

Inner product

Definition

The inner product (or scalar product) of two vectors, \mathbf{x} and \mathbf{y} having the same dimensionality, d , will be denoted as $\mathbf{x}^t \mathbf{y}$ and the result is a scalar,

$$\mathbf{x}^t \mathbf{y} = \sum_{i=1}^d \mathbf{x}_i \mathbf{y}_i = \mathbf{y}^t \mathbf{x}$$

Definition

The Euclidean norm or length of a vector \mathbf{x} is,

$$\|\mathbf{x}\| = \sqrt{\mathbf{x}^t \mathbf{x}}$$

The vector is normalized if $\|\mathbf{x}\| = 1$.

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Definition

The angle, θ , between two d -dimensional vectors \mathbf{x} and \mathbf{y} is given by,

$$\cos\theta = \frac{\mathbf{x}^t\mathbf{y}}{\|\mathbf{x}\|\|\mathbf{y}\|}$$

Inner product is a measure of the colinearity of two vectors; an indication of similarity (to within a scale factor).

Inner product

Example (Inner product)

The inner product of two vectors $\mathbf{x} = \begin{bmatrix} 1 \\ 3 \\ 4 \\ 6 \\ 8 \end{bmatrix}$ and $\mathbf{y} = \begin{bmatrix} 5 \\ 1 \\ 0 \\ 2 \\ 7 \end{bmatrix}$ is

$$1 \times 5 + 3 \times 1 + 4 \times 0 + 6 \times 2 + 8 \times 7 = 76$$

Inner product

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Example (Magnitude)

The magnitude of vector \mathbf{x} is

$$\|\mathbf{x}\| = \sqrt{1 \times 1 + 3 \times 3 + 4 \times 4 + 6 \times 6 + 8 \times 8} = \sqrt{126} = 11.23.$$

Inner product

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Example (Magnitude)

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$$\|\mathbf{x}\| = \sqrt{1 \times 1 + 3 \times 3 + 4 \times 4 + 6 \times 6 + 8 \times 8} = \sqrt{126} = 11.23.$$

Example (Magnitude)

The magnitude of vector \mathbf{y} is

$$\|\mathbf{y}\| = \sqrt{5 \times 5 + 1 \times 1 + 0 + 2 \times 2 + 7 \times 7} = \sqrt{79} = 8.89.$$

Inner product

Example (Inner product)

The inner product of two vectors $\mathbf{x} = \begin{bmatrix} 1 \\ 3 \\ 4 \\ 6 \\ 8 \end{bmatrix}$ and $\mathbf{y} = \begin{bmatrix} 5 \\ 1 \\ 0 \\ 2 \\ 7 \end{bmatrix}$ is

$$1 \times 5 + 3 \times 1 + 4 \times 0 + 6 \times 2 + 8 \times 7 = 76$$

Example (Angle)

The angle between the vectors \mathbf{x} and \mathbf{y} is

$$\theta = \arccos \frac{\mathbf{x}^t \mathbf{y}}{\|\mathbf{x}\| \|\mathbf{y}\|} = \arccos \frac{76}{11.23 \times 8.89} = 0.707 \text{ radians}$$

Orthogonal vectors

Definition

If $\mathbf{x}^t \mathbf{y} = 0$, the vectors are orthogonal.

Definition

If $\mathbf{x}^t \mathbf{y} = \|\mathbf{x}\| \|\mathbf{y}\|$, the vectors are colinear.

Definition

The Cauchy-Schwartz inequality follows directly from previous definition of an angle between two vectors. In other words,

$$\mathbf{x}^t \mathbf{y} \leq \|\mathbf{x}\| \|\mathbf{y}\|.$$

Matrices

Definition

The trace of a square matrix \mathbf{A} , denoted $Tr\{\mathbf{A}\}$, is the sum of its diagonal elements,

$$Tr\{A\} = \sum_{i=1}^d a_{ii}$$

and $Tr\{\mathbf{CD}\} = Tr\{\mathbf{DC}\}$ provided that the product \mathbf{CD} is a square matrix. Neither \mathbf{C} nor \mathbf{D} need be square.

Example

The trace of the matrix $\mathbf{D} = \begin{bmatrix} 4 & 2 & 9 & 0 \\ 1 & 6 & 7 & 1 \\ 4 & 9 & 2 & 3 \\ 2 & 7 & 3 & 1 \end{bmatrix}$ is $4 + 6 + 2 + 1 = 13$.

Definition

The determinant of a square ($d \times d$) matrix \mathbf{M} , written as $|\mathbf{M}|$ is the sum,

$$|\mathbf{M}| = \sum_{j=1}^d m_{ij} M_{ij} \text{ for } i = 1 \dots, d$$

where the cofactor, M_{ij} , is the determinant of the matrix formed by deleting the i th row and the j th column of \mathbf{M} , multiplied by $(-1)^{i+j}$.

Definition

The transpose of the matrix of the cofactors, $\mathbf{C}(c_{ij} = M_{ij})$, is called the adjoint of \mathbf{M} , $\text{Adj}[\mathbf{M}]$.

Matrices

Definition

The inverse of a $d \times d$ matrix, \mathbf{M} is that unique matrix \mathbf{M}^{-1} of dimension $d \times d$, with entries such that,

$$\mathbf{M}^{-1}\mathbf{M} = \mathbf{M}\mathbf{M}^{-1} = \mathbf{I}$$

Definition

We can obtain the inverse of a matrix from,

$$\mathbf{M}^{-1} = \frac{Adj[\mathbf{M}]}{|\mathbf{M}|}$$

Definition

If the inverse exists the matrix is said to be nonsingular otherwise it is singular and $|\mathbf{M}| = 0$

Note that, $(\mathbf{M}^t)^{-1} = (\mathbf{M}^{-1})^t$ and $(\mathbf{AB})^{-1} = \mathbf{B}^{-1}\mathbf{A}^{-1}$.

Matrices

Example

Consider the matrix $\mathbf{A} = \begin{bmatrix} 2 & 4 & 6 \\ -1 & 2 & 3 \\ 1 & 4 & 9 \end{bmatrix}$. The determinant is,

$$2 \times (2 \times 9 - 3 \times 4) - 4 \times (-1 \times 9 - 1 \times 3) + 6 \times (-1 \times 4 - 2 \times 1) = 24$$

Example

The adjoint of \mathbf{A} is $\begin{bmatrix} 6 & 12 & -6 \\ -12 & 12 & -4 \\ 0 & -12 & 8 \end{bmatrix}^t = \begin{bmatrix} 6 & -12 & 0 \\ 12 & 12 & -12 \\ -6 & -4 & 8 \end{bmatrix}$.

Example

The inverse of \mathbf{A} is

$$\mathbf{A}^{-1} = \frac{1}{24} \begin{bmatrix} 6 & -12 & 0 \\ 12 & 12 & -12 \\ -6 & -4 & 8 \end{bmatrix}.$$

Linear independence and rank

Definition

A set of k vectors of equal dimension are linearly dependent if there exists a set of scalars c_1, c_2, \dots, c_k , not all zero, such that,

$$c_1 \mathbf{x}_1 + \dots + c_k \mathbf{x}_k = 0$$

If it is impossible to find such a set of c_1, c_2, \dots, c_k , then the vectors, $\mathbf{x}_1, \dots, \mathbf{x}_k$ are said to be linearly independent.

Definition

The rank of a matrix is the maximum number of linearly independent rows (or equivalently, the maximum number of linearly independent columns).

Definition

A $d \times d$ matrix is of full rank if the rank is equal to d . It will also be true that the determinant is non-zero and it will possess an inverse.

Linear independence and rank

Definition

For a rectangular matrix (i.e. non-square) \mathbf{M} of dimension $d \times n$, the rank of \mathbf{M} , denoted $\text{rank}(\mathbf{M})$ is such that $\text{rank}(\mathbf{M}) \leq \min(d, n)$.

We have that,

$$\text{rank}(\mathbf{M}) = \text{rank}(\mathbf{M}^t) = \text{rank}(\mathbf{M}^t \mathbf{M}) = \text{rank}(\mathbf{M} \mathbf{M}^t)$$

Orthogonal matrix

Definition

A square matrix \mathbf{M} , is orthogonal if,

$$\mathbf{M}\mathbf{M}^t = \mathbf{M}^t\mathbf{M} = \mathbf{I}$$

- The rows and columns are orthonormal, $\mathbf{x}^t\mathbf{y} = 0$ and $\mathbf{x}^t\mathbf{x} = 1$, $\mathbf{y}^t\mathbf{y} = 1$ for any two different columns \mathbf{x} and \mathbf{y} .
- An orthogonal matrix represents a linear transformation that preserves distances and angles, consisting of a rotation and/or reflection
- An orthogonal matrix is nonsingular and the inverse is its transpose, $\mathbf{M}^{-1} = \mathbf{M}^t$
- The determinant of an orthogonal matrix is ± 1 , with -1 indicating a reflection and $+1$ indicating pure rotation.

Positive definiteness

Definition

A square matrix \mathbf{M} is positive definite if the quadratic form, $\mathbf{x}^t \mathbf{M} \mathbf{x} > 0$ for all vectors $\mathbf{x} \neq \mathbf{0}$.

Definition

A square matrix \mathbf{M} is positive semidefinite if the quadratic form, $\mathbf{x}^t \mathbf{M} \mathbf{x} \geq 0$ for all vectors $\mathbf{x} \neq \mathbf{0}$. A positive definite matrix will have a full rank.

Eigenvalue problem

- Given a $d \times d$ matrix \mathbf{M} , an important class of linear equations is of the form,

$$\mathbf{M}\mathbf{x} = \lambda\mathbf{x}$$

or

$$(\mathbf{M} - \lambda\mathbf{I})\mathbf{x} = \mathbf{0}$$

for a scalar λ .

- The solution to the characteristic equation,

$$|\mathbf{M} - \lambda\mathbf{I}| = 0$$

gives the eigenvalues or characteristic roots of the $d \times d$ matrix.

Eigenvalue problem

- The characteristic equation is a d th order polynomial in λ . There are d solutions, $\lambda_1, \lambda_2, \dots, \lambda_d$. They are not necessarily distinct and may be real or complex.
- Associated with each eigenvalue, λ_i is an eigenvector, \mathbf{u}_i , such that,

$$\mathbf{M}\mathbf{u}_i = \lambda_i\mathbf{u}_i$$

- The eigenvectors are not unique because any scalar multiple of \mathbf{u}_i satisfies $\mathbf{M}\mathbf{u}_i = \lambda_i\mathbf{u}_i$. Eigenvectors are usually normalized so that, $\mathbf{u}_i^t\mathbf{u}_i = 1$.

Eigenvalue problem

Properties of eigenvalues and eigenvectors

- 1 The product of the eigenvalues is equal to the determinant of \mathbf{M} .
 $\prod_{i=1}^d \lambda_i = \det(\mathbf{M})$. For a given matrix, if the eigenvalues are all non-zero, then the inverse of \mathbf{M} exists.
- 2 The sum of the eigenvalues is equal to the trace of the matrix. $\sum_{i=1}^d \lambda_i = \text{Tr}(\mathbf{M})$.
- 3 If \mathbf{M} is a real symmetric matrix, the eigenvalues and eigenvectors are all real.
- 4 If \mathbf{M} is positive definite, the eigenvalues are all greater than zero.
- 5 If \mathbf{M} is positive semidefinite of rank m , then there will be m non-zero eigenvalues and $d - m$ eigenvalues with the value of zero.

Eigenvalue problem

Properties of eigenvalues and eigenvectors (continued)

- 6 Every real symmetric matrix has a set of orthonormal characteristic vectors. The matrix, \mathbf{U} , whose columns are the eigenvectors of the real symmetric matrix is orthogonal. $\mathbf{U} = \{\mathbf{u}_1, \dots, \mathbf{u}_d\}$. We have $\mathbf{U}^t \mathbf{U} = \mathbf{U} \mathbf{U}^t = \mathbf{I}$.
- 7 The matrix \mathbf{U} diagonalizes \mathbf{M} ,

$$\mathbf{U}^t \mathbf{M} \mathbf{U} = \Lambda$$

where $\Lambda = \text{diag}(\lambda_1, \dots, \lambda_d)$ is a diagonal matrix whose entries are the eigenvalues of \mathbf{M} .

$$\mathbf{M} = \mathbf{U} \Lambda \mathbf{U}^t = \sum_{i=1}^d \lambda_i \mathbf{u}_i \mathbf{u}_i^t$$

Eigenvalue problem

Properties of eigenvalues and eigenvectors (continued)

- 8 If \mathbf{M} is positive definite, then $\mathbf{M}^{-1} = \mathbf{U}\mathbf{\Lambda}^{-1}\mathbf{U}^t$. Here, $\mathbf{\Lambda}^{-1} = \text{diag}(1/\lambda_1, \dots, 1/\lambda_d)$

There are more linear algebraic results and they will be introduced as we need them!

Machine Learning: Algorithms and Applications

Advanced Multimedia Research Lab
University of Wollongong

Math Review - Elementary Probability Theory

Outline of Topics

- 1 Context
- 2 Events and probabilities
- 3 Conditional probability
- 4 Independent and mutually exclusive events
- 5 Joint and marginal probability distribution
- 6 Random variables
- 7 Normal distribution
- 8 Summary of key points

Pattern recognition and probability

- Machine learning algorithms help us to analyse data and extract or recognize inherent patterns.
- Real life data displays uncertainty and variability.
- Probability theory allows us to model the uncertainty and variability in data.
- The models are powerful and provide the basis for statistical/probabilistic reasoning.

Introduction to Probability

- Probability is the study of **non-deterministic** events.
- We represent the possible outcomes of an experiment as a set S (also called *sample space*).
- For example if,

$$S = \{f_1, f_2, f_3, \dots, f_6\},$$

is the set of outcomes of rolling a die. The probability of each outcome is,

$$p(f_i) = 1/6, \quad i = 1, 2, \dots, 6$$

- **Events** are a subsets of sample space. For example the sets:

$$\{even\} = \{f_2, f_4, f_6\} \text{ and } \{odd\} = \{f_1, f_3, f_5\}$$

are events in the rolling of the die.

Introduction to Probability

- Two events, A and B , are **mutually exclusive** if $A \cap B = \emptyset$.
- Suppose we have a **finite sample space** and a set of events. A **probability function**, P gives a real value for each event such that:
 - for every event A we have $0 \leq P(A) \leq 1$.
 - $P(S) = 1$.
 - If A and B are mutually exclusive, then $P(A \cup B) = P(A) + P(B)$.
- If each sample point in the space has the same probability, we refer to this as **equiprobable** or **uniform space**.
- We can use **the number of points in A** to find its probability.

Introduction to Probability

Example 1:

Suppose we have a loaded die with

$$p(f_i) = 1/4, \quad \text{for } i = 1, 2, 3$$

$$p(f_i) = 1/12, \quad \text{for } i = 4, 5, 6$$

Then,

$$p(\text{even}) = \frac{1}{4} + 2 \times \frac{1}{12} = \frac{5}{12}$$

It can be shown that for two events A and B

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

Use Venn diagram to show this to yourself!

Introduction to Probability

Example 2:

Consider the events:

A= even; B= multiple of 3
such that,

$$P(A) = \frac{1}{2}, \quad P(B) = \frac{1}{3}$$

$$P(A \cap B) = P(\text{even and multiple of 3}) = \frac{1}{6}$$

and

$$P(A \cup B) = \frac{1}{3} + \frac{1}{2} - \frac{1}{6} = \frac{2}{3}$$

Introduction to Probability

Example 3:

Consider the toss of two dice.

Each point of the sample space can be represented by a two tuple,

$$(i, j), \quad i = 1, 2, \dots, 6, \quad j = 1, 2, \dots, 6$$

For an un-biased die every outcome has equal probability, $p(i, j) = \frac{1}{36}$ of occurrence.

For illustrative purposes, let us define some events. (See next slide)

Introduction to Probability

We define 12 different events v_1, \dots, v_{12} where v_u is the event that $i + j = u$. So, we have:

$$v_1 \Rightarrow p(v_1) = 0$$

$$v_2 (1, 1) \Rightarrow p(v_2) = 1/36$$

$$v_3 (1, 2), (2, 1) \Rightarrow p(v_3) = 1/18$$

$$v_4 (1, 3), (2, 2) (3, 1) \Rightarrow p(v_4) = 1/12$$

$$v_5 (1, 4), (2, 3) (3, 2) (4, 1) \Rightarrow p(v_5) = 1/9$$

$$v_6 (1, 5), (2, 4) (3, 3) (4, 2) (5, 1) \Rightarrow p(v_6) = 5/36$$

$$v_7 (1, 6), (2, 5) (3, 4) (4, 3) (5, 2) (6, 1) \Rightarrow p(v_7) = 1/6$$

$$v_8 (2, 6) (3, 5) (4, 4) (5, 3) (6, 2) \Rightarrow p(v_8) = 5/36$$

$$v_9 (3, 6) (4, 5) (5, 4) (6, 3) \Rightarrow p(v_9) = 1/9$$

$$v_{10} (4, 6) (5, 5) (6, 4) \Rightarrow p(v_{10}) = 1/12$$

$$v_{11} (5, 6) (6, 5) \Rightarrow p(v_{11}) = 1/18$$

$$v_{12} (6, 6) \Rightarrow p(v_{12}) = 1/36$$

Conditional Probability

- Events may be dependent; that is, occurrence of one influences the probability of the other.
- The probability of event A given that event E has occurred is given by

$$P(A|E) = \frac{P(A \cap E)}{P(E)}$$

Example 4:

Consider two events described as follows:

$$\begin{aligned} A &= \{ \text{points with } i + j = 2u \} \\ &= \{v_2, v_4, v_6, v_8, v_{10}, v_{12}\} \\ E &= \{ \text{points with } i = j \} \\ &= \{(1, 1), (2, 2), (3, 3), (4, 4), (5, 5), (6, 6)\} \end{aligned}$$

Conditional Probability

Example 4 cont.

Then,

- $P(A) = 18/36 = 1/2$
- $P(E) = 6/36 = 1/6$
- E is a subset of A:
 - $a \in E \Rightarrow a \in A.$
- We can then infer that $P(A \cap E) = P(E) = 6/36 = 1/6$
So $P(E|A) = P(A \cap E)/P(A) = (1/6)/(1/2) = 1/3$

Independent Events

- Two events E and A are **independent** if

$$P(A|E) = P(A)$$

That is, occurrence of E does not influence probability of occurrence of A .

- Equivalently we can write,

$$P(A \cap E) = P(A|E)P(E) = P(A)P(E)$$

Example 5:

A fair coin is tossed three times. Consider two events:

$A = \{\text{first toss is heads}\}$

$B = \{\text{second toss is heads}\}$

We have

$$P(A) = P(\{hhh, hht, hth, htt\}) = 1/2$$

$$P(B) = P(\{hhh, hht, thh, tht\}) = 1/2$$

$$P(A \cap B) = P(\{hht, hhh\}) = 1/4$$

So $P(A \cap B) = P(A) \times P(B)$ since the two events are independent.

Mutually Exclusive Events

- If events E and A are **mutually exclusive**, occurrence of one excludes the occurrence of the other and so,

$$P(A|E) = 0$$

- Note that this means that mutually exclusive events are *not* independent. Recall that independence requires that $P(A \cap E) = P(A|E)P(E) = P(A)P(E)$.
- If $A \subset E$ then

$$P(A \text{ and } E) = P(A)$$

Mutually Exclusive Events

- S consists of the outcomes of rolling a die.

$$P(f_i) = \frac{1}{6} \quad i = 1, 2, \dots, 6$$

- Let $E = \text{even}$ and $A = \{f_2\}$.

Then

$$P(f_2|E) = \frac{P(f_2 \text{ and } E)}{P(E)} = \frac{1/6}{1/2} = \frac{1}{3}$$

- Suppose we are interested in the event $A = \{ \text{even and divisible by 3} \}$.

Then

$$A \subset E$$

and

$$P(A \text{ and } E) = P(A) = 1/6.$$

Joint Probability Distribution

We may have **joint probability** distribution: In this case the sample space consists of pairs of outcomes of two experiments. We illustrate by way of example.

Example 6:

Joint probability of (temperature, humidity) in summer:

	Low	High
50°F	1/6	1/12
60°F	5/18	1/9
70°F	7/36	2/12

Marginal Distribution

- Given a joint distribution for (X, Y) , one can derive distribution of X and Y .
- Let $P(x, y)$ denote probability of the sample point (x, y) .

$$P(x) = \sum_{\text{all values of } y} P(x, y)$$

$$P(y) = \sum_{\text{all values of } x} P(x, y)$$

- $P(x)$ and $P(y)$ are *marginal distributions*.

Example 7:

In the temperature and humidity example (6) we have marginal distributions as follows:

$$P(\text{Temperature}) = \left(\frac{9}{36}, \frac{14}{36}, \frac{13}{36} \right)$$

$$P(\text{Humidity}) = \left(\frac{23}{36}, \frac{13}{36} \right)$$

Conditional independence

- Recall the joint probability of two events A and B ,
 $P(A, B) = P(A)P(B|A) = P(B)P(A|B)$
- If events A and B are independent we write:

$$P(A, B) = P(A)P(B)$$

- If we have three events A, B, C , the joint probability is written as:

$$P(A, B, C) = P(A)P(B|A)P(C|A, B)$$

- Of course if the three events are independent we write:

$$P(A, B, C) = P(A)P(B)P(C)$$

- Can we have conditional independence? What does this mean?
How useful is it?

Conditional independence (adapted from QMC London slides)

Example of independence

Suppose Jane and John each toss separate coins. Let A represent the variable “Jane’s toss outcome”, and B represent the variable “John’s toss outcome”. Both A and B have two possible values (Heads and Tails). Without any controversy we can assume that A and B are independent. Evidence about B will not change our belief in A .

Conditional independence (adapted from QMC London slides)

Example of conditional independence

Now suppose both Jane and John toss the same coin. Again let A represent the variable “Jane’s toss outcome”, and B represent the variable “John’s toss outcome”. Assume also that there is a possibility that the coin is biased towards heads but we do not know this for certain. In this case A and B are not independent. For example, observing that A is Heads causes us to increase our belief in B being Heads (i.e. $P(b|a) > P(a)$ in the case when $a = \text{Heads}$ and $b = \text{Heads}$).

Conditional independence (adapted from QMC London slides)

In the conditional independence example the variables A and B are both dependent on a separate variable C , “the coin is biased towards Heads” (which has the value True or False). Although A and B are not independent, it turns out that once we know for certain the value of C then any evidence about A cannot change our belief about B . Specifically:

$$P(B|C) = P(B|A, C)$$

Conditional independence (adapted from QMC London slides)

Another example

Suppose that Norman and Martin live on opposite sides of the City of London and come to work by completely different means, say Norman comes by train while Martin drives. Let A represent the variable “Norman late” (which has value true or false) and similarly let B represent the variable “Martin late”. It would be tempting in these circumstances to assume that A and B must be independent. However, even if Norman and Martin lived and worked in different countries there may be factors (such as an international fuel shortage) which could mean that A and B are not independent.

In practice any model of uncertainty should take account of all reasonable factors.

Thus while, say, a meteorite hitting the Earth might be reasonably excluded it does not seem reasonable to exclude the fact that both A and B may be affected by a Train strike (C). Clearly $P(A)$ will increase if C is true; but $P(B)$ will also increase because of extra traffic on the roads.

Random Variable

- A **random variable** assigns **values** to events in a probability space.

Example 8:

X is a random variable that assigns 0 or 1 to the outcome of the experiment in example 3.

If the sum is divisible by 4, then $X = 0$ otherwise it is 1.

We have

$$\begin{aligned} p(X = 0) &= \sum p(\text{sum divisible by 4}) \\ &= 1/12 + 5/36 + 1/36 = 9/36 = 1/4 \\ p(X = 1) &= 1 - 1/4 = 3/4 \end{aligned}$$

Example 8 contd.:

- Let a random variable X take values x_1, x_2, \dots, x_N .
- **Expected value** of X is

$$E(X) = \sum x_n p(x_n)$$

Expected value can be regarded as the average value of the random variable.

In example 8

$$E(X) = 0 \times 1/4 + 1 \times 3/4 = 3/4$$

Random Variables

Example 9: We define a random variable (RV) \mathbf{X} as a real-valued function of the sample space \mathbf{S} , with a range denoted by \mathbf{R}_X . Further, we say that a RV is discrete if the range \mathbf{R}_X is a discrete set. Now, the probability function for a discrete RV \mathbf{X} is,

$$p_X(x) = P(\mathbf{X} = x), \quad \text{for } x \in \mathbf{R}_X$$

A class in information theory contains 10 students, three of whom are 19, 4 are 20, one is 21, one is 24 and one is 26. Two students are selected at random without replacement from the class.

Let \mathbf{X} be the random variable denoting the average age of the 2 selected students; derive the probability function for \mathbf{X} .

Solution

The random variable denoting the average age of two students selected randomly without replacement is **X**. Total number of ways of selecting 2 from 10 is, $\binom{10}{2}$. That is,

$$\frac{10!}{(10-2)! \times 2!} = 45$$

The number of “average of 2 ages” that can be made from $\{19, 20, 21, 24, 26\}$ is $\binom{5}{2} = 10$. The “average of 2 ages” are: $\{19, 19.5, 20, 20.5, 21.5, 22, 22.5, 23, 23.5, 25\}$.

Random Variables

Solution (contd.)

There are 3 out of 45 ways of making an average age of 19 and $P(\mathbf{X} = 19) = \frac{3}{45}$. There are $3 \times 4 = 12$ out of 45 ways of making an average age of 19.5 and $P(\mathbf{X} = 19.5) = \frac{12}{45}$. There are $6 + 3 = 9$ out of 45 ways of making an average age of 20 and $P(\mathbf{X} = 20) = \frac{9}{45}$. We can compute the rest of the probabilities using similar reasoning.

The required probability function for \mathbf{X} is as follows:

\mathbf{X}	19	19.5	20	20.5	21.5	22	22.5	23	23.5	25
$P(\mathbf{X})$	$\frac{3}{45}$	$\frac{12}{45}$	$\frac{9}{45}$	$\frac{4}{45}$	$\frac{3}{45}$	$\frac{4}{45}$	$\frac{4}{45}$	$\frac{4}{45}$	$\frac{1}{45}$	$\frac{1}{45}$

Second moment and Variance

- The **second moment** of a random variable X that takes on values x_i is

$$E[X^2] = \sum_i x_i^2 P(x_i) \quad (1)$$

- If we denote the expected value of X as $\mu = E[X]$, the **variance** of X is,

$$\text{Var}\{X\} = \sigma^2 = E[(X - \mu)^2] = \sum_i (x_i - \mu)^2 P(x_i) \quad (2)$$

σ is the standard deviation of the X .

Functions of two random variables

- If we have a function of two random variables $f(x, y)$ that are jointly distributed as $P(x, y)$. The expected value of the function is

$$E[f(x, y)] = \sum_{R_X} \sum_{R_Y} f(x, y)P(x, y)$$

R_X and R_Y are the ranges of X and Y respectively.

- The expectation operator is a linear function since,

$$E[\alpha_1 f_1(x, y) + \alpha_2 f_2(x, y)] = \alpha_1 E[f_1(x, y)] + \alpha_2 E[f_2(x, y)]$$

for any two functions $f_1(x, y)$ and $f_2(x, y)$; and scalars α_1 and α_2 .

Functions of two random variables

- The expected value of X and Y can be written as,

$$\mu_x = E[X] = \sum_{x \in R_X} \sum_{y \in R_Y} xP(x, y)$$

$$\mu_y = E[Y] = \sum_{x \in R_X} \sum_{y \in R_Y} yP(x, y)$$

The variances are,

$$\sigma_x^2 = E[(X - \mu_x)^2] = \sum_{x \in R_X} \sum_{y \in R_Y} (x - \mu_x)^2 P(x, y)$$

and

$$\sigma_y^2 = E[(Y - \mu_y)^2] = \sum_{x \in R_X} \sum_{y \in R_Y} (y - \mu_y)^2 P(x, y)$$

Functions of two random variables

- The **covariance** of X and Y is defined as,

$$\sigma_{xy} = E[(X - \mu_x)(Y - \mu_y)] = \sum_{x \in R_X} \sum_{y \in R_Y} (x - \mu_x)(y - \mu_y)P(x, y)$$

- Covariance is a measure of the degree of statistical dependence between X and Y . If X and Y are statistically independent $\sigma_{xy} = 0$.
- If $\sigma_{xy} = 0$, then the random variables are said to be **uncorrelated**.
- It does not follow that uncorrelated variables must be statistically independent.

Functions of two random variables

- Note that uncorrelated variables are statistically independent if they have multivariate normal distribution.
- If we have the relation, $Y = \alpha X$, which denotes that Y depends strongly on X , then $\sigma_{xy} = \alpha\sigma_x$.
- Covariance is **positive** if X and Y increase or decrease together and **negative** if Y decreases when X increases.
- The **correlation coefficient** defined as

$$\rho = \frac{\sigma_{xy}}{\sigma_x\sigma_y}$$

is normalized covariance and takes values between -1 and +1

Vector random variables

- If we have more than two variables it is convenient to represent them in a vector notation.
- \mathbf{x} denotes a random vector with components x_1, x_2, \dots, x_d (d-dimensional vector).
- The joint distribution can be obtained and will in general be complicated. If the variables are statistically independent, the joint probability mass function, $P(\mathbf{x})$ is the product,

$$\begin{aligned} P(\mathbf{x}) &= P_{x_1}(x_1)P_{x_2}(x_2) \cdots P_{x_d}(x_d) \\ &= \prod_{i=1}^d P_{x_i}(x_i) \end{aligned}$$

The marginal distributions, $P_{x_i}(x_i)$ are obtained by summing the joint distribution over the other variables.

Expectation, Mean Vectors and Covariance Matrices

- The expected value of a vector is defined as the vector whose components are the expected values of the original components.
- The d -dimensional mean vector μ is defined by,

$$\mu = E[\mathbf{x}] = \begin{bmatrix} E[x_1] \\ E[x_2] \\ \vdots \\ E[x_d] \end{bmatrix} = \begin{bmatrix} \mu_1 \\ \mu_2 \\ \vdots \\ \mu_d \end{bmatrix} = \sum_x \mathbf{x}P(\mathbf{x})$$

Expectation, Mean Vectors and Covariance Matrices

- The **covariance matrix** denoted by Σ , is defined as

$$\sigma_{ij} = \sigma_{ji} = E[(x_i - \mu_i)(x_j - \mu_j)] \quad i, j = 1, \dots, d$$

$$\begin{aligned} \Sigma &= \begin{bmatrix} E[(x_1 - \mu_1)(x_1 - \mu_1)] & \dots & E[(x_1 - \mu_1)(x_d - \mu_d)] \\ E[(x_2 - \mu_2)(x_1 - \mu_1)] & \dots & E[(x_2 - \mu_2)(x_d - \mu_d)] \\ \vdots & \ddots & \vdots \\ E[(x_d - \mu_d)(x_1 - \mu_1)] & \dots & E[(x_d - \mu_d)(x_d - \mu_d)] \end{bmatrix} \\ &= \begin{bmatrix} \sigma_{11} & \sigma_{12} & \dots & \sigma_{1d} \\ \sigma_{21} & \sigma_{22} & \dots & \sigma_{2d} \\ \vdots & \vdots & \ddots & \vdots \\ \sigma_{d1} & \sigma_{d2} & \dots & \sigma_{dd} \end{bmatrix} = \begin{bmatrix} \sigma_1^2 & \sigma_{12} & \dots & \sigma_{1d} \\ \sigma_{21} & \sigma_2^2 & \dots & \sigma_{2d} \\ \vdots & \vdots & \ddots & \vdots \\ \sigma_{d1} & \sigma_{d2} & \dots & \sigma_d^2 \end{bmatrix} \end{aligned}$$

Expectation, Mean Vectors and Covariance Matrices

- The covariance matrix, Σ , can be written as the vector outer product,

$$\Sigma = E[(\mathbf{x} - \mu)(\mathbf{x} - \mu)^t]$$

- The covariance matrix is symmetric and the diagonal elements are the variances of the individual elements of \mathbf{x} . They are always positive.
- The off-diagonal entries of Σ are the covariances and they are positive or negative.

Expectation, Mean Vectors and Covariance Matrices

- If the variables are statistically independent, the covariances are zero and the covariance matrix is diagonal,

$$\Sigma = \begin{bmatrix} \sigma_1^2 & 0 & \dots & 0 \\ 0 & \sigma_2^2 & & 0 \\ \vdots & & \ddots & \vdots \\ 0 & 0 & \dots & \sigma_d^2 \end{bmatrix}$$

- For any d -dimensional vector \mathbf{w} , if the quadratic form

$$\mathbf{w}^t \Sigma \mathbf{w} > 0$$

Σ is positive definite. It is positive semi-definite if,

$$\mathbf{w}^t \Sigma \mathbf{w} \geq 0$$

Continuous Random Variables

- If the random variable can take values in the continuum we refer to it as a continuous random variable and describe it in terms of its probability density function.
- Rather than the probability mass function $P(X)$ we define the probability density function $p(x)$ with the property,

$$Pr[x \in (a, b)] = \int_a^b p(x) dx.$$

This is the probability that the continuous random variable, x , has value lying in the interval (a, b)

- The following is true:

$$p(x) \geq 0 \quad \text{and} \quad \int_{-\infty}^{\infty} p(x) = 1$$

Normal Distributions

- The normal or Gaussian probability density function is defined (for a single random variable) by,

$$p(x) = \frac{1}{\sqrt{2\pi}\sigma} \exp^{-1/2((x-\mu)^2/\sigma^2)}$$

- The following expected values can be computed,

$$E[1] = \int_{-\infty}^{\infty} p(x) dx = 1$$

$$E[x] = \int_{-\infty}^{\infty} xp(x) dx = \mu$$

$$E[(x - \mu)^2] = \int_{-\infty}^{\infty} (x - \mu)^2 p(x) dx = \sigma^2$$

Summary of key points and relations - I

- 1 Sample space, (S) - the collection of possible outcomes for an experiment
- 2 Event - A subset of the sample space
- 3 Mutually exclusive events - $A \cap B = \emptyset$
- 4 Probability of event - real number associated with the event that represents relative frequency of occurrence of that event
- 5 $P(S) = 1$
- 6 $P(A) \geq 0$
- 7 $P(A \cup B) = P(A) + P(B)$ if $A \cap B = \emptyset$
- 8 $P(B_1 \cup B_2 \cup \dots \cup B_n \dots) = \sum_{i=1}^{\infty} P(B_i)$ if $B_i \cap B_j = \emptyset$ for all $i \neq j$
 - $P(\emptyset) = 0$
 - $P(\bar{A}) = 1 - P(A)$
 - $P(A \cup B) = P(A) + P(B) - P(A \cap B)$
- 9 Equally likely rule : $P(A) = \frac{\text{number of elements in } A}{\text{number of elements in } S}$
- 10 Conditional probability: $P(A|B) = \frac{P(A \cap B)}{P(B)}$
- 11 $P(A \cap B \cap C) = P(A)P(B|A)P(C|A \cap B)$

Summary of key points and relations - II

- 12 Partition: $E_1 \cup E_2 \cup \dots \cup E_n = S$ and $E_i \cap E_j = \emptyset$ if $i \neq j$
- 13 Bayes Theorem: $P(E_i|A) = \frac{P(E_i)P(A|E_i)}{\sum_{j=1}^n P(E_j)P(A|E_j)}$ if E_1, E_2, \dots, E_n is a partition of S .
- 14 Conditional independence: $P(A|C) = P(A|B, C)$;
 $P(A, B|C) = P(A|C) \times P(B|C)$; A and B are conditionally independent given C .
A typical situation where conditional independence holds between events, is when the events share a common cause.
- 15 Chain rule with conditional independence: Consider the usual chain rule

$$P(E_i, E_2, \dots, E_n|A = a) = P(E_1|A = a) \times P(E_2, E_1|A = a) \\ \times P(E_n, E_{n-1}, \dots, E_1|A = a)$$

Under the assumption of **conditional independence** (i.e. the E_i have a common cause in the event $A = a$;

$$P(E_i, E_2, \dots, E_n|A = a) = P(E_1|A = a) \times P(E_2|A = a) \times \dots \times P(E_n|A = a)$$