 12 million kilowatt-hours, what is the probability that this power supply will be inadequate on any given day?

12. (a) (i) Let  $X$  and  $Y$  be two random variables having the joint probability function  $f(x, y) = k(x + 2y)$  where  $x$  and  $y$  can assume only the integer values 0, 1 and 2. Find the marginal and conditional distributions.  
 (ii) Two random variables  $X$  and  $Y$  have the joint probability density function

$$f(x, y) = \begin{cases} c(4 - x - y), & 0 \leq x \leq 2, 0 \leq y \leq 2 \\ 0, & \text{elsewhere} \end{cases}. \text{ Find } \text{cov}(X, Y)$$

Or

- (b) (i) Two dimensional random variables  $(X, Y)$  have the joint probability density function  $f(x, y) = 8xy$ ,  $0 < x < y < 1$   
 $= 0$ , elsewhere

$$(1) \text{ Find } P\left(X < \frac{1}{2} \cap Y < \frac{1}{4}\right)$$

- (2) Find the marginal and conditional distributions.  
 (3) Are  $X$  and  $Y$  independent?

- (ii) Suppose that in a certain circuit, 20 resistors are connected in series. The mean and variance of each resistor are 5 and 0.20 respectively. Using the central limit theorem, find the probability

that the total resistance of the circuit will exceed 98 ohms assuming independence.

13. (a) (i) The process  $\{X(t)\}$  whose probability distribution under certain condition is given by

$$P[X(t) = n] = \frac{(at)^{n-1}}{(1+at)^{n+1}}, \quad n = 1, 2, 3, \dots$$

$$= \frac{at}{1+at}, \quad n = 0$$

stationary.

- (ii) A salesman territory consists of three cities,  $A$ ,  $B$  and  $C$ . He never sells in the same city on successive days. If he sells in city- $A$ , then the next day he sells in city- $B$ . However, if he sells in either city- $B$  or city- $C$ , the next day he is twice as likely to sell in city- $A$  as in the other city. In the long run, how often does he sell in each of the cities?

Or

- (b) (i) The transition probability matrix of a Markov chain  $\{X(t)\}$ ,  $n = 1, 2, 3, \dots$ , having three states 1, 2 and 3 is  $P = \begin{bmatrix} 0.1 & 0.5 & 0.4 \\ 0.6 & 0.2 & 0.2 \\ 0.3 & 0.4 & 0.3 \end{bmatrix}$  and

the initial distribution is  $p^{(0)} = (0.7 \ 0.2 \ 0.1)$ .

Find (1)  $p[X_2 = 3]$

(2)  $p[X_3 = 2, X_2 = 3, X_1 = 3, X_0 = 2]$ .

- (ii) Suppose that customers arrive at a bank according to a Poisson process with mean rate of 3 per minute. Find the probability that during a time interval of two minutes,
- (1) exactly 4 customers arrive
  - (2) greater than 4 customers arrive
  - (3) fewer than 4 customers arrive.

14. (a) (i) A T.V. repairman finds that the time spent on his job has an exponential distribution with a mean of 30 minutes. The repair sets in the order in which they came in and if the arrival of sets is approximately Poisson with an average rate of 10 per 8 hour day.

(1) What is the repairman's expected idle time each day?

(2) How many jobs are ahead of average set just brought?

- (ii) A supermarket has 2 girls running up sales at the counters. If the service time for each customer is exponential with mean 4 minutes and if people arrive in Poisson fashion at the rate of 10 per hour, find the following:

(1) What is the probability of having to wait for service?

(2) What is the expected percentage of idle time for each girl?

(3) What is the expected length of a customer's waiting time?

Or

- (b) (i) Trains arrive at the yard every 15 minutes and the service time is 33 minutes. If the line capacity of the yard is limited to 5 trains, find the probability that the yard is empty and the average number of trains in the system, given that the inter-arrival time and service time are following exponential distribution.
- (ii) There are three typists in an office. Each typist can type an average of 6 letters per hour. If letters arrive for being typed at the rate of 15 letters per hour, what fraction of times will all the typists be busy? What is the average number of letters waiting to be typed?
15. (a) (i) An automatic car-wash facility operates with only one bay. Car arrive according to a Poisson distribution with a mean of 4 cars per hour and may wait in the facility's parking lot if the bay is busy. The parking lot is large enough to accommodate any number of cars. If the service time for all cars is constant and equal to 10 minutes, determine
- (1) mean number of customers in the system,  $L_s$
  - (2) mean number of customers in the queue,  $L_q$
  - (3) mean waiting time of a customer in the system,  $W_s$
  - (4) mean waiting time of a customer in the queue,  $W_q$
- (ii) An average of 120 students arrive each hour (inter-arrival times are exponential) at the controller office to get their hall tickets. To complete the process, a candidate must pass through three counters. Each counter consists of a single server, service times at each counter are exponential with the following mean times: counter 1, 20 seconds; counter 2, 15 seconds and counter 3, 12 seconds. On the average, how many students will be present in the controller's office?

Or

- (b) (i) Derive the  $P-K$  formula for  $(M/G/1):(GD/\infty/\infty)$  queueing model and hence deduce that with the constant service time the  $P-K$  reduces to  $L_s = \rho + \frac{\rho^2}{2(1-\rho)}$  where  $\mu = \frac{1}{E(T)}$  and  $\rho = \frac{\lambda}{\mu}$ .
- (ii) For an open queueing network with three nodes 1, 2 and 3, let customers arrive from outside the system to node  $j$  according to a Poisson input process with parameters  $r_j$  and let  $P_{ij}$  denote the proportion of customers departing from facility  $i$  to facility  $j$ . Given  $(r_1, r_2, r_3) = (1, 4, 3)$  and  $P_{ij} = \begin{bmatrix} 0 & 0.6 & 0.3 \\ 0.1 & 0 & 0.3 \\ 0.4 & 0.4 & 0 \end{bmatrix}$ , determine the average rate  $\lambda_j$  to the node  $j$  for  $j = 1, 2, 3$ .

## Solutions

### Part-A

1.  $f(x) = \lambda e^{-\lambda x} \geq 0$ , for all  $x \geq 0$ , since  $\lambda > 0$

$$\text{and } \int_0^{\infty} f(x) dx = \left( \lambda \frac{e^{-\lambda x}}{-\lambda} \right)_0^{\infty} = 1 \quad \therefore f(x) \text{ is a genuine pdf.}$$

$$2. M_X(t) = \frac{2}{2\left(1 - \frac{t}{2}\right)} = 1 + \frac{t}{2} + \frac{t^2}{4} + \dots \infty$$

$$E(X) = \text{coefficient of } \frac{t}{1} = \frac{1}{2}; E(X^2) = \text{coefficient of } \frac{t^2}{2} = \frac{1}{2}$$

$$\therefore V(X) = E(X^2) - E^2(X) = \frac{1}{2} - \frac{1}{4} = \frac{1}{4}$$

3.  $(\bar{X}, \bar{Y})$  is the point of intersection of the two regression lines.

$$5\bar{X} - \bar{Y} = 22 \text{ or } 225\bar{X} - 45\bar{Y} = 990 \dots (1); 64\bar{X} - 45\bar{Y} = 24 \dots (2)$$

Solving of (1) and (2);  $\bar{X} = 6$  and  $\bar{Y} = 8$ .

4. Statement of CLT available in the book.

5. Definition of WSS available in the book.

$$6. p^{(1)} = p^{(0)}p = \left( \frac{5}{6}, \frac{1}{6} \right) \begin{pmatrix} 0 & 1 \\ \frac{1}{2} & \frac{1}{2} \end{pmatrix} = \left( \frac{1}{12}, \frac{11}{12} \right)$$

$$p^{(2)} = p^{(1)}p = \left( \frac{1}{12}, \frac{11}{12} \right) \begin{pmatrix} 0 & 1 \\ \frac{1}{2} & \frac{1}{2} \end{pmatrix} = \left( \frac{11}{24}, \frac{13}{24} \right)$$

7. Statements of Little's formulas are available in the book.

8. Definitions are available in the book.

9. In the (M/M/1) model, the number of customers serviced follows a Poisson process or the service time follows an exponential distribution, whereas in the (M/G/1) model, the service time follows an arbitrary (general) distribution.

10. Pollaczek–Khinchine formula for  $L_S$  or  $E[N_S]$  is available in the book.

### Part-B

11. (a) (i) This is a worked example in the book.

Part (2) of the problem is not given in the book, which is given below:

$$(2) P(X < 6) = P(X = 0 \text{ or } 1 \text{ or } \dots \text{ or } 5)$$

$$= 8k + k^2 = \frac{8}{10} + \frac{1}{100} = \frac{81}{100} \text{ and } P(X \geq 6) = 1 - \frac{81}{100} = \frac{19}{100}$$

(ii) Pdf of the exponential RV  $X$  is  $f(x) = \lambda e^{-\lambda x}; x \geq 0; \lambda > 0$

$$M_X(t) = \int_0^{\infty} e^{tx} \lambda e^{-\lambda x} dx = \lambda \left\{ \frac{e^{-(\lambda-t)x}}{-(\lambda-t)} \right\}_0^{\infty} = \frac{\lambda}{\lambda-t}$$

$$M_X(t) = \left( 1 - \frac{t}{\lambda} \right)^{-1} = 1 + \frac{t}{\lambda} + \frac{t^2}{\lambda^2} + \dots \infty \dots (1)$$

$$E(X) = \text{coefficient of } \frac{t}{[1]} \text{ in (1)} = \frac{1}{\lambda}$$

$$E(X^2) = \text{coefficient of } \frac{t^2}{[2]} \text{ in (1)} = \frac{2}{\lambda^2}$$

$$V(X) = E(X^2) - E^2(X) = \frac{2}{\lambda^2} - \frac{1}{\lambda^2} = \frac{1}{\lambda^2}$$

11. (b) (i) Let  $X$  follow a  $P(\lambda)$   $\therefore P(X=r) = \frac{e^{-\lambda} \cdot \lambda^r}{[r]}$ ;  $\lambda > 0$

$$\text{Given: } P(X=2) = 9P(X=4) + 90 P(X=6)$$

$$\therefore \frac{e^{-\lambda} \lambda^2}{[2]} = \frac{9e^{-\lambda} \lambda^4}{[4]} + \frac{90e^{-\lambda} \lambda^6}{[6]}; \text{ i.e., } 360 \lambda^2 = 270 \lambda^4 + 90 \lambda^6$$

$$\text{Since } \lambda^2 \neq 0, \lambda^4 + 3\lambda^2 - 4 = 0; \text{ i.e., } (\lambda^2 - 1)(\lambda^2 + 4) = 0$$

$$\text{Since } \lambda^2 \neq -4, \text{ we get } \lambda^2 = 1; \text{ Since } \lambda > 0, \lambda = 1$$

$$(1) E(X) = \lambda = 1; V(X) = \lambda = 1; \text{ i.e., } E(X^2) - E^2(X) = 1 \quad \therefore E(X^2) = 2$$

$$(2) P(X \geq 2) = 1 - P(X < 2) = 1 - \{P(X=0) + P(X=1)\} \\ = 1 - \{e^{-\lambda} + \lambda e^{-\lambda}\} = 1 - 2e^{-\lambda} \\ = 1 - \frac{2}{e}$$

(ii) This is a worked example in the book.

12. (a) (i)

$X \backslash Y$	0	1	2
0	0	$2k$	$4k$
1	$k$	$3k$	$5k$
2	$2k$	$4k$	$6k$

$$\sum_{y_j=0}^2 \sum_{x_j=0}^2 p_{ij} = 1; \text{ i.e., } 27 k = 1$$

$$\therefore k = \frac{1}{27}$$

Marginal distribution of  $X$ 

$X$	$P$
0	6/27
1	9/27
2	12/27

 Marginal distribution of  $Y$ 

$Y$	$P$
0	3/27
1	9/27
2	15/27

 CPD of  $X/Y = 0$ 

$X$	$P(X/Y) = 0$
0	0
1	$k / 3k = \frac{1}{3}$
2	$2k/3k = \frac{2}{3}$

 CPD of  $X/Y = 1$ 

$X$	$P(X/Y) = 1$
0	$\frac{2}{9}$
1	$\frac{3}{9}$
2	$\frac{4}{9}$

 CPD of  $X/Y = 2$ 

$X$	$P(X/Y) = 2$
0	$\frac{4}{15}$
1	$\frac{5}{15}$
2	$\frac{6}{15}$

 CPD of  $Y/X = 0$ 

$Y$	$P(Y/X = 0)$
0	0
1	$\frac{2}{6}$
2	$\frac{4}{6}$

 CPD of  $Y/X = 1$ 

$Y$	$P(Y/X = 1)$
0	$\frac{1}{9}$
1	$\frac{3}{9}$
2	$\frac{5}{9}$

 CPD of  $Y/X = 2$ 

$Y$	$P\{Y(X = 2)\}$
0	$\frac{2}{12}$
1	$\frac{4}{12}$
2	$\frac{6}{12}$

$$(ii) \quad \iint_{R_{XY}} f(x, y) dx dy = 1; \text{ i.e., } \int_0^2 \int_0^2 (4 - x - y) dx dy = 1$$

$$\text{i.e., } c \cdot \int_0^2 \left( 4x - \frac{x^2}{2} - yx \right)_0^2 dy = 1; \text{ i.e., } c \int_0^2 (6 - 2y) dy = 1$$

$$\text{i.e., } c(6y - y^2)_0^2 = 1; \text{i.e., } 8C = 1 \quad \therefore C = \frac{1}{8}$$

$$\begin{aligned} E(X) &= \int_0^2 \int_0^2 x f(x,y) dx dy = \frac{1}{8} \int_0^2 \int_0^2 (4x - x^2 - yx) dx dy \\ &= \frac{1}{8} \int_0^2 \left( 2x^2 - \frac{x^3}{3} - y \frac{x^2}{2} \right)_0^2 dy = \frac{1}{8} \int_0^2 \left( \frac{16}{3} - 2y \right) dy = \frac{5}{6} \\ &= E(Y), \text{ by symmetry} \end{aligned}$$

$$\begin{aligned} E(XY) &= \int_0^2 \int_0^2 x y f(x,y) dx dy = \frac{1}{8} \int_0^2 \int_0^2 (4xy - x^2 y - xy^2) dx dy \\ &= \frac{1}{8} \int_0^2 \left( 2yx^2 - y \frac{x^3}{3} - y^2 \frac{x^2}{2} \right)_0^2 dy = \frac{1}{8} \int_0^2 \left( 8y - \frac{8}{3} - 2y^2 \right) dy = \frac{2}{3} \\ \therefore \text{Cov}(X, Y) &= E(XY) - E(X) \cdot E(Y) = \frac{2}{3} - \frac{25}{36} = -\frac{1}{36} \end{aligned}$$

12. (b) (i)

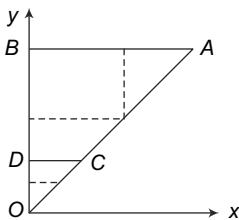


Fig. 1

$$(1) P\left(X < \frac{1}{2} \cap Y < \frac{1}{4}\right) = \int_0^{\frac{1}{4}} \int_0^y 8xy dx dy = \int_0^{\frac{1}{4}} 4y^3 dy = \frac{1}{256}$$

$$(2) f_X(x) = \int_x^1 8xy dy = 4x(y^2)_x^1 = 4x - 4x^3; 0 < x < 1$$

$$f_Y(y) = \int_0^y 8xy dx = 4y(x^2)_0^y = 4y^3; 0 < y < 1$$

$$(3) f_X(x) \cdot f_Y(y) \neq f(x, y) \quad \therefore X \text{ and } Y \text{ are not independent.}$$

(ii)  $E(X_i) = 5$  and  $V(X_i) = 0.20$

By the central limit theorem,  $S_{20} = X_1 + X_2 + \dots + X_{20}$  follows  $N(n\mu, \sigma\sqrt{n})$

i.e.,  $S_{20}$  follows  $N(100, 2)$

$$\begin{aligned}\therefore P(S_{20} > 98) &= P\left\{\frac{S_{20} - 100}{2} > \frac{98 - 100}{2}\right\} \\ &= P(Z > -1) = P(0 < Z < 1) + 0.5 = 0.8413\end{aligned}$$

13. (a) (i) This is a worked example in the book.  
(ii) The one-step tpm of the corresponding Markov chain is

$$\begin{array}{c} A \quad B \quad C \\ \begin{matrix} A \\ B \\ C \end{matrix} \end{array} \begin{pmatrix} 0 & 1 & 0 \\ \frac{2}{3} & 0 & \frac{1}{3} \\ \frac{2}{3} & \frac{1}{3} & 0 \end{pmatrix}. \text{ If } (\pi_1, \pi_2, \pi_3) \text{ is the stationary state probability distribution, then } (\pi_1, \pi_2, \pi_3) \begin{pmatrix} 0 & 1 & 0 \\ \frac{2}{3} & 0 & \frac{1}{3} \\ \frac{2}{3} & \frac{1}{3} & 0 \end{pmatrix} = (\pi_1, \pi_2, \pi_3)$$

$$\text{i.e., } 3\pi_1 - 2\pi_2 - 2\pi_3 = 0$$

$$3\pi_1 - 3\pi_2 + \pi_3 = 0$$

$$\pi_2 = 3\pi_3$$

Solving these equations, we get

$$\frac{\pi_1}{-8} = \frac{\pi_2}{-9} = \frac{\pi_3}{-3} = k \text{ and } \pi_1 + \pi_2 + \pi_3 = 1$$

$$\text{i.e., } \pi_1 = \frac{2}{5}, \pi_2 = \frac{9}{20} \text{ and } \pi_3 = \frac{3}{20}$$

i.e., the salesman sells in the cities  $A, B, C$ , 40%, 45% and 15% of the days respectively.

13. (b) (i) This is a worked example in the book.  
(ii) This is a worked example in the book. Solution for part (3) is not available in the book, which is given below:  
(3)  $P\{X(2) < 4\} = 1 - P\{X(2) \geq 4\} = 1 - (0.133 + 0.715) = 0.152$

14. (a) (i)  $\frac{1}{\mu} = 30 \text{ min}; \therefore \mu = \frac{1}{30} / \text{min or } 2/\text{hour}; \lambda = \frac{10}{8} \text{ or } \frac{5}{4} / \text{hour}$

(1)  $P(\text{the repairman is idle}) = P(\text{no customer in the system})$

$$= P_0 = 1 - \frac{\lambda}{\mu} = 1 - \frac{5}{8} = \frac{3}{8}$$

∴ Expected idle time each 8 hour day =  $\frac{3}{8} \times 8 = 3$  hours

$$(2) E(\text{No. of customers in the system}) = E(N_S) = \frac{\lambda}{\mu - \lambda} = \frac{5}{3}$$

(ii) This is a worked example in the book.

14. (b) (i)  $\frac{1}{\lambda} = 15$  min or  $\lambda = \frac{1}{15}$  /min or 4/hour;  $\frac{1}{\mu} = 33$  min or  $\mu = \frac{1}{33}$  /min

or  $\frac{20}{11}$  /hour

$$k = 5, \quad P_0 = \frac{1 - \frac{\lambda}{\mu}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} = \frac{1 - (2.2)}{1 - (2.2)^6} = 0.0187$$

$$E(N_S) = \frac{\lambda}{\mu - \lambda} - \frac{(k+1)\left(\frac{\lambda}{\mu}\right)^{k+1}}{1 - \left(\frac{\lambda}{\mu}\right)^{k+1}} = \frac{132}{131} + \frac{6 \times 113 \cdot 380}{112 \cdot 380} \approx 5$$

(ii) This is a worked example in the book.

15. (a) (i)  $\lambda = 4$ /hour or  $\frac{1}{15}$  /min  $E(T) = 10$  and  $V(T) = 0$

$$(1) \text{ By P.K. formula, } E(N_S) = L_S = \lambda \cdot E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}}$$

$$\text{i.e., } L_S = \frac{1}{15} \times 10 + \frac{\frac{1}{225}(0+100)}{2\left\{1 - \frac{1}{15} \times 10\right\}} = \frac{4}{3}$$

$$(2) L_Q = E(N_Q) = E(N_S) - \frac{\lambda}{\mu} = \frac{4}{3} - \frac{\frac{1}{15}}{\frac{1}{10}} = \frac{2}{3}$$

$$(3) E(W_S) = \frac{1}{\lambda} E(N_S) = 15 \times \frac{4}{3} = 20 \text{ min}$$

$$(4) E(W_Q) = E(W_S) - \frac{1}{\mu} = 20 - 10 = 10 \text{ min}$$

(ii)  $\lambda = 120$ /hour or  $\frac{1}{30}$  /second;  $\mu_1 = \frac{1}{20}$  /second;  $\mu_2 = \frac{1}{15}$  /second;

$$\mu_3 = \frac{1}{12} \text{ /second}$$

$$\begin{aligned}
 E(N_S) &= \frac{\lambda}{\mu_1 - \lambda} + \frac{\lambda}{\mu_2 - \lambda} + \frac{\lambda}{\mu_3 - \lambda} \\
 &= \frac{1}{30} \left\{ \frac{1}{\frac{1}{20} - \frac{1}{30}} + \frac{1}{\frac{1}{15} - \frac{1}{30}} + \frac{1}{\frac{1}{12} - \frac{1}{30}} \right\} \\
 &= \frac{1}{30} \left\{ \frac{600}{10} + \frac{450}{15} + \frac{360}{18} \right\} = \frac{60 + 30 + 20}{30} = \frac{11}{3}
 \end{aligned}$$

15. (b) (i) Derivation of P.K. formula is available in the book.

$$\text{P.K. formula is a } L_S = \lambda E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}}, \quad E(T) = \frac{1}{\mu} \text{ and} \\
 V(T) = 0$$

$$= \frac{\lambda^2 \cdot \frac{1}{\mu^2}}{2\left(1 - \frac{\lambda}{\mu}\right)} = \rho + \frac{\rho^2}{2(1-\rho)}$$

(ii) Jackson's flow balance equations are  $\lambda_j = r_j + \sum_{i=1}^3 \lambda_i \cdot P_{ij}$ ; ( $j = 1, 2, 3$ )

$$\begin{aligned}
 \text{i.e., } \lambda_1 &= 1 + \lambda_1 P_{11} + \lambda_2 P_{21} + \lambda_3 P_{31} \text{ (when } j = 1) \\
 \text{i.e., } \lambda_1 &= 1 + 0.1\lambda_2 + 0.4\lambda_3
 \end{aligned} \tag{1}$$

When  $j = 2$ ,

$$\begin{aligned}
 \lambda_2 &= 4 + \lambda_1 P_{12} + \lambda_2 P_{22} + \lambda_3 P_{32} \\
 \text{i.e., } \lambda_2 &= 4 + 0.6\lambda_1 + 0.4\lambda_3
 \end{aligned} \tag{2}$$

When  $j = 3$ ,

$$\begin{aligned}
 \lambda_3 &= 3 + \lambda_1 P_{13} + \lambda_2 P_{23} + \lambda_3 P_{33} \\
 \text{i.e., } \lambda_3 &= 3 + 0.3\lambda_1 + 0.3\lambda_2
 \end{aligned} \tag{3}$$

Solving (1), (2) and (3);  $\lambda_1 = 5$ ,  $\lambda_2 = 10$ ,  $\lambda_3 = \frac{15}{2}$



**BE/BTech Degree Examination, May/June 2013**

**Fourth Semester**

**Electronics and Communication Engineering**

**MA2261/MA45/MA1253/10177 PR 401/080380009**

**Probability and Random Processes**

**(Common to Biomedical Engineering)**

**(Regulation 2008/2010)**

**Time: Three hours**

**Maximum: 100 marks**

**Use of Statistical Tables is permitted.**

**Answer ALL Questions.**

**Part-A (10 × 2 = 20 marks)**

1. A random variable  $X$  has cdf

$$F_X(x) = \begin{cases} 0; & x < 1 \\ 1/2(x-1); & 1 \leq x < 3 \\ 1; & x \geq 3 \end{cases}$$

Find the pdf of  $X$  and the expected value of  $X$ .

2. Find the moment-generating function of a binomial distribution.
3. The joint pmf of two random variables  $X$  and  $Y$  is given by

$$P_{X, Y}(x, y) = \begin{cases} kxy, & x = 1, 2, 3; y = 1, 2, 3 \\ 0, & \text{otherwise} \end{cases}$$

Determine the value of the constant  $k$ .

4. The joint pdf of a random variable  $(X, Y)$  is  $f_{xy}(x, y) = xy^2 + \frac{x^2}{8}$ ,  $0 \leq x \leq 2$ ,  $0 \leq y \leq 1$ . Find  $P\{X < Y\}$ .
5. Define wide sense stationary process.
6. Show that a binomial process is Markov.
7. A random process  $X(t)$  is defined by  $X(t) = K \cos \omega t$ ,  $t \geq 0$ , where  $\omega$  is a constant and  $K$  is uniformly distributed over  $(0, 2)$ . Find the autocorrelation function of  $X(t)$ .
8. Define cross correlation function of  $X(t)$  and  $Y(t)$ . When do you say that they are independent?
9. Define a linear time invariant system.
10. State the convolution form of the output of a linear time invariant system.

**Part-B (5 × 16 = 80 marks)**

11. (a) (i) A random variable  $X$  has pdf

$$f_X(x) = \begin{cases} kx^2 e^{-x}; & x > 0 \\ 0; & \text{otherwise} \end{cases}$$

Find the  $r^{\text{th}}$  moment of  $X$  about the origin. Hence, find the mean and variance. (8)

- (ii) A random variable  $X$  is uniformly distributed over  $(0, 10)$ . Find

$$(1) P(X < 3), P(X > 7) \text{ and } P(2 < X < 5)$$

$$(2) P(X = 7) \quad (8)$$

Or

- (b) (i) An office has four phone lines. Each is busy about 10% of the time. Assume that the phone lines act independently.

$$(1) \text{ What is the probability that all four phones are busy?}$$

$$(2) \text{ What is the probability that at least two of them are busy?} \quad (6)$$

- (ii) Describe gamma distribution. Obtain its moment generating function. Hence, compute its mean and variance. (10)

12. (a) (i) Two independent random variables  $X$  and  $Y$  are defined by

$$f_X(x) = \begin{cases} 4ax; & 0 < x < 1 \\ 0; & \text{otherwise} \end{cases} \text{ and } f_Y(y) = \begin{cases} 4by; & 0 < y < 1 \\ 0; & \text{otherwise} \end{cases}$$

Show that  $U = X + Y$  and  $V = X - Y$  are uncorrelated. (8)

- (ii) State and prove the central limit theorem for the case of iid random variables. (8)

Or

- (b) (i) The equations of two regression lines are  $3x + 12y = 19$  and  $3y + 9x = 46$ . Find  $\bar{x}, \bar{y}$  and the correlation coefficient between  $X$  and  $Y$ . (8)

- (ii) Give the joint pdf of  $X$  and  $Y$

$$f_{X, Y}(x, y) = \begin{cases} CX(x - y); & 0 < x < 2, -x < y < x \\ 0; & \text{otherwise} \end{cases}$$

$$(1) \text{ Evaluate } C.$$

$$(2) \text{ Find the marginal pdf of } X.$$

$$(3) \text{ Find the conditional density of } Y|X. \quad (8)$$

13. (a) (i) Define a semirandom telegraph signal process and prove that it is evolutionary. (10)

- (ii) Mention any three properties each of autocorrelation and of cross correlation functions of a wide sense stationary process. (6)

Or

- (b) (i) A random process  $X(t)$  is defined by

$$X(t) = A \cos t + B \sin t; -\infty < t < \infty$$

where  $A$  and  $B$  are independent random variables, each of which has a value of  $-2$  with a probability of  $1/3$  and a value of  $1$  with

- a probability of 2/3. Show that  $X(t)$  is a wide sense stationary process. (8)
- (ii) Define a Poisson process. Show that the sum of two Poisson processes is a Poisson process. (8)
14. (a) (i) Define spectral density of a stationary random process  $X(t)$ . Prove that for a real random process  $X(t)$ , the power spectral density is an even function. (8)
- (ii) Two random processes  $X(t)$  and  $Y(t)$  are defined as follows:  
 $X(t) = A \cos(\omega t + \theta)$  and  $Y(t) = B \sin(\omega t + \theta)$  where  $A, B$  and  $\omega$  are constants;  $\theta$  is a uniform random variable over  $(0, 2\pi)$ . Find the cross correlation function of  $X(t)$  and  $Y(t)$ . (8)
- Or
- (b) (i) State and prove the Wiener–Khintchine theorem. (8)
- (ii) If the cross power spectral density of  $X(t)$  and  $Y(t)$  is
- $$S_{XY}(\omega) = \begin{cases} a + \frac{ib\omega}{\alpha}; & -\alpha < \omega < \alpha, \alpha > 0 \\ 0; & \text{otherwise} \end{cases}$$
- where  $a$  and  $b$  are constants. Find the cross correlation function. (8)
15. (a) (i) A random process  $X(t)$  is the input to a linear system whose impulse function is  $h(t) = 2e^{-t}$ ;  $t \geq 0$ . The autocorrelation function of the process is  $R_{XX}(\tau) = e^{-2|\tau|}$ . Find the power spectral density of the output process  $Y(t)$ . (8)
- (ii) A wide sense stationary noise process  $N(t)$  has an autocorrelation function  $R_{NN}(\tau) = Pe^{-3|\tau|}$  where  $P$  is a constant. Find its power spectrum.
- Or
- (b) (i) If the input to a time invariant stable, linear system is a wide sense stationary process, prove that the output will also be a wide sense stationary process. (8)
- (ii) Let  $X(t)$  be a wide sense stationary process which is the input to a linear time invariant system with unit impulse  $h(t)$  and output  $Y(t)$ . Then prove that  $S_{YY}(\omega) = |H(\omega)|^2 S_{XX}(\omega)$  where  $H(\omega)$  is the Fourier transform of  $h(t)$ . (8)

## Solutions

### Part-A

1.  $f(x) = \frac{d}{dx} F(x) = \begin{cases} 0, & \text{in } x < 1 \\ \frac{1}{2}, & \text{in } 1 \leq x < 3; E(X) = \int_1^3 \frac{x}{2} dx = \left[ \frac{x^2}{4} \right]_1^3 = 2. \\ 0, & \text{in } x \geq 3 \end{cases}$

2.  $P(X = r) = p_r = nC_r p^r q^{n-r}; r = 0, 1, \dots, n$

$$M_X(t) = \sum_{r=0}^n e^{tr} \cdot nC_r p^r q^{n-r} = \sum_{r=0}^n nC_r (pe^t)^r q^{n-r} = (q + pe^t)^n$$

3.

x/y	1	2	3
1	k	2k	3k
2	2k	4k	6k
3	3k	6k	9k

$$\sum_{y=1}^3 \sum_{x=1}^3 p_{xy} = 1; \text{i.e., } 36k = 1$$

$$\therefore k = \frac{1}{36}$$

4. This is a part of a worked example in the book.  
 5. Definition of WSS process—available in the book.  
 6. When  $P(X_n = 1) = p$  and  $P(X_n = 0) = q$ , where  $p + q = 1$ , then  $\{X_n\}$  is a Bernoulli's process.

If  $S_n = \sum_{n=1}^n X_n$ , then  $\{S_n; n \geq 1\}$  is a Binomial process.

Now,  $P(S_{n+1} = k + 1 | S_n = k) = p$  and  $P(S_{n+1} = k | S_n = k) = q$   
 $\therefore \{S_n\}$  is a Markov process.

7.  $K$  is  $U(0, 2)$        $\therefore \text{Var}(K) = \frac{(2-0)^2}{12} = \frac{1}{3}; E(K) = \frac{b+a}{2} = 1$

$$\text{Var}(K) = E(K)^2 - E^2(K); E(K^2) = \frac{4}{3}$$

$$R(t_1, t_2) = E\{X(t_1) \cdot X(t_2)\} = E\{K^2 \cos \omega t_1 \cdot \cos \omega t_2\} = \frac{4}{3} \cos \omega t_1 \cos \omega t_2$$

8.  $R_{XY}(\tau) = E\{X(t) \cdot Y(t - \tau)\}$   
 $\{X(t)\}$  and  $\{Y(t)\}$  are independent, if  $R_{XY}(\tau) = \mu_X \cdot \mu_Y$   
 9. Definition is available in the book.

10.  $Y(t) = \int_{-\infty}^{\infty} h(u) X(t-u) du$ , where  $X(t)$  and  $Y(t)$  are the input and output process  
 and  $h(t)$  is the system weighting function.

### Part-B

11. (a) (i)  $\int_0^{\infty} kx^2 e^{-x} dx = 1$ ; i.e.,  $k[x^2(-e^{-x}) - 2x(e^{-x}) + 2(-e^{-x})]_0^{\infty} = 1$   
 i.e.,  $2k = 1 \quad \therefore k = 1/2$

$$\mu'_r = \frac{1}{2} \int_0^{\infty} x^{r+2} e^{-x} dx = \frac{-1}{2} [\{x^{r+2} + (r+2)x^{r+1} + (r+2)(r+1)x^r + \dots (r+2)(r+1)\dots 2.1\}e^{-x}]_0^{\infty} = \frac{1}{2} [r+2]$$

$$E(X) = \mu'_1 = \frac{1}{2} \lfloor 3 \rfloor = 3; \quad E(X^2) = \mu'_2 = \frac{1}{2} \lfloor 4 \rfloor = 12$$

$$V(X) = E(X^2) - E^2(X) = 12 - 9 = 3.$$

(ii)  $X$  is  $U(0, 10)$        $\therefore f_X(x) = \frac{1}{10}$

$$(1) P(X < 3) = P(0 < X < 3) = \int_0^3 \frac{1}{10} dx = \frac{3}{10}$$

$$P(X > 7) = P(7 < X < 10) = \int_7^{10} \frac{1}{10} dx = \frac{3}{10}$$

$$P(2 < X < 5) = \int_2^5 \frac{1}{10} dx = \frac{3}{10}$$

(2) Since  $X$  is a continuous RV,  $P(X = 7) = 0$

11. (b) (i)  $X \equiv$  Number of busy phones;  $p = 0.1, q = 0.9, n = 4$

$$(1) P(X = 4) = 4C_4 \cdot (0.1)^4 \cdot (0.9)^0 = 0.0001$$

$$\begin{aligned} (2) P(X = 2, 3 \text{ or } 4) &= 1 - \{P(X = 0) + P(X = 1)\} \\ &= 1 - \{4C_0 \cdot (0.1)^0 \cdot (0.9)^4 + 4C_1 \cdot (0.1)^1 \cdot (0.9)^3\} \\ &= 1 - (0.6561 + 0.2916) = 0.0523 \end{aligned}$$

(ii) If  $X$  is a continuous RV with  $f(x) = x^{n-1} e^{-x} / \lceil(n\rceil); x \geq 0$ , then  $X$  is said to follow Gamma distribution with parameter  $n > 0$

$$\begin{aligned} M_X(t) &= \int_0^\infty e^{tx} x^{n-1} e^{-x} / \lceil(n\rceil) dx = \frac{1}{\lceil(n\rceil)} \int_0^\infty x^{n-1} e^{-(1-t)x} dx \\ &= \frac{1}{\lceil(n\rceil)} \int_0^\infty \left( \frac{y}{1-t} \right)^{n-1} e^{-y} \frac{dy}{1-t} \\ &= \frac{1}{\lceil(n\rceil)} \cdot \frac{1}{(1-t)^n} \int_0^\infty y^{n-1} e^{-y} dy = (1-t)^{-n} \end{aligned}$$

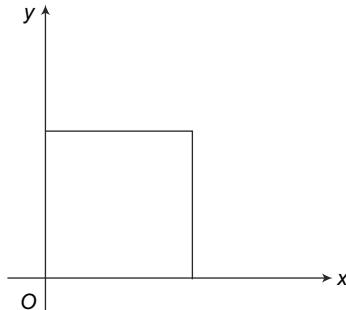
$$M_X(t) = 1 + \frac{n}{1} t + \frac{n(n+1)}{2} t^2 + \dots \quad (1)$$

$$E(x) = \text{coefficient of } \frac{t}{1} \text{ in (1)} = n$$

$$E(X^2) = \text{coefficient of } \frac{t^2}{2} \text{ in (1)} = n(n+1)$$

$$\therefore \text{Var}(X) = E(X^2) - E^2(X) = n$$

12. (a) (i)



**Fig. 1**

$$f_{XY}(x, y) = 16 abxy; \quad 0 < x, y < 1$$

$$\begin{aligned} E(U) &= E(X + Y) = \int_0^1 \int_0^1 (x + y) 16 abxy dx dy \\ &= 16 ab \int_0^1 \left( \frac{x^3}{3} y + \frac{x^2}{2} y^2 \right)_0^1 dy \\ &= 16 ab \int_0^1 \left( \frac{y}{3} + \frac{y^2}{2} \right) dy = 16 ab \left( \frac{y^2}{6} + \frac{y^3}{6} \right)_0^1 = \frac{16}{3} ab \end{aligned}$$

$$\begin{aligned} E(V) &= E(X - Y) = \int_0^1 \int_0^1 (x - y) 16 abxy dx dy \\ &= 16 ab \int_0^1 \left( \frac{x^3}{3} y - \frac{x^2}{2} y^2 \right)_0^1 dy = 16 ab \int_0^1 \left( \frac{y}{3} - \frac{y^2}{2} \right) dy \\ &= 16 ab \left( \frac{y^2}{6} - \frac{y^3}{6} \right)_0^1 = 0 \end{aligned}$$

$$\begin{aligned} E(UV) &= E(X^2 - Y^2) = \int_0^1 \int_0^1 (x^2 - y^2) 16 abxy dx dy \\ &= 16 ab \int_0^1 \left( \frac{x^4}{4} y - \frac{x^2}{2} y^3 \right)_0^1 dy \\ &= 16 ab \int_0^1 \left( \frac{y}{4} - \frac{y^3}{2} \right) dy = 16 ab \left( \frac{y^2}{8} - \frac{y^4}{8} \right)_0^1 = 0 \end{aligned}$$

$$C(U, V) = E(UV) - E(U) \cdot E(V) = 0; \quad \therefore r_{UV} = 0;$$

$\therefore U$  and  $V$  are uncorrelated.

- (ii) Statement of Lindeberg-Levy's form of CLT—available in the book. Proof is beyond the scope of the syllabus and the book.
12. (b) (i)  $(\bar{X}, \bar{Y})$  is the point of intersection of the two regression lines.

$$3\bar{x} + 12\bar{y} = 19 \quad (1)$$

$$9\bar{x} + 3\bar{y} = 46 \quad (2)$$

Solving (1) and (2), we get  $\bar{X} = 5$  and  $\bar{Y} = \frac{1}{3}$

Rewriting the regression equations, we have  $y = -\frac{1}{4}x + \frac{19}{12}$  (1)

and  $x = -\frac{1}{3}y + \frac{46}{9}$  (2)

$$\therefore b_{YX} = -\frac{1}{4} \text{ and } b_{XY} = -\frac{1}{3}$$

$$r_{XY}^2 = b_{XY} \cdot b_{YX} = \frac{1}{12} \quad \therefore r_{XY} = \frac{-1}{\sqrt{12}} \text{ or } -\frac{1}{2\sqrt{3}}$$

(– sign is taken, since both  $b_{XY}$  and  $b_{YX}$  are –ve).

- (ii) This is a worked example in the book.
13. (a) (i) This is a worked example in the book.
- (ii) Properties of ACF and CCF are available in the book.
13. (b) (i)  $E\{X(t)\} = \cos t E(A) + \sin t E(B)$ , where the probability distribution of  $A$  or  $B$  is given below:

$A$ or $B$	-2	1
$P(A)$ or $P(B)$	$\frac{1}{3}$	$\frac{2}{3}$

$$\therefore E(A) = E(B) = -\frac{2}{3} + \frac{2}{3} = 0 \text{ and}$$

$$E(A^2) = E(B^2) = \frac{4}{3} + \frac{2}{3} = 2$$

$E(AB) = E(A) \cdot E(B) = 0$ , since  $A$  and  $B$  are independent.

$$\therefore E\{X(t)\} = 0$$

$$E\{X^2(t)\} = \cos^2 t \cdot E(A^2) + \sin^2 t \cdot E(B^2) + 2 \sin t \cos t \cdot E(AB) = 2$$

$$\therefore V\{X(t)\} = 2$$

Since  $E\{X(t)\}$  and  $V\{X(t)\}$  are constants,  $\{X(t)\}$  is a WSS process.

- (ii) Definition and property of Poisson process are available in the book.
14. (a) (i) Definition and property of spectral density are available in the book.

$$(ii) f_\theta(\theta) = \frac{1}{2\pi} \text{ in } 0 < \theta < 2\pi$$

$$\begin{aligned}
 E\{X(t_1)\} \cdot Y(t_2) &= E\{A \cos(\omega t_1 + \theta) \cdot B \sin(\omega t_2 + \theta)\} \\
 &= AB \cdot \frac{1}{2\pi} \int_0^{2\pi} \cos(\omega t_1 + \theta) \cdot \sin(\omega t_2 + \theta) d\theta \\
 &= \frac{AB}{4\pi} \int_0^{2\pi} [\sin\{\omega(t_1 + t_2) + 2\theta\} - \sin\{\omega(t_1 - t_2)\}] d\theta \\
 &= \frac{AB}{4\pi} \left[ -\frac{1}{2} \cos\{\omega(t_1 + t_2) + 2\theta\} - \sin\{\omega(t_1 - t_2)\}\theta \right]_0^{2\pi} \\
 &= -\frac{AB}{2} \sin\omega(t_1 - t_2) \\
 \text{i.e., } R_{XY}(t_1, t_2) &= R_{XY}(t_1 - t_2) = -\frac{AB}{2} \sin\omega(t_1 - t_2)
 \end{aligned}$$

$$\therefore R_{XY}(\tau) = -\frac{AB}{2} \sin\omega\tau$$

14. (b) (i) Wiener–Khinchine theorem—statement and proof are available in the book.

$$\begin{aligned}
 \text{(ii)} \quad R_{XY}(\tau) &= F^{-1}\{S_{XY}(\omega)\} = \frac{1}{2\pi} \int_{-\infty}^{\infty} \left( a + i\frac{b}{a}\omega \right) e^{i\tau\omega} d\omega \\
 &= \frac{1}{2\pi} \int_{-a}^a \left( a + i\frac{b\omega}{a} \right) (\cos\tau\omega + i\sin\tau\omega) d\omega \\
 &= \frac{1}{\pi} \int_0^a \left( a \cos\tau\omega - \frac{b}{a}\omega \sin\tau\omega \right) d\omega \\
 &= \frac{1}{\pi} \left[ a \left( \frac{\sin\tau\omega}{\tau} \right)_0^a - \frac{b}{a} \left\{ \omega \left( -\frac{\cos\tau\omega}{\tau} \right) + \frac{\sin\tau\omega}{\tau^2} \right\}_0^a \right] \\
 &= \frac{1}{\pi} \left[ a \frac{\sin a\tau}{\tau} + b \frac{\cos a\tau}{\tau} - \frac{b}{a} \frac{\sin a\tau}{\tau^2} \right]
 \end{aligned}$$

$$\begin{aligned}
 \text{15. (a) (i)} \quad H(\omega) &= F\{h(t)\} = \int_{-\infty}^{\infty} h(t) e^{-i\omega t} dt = \int_0^{\infty} 2e^{-t} \cdot e^{-i\omega t} dt \\
 &= 2 \left[ \frac{e^{-(1+i\omega)t}}{-(1+i\omega)} \right]_0^{\infty} = \frac{2}{1+i\omega}
 \end{aligned}$$

$$S_{XX}(\omega) = F\{R_{XX}(\tau)\} = \int_{-\infty}^{\infty} e^{-2|\tau|} e^{-i\omega\tau} d\tau$$

$$\begin{aligned}
 &= \int_{-\infty}^0 e^{2\tau} \cdot e^{-i\omega\tau} d\tau + \int_0^\infty e^{-2\tau} \cdot e^{-i\omega\tau} d\tau \\
 &= \left[ \frac{e^{(2-i\omega)\tau}}{2-i\omega} \right]_{-\infty}^0 + \left[ \frac{e^{-(2+i\omega)\tau}}{(2+i\omega)} \right]_0^\infty \\
 &= \frac{1}{2-i\omega} + \frac{1}{2+i\omega} = \frac{4}{4+\omega^2}
 \end{aligned}$$

$$S_{YY}(\omega) = |H(\omega)|^2 \quad S_{XX}(\omega) = \frac{4}{1+\omega^2} \cdot \frac{4}{4+\omega^2} = \frac{16}{(\omega^2+1)(\omega^2+4)}$$

$$\begin{aligned}
 \text{(ii)} \quad S_{NN}(\omega) &= F\{R_{NN}(\tau)\} = P \left[ \int_{-\infty}^0 e^{3\tau} e^{-i\omega\tau} d\tau + \int_0^\infty e^{-3\tau} e^{-i\omega\tau} d\tau \right] \\
 &= P \left[ \left\{ \frac{e^{(3-i\omega)\tau}}{3-i\omega} \right\}_{-\infty}^0 + \left\{ \frac{e^{-(3+i\omega)\tau}}{-(3+i\omega)} \right\}_0^\infty \right] \\
 &= P \left[ \frac{1}{3-i\omega} + \frac{1}{3+i\omega} \right] = \frac{6P}{\omega^2+9}
 \end{aligned}$$

15. (b) (i) A standard property; proof is available in the book.  
(ii) A standard result, the derivation is available in the book.

**BE/BTech Degree Examination, May/June 2013**

**Fourth Semester**

**Computer Science and Engineering**

**MA2262/MA44/MA1252/10177 PQ 401/080250008**

**Probability and Queueing Theory**

**(Common to Information Technology)**

**(Regulation 2008/2010)**

**Time: Three hours**

**Maximum: 100 marks**

**Answer ALL Questions.**

**Part-A (10 × 2 = 20 marks)**

1. If  $X$  and  $Y$  are two independent random variables with variances 2 and 3, find the variance of  $3X + 4Y$ .
2. State memoryless property of exponential distribution.
3. If the joint pdf of  $(X, Y)$  is given by  $f(x, y) = 2$ , in  $0 \leq x \leq y \leq 1$ , find  $E(X)$ .
4. State the Central Limit theorem.
5. Define wide sense stationary process.
6. If the transition probability matrix (tpm) of a Markov chain is  $\begin{pmatrix} 0 & 1 \\ 1 & 1 \\ 2 & 2 \end{pmatrix}$ , find the steady-state distribution of the chain.
7. What are the characteristics of a queueing system?
8. What is the probability that a customer has to wait more than 15 minutes to get his service completed in a M/M/1 queueing system, if  $\lambda = 6$  per hour and  $\mu = 10$  per hour?
9. State Polaczek–Khinchine formula.
10. Define closed network of a queueing system.

**Part-B (5 × 16 = 80 marks)**

11. (a) (i) A continuous random variable has the pdf  $f(x) = kx^4$ ,  $-1 < x < 0$ . Find the value of  $k$  and also  $P\left\{X > \left(-\frac{1}{2}\right)\middle/\left|X < \left(-\frac{1}{4}\right)\right.\right\}$ . (8)  
(ii) Find the moment-generating function of uniform distribution. Hence, find its mean and variance. (8)

Or

- (b) (i) Find the moment-generating function and  $r^{\text{th}}$  moment for the distribution whose pdf is  $f(x) = Ke^{-x}$ ,  $0 \leq x \leq \infty$ . Hence, find the mean and variance. (8)
- (ii) In a large consignment of electric bulbs, 10 percent are defective. A random sample of 20 is taken for inspection. Find the probability that (1) all are good bulbs (2) at most there are 3 defective bulbs, and (3) exactly there are 3 defective bulbs. (8)
12. (a) (i) The joint probability density function of a two-dimensional random variable  $(X, Y)$  is  $f(x, y) = \frac{1}{8}(6 - x - y)$ ,  $0 < x < 2$ ,  $2 < y < 4$ . Find (1)  $P(X < 1 \cap Y < 3)$  (2)  $P(X + Y < 3)$  (3)  $P(X < 1/Y < 3)$ . (8)
- (ii) If  $X$  and  $Y$  each follow an exponential distribution with parameter  $I$  and the independent, find the pdf of  $U = X - Y$ . (8)

Or

- (b) (i) The marks obtained by 10 students in Mathematics and Statistics are given below. Find the correlation coefficient between the two subjects.

Marks in mathematics	75	30	60	80	53	35	15	40	38	48
Marks in statistics	85	45	54	91	58	63	35	43	45	44

- (ii) A distribution with unknown mean  $\mu$  has variance equal to 1.5. Use the central limit theorem to find how large a sample should be taken from the distribution in order that the probability will be at least 0.95 that the sample mean will be within 0.5 of the population mean. (8)
13. (a) (i) Show that the process  $X(t) = A \cos \lambda t + B \sin \lambda t$  is wide sense stationary, if  $E(A) = E(B) = 0$ ,  $E(A^2) = E(B^2)$  and  $E(AB) = 0$ , where  $A$  and  $B$  are random variables.
- (ii) A gambler has Rs. 2. He bets Re. 1 at a time and wins Re. 1 with probability of 1/2. He stops playing if he loses Rs. 2 or wins Rs. 4. (1) What is the tpm of the related Markov chain? (2) What is the probability that he has lost his money at the end of 5 plays? (8)

Or

- (b) (i) Find the nature of the states of the Markov chain with the tpm
- $$P = \begin{bmatrix} 0 & 1 & 0 \\ \frac{1}{2} & 0 & \frac{1}{2} \\ 0 & 1 & 0 \end{bmatrix} \quad (8)$$
- (ii) Prove that the differences of two independent Poisson processes is not a Poisson process. (4)
- (iii) Prove that the Poisson process is a Markov process. (4)

14. (a) (i) Derive (1)  $L_s$ , average number of customers in the system (2)  $L_q$ , average number of customers in the queue for the queueing model (M/M/1): (N/FIFO). (8)
- (ii) There are three typists in an office. Each typist can type an average of 6 letters per hour. If letters arrive for being typed at the rate of 15 letters per hour, what fraction of time will all the typists be busy? What is the average number of letters waiting to be typed? (Assume Poisson arrivals and exponential service times) (8)
- Or
- (b) Customers arrive at a one-man barber shop according to a Poisson process with a mean inter-arrival time of 20 minutes. Customers spend an average of 15 minutes in the barber chair. The service time is exponentially distributed. If an hour is used as a unit of time, then
- (i) What is the probability that a customer need not wait for a hair cut?
  - (ii) What is the expected number of a customer in the barber shop and in the queue?
  - (iii) How much time can a customer expect to spend in the barber shop?
  - (iv) Find the average time that a customer spends in the queue.
  - (v) Estimate the fraction of the day that the customer will be idle.
  - (vi) What is the probability that there will be 6 or more customers?
  - (vii) Estimate the percentage of customers who have to wait prior to getting into the barber's chair. (16)
15. (a) An automatic car-wash facility operates with only one bay. Cars arrive according to a Poisson process at the rate of 4 cars per hour and may wait in the facility's parking lot if the bay is busy. The service time for all cars is constant and equal to 10 minutes. Determine  $L_s$ ,  $L_q$ ,  $W_s$  and  $W_q$ . (16)
- (b) Consider a system of two servers where customers from outside the system arrive at server 1 at a Poisson rate of 4 and at server 2 at a Poisson rate of 5. The service rates for servers 1 and 2 are 8 and 10 respectively. A customer upon completion of service at server 1 is likely to go to server 2 or leave the system; whereas a departure from server 2 will go 25 percent of the time to server 1 and will depart the system otherwise. Determine the limiting probabilities  $L_s$  and  $W_s$ . (16)

## Solutions

### Part-A

1. Let  $Z = 3X + 4Y \quad \therefore Z - E(Z) = 3\{X - E(X)\} + 4\{Y - E(Y)\}$
- $$\therefore V(Z) = E\{Z - E(Z)\}^2 = 9E\{X - E(X)\}^2 + 16E\{Y - E(Y)\}^2 + 24[E\{X - E(X)\}\{Y - E(Y)\}^2]$$

$$= 9 \times 2 + 16 \times 3 + 24 \times 0 \quad (\because X \text{ and } Y \text{ are independent}) \\ = 66$$

2. Statement available in the book.  
3.

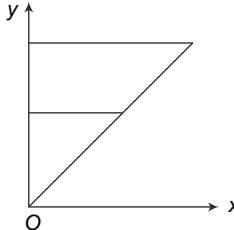


Fig. 1

$$E(X) = \int_0^1 \int_0^y 2x \, dx \, dy = \int_0^1 y^2 \, dy = \frac{1}{3}$$

4. Statement of CLT available in the book.  
5. Definition of WSS available in the book.  
6. Let the steady-state distribution be  $(\pi_1, \pi_2)$ .

$$\text{Then } (\pi_1, \pi_2) = \begin{bmatrix} 0 & 1 \\ \frac{1}{2} & \frac{1}{2} \end{bmatrix}(\pi_1, \pi_2); \text{ viz., } \frac{1}{2}\pi_2 = \pi_1 \text{ and } \pi_1 + \frac{1}{2}\pi_2 = \pi_2 \\ \text{or } \pi_1 = \frac{1}{2}\pi_2 \quad (1)$$

$$\text{Also, } \pi_1 + \frac{\pi_2}{2} = 1 \quad (2); \text{ solving (1) and (2); } \pi_1 = \frac{1}{3} \text{ and } \pi_2 = \frac{2}{3}$$

$$\therefore \text{ the required distribution is } \left( \frac{1}{3}, \frac{2}{3} \right)$$

7. Answer for this question is available in the book.

8.  $\lambda = 6/\text{hour}$  and  $\mu = 10/\text{hour}$

$$P(W_S > t) = \int_t^\infty (\mu - \lambda) e^{-(\mu - \lambda)w} dw = e^{-(\mu - \lambda)t}$$

$$\therefore P\left(W_S > \frac{1}{4}\right) = e^{-\frac{(10-6)}{4}} = \frac{1}{e}$$

9. Statement of P.K. formula is available in the book.  
10. Definition of closed network of a queueing system is available in the book.

### Part-B

11. (a) (i)  $\int_{-1}^0 kx^4 dx = 1$ , i.e.,  $k \left( \frac{x^5}{5} \right) \Big|_{-1}^0 = 1$ ; i.e.,  $\frac{k}{5} = 1 \quad \therefore k = 5$

$$\begin{aligned}
 P\left\{x > -\frac{1}{2} \middle/ x < -\frac{1}{4}\right\} &= \frac{P\left\{-\frac{1}{2} < X < -\frac{1}{4}\right\}}{P\left\{-1 < X < -\frac{1}{4}\right\}} = \int_{-\frac{1}{2}}^{-\frac{1}{4}} 5x^4 dx \div \int_{-1}^{-\frac{1}{4}} 5x^4 dx \\
 &= (x^5)_{-1/2}^{-1/4} \div (x^5)_{-1}^{-1/4} \\
 &= \left(-\frac{1}{1024} + \frac{1}{32}\right) \div \left(-\frac{1}{1024} + 1\right) = \frac{1}{33}
 \end{aligned}$$

(ii) Let  $X$  follow  $U(a, b)$   $\therefore f(x) = \frac{1}{b-a}$

$$\begin{aligned}
 M_X(t) &= \int_a^b \frac{1}{b-a} e^{tx} dx = \frac{1}{b-a} \frac{(e^{bt} - e^{at})}{t} \\
 &= \frac{1}{b-a} \cdot \frac{1}{t} \left[ \begin{aligned} &\left\{ 1 + \frac{bt}{1} + \frac{b^2 t^2}{2} + \frac{b^3 t^3}{3} + \dots \right\} \\ &- \left\{ 1 + \frac{at}{1} + \frac{a^2 t^2}{2} + \frac{a^3 t^3}{3} + \dots \right\} \end{aligned} \right] \\
 &= 1 + \left( \frac{b+a}{2} \right) t + \frac{b^2 + ba + a^2}{6} \cdot t^2 + \dots \tag{1}
 \end{aligned}$$

$$E(X) = \text{coefficient of } \frac{t}{1} \text{ in (1)} = \frac{b+a}{2}$$

$$E(X^2) = \text{coefficient of } \frac{t^2}{2} \text{ in (1)} = \frac{b^2 + ba + a^2}{3}$$

$$\begin{aligned}
 V(X) &= E(X^2) - E^2(X) \\
 &= \left( \frac{b^2 + ba + a^2}{3} \right) - \left( \frac{b^2 + 2ba + a^2}{4} \right) = \frac{b^2 - 2ba + a^2}{12} \\
 &= \frac{(b-a)^2}{12}
 \end{aligned}$$

11. (b) (i)  $\int_{R_X} f(x) dx = 1$ ; i.e.,  $k \int_0^\infty e^{-x} dx = 1 \quad \therefore k = 1$ ; i.e.,  $f(x) = e^{-x}; x \geq 0$

$$M_X(t) = \int_0^\infty e^{-x} e^{tx} dx = \left\{ \frac{e^{-(1-t)x}}{-(1-t)} \right\}_0^\infty = \frac{1}{1-t}$$

$$M_X(t) = 1 + t + t^2 + \dots + t^r + \dots \quad (1)$$

$$\therefore \mu_r' = \text{coefficient of } \frac{t^r}{|r|} \text{ in (1)} = |r|$$

$$E(X) = \mu_1' = 1; E(X^2) = \mu_2' = 2 \quad \therefore V(X) = E(X^2) - E^2(X) = 1$$

(ii)  $p = P(\text{a bulb is defective}) = 0.1$  and  $q = 0.9$ ;  $n = 20$

(i)  $P(\text{all the bulbs are good}) = P(\text{no defective bulb})$

$$= P(X = 0) = 20C_0 p^0 q^{20-0} = (0.9)^{20}$$

(ii)  $P(\text{at most 3 bulbs are defective}) = P(X = 0, 1, 2, 3)$

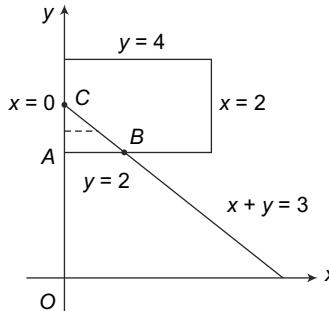
$$= \sum_{r=0}^3 20C_r p^r q^{20-r}$$

$$= [(0.9)^{20} + 20 \times (0.1) \times (0.9)^{19} + 190 \times (0.1)^2 \times (0.9)^{18} \\ + 1140 \times (0.1)^3 \times (0.9)^{17}]$$

(iii)  $P(\text{exactly 3 bulbs are defective})$

$$= P(X = 3) = 20C_3 \times (0.1)^3 \times (0.9)^{17}.$$

12. (a) (i)



**Fig. 2**

$$(1) P(X < 1 \cap Y < 3) = P(0 < X < 1 \text{ and } 2 < Y < 3)$$

$$= \int_2^1 \int_0^1 \frac{1}{8}(6-x-y) dx dy = \frac{1}{8} \int_2^1 \left( \frac{11}{2} - y \right) dy$$

$$= \frac{1}{8} \left( \frac{11}{2}y - \frac{y^2}{2} \right)_2^3 = \frac{3}{8}$$

$$(2) P(X + Y < 3) = \iint_{\Delta ABC} f(x, y) dx dy = \frac{1}{8} \int_2^3 \int_0^{3-y} (6-x-y) dx dy$$

$$= \frac{1}{16} \int_2^3 (3-y)(9-y) dy = \frac{1}{16} \left( 27y - 6y^2 + \frac{1}{3}y^3 \right)_2^3 = \frac{5}{24}$$

$$(3) P(X < 1 \cap Y < 3) = \frac{P(X < 1 \cap Y < 3)}{P(2 < Y < 3)} = \frac{3}{8} \div \int_{2/0}^{3/2} \int \frac{1}{8} (6 - x - y) dx dy \\ = \frac{3}{8} \div \int_{2/2}^{3/8} \frac{1}{8} (10 - 2y) dy = 3 \div (10y - y^2)_2^3 = \frac{3}{5}$$

- (ii) This is a worked example in the book.
12. (b) (i)  $x$  = Marks in Mathematics;  $y$  = Marks in Statistics  
 $\Sigma x = 471$ ;  $\Sigma y = 563$ ;  $\Sigma x^2 = 26,132$ ;  $\Sigma y^2 = 34,815$ ;  $\Sigma xy = 29,591$

$$r_{XY} = \frac{n\Sigma xy - \Sigma x \cdot \Sigma y}{\sqrt{\{n\Sigma x^2 - (\Sigma x)^2\} \{n\Sigma y^2 - (\Sigma y)^2\}}} \\ = \frac{295910 - 265173}{\sqrt{(261320 - 221841)(348150 - 316969)}} \\ = \frac{30737}{\sqrt{39479 \times 31181}} = 0.876$$

- (ii) This is a worked example in the book.
13. (a) (i) This is a worked example in the book.  
(ii) This is a worked example in the book.
13. (b) (i) This is a worked example in the book.  
(ii) This is a standard property, proof is available in the book.  
(iii) This is a standard property, proof is available in the book.
14. (a) (i) Derivation is available for this model in the book.  
(ii) This is a worked example in the book.

14. (b)  $\frac{1}{\lambda} = 20 \quad \therefore \quad \lambda = \frac{1}{20} / \text{min}; \quad \frac{1}{\mu} = 15 \quad \therefore \quad \mu = \frac{1}{15} / \text{min}$

(i)  $P(\text{the customers need not wait for a haircut})$   
 $= P(\text{No customer in the system})$

$$= P_0 = 1 - \frac{\lambda}{\mu} = 1 - \frac{1}{20} \times 15 = \frac{1}{4}$$

(ii)  $E(\text{Number of customers in the shop}) = E(N_S) = \frac{\lambda}{\mu - \lambda} = \frac{\frac{1}{20}}{\frac{1}{15} - \frac{1}{20}} = 3$

$E(\text{Number of customers in the queue}) =$

$$E(N_Q) = \frac{\lambda^2}{\mu(\mu - \lambda)} = \frac{\frac{1}{400}}{\frac{1}{15} \left( \frac{1}{15} - \frac{1}{20} \right)} = 2.25$$

$$(iii) E(\text{waiting time in the shop}) = E(W_S) = \frac{1}{\mu - \lambda} = \frac{1}{\frac{1}{15} - \frac{1}{20}} = 60 \text{ min}$$

$$(iv) E(\text{waiting time in the queue}) =$$

$$E(W_Q) = \frac{\lambda}{\mu(\mu - \lambda)} = \frac{\frac{1}{20}}{\frac{1}{15}\left(\frac{1}{15} - \frac{1}{20}\right)} = 45 \text{ min}$$

$$(v) \text{ The fraction of the day when the barber (it is not customer as given in the problem) will be idle} = P_0 = \frac{1}{4}$$

$$(vi) P(N_S \geq 6) = 1 - \{P_0 + P_1 + P_2 + P_3 + P_4 + P_5\}$$

$$\begin{aligned} &= 1 - \left\{1 - \left(\frac{\lambda}{\mu}\right)\right\} \left\{1 + \frac{\lambda}{\mu} + \left(\frac{\lambda}{\mu}\right)^2 + \dots + \left(\frac{\lambda}{\mu}\right)^5\right\} \\ &= 1 - \left\{1 - \frac{\lambda}{\mu}\right\} \times \frac{1 - \left(\frac{\lambda}{\mu}\right)^6}{1 - \frac{\lambda}{\mu}} = \left(\frac{\lambda}{\mu}\right)^6 = \left(\frac{3}{4}\right)^6 = 0.178 \end{aligned}$$

$$(vii) \text{ Fraction of customers who have to wait in the queue}$$

$$= \frac{E(N_Q)}{E(N_S)} = \frac{2.25}{3} = 0.75$$

$$\therefore \% \text{ of customers who have to wait in the queue} = 75.$$

15. (a) We shall assume that the parking lot is large enough to accommodate any number of cars.

$$\lambda = 4/\text{hour} \text{ or } \frac{1}{15}/\text{min} ; E(T) = 10 \text{ and } V(T) = 0$$

$$\begin{aligned} (1) \text{ By P.K. formula, } L_S &= E(N_S) = \lambda \cdot E(T) + \frac{\lambda^2 \{V(T) + E^2(T)\}}{2\{1 - \lambda E(T)\}} \\ &= \frac{1}{15} \times 10 + \frac{\frac{1}{225} \{0 + 100\}}{2 \left\{1 - \frac{1}{16} \times 10\right\}} = \frac{4}{3} \end{aligned}$$

$$(2) L_Q = E(N_Q) = E(N_S) - \frac{\lambda}{\mu} = \frac{2}{3}$$

$$(3) E(W_S) = \frac{1}{\lambda} E(N_S) = 20 \text{ min}$$

$$(4) \quad E(W_Q) = \frac{1}{\lambda} E(N_Q) = 10 \text{ min}$$

15. (b)  $r_1 = 4, r_2 = 5; P_{11} = 0, P_{12} = \frac{1}{2}$  [Note: It is to be assumed that a customer upon completion of service at  $S_1$  is *equally* likely to go to  $S_2$  or leave the system. The word ‘equally’ is missing in the question.]

$$P_{21} = \frac{1}{4}, P_{22} = 0 \text{ and } \mu_1 = 8, \mu_2 = 10$$

Jackson's flow balance equation for this open network is

$$\lambda y = r_j + \sum_{i=1}^2 \lambda_i P_{ij} \quad (j = 1, 2); \quad \therefore \quad \lambda_1 = 4 + \frac{1}{4} \lambda_2 \text{ and } \lambda_2 = 5 + \frac{1}{2} \lambda_1$$

Solving these equations,  $\lambda_1 = 6$  and  $\lambda_2 = 8$

$$L_S = E(N_S) = \frac{\lambda_1}{\mu_1 - \lambda_1} + \frac{\lambda_2}{\mu_2 - \lambda_2} = \frac{6}{8 - 6} + \frac{8}{10 - 8} = 7$$

$$E(W_S) = \frac{1}{\lambda} E(N_S) = \frac{1}{r_1 + r_2} E(N_S) = \frac{7}{9}$$

**BE/BTech Degree Examination, May/June 2014**  
**Fourth Semester**  
**Electronics and Communication Engineering**  
**MA2261/MA45/MA1253/10177 PR 401/080380009—Probability and  
Random Processes**  
**(Common to Biomedical Engineering)**  
**(Regulation 2008/2010)**

**Time: Three hours**      **Maximum: 100 marks**  
**Use of Statistical Tables is permitted.**  
**Answer ALL Questions.**

**Part-A (10 × 2 = 20 marks)**

1.  $X$  and  $Y$  are independent random variables, with variances of 2 and 3. Find the variance of  $3X + 4Y$ .
  2. A continuous random variable  $X$  has a probability density function (pdf)  
$$f(x) = \begin{cases} 3x^2; & 0 \leq x \leq 1 \\ 0; & \text{otherwise} \end{cases}$$
. Find  $k$  such that  $P(X > k) = 0.5$ .
  3. State the Central Limit theorem for iid random variables.
  4. State the basic properties of joint distribution of  $(X, Y)$  when  $X$  and  $Y$  are random variables.
  5. State the properties of an ergodic process.
  6. Explain any two applications of a binomial process.
  7. Define cross-correlation function and state any two of its properties.
  8. Find the variance of the stationary ergodic process  $\{X(t)\}$  whose autocorrelation function is given by  $R_{XX}(\tau) = 25 + 4/(1 + 6\tau^2)$ .
  9. Define a system. When is it called a linear system?
  10. Define band-limited white noise.

## **Part-B ( $5 \times 16 = 80$ marks)**

11. (a) (i) Define the moment-generating function (MGF) of a random variable. Derive the MGF, mean, variance and the first four moments of a Gamma distribution. (8)

(ii) Describe binomial  $B(n, p)$  distribution and obtain the moment-generating function. Hence, compute (1) the first four moments, and (2) the recursion relation for the central moments. (8)

Or

- (b) (i) A random variable  $X$  has the following probability distribution:

$X:$	0	1	2	3	4	5	6	7
$P(x):$	0	$K$	$2K$	$2K$	$3K$	$K^2$	$2K^2$	$7K^2 + K$

Find

- (1) the value of  $K$
- (2)  $P(1.5 < X < 4.5/X > 2)$ , and
- (3) the smallest value of  $n$  for which  $P(X \leq n) > 1/2$ . (8)

- (ii) Find the MGF of a random variable  $X$  having the pdf

$$f(x) = \begin{cases} \frac{x}{4e^{+x/2}}; & x > 0 \\ 0; & \text{otherwise} \end{cases}$$

Also deduce the first four moments about the origin. (8)

12. (a) If the joint pdf of a two-dimensional random variable  $(X, Y)$  is given by

$$f(x, y) = \begin{cases} x^2 + \frac{xy}{3}, & 0 < x < 1; 0 < y < 2 \\ 0, & \text{otherwise} \end{cases}$$

Find

- (i)  $P\left(X > \frac{1}{2}\right)$ ,
- (ii)  $P(Y < X)$ ,
- (iii)  $P[X + Y \geq 1]$ , and
- (iv) The conditional density functions. (16)

Or

- (b) (i) The joint pdf of the random variable  $(X, Y)$  is

$f(x, y) = 3(x + y)$   $0 \leq x \leq 1, 0 \leq y \leq 1, x + y \leq 1$ , find  $\text{Cov}(X, Y)$ . (8)

- (ii) Marks obtained by 10 students in Mathematics ( $x$ ) and statistics ( $y$ ) are given below:

$x:$	60	34	40	50	45	40	22	43	42	64
$y:$	75	32	33	40	45	33	12	30	34	51

Find the two regression lines. Also find  $y$  when  $x = 55$ . (8)

13. (a) (i) A process  $\{X(t)\}$  has a probability distribution under a certain

$$\text{condition given by } P\{X(t) = n\} = \begin{cases} \frac{(at)^{n-1}}{(1+at)^{n+1}}, & n = 1, 2 \\ \frac{at}{1+at}, & n = 0 \end{cases}.$$

Find the mean and variance of the process. Is the process first-order stationary?

- (ii) If the WSS process  $\{X(t)\}$  is given by  $X(t) = 10 \cos(100t + \theta)$ , where  $\theta$  is uniformly distributed over  $(-\pi, \pi)$ , prove that  $\{X(t)\}$  is correlation ergodic.

Or

- (b) (i) If the process  $\{X(t); t \geq 0\}$  is a Poisson process with parameter  $\lambda$ , obtain  $P\{X(t) = n\}$ . Is the process first-order stationary? (10)
- (ii) Prove that a random telegraph signal process  $Y(t) = aX(t)$  is a Wide Sense Stationary Process where  $a$  is a random variable which is independent of  $X(t)$  and assumes values  $-1$  and  $+1$  with equal probability and  $R_{XX}(t_1, t_2) = e^{-2\lambda|t_1 - t_2|}$ . (6)
14. (a) (i) Find the mean and autocorrelation of the Poisson process. (8)
- (ii) Prove that the random processes  $X(t)$  and  $Y(t)$  defined by  $X(t) = A \cos \omega t + B \sin \omega t$  and  $Y(t) = B \cos \omega t - AB \sin \omega t$  are jointly wide sense stationary. (8)

Or

- (b) State and prove the Weiner–Khintchine Theorem. (16)
15. (a) (i) Show that if the input  $\{X(t)\}$  is a WSS process for a linear system then output  $\{Y(t)\}$  is a WSS process. Also find  $R_{YY}(\tau)$ . (8)
- (ii) If  $\{X(t)\}$  is the input voltage to a circuit and  $\{Y(t)\}$  is the output voltage,  $\{X(t)\}$  is a stationary random process with  $\mu_X = 0$  and  $R_{xx}(\tau) = e^{-\alpha|\tau|}$ . Find the mean  $\mu_Y$  and power spectrum  $S_{YY}(\omega)$  of the output if the power transfer function is given by  $H(\omega) = \frac{R}{R + iL\omega}$ . (8)

Or

- (b) (i) If  $Y(t) = A \cos(\omega t + \theta) - N(t)$ , where  $A$  is a constant,  $\theta$  is a random variable with a uniform distribution in  $(-\pi, \pi)$  and  $\{N(t)\}$  is a bandlimited Gaussian white noise with power spectral density

$$S_{NN}(\omega) = \begin{cases} \frac{N_0}{2}, & \text{for } |\omega_0 - \omega| < \omega_B \\ 0, & \text{elsewhere} \end{cases}$$

Find the power spectral density  $S_{YY}(\omega)$ . Assume that  $\{N(t)\}$  and  $\theta$  are independent. (10)

- (ii) A system has an impulse response  $h(t) = e^{-\beta t}U(t)$ , find the power spectral density of the output  $Y(t)$  corresponding to the input  $X(t)$ . (6)

## Solutions

### Part-A

1.  $Z = 3X + 4Y; E(Z) = 3E(X) + 4E(Y);$   
 $\therefore Z - E(Z) = 3\{X - E(X)\} + 4\{Y - E(Y)\}$

$$\begin{aligned}\text{Var}(Z) &= E\{Z - E(Z)\}^2 = 9 \text{ Var}(X) + 16 \text{ Var}(Y) + 24 E\{X - E(X)\} \{Y - E(Y)\} \\ &= 9 \times 2 + 16 \times 3 + 24 \times 0 \quad (\because X \text{ and } Y \text{ are independent}) \\ &= 66\end{aligned}$$

2.  $P(X > k) = \int_k^1 3x^2 dx = 0.5; [x^3]_k^1 = 0.5; \text{ i.e., } 1 - k^3 = 0.5 \text{ or } k^3 = 0.5 \text{ and } k = 0.794.$
3. Statement of CLT for iid RV's is available in the book.
4. Properties of  $F(x, y)$  is available in the book.
5. Definitions of ergodic process, mean ergodic, correlation ergodic and distribution ergodic processes will form the answer for this question.
6. As binomial process is a Markov process, viz., a discrete time Markov chain, we can consider (1) random walk with reflecting barriers, and (2) random walk with absorbing barriers as applications of binomial process.
7. Definition and properties of cross-correlation function  $R_{XY}(\tau)$  are available in the book.
8. This is a worked example in the book.
9. Definitions of a system and a linear system are available in the book.
10. Definition is available in the book.

## Part-B

11. (a) (i) Definition of MGF of a random variable is available in the book.

Gamma distribution is defined by  $f_x(x) = \frac{1}{(n)} x^{n-1} e^{-x}; x \geq 0, n > 0.$

$$\begin{aligned}M_X(t) &= \int_0^\infty e^{tx} \cdot \frac{1}{(n)} x^{n-1} e^{-x} dx = \frac{1}{(n)} \int_0^\infty x^{n-1} \cdot e^{-(1-t)x} dx \\ &= \frac{1}{(n)} \cdot \frac{1}{(1-t)^n} \int_0^\infty y^{n-1} e^{-y} dy, \text{ on putting } (1-t)x = y \\ &= \frac{1}{(1-t)^n} \text{ or } (1-t)^{-n}\end{aligned}$$

$$\begin{aligned}M_X(t) &= 1 + \underbrace{\frac{n}{1} t}_{[1]} + \underbrace{\frac{n(n+1)}{2} t^2}_{[2]} + \underbrace{\frac{n(n+1)(n+2)}{3} t^3}_{[3]} \\ &\quad + \underbrace{\frac{n(n+1)(n+2)(n+3)}{4} t^4}_{[4]} + \dots \infty\end{aligned}$$

$$\begin{aligned}\mu'_1 &= E(X) = n; \mu'_2 = E(X^2) = n(n+1); \text{Var}(X) = \mu'_2 - \mu'^2_1 = n; \\ \mu'_3 &= n(n+1)(n+2); \mu'_4 = n(n+1)(n+2)(n+3)\end{aligned}$$

(ii) Binomial distribution is defined by  $p_r = P(X = r) = nC_r p^r q^{n-r}$

$$M_X(t) = \sum_{r=0}^n e^{tr} nC_r p^r q^{n-r} = \sum_{r=0}^n nC_r (pe^t)^r q^{n-r}$$

$$= (pe^t + q)^{n(r=0,1,2,\dots,n)}$$

Derivation of the first four moments.

$$\begin{aligned} M_X(t) &= \left[ p \left( 1 + \frac{t}{1} + \frac{t^2}{2} + \frac{t^3}{3} + \frac{t^4}{4} + \dots \right) + q \right]^n \\ &= \left\{ 1 + pt \left( 1 + \frac{t}{2} + \frac{t^2}{6} + \frac{t^3}{24} + \dots \right) \right\}^n \\ &= 1 + \frac{n}{1} pt \left( 1 + \frac{t}{2} + \dots \right) + \frac{n(n-1)}{2} p^2 t^2 \left( 1 + \frac{t}{2} + \frac{t^2}{6} \right)^2 \\ &\quad + \frac{n(n-1)(n-2)}{3} p^3 t^3 \left( 1 + \frac{t}{2} + \frac{t^2}{6} + \frac{t^3}{24} \right)^3 \\ &\quad + \frac{n(n-1)(n-2)(n-3)}{4} p^4 t^4 \left( 1 + \frac{t}{2} + \frac{t^2}{6} + \frac{t^3}{24} \right)^4 + \dots \end{aligned}$$

$$(1) \mu'_1 = \text{coefficient of } \frac{t}{1} = np;$$

$$\mu'_2 = \text{coefficient of } \frac{t^2}{2} = np + n(n-1)p^2;$$

$$\mu'_3 = \text{coefficient of } \frac{t^3}{3} = np + 3n(n-1)p^2 + n(n-1)(n-2)p^3$$

$$\mu'_4 = \text{coefficient of } \frac{t^4}{4} = np + 7n(n-1)p^2 + 4n(n-1)(n-2)p^3 + n(n-1)(n-2)(n-3)$$

(2) Recurrence formula for the central moments is given in the book.

11. (b) (i) This is a worked example in the book.

$$(ii) f(x) = \frac{1}{4} xe^{-x/2}; x > 0$$

$$M_X(t) = \int_0^\infty e^{tx} f(x) dx = \frac{1}{4} \int_0^\infty x e^{-\left(\frac{1}{2}-t\right)x} dx$$

$$= \frac{1}{4} \left[ x \left\{ \frac{e^{-\left(\frac{1}{2}-t\right)x}}{-\left(\frac{1}{2}-t\right)} \right\} - \frac{e^{-\left(\frac{1}{2}-t\right)x}}{\left(\frac{1}{2}-t\right)^2} \right]_0^{\infty} = \frac{1}{4} \cdot \frac{1}{\left(\frac{1}{2}-t\right)^2} = (1-2t)^{-2}$$

$$M_X(t) = 1 + 2(2t) + 3(2t)^2 + 4(2t)^3 + 5(2t)^4 + \dots \infty \quad (1)$$

$\mu'_1$  = coefficient of  $\frac{t}{1}$  in (1) = 4;  $\mu'_2$  = 24;  $\mu'_3$  = 192;  $\mu'_4$  = 1920.

12. (a) (i)

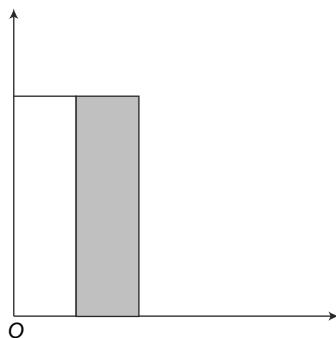


Fig. 1

$$\begin{aligned} P\left(X > \frac{1}{2}\right) &= \int_0^2 \int_{1/2}^1 \left(x^2 + \frac{xy}{3}\right) dx dy = \int_0^2 \left(\frac{x^3}{3} + y \frac{x^2}{6}\right)_{1/2}^1 dy \\ &= \int_0^2 \left(\frac{7}{24} + \frac{y}{8}\right) dy = \left(\frac{7}{24}y + \frac{y^2}{16}\right)_0^2 = \frac{5}{6} \end{aligned}$$

(ii)

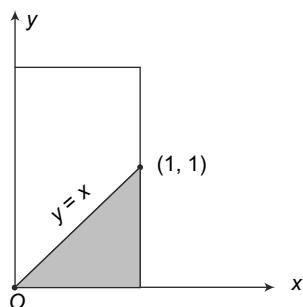
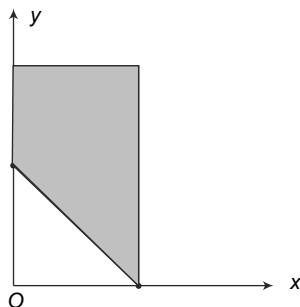


Fig. 2

$$\begin{aligned}
 P(Y < X) &= \int_0^1 \int_y^1 \left( x^2 + \frac{xy}{3} \right) dx dy = \int_0^1 \left( \frac{x^3}{3} + \frac{yx^2}{6} \right)_y^1 dy \\
 &= \int_0^1 \left( \frac{1}{3} + \frac{y}{6} - \frac{y^3}{3} - \frac{y^3}{6} \right) dy = \left( \frac{y}{3} + \frac{y^2}{12} - \frac{y^4}{8} \right)_0^1 = \frac{7}{24}
 \end{aligned}$$

(iii)

**Fig. 3**

$$\begin{aligned}
 P(X + Y \geq 1) &= \int_0^1 \int_{1-x}^2 \left( x^2 + \frac{xy}{3} \right) dy dx = \int_0^1 \left( x^2 y + \frac{xy^2}{6} \right)_{1-x}^2 dx \\
 &= \int_0^1 \left\{ 2x^2 + \frac{2}{3}x - x^2(1-x) - \frac{x}{6}(1-x)^2 \right\} dx \\
 &= \int_0^1 \left( \frac{1}{2}x + \frac{4}{3}x^2 + \frac{5}{6}x^3 \right) dx \\
 &= \frac{1}{4} + \frac{4}{9} + \frac{5}{24} = \frac{65}{72}
 \end{aligned}$$

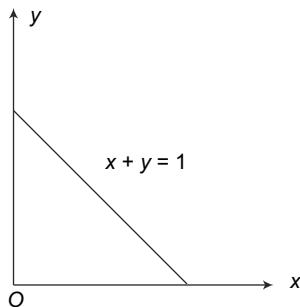
$$(iv) f_X(x) = \int_0^2 \left( x^2 + \frac{xy}{3} \right) dy = 2x^2 + \frac{2}{3}x; 0 < x < 1$$

$$\therefore f_{Y/X}(y) = \frac{f(x, y)}{f_X(x)} = \frac{x^2 + \frac{xy}{3}}{2x^2 + \frac{2}{3}x}; 0 < y < 2$$

$$f_Y(y) = \int_0^1 \left( x^2 + \frac{xy}{3} \right) dx = \frac{1}{3} + \frac{y}{6}; 0 < y < 2$$

$$\therefore f_{X/Y}(x) = \frac{f(x, y)}{f_Y(y)} = \frac{x^2 + \frac{xy}{3}}{\frac{1}{3} + \frac{t}{6}}; 0 < x < 1$$

12. (b) (i)



**Fig. 4**

$$\begin{aligned} E(X) &= \int_0^{1-y} \int_0^y 3x(x+y) dx dy = \int_0^1 \left( x^3 + \frac{3}{2}yx^2 \right)_{0}^{1-y} dy \\ &= \int_0^1 \left[ (1-y)^3 + \frac{3}{2}y(1-y)^2 \right] dy = \int_0^1 \left( 1 - \frac{3}{2}y + \frac{y^3}{2} \right) dy = \frac{3}{8} \end{aligned}$$

$$\text{By symmetry, } E(Y) = \frac{3}{8}$$

$$\begin{aligned} E(XY) &= \int_0^{1-y} \int_0^y 3xy(x+y) dx dy = \int_0^1 \left( yx^3 + \frac{3}{2}y^2x^2 \right)_{0}^{1-y} dy \\ &= \int_0^1 \left[ y(1-y)^3 + \frac{3}{2}y^2(1-y)^2 \right] dy = \int_0^1 \left( y - \frac{3}{2}y^2 + \frac{y^4}{2} \right) dy = \frac{1}{10} \end{aligned}$$

$$\text{Cov}(X, Y) = E(XY) - E(X) \cdot E(Y) = \frac{1}{10} - \frac{9}{64} = -\frac{13}{320}$$

(ii)  $\Sigma x = 440; \Sigma y = 385; \Sigma x^2 = 20,674; \Sigma y^2 = 17,253; \Sigma xy = 18,499$

Let  $y = ax + b$  be the regression line of  $Y$  on  $X$ .

The normal equations are  $440a + 10b = 385$  and  $20,674a + 440b = 18,499$ .

Solving these equations,  $a = 1.186$  and  $b = -13.684$

$\therefore$  regression line of  $Y$  on  $X$  is  $y = 1.186x - 13.684$

$$\therefore [y]_{x=55} = 51.546$$

Let  $x = cy + d$  be the regression line of  $X$  on  $Y$ .

The normal equations are  $385c + 10d = 440$  and  $17253c + 385d = 18499$ .

Solving these equations,  $c = 0.520$  and  $d = 23.990$

$\therefore$  regression line of  $X$  on  $Y$  is  $x = 0.520y + 23.990$ .

13. (a) (i) This is a worked example in the book.

Since  $E\{X(t)\} = 1$ , a constant,  $\{X(t)\}$  is first order stationary.

(ii) This is a worked example in the book.

13. (b) (i) This is a standard book-work available in the book.  
 Since  $E\{X(t)\} = \lambda t$ , a function of  $t$ ,  $\{X(t)\}$  is not first order stationary.  
 (ii) This is a worked example in the book.
14. (a) (i) This is a book-work, available in the book.  
 (ii) This is a worked example in the book.
14. (b) Statement and proof of Weiner–Kinchine theorem available in the book.
15. (a) (i) A standard property, derivation of which is available in the book.  
 To find  $R_{XX}(\tau)$  without the definition of  $\{X(t)\}$  is not possible.  
 (ii) This is a worked example in the book.
15. (b) (i) This is a worked example in the book.

$$\begin{aligned}(ii) \quad H(\omega) &= F\{h(t)\} = \int_{-\infty}^{\infty} h(t) e^{-i\omega t} dt = \int_0^{\infty} e^{-\beta t} \cdot e^{-i\omega t} dt \\ &= \left\{ \frac{e^{-(\beta+i\omega)t}}{-(\beta+i\omega)} \right\}_0^{\infty} = \frac{1}{\beta+i\omega}\end{aligned}$$

$$S_{YY} = |H(\omega)|^2 \quad S_{XX}(\omega) = \frac{1}{\beta^2 + \omega^2} S_{XX}(\omega), \text{ where } S_{XX}(\omega) \text{ and } S_{YY}(\omega) \text{ are the power spectral densities of } \{X(t)\} \text{ and } \{Y(t)\}.$$

