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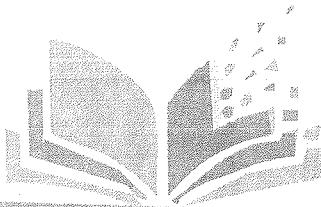
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YOUVA

SUBJECT: Multivariate Analysis

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Multivariate Analysis

In the study of statistics the variable under consideration so far were one or two or three and the objective of study was some characteristics of a group of individuals or objects. Then we have discussed descriptive and inferential statistics.

We may study more than three variables together so that it will generate a multivariate data in the form of a vector

$X_{p \times 1} = \begin{bmatrix} x_1 \\ \vdots \\ x_p \end{bmatrix}$, suppose we have collected the data on 'n' individuals / objects then we shall have 'n' 'p' components vectors as the observations and we have the observation matrix $X_{p \times n}$

The analysis of Multivariate data is known as the Multivariate Analysis.

The Variate

It is a linear combination of the variables with empirically determined weights. The variables are specified by the researcher whereas the weights are determined by Multivariate techniques to meet the specific objectives. A variate of the 'n' variables x_1, x_2, \dots, x_n can be stated Mathematically as variate value = $w_1x_1 + w_2x_2 + \dots + w_nx_n$ where x_n is the observed value and w_n is the weight determined by Multivariate techniques

* Analysis

$$\underline{x} = \begin{bmatrix} x_{11} & x_{12} & \dots & x_{1n} \\ x_{21} & x_{22} & \dots & x_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ x_{p1} & x_{p2} & \dots & x_{pn} \end{bmatrix}_{p \times n}$$

Define : $E(\underline{x}) = \begin{bmatrix} E(x_1) \\ E(x_2) \\ \vdots \\ E(x_n) \end{bmatrix}_{n \times 1} = \begin{bmatrix} \mu_1 \\ \mu_2 \\ \vdots \\ \mu_n \end{bmatrix}_{n \times 1}$

where. $E(x_i) = \mu_i$

We have $\text{Var}(\underline{x}) = E[(\underline{x} - E(\underline{x}))^2]$

and Multivariate case

$$\text{Var}(\underline{x}) = E[(\underline{x} - E(\underline{x}))(\underline{x} - E(\underline{x}))']$$

Let $\Sigma_{p \times p}$ = be the Variance - covariance matrix
of \underline{x} then (Dispersion matrix)

$$\Sigma = E \begin{bmatrix} x_1 - \mu_1 \\ x_2 - \mu_2 \\ \vdots \\ x_p - \mu_p \end{bmatrix} \begin{bmatrix} x_1 - \mu_1 & x_2 - \mu_2 & \dots & x_p - \mu_p \end{bmatrix}'$$

$$= E \begin{bmatrix} (x_1 - \mu_1)^2 & (x_1 - \mu_1)(x_2 - \mu_2) & \dots & (x_1 - \mu_1)(x_p - \mu_p) \\ (x_2 - \mu_2)(x_1 - \mu_1) & (x_2 - \mu_2)^2 & & \\ \vdots & & & \\ (x_p - \mu_p)(x_1 - \mu_1) & (x_p - \mu_p)(x_2 - \mu_p) & \dots & (x_p - \mu_p)^2 \end{bmatrix}$$

OR

$$\Sigma = \begin{bmatrix} \sigma_1^2 & \sigma_{12} & \dots & \sigma_{1p} \\ \sigma_{12} & \sigma_{22} & & \\ \vdots & & & \\ \sigma_{1p} & \sigma_{2p} & \dots & \sigma_{pp}^2 \end{bmatrix}$$

where $\sigma_i^2 = E(x_i - \bar{x}_i)^2 \quad i=1, \dots, p$

$\sigma_{ij} = E(x_i - \bar{x}_i)(x_j - \bar{x}_j) \quad i \neq j$

if x_1, x_2, \dots, x_p are independent then

$$\Sigma_{p \times p} = \begin{bmatrix} \sigma_1^2 & 0 & \cdots & 0 \\ 0 & \sigma_2^2 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & \sigma_p^2 \end{bmatrix}$$

Theorem: The variance-covariance matrix $\Sigma_{p \times p}$ is at least positive semidefinite.

Proof :- By definition we have $\Sigma = E(x - \bar{x})(x - \bar{x})'$
Define

$$Y = \begin{bmatrix} y_1 \\ y_2 \\ \vdots \\ y_p \end{bmatrix} \quad \text{and} \quad d = \begin{bmatrix} d_1 \\ d_2 \\ \vdots \\ d_p \end{bmatrix}$$

Also transform $Y = (x - \bar{x})$ then

$$E(Y) = E(x - \bar{x}) = 0$$

$$\text{then } d' \Sigma d = d'E(Y Y')d = E(d' Y Y' d)$$

Also

$$d' Y = Z = d_1 y_1 + d_2 y_2 + \dots + d_p y_p$$

$$E(Z) = E(d' Y) = d'E(Y) = 0$$

$$\text{Var}(Z) = E(Z^2) = (E(Z))^2 = \text{Var}(Z) = E(Z^2) \geq 0$$

$\therefore \Sigma$ is at least positive semidefinite

Note :- $\Sigma_{p \times p}$ is square symmetric and positive semidefinite Matrix

* Multivariate Normal distribution

For a random variable X the univariate normal density function can be written as

$$f(x; \mu, \sigma^2) = \frac{1}{\sigma \sqrt{2\pi}} e^{-\frac{1}{2} \left(\frac{x-\mu}{\sigma}\right)^2}$$

which can be written in a general form

$$f(x; \mu, \sigma^2) = k e^{-\frac{1}{2} (x-\mu)^2} \quad \text{--- (1)}$$

where σ is positive and the constant ' k ' is determined such that the integral of eq (1) is unity.

The density function of Multivariate Normal Distribution (MVND) has an analogous form as that of a univariate Normal distribution.

Let x_1, x_2, \dots, x_p are p -independent normal variate and $x_i \sim N(\mu_i, \sigma_i^2)$ then the p.d.f of (x_1, x_2, \dots, x_p) is given by

$$\begin{aligned} f(x_1, x_2, \dots, x_p) &= \prod_{i=1}^p \frac{1}{\sqrt{2\pi} \sigma_i} e^{-\frac{1}{2} \left(\frac{x_i - \mu_i}{\sigma_i}\right)^2} \quad \text{--- (*)} \\ &= \frac{1}{(2\pi)^{p/2} \sigma_1 \sigma_2 \dots \sigma_p} e^{-\frac{1}{2} \sum_{i=1}^p \left(\frac{x_i - \mu_i}{\sigma_i}\right)^2} \end{aligned}$$

As x_i 's are independent $\Sigma = \begin{bmatrix} \sigma_1^2 & 0 & \dots & 0 \\ 0 & \sigma_2^2 & \dots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \dots & \sigma_p^2 \end{bmatrix}$

and $\Sigma^{-1} = \text{diag} \begin{pmatrix} 1 & \dots & 1 \\ \sigma_1^2 & & \sigma_p^2 \end{pmatrix}$

as x is replaced by \underline{x} consider the quantity

$$= (\underline{x} - \underline{\mu})' \Sigma^{-1} (\underline{x} - \underline{\mu})$$

$$= (\underline{x}_1 - \mu_1, \underline{x}_2 - \mu_2, \dots, \underline{x}_p - \mu_p)' \begin{bmatrix} \frac{1}{\sigma_1^2} & 0 & \dots & 0 \\ 0 & \frac{1}{\sigma_2^2} & & \\ \vdots & & \ddots & \\ 0 & 0 & \dots & \frac{1}{\sigma_p^2} \end{bmatrix} \begin{bmatrix} \underline{x}_1 - \mu_1 \\ \underline{x}_2 - \mu_2 \\ \vdots \\ \underline{x}_p - \mu_p \end{bmatrix}$$

$$= \sum_{i=1}^p (\underline{x}_i - \mu_i)^2$$

$$\text{Also, } |\Sigma|^{1/2} = \sigma_1 \cdot \sigma_2 \cdots \sigma_p$$

Substituting the above equation quantities in (*) we get

$$f(\underline{x}) = \frac{1}{(2\pi)^{p/2} |\Sigma|^{1/2}} e^{-(\underline{x} - \underline{\mu})' \Sigma^{-1} (\underline{x} - \underline{\mu})}$$

Notation : $\underline{x} \sim N_p(\underline{\mu}, \Sigma)$

* Derivation (General form)

In general the exponential family of distribution in a Multivariate form can be written as

$$f(\underline{x}) = K e^{-\frac{1}{2} (\underline{x} - \underline{B})' A (\underline{x} - \underline{B})}$$

where \underline{B} : Vector, A : the definite matrix (symmetric)

as A is $(\underline{x} - \underline{B})' A (\underline{x} - \underline{B}) > 0$

$e^{-\frac{1}{2} (\underline{x} - \underline{B})' A (\underline{x} - \underline{B})} \leq 1 \Rightarrow f(\underline{x})$ will be bounded fn.

* Determination of K

Consider the transformation $(x - B) = CY$
 where 'C' is a Non-Singular matrix (NSM)
 and is determined by using the result
 that : if 'A' is a positive definite matrix
 then there exists a NSM 'C' such that

$$C^T A C = I \Rightarrow |C| |A| |C| = |I|$$

Since $(x - B) = CY$

\therefore Jacobian of the transformation is $|J| = |C|$

$$\text{As } f(x) = k e^{-\frac{1}{2} (x - B)^T A (x - B)}$$

$$\therefore f(Y) = k e^{-\frac{1}{2} Y^T C^T A C Y} \cdot |J|$$

$$\text{Or, } f(Y) = k e^{-\frac{1}{2} Y^T Y} \cdot |J| \quad \text{as } C^T A C = I$$

$$\text{Or, } f(Y) = k e^{-\frac{1}{2} \sum_{i=1}^p y_i^2} |C|$$

$$\therefore \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \cdots \int_{-\infty}^{\infty} k e^{-\frac{1}{2} \sum_{i=1}^p y_i^2} |C| dy_1 \cdots dy_p = 1$$

$$\text{Or, } k |C| \prod_{i=1}^p \int_{-\infty}^{\infty} e^{-\frac{1}{2} y_i^2} dy_i = 1$$

$$\therefore k |C| (\sqrt{2\pi})^p = 1$$

As

$$\frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-\frac{1}{2} y_i^2} dy_i = 1$$

$$\Rightarrow \int_{-\infty}^{\infty} e^{-\frac{1}{2} y_i^2} dy_i = \sqrt{2\pi}$$

$$\therefore k = \frac{1}{(2\pi)^{p/2} |c|}$$

$$\text{Again } f(y) = \frac{1}{(2\pi)^{p/2} |c|} e^{-\frac{1}{2} y^T y} \cdot \text{let}$$

$$= \frac{1}{(2\pi)^{p/2}} e^{-\frac{1}{2} y^T y}$$

$$\text{as } C^T A C = I \Rightarrow |C|^2 |A| |C| = 1$$

$$\Rightarrow |C|^2 |A| = 1 \quad (\text{as } |C|^2 = |C|)$$

$$\Rightarrow |C| = \frac{1}{|A|^{1/2}}$$

$$(k = \frac{1}{(2\pi)^{p/2}} \frac{1}{|A|^{1/2}})$$

* Determination of 'B'

$$\text{we have } f(y) = \left(\frac{1}{2\pi}\right)^{p/2} e^{-\frac{1}{2} y^T y}$$

Consider

$$E(y_i) = \int_{-\infty}^{\infty} y_i \frac{1}{\sqrt{2\pi}} e^{-\frac{1}{2} y_i^2} dy = \frac{1}{(2\pi)^{p/2}} \sum_{j=1}^p \int_{-\infty}^{\infty} y_j \frac{1}{\sqrt{2\pi}} e^{-\frac{1}{2} y_j^2} dy$$

$$= 0 \times 1 = 0$$

$$\therefore E(y_i) = 0 \quad \forall i = 1, \dots, p.$$

$$\therefore E(Y) = 0 \quad (\text{Null vector})$$

Since $(\underline{x} - \underline{B}) = C\underline{Y} \Rightarrow \underline{x} = \underline{B} + C\underline{Y}$

$$\therefore E(\underline{x}) = \underline{B} + C E(\underline{Y}) = \underline{B}$$

But $E(\underline{x})$ is the mean vector \underline{x} denoted by

\underline{u}

$$\therefore \underline{B} = \underline{u}$$

* Determination of 'A'

By defⁿ we have.

$$\frac{1}{2} = E(\underline{x} - \underline{u})(\underline{x} - \underline{u})'$$

$$= E(\underline{x} - \underline{B})(\underline{x} - \underline{B})'$$

$$= E(C\underline{Y}\underline{Y}'C) \quad \text{OR} \quad \frac{1}{2} = CE(\underline{Y}\underline{Y}')C.$$

in order to obtain ^{this} consider the 2nd moment
of \underline{Y} and obtain $E(Y_j, Y_k)$ $j \neq k$

Now

$$E(Y_j, Y_k) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} Y_j Y_k \perp \frac{e^{-\frac{1}{2} \underline{y}^T \underline{y}}}{(2\pi)^{p/2}} dy_p \dots dy_1$$

$$= \left[\int_{-\infty}^{\infty} Y_j \perp \frac{1}{\sqrt{2\pi}} e^{-\frac{y_j^2}{2}} dy_j \right] \left[\int_{-\infty}^{\infty} Y_k \perp \frac{1}{\sqrt{2\pi}} e^{-\frac{y_k^2}{2}} dy_k \right]$$

$$\prod_{\substack{j=1 \\ j \neq j+k}}^p \left[\int_{-\infty}^{\infty} \frac{1}{(2\pi)^{p/2}} e^{-\frac{1}{2} \underline{y}^T \underline{y}} dy \right]$$

$\therefore \text{or } E(y_j, y_k) = 0 \text{ for } j \neq k$

if $j = k$,

$$E(y_j^2) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} y_j^2 \frac{1}{\sqrt{2\pi}} e^{-\frac{1}{2} y_j^2} dy_j$$

$$\frac{P}{\prod_{k=1}^p} \left[\int_{-\infty}^{\infty} \frac{1}{(2\pi)^{p-1}} e^{-\frac{1}{2} y_k^2} dy_k \right]$$

$$\therefore E(y_j, y_k) = 0, j \neq k \\ = 1, j = k$$

$$\begin{aligned} \text{By Defn } VCM(\mathbf{Y}) &= E[\mathbf{Y} - E(\mathbf{Y})][\mathbf{Y} - E(\mathbf{Y})]^T \\ &= E[\mathbf{Y} - \mathbf{0}][\mathbf{Y} - \mathbf{0}]^T \\ &= E[\mathbf{Y}'\mathbf{Y}] \end{aligned}$$

(VCM = Variance - covariance Matrix)

$$\text{Again } (\mathbf{x} - \mathbf{B}) = \mathbf{C}\mathbf{Y} \text{ or } (\mathbf{x} - \mathbf{u}) = \mathbf{C}\mathbf{Y} \\ \Rightarrow \mathbf{Y} = \mathbf{C}^{-1}(\mathbf{x} - \mathbf{u})$$

so,

$$\begin{aligned} E[\mathbf{C}^{-1}(\mathbf{x} - \mathbf{u})(\mathbf{x} - \mathbf{u})'(\mathbf{C}^{-1})'] &= I_p \\ \text{or } \mathbf{C}^{-1} E(\mathbf{x} - \mathbf{u})(\mathbf{x} - \mathbf{u})'(\mathbf{C}^{-1})' &= I_p \\ \text{and } E(\mathbf{x} - \mathbf{u})(\mathbf{x} - \mathbf{u})' &= \mathbf{S} \end{aligned}$$

$$\therefore \mathbf{C}^{-1}(\mathbf{C}^{-1})' = I_p \Rightarrow \mathbf{C}^{-1} = \mathbf{C}\mathbf{C}' \rightarrow (1)$$

$$\begin{aligned} \text{Also } \mathbf{C}'\mathbf{A}\mathbf{C} &= I_p \Rightarrow \mathbf{A} = (\mathbf{C}\mathbf{C}')^{-1} I_p \\ &\Rightarrow \mathbf{A} = \mathbf{C}' I_p \mathbf{C}^{-1} \rightarrow (2) \end{aligned}$$

using (1) and (2) $\boxed{\mathbf{A} = \mathbf{C}^{-1}}$

recall that $K \in |A|^{1/2}$ $\therefore K = \frac{1}{(2\pi)^{p/2}}$ $(2\pi)^{p/2} |A|^{1/2}$

Using the values of K , A , B in $f(x)$
we obtain

$$f(x) = \frac{1}{(2\pi)^{p/2} |A|^{1/2}} e^{-\frac{1}{2}(x-u)^T A^{-1}(x-u)}$$

i.e. the p.d.f of random vector $X \sim N_p(\mu, \Sigma)$

the parameters are mean vector: μ_{px}

Var-Covariance Matrix: Σ_{pxp}

Note: If X is multivariate than each component x_i should be univariate.

* Bivariate Normal Distribution

$$\text{Univariate } f(x) = \frac{1}{\sqrt{2\pi \sigma^2}} e^{-\frac{1}{2} \left(\frac{x-\mu}{\sigma}\right)^2}$$

Case 1: When x_1 and x_2 are independent

If x_1 and x_2 are independent

$$f(x_1, x_2), x = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}, \mu = \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix}$$

$$\Sigma = \begin{bmatrix} \sigma_1^2 & \sigma_{12} \\ \sigma_{12} & \sigma_2^2 \end{bmatrix} = \begin{bmatrix} \sigma_1^2 & 0 \\ 0 & \sigma_2^2 \end{bmatrix}$$

($\because x_1$ and x_2 are independent their covariance is zero)

Now,

$$F(x_1, x_2) = f(x_1) \cdot f(x_2)$$

$$= \frac{1}{\sqrt{2\pi\sigma_1^2}} e^{-\frac{1}{2}\left(\frac{x_1-\mu_1}{\sigma_1}\right)^2} \times \frac{1}{\sqrt{2\pi\sigma_2^2}} e^{-\frac{1}{2}\left(\frac{x_2-\mu_2}{\sigma_2}\right)^2}$$

$$= \frac{1}{(2\pi)^{1/2} (\frac{1}{2})^{1/2}} e^{-\frac{1}{2}\left[\left(\frac{x_1-\mu_1}{\sigma_1}\right)^2 + \left(\frac{x_2-\mu_2}{\sigma_2}\right)^2\right]}$$

$$= \frac{1}{(2\pi)^{1/2} (\frac{1}{2})^{1/2}} e^{-\frac{1}{2}(x-\mu_1)^2 - \frac{1}{2}(x-\mu_2)^2}$$

→ Bivariate case

$$\text{Now } \frac{1}{2} = \sigma_1^2 \sigma_2^2$$

$$\underline{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}, \underline{\mu} = \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix}, \frac{1}{2} = \begin{bmatrix} \sigma_1^2 & 0 \\ 0 & \sigma_2^2 \end{bmatrix}$$

Now deriving exponential part

$$= -\frac{1}{2} [(x_1-\mu_1)(x_2-\mu_2)] \begin{bmatrix} \sigma_1^2 & 0 \\ 0 & \sigma_2^2 \end{bmatrix}^{-1} \begin{bmatrix} x_1-\mu_1 \\ x_2-\mu_2 \end{bmatrix}$$

$$= -\frac{1}{2} [(x_1-\mu_1)(x_2-\mu_2)] \frac{1}{\sigma_1^2 \sigma_2^2} \begin{bmatrix} \sigma_2^2 & 0 \\ 0 & \sigma_1^2 \end{bmatrix} \begin{bmatrix} (x_1-\mu_1) \\ (x_2-\mu_2) \end{bmatrix}$$

$$= -\frac{1}{2\sigma_1^2 \sigma_2^2} [(x_1-\mu_1)(x_2-\mu_2)] \begin{bmatrix} \sigma_2^2 (x_1-\mu_1) \\ \sigma_1^2 (x_2-\mu_2) \end{bmatrix}$$

$$= -\frac{\sigma_1^2 \sigma_2^2}{2\sigma_1^2 \sigma_2^2} \left[\left(\frac{x_1-\mu_1}{\sigma_1}\right)^2 + \left(\frac{x_2-\mu_2}{\sigma_2}\right)^2 \right]$$

{Multiplying and dividing $\sigma_1^2 \sigma_2^2$ }

$$= -\frac{1}{2} \left[\left(\frac{x_1 - \mu_1}{\sigma_1} \right)^2 + \left(\frac{x_2 - \mu_2}{\sigma_2} \right)^2 \right]$$

$$f(x_1, x_2) = \frac{1}{(2\pi)^{1/2} |\Sigma|^{1/2}} e^{-\frac{1}{2} \left[\left(\frac{x_1 - \mu_1}{\sigma_1} \right)^2 + \left(\frac{x_2 - \mu_2}{\sigma_2} \right)^2 \right]}$$

Case 2 :- At $p=2$ define the Bivariate Normal Distribution using $N_p(\mu, \Sigma)$

$$\underline{x} = \begin{pmatrix} x_1 \\ x_2 \end{pmatrix}, \underline{\mu} = \begin{pmatrix} \mu_1 \\ \mu_2 \end{pmatrix}, \rho_{12} \text{ cov}(x_1, x_2)$$

Here x_1 and x_2 are dependent

$$\rightarrow \underline{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}, \underline{\mu} = \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix}, \Sigma = \begin{bmatrix} \sigma_1^2 & \rho_{12} \\ \rho_{12} & \sigma_2^2 \end{bmatrix}$$

$$f(x_1, x_2) = \frac{1}{(2\pi)^{1/2} |\Sigma|^{1/2}} e^{-\frac{1}{2} \left[(\underline{x} - \underline{\mu})^\top \Sigma^{-1} (\underline{x} - \underline{\mu}) \right]} \quad (1)$$

Now deriving exponential term

$$\Sigma = \begin{bmatrix} \sigma_1^2 & \rho_{12} \\ \rho_{12} & \sigma_2^2 \end{bmatrix}$$

$$\begin{aligned} |\Sigma| &= \sigma_1^2 \sigma_2^2 - \rho_{12}^2 \\ &= \sigma_1^2 \sigma_2^2 - (\rho_{12} \sigma_1 \sigma_2)^2 \quad \left(\because \rho = \frac{\rho_{12}}{\sigma_1 \sigma_2} \right) \\ &= \sigma_1^2 \sigma_2^2 (1 - \rho_{12}^2) \quad (2) \end{aligned}$$

$$\Sigma^{-1} = \frac{1}{\sigma_1^2 \sigma_2^2 - \rho_{12}^2} \begin{bmatrix} \sigma_2^2 & -\rho_{12} \\ -\rho_{12} & \sigma_1^2 \end{bmatrix}$$

$$\text{Now, } -\frac{1}{2} \left[(\underline{x} - \underline{u}) \cdot \underline{\$}^{-1} (\underline{x} - \underline{u}) \right]$$

Substituting value of $\underline{\$}^{-1}$

$$-\frac{1}{2} \left[(\underline{x} - \underline{u})' \right], \quad \frac{1}{6_1^2 6_2^2 - 6_{12}^2} \begin{bmatrix} 6_2^2 & -6_{12} \\ -6_{12} & 6_1^2 \end{bmatrix} \left[(\underline{x} - \underline{u}) \right]$$

$$= \frac{1}{2 6_1^2 6_2^2 - 6_{12}^2} \left[(x_1 - u_1)(x_2 - u_2) \right] \begin{bmatrix} 6_2^2 & -6_{12} \\ -6_{12} & 6_1^2 \end{bmatrix} \begin{bmatrix} (x_1 - u_1) \\ (x_2 - u_2) \end{bmatrix}$$

$$= \frac{1}{2 6_1^2 6_2^2 - 6_{12}^2} \left[(x_1 - u_1)(x_2 - u_2) \right] \begin{bmatrix} 6_2^2(x_1 - u_1) - 6_{12}(x_2 - u_2) \\ 6_1^2(x_2 - u_2) - 6_{12}(x_1 - u_1) \end{bmatrix}$$

$$= \frac{1}{2 6_1^2 6_2^2 - 6_{12}^2} \left[(x_1 - u_1)^2 6_2^2 - 2 6_{12}(x_1 - u_1)(x_2 - u_2) + 6_1^2(x_2 - u_2)^2 \right]$$

$$= \frac{-1}{2} \frac{6_1^2 6_2^2}{6_1^2 6_2^2 - 6_{12}^2} \left[\left(\frac{x_1 - u_1}{6_1} \right)^2 - 2 \frac{6_{12}}{6_1 6_2} \left(\frac{x_1 - u_1}{6_1} \right) \left(\frac{x_2 - u_2}{6_2} \right) + \frac{(x_2 - u_2)^2}{6_2^2} \right]$$

[Multiplying $6_1 6_2$ to numerator and denominator]

$$= -\frac{1}{2} \frac{6_1^2 6_2^2}{6_1^2 6_2^2 - 6_{12}^2 6_1^2 6_2^2}$$

$$\left[\left(\frac{x_1 - u_1}{6_1} \right)^2 - 2 \frac{6_{12}}{6_1} \left(\frac{x_1 - u_1}{6_1} \right) \left(\frac{x_2 - u_2}{6_2} \right) + \left(\frac{x_2 - u_2}{6_2} \right)^2 \right]$$

By using eqn(2)

$$= \frac{-1}{2(1 - \beta_{12}^2)} \left[\left(\frac{x_1 - u_1}{6_1} \right)^2 - 2 \beta_{12} \left(\frac{x_1 - u_1}{6_1} \right) \left(\frac{x_2 - u_2}{6_2} \right) + \left(\frac{x_2 - u_2}{6_2} \right)^2 \right]$$

Thus,

$$f(x_1, x_2) = \frac{1}{(2\pi)^{2/2} |\frac{1}{2}|^{1/2}} e^{-\frac{1}{2(1-\rho_{12}^2)} \left[\left(\frac{x_1 - \mu_1}{\sigma_1}\right)^2 - 2\rho_{12} \left(\frac{x_1 - \mu_1}{\sigma_1}\right) \left(\frac{x_2 - \mu_2}{\sigma_2}\right) + \left(\frac{x_2 - \mu_2}{\sigma_2}\right)^2 \right]}$$

* Properties of MVND:- (17/8/22)

Scale & ORIGIN.

① Given that a random vector $\underline{X} \sim N_p(\underline{\mu}, \frac{1}{2})$

Consider $\underline{Y} = C\underline{X}$ where $C_{p \times p}$ NSM then

$\underline{Y} \sim N_p(C\underline{\mu}, C \frac{1}{2} C')$

Proof Since, $\underline{Y} = C\underline{X}$ is a non-singular transformation
 $\therefore \underline{X} = C^{-1}\underline{Y}$ and for such type of
transformation $|J| = |C^{-1}|$

$$\therefore f(y) = \frac{1}{(2\pi)^{p/2} |\frac{1}{2}|^{1/2}} e^{-\frac{1}{2} (C^{-1}y - \underline{\mu})^T C^{-1} (C^{-1}y - \underline{\mu})} \cdot |C^{-1}|$$

$$\therefore f(y) = \frac{1}{(2\pi)^{p/2} |\underline{\mu}|^{1/2} |\frac{1}{2}|^{1/2} |C|^2}$$

$$\cdot e^{-\frac{1}{2} \frac{1}{2} C^{-1} (y - C\underline{\mu})^T C^{-1} (y - C\underline{\mu})}$$

$$\text{OR } f(y) = \frac{1}{(2\pi)^{p/2} |C \underline{\mu} C'|^{1/2}} e^{-\frac{1}{2} (y - C\underline{\mu})^T (C \underline{\mu} C')^{-1} (y - C\underline{\mu})}$$

\therefore Comparing $f(y)$ with standard form of $f(\underline{y})$, we note that

$$\underline{Y} \sim N_p(C\underline{\mu}, C \underline{\mu} C')$$

(2)

Characteristic function :-

If $\underline{x} \sim N_p(\underline{\mu}, \Sigma)$ then (show that)

$$\Phi_{\underline{x}}(t) = e^{it' \underline{\mu} - \frac{1}{2} t' \Sigma t} \quad \text{where } t \text{ px, real vector}$$

Proof

As Σ assumed to a NSM,

Σ^{-1} exists and hence there exists a NSM

'C' such that $C' \Sigma^{-1} C = I_p$

$$\Rightarrow \Sigma^{-1} = (C')^{-1} C^{-1} \quad \text{or} \quad \Sigma^{-1} = (C C')^{-1}$$

$$\Rightarrow \boxed{\Sigma = C C'}. \rightarrow (\text{Imp for 2marks}) \text{ Prove that } \Sigma = C C'$$

Again consider transformation

$$(\underline{x} - \underline{\mu}) = C \underline{y} \Rightarrow \underline{x} = \underline{\mu} + C \underline{y}$$

$$\text{Also } E(\underline{x} - \underline{\mu}) = C E(\underline{y}) = \underline{0}$$

$$\Rightarrow E(\underline{y}) = \underline{0} \quad \text{as} \quad E(\underline{x} - \underline{\mu}) = \underline{0}$$

Now the var-cov matrix of \underline{y} is

$$E(\underline{y} \underline{y}') = E[C^{-1}(\underline{x} - \underline{\mu})(\underline{x} - \underline{\mu})' C (C^{-1})']$$

or

$$\begin{aligned} E(\underline{y}' \underline{y}) &= C^{-1} E[(\underline{x} - \underline{\mu})(\underline{x} - \underline{\mu})'] (C^{-1})' \\ &= C^{-1} \Sigma (C^{-1})' \\ &= C^{-1} (C C') (C^{-1})' \\ &= I_p \end{aligned}$$

$$\therefore \underline{y} \sim N_p(\underline{0}, I_p)$$

as

$$\underline{y} \sim N_p(\underline{0}, I_p) \text{ it } \Rightarrow y_i \sim N(0, 1)$$

By defⁿ $\phi_{\underline{x}}(t) = E[e^{it' \underline{x}}]$

$$= E[e^{it'(\mu + cY)}]$$

$$= e^{it'\mu} \cdot E[e^{it'cY}]$$

OR

$$\phi_{\underline{x}}(t) = e^{it'\mu} \cdot E[e^{it'cY}]$$

OR

$$\phi_{\underline{x}}(t) = e^{it'\mu} \cdot E[e^{i \sum_{j=1}^p u_j Y_j}]$$

$$= e^{it'\mu} \prod_{j=1}^p e^{iu_j Y_j}$$

real variable
real constant

$$= e^{it'\mu} \prod_{j=1}^p e^{iu_j Y_j}$$

$$= e^{it'\mu} \prod_{j=1}^p e^{-\frac{1}{2} u_j^2}$$

$$= e^{it'\mu} e^{-\frac{1}{2} \sum_j u_j^2}$$

$$= e^{it'\mu} e^{-\frac{1}{2} \underline{u} \underline{u}}$$

$$= e^{it'\mu} e^{-\frac{1}{2} t' c c t}$$

$$= e^{it'\mu - \frac{1}{2} t' c c t}$$

$$= e^{-\frac{1}{2} t' c c t} \rightarrow N(0, \frac{1}{2} c c t)$$

(Now $t' c = \underline{u}$ as $t'_{1:p} = c_p x_p \Rightarrow t' c = \underline{u}$ (say))

INDEPENDENT

③ If x_1, x_2, \dots, x_p have joint multivariate Normal distⁿ then Necessary and sufficient condition that one subset $x^{(1)}$ is independent of other subset $x^{(2)}$ containing the remaining variables is that each covariance of a variable from $x^{(1)}$ and a variable from $x^{(2)}$ be zero.

Mathematically, let

$$\underline{x}_{px1} = \begin{bmatrix} \underline{x}^{(1)} & \xrightarrow{\quad q} \\ \hline \underline{x}^{(2)} & \xrightarrow{(p-q)} \end{bmatrix}$$

accordingly $\underline{u}_{px1} = \begin{bmatrix} \underline{u}^{(1)} & \xrightarrow{\quad q} \\ \hline \underline{u}^{(2)} & \xrightarrow{(p-q)} \end{bmatrix}$

$$\underline{\$}_{q \times q} = \begin{bmatrix} \$_{11} & \$_{12} \\ \hline \$_{21} & \$_{22} \end{bmatrix}$$

where $\Sigma_{11} \Rightarrow q \times q = E[(\underline{x}^{(1)} - \underline{u}^{(1)}) (\underline{x}^{(1)} - \underline{u}^{(1)})']$

$$\Sigma_{12} = E[(\underline{x}^{(1)} - \underline{u}^{(1)}) (\underline{x}^{(2)} - \underline{u}^{(2)})']$$

$$\Sigma_{22} = E[(\underline{x}^{(2)} - \underline{u}^{(2)}) (\underline{x}^{(2)} - \underline{u}^{(2)})']$$

Proof Necessary Part :-

let $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ be independent and let x_i be a variable from the i^{th} 1st subset $\underline{x}^{(1)}$ and x_j be any variable from $\underline{x}^{(2)}$, then we have to show $\text{cov}(x_i, x_j) = 0$

By defⁿ

$$\text{cov}(x_i, x_j) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} (x_i - \bar{u}_i)(x_j - \bar{u}_j) f(x) dx_1 dx_2 \dots dx_p$$

Since $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ are independent. So the joint density of $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ will be the product of their marginal

i.e $f(\underline{x}) = f_1(\underline{x}^{(1)}) \cdot f_2(\underline{x}^{(2)})$ (say)

$$\therefore \text{Cov}(x_i, x_j) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} (x_i - \mu_i) f_1(x^{(1)}) dx^{(1)}$$

$$\int_{-\infty}^{\infty} (x_j - \mu_j) f_2(x^{(2)}) dx^{(2)}$$

$$\rightarrow E(x_i - \mu_i) \cdot E(x_j - \mu_j) = 0$$

$$\text{i.e } \text{Cov}(x_i, x_j) = 0 \quad \text{for } i = 1, 2, \dots, q \\ j = q+1, \dots, p$$

this means all the elements

of Σ_{12} are zero i.e $\Sigma_{12} = 0$

Sufficient Part

let $\Sigma_{12} = 0$ then we have to show that
 $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ are independent. Here

$$\underline{\Sigma} = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix} = \begin{bmatrix} \Sigma_{11} & 0 \\ 0 & \Sigma_{22} \end{bmatrix}$$

$$\Rightarrow \Sigma^{-1} = \begin{bmatrix} \Sigma_{11}^{-1} & 0 \\ 0 & \Sigma_{22}^{-1} \end{bmatrix}$$

$$\text{we have } f(\underline{x}) = \frac{1}{(2\pi)^{q/2} |\underline{\Sigma}|^{1/2}} e^{-\frac{1}{2} (\underline{x} - \underline{\mu})' \underline{\Sigma}^{-1} (\underline{x} - \underline{\mu})}$$

after partition consider the exponential terms:

$$(\underline{x} - \underline{\mu})' \underline{\Sigma}^{-1} (\underline{x} - \underline{\mu})$$

$$\begin{bmatrix} \underline{x}^{(1)} - \underline{\mu}^{(1)} \\ \underline{x}^{(2)} - \underline{\mu}^{(2)} \end{bmatrix}^T \begin{bmatrix} \Sigma_{11}^{-1} & 0 \\ 0 & \Sigma_{22}^{-1} \end{bmatrix} \begin{bmatrix} \underline{x}^{(1)} - \underline{\mu}^{(1)} \\ \underline{x}^{(2)} - \underline{\mu}^{(2)} \end{bmatrix}$$

$$(\underline{x}^{(1)} - \underline{\mu}^{(1)})^T \underline{\$}_{11}^{-1} (\underline{x}^{(1)} - \underline{\mu}^{(1)}) + (\underline{x}^{(2)} - \underline{\mu}^{(2)})^T \underline{\$}_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})$$

$$\therefore f_1(\underline{x}^{(1)}) = \frac{1}{(2\pi)^{q/2} |\underline{\$}|_{11}^{1/2}} e^{-\frac{1}{2} (\underline{x}^{(1)} - \underline{\mu}^{(1)})^T \underline{\$}_{11}^{-1} (\underline{x}^{(1)} - \underline{\mu}^{(1)})}$$

$$f_2(\underline{x}^{(2)}) = \frac{1}{(2\pi)^{(p-q)/2} |\underline{\$}|_{22}^{1/2}} e^{-\frac{1}{2} (\underline{x}^{(2)} - \underline{\mu}^{(2)})^T \underline{\$}_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})}$$

$$= N_q(\underline{\mu}^{(1)}, \underline{\$}_{11}) \times N_{p-q}(\underline{\mu}^{(2)}, \underline{\$}_{22})$$

$$\therefore f(\underline{x}) = f_1(\underline{x}^{(1)}) f_2(\underline{x}^{(2)})$$

$\Rightarrow \underline{x}^{(1)}$ and $\underline{x}^{(2)}$ are independent

④

Marginal Distribution

Let $\underline{x}_{px1} \sim N_p(\underline{\mu}, \underline{\$})$ and $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ are two subsets of \underline{x}_{px1} where $\underline{x}^{(1)}_{(p-q) \times 1}$ and $\underline{x}^{(2)}_{(q \times 1)}$. Then the Marginal dist'n of any subset of \underline{x}_{px1} is also MVND with mean vector and var-cov matrix obtained by taking proper components of $\underline{\mu}$ and $\underline{\$}$.

Proof: we have $\underline{x}_{px1} = \begin{bmatrix} \underline{x}^{(1)} \\ \underline{x}^{(2)} \end{bmatrix}$ acc. $\underline{\mu} = \begin{bmatrix} \underline{\mu}^{(1)} \\ \underline{\mu}^{(2)} \end{bmatrix}$

$$\underline{\$} = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ - & - \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix} \quad \text{consider the transformation}$$

$y^{(1)} = x^{(1)} + M(x^{(2)})^{(2)}$

$\& \quad y^{(2)} = x^{(2)}$

i.e. $\underline{Y} = \begin{bmatrix} \underline{I} & M \\ 0 & \underline{I} \end{bmatrix} \underline{X}$ where M is a NSM and
is chosen such that

$\underline{Y}^{(1)}$ and $\underline{Y}^{(2)}$ are independent

$\therefore \text{cov} [\underline{Y}^{(1)}, \underline{Y}^{(2)}] = 0$ determine M

$$\text{Now } \text{cov} [\underline{Y}^{(1)}, \underline{Y}^{(2)}] = E[\underline{Y}^{(1)} - E(\underline{Y}^{(1)})][\underline{Y}^{(2)} - E(\underline{Y}^{(2)})]$$

$$\text{Also, } E(\underline{Y}) = \begin{bmatrix} \underline{I} & M \\ 0 & \underline{I} \end{bmatrix} E(\underline{X})$$

$$\begin{aligned} \text{i.e. } E \begin{bmatrix} \underline{Y}^{(1)} \\ \underline{Y}^{(2)} \end{bmatrix} &= \begin{bmatrix} \underline{I} & M \\ 0 & \underline{I} \end{bmatrix} \begin{bmatrix} \underline{u}^{(1)} \\ \underline{u}^{(2)} \end{bmatrix} \\ &= \begin{bmatrix} \underline{u}^{(1)} + M\underline{u}^{(2)} \\ \underline{u}^{(2)} \end{bmatrix} \end{aligned}$$

So that

$$\begin{aligned} \text{cov} [\underline{Y}^{(1)}, \underline{Y}^{(2)}] &= E[(\underline{x}^{(1)} + M\underline{x}^{(2)} - \underline{u}^{(1)} - M\underline{u}^{(2)})' (\underline{x}^{(2)} - \underline{u}^{(2)})] \\ &= E[(\underline{x}^{(1)} - \underline{u}^{(1)}) (\underline{x}^{(2)} - \underline{u}^{(2)})'] + M E[\underline{x}^{(2)} - \underline{u}^{(2)}]' (\underline{x}^{(2)} - \underline{u}^{(2)}) \\ &= \Sigma_{12} + M \Sigma_{22} \end{aligned}$$

$$\text{Now, } \text{cov} (\underline{Y}^{(1)}, \underline{Y}^{(2)}) = 0$$

$$\Rightarrow \Sigma_{12} + M \Sigma_{22} = 0 \Rightarrow M = -\Sigma_{12} \Sigma_{22}^{-1}$$

$$\underline{Y}^{(1)} = \underline{x}^{(1)} - \Sigma_{12} \Sigma_{22}^{-1} \underline{x}^{(2)} \quad \text{OR} \quad \underline{Y} = \begin{bmatrix} \underline{I} & -\Sigma_{12} \Sigma_{22}^{-1} \\ 0 & \underline{I} \end{bmatrix} \underline{X}$$

$$\underline{Y}^{(2)} = \underline{x}^{(2)}$$

$$\text{OR } \underline{Y} = C \underline{X} \text{ (say) 'C' say being a NSM}$$

∴ Using property (1) if $\underline{x} \sim N_p(\mu, \Sigma)$

$\underline{y} \sim N_p(C\mu, C\Sigma C')$ we obtain

$$E(\underline{y}) = \begin{bmatrix} \underline{\mu}^{(1)} - \Sigma_{12} \Sigma_{22}^{-1} \underline{\mu}^{(2)} \\ \vdots \\ \underline{\mu}^{(2)} \end{bmatrix} = \begin{bmatrix} \underline{\sigma}_1 \\ \vdots \\ \underline{\sigma}_2 \end{bmatrix} = \underline{\sigma} \text{ (say)}$$

and

$$C \& C' = \begin{bmatrix} I & -\Sigma_{12} \Sigma_{22}^{-1} \\ 0 & I \end{bmatrix} \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix} \begin{bmatrix} I & -\Sigma_{12} \Sigma_{22}^{-1} \\ 0 & I \end{bmatrix}$$

$$= \begin{bmatrix} \Sigma_{11} - \Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21} & 0 \\ 0 & \Sigma_{22} \end{bmatrix}$$

$$= \begin{bmatrix} \Sigma_{11.2} & 0 \\ 0 & \Sigma_{22} \end{bmatrix} = \Sigma \rightarrow \text{Partial Var-Covariance Matrix}$$

$$\therefore f(\underline{y}) = \frac{1}{(2\pi)^{p/2} |\Sigma|^{1/2}} e^{-\frac{1}{2} (\underline{y} - \underline{\sigma})^T \Sigma^{-1} (\underline{y} - \underline{\sigma})}$$

$$= \frac{1}{(2\pi)^{q/2} |\Sigma_{11.2}|^{1/2}} e^{-\frac{1}{2} (\underline{y}^{(1)} - \underline{\sigma}^{(1)})^T \Sigma_{11.2}^{-1} (\underline{y}^{(1)} - \underline{\sigma}^{(1)})}$$

$$\times \frac{1}{(2\pi)^{p-q/2} |\Sigma_{22}|^{1/2}} e^{-\frac{1}{2} (\underline{y}^{(2)} - \underline{\sigma}^{(2)})^T \Sigma_{22}^{-1} (\underline{y}^{(2)} - \underline{\sigma}^{(2)})}$$

$$= f(y^{(1)}) \times f(y^{(2)}) \text{ as } y^{(1)} \text{ and } y^{(2)}$$

are independent

- (*)

As per the transformation $\underline{Y}^{(2)} = \underline{X}^{(2)}$

\therefore The marginal density of $\underline{Y}^{(2)}$ will be obtained by integrating (\underline{x}) over $\underline{Y}^{(1)}$ also it is to be noted that $f(\underline{Y}^{(1)})$ is a q component MVND.

So the marginal density of $\underline{Y}^{(2)}$ is given by

$$g(\underline{y})^{(2)} = \frac{1}{(2\pi)^{\frac{p-q}{2}} |\underline{\varphi}_{22}|^{1/2}} e^{-\frac{1}{2} (\underline{y}^{(2)} - \underline{\mu}^{(2)})^T \underline{\varphi}_{22}^{-1} (\underline{y}^{(2)} - \underline{\mu}^{(2)})}$$

as $\underline{Y}^{(2)} = \underline{X}^{(2)}$

\therefore the marginal density of $\underline{X}^{(2)}$ will be given by

$$f(\underline{x}^{(2)}) = \frac{1}{(2\pi)^{\frac{p-q}{2}} |\underline{\varphi}_{22}|^{1/2}} e^{-\frac{1}{2} (\underline{x}^{(2)} - \underline{\mu}^{(2)})^T \underline{\varphi}_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)}) / I}$$

(for transformation
consider here $|I|=1$)

i.e $\underline{X}^{(2)} \sim N_{p-q}(\underline{\mu}^{(2)}, \underline{\varphi}_{22})$

(5) Conditional Distribution

Proof Let $\underline{X} \sim N_p(\underline{\mu}, \underline{\varphi})$ and $\underline{X} = \begin{bmatrix} \underline{X}^{(1)} \\ \underline{X}^{(2)} \end{bmatrix}$ accordingly

$$\underline{\mu} = \begin{bmatrix} \underline{\mu}_{11} \\ \underline{\mu}_{12} \end{bmatrix}, \quad \underline{\varphi} = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ - & - \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix}$$

We wish to obtain the conditional density of $\underline{X}^{(1)}$ given $\underline{X}^{(2)} = \underline{x}^{(2)} \rightarrow$ specified vector

Consider the transformation : $\underline{y}^{(1)} = \underline{x}^{(1)} + M \underline{x}^{(2)}$
 $\underline{y}^{(2)} = \underline{x}^{(2)}$

where M is chosen such that $\underline{y}^{(1)}$ and $\underline{y}^{(2)}$ are independent (we know that)

$$M = -\Sigma_{12} \Sigma_{22}^{-1}, \text{ Also,}$$

$$\underline{Y} = \begin{bmatrix} \underline{y}^{(1)} \\ -\underline{y}^{(2)} \end{bmatrix} \sim N_p \left[\begin{bmatrix} \underline{y}^{(1)} \\ \underline{y}^{(2)} \end{bmatrix} \middle| \begin{pmatrix} \Sigma_{11,2} & 0 \\ 0 & \Sigma_{22} \end{pmatrix} \right]$$

$$\Rightarrow \underline{y}^{(1)} \sim N_q \left[\underline{\mu}^{(1)} - \Sigma_{12} \Sigma_{22}^{-1}, \Sigma_{11,2} \right]$$

and

$$\Rightarrow \underline{y}^{(2)} \sim N_{p-q} \left[\underline{\mu}^{(2)}, \Sigma_{22} \right]$$

Now the conditional density of $\underline{x}^{(1)} | \underline{x}^{(2)} = \underline{x}^{(2)}$
 is (say)

$$h \left[\underline{x}^{(1)} \mid \underline{x}^{(2)} = \underline{x}^{(2)} \right] = \frac{\phi(\underline{x}^{(1)}, \underline{x}^{(2)})}{g(\underline{x}^{(2)} = \underline{x}^{(2)})}$$

↑ ↑
R.V Value of
R.V

We find the joint density $\phi(\underline{x}^{(1)}, \underline{x}^{(2)})$
 from it by $|J|$ here

$$|J| = \begin{vmatrix} I & -\Sigma_{12} \Sigma_{22}^{-1} \\ 0 & I \end{vmatrix} = 1$$

$$\therefore \phi(\underline{x}^{(1)}, \underline{x}^{(2)}) = \frac{1}{(2\pi)^{q/2} |\Sigma_{11,2}|}$$

$$e^{-\frac{1}{2} [\underline{x}^{(1)} - \Sigma_{12} \Sigma_{22}^{-1} \underline{x}^{(2)} - \underline{\mu}^{(1)} + \Sigma_{12} \Sigma_{22}^{-1} \underline{\mu}^{(2)}]^T \Sigma_{11,2}^{-1}}$$

$$[\underline{x}^{(1)} - \Sigma_{11} \Sigma_{22}^{-1} \underline{x}^{(2)} - \underline{\mu}^{(1)} + \Sigma_{12} \Sigma_{22}^{-1} \underline{\mu}^{(2)}]$$

$$\times \frac{1}{(2\pi)^{q/2} |\Sigma_{22}|^{1/2}} e^{-\frac{1}{2} (\underline{x}^{(2)} - \underline{\mu}^{(2)})^T \Sigma_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})}$$

→ * *

$\therefore h = [\underline{x}^{(1)} \mid \underline{x}^{(2)} = \underline{x}^{(2)}] = \text{the first part}$
of $(\ast \ast)$

or $h[\underline{x}^{(1)} \mid \underline{x}^{(2)}] \sim N_q [\underline{\mu}^{(1)} + \Sigma_{12} \Sigma_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})]$
 $\Sigma_{11+2}]$

where $\Sigma_{11+2} = \Sigma_{11} - \Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21}$.

So the conditional mean will be

$E[\underline{x}^{(1)} \mid \underline{x}^{(2)} = \underline{x}^{(2)}] = \underline{\mu}^{(1)} + \Sigma_{12} \Sigma_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})$
is also known as Regression fⁿ.

The matrix $\Sigma_{12} \Sigma_{22}^{-1}$ is called matrix of Reg.
coeff and is denoted by \bar{P}

$$\Sigma_{12} \Sigma_{22}^{-1} = \bar{P}$$

$$\underline{x}^{(1)} = \underline{\mu}^{(1)} + \Sigma_{12} \Sigma_{22}^{-1} (\underline{x}^{(2)} - \underline{\mu}^{(2)})$$

$$(\underline{x})^{(1)} = \bar{x} + bxy (y - \bar{y})$$

Note :- Regression is nothing but it is
conditional mean where we find conditional
distⁿ \rightarrow To get conditional mean.

SINGULAR DISTRIBUTION

Given that $\underline{x} \sim N_p(\underline{\mu}, \underline{\delta})$ where $\underline{\delta} p \times p \rightarrow NSM$.
then the random vector \underline{x} is said to have
a Non-Singular or Proper distribution,
on the other hand if the matrix $\underline{\delta}$ is
a singular matrix i.e. rank of $\underline{\delta}$ is 'r' (say)

where $r < p$ then the random vector x is said to have a singular / degenerated dist' of rank r .

→ In this case Σ^{-1} does not exists but it is possible to find 'g-inverse' of Σ and if $\lambda_1, \lambda_2, \dots, \lambda_k$ are the k eigen values of Σ in such a case the density $f(x)$ of x is given by

$$f(x) = \frac{1}{(2\pi)^{k/2} (\lambda_1, \dots, \lambda_k)} e^{-\frac{1}{2} \underline{x}^T (\underline{x} - \underline{\mu})^T (\lambda_1, \lambda_2, \dots, \lambda_k)^{-1} (\underline{x} - \underline{\mu})}$$

Thm: Let $x \sim N_p(\underline{\mu}, \Sigma)$ consider the transformation $z = Dx$ where 'D' is a $t \times p$ matrix ($t \leq p$) of rank t then $z \sim N_t(D\underline{\mu}, D\Sigma D^T)$ (Marginal density)

Proof: If $t < p$ we can find a matrix 'H' s.t $C_{p \times p}^{\frac{1}{2}} = \begin{bmatrix} D \\ H \end{bmatrix}$ is a NSM consider the transformation

$$z^* = C_{p \times p}^{\frac{1}{2}} x = \begin{bmatrix} D \\ H \end{bmatrix} x = \begin{bmatrix} Dx \\ Hx \end{bmatrix} = \begin{bmatrix} z \\ w \end{bmatrix} \text{ (say)}$$

Since z^* is a Non-Singular transformation

$$\therefore z^* \sim N_p(C\underline{\nu}, C\Sigma C^T) \text{ i.e.}$$

$$z^* = \begin{pmatrix} z \\ w \end{pmatrix} \sim N_p \left[\begin{pmatrix} D\underline{\mu} \\ H\underline{\mu} \end{pmatrix}, \begin{pmatrix} D\Sigma D^T & D\Sigma H^T \\ H\Sigma D^T & H\Sigma H^T \end{pmatrix} \right]$$

∴ the marginal density of z will be $N_t(D\underline{\mu}, D\Sigma D^T)$

Thm A random vector \underline{X}_{px1} with $E(\underline{X}) = \underline{\mu}$ and $\text{Var}(\underline{X}) = \Sigma$ is said to have a MVNID with mean vector $\underline{\mu}$ and Variance-Covariance Matrix Σ iff there is a linear transformation $\underline{X} = A\underline{Y} + \underline{\lambda}$ (1), where A is a $p \times p$ matrix \underline{X}_{px1} is of rank ' r ' and $\underline{Y}_{rx1} \sim N_r(\underline{\gamma}, \Sigma)$

NECESSARY PART

Let (1) holds i.e. $\underline{X} = A\underline{Y} + \underline{\lambda}$ } as $\underline{Y} \sim N_r(\underline{\gamma}, \Sigma)$
 $\Rightarrow E(\underline{X}) = A E(\underline{Y}) + \underline{\lambda} = A\underline{\gamma} + \underline{\lambda}$

and

$$E[\underline{X} - E(\underline{X})][\underline{X} - E(\underline{X})'] = [A(\underline{Y} - \underline{\gamma})][(Y - \underline{\gamma})' A]$$

$A E(Y - \underline{\gamma})(Y - \underline{\gamma})' A' = A A' = \Sigma$ (say). If $p > r$, then Σ is a singular matrix. Since, \underline{X} is a linear combination of \underline{Y} , so \underline{X} also has a MVNID, now as the rank of Σ is r which is less than ' p ' so the 'dist' of \underline{X} will be singular as Σ^{-1} does not exist, so we can not write the expression of p.d.f of \underline{X} (Density free) approach.

SUFFICIENCY PART

Let \underline{X} has a singular MVNID of rank ' r ' then there exists a NSM B s.t $B \Sigma B' = \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix}$

now consider the transformation $B\underline{X} = \underline{V}$ (say) and consider the partitioning $\underline{V} = \begin{bmatrix} V^{(1)} \\ V^{(2)} \end{bmatrix}$

So $\underline{V} = B \underline{X}$ defines it having the mean vector $E(\underline{V}) = B E(\underline{X}) = B \underline{\mu}$

Or,

$$E(\underline{V}) = \begin{pmatrix} \underline{V}^{(1)} \\ \underline{V}^{(2)} \end{pmatrix} = \underline{\gamma} \text{ (say) and } B \neq B^T = \begin{pmatrix} I_2 & 0 \\ 0 & 0 \end{pmatrix}$$

\Rightarrow VCM of $V^{(1)} = I_2$ and VCM of $V^{(2)} = 0$

But VCM of $V^{(2)} = 0 \Rightarrow V^{(2)} = \underline{\gamma}^{(2)}$

$$\therefore B \underline{X} = \begin{bmatrix} \underline{V}^{(1)} \\ \underline{V}^{(2)} \end{bmatrix} = \begin{bmatrix} \underline{V}^{(1)} \\ \underline{\gamma}^{(2)} \end{bmatrix} \Rightarrow \underline{X} = B^{-1} \begin{bmatrix} \underline{V}^{(1)} \\ \underline{\gamma}^{(2)} \end{bmatrix}$$

Now partition B^{-1} as $[C \ D]$

so that $\underline{X} = C \underline{V}^{(1)} + D \underline{\gamma}^{(2)} = C \underline{V}^{(1)} + \underline{\lambda} \text{ (say)}$

Or $\underline{X} = A \underline{Y} + \underline{\lambda}$ from

\therefore If \underline{X} has a singular MVND it can be expressed as above.

07/09/2022

Result :- A p-component vector \underline{U} is said to have a MVND iff every linear combination of \underline{U} has a univariate normal dist'

* NECESSARY PART

Consider an arbitrary vector $\underline{L} p_{xi}$ and let \underline{U}_{pxi} has the MVND then since \underline{L} is assumed to be NON - NULL, \underline{L}' will be of rank 1 so that there exists a matrix B such that $B = (\underline{U}')$ and B' is a NSM consider the transform D) transformation : $Z = B \underline{U}$ or $Z = \begin{pmatrix} L' U \\ D U \end{pmatrix}$

Also,

$$Z \sim N \left[\begin{pmatrix} L' U \\ D U \end{pmatrix}, \begin{pmatrix} L' \Sigma L & L' \Sigma D \\ D \Sigma L & D \Sigma D' \end{pmatrix} \right]$$

$$\Leftrightarrow \underline{L'U} \sim N_1 (\underline{L'\mu}, \underline{L'L})$$

$\underline{\underline{L'x_1}}$ $\underline{(L'x_1)}$ $\underline{(x_1)}$

i.e every linear combination of \underline{U} is a univariate normal.

SUFFICIENT PART

let $\underline{L'x_1}$ be any arbitrary vector and \underline{U} has a univariate normal dist with mean $\underline{L'\mu}$ and variance $\underline{L'L}$ then the characteristic fn of \underline{U} will be

$$\phi_t(\underline{L'U}) = E[e^{it\underline{L'U}}] = e^{it\underline{L'\mu} - \frac{1}{2}t^2\underline{L'L}}$$

now iff $t = 1$ then

$$\phi_1(\underline{L'U}) = e^{i\underline{L'\mu} - \frac{1}{2}\underline{L'L}} \quad (\text{By defn})$$

i.e $\phi_1(\underline{L'U}) = E[e^{i\underline{L'U}}]$ which is the characteristic fn

since \underline{L} is any arbitrary vector so

$\phi(\underline{L'U}) = E[e^{i\underline{L'U}}]$ which is the ch.f of a MVND

$\therefore \underline{U}$ has a MVND using the uniqueness property of ch.f

R

REPRODUCTIVE PROPERTY

Statement:- Let the i^{th} random vector $\underline{U}_i \sim N_p(\underline{\mu}_i, \underline{\Sigma}_i)$ $i=1, 2, \dots, n$ then for fixed constants A_1, A_2, \dots, A_n Consider $Y = \sum_{i=1}^n A_i \underline{U}_i$; Then

$$Y \sim N_p \left[\sum_i A_i \underline{\mu}_i, \sum_{i=1}^n A_i^2 \underline{\Sigma}_i \right]$$

Proof

We have $Y = A_1 \underline{U}_1 + \dots + A_n \underline{U}_n$. Consider an arbitrary linear combination $L' Y = A_1 L' \underline{U}_1 + \dots + A_n L' \underline{U}_n$. Since $\underline{U}_i \sim N_p(\underline{\mu}_i, \underline{\Sigma}_i)$ $L' \underline{U}_i$ is a univariate Normal. Again

$$\underline{L}' Y = \sum_{i=1}^n \underline{L}' A_i \underline{U}_i \quad \text{or.}$$

$$\underline{L}' Y = \sum_{i=1}^n A_i Z_i \quad \text{where } Z_i = L' \underline{U}_i$$

Since it is a linear combination of univariate normal variable, so that $\underline{L}' Y$ also has a univariate normal with MEAN = $E[\underline{L}' Y]$

OR

$$\begin{aligned} \text{MEAN} &= E[A_1 L' \underline{U}_1 + \dots + A_n L' \underline{U}_n] \\ &= L' [A_1 \underline{U}_1 + \dots + A_n \underline{U}_n] \\ &= L' \sum_{i=1}^n A_i \underline{U}_i \end{aligned}$$

And

$$\begin{aligned} \text{Variance} &= \text{Var}[A_1 L' \underline{U}_1 + \dots + A_n L' \underline{U}_n] \\ &= A_1^2 \text{Var}(L' \underline{U}_1) + \dots + A_n^2 \text{Var}(L' \underline{U}_n) \\ &= A_1^2 L' \underline{\Sigma}_1 L + \dots + A_n^2 L' \underline{\Sigma}_n L \\ &= L' \sum_{i=1}^n A_i^2 \underline{\Sigma}_i L \quad \text{i.e.} \end{aligned}$$

$$\text{i.e. } L' Y \sim N_p \left(\sum_{i=1}^n A_i \underline{\mu}_i, \sum_{i=1}^n A_i^2 \underline{\Sigma}_i \right)$$

Since L is any arbitrary vector and $L^T Y$ vector combination and $L^T Y$ is a linear combination of Y which has a MVNID

$$\therefore A Y \sim N_p \left(\sum_{i=1}^m A_i \mu_i, \sum_{i=1}^m A_i \Sigma_i \right)$$

* MLE's of the parameters of $N_p(\underline{\mu}, \underline{\Sigma})$:

Rules of differentiation:-

$$(1) \frac{\partial |X|}{\partial X} = |X| X^{-1}$$

$$(2) \frac{\partial \text{tr} \cdot AX}{\partial X} = A \quad [\text{tr} - \text{trace}]$$

$$(3) \frac{\partial \underline{x}' A \underline{x}}{\partial \underline{x}} = 2 \underline{x}' A = 2 A' \underline{x}$$

$\underline{x} \sim N_p(\underline{\mu}, \underline{\Sigma})$ where $\underline{\mu}$ is mean vector and $\underline{\Sigma}$ is var. cov Matrix of the Random Vector X

If $\underline{\mu}$ and $\underline{\Sigma}$ are known then the distribution of X is completely specified

However if they are not known then for specifying the distribution of X we have to estimate $\underline{\mu}$ and $\underline{\Sigma}$ with the help of a random sample of vectors from $N_p(\underline{\mu}, \underline{\Sigma})$

In the estimation procedure we will assume that $\underline{\Sigma}$ is a non-singular matrix

Derivation :-

Let $\underline{x}_1, \dots, \underline{x}_N$ be a set of N random vectors from $N_p(\underline{\mu}, \underline{\Sigma})$.

Define:- The sample mean vector

$$\bar{\underline{x}} = \frac{1}{N} \sum_{\alpha=1}^N \underline{x}_{\alpha} = \begin{pmatrix} \bar{x}_1 \\ \bar{x}_2 \\ \vdots \\ \bar{x}_p \end{pmatrix} = \frac{1}{N} \begin{pmatrix} \sum_{\alpha} \bar{x}_{1\alpha} \\ \sum_{\alpha} \bar{x}_{2\alpha} \\ \vdots \\ \sum_{\alpha} \bar{x}_{p\alpha} \end{pmatrix}$$

and let the sample VCM be,

$$\underline{S} = \frac{1}{N} \sum_{\alpha=1}^N (\bar{\underline{x}}_{\alpha} - \bar{\underline{x}})(\bar{\underline{x}}_{\alpha} - \bar{\underline{x}})^T \rightarrow (p \times p) \text{ matrix}$$

The likelihood function (LF) is given by

$$L = \prod_{i=1}^N f(\underline{x}_i, \underline{\mu}, \underline{\Sigma}) = e^{-\frac{1}{2} \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{\mu})^T \underline{\Sigma}^{-1} (\underline{x}_{\alpha} - \underline{\mu})} (2\pi)^{\frac{Np}{2}} |\underline{\Sigma}|^{N/2}$$

Since, $\underline{\Sigma}$ is NSM, so $\underline{\Sigma}^{-1}$ will exist.

∴ we find the MLE of $\underline{\Sigma}^{-1}$ and then using the invariance property of MLE we obtain the MLE of $\underline{\Sigma}$.

Now

$$\log L = k - \frac{N}{2} \log |\underline{\Sigma}| - \frac{1}{2} \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{\mu})^T \underline{\Sigma}^{-1} (\underline{x}_{\alpha} - \underline{\mu})$$

$$\log L = k + \frac{N}{2} \log |\underline{\Sigma}^{-1}| - \frac{1}{2} \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{\mu})^T \underline{\Sigma}^{-1} (\underline{x}_{\alpha} - \underline{\mu})$$

Consider the exponent

$$\sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{\mu})^T \underline{\Sigma}^{-1} (\underline{x}_{\alpha} - \underline{\mu})$$

$$\sum_{\alpha=1}^N (\underline{x}_{\alpha} - \bar{\underline{x}} + \bar{\underline{x}} - \underline{\mu})^T \underline{\Sigma}^{-1} (\underline{x}_{\alpha} - \bar{\underline{x}} + \bar{\underline{x}} - \underline{\mu})$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) + \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

$$+ \sum_{\alpha=1}^N (\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) + \sum_{\alpha=1}^N (\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) + \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) + N(\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

Again,

$$\sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \underline{y})$$

$$= \text{tr} \left[\sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) \right] + N(\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

$$= \sum_{\alpha=1}^N \text{tr} \underline{\Sigma}^{-1} (\underline{x}_\alpha - \bar{\underline{x}}) (\underline{x}_\alpha - \bar{\underline{x}})^T + N(\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

Let $A = \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}}) (\underline{x}_\alpha - \bar{\underline{x}})^T$ then,

$$\sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y})^T \underline{\Sigma}^{-1} (\underline{x}_\alpha - \underline{y})$$

$$= \text{tr}(A \underline{\Sigma}^{-1}) + N(\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

$$\log L = k + \frac{N}{2} \log |\underline{\Sigma}^{-1}| - \frac{1}{2} \left[\text{tr}(\underline{\Sigma}^{-1} A) + N(\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y}) \right]$$

$$= k + \frac{N}{2} \log |\underline{\Sigma}^{-1}| - \frac{1}{2} \text{tr} \cdot \underline{\Sigma}^{-1} A - \frac{N}{2} (\bar{\underline{x}} - \underline{y})^T \underline{\Sigma}^{-1} (\bar{\underline{x}} - \underline{y})$$

— (*)

Differentiating (*) w.r.t \underline{u} and $\mathbf{\Sigma}^{-1}$ respectively we obtain:-

$$\frac{\partial \log L}{\partial \underline{u}} = \frac{N}{2} \mathbf{\Sigma}^{-1} (\bar{x} - \underline{u}) \geq 0 \quad \text{--- (1)}$$

$$\Rightarrow (\bar{x} - \hat{\underline{u}}) = 0 \Rightarrow \boxed{\hat{\underline{u}} = \bar{x}}$$

$$\text{Also, } \frac{\partial \log L}{\partial \mathbf{\Sigma}^{-1}} = N \left[\frac{1}{2} \mathbf{\Sigma}^{-1} \mathbf{\Sigma}^{-1} - \frac{A}{2} - \frac{N}{2} (\bar{x} - \underline{u})(\bar{x} - \underline{u})' \right] = 0$$

$$\Rightarrow \mathbf{\Sigma}^{-1} = (A)^{-1} = N A^{-1}$$

$$\therefore \hat{\Sigma} = \frac{1}{N} \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \bar{x})(\underline{x}_{\alpha} - \bar{x})'$$

Note It has been shown that,

$$\sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{u})(\underline{x}_{\alpha} - \underline{u})' = A + N(\bar{x} - \underline{u})(\bar{x} - \underline{u})'$$

$$\text{and } \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{u})' \mathbf{\Sigma}^{-1} (\underline{x}_{\alpha} - \underline{u})$$

$$= \text{tr}(\Sigma^{-1} A) + N(\bar{x} - \underline{u})' \Sigma^{-1} (\bar{x} - \underline{u})$$

Thus the joint p.d.f. of x_1, \dots, x_N can be written as:

$$K e^{-\frac{1}{2} N (\bar{x} - \underline{u})' \mathbf{\Sigma}^{-1} (\bar{x} - \underline{u})} \cdot \frac{1}{2} e^{-\frac{1}{2} \text{tr} \mathbf{\Sigma}^{-1} A}$$

Thus \bar{x} and A form a set of jointly sufficient statistic for \underline{u} and $\mathbf{\Sigma}$. If $\mathbf{\Sigma}$ is known then \bar{x} is sufficient for \underline{u} however if \underline{u} is known then A/N is not sufficient for $\mathbf{\Sigma}$ but $\frac{1}{N} \sum_{\alpha=1}^N (\underline{x}_{\alpha} - \underline{u})(\underline{x}_{\alpha} - \underline{u})'$ is sufficient for $\mathbf{\Sigma}$.

* Elliptical Distribution

Elliptical distⁿs provide a generalized form for a multivariate normal distⁿ. In the 2 and 3 dimension cases the joint distⁿ function of the variables for ellips or ellipsoids.

It is to be noted that the multivariate normal distⁿ is constant on ellipsoids i.e.

$$(\underline{x} - \underline{y})' \underline{\Sigma}^{-1} (\underline{x} - \underline{y}) = c \quad \text{--- (1)}$$

For every positive value of 'c' in p-dimensional space, the center of each ellipsoid is at the point \underline{y} . The shape and orientation is determined by $\underline{\Sigma}$.

because (1) is sphere if $\underline{\Sigma} = \sigma^2 I$ then $N_p(\underline{y}, \sigma^2 I)$ is known as spherical normal density.

It is noted that MVND with mean \underline{y} and VCM $\underline{\Sigma}$ is constant on concentric ellipsoids

A general class of distⁿ with this property is the class of elliptically contoured distⁿ with the density.

$$|\Lambda|^{-\frac{1}{2}} g[(\underline{x} - \underline{\delta})' \Lambda (\underline{x} - \underline{\delta})] \text{ where } \Lambda \text{ is a tve definite matrix}$$

If C is a NSM s.t. $C' \Lambda^{-1} C = I$ then the transformation $(\underline{x} - \underline{\delta}) = C \underline{y}$ carries the general density to Multivariate normal density then the contours of constants

density these are the spheres centered at origin. the class of such density is known as spherically centered dist?

Applications :-

Elliptical dist? were important in the study of economics and financials situations where various factors of portfolio analysis including mutual fan separation theorem and capital assets pricing model hold good for elliptical dist?

UNIT-II

Wishart Distribution (1928)

Another name is multivariate χ^2 or generalization of χ^2

The wishart "dist" is the "dist" of the sample variance as w variance matrix it is also called as "dist" of the moment matrix

Wishart "dist" is a multivariate generalization of χ^2 -dist

The Sample VCM is given by

$$S = \frac{1}{(N-1)} \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})(\underline{x}_\alpha - \bar{\underline{x}})'$$

OR

$$(N-1)S = \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{\underline{x}})(\bar{\underline{x}}_\alpha - \bar{\underline{x}})$$

Let $z_\alpha = (\underline{x}_\alpha - \bar{\underline{x}})$ then:

$z_\alpha \sim N_p(0, \Sigma)$ Now the p.d.f of z_α is given by

$$\frac{1}{(2\pi)^{p/2} |\Sigma|^{1/2}} e^{-\frac{1}{2} z_\alpha' \Sigma^{-1} z_\alpha}$$

∴ the p.d.f of z_1, \dots, z_N (say) $f(z)$ is given by

$$f(z) = \frac{1}{(2\pi)^{(N-p)/2} |z|^{N/2}}$$

$$e^{-\frac{1}{2} \sum_{\alpha=1}^N z_\alpha' \Sigma^{-1} z_\alpha}$$

Exponent is $-\frac{1}{2} \sum_{\alpha=1}^N z_\alpha' \Sigma^{-1} z_\alpha$

which is a scalar quantity & since $\sum_{\alpha=1}^N \mathbf{z}_{\alpha}' \mathbf{S}^{-1} \mathbf{z}_{\alpha}$ is scalar.

$$\therefore \sum_{\alpha=1}^{N-1} \mathbf{z}_{\alpha}' \mathbf{S}^{-1} \mathbf{z}_{\alpha} = \text{tr} \sum_{\alpha=1}^{N-1} \mathbf{z}_{\alpha}' \mathbf{S}^{-1} \mathbf{z}_{\alpha}$$

$$= \sum_{\alpha=1}^{N-1} \text{tr} \cdot \mathbf{z}_{\alpha}' \mathbf{S}^{-1} \mathbf{z}_{\alpha} = \sum_{\alpha=1}^{N-1} \text{tr} \mathbf{S}^{-1} \mathbf{z}_{\alpha} \mathbf{z}_{\alpha}'$$

$$= 2 \text{tr} \mathbf{S}^{-1} \sum_{\alpha=1}^{N-1} \mathbf{z}_{\alpha} \mathbf{z}_{\alpha}'$$

$$= \text{tr} \cdot (\mathbf{S}^{-1} \mathbf{A}) \text{ as } \mathbf{A} = \sum_{\alpha} (\mathbf{x}_{\alpha} - \bar{\mathbf{x}})(\mathbf{x}_{\alpha} - \bar{\mathbf{x}})'$$

$$f(\underline{z}) = \frac{1}{(2\pi)^{(N-1)p/2} |\mathbf{S}|^{(N-1)/2}} e^{-\frac{1}{2} \text{tr} \cdot \mathbf{S}^{-1} \mathbf{A}}$$

$$\text{As } f(\underline{z}) \text{ is p.d.f } \int f(\underline{z}) d\underline{z} = 1$$

$$\text{i.e. } \int e^{-\frac{1}{2} \text{tr} \cdot \mathbf{S}^{-1} \mathbf{A}} = (2\pi)^{(N-1)p/2} |\mathbf{S}|^{(N-1)/2}$$

— (1)

Since the Matrix \mathbf{A} is a symmetric matrix
the no of distinct elements of the
matrix will be

$\frac{P(P+1)}{2}$

and let these are $\begin{bmatrix} a_{11} & a_{12} & \dots & a_{1p} \\ a_{21} & a_{22} & \dots & a_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ \vdots & \vdots & \ddots & a_{pp} \end{bmatrix}$

Now the ch.f of \mathbf{A} is given by

$\Phi(\mathbf{A}) = E [e^{\text{tr} \cdot i \mathbf{U} \mathbf{A}}]$ where \mathbf{U} is a matrix such that

$$U_{ij} = U_{ji} = \frac{t_{ij}}{2}$$

$$\text{i.e. } U = \begin{bmatrix} t_{11} & t_{12}/2 & \dots & t_{1p}/2 \\ t_{21}/2 & t_{22} & \dots & t_{2p}/2 \\ \vdots & \vdots & \ddots & \vdots \\ t_{p1}/2 & t_{p2}/2 & \dots & t_{pp} \end{bmatrix} \quad \text{and} \quad \begin{bmatrix} a_{11} & a_{12} & \dots & a_{1p} \\ a_{21} & a_{22} & \dots & a_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ a_{p1} & a_{p2} & \dots & a_{pp} \end{bmatrix}$$

$$\therefore \text{tr}(UA) = t_{11}a_{11} + t_{12}a_{12} + \dots + t_{1p}a_{1p} \\ t_{21}a_{21} + \dots + \dots + t_{2p}a_{2p} \\ + \dots + \dots + \dots + \dots + \dots + \\ + \dots + \dots + \dots + t_{pp}a_{pp} \\ - (\ast \ast)$$

$$\therefore \Phi_A(U) = E \left[\exp \{ i(t_{11}a_{11} + t_{12}a_{12} + \dots + t_{pp}a_{pp}) \} \right]$$

Again

$$\Phi_A(U) = E \left[e^{\text{tr}.iUA} \right] = \int e^{\text{tr}.iUA} f(z) dz$$

$$= \frac{1}{(2\pi)^{\frac{(N-1)p}{2}}} \int_{\mathbb{C}} e^{\text{tr}.iUA} \cdot e^{-\frac{1}{2}\text{tr} \cdot z^{-1}A} dz$$

$$= \frac{1}{(2\pi)^{\frac{(N-1)p}{2}}} \int_{\mathbb{C}} e^{-\frac{1}{2}\{ \text{tr} \cdot (z^{-1} - 2iU) A \}} dz \quad (2)$$

But from the analogy of (1) we have

$$\int e^{-\frac{1}{2}\text{tr} \cdot (z^{-1} - 2iU) A} dz = |z^{-1} - 2iU|^{\frac{-(N-1)}{2}} (2\pi)^{\frac{(N-1)p}{2}}$$

$$\Phi_A(U) = \frac{|z^{-1} - 2iU|^{\frac{-(N-1)}{2}}}{(2\pi)^{\frac{(N-1)p}{2}}} \times (2\pi)^{\frac{(N-1)p}{2}}$$

using $(*)$

$$\Phi_A(v) = \frac{1}{|I - 2iv|^{\frac{(N-1)}{2}}} |S|^{\frac{(N-1)}{2}}$$

OR

$$\Phi_A(v) = \frac{1}{|I - 2iv|^{\frac{(N-1)}{2}}}$$

Now suppose that the pdf of A is given by

$$P(A) = C |S|^{\frac{-1}{2}} |A|^{\frac{(N-p-2)}{2}} \exp\left\{-\frac{1}{2} \operatorname{tr} S^{-1} A^2\right\}$$

where

$$C = \frac{1}{(2\pi)^{\frac{(N-p)}{2}} \frac{P}{\prod_{j=1}^p \sqrt{\frac{N-j}{2}}}}$$

Again $\int P(A) dA = 1$ gives

$$\int |A|^{\frac{(N-p-2)}{2}} C |S|^{\frac{-1}{2}} \exp\left\{-\frac{1}{2} \operatorname{tr} S^{-1} A^2\right\} dA = 1$$

$$\Rightarrow \int |A|^{\frac{(N-p-2)}{2}} \exp\left\{-\frac{1}{2} \operatorname{tr} S^{-1} A^2\right\} dA = |S|^{\frac{(N-1)}{2}}$$

and $\Phi_A(v) = E[e^{trivA}]$

$$= C |S|^{\frac{-1}{2}} \int |A|^{\frac{(N-p-2)}{2}} e^{-\frac{1}{2} \operatorname{tr}(S^{-1} - 2iv) A} dA,$$

$$= C |S|^{\frac{-1}{2}} |S^{-1} - 2iv|^{\frac{-(N-1)}{2}}$$

$$= |I - 2iv|^{\frac{-(N-1)}{2}} = |I - 2iv|^{\frac{-(N-1)}{2}}$$

(4)

Hence, it is observed that the ch.fn 'A' is same as that given by (4) which is obtained by using the proposed density $P(A)$. Using the uniqueness theorem ch.fn we conclude that the pdf of A is the proposed one given by (4). This distⁿ of 'A' is called the Wishart distⁿ and is denoted by $W_p(n, \mathbf{S})$ where $n = (N-1)$

Remark :- Wishart distⁿ is a Multivariate analogue of χ^2 density i.e. for $p=1$ we obtain χ^2_n from $W_p(n, \mathbf{S})$

Proof for $p=1$, $e^{-\frac{1}{2}x_1} = \sigma^2$ (say)

$$A = \sum_{k=1}^N (x_k - \bar{x})^2 = (N-1)s^2 = n s^2$$

$$\text{For } p=1, C = \frac{1}{2^{(N-1)/2} \pi^{N/2} \left(\frac{N-1}{2}\right)^{(N-1)/2}} \cdot \frac{2^{n/2}}{\sqrt{n}}$$

$$\begin{aligned} P(A) &= C |\sigma^2|^{-\frac{N-1}{2}} (n s^2)^{\frac{N-1-2}{2}} \exp \left(-\frac{1}{2} \text{tr} \cdot \mathbf{S}^{-1} \mathbf{A} \mathbf{S} \right) \\ &= C \cdot \frac{1}{\sigma^2} \left(\frac{n s^2}{\sigma^2} \right)^{\frac{n-1}{2}-1} e^{-\frac{n s^2}{2 \sigma^2}} \end{aligned}$$

OR

$$P(A) = \frac{1}{2^{n/2} \sqrt{\frac{n}{2}} \sigma^2} \left(\frac{n s^2}{\sigma^2} \right)^{\frac{n-1}{2}} e^{-\frac{n s^2}{2 \sigma^2}} dA$$

$$\text{let } A = n s^2 \Rightarrow dA = n d s^2$$

$$\text{putting } \frac{n s^2}{\sigma^2} = u \Rightarrow \frac{n d s^2}{\sigma^2} = du$$

$$P(V) dV = \frac{1}{2^{n/2} \sqrt{\pi}} V^{\frac{n}{2}-1} e^{-\frac{v}{2}} dV \quad \left| \begin{array}{l} \\ \\ \end{array} \right. \chi_m^2$$

Properties of Wishart distribution

(1) The chfn of A (Wishart distn) is

$$\phi_A = |I - 2V|^{-\frac{(N-1)}{2}}$$

(2) Reproductive Property

If $A_i \sim W_p(n_i, \mathbf{I})$, $i=1, 2$. then
 $A_1 + A_2 \sim W_p(n_1 + n_2, \mathbf{I})$

Proof By defⁿ $\phi_A(v) = E[e^{tr.vY}]$

$$\begin{aligned} &= E[e^{tr.v(A_1 + A_2)}] \\ &= E[e^{tr.vA_1}] \cdot E[e^{tr.vA_2}] \\ &= \phi_{A_1}(v) \cdot \phi_{A_2}(v) \end{aligned}$$

$$\begin{aligned} &= |I - 2v\$|^{-\frac{n_1}{2}} \cdot |I - 2v\$|^{-\frac{n_2}{2}} \\ &= |I - 2v\$|^{-\frac{(n_1+n_2)}{2}} \end{aligned}$$

$$\Rightarrow Y \sim W_p(n_1 + n_2, \mathbf{I})$$

In general, if A_i ($i=1, \dots, k$) are independently distributed as $W_p(n_i, \mathbf{I})$ then

$$\sum_{i=1}^k A_i \sim W_p\left(\sum_{i=1}^k n_i, \mathbf{I}\right)$$

(3) If $A \sim W_p(n, \mathbf{\Sigma})$ consider CA , where C is a scalar constant then $CA \sim W_p(n, C\mathbf{\Sigma})$

proof let $Y = CA$ then $\Phi_Y(u) = E[e^{t^T u Y}]$
 $= E[e^{t^T u CA}]$,

let $U^* = UC$ so $\Phi_Y(u) = E[e^{t^T u^* A}]$
 $= \Phi_A(u^*)$

$$\text{or } \Phi_Y(u) = |I - 2iuC\mathbf{\Sigma}|^{-n/2}$$

$$\Rightarrow Y \sim W_p(n, C\mathbf{\Sigma})$$

(4) Marginal Distribution :-

Let A and $\mathbf{\Sigma}$ be partitioned as :

$$A = \begin{bmatrix} A_{11} & A_{12} \\ A_{21} & A_{22} \end{bmatrix} \text{ and } \mathbf{\Sigma} = \begin{bmatrix} \mathbf{\Sigma}_{11} & \mathbf{\Sigma}_{12} \\ \mathbf{\Sigma}_{21} & \mathbf{\Sigma}_{22} \end{bmatrix}$$

$$\text{where } A_{11} = q \times q \quad A_{21} = (p-q) \times q$$

$$A_{12} = q \times (p-q) \quad A_{22} = (p-q) \times (p-q)$$

$$\Sigma_{11} = q \times q \quad \Sigma_{21} = (p-q) \times q$$

$$\Sigma_{12} = q \times (p-q) \quad \Sigma_{22} = (p-q) \times (p-q)$$

If $A \sim W_p(n, \mathbf{\Sigma})$ then $A_{11} \sim W_q(n, \mathbf{\Sigma}_{11})$

proof Given that $A \sim W_p(n, \mathbf{\Sigma})$ we also write :

$A = \sum_{\alpha=1}^{N_p} Z_\alpha Z_\alpha'$ where $Z_\alpha \sim N_p(0, \mathbf{\Sigma})$ and are iid.

Consider the partitioning of $Z_\alpha = \begin{bmatrix} Z_\alpha^{(1)} \\ Z_\alpha^{(2)} \end{bmatrix} \rightarrow \begin{cases} q \times 1 \\ (p-q) \times 1 \end{cases}$

such that

$$\underline{z}_\alpha^{(1)} \sim N_q(0, \Sigma_{11})$$

$$\text{Thus, } A = \sum_{k=1}^{n-1} \underline{z}_\alpha \underline{z}_\alpha' = \sum_{\alpha=1}^n \begin{bmatrix} \underline{z}_\alpha^{(1)} \\ \underline{z}_\alpha^{(2)} \end{bmatrix} \begin{bmatrix} \underline{z}_\alpha^{(1)'} & \underline{z}_\alpha^{(2)'} \end{bmatrix}$$

$$\text{Or, } A = \begin{bmatrix} \sum_{\alpha=1}^n \underline{z}_\alpha^{(1)} \underline{z}_\alpha^{(1)'} & \sum_{\alpha=1}^n \underline{z}_\alpha^{(1)} \underline{z}_\alpha^{(2)'} \\ \sum_{\alpha=1}^n \underline{z}_\alpha^{(2)} \underline{z}_\alpha^{(1)'} & \sum_{\alpha=1}^n \underline{z}_\alpha^{(2)} \underline{z}_\alpha^{(2)'} \end{bmatrix} = \begin{bmatrix} A_{11} & A_{12} \\ A_{21} & A_{22} \end{bmatrix}$$

$$\therefore A_{11} = \sum_{\alpha=1}^n \underline{z}_\alpha^{(1)} \underline{z}_\alpha^{(1)'}, \quad \underline{z}_\alpha \sim N_q(0, \Sigma_{11}) \text{ and iid}$$

$$\therefore A_{11} \sim W_q(n, \Sigma_{11})$$

5. If $A \sim W_p(n, \frac{1}{2})$, and L px is any arbitrary
find vector then $L'AL \sim \chi^2_{n \cdot 6^2}$

Proof

Given that $A \sim W_p(n, \frac{1}{2})$ and $A = \sum_{\alpha=1}^n \underline{z}_\alpha \underline{z}_\alpha'$

$\underline{z}_\alpha \sim N_p(0, \frac{1}{2})$ and are iid

$$\therefore L'AL = L' \sum_{\alpha=1}^n \underline{z}_\alpha \underline{z}_\alpha' L = \sum_{\alpha=1}^n (L'\underline{z}_\alpha)(L'\underline{z}_\alpha)'$$

Since $L'\underline{z}_\alpha$ is a linear combination of \underline{z}_α also
it is scalar $\therefore L'AL = \sum_{\alpha=1}^n (L'\underline{z}_\alpha)^2$

Again, $\underline{z}_\alpha \sim N_p(0, \frac{1}{2}) \therefore L'\underline{z}_\alpha \sim N_1(0, L'L)$
or,

$$L'\underline{z}_\alpha \sim N_1(0, \sigma^2) \quad \text{where } L'L = \sigma^2 \text{ (say)}$$

$\therefore L'AL = \sum_{\alpha=1}^n (L'\underline{z}_\alpha)^2$ is the sum of square
of SNV's

$$\therefore \frac{\text{E}[AL]}{\sigma^2} \sim \chi_m^2$$

or,

$$\boxed{\text{E}[AL] \sim \sigma^2 \chi_m^2}$$

6. If $A \sim W_p(n, \mathbf{I})$ then $B^T A B \sim W_q(n, B^T \mathbf{I} B)$
where B is a $(p \times p)$ matrix

Proof: Given that $A \sim W_p(n, \mathbf{I})$ where $A = \sum_{\alpha=1}^n Z_\alpha Z_\alpha^T$,

$$Z_\alpha \sim N_p(0, I)$$

$$\text{now } B^T A B = B^T \sum_{\alpha=1}^n Z_\alpha Z_\alpha^T B = \sum_{\alpha=1}^n B^T Z_\alpha Z_\alpha^T B = \sum_{\alpha=1}^n W_\alpha W_\alpha^T \quad (\text{say})$$

where

$$W_\alpha = B^T Z_\alpha \text{ and are indep. for } \alpha = 1, 2, \dots, n$$

$$\therefore B^T A B \sim W_q(n, B^T \mathbf{I} B)$$

7. CONDITIONAL DISTRIBUTION

Let $A \sim W_p(n, \mathbf{I})$ and A and \mathbf{I} are partitioned

$$\text{as: } A = \begin{bmatrix} A_{11} & A_{12} \\ \bar{A}_{21} & A_{22} \end{bmatrix} \quad \mathbf{I} = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix}$$

then $A_{11|2}$ is distributed as $W_p[n - (p-q), \Sigma_{11|2}]$

$$\text{where } A_{11|2} = A_{11} - A_{12} A_{22}^{-1} A_{21}$$

$$\text{and } \Sigma_{11|2} = \Sigma_{11} - \Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21}$$

*

DISTRIBUTION OF SAMPLE MEAN VECTOR AND SAMPLE VCM :-

Let $x \sim N_p(\mu, \Sigma)$ then L.F of $\bar{x}_1, \bar{x}_2, \dots, \bar{x}_N$ is:

$$L = \frac{1}{(2\pi)^{\frac{Np}{2}} |\Sigma|^{N/2}} e^{-\frac{1}{2} \text{tr. } \Sigma^{-1} A} e^{-\frac{1}{2} (\bar{x} - \mu)' (\bar{x} - \mu)} \quad \text{--- (1)}$$

The MLE's are:

$$\hat{\mu} = \bar{x}, \quad \hat{\Sigma} = A/N \quad \text{or,} \quad \hat{\Sigma} = S! \text{ Sample VCM}$$

Consider the transformation $y = p x$

where p is an orthogonal matrix
and its last row is taken as $\frac{1}{\sqrt{n}}$

in this case $|J|=1$ and consider the partitioning of $y = \begin{bmatrix} z \\ \sqrt{n} \bar{x} \end{bmatrix}$ where $z_{(n-1) \times p}$

We can write $(\bar{x} - \mu)' (\bar{x} - \mu) = (\bar{x} - \mu)' P P' (\bar{x} - \mu)$

$$\text{as } P P' = I$$

$$= \begin{bmatrix} z \\ \sqrt{n} (\bar{x} - \mu) \end{bmatrix}' \begin{bmatrix} P & z \\ \sqrt{n} (\bar{x} - \mu) \end{bmatrix}$$

$$= z z' + N (\bar{x} - \mu)' (\bar{x} - \mu)$$

$$= \frac{1}{(2\pi)^{\frac{Np}{2}} |\Sigma|^{N/2}} e^{-\frac{1}{2} \text{tr. } \Sigma^{-1} [z z' + N (\bar{x} - \mu)' (\bar{x} - \mu)]}$$

which can be partitioned as $z z'$ and $N (\bar{x} - \mu)' (\bar{x} - \mu)$. Again z is a matrix of $(N-1)$ observations from $(p-1)$ variables.
 $\therefore z' z \sim W_{p-1}(n-1, \Sigma)$ as $z' z = Y' Y / N \propto 1$

$\therefore Z'Z \sim \text{W}_{p-1} (m-1, \frac{1}{n})$

$$\begin{aligned} \text{as } Z'Z &= Y'Y - n\bar{S}' \\ &= X'P'P X - n\bar{S}' \\ &= nS = n\hat{\Sigma}, \end{aligned}$$

further $\bar{S} \sim N_p (\underline{0}, \frac{1}{n} I)$

* Problem of classification and discrimination :-

INTRODUCTION:-

In many real life situations sometimes we have a group of individuals or objects or some conditions and we wish to classify them into different categories.

This may be done by taking into account certain observations on a given set of individuals or objects the condition and on the basis of their characteristics we wish to decide to which group a particular individual to which group belongs. This is the case where all the individual / of a objects / condition to be classified.

In the problem of classification we do not know the number of categories into which individuals / objects / conditions are to be classified & here we have to make our own categories according to our satisfaction. However in case where the no. of categories are known and we have to only separate the observation into the known groups it is known as the problem of discrimination.

For the both purposes we measure p-characteristic on each individual / object / condition and thus we have a vector \mathbf{x}_{px} as the observation. Corresponding to k-categories, we can

assume k -underlying pop ns and when we say a vector belongs to a category it is equivalent to saying that the individual / objects / conditions belongs to corresponding pop n .

* AREAS OR APPLICATIONS:

1. Anthropological/ Archaeological studies (when the information is lost)
2. Medical diagnosis of some disease (like - brain cancer etc)
3. Financial Sector
4. In general where the measurement of identification variable is difficult or expensive.

* Mathematical formulation:

Let $\mathbf{x} = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_p \end{pmatrix}$ for the sake of simplicity we consider only two categories into which x_{px} is to be classified

then any random vector x_{px} may come from either the pop n π_1 or π_2

For taking our decision we divide (?) our p -dimensional space into two regions R_1 and R_2 .

and the decision is made as: if the observed value of $x \in R_1$, we say that the individual belongs to the Pop². Π_1 , otherwise it belongs to Π_2 . R_1 and R_2 must be mutually exclusive.

Decision Taken	Actual Population	
	Π_1	Π_2
Π_1	No Error	Second type of misclassification (m.c)
Π_2	First type of m.c	No Error

In any problem while choosing R_1 and R_2 , there is same associated loss, so our purpose is to choose R_1 and R_2 in such a way that the average loss or cost is minimum.

For this we define:

$C(2|1)$: Cost of 1st type of m.c i.e. cost of misclassifying an order as coming from Π_2 when actually it belongs to Π_1 .

$C(1|2)$: Cost of 2nd type of m.c i.e. cost of misclassifying an order as coming from Π_1 when actually it belongs to Π_2 .

We also assume that $C(2|1) \geq 0$, $C(1|2) \geq 0$ and $C(1|1) = 0$ & $C(2|2) = 0$.

Let $f_i(x)$ be the p.d.f of the i^{th} population
(Π_i)

define: $P(j|i) = \text{Prob. of misclassifying an}\br/> \text{obsr. into } j^{th} \text{ population}\br/> \text{when actually it belongs to}\br/> i^{th} \text{ pop^n}\br/> = \Pr[x \in \Pi_j | \Pi_i \text{ is actual pop^n}]$

Some Definitions

1. Average cost (A.C), it is defined as.

$$\sum_{j=1}^k P(j|i) c(j|i) = \bar{x}_i(R) \text{ (say)}$$

now, if we have different rules of classification
(say) : R' , R'' , etc. we will have
different A.C like $\bar{x}_i(R)$, $\bar{x}_i(R')$, $\bar{x}_i(R'')$ etc.

2. Among the rules of classifications (say)

R_i and R_i^* , R_i is said to be better than

$$R^* \text{ if } \bar{x}_i(R_i) < \bar{x}_i(R^*) \quad i=1, \dots, k$$

with atleast one inequality

R_i is said to be ADMISSIBLE if there

does not exist any R^* which is
better than R_i .

3 BAYES RULE

Let q_i be the prob. that an
individual belongs to Π_i (i.e. a prior prob).
then the A.C. of m.c is defined

as $\sum_{i=1}^K q_i \delta_i(R)$ w.r.t a prior prob.
 (q_1, q_2, \dots, q_K)

A Bayes rule is defined as :

Given $q : (q_1, q_2, \dots, q_K)$ a classification rule
 which minimizes the a.c. of m.c is
 called the Bayes rule w.r.t. q .

Note This may result in large prob. of m.c
 though it is optimum.

* Problem of classification into two populations
 (General theory)

Let q_1 be the probability that an observation selected at random belongs to Π_1 and q_2 be the probability that it belongs to Π_2 . Again let us suppose that $p_1(x)$ be the pdf of Π_1 and $p_2(x)$ be the pdf of Π_2 define

$P(2|1)$: Prob of misclassification of an observation Π_2 when actually it belongs to Π_1 .

= $P[x \in R_2 | \text{when } p_1(x) \text{ is the actual pdf}]$

i.e

$$P(2|1) = \int_{R_2} p_1(x) dx, \text{ similarly } P(1|2) = \int_{R_1} p_2(x) dx$$

\therefore the prob of 1st type of m.c = $q_1 P(2|1)$
 and the prob of 2nd type of m.c = $q_2 P(1|2)$

So the Average cost (A.C)

$$= C(1/2) \cdot q_2 P(1|2) + C(2|1) q_1 P(2|1)$$

$$= C(1/2) \cdot q_2 \int_{R_1} p_2(x) dx + C(2|1) \int_{R_2} p_1(x) dx$$

now we wish to obtain

R_1 and R_2 for a good classification in such a way that A.C (R_1, R_2) is minimum.

i.e If R_1^* & R_2^* be any other regions of classification than $A.C(R_1, R_2) \leq A.C(R_1^*, R_2^*)$

i.e

$$q_1 c(2|1) \int_{R_2} P_1(x) dx + q_2 c(1|2) \int_{R_1} P_2(x) dx \leq$$

$$q_1 c(2|1) \cdot \int_{R_2^*} P_1(x) dx + q_2 c(1|2) \int_{R_1^*} P_2(x) dx$$

— (*)

as $P_1(x)$ and $P_2(x)$ are the pdf

$$\int_{R_1 \cup R_2} P_i(x) dx = 1 \quad \text{and} \quad \int_{R_1^* \cup R_2^*} P_i(x) dx = 1$$

i.e

$$\int_{R_2} P_i(x) dx = 1 - \int_{R_1} P_i(x) dx \quad \text{and}$$

$$\int_{R_2^*} P_i(x) dx = 1 - \int_{R_1^*} P_i(x) dx \quad \text{for } i=1,2$$

Substituting all the above values in the expression (*) we get

$$q_1 c(2|1) \left[1 - \int_{R_1} P_1(x) dx \right] + q_2 c(1|2) \int_{R_1} P_2(x) dx$$

$$\leq q_1 c(2|1) \left[1 - \int_{R_1^*} P_1(x) dx \right] + q_2 c(1|2) \int_{R_1^*} P_2(x) dx$$

i.e

$$\int_{R_1} [q_2 c(1|2) P_2(x) - q_1 c(2|1) P_1(x)] dx$$

$$- \int_{R_1^*} [q_2 c(1|2) P_2(x) - q_1 c(2|1) P_1(x)] dx \leq 0$$

this condition will be satisfied, if we choose R_1 s.t

$q_2 \cdot C(2|2) P_2(x) - q_1 \cdot C(2|1) P_1(x) \leq 0$ inside R_1
and $q_2 \cdot C(1|2) P_2(x) - q_1 \cdot C(2|1) P_1(x) \geq 0$ outside R_1

So, R_1 is the region where

$$q_1 \cdot C(2|1) P_1(x) \geq q_2 \cdot C(1|2) P_2(x)$$

or,
$$\frac{P_1(x)}{P_2(x)} \geq \frac{q_2}{q_1} \cdot \frac{C(1|2)}{C(2|1)}$$

and
$$\frac{P_1(x)}{P_2(x)} < \frac{q_2}{q_1} \cdot \frac{C(1|2)}{C(2|1)}$$

If q_1, q_2 are known and $C(1|2), C(2|1)$ are also known we can obtain R_1 and R_2 if $P_1(x)$ and $P_2(x)$ are completely specified.

SPECIAL CASE : If $q_1 = q_2 = \frac{1}{2}$ and

$$C(1|2) = C(2|1) \text{ then RHS} = 1$$

(*) Problem of classification (in to two MND) :-

Let \underline{x}_{px1} be an observations vector and suppose that we have to classify it into two categories i.e we assume that an observation can come either from the pop¹ $\pi_1 : N_p(\underline{\mu}^{(1)}, \underline{\Sigma})$ or from the pop² $\pi_2 : N_p(\underline{\mu}^{(2)}, \underline{\Sigma})$ equality of

$$f_1(x) = \frac{1}{(2\pi)^{p/2} |\underline{\Sigma}|^{1/2}} e^{-\frac{1}{2} (\underline{x} - \underline{\mu}^{(1)})^T \underline{\Sigma}^{-1} (\underline{x} - \underline{\mu}^{(1)})}$$

$$f_2(\underline{x}) = \frac{-\frac{1}{2}(\underline{x} - \underline{\mu}^{(2)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(2)})}{(2\pi)^{D/2} |\Sigma|^{1/2}}$$

$$\frac{f_1(\underline{x})}{f_2(\underline{x})} = e^{-\frac{1}{2}(\underline{x} - \underline{\mu}^{(1)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(1)})}$$

$$e^{-\frac{1}{2}(\underline{x} - \underline{\mu}^{(2)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(2)})}$$

$$= e^{-\frac{1}{2}(\underline{x} - \underline{\mu}^{(1)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(1)}) + \frac{1}{2}(\underline{x} - \underline{\mu}^{(2)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(2)})}$$

the Region of classification are given by

$$R_1 : \frac{f_1(\underline{x})}{f_2(\underline{x})} \geq k \quad \text{and} \quad R_2 : \frac{f_1(\underline{x})}{f_2(\underline{x})} < k$$

so, hence

$$R_1 = e^{-\frac{1}{2}[(\underline{x} - \underline{\mu}^{(2)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(2)}) - (\underline{x} - \underline{\mu}^{(1)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(1)})]} \geq k$$

OR

$$R_1 = \frac{1}{2} [(\underline{x} - \underline{\mu}^{(2)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(2)}) - (\underline{x} - \underline{\mu}^{(1)})^T \Sigma^{-1}(\underline{x} - \underline{\mu}^{(1)})] \geq \log k$$

OR

$$R_1 = \frac{1}{2} \left[\frac{1}{2} \underline{x}' \Sigma^{-1} \underline{x} - \underline{x}' \Sigma^{-1} \underline{\mu}^{(2)} - \underline{\mu}^{(2)'} \Sigma^{-1} \underline{x} + \underline{\mu}^{(2)'} \Sigma^{-1} \underline{\mu}^{(2)} \right. \\ \left. - \frac{1}{2} \underline{x}' \Sigma^{-1} \underline{x} - \underline{x}' \Sigma^{-1} \underline{\mu}^{(1)} - \underline{\mu}^{(1)'} \Sigma^{-1} \underline{x} + \underline{\mu}^{(1)'} \Sigma^{-1} \underline{\mu}^{(1)} \right]$$

$$\geq \log k$$

$$\text{OR } R_1 = \frac{1}{2} \text{ Since } \underline{x}' \Sigma^{-1} \underline{\mu}^{(2)} = \underline{\mu}^{(2)'} \Sigma^{-1} \underline{x}$$

$$\text{and } \underline{x}' \Sigma^{-1} \underline{\mu}^{(1)} = \underline{\mu}^{(1)'} \Sigma^{-1} \underline{x}$$

$$\therefore R_1 = \frac{1}{2} [2\underline{x}' \Sigma^{-1} \underline{\mu}^{(1)} - 2\underline{x}' \Sigma^{-1} \underline{\mu}^{(2)} - \underline{\mu}^{(1)'} \Sigma^{-1} \underline{\mu}^{(1)} + \underline{\mu}^{(2)'} \Sigma^{-1} \underline{\mu}^{(2)}] \geq \log k$$

OR

$$R_1 = \underline{x}^T \underline{g}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) - \frac{1}{2} \left\{ (\underline{\mu}^{(1)} + \underline{\mu}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) \right\}$$

$\geq \log k$

OR

$$R_1 = \underline{x}^T \underline{g}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \geq \frac{1}{2} \left\{ (\underline{y}^{(1)} + \underline{y}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) \right\}$$

$+ \log k$

$$R_2 = \underline{x}^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) \leq \frac{1}{2} \left\{ (\underline{\mu}^{(1)} + \underline{\mu}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) \right\}$$

$+ \log k$

Under the assumption of the two populations equally likely i.e. $q_1 = q_2 = \frac{1}{2}$ and if the two costs of misclassifications are SAME i.e. $C(1|2) = C(2|1)$, $\log k = 0$. Also,

Also, let $\underline{g}^{-1}(\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) = \underline{q}_{px_1}$ then

$$R_1 : \underline{x}^T \underline{d} \geq \frac{1}{2} (\underline{\mu}^{(1)} + \underline{\mu}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)})$$

$$R_2 : \underline{x}^T \underline{d} \leq \frac{1}{2} (\underline{\mu}^{(1)} + \underline{\mu}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)})$$

$$C(1|2) P(1|2) q_{12} = C(2|1) P(2|1) q_{21} = C \text{ (say)}$$

$$\text{let } U = \underline{x}^T \underline{g}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) - \frac{1}{2} \left\{ (\underline{y}^{(1)} + \underline{y}^{(2)})^T \underline{g}^{-1} (\underline{\mu}^{(1)} - \underline{\mu}^{(2)}) \right\}$$

then $R_1 = U > c$ and $R_2 : U \leq c$,

In order to find out the value of 'c' we must know the dist² of U. If X comes from II, then $P_2(U)$ can be obtained.

We have

$$P(1|2) = P_{\gamma} \{ U \geq c | X \in \Pi_2 \} = P_{\gamma} \{ U \geq c | X \sim N_p(\underline{y}^{(2)}, \underline{\Sigma}) \}$$

$$\text{and } P(2|1) = P_{\gamma} \{ U < c | X \in \Pi_1 \} = P_{\gamma} \{ U < c | X \sim N_p(\underline{y}^{(1)}, \underline{\Sigma}) \}$$

it is observed that U is a linear combination of γ which is normally distributed
So U is NORMALLY distributed

1. When $X \sim N_p(\underline{y}^{(1)}, \underline{\Sigma})$,

$$E_1(U) = \underline{y}^{(1)}' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) - \frac{1}{2} \{ (\underline{y}^{(1)} + \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \}$$

OR

$$E_2(U) = -\frac{1}{2} (\underline{y}^{(1)} - \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)})$$

Also, Mahalanobis - D^2 (Distance) is given by
 $(\underline{y}^{(1)} - \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) = \alpha$ (say)

$$\text{i.e. } E_1(U) = \alpha/2$$

$$\text{Again, } \text{Var}_1(U) = E[U - E_1(U)]^2$$

OR

$$\text{Var}_1(U) = [X' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) - \frac{1}{2} \{ (\underline{y}^{(1)} + \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \}]^2$$

$$= E[(\underline{y}^{(1)} - \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)})]^2 / 2 \{ (\underline{y}^{(1)} - \underline{y}^{(2)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \}$$

$$= (\underline{y}^{(1)} - \underline{y}^{(2)})' \underline{\Sigma}^{-1} E(X - \underline{y}^{(1)}) (X - \underline{y}^{(1)})' \underline{\Sigma}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)})$$

$$= (\underline{y}^{(1)} - \underline{y}^{(2)})' (\underline{y}^{(1)} - \underline{y}^{(2)})$$

$$= \alpha$$

$$\text{i.e. } \text{Var}_1(U) = \alpha$$

2. If $X \sim N_p(\underline{\mu}(1), \underline{\Sigma})$ then

$$P_2(v) = -\alpha/2 \quad \text{and} \quad V_2(v) = \alpha$$

So,

If $X \in \Pi_1$, $v \sim N(\alpha/2, \alpha)$ and

If $X \in \Pi_2$, $v \sim N(-\alpha/2, \alpha)$

$$X \in \Pi_1, P(v) = \frac{1}{\sqrt{2\pi}\alpha} e^{-\frac{1}{2\alpha}(v-\frac{\alpha}{2})^2}$$

$$X \in \Pi_2, P(v) = \frac{1}{\sqrt{2\pi}\alpha} e^{-\frac{1}{2\alpha}(v+\frac{\alpha}{2})^2}$$

$$\therefore C(2)_{12} = \int_{-\infty}^{\frac{\alpha}{2}} \frac{1}{\sqrt{2\pi}\alpha} e^{-\frac{1}{2\alpha}(v-\frac{\alpha}{2})^2} dv$$

and

$$\therefore C(1|2) = \int_{-\frac{\alpha}{2}}^{\infty} \frac{1}{\sqrt{2\pi}\alpha} e^{-\frac{1}{2\alpha}(v+\frac{\alpha}{2})^2} dv$$

$$\text{let us take } v - \frac{\alpha}{2} = z \quad \text{and} \quad v + \frac{\alpha}{2} = t$$

$$\text{then } C(2|1) q_1 \int_{-\infty}^{\frac{\alpha}{2} - \frac{1}{2}\sqrt{\alpha}} \phi(z) dz$$

$$= C(1|2) q_2 \int_{\frac{\alpha}{2} + \frac{1}{2}\sqrt{\alpha}}^{\infty} \phi(t) dt \quad --- (*)$$

\therefore knowing the values of $q_1, q_2, C(1|2)$ and $C(2|1)$ ($*$) can be evaluated using the tables of Standard Normal dist² and we can obtain the value of ' C '.

Hence

$$R_1 : \underline{x}' \underline{\delta}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \geq \frac{1}{2} \left[(\underline{y}^{(1)} + \underline{y}^{(2)})' \underline{\delta}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \right] + C$$

$$R_2 : \underline{x}' \underline{\delta}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \leq \frac{1}{2} \left[(\underline{y}^{(1)} + \underline{y}^{(2)})' \underline{\delta}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)}) \right] + C$$

* Classification when the Population Parameters are estimated :-

(Fisher's Criterion 1936)

Let \underline{x}_{pxq} be a vector it has to be discriminated into two population $N_p(\underline{y}^{(1)}, \underline{\delta})$ and $N_p(\underline{y}^{(2)}, \underline{\delta})$ consider another vector \underline{l}_{pxq} then Fisher suggested a linear combination $\underline{l}' \underline{x}$ for the discrimination where \underline{l} is to be chosen such that $[E_1(\underline{l}' \underline{x}) - E_2(\underline{l}' \underline{x})]^2$ is maximum such that consider $\text{var}(\underline{l}' \underline{x}) = k$ (constant) generally taken to be unity. Now, in vector to discriminate \underline{x} we have to maximize $(\underline{l}' \underline{y}^{(1)} - \underline{l}' \underline{y}^{(2)})$ such that consider $\text{var}(\underline{l}' \underline{x}) = \underline{l}' \underline{\delta} \underline{l} = 1$

Again, $[(\underline{l}' \underline{y}^{(1)} - \underline{l}' \underline{y}^{(2)})]^2 = (\underline{l}' \underline{d})^2 = \underline{l}' \underline{d} \underline{d}' \underline{l}$
where $\underline{d} = (\underline{y}^{(1)} - \underline{y}^{(2)})$

∴ in order to maximize we consider a f^n ϕ using Lagrange's multiplier ' λ ' i.e

$$\begin{aligned}\phi &= \underline{l}' \underline{d} \underline{d}' \underline{l} - \lambda [\underline{l}' \underline{\delta} \underline{l} - k] \\ &= \underline{l}' \underline{d} \underline{d}' \underline{l} - \lambda \underline{l}' \underline{\delta} \underline{l} - \lambda k\end{aligned}$$

differentiating ϕ w.r.t \underline{d} , we have :

$$\frac{\partial \phi}{\partial \underline{d}} = 2 \underline{d} \cdot \underline{d}' \underline{l} - 2 \lambda \sum \underline{l} = 0$$

$\Rightarrow \underline{d}' \underline{d} \cdot \underline{l} = \lambda \sum \underline{l}$, premultiplying it by \underline{l}'

we have

$$\underline{l}' \underline{d}' \underline{d} \cdot \underline{l} = \lambda \underline{l}' \cdot \underline{l} = \lambda k \quad | \quad \underline{d}' \underline{l}' \underline{l} = \underline{k}$$

i.e.

$$(\underline{l}' \underline{d})^2 = \lambda k \Rightarrow \underline{l}' \underline{d} = \sqrt{\lambda k} = \underline{d}' \underline{l}$$

we have

$$\underline{d}' \sqrt{\lambda k} = \lambda \sum \underline{l} \Rightarrow \underline{l} = \frac{\sqrt{\lambda k} \underline{d}'}{\lambda} = \frac{\sqrt{\lambda} \underline{d}'}{\sqrt{k}}$$

i.e.

$$\underline{l} = \frac{\sqrt{\lambda} \underline{d}'}{\sqrt{k}} (\underline{y}^{(1)} - \underline{y}^{(2)})$$

so, the best linear discriminant function

$$\text{is } \underline{l}' \underline{x} \text{ where } \underline{l} \propto \underline{d}^{-1} (\underline{y}^{(1)} - \underline{y}^{(2)})$$

If the parameters $\underline{y}^{(1)}$, $\underline{y}^{(2)}$ and \underline{d} are unknown

$$\bar{\underline{x}}^{(1)} = \frac{1}{N_1} \sum_{\alpha=1}^{N_1} \underline{x}_{\alpha}^{(1)}, \quad \bar{\underline{x}}^{(2)} = \frac{1}{N_2} \sum_{\alpha=1}^{N_2} \underline{x}_{\alpha}^{(2)}$$

and

$$S = \frac{(N_1 - 1) S_1 + (N_2 - 1) S_2}{(N_1 + N_2 - 2)}$$

$$\text{where } S_i = \frac{1}{(N_i - 1)} \sum_{\alpha=1}^{N_i} (\underline{x}_{\alpha}^{(i)} - \bar{\underline{x}}^{(i)}) (\underline{x}_{\alpha}^{(i)} - \bar{\underline{x}}^{(i)})'$$

$$\therefore R_1 = \mathbf{x}^T \mathbf{S}^{-1} (\bar{\mathbf{x}}^{(1)} - \bar{\mathbf{x}}^{(2)})^T \geq \frac{1}{2} (\bar{\mathbf{x}}^{(1)} + \bar{\mathbf{x}}^{(2)})^T \mathbf{S}^{-1} (\bar{\mathbf{x}}^{(1)} - \bar{\mathbf{x}}^{(2)})$$

$$R_2 = \mathbf{x}^T \mathbf{S}^{-1} (\bar{\mathbf{x}}^{(1)} - \bar{\mathbf{x}}^{(2)})^T \leq \frac{1}{2} (\bar{\mathbf{x}}^{(1)} + \bar{\mathbf{x}}^{(2)})^T \mathbf{S}^{-1} (\bar{\mathbf{x}}^{(1)} - \bar{\mathbf{x}}^{(2)})$$

→ Remark :- When $\mathbf{S}_1 + \mathbf{S}_2$ then FISHER suggested to use $(\mathbf{x} - \mathbf{y}^{(1)})^T \mathbf{S}_1^{-1} (\mathbf{x} - \mathbf{y}^{(1)}) + (\mathbf{x} - \mathbf{y}^{(2)})^T \mathbf{S}_2^{-1} (\mathbf{x} - \mathbf{y}^{(2)})$ as the discriminant if it is of a Q.F. Quadratic form distⁿ function.

Obviously it is not a linear discriminant but it is of a quadratic form therefore we shall called it a quadratic discriminant function.

* DIFFICULTIES :-

the two regions of classifications are given by R_1 and R_2 but if in practical situations we get an equality then what to do?

Then it becomes subjective and a difficult problem to classify the random vector \mathbf{x} either in Π_1 or Π_2 .

Hence instead of two regions we must have a third region also in which the observation may fall. Otherwise we must have some more observations to classify properly the random vector \mathbf{x} .

* Canonical Correlations Analysis (CCA)

Introduction (CCA)

The most commonly used dependence analysis in multivariate is the multiple regression analysis. The technique deals with the study of one variable on a set of predictor variables. But MRA does n't provide any information about the correlation of the two sets of variables. taking into the consideration of the interrelationship among the variables in the wide set, Y-set.

CCA is an alternative multivariate method use to study the correlation of the variables in the Y-set and the X-set taking consideration of the interdependence of the variables in the Y-set.

For Example:- In biological experiments we may be interested in studying the impact of the amount of insecticide and its concentration on the number of dead insects and on the no. of days required to kill the insects.

Criterion sets include the no. of dead insects and the no. of days required to kill the insect. The predictor set includes the amount of insecticide and its concentration. The former set is usually denoted by Y-set and the later set is denoted by X-set.

The MCA deals with one variable at a time from Y set.

However CCA is appropriate when the criterion variables themselves are correlated. The main purpose of CCA is to study the structural relationship of two sets of variables simultaneously.

* Canonical Variables and Canonical Correlation

Let \underline{x}_{px} be a random vector with μ and Σ as the parameters of the popⁿ under consideration. Let \underline{x} be partitioned into two subset $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ having p_1 and p_2 components respectively ($p_1 < p_2 \rightarrow (?)$) the Σ is also partitioned as:

$$\Sigma = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix}$$

Let α and β be any two $(p_1 \times 1)$ and $(p_2 \times 1)$ vectors and let $U = \alpha^T \underline{x}^{(1)}$ and $V = \beta^T \underline{x}^{(2)}$ be two linear combination of $\underline{x}^{(1)}$ and $\underline{x}^{(2)}$ respectively. We have

$$\text{Var}(U) = \text{Var}(\alpha^T \underline{x}^{(1)}) = \alpha^T \Sigma_{11} \alpha \text{ and}$$

$$\text{Var}(V) = \text{Var}(\beta^T \underline{x}^{(2)}) = \beta^T \Sigma_{22} \beta \text{ and}$$

$$\text{Cov}(U, V) = \alpha^T \Sigma_{12} \beta.$$

We wish to choose α and γ in such a way that the correlation between U and V is maximum such that consider that their variance are constant.

We may normalize the linear combination so that $\text{Var}(U) = \text{Var}(V) = 1$

$$\text{i.e. } \alpha^T \Sigma_{11} \alpha = \gamma^T \Sigma_{22} \gamma = 1 \quad \text{--- (1)}$$

Now, the correlation between U and V is

$$\begin{aligned} \rho_{UV} &= \text{Cov}(U, V) = \frac{\alpha^T \Sigma_{12} \gamma}{\sqrt{\text{Var}(U)} \cdot \sqrt{\text{Var}(V)}} \\ &= \frac{\alpha^T \Sigma_{12} \gamma}{\sqrt{\alpha^T \Sigma_{11} \alpha} \sqrt{\gamma^T \Sigma_{22} \gamma}} \end{aligned}$$

if (1) holds

$$\text{i.e. } \rho_{UV} = \frac{\alpha^T \Sigma_{12} \gamma}{\sqrt{\alpha^T \Sigma_{11} \alpha}} \quad \text{--- (2)}$$

Obviously maximizing ρ_{UV} is equivalent to maximizing $\alpha^T \Sigma_{12} \gamma$ such that consider (1). For this consider the function

$$\phi = \frac{\alpha^T \Sigma_{12} \gamma}{2} - \frac{\lambda_1}{2} (\alpha^T \Sigma_{11} \alpha - 1) - \frac{\lambda_2}{2} (\gamma^T \Sigma_{22} \gamma - 1)$$

where λ_1 and λ_2 are the lagrange's multipliers. Now, differentiating ϕ w.r.t α and γ respectively,

we have

$$\frac{\partial \phi}{\partial \alpha} = \Sigma_{12} \gamma - \frac{\lambda_1}{2} 2 \Sigma_{11} \alpha = 0$$

$$\Rightarrow \Sigma_{12} \gamma = \lambda_1 \Sigma_{11} \alpha \quad \text{--- (3)}$$

$$\frac{d\Phi}{d\gamma} = \underline{\alpha}' \Sigma_{12} - \lambda_2 \cdot 2 \underline{\gamma}' \Sigma_{22} = 0$$

$$\Rightarrow \underline{\alpha}' \Sigma_{12} - \lambda_2 \underline{\gamma}' \Sigma_{22} = 0 \quad \text{or.}$$

$$\Rightarrow \Sigma_{21} \underline{\alpha} - \lambda_2 \Sigma_{22} \underline{\gamma} = 0 \quad \rightarrow (4)$$

(transforming the equation) (transposing the eqn)

premultiplying (3) by $\underline{\alpha}$ and (4) by $\underline{\gamma}'$ we have,

$$(\underline{\alpha}' \Sigma_{12} \underline{\gamma} - \lambda_1 \underline{\alpha}' \Sigma_{11} \underline{\alpha} = 0 \quad \text{i.e.} \quad \underline{\alpha}' \Sigma_{12} \underline{\gamma} = \lambda_1 \underline{\alpha}' \Sigma_{11} \underline{\alpha}) \quad \underline{\alpha}' \Sigma_{11} \underline{\alpha} = 1$$

$$\text{and} \quad (\underline{\gamma}' \Sigma_{21} \underline{\alpha} - \lambda_2 \underline{\gamma}' \Sigma_{22} \underline{\gamma} = 0 \quad \text{i.e.} \quad \underline{\gamma}' \Sigma_{22} \underline{\gamma} = \lambda_2 \underline{\gamma}' \Sigma_{22} \underline{\gamma} = 1)$$

as $\underline{\alpha}' \Sigma_{12} \underline{\gamma}$ and $\underline{\gamma}' \Sigma_{21} \underline{\alpha}$ are same seeing the correlation between $\underline{\alpha}$ & $\underline{\gamma}$

$$\therefore \lambda_1 = \lambda_2 = \rho_{UV} = \rho \text{ (say)}$$

Again from (3) and (4) we have.

$$\Sigma_{12} \underline{\gamma} - \rho \Sigma_{11} \underline{\alpha} = 0 \quad \rightarrow (5)$$

$$\Sigma_{21} \underline{\alpha} - \rho \Sigma_{22} \underline{\gamma} = 0 \quad \rightarrow (6)$$

$$\text{Or, } -\rho \Sigma_{11} \underline{\alpha} + \Sigma_{12} \underline{\gamma} = 0$$

$$\Sigma_{21} \underline{\alpha} - \rho \Sigma_{22} \underline{\gamma} = 0$$

which can be written in matrix form as:

$$\begin{bmatrix} -\rho \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & -\rho \Sigma_{22} \end{bmatrix} \begin{bmatrix} \underline{\alpha}' \\ \underline{\gamma}' \end{bmatrix} = 0 \quad \rightarrow (7)$$

Since, $\underline{\alpha}$ and $\underline{\gamma}$ are NON-zero vectors

\therefore for a non-zero solution of (7) will exist only if the coefficient matrix is singular i.e.

$$\begin{vmatrix} -\rho \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & -\rho \Sigma_{22} \end{vmatrix} = 0 \quad \rightarrow (8)$$

Since, the LHS of (8) is a polynomial of order ' p' so there will be ' p ' roots of φ .

Again, premultiplying (5) by $\Sigma_{21}\Sigma_{11}^{-1}$ and (6) by φ we get

$$\Sigma_{21}\Sigma_{11}^{-1}\Sigma_{12}\varphi - \varphi\Sigma_{21}\Sigma_{11}^{-1}\Sigma_{11}\varphi = 0 \quad \text{OR}$$

$$\Sigma_{21}\Sigma_{11}^{-1}\Sigma_{12}\varphi - \varphi\Sigma_{21}\varphi = 0$$

$$\varphi\Sigma_{21}\varphi - \varphi^2\Sigma_{22}\varphi = 0 \quad \text{OR} \quad \varphi\Sigma_{21}\varphi - \varphi^2\Sigma_{22}\varphi = 0$$

Adding both we get.

$$(\Sigma_{21}\Sigma_{11}^{-1}\Sigma_{12}\varphi - \varphi^2\Sigma_{22}\varphi) = 0 \quad \text{--- (9)}$$

Again, φ is a non-zero vector, so we must have $(\Sigma_{21}\Sigma_{11}^{-1}\Sigma_{12}\varphi - \varphi^2\Sigma_{22}\varphi) = 0 \quad \text{--- (10)}$

Since, the LHS of (10) is a polynomial of order p_2 in φ^2 , so it gives p_2 roots of φ^2 (say) $\varphi_1^2, \varphi_2^2, \dots, \varphi_{p_2}^2$.

Now premultiplying (5) by φ and (6) by $\Sigma_{12}\Sigma_{21}^{-1}$ we have (get)

$$\varphi\Sigma_{12}\varphi - \varphi^2\Sigma_{11}\varphi = 0$$

$$\Sigma_{12}\Sigma_{21}^{-1}\Sigma_{21}\varphi - \varphi\Sigma_{12}\Sigma_{21}^{-1}\Sigma_{22}\varphi = 0$$

Adding these

$$(\Sigma_{12}\Sigma_{21}^{-1}\Sigma_{21} - \varphi^2\Sigma_{11})\varphi = 0 \quad \text{--- (11)} \quad \text{Again, } \varphi \neq 0$$

$$|\Sigma_{12}\Sigma_{21}^{-1}\Sigma_{22}| - \varphi^2|\Sigma_{11}| = 0 \quad \text{--- (12)}$$

which is a polynomial of order p_1 in φ^2 .

which gives p_1 roots as $\varphi_1^2, \varphi_2^2, \dots, \varphi_{p_1}^2$.

Again

As the equations (9) and (11) have been obtained from (7) hence their non-zero roots should be equal and there must be $(p_2 - p_1)$ ^{zero} roots, so the number of non-zero roots of S is:

$$\underline{p} = (p_2 - p_1) = (p_1 + p_2) - (p_2 - p_1) = 2p_1$$

$\therefore S^2$ has p_1 non-zero roots and S has $[2p_1]$ non-zero roots of the form $= S = \pm e_i$ ($i=1, 2, \dots, p_1$)

now, the positive roots are arranged as:

$$S(1) \geq S(2) \geq \dots \geq S(p_1)$$

Here $S(1)$ is called the FIRST CANONICAL CORR and substituting $S = S(1)$ in (9) and (11) and solving them for α and β

Let $\underline{\alpha}^{(1)}$ and $\underline{\gamma}^{(1)}$ be the solution then

$$[V_1 = \underline{\alpha}^{(1)}' \underline{X}^{(1)}] \text{ and } [V_1 = \underline{\gamma}^{(1)}' \underline{X}^{(2)}]$$

are the 1st Canonical Variables

Similarly $S(2)$ representing the SECOND CANONICAL CORRELATION and

$$[V_2 = \underline{\alpha}^{(2)}' \underline{X}^{(1)}] \text{ and } [V_2 = \underline{\gamma}^{(2)}' \underline{X}^{(2)}]$$

are the 2nd Canonical variables



CCA Properties :-

(1) For different values of canonical correlation we can obtain the vectors $\alpha^{(1)}, \dots, \alpha^{(P_1)}$ and $\gamma^{(1)}, \gamma^{(2)}, \dots, \gamma^{(P_2)}$ etc. define the matrices

$$A = [\alpha^{(1)} \dots \alpha^{(P_1)}] \text{ and } \Gamma = [\gamma^{(1)} \gamma^{(2)} \dots \gamma^{(P_2)}]$$

then we have

$$A^T \Sigma_{11} A = I_{P_1},$$

$$A^T \Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21} A = R^* = \log(\beta_1^2, \beta_2^2, \dots, \beta_{P_2}^2)$$

and

$$\Gamma^T \Sigma_{22} \Gamma = I_{P_2}$$

$$\Gamma^T \Sigma_{21} \Sigma_{11}^{-1} \Sigma_{21} \Gamma = \text{diag}(\beta_1^2, \beta_2^2, \dots, \beta_{P_2}^2)$$

(2) The number of non-zero roots of

$$|\Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21} - \beta^2 \Sigma_{11}| = 0 \text{ or }$$

$|\Sigma_{21} \Sigma_{11}^{-1} \Sigma_{12} - \beta^2 \Sigma_{22}| = 0$ is equal to the roots of Σ_{12}



Method of TE81 :-

The principle of CCA is to find out two linear combination $\alpha^1 x^{(1)}$ of the variables in $x^{(1)}$ and $x^{(2)}$ so that the linear combinations are highly correlated.

this is possible only if the variables in $x_1^{(1)}$ and $x_1^{(2)}$.

However, the CCA will not be useful if the variables in the two subsets are not related. therefore, before performing CCA we should test

$$H_0: \Sigma_{12} = 0$$

$$\text{eg. } H_1: \Sigma_{12} \neq 0.$$

Bartlett (1951) proposed a criterion which was further developed by WILK so it is known as WILK'S criterion, and the proposed test statistics is

$$\Delta = \frac{|S_{22}^{-1} S_2 S_{11}^{-1} S_{12}|}{|S_{22}|}$$

it is present that $\Delta \sim \chi^2_{p_1, p_2, n}$

$$\text{where } \chi^2 = -[(n-1) - \frac{1}{2}(p_1 + p_2 + 1)] \log \Delta$$

If $\chi^2_{\text{cal}} > \chi^2_{p_1, p_2, \alpha}$, H_0 is rejected

\therefore Proceed for CCA



BEST REGRESSION EQUATION

We have $\underline{x}_{px} = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_p \end{pmatrix}$

Assume $E(\underline{x}) = 0$

Consider the partitioning of \underline{x}_{px} as

$$\underline{x}_{px} = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_p \end{pmatrix} \text{ accordingly}$$

$$\underline{\sigma} = \begin{bmatrix} \text{var}(x_1) & \text{cov}(x_1, x^{(2)}) \\ \text{cov}(x^{(2)}, x_1) & \text{var}(x^{(2)}) \end{bmatrix}$$

where $x^{(2)} = (\bar{x}_1, x_2)$

denote $\text{var}(x_1) = \sigma_{11}$ (say)

$$\text{cov}(x_1, x^{(2)}) = \underline{\sigma}$$

$$\text{var}(x^{(2)}) = \Sigma_{22}$$

Now, let 'predicted' value of x_1 be $\underline{\beta}^1 x^{(2)}$

where $\underline{\beta} = \begin{bmatrix} \beta_2 \\ \beta_3 \\ \vdots \\ \beta_p \end{bmatrix}$ is the $(p-1) \times 1$ vector

We wish to obtain $\underline{\beta}$ in such a way that it minimizes the residual sum of squares let it be denoted by Ω . Then

$$\Omega = E[(x_1 - \underline{\beta}^1 x^{(2)})]^2$$

$$= E[(x_1 - \underline{\beta}^1 x^{(2)})] [(x_1 - \underline{\beta}^1 x^{(2)})]'$$

Or

$$Q = E[x_1^2 - x_1 \underline{x}^{(2)} \beta - \beta^T \underline{x}^{(2)} x_1 + \beta^T \underline{x}^{(2)} \underline{x}^{(2)} \beta]$$

$$= \sigma_{11} - 2\beta^T \sigma + \beta^T \Sigma_{22} \beta$$

Now differentiating Q with respect to β

$$\frac{\partial Q}{\partial \beta} = -2\sigma + 2\beta^T \Sigma_{22} = 0$$

$$\Rightarrow \beta^T \Sigma_{22} = \sigma \Rightarrow \boxed{\beta = \Sigma_{22}^{-1} \sigma}$$

So the predicted value of x_1 called Regression is given by

$$\boxed{\hat{x}_1 = \beta^T \underline{x}^{(2)} = \sigma^T \Sigma_{22}^{-1} \underline{x}^{(2)}}$$

↳ Regr of x_1 on (x_2, x_3, \dots, x_p)

$\therefore [x_1 - \sigma^T \Sigma_{22}^{-1} \underline{x}^{(2)}]$ is the residual denoted by $x_{1,23\dots p}$.

Note that $\text{cov}[x_{1,23\dots p}, \underline{x}^{(2)}] = 0$

i.e. residual on $\underline{x}^{(2)}$ are uncorrelated

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Multivariate Analysis - II

Theorem If an m -component vector $\mathbf{y} \sim N_m(\mathbf{0}, \mathbf{I})$

where T is a Non Singular Matrix then
 $\mathbf{y}' T^{-1} \mathbf{y} \sim \chi^2_m$ d.f

Proof Since T is NSM we can find another Matrix C
such that $C T C' = I$

Now consider the transformation

$$\mathbf{z} = C \mathbf{y}$$

then

$$\mathbf{z} \sim N_m(\mathbf{0}, CTC') \text{ or } \mathbf{z} \sim N_m(\mathbf{0}, \mathbf{I})$$

the components of \mathbf{z} are iid $N_1(0, 1)$

Now

$$\begin{aligned} \mathbf{y}' T^{-1} \mathbf{y} &= (C^{-1} \mathbf{z})' T^{-1} (C^{-1} \mathbf{z}) \\ &= \mathbf{z}' (C' T C)^{-1} \mathbf{z} = \mathbf{z}' \mathbf{I} \mathbf{z} \\ &= \mathbf{z}' \mathbf{z} = \sum_{i=1}^m z_i^2 \end{aligned}$$

Since z_i^2 are iid $N(0, 1)$

$$\mathbf{y}' T^{-1} \mathbf{y} = \sum_{i=1}^m z_i^2 \sim \chi^2_m$$

Application:- Then it is used in testing $H_0: \mathbf{y} = \mathbf{y}_0$
when σ^2 is known. As we know

that $\bar{\mathbf{x}} \sim N_p(\mathbf{u}, \frac{\sigma^2}{N} \mathbf{I}) \Rightarrow (\bar{\mathbf{x}} - \mathbf{y}) \sim N_p(\mathbf{0}, \frac{\sigma^2}{N} \mathbf{I})$

$$\Rightarrow \sqrt{N}(\bar{\mathbf{x}} - \mathbf{y}) \sim N_p(\mathbf{0}, \frac{\sigma^2}{N} \mathbf{I})$$

Or,

$$N(\bar{\mathbf{x}} - \mathbf{y})' \Sigma^{-1} (\bar{\mathbf{x}} - \mathbf{y}) \sim \chi^2_{p; \alpha, y}$$

calculated tabulated

So reject H_0 if $\chi^2_{\text{cal}} > \chi^2_{\text{tab}}$

In univariate case for testing the hypothesis

$H_0: \mathbf{y} = \mathbf{y}_0$ when σ^2 is unknown we use

$$t = (\bar{\mathbf{x}} - \mathbf{u}) \sim t_{n-1}$$

$$\frac{s}{\sqrt{n}}$$

Now $t^2 = n(\bar{x} - \mu)/(s^2) \sim f_{n-1}$

On the same analogy Hotelling defined a test for one sample / two sample as :

Let $x_1, x_2, \dots, x_n \sim N_p(\mu, \Sigma)$, Σ unknown,
if we wish to test $H_0: \mu = \mu_0$ (specified)
he proposed : $T^2 = N(\bar{x} - \mu_0)' S^{-1} (\bar{x} - \mu_0)$

$$\text{where } S = \frac{1}{(N-1)} \sum_{\alpha=1}^N (x_\alpha - \bar{x})(x_\alpha - \bar{x})'$$

* Likelihood Ratio Test for the Mean (Σ unknown)

We wish to test $H_0: \mu = \mu_0$, Σ being unknown
let x_1, x_2, \dots, x_n be a g.s of size N from
 $N_p(\mu, \Sigma)$ then its likelihood of Σ will be :

$$L = \frac{1}{(2\pi)^{\frac{Np}{2}} |\Sigma|^{N/2}} e^{-\frac{1}{2} \sum_{\alpha=1}^N (x_\alpha - \mu)' \Sigma^{-1} (x_\alpha - \mu)} \quad (11)$$

we define the LRT $\lambda = \frac{\max_{H_0} L}{\max_{H_1} L}$

The L.F will be max if the pmf are replaced by their MLE's, we have $\hat{\mu} = \bar{x}$, $\hat{\Sigma} = \hat{\Sigma}_0 = A/N$
where

$$A = \sum_{\alpha=1}^N (x_\alpha - \bar{x})(x_\alpha - \bar{x})' \rightarrow \text{under } H_1$$

and

$$\mu = \mu_0, \hat{\Sigma} = \hat{\Sigma}_0 = \frac{1}{N} \sum_{\alpha=1}^N (x_\alpha - \mu_0)(x_\alpha - \mu_0)' \text{ is under } H_0$$

Under H_0 we have

$$\max_{H_0} L = \frac{1}{(2\pi)^{\frac{Np}{2}} |\hat{\Sigma}_0|^{N/2}} e^{-\frac{1}{2} \sum_{\alpha=1}^N (x_\alpha - \mu_0)' \hat{\Sigma}_0^{-1} (x_\alpha - \mu_0)}$$

Consider the Exponent: $\sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)^T \hat{\Sigma}_0^{-1} (\underline{x}_\alpha - \underline{y}_0)$

$$\therefore \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)^T \hat{\Sigma}_0^{-1} (\underline{x}_\alpha - \underline{y}_0) = \text{tr} \cdot \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)^T \hat{\Sigma}_0^{-1} (\underline{x}_\alpha - \underline{y}_0) \xrightarrow{\text{→ } N \times 1 \text{ scales}}$$

$$= \sum_{\alpha=1}^N \text{tr} \cdot (\underline{x}_\alpha - \underline{y}_0)^T \hat{\Sigma}_0^{-1} (\underline{x}_\alpha - \underline{y}_0)$$

$$= \sum_{\alpha=1}^N \text{tr} \cdot \hat{\Sigma}_0^{-1} (\underline{x}_\alpha - \underline{y}_0)^T (\underline{x}_\alpha - \underline{y}_0)$$

$$= \text{tr} \cdot \hat{\Sigma}_0^{-1} \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)^T (\underline{x}_\alpha - \underline{y}_0)$$

$$= \text{tr} \cdot \hat{\Sigma}_0^{-1} (N) \hat{\Sigma}_0$$

$$= N \text{tr} I_p$$

$$= Np$$

$$\therefore \max L = \frac{1}{H_0} e^{-\frac{1}{2} NP} \quad (2) \quad H_0 = (2\pi)^{\frac{NP}{2}} |\hat{\Sigma}_0|^{N/2}$$

Now

$$\max L = \frac{1}{H_1} e^{-\frac{1}{2} \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{x})^T \hat{\Sigma}_1^{-1} (\underline{x}_\alpha - \bar{x})} \quad (3) \quad H_1 = (2\pi)^{\frac{NP}{2}} |\hat{\Sigma}_1|^{N/2}$$

Again it can be proved that the exponent of (3) is $-\frac{NP}{2}$

∴ Using (2) and (3)

$$\lambda = \frac{|\hat{\Sigma}_1|^{N/2}}{|\hat{\Sigma}_0|^{N/2}} = \frac{|A|^{\frac{N}{2}}}{\left| \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)(\underline{x}_\alpha - \underline{y}_0)^T \right|^{N/2}}$$

$$\text{Now } \sum_{\alpha=1}^N (\underline{x}_\alpha - \underline{y}_0)(\underline{x}_\alpha - \underline{y}_0)^T$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{x} + \bar{x} - \underline{y}_0)(\underline{x}_\alpha - \bar{x} + \bar{x} - \underline{y}_0)^T$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{x})(\underline{x}_\alpha - \bar{x})' + \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{x})(\bar{x} - \underline{y}_0)' +$$

$$= \sum_{\alpha=1}^N (\bar{x} - \underline{y}_0)(\underline{x}_\alpha - \bar{x})' + \sum_{\alpha=1}^N (\bar{x} - \underline{y}_0)(\bar{x} - \underline{y}_0)'$$

$$= \sum_{\alpha=1}^N (\underline{x}_\alpha - \bar{x})'(\underline{x}_\alpha - \bar{x}) + \sum_{\alpha=1}^N (\bar{x} - \underline{y}_0)(\bar{x} - \underline{y}_0)'$$

$$= A + N(\bar{x} - \underline{y}_0)(\bar{x} - \underline{y}_0)'$$

$$\therefore \lambda^2/N = \frac{|A|}{|A + N(\bar{x} - \underline{y}_0)(\bar{x} - \underline{y}_0)'|}$$

Consider the denominator

$$|A + N(\bar{x} - \underline{y}_0)(\bar{x} - \underline{y}_0)'| = \begin{vmatrix} A & \sqrt{N}(\bar{x} - \underline{y}_0) \\ -\sqrt{N}(\bar{x} - \underline{y}_0)' & 1 \end{vmatrix}$$

(*)

Now using the result : if $|B| \neq 0$ and E is a square matrix then $\begin{vmatrix} B & C \\ D & E \end{vmatrix} = |B| |E - DB^{-1}C|$

applying it to (*) we have

$$\begin{vmatrix} A & \sqrt{N}(\bar{x} - \underline{y}_0) \\ -\sqrt{N}(\bar{x} - \underline{y}_0)' & 1 \end{vmatrix} = |A| \begin{vmatrix} 1 + N(\bar{x} - \underline{y}_0)' A^{-1}(\bar{x} - \underline{y}_0) \end{vmatrix}$$

$$\lambda^2/N = \frac{|A|}{|A| \begin{vmatrix} 1 + N(\bar{x} - \underline{y}_0)' A^{-1}(\bar{x} - \underline{y}_0) \end{vmatrix}} = \frac{1}{1 + N(\bar{x} - \underline{y}_0)' A^{-1}(\bar{x} - \underline{y}_0)}$$

$$\text{Again } A = (N-1)S \Rightarrow A^{-1} = \frac{S^{-1}}{(N-1)}$$

$$\text{Let } T^2 = N(\bar{x} - \underline{y}_0)' S^{-1}(\bar{x} - \underline{y}_0)$$

then

$$\therefore \lambda^2/N = \frac{1}{1 + N(\bar{x} - \underline{y}_0)' S^{-1}(\bar{x} - \underline{y}_0)} = \frac{1}{1 + T^2}$$

$$\frac{1 + N(\bar{x} - \underline{y}_0)' S^{-1}(\bar{x} - \underline{y}_0)}{(N-1)} = \frac{1 + T^2}{(N-1)}$$

$$\frac{1}{1+\frac{T^2}{N-1}} \leq \lambda_0^{2/N} \Rightarrow 1 + \frac{T^2}{(N-1)} \geq \lambda_0^{-2/N}$$

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The LRT is to Reject H_0 if $\lambda \leq \lambda_0$ (specified)

OR

$$\lambda^{2/N} \leq \lambda_0^{2/N} \quad \text{OR} \quad 1 + \frac{T^2}{N-1} \leq \lambda_0^{-2/N}$$

i.e. $T^2 \geq [\lambda_0^{-2/N} - 1](N-1) = T_0^2$ (say) is a
another constant to be determined in such
a way that

$$P_{\lambda} [T^2 \geq T_0^2 | H_0] = \alpha \quad (\text{Size of the test})$$

Hence, for a desired size α test to obtain T_0^2 we
must know that the distⁿ of T^2

* Distribution of T^2

In general Hotline T^2 is defined as

$$T^2 = \frac{Y^T S^{-1} Y}{N-1} \quad \text{where } Y \sim N_p(\bar{Y}, I)$$

and

$$S = \frac{1}{N-1} \sum_{\alpha=1}^N Z_\alpha Z_\alpha^T; \quad Z_\alpha \sim N_p(0, I) \quad \text{and}$$

are independent for $\alpha = 1, 2, \dots, (N-1)$ i.e.

$$S \sim (N-1) \quad \text{or} \quad S = \sum_{\alpha=1}^N Z_\alpha Z_\alpha^T$$

Again,

$$T^2 = N(\bar{Y} - \mu_0)^T S^{-1} (\bar{Y} - \mu_0) \quad | \text{ under } H_0$$

let D be a NSM such that $D^T D = I$ — (1)

$$\text{Define: } \begin{aligned} Y^* &= DY \\ V^* &= D\bar{Y} \end{aligned} \quad | \quad - (2)$$

$$S^* = DSD^T$$

$$\text{then } Y^* \sim N_p(D\bar{Y}, DSD^T) \quad \text{or} \quad Y^* \sim N_p(V^*, I) \quad - (3)$$

Again in general : $T^2 = Y' S^{-1} Y$

$$= (D^{-1} Y^*)' S^{-1} (D^{-1} Y^*)$$

$$= Y^* (D S D')^{-1} Y^*$$

OR $T^2 = Y^* (S^*)^{-1} Y^* - (4)$

Also $\bar{Y}' \bar{S}^{-1} \bar{Y} = (D^{-1} V^*)' \bar{S}^{-1} (D^{-1} V^*)$

$$= (V^*)' (D \bar{S} D')^{-1} V^*$$

$$= V^* I V^* = V^* V^* - (5)$$

We see that $(N-1)S^* \sim \frac{1}{\lambda} \sum_{\alpha=1}^{N-1} Z_\alpha^* Z_\alpha^{*\top}$

where $Z_\alpha \sim D Z_\alpha \sim N_p(0, I)$ as $D \Sigma D' = I$
and each Z_α^* is indep for $\alpha = 1, 2, \dots, N$

Consider an orthogonal matrix $Q_{p \times p}$ s.t its
1st row is given by the element

$$q_{ij} = \frac{y_j^*}{\sqrt{y^* y^*}} \quad \text{for } j=1, 2, \dots, p \text{ where}$$

y^* is j th component of $y^* = \begin{pmatrix} y_1^* \\ \vdots \\ y_p^* \end{pmatrix}$

now if we can show that $\sum q_{ij}^2 = 1$ then
we shall be able to define such a matrix
we have

$$\sum_{j=1}^p q_{ij}^2 = \sum_{j=1}^p \frac{y_j^{*2}}{y^* y^*} = \frac{\sum_{j=1}^p y_j^* y_j^*}{y^* y^*} = \frac{y^* y^*}{y^* y^*} = 1$$

Since $\sum_{j=1}^p q_{ij}^2 = 1$

Hence the choice of the matrix of the orthogonal
matrix with q_{ij} as the elements in the
1st row is permissible

Now, we consider another transformation.

$$U_{px_1} = Q Y^* = \begin{pmatrix} U_1 \\ U_2 \\ \vdots \\ U_p \end{pmatrix}$$

$$\text{where } U_\alpha = \sum_{j=1}^p q_{\alpha j} Y_j^*$$

Again, since Q is an orthogonal matrix so

$$\sum_{j=1}^p q_{ij} q_{\alpha j} = 0 \text{ if } \alpha \neq 1$$

$$= 1 \text{ if } \alpha = 1$$

$$\text{Thus, } \frac{\sum_{j=1}^p Y_j^*}{\sqrt{Y^{*1} Y^*}} \cdot q_{\alpha j} = 0 \text{ if } \alpha \neq 1$$

$$\Rightarrow \sum_{j=1}^p Y_j^* q_{\alpha j} = 0 \text{ if } \alpha \neq 1$$

$$\Rightarrow U_\alpha = 0 \text{ if } \alpha \neq 1 \quad \text{--- (6)}$$

$$U_{px_1} = \begin{pmatrix} U_1 \\ U_2 \\ \vdots \\ U_p \end{pmatrix} \quad \text{Again we have.}$$

$$U_1 = \sum_{j=1}^p q_{1j} Y_j^*$$

$$= \frac{\sum_{j=1}^p Y_j^{*1} Y_j^*}{\sqrt{Y^{*1} Y^*}} = \frac{Y^{*1} Y^*}{\sqrt{Y^{*1} Y^*}}$$

$$\text{So } U_1 = \sqrt{Y^{*1} Y^*} \quad \text{--- (7)}$$

$$\therefore U = Q Y^* = \begin{pmatrix} \sqrt{Y^{*1} Y^*} \\ \vdots \\ 0 \end{pmatrix}$$

$$\text{Again we have } T^2 = Y^{*1} S^{*1} Y^*$$

$$= (Q^{-1} U)^* S^{*1} (Q^{-1} U)$$

$$= U^* (Q^{-1})^* S^* (Q^{-1}) U$$

$$\text{OR } T^2 = U^* (Q S^* Q^*)^{-1} U$$

Also $(N-1)S^* = A^*$ $\Rightarrow S^* = \underline{A^*}$
 $(N-1)$

$$T^2 = \underline{U}^1 \underline{Q} \underline{A^*} \underline{Q}^{-1} \underline{U}$$

$$\text{OR } T^2 = \underline{U}^1 B^{-1} \underline{U} \quad \text{--- (8)}$$

$$\text{where } B = \underline{Q} A^* \underline{Q}^{-1}$$

$$\therefore T^2 = \underline{U}^1 B^{-1} \underline{U}$$

$$\text{Now let } B = \begin{bmatrix} b_{11} & b_{12} \\ b_{21} & B_{22} \end{bmatrix}, B^{-1} = \begin{bmatrix} b^{(1)} & b^{(1)} \\ b^{(1)\top} & B^{22} \end{bmatrix}$$

As $BB^{-1} = I$ we have

$$B = \begin{bmatrix} b_{11} & b_{12} & \dots & b_{1p} \\ \vdots & & & \\ b_{21} & & & \\ b_{p1} & b_{p2} & \dots & b_{pp} \end{bmatrix}_{B_{22} (p-1) \times (p-1)}$$

$$I = \begin{bmatrix} 1 & 0 & 0 & \dots & 0 \\ 0 & 1 & 0 & \dots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & 0 & \dots & 1 \end{bmatrix}_{\text{Matrix}}$$

$$b_{11} b^{(1)} + b_{12} \underline{b^{(1)}} = 1 \rightarrow (A)$$

$$b_{11} b^{(1)} + b_{12} B^{22} = 0 \rightarrow (B)$$

$$b_{21} b^{(1)} + B_{22} \underline{b^{(1)}} = 0 \rightarrow (C)$$

$$b_{21} b^{(1)} + B_{22} B^{22} = I_{p-1} \rightarrow (D)$$

Using C we have

$$\underline{b^{(1)}} = -B_{22}^{-1} \underline{b_{21}} b^{(1)}$$

$$\underline{b^{(1)}} = -B_{22}^{-1} b_{21} b^{(1)}$$

Substituting this in A) we have

$$b_{11} b'' = b_{11} B_{22}^{-1} b_{11}' b'' = 1$$

OR $[b_{11} - b_{11} B_{22}^{-1} b_{11}'] b'' = 1$

$$\Rightarrow \boxed{b_{11,2} = 1}$$

where $b_{11,2} = b_{11} - b_{11} B_{22}^{-1} b_{11}'$

$$\frac{T^2}{(N-1)} = (U_1, 0, \dots, 0) \begin{bmatrix} b'' & b'' \\ b''' & B_{22} \end{bmatrix} \begin{pmatrix} U_1 \\ 0 \\ \vdots \\ 0 \end{pmatrix}$$

$$= U_1 b'' U_1$$

$$= U_1^2 b''$$

Again $\frac{T^2}{(N-1)} = U^T B^{-1} U$

$$= \begin{pmatrix} U_1 \\ 0 \\ \vdots \\ 0 \end{pmatrix} B^{-1} (U_1, 0, \dots, 0)$$

OR $\frac{T^2}{N-1} = \frac{Y^* Y^T}{b_{11,2}} \quad \text{--- (9)}$

as $b_{11,2} = 1/b''$

As $Y^* \sim N_p (Y', I)$ and its components are independent $(\alpha=1, 2, \dots)$. also if $Y^* \neq 0$ then

$$Y^* Y^T \sim \chi^2_{(s, p)}$$

where s is Non central Parameter and

$$s = \underline{\gamma^*} \underline{\gamma^*}$$

Thus the numerator of (9) has $\chi^2_{(s, p)}$ distⁿ

Now we wish to obtain distⁿ of denominator

For that we have $A^* = \sum_{\alpha=1}^{N-1} \underline{Z}_{\alpha}^{**} \underline{Z}_{\alpha}^{*}$

where each $\underline{Z}_{\alpha} \sim N_p(0, I)$ and are independent for $\alpha = 1, \dots, (N-1)$

matrix Q containing elements of \underline{Y}^* is a random matrix

$\therefore B = Q A^* Q'$ (also become random)

Hence we consider the conditional of B for given Q (i.e. \underline{Y}^*) conditionally.

$$B \sim \sum_{\alpha=1}^{N-1} Q \underline{Z}_{\alpha}^{*} \underline{Z}_{\alpha}^{*'} Q' = \sum_{\alpha=1}^{N-1} \underline{U}_{\alpha} \underline{U}_{\alpha}' \text{ (say)}$$

where $\underline{U}_{\alpha} = Q \underline{Z}_{\alpha}^{*}$ and $\underline{U}_{\alpha} \sim N_p(0, I)$ as $Q Q' = I$

and are independent

$$\text{Consider } A = \begin{bmatrix} A_{11} & | & A_{12} \\ \hline A_{21} & | & A_{22} \end{bmatrix} = \sum_{\alpha=1}^{N-1} \underline{Z}_{\alpha} \underline{Z}_{\alpha}'$$

$$\underline{Z}_{\alpha} \sim N_p(0, \frac{1}{2}I)$$

$$\text{then } A_{11.2} = A_{11} - A_{12} A_{22}^{-1} A_{21}$$

$$= \sum_{\alpha=1}^{(N-1)-(p-q)} \underline{w}_{\alpha} \underline{w}_{\alpha}' \quad \text{CONDITIONAL Distr.} \\ \text{OF } w_{11.2} \quad \text{OF } w_{12} \dots w_{(N-1)2}$$

$w_{11.2} \sim N_p(0, \frac{1}{2}I_{11.2})$ are independent $\alpha=1, 2, \dots, (N-1)$
applying this result the conditional distn of
 $w_{11.2}$ is $\sum_{\alpha=1}^{(N-1)-(p-q)} \underline{w}_{\alpha} \underline{w}_{\alpha}' \quad w_{\alpha} \sim N_p(0, \frac{1}{2}I)$

and are independent for $\alpha = 1, 2, \dots, (N-1)$
we have

$$I_p = \begin{bmatrix} 1 & 0 \\ 0 & I_{p-1} \end{bmatrix}$$

$$\therefore I_{11.2} = I_1 + I_{12} I_{p-1}^{-1} I_{12}' = 1$$

$\therefore w_\alpha \sim N(0,1)$

So the conditional distribution of

$$b_{11,2} = \sum_{\alpha=1}^{(N-p)} w_\alpha^2$$

where w_α^2 are Standard Normal variables and are independent $\alpha = 1, 2, \dots, (N-1)$ $\chi^2_{(N-1)}$ df which implies the conditional distributions of T^2 given Ω (or Y^*) is same as $\chi^2_{(N-1)}$ that of

$\chi^2_{(s,p)}$ as this does not involve Y^* or Ω matrix, so this become the unconditional dist' of T^2 also

$$\text{Thus } T^2 \sim \frac{\chi^2_{(s,p)}}{\chi^2_{(N-1)}}$$

OR

$$\left[\begin{array}{cc} T^2 & (N-p) \\ (N-1) & p \end{array} \right] \sim F_{p, (N-p), s}$$

$\rightarrow T^2$ - distribution

$$\text{Here } S = Y^* S^{-1} Y$$

if $s = 0$ T^2 becomes central distribution

Applications of T^2

① One Sample Problem:-

Suppose we are interested in testing the Hypothesis

$$H_0: \mu = \mu_0, \sigma^2 \text{ being unknown}$$

Let x_1, x_2, \dots, x_N be a random sample of size N from $N_p(\mu, \Sigma)$

Then the proposed Test statistics is

$$T^2 = N(\bar{X} - \mu_0)' S^{-1} (\bar{X} - \mu_0)$$

where $\bar{X} = \frac{1}{N} \sum_{i=1}^N x_i$ and $S = \frac{1}{(N-1)} \sum_{\alpha=1}^{N-1} (\bar{x}_\alpha - \bar{X})(\bar{x}_\alpha - \bar{X})'$

So,

$$\frac{T^2}{(N-p)} \sim F_{p, N-p; \alpha/2}$$

and ∴ the critical region is : $F_{\text{cal}} > F_{p, N-p; \alpha/2}$.

(2) Two Sample Problem $(H_0: \mu_1^{(1)} = \mu_2^{(1)})$

Suppose that we wish to test $\mu_1^{(1)} = \mu_2^{(1)}$ from 2 populations $N_p(\mu_1^{(1)}, \sigma_1)$ & $N_p(\mu_2^{(1)}, \sigma_2)$ as

$$\sigma_1 = \sigma_2 = \sigma \text{ say}$$

let $X_\alpha^{(1)} (\alpha = 1, 2, \dots, N_1)$ and $X_\beta^{(2)} (\beta = 1, 2, \dots, N_2)$ be two independent random sample from these popns respectively. We calculate

$$\bar{X}^{(1)} = \frac{1}{N_1} \sum_{\alpha=1}^{N_1} X_\alpha^{(1)} \quad \text{and} \quad \bar{X}^{(2)} = \frac{1}{N_2} \sum_{\beta=1}^{N_2} X_\beta^{(2)}$$

and

$$S_i = \frac{1}{(N_i-1)} \sum_{\alpha=1}^{N_i} (X_\alpha^{(i)} - \bar{X}^{(i)}) (X_\alpha^{(i)} - \bar{X}^{(i)})', \quad i=1, 2$$

also calculated the pooled Variance covariance matrix given by :

$$S = \frac{(N_1-1)S_1 + (N_2-1)S_2}{(N_1+N_2-2)}$$

We know that $\bar{X} \sim N_p(\bar{\mu}, \frac{S}{N})$. Using this

We have

$$\bar{X}^{(1)} \sim N_p(\mu^{(1)}, \frac{\Sigma}{N}), \quad \bar{X}^{(2)} \sim N_p(\mu^{(2)}, \frac{\Sigma}{N})$$

$$\Rightarrow \sqrt{\frac{N_1 N_2}{N_1 + N_2}} (\bar{X}^{(1)} - \bar{X}^{(2)}) \sim N_p(0, \frac{\Sigma}{N}) \text{ under } H_0$$

∴ the proposed test statistics is :

$$T^2 = \frac{N_1 N_2}{(N_1 + N_2)} (\bar{X}^{(1)} - \bar{X}^{(2)})' S^{-1} (\bar{X}^{(1)} - \bar{X}^{(2)})$$

$$\text{Or } T^2 = \frac{T^2}{(N_1 + N_2 - p)} \quad , \quad N_1 + N_2 - p \text{ f.d.} \sim F_{p, N_1 + N_2 - p - 3, \alpha}$$

③ q-Sample problem

Let $X_\alpha^{(i)}$ ($\alpha = 1, 2, \dots, N_i$), $i = 1, \dots, q$ be independent random sample from $N_p(\mu^{(i)}, \frac{\Sigma}{N_i})$ equality of

$\sum \beta_i \mu^{(i)}$ assumed and we wish to test $H_0: \sum \beta_i \mu^{(i)} = \mu_0$

$$H_0: \sum_{i=1}^q \beta_i \mu^{(i)} = \mu_0 \text{ (Specified) where } \beta_i \text{ are scalars}$$

and VCM as

$$\bar{X}^{(i)} = \frac{1}{N_i} \sum_{\alpha=1}^{N_i} X_\alpha^{(i)} \quad \text{and} \quad S_i = \frac{1}{(N_i - 1)} \sum_{\alpha=1}^{N_i} (\bar{X}_\alpha^{(i)} - \bar{X}^{(i)})' (\bar{X}_\alpha^{(i)} - \bar{X}^{(i)})'$$

and the pooled VCM

$$S = \frac{1}{\sum_{i=1}^q N_i - q} \sum_{i=1}^q \sum_{\alpha=1}^{N_i} (\bar{X}_\alpha^{(i)} - \bar{X}^{(i)})' (\bar{X}_\alpha^{(i)} - \bar{X}^{(i)})'$$

then for testing H_0 the proposed test statistic is:

$$T^2 = C \left[\sum_{i=1}^q \beta_i \bar{X}^{(i)} - \mu_0 \right]' S^{-1} \left[\sum_{i=1}^q \beta_i \bar{X}^{(i)} - \mu_0 \right] \sim F_{\sum_{i=1}^q N_i - q, (p-1) df}$$

So under H_0 :

$$\frac{T^2}{\sum_{i=1}^q N_i - q} \sim F_{p, \sum_{i=1}^q N_i - q - (p-1), \alpha}$$

(4)

Problem of Symmetry

Let $\mathbf{X} \in \mathbb{R}^{N \times p}$ ($N, p \in \mathbb{N}, p \geq 1$), & assumed to known. We may be interested in testing $H_0: \mu_1 = \mu_2 = \dots = \mu_p$ (ie all the components of the mean vectors are same). $H_0: (\mu_1 - \mu_p) = (\mu_2 - \mu_p) = \dots = (\mu_{p-1} - \mu_p) = 0$ which can be written in matrix form as:

$$\begin{bmatrix} 1 & 0 & 0 & \dots & 1 \\ 0 & 1 & 0 & \dots & 1 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & 0 & \dots & 1 \end{bmatrix} \mathbf{y} = 0$$

In order to test this H_0 we choose a matrix C such that it is of order $(p-1) \times p$ and it is of order rank $(p-1)$ and $C \cdot \mathbf{1} = 0$ where $\mathbf{1}_{p \times 1}$ is a unit vector.

$$C = \begin{bmatrix} C_{11} & C_{12} & \dots & C_{1p} \\ C_{21} & C_{22} & \dots & C_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ C_{p-1,1} & \dots & \dots & C_{p-1,p} \end{bmatrix} \quad \begin{bmatrix} 1 \\ 1 \\ \vdots \\ 1 \end{bmatrix} = 0$$

$$C \cdot \mathbf{1} = \begin{bmatrix} C_{11} + C_{12} + \dots + C_{1p} \\ C_{21} + C_{22} + \dots + C_{2p} \\ \vdots \\ C_{p-1,1} + \dots + C_{p-1,p} \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \\ \vdots \\ 0 \end{bmatrix} = \mathbf{0}$$

Consider the transformation $\mathbf{y} = C\mathbf{x}$ then

Since $\mathbf{x} \sim N_p(\mathbf{0}, \mathbf{I})$, $\mathbf{y} \sim N_p(C\mathbf{0}, \mathbf{C} \otimes \mathbf{C}^T)$

Suppose $\mathbf{y} \sim N_p(\boldsymbol{\eta}, \boldsymbol{\Phi})$ where $\boldsymbol{\eta} = C\mathbf{0}$.

Then $\mathbf{E}\mathbf{y} = \boldsymbol{\eta}$ and $\boldsymbol{\Phi} = \mathbf{C} \otimes \mathbf{C}^T$

If H_0 is true all the components of \mathbf{y} will be equal and then

$\eta = C Y = C \begin{bmatrix} u_1 \\ u_2 \\ \vdots \\ u_p \end{bmatrix} = 0$ is $H_0 \rightarrow H_0^1 : \eta = 0$
 in this case the proposed test statistics is

$T^2 = N \bar{Y}' S^{-1} \bar{Y}$ and then we can proceed with as usual testing for one sample problem.

Ex: Let $X = \begin{pmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \\ x_5 \end{pmatrix} \rightarrow \begin{array}{l} \text{Score in Physics} \\ \text{Maths} \\ \text{English} \\ \text{Chemistry} \\ \text{Social science} \end{array}$

Suppose we wish to test the hypothesis that the student is good in all the subjects then we can have $E(X) = \mu$ and we can propose H_0 as $\mu = 0$ and again we can test the hypothesis.

(*) Properties of T^2 .

T^2 -test is (i) UMP (ii) UMPU - Unbiased
 (iii) UMPS - Similar (iv) UMPI - Invariance

T^2 -test has a monotonic power function.

(*) Optimum property

T^2 is unaffected by change of origin and scale of origin and scale of the response variable. It is Invariant w.r.t. all the transformation of type

$$Y = CX + d$$

where C is NSM and d is non-random vector of component

proof we have $\bar{X} \sim N_p(\mu, \Sigma)$.

Consider $Y = C\bar{X} + d$ and suppose we wish to test $H_0: \mu = \mu_0$ then $E(Y) = C\mu + d = C\mu_0 + d$ and suppose we wish to test

$$H_0: \mu = \mu_0 \text{ then } E(Y) = C\mu + d \\ = C\mu_0 + d \text{ (under } H_0)$$

for this we have $T_x^2 = N(\bar{X} - \mu_0)' S_x^{-1} (\bar{X} - \mu_0)$

$$\text{Now } T_y^2 = N(\bar{Y} - \mu_0)' S_y^{-1} (\bar{Y} - \mu_0)$$

OR

$$\begin{aligned} T_y^2 &= N[(C\bar{X} + d - C\mu_0 - d)' (S_y^{-1}) (C\bar{X} + d - C\mu_0 - d)] \\ &= N[C(C\bar{X} + d - C\mu_0 - d)' (S_x^{-1} C') (C\bar{X} + d - C\mu_0 - d)] \\ &= N[(\bar{X} - \mu_0)' C' (S_x^{-1} C') (\bar{X} - \mu_0)] \\ &= N[(\bar{X} - \mu_0)' C' (C')^{-1} S_x^{-1} (C^{-1}) C (\bar{X} - \mu_0)] \\ &= N[(\bar{X} - \mu_0)' S_x^{-1} (\bar{X} - \mu_0)] \end{aligned}$$

$$\therefore T_y^2 = T_x^2$$

$\therefore T_x^2$ is invariant under a group of transformation.



Repeated Measurement (Repeated)

Frequently the experimental observations are collected at different times but on the same sample units. The simplest example of this repeated measurement design consists of a controlled treatment (a placebo, standard drug or any other null treatment) on one occasion and the treatment of interest on other occasions.

It is reasonable to assume that the treatment has an additive effect than we may write

the responses of the i th subject under control and treatment conditions as:-

$$x_{i1} = \mu + c_{i1} \quad (1)$$

$$x_{i2} = \mu + T + c_{i2}$$

where μ is the general mean effect common to all the subjects and T is the treatment effect of the experimental condition.

And c_{i1}, c_{i2} , are the random disturbance terms.

Let us extend the paired observation case to the general setup of ' k ' (say) responses collected on the same experimental unit at successive times or under a variety of experimental conditions.

We shall assume that these ' k ' responses constitute a fixed set of observations. Let the observations on N independent units under ' k ' conditions are arranged as:

Sample Units	Condition	Observations from a repeated Measurements Experiment.
1	$x_{11} x_{12} x_{13} \dots x_{1k}$	
2	$x_{21} x_{22} \dots x_{2k}$	
3	$x_{31} x_{32} \dots x_{3k}$	
\vdots	\vdots	
N	$x_{N1} x_{N2} \dots x_{Nk}$	

The mathematical model for the j th response in the i th sampling unit may be written as:

$$x_{ij} = \mu + \alpha_j + c_{ij}, \quad i = 1, 2, \dots, N$$

$$j = 1, 2, \dots, k$$

-(1)

where μ_j : general mean effect (or general level of the jth common to all observations)

α_{ij} : effects of jth specification and
 e_{ij} : random disturbance

e_{ij} reflects both the interaction of the ith unit with the jth response and the experimental error.

To obtain the test of hypothesis for μ_j , it will be necessary to assume that the error vector $(e_{i1}, e_{i2}, \dots, e_{ik}) = (0, 0, \dots, 0) = \underline{0}$ should be a null vector and the covariance matrix

$$\Sigma = E \begin{bmatrix} (c_{11})(e_{i1} e_{i2} \dots e_{ik}) \\ c_{12} \\ \vdots \\ c_{ik} \end{bmatrix} \quad - (3)$$

write :-

$$y = [\mu_1 + e_1, \mu_2 + e_2, \dots, \mu_k + e_k] \quad - (4)$$

then the null hypothesis of equal responses effect is :-

$$H_0: \mu_1 = \mu_2 = \dots = \mu_k$$

vs $H_1: \mu_i \neq \mu_j$ for at least one pair of i and j

the vector statements of H_0 and H_1 are :-

$$\begin{aligned}
 H_0: & (\mu_1 - \mu_2) = (0) \quad \text{vs} \quad H_1: (\mu_1 - \mu_2) \neq 0 \\
 \text{thus for } & \begin{pmatrix} \mu_1 - \mu_3 \\ \mu_2 - \mu_3 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix} \neq 0 \quad \text{and} \quad \begin{pmatrix} \mu_1 - \mu_3 \\ \mu_2 - \mu_3 \\ \vdots \\ \mu_{k-1} - \mu_k \end{pmatrix} \neq 0
 \end{aligned}
 \quad - (5)$$

The test of H_0 stated above can be now carried out with the help of Hotunes T^2 statistics computed from the mean vector and VCM of the differences.

$$Y_{ij} = (x_{ij} - x_{i,j+1}) \quad \text{--- (6)}$$

if the observations are adjacent responses. If we denote \bar{Y}' the $(k-1)$ component vector of differences of successive response means :-

$$\bar{Y}' = [\bar{y}_1, \bar{y}_2, \dots, \bar{y}_{k-1}] = [\bar{x}_1 - \bar{x}_2, \dots, \bar{x}_{k-1} - \bar{x}_k] \quad \text{--- (7)}$$

and let $S_{YY'}$ denotes the covariance matrix of the Y_{ij} difference then the test statistics for testing H_0 of equal response effect is :-

$$\text{test statistic } T^2 = N \bar{Y}' S_{YY'}^{-1} \bar{Y}' \text{ and under } H_0 \text{ it is}$$

$$\frac{T^2 (N-k+1)}{(N-1)(k-1)} \sim F_{(k-1), (N-k+1)}, \alpha$$

then H_0 may be rejected if $F_{\text{cal}} > F_{\text{tab}}, \alpha$.

or H_0 may be accepted.

3) Multiple Comparisons:

If the hypothesis of equal treatment effects is rejected consider any linear combination

$$\sum_{j=1}^k b_j Y_j \quad \text{--- (1)} \quad \text{of the treat effects}$$

When the coefficients b_j 's sum to zero

We have to choose a matrix c so that the element generate the contrast of adjacent treatments. Now let a' be any $(k-1)$ component row vector then by the UNION-INTERSECTION principle it follows that all the inequalities like :-

$$[\underline{a}' c (\bar{x} - \bar{y})]^2 \leq \frac{1}{N} [\underline{a}' c S_{YY'}^{-1} \underline{a}]^2 \quad (k-1)(N-k+1), \alpha$$

formed for the different choices of the elements of α held simultaneously with prob $(1-\alpha)$. But since the row sums of $\alpha'c$ of each row is zero, so must the sum of elements of $\alpha'c$ be zero. $b'y = \alpha'c$ \Rightarrow $b'y = 0$. This is in contrast of the treatment effects.

So, here $b'y = 0$.

$$b'y = \sum_{j=1}^k b_j (y_1 + y_2) = \sum_{j=1}^k b_j d_{1j}$$

Expansion of (*) leads to a set of $100(1-\alpha)\%$

terms simultaneously C.I.s whose general member is :

$$\text{given } b'\bar{x}_i - \frac{b'Sb}{N} T^2 \leq b'y \leq b'\bar{x}_i + \frac{b'Sb}{N} T^2$$

If the interval includes zero, the hyp.

$\sum b_j d_{1j} = 0$, is acceptable in the sense of multiple comparison at $\alpha\%$ level.

Profile Analysis

Profile analysis is an equivalent of RM MANOVA mainly concerned with the test scores obtained by a group of persons from educational and psychological assessment.

Profile also refers to an outline or a short description of someone's life or achievements.

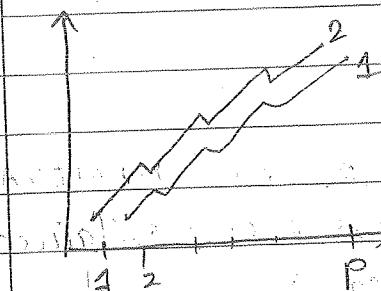
In a battery of say 'p' psychological tests administered to a group. Let y_1, y_2, \dots, y_p be the expected scores in these 'p' tests, then the PROFILE of the group is defined as the GRAPH obtained by joining the points $(1, y_1), (2, y_2), \dots, (p, y_p)$ successively.

In practice, the μ_i 's value were estimated from the average scores in a sample.

If we have two possibly different groups we will have two profiles, one determined by joining the points (i, \bar{x}_i) $i=1, 2, \dots, p$ and the other by (i, \bar{s}_i) , $i=1, 2, \dots, p$ where $\bar{s}_1, \bar{s}_2, \dots, \bar{s}_p$ are the true mean scores of the other groups.

(i) The profiles are set to be similar if the segment of line joining $(1, \bar{x}_1)$ to $(i+1, \bar{x}_{i+1})$ is PARALLEL to the corresponding segment of the line joining $(1, \bar{s}_1)$ to $(i+1, \bar{s}_{i+1})$ of the other profiles each.

Obviously this will be so if the hypothesis $H_0: (\bar{x}_{i+1} - \bar{x}_i) = (\bar{s}_{i+1} - \bar{s}_i)$ for $i=1, 2, \dots, p$ is TRUE



Let us assume that the test scores x_{pxi} are normally distributed and their VCM is same after \rightarrow both the group.

Let X be $(p \times n_1)$ matrix of scores of n_1 individual from the first group and similarly Y be $(p \times n_2)$ matrix of scores of n_2 individual from the 2nd group. Then the matrices X and Y represent the samples of the sizes n_1 and n_2 from $N_p(\mu_1, \Sigma)$ and $N_p(\mu_2, \Sigma)$ resp.

The hypothesis H_0 can be re-written as $H_0: E(CC^T\bar{X})$ is same in both the groups where C is the contrast matrix given by

$$C_1 = \begin{bmatrix} 1 & -1 & 0 & \dots & 0 \end{bmatrix} \quad \text{Note that } C_1 \bar{x} \text{ has}$$

the distribution $N_{p-1}(C_1 \bar{x}, E_1 \Sigma G^1)$ in

samples from the 1st population in the 1st group and

$N_{p-1}(G_2 \bar{x}, G_2 \Sigma G^1)$ in the 2nd group and G_1, G_2

haven't the samples these divide

The problem of testing of H_0 is same as the test of equality of two mean vectors from two multivariate Normal dist' with the same unknown Σ vct', then the test is provided by the following test statistics:

TEST STATISTICS:

$$\frac{(n_1 + n_2) - (p-1) - 1}{(p-1)} \cdot \frac{n_1 n_2}{(n_1 + n_2)(n_1 + n_2 - 2)} \cdot D^2 \sim F_{p-1, n_1 + n_2 - (p-1) - 1, \alpha}$$

$$D^2 = (n_1 + n_2 - 2) (G_1 d) (G_2 S G^1)^{-1} (G_1 d)$$

$d = (\bar{x}_1 - \bar{x}_2)$, \bar{x}_1 : Sample mean vector of population 1

\bar{x}_2 : Sample mean vector of population 2

$$S = X \left[I - \frac{1}{n_1} E_{n_1, n_1} \right] X^T + Y \left[I - \frac{1}{n_2} E_{n_2, n_2} \right] Y^T$$

where: $E_{p,q}$ ($p \times q$) matrix of all the elements being unit

If we decide to accept the hypothesis of similarity of two profiles on the basis of the test given by the \star , the next logical question will be whether the average true scores of the two group is same i.e. whether

$$H_0: \Sigma (d_1 + d_2 + \dots + d_p) = \Sigma (\delta_1 + \delta_2 + \dots + \delta_p) \text{ OR } E_{\text{Ip}} \Sigma = E_{\text{Ip}} \delta$$

Observe that $E_{\text{Ip}} \Sigma$ has a univariate normal distribution with variance $E_{\text{Ip}} \Sigma^2 - E_{\text{Ip}} \Sigma$ and the means $E_{\text{Ip}} \delta$, $E_{\text{Ip}} \Sigma$ is the two groups and hence the hypo H_1 can be tested by the usual t-test for testing the equality of the mean of two normal populations with the same variance.

If we square t for testing H_1 we will get the F-distribution or F-test statistics.

If H_0 is true, and if in addition H_1 is also true then we would like to specify that $d_1 = \delta_1, d_2 = \delta_2, \dots, d_p = \delta_p$. i.e. BOTH THE PROFILES ARE EXACTLY SAME.

If H_0, H_1 and H_2 all hold then the distribution of \bar{X} is same in both the groups and either X and Y are same samples from the same group so now we can POOL these giving the GRAND MEAN VECTOR :

$$\bar{g} = \frac{n_1 \bar{x} + n_2 \bar{x}}{(n_1 + n_2)} \text{ is an estimate}$$

of the common value of μ and δ and it's dist' is $N_p(\bar{g}, \mathbf{I}/N)$ where $N = (n_1 + n_2)$

The hypo H_2 is equivalent to $E(Cg) = 0$, where C is the contrast matrix and so the test reduce to the test of sign of contrasts which can be accomplished by

$$\frac{(n_1 + n_2 - 2) - (p-1) - 1}{(p-1)} \sim \chi^2_{n_1 + n_2 - p - 1}$$

$$T^2 = N(c_g)^T (csc)^{-1} (c_g)$$

Note: If H_0 is accepted, the common profile of two groups consists of parallel lines to the x -axis at a distance equal to the common value of μ and γ .

Notations:

- $\bar{X} = \bar{X}_{E_{n_1}} / n_1 \rightarrow$ Vector of sample means for 1st Sample
- $\bar{Y} = \bar{Y}_{E_{n_2}} / n_2 \rightarrow$ vector of sample means for 2nd Sample
- $S_x = X \left(I - \left(\frac{1}{n} \right) E_{n, n_2} \right) X^T$: Matrix of corrected S.S and of SP of the observations from the 1st sample
- $S_y = Y \left[I - \left(\frac{1}{n} \right) E_{n, n_2} \right] Y^T$: Matrix of corrected S.S and of SP of the observations from the 2nd sample

Also

$$\bar{X} \sim N_p(\mu, \Sigma_1; \frac{1}{n}) \quad S_x \sim W_p(n_1 - 1; \frac{1}{n})$$

$$\bar{Y} \sim N_p(\gamma, \Sigma_2; \frac{1}{n}) \quad S_y \sim W_p(n_2 - 1; \frac{1}{n})$$

* Fisher's Between's Problem:-

Suppose we wish to test $H_0: \mu = \gamma$ — (1)
 where μ is the mean vector of $N_p(\mu, \Sigma_1)$ and
 γ is the mean vector of $N_p(\gamma, \Sigma_2)$, Σ_1 and Σ_2
 are unknown and unequal

Consider that we have random samples of sizes n_1 and n_2 from the two multivariate normal population and let X and Y be (pxn_1) and (pxn_2) matrices of sample observations.

We also assume that $n_1 < n_2$ and now partition the Y matrix as

$$Y = \begin{bmatrix} Y_1 \\ Y_2 \end{bmatrix} \quad (2) \quad \text{where } m = (n_2 - n_1)$$

Note that $E(X) = \Sigma E_{1,n_1}$, $\text{Var}(X) = \Sigma (X) I_n$ — (3)

$A \otimes B$: represent the kronecker product of A & B
Similarly,

$$E(Y_1) = \Sigma E_{1,n_1}, \quad E(Y_2) = \Sigma E_m \quad (4)$$

$$\text{var}(Y_1) = \Sigma (Y_1) I_{n_1}, \quad \text{var}(Y_2) = \Sigma (Y_2) I_m. \quad (5)$$

further, $\text{Cov}(Y_1, Y_2) = 0 \quad (6)$

As all the col of Y are indep define a (pxn_1) matrix Z as:

$$Z = X - \left(\frac{n_1}{n_2} \right) Y_2 \left(I - \frac{1}{n_1} E_{n_1, n_1} \right) - \frac{1}{n_2} Y E_{n_2, n_1} \quad (7)$$

It can be observed from (3) & (4) that

$$E(Z) = (\Sigma - \Sigma) E_{1, n_1} \quad (8)$$

Or,

$$Z = X - Y_1 P - \frac{1}{n_2} Y_2 E_{m, n_1} \quad (9)$$

where

$$P = \left(\frac{n_1}{n_2} \right) Y_2 \left(I - \frac{1}{n_1} E_{n_1, n_1} \right) + \frac{1}{n_2} E_{n_1, n_1} \quad (10)$$

using (2) Z can be written as

$$Z = X - \left(\frac{n_1}{n_2} \right) Y_2 Y_1 \left(I - \frac{1}{n_1} E_{n_1, n_1} \right) - \frac{1}{n_2} (Y_1 E_{n_1, n_1} + Y_2 E_{m, n_1})$$

then using (3), (4) and (5) we have obtain:-

$$\text{Var}(z) = \sum_1 (x) I_{n_1} + \sum_2 (x) (P' E_{n_1, p}) + \frac{1}{n_2} (E_{n_1, n_1} - E_{m_1, m_1})$$

$$= (\sum_1 (x) I_{n_1} + \sum_2 (x) \left(P' P + \frac{m_1}{n_2} E_{n_1, n_1} \right))$$

$$= \left(\sum_1 + \frac{n_1 \cdot \sum_2}{n_2} \right) (x) I_{n_1} \quad (1)$$

OR

$$\text{Var}(z) = \frac{1}{2} (x) I_{n_1} \text{ where } \frac{1}{2} = \sum_1 + \left(\frac{n_1}{n_2} \right) \sum_2 \quad (12)$$

$\frac{1}{2}$ is non singular also z is a linear combination of the variables from x and y , it is observed that each column of z will have a multivariate normal distribution with mean $(\bar{x} - \bar{y})$ and a var-cov matrix $\frac{1}{2}$ given by eq (12)

Further, as all the cols of z are independent then if the hypothesis $H_0: \bar{y} = \bar{x}$ of the z will be testing that $z \perp \perp 0$ (null vector)

which can be again tested by using hotline's T^2 statistic given by

$$T^2 = n_1 \bar{z}' S_z^{-1} \bar{z} \text{ where } \bar{z} = \frac{1}{n_1} \sum z E_{n_1, 1} \quad (13)$$

mean vector of obs in z and $S_z = z (I - \frac{1}{n_1} E_{n_1, n_1}) z'$ — (14)

∴ the proposed test statistics is

$$\frac{(n_1 - p)}{p} \frac{T^2}{(n_1 - 1)} \sim F_{p, n_1 - p}, \text{ i.e., } \chi^2_{p, n_1 - p}$$

Note that $\bar{z} = (\bar{x} - \bar{y})$ and

$$S_z = V \left(I - \frac{1}{n_1} E_{n_1, n_1} \right) V' \text{ where } V = \begin{pmatrix} n_1 \\ n_2 \end{pmatrix}^{1/2} Y_1 \quad (15)$$

it is easier to form the matrix U by subtracting from every col of X , $(\frac{n_1}{n_2})Y_2$ times the corresponding col. of Y . Then we may write:

$$U^T U = n_1 (X - Y)' S_0^{-1} (X - Y) \quad (16)$$

$(n_1, -1)$

it is to be noted that all the n_2 observations are utilized for finding out Y but only n_1 observations are utilized in S_0 so which of the $(n_2 - n_1)$ columns of Y are ignored is arbitrary in this test and therefore this is the weakness of this test.

* Test of Symmetry of ORGANs

In many biological and anthropological problem let us consider the measurement characteristics on the LEFT side are $x_{11}, x_{12}, \dots, x_{1p}$ and x_{p+1}, \dots, x_{2p} represent those same measurements on the right side. or $x_{i1} (i=1, \dots, p)$ represents the characteristics of a twin pair and $x_{pi} (i=1, \dots, p)$ are those for the other members and we are interested in testing the hypothesis of symmetry of the left and right sides or the equality of the measurement of the twins.

Mathematically the hypothesis to be tested can be expressed as

$$H_0: \mu_1 = \mu_{p+1} \quad (i=1, \dots, p) \quad (1)$$

where $\mu_1, \mu_2, \dots, \mu_{2p}$ are the true means of x_{11}, \dots, x_{2p} .

Let \bar{X} and S represent the mean vector and the matrix of corrected SS and S.P. of 'm' samples observation from this $2p$ -variate normal pop.

Now, consider the partitioning of \bar{X} and S as

$$\bar{X} = \begin{bmatrix} \bar{X}^{(1)} \\ \bar{X}^{(2)} \end{bmatrix} \xrightarrow{p} \text{ and } S = \begin{bmatrix} S_1 & S_2 \\ S_2 & S_3 \end{bmatrix} \xrightarrow{p} \begin{bmatrix} p \\ p \end{bmatrix} \quad (2)$$

then we know that

$$\sqrt{n} \bar{X} \sim N_p(\sqrt{n}\mu, \Sigma) \text{ and } S \sim W_{2p}(n-1, \Sigma)$$

where $\mu = \begin{bmatrix} \mu_1 \\ \mu_2 \\ \vdots \\ \mu_{2p} \end{bmatrix}$ and Σ is $(2p \times 2p)$ VCM of $2p$ -variate

$$\therefore Y = \sqrt{n}(\bar{X}^{(1)} - \bar{X}^{(2)}) = \sqrt{n} A \bar{X} \quad (3)$$

that $N_p\left[\sqrt{n}(\bar{X}^{(1)} - \bar{X}^{(2)}), A \Sigma A'\right]$ where $\bar{\mu} = \begin{bmatrix} \bar{\mu}^{(1)} \\ \bar{\mu}^{(2)} \end{bmatrix} \xrightarrow{p}$

then the dist' of $A \Sigma A' = (S_1 - S_2' - S_2 + S_3)$

is $W_p(A \Sigma A', n-1)$ — (4)

and it is independent of Y .

Hence, if H_0 is true i.e. $\bar{\mu}^{(1)} = \bar{\mu}^{(2)}$, then

$$\frac{(n-1)-(p-1)}{p} \cdot \frac{T^2}{(n-1)} \sim F_{p, (n-p)}, \alpha, v.$$

$$\text{where } \frac{T^2}{(n-1)} = \bar{Y}' (A \Sigma A')^{-1} \bar{Y}$$

$$\text{OR, } \frac{T^2}{(n-1)} = n (\bar{X}^{(1)} - \bar{X}^{(2)})' (S_1 - S_2' - S_2 + S_3)^{-1} (\bar{X}^{(1)} - \bar{X}^{(2)})$$

* The POWER of TESTS ON MEAN VECTOR(S) :-

So far we have been concerned with the size of the test, however in a properly designed experiment the rate of type -2 error should also be controlled.

We know $T^2 = \mathbf{Y}' \mathbf{S}^{-1} \mathbf{Y} - \mathbf{O}$ is calculated from the p dimensional which has a multivariate normal distⁿ $N_p(\mathbf{y}, \mathbf{I})$ and we have also estimated the parameters μ and σ . Based on n degrees of freedom then we know that quantity

$f = T^2 / (n-p+1)$, has a NON-CENTRAL F-distⁿ with d.f. $(p, n-p+1)$ and ncp $\delta^2 = \mathbf{y}' \mathbf{I}^{-1} \mathbf{y}$
in particular if we are testing

$$H_0: \mu = \mu_0 \text{ vs } H_1: \mu \neq \mu_0$$

then the T^2 -statistics, then ncp $\delta^2 = N(0)$
 $n \text{cp } \delta^2 = N(\mathbf{y} - \mathbf{y}_0)' \mathbf{I}^{-1} (\mathbf{y} - \mathbf{y}_0)$

Again, for the two sample problem, Testing
 $H_0: \mu^{(1)} = \mu^{(2)}$ vs $H_1: \mu^{(1)} \neq \mu^{(2)}$

the test statistic has a NON-CENTRAL F-dist with p.m

$$\frac{N_1 N_2}{(N_1 + N_2)} (\mu^{(1)} - \mu^{(2)})' \mathbf{I}^{-1} (\mu^{(1)} - \mu^{(2)})$$

and the df $p, N_1 + N_2 - (p - 1)$

In either case the POWER function of the test is

$$1 - \beta(\delta^2) = P[F_p > f_{p, n-p+1; \alpha/2}]$$

where F' has non-central F -distribution with prescribed parameters

PEARSON and HEARTLY have prepared the power charts of the power function for $\alpha = 1\%, 2.5\%$.

References:-

- (1) Page - 149 (MORRISON Book)
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PRINCIPLE COMPONENT ANALYSIS

INTRODUCTION

In Multivariate studies if the number of variables on which the data is collected is too large and we are interested in only few of them, then a way of reducing a number of variable is to consider those linear combination of the variables which have similar or smaller variance.

PCA is a technique to help in this reduction as it tries to find few linear combinations of the original variables which can be used to summarize the data without losing much information.

It is a technique to transform the original variables in such a way that most of the variation is explained by these linear combinations. The linear combination is so selected are called the principle components and the major objective of PCA is to reduce the dimension of the data into few components.

A PCA is concerned with explaining the variance covariance structure of the given multivariate data.

It is important to note that the total variation remains the same even after the transformation of original variables into principle components.



Mathematics of PCA

Let $X = \begin{pmatrix} X_1 \\ \vdots \\ X_p \end{pmatrix}$ be a random vector of p -component with its dispersion matrix Σ , Since we are in the var-cov structure, we assume that $E(X) = 0$. Let β be a $(p \times 1)$ comp. vector s.t. $\beta^T \beta = 1$ and let $U = \beta^T X$ be a linear combination of the components of X we wish to find out β in such a way that $\text{var}(\beta^T X)$ is maximum. i.e

$$\text{Var}(U) = \text{Var}(\beta^T X) = \beta^T \text{Var}(X) \beta = \beta^T \Sigma \beta \quad (1)$$

where β is a Normal vector defines a fⁿ:

$$\Phi = \beta^T \Sigma \beta - \lambda (\beta^T \beta - 1);$$

λ is the language's multiplier now,

$$\frac{\partial \Phi}{\partial \beta} = 2 \Sigma \beta - 2 \lambda \beta = 0 \quad (\text{equating it to zero})$$

$$\Rightarrow \Sigma \beta - \lambda \beta = 0 \quad \text{i.e } (\Sigma - \lambda I) \beta = 0$$

Since $\beta \neq 0$

for a non-zero solution of the above equations we must have $(\Sigma - \lambda I) = 0 \quad (2)$

Thus, λ is the ch. root of Σ , as Σ is a $(p \times p)$ matrix, $(\Sigma - \lambda I) = 0$, will be a polynomial of order 'p' which will give 'p' ch. roots for Σ (say), $\lambda_1, \lambda_2, \dots, \lambda_p$.

Now, let β_i be the ch. vector of Σ corresponding to the ch. root λ_i ,

∴ it must satisfy the original equations and we must have

$$\beta_i^T \beta_i - \lambda_i \beta_i = 0,$$

now premultiplying it by β_i , we get

$$\beta_i^T \beta_i - \beta_i^T \lambda_i \beta_i = 0 \text{ or } \beta_i^T \beta_i - \lambda_i = 0$$

$$\text{as } \beta_i^T \beta_i = 1 \Rightarrow \beta_i^T \beta_i = 1$$

$$\therefore \beta_i^T \beta_i = \lambda_i \text{ or, } \text{Var}(\beta_i^T \mathbf{x}) = \lambda_i,$$

We wish to maximize $\text{Var}(\beta_i^T \mathbf{x})$

so, we should choose λ_i to be Maximum

now, arranging the 'p' roots as: $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_p$

If β_1 is the ch. vector corresponding to λ_1 (largest ch. root), then $(\beta_1^T \mathbf{x})$ is the

NORMALIZED linear combination of \mathbf{x} with this max. vari = λ_1

$$\text{So, } U_1 = \beta_1^T \mathbf{x} \rightarrow \text{1st PC (Principal component)}$$

for Second principle component again we find out in such a way that $(\beta_2^T \mathbf{x})$. In such a way that

$$\text{Var}(\beta_2^T \mathbf{x}) = \beta_2^T \beta_2 \text{ is maximum}$$

Subjected to the conditions that

$$(i) \beta_2^T \beta_2 = 1 \text{ and}$$

(ii) $\beta_2^T \mathbf{x}$ is uncorrelated to U_1 i.e.

$$\text{cov}(U_1, \beta_2^T \mathbf{x}) = 0 \text{ or } \beta_2^T \beta_1 = 0$$

As β_1 is vector corresponding to λ_1

$\therefore \beta_1 = \lambda_1 \beta_1$, now premultiplying it by β_1' ,
we have $\beta_1' \beta_1 = \beta_1' \lambda_1 \beta_1 \Rightarrow \lambda_1 \beta_1' \beta_1 = \beta_1' \beta_1 = 0$ under (ii)

$$\beta_1' \beta_1 = \beta_1' \lambda_1 \beta_1 \quad \text{i.e. } \lambda_1 \beta_1' \beta_1 = \beta_1' \beta_1 = 0$$

$$\text{i.e. } \lambda_1 \beta_1' \beta_1 = 0 \text{ but } \lambda_1 \neq 0$$

$$\therefore \beta_1' \beta_1 = 0,$$

$$\text{Again redefine } \Phi = \beta' \beta - \lambda (\beta' \beta - 1) - \mu (\beta' \beta_1 - 0)$$

where λ and μ are the Lagrange's multipliers.

Now,

$$\frac{\partial \Phi}{\partial \beta} = 2 \beta' \beta - 2 \lambda \beta - \mu \beta_1 = 0,$$

premultiplying it by β_1' , we have

$$2 \beta_1' \beta - 2 \beta_1' \lambda \beta - \mu \beta_1' \beta_1 = 0$$

$$\text{Now, } \text{cov}(\beta_1' \beta, \beta_1' \beta) = \beta_1' \beta_1 = \beta_1' \beta = 0$$

$$\therefore \frac{\partial \Phi}{\partial \beta} = 2 \beta' \beta - 2 \lambda \beta = 0 \text{ or, } (\Sigma - \lambda I) \beta = 0$$

$$\text{again, } \beta' \beta = 1 \quad \therefore \beta \neq 0 \quad \therefore (\Sigma - \lambda I) = 0$$

$\therefore \lambda$ is the ch. root of Σ and β is the corresponding ch. vector. Now, to maximize var($\beta' \underline{x}$) we choose 2nd largest ch. root of Σ which is the λ_2 and let β_2 be the corresponding ch. vector then

$\boxed{U_2 = \beta_2' \underline{x}}$ is the SECOND PRINCIPLE COMPONENT
which is uncorrelated with U_1

Reiner. It is to be noted that $\beta_1 \perp \beta_2$ satisfied $\beta_2^T \beta_1 = 0$ and $\beta_2^T \beta_2 = 0$ which $\Rightarrow \beta_2$ and β_1 are orthogonal. Now, let $\beta = [\beta_1, \beta_2, \dots, \beta_p]$ then β represent a $(p \times p)$ orthogonal matrix.

i.e $\beta^T \beta = \text{diag}(\lambda_1, \lambda_2, \dots, \lambda_p)$. Proceeding like this at the $(i+1)^{\text{th}}$ steps, we have to determine β_i s.t. $\text{var}(\beta_i^T x) = \beta_i^T \beta_i$ is maximum subjected to conditions.

$$(i) \beta_i^T \beta = 1 \text{ and}$$

(ii) $\beta_i^T x$ is uncorrelated with 1st, 2nd, ..., ith PC's

Now, let the ith P.C $v_i = \beta_i^T x$ and β_i the ch vectors corresponding to λ_i and (ii) implies that $\text{cov}(\beta_i^T x, \beta_j^T x) = 0$ for $j=1, 2, \dots, i$

i.e $\beta_i^T \beta_j = 0$ for $j=1, 2, \dots, i$ also for $j < i$ the original equation must be satisfied

$$\text{i.e } \beta_i^T \beta_j = \lambda_j \beta_j$$

now, premultiplying it by β_i^T , we get

$$\beta_i^T \beta_j = \beta_i^T \lambda_j \beta_j = \lambda_j \beta_i^T \beta_j \Rightarrow \beta_i^T \beta_j = 0, j=1, 2, \dots, i$$

as $\beta_i^T \beta_j = 0$ for $j < i$

Now redefine a function

$$\Phi = \beta^T \Sigma \beta - \lambda (\beta^T \beta - 1) = \sum_{j=1}^i \lambda_j \beta_j^T \Sigma \beta_j$$

$$\therefore \frac{\partial \Phi}{\partial \beta_j} = 2 \sum \beta_j - 2 \lambda_j \beta_j - \sum_{j=1}^i \lambda_j \beta_j^T \Sigma \beta_j$$

again premultiplying it by β_i^T ($j \leq i$), we get

$$2 \beta_i^T \Sigma \beta_j - 2 \lambda_j \beta_i^T \beta_j - \sum_{j=1}^i \lambda_j \beta_j^T \Sigma \beta_j$$

$$\Rightarrow \lambda_j = 0 \text{ for } j \leq i$$

∴ the equation reduces to the form

$$\sum \beta_i - \lambda \beta = 0 \quad \text{or} \quad (\Sigma - \lambda I) \beta = 0 \Rightarrow |\Sigma - \lambda I| = 0$$

Then, we choose $\lambda_{i+1}^{\text{th}}$ ch. roots and let β_{i+1} be the $(i+1)^{\text{th}}$ corresponding ch. vector. Then,

$U_{i+1} = \beta_{i+1}^T X$ $(i+1)^{\text{th}}$ P.C. with which is independent of all the other i P.C.s

* MLE's of Principle Components

For a given multivariate data involving the var-cov matrix Σ , it may happen that Σ is unknown. Hence we proceed with an estimate of Σ .

Let x_1, x_2, \dots, x_N ($N > p$) be random sample of size N and we find the MLE of Σ as $\hat{\Sigma} = S$. Then we have to find out the characteristic roots of S which are given by $|S - \lambda I| = 0$ eqn.

In general, the matrix S is non-singular and let $\lambda_1, \lambda_2, \dots, \lambda_p$ be the ch. roots of S such that $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_p$. The ch. vectors b_i corresponding to λ_i will be obtained by solving $|S - \lambda_i I| = 0$ for b_i clearly b_i are such that $b_i^T b_i = 1$.

Let b_1, b_2, \dots, b_p be the ch. vectors then the estimate of p PCA's are $b_1^T x, b_2^T x, \dots, b_p^T x$ or $\hat{U} = \begin{pmatrix} b_1 \\ b_2 \\ \vdots \\ b_p \end{pmatrix} X$ we may write $B = [b_1, b_2, \dots, b_p]$ then B' will be an orthogonal matrix.

then, $B' S B = \text{diag}(l_1, l_2, \dots, l_p) = D$ (say)

then,

$$S = (B')^{-1} D B^{-1} = B D B^{-1} \text{ (as } B \text{ is orthogonal)}$$

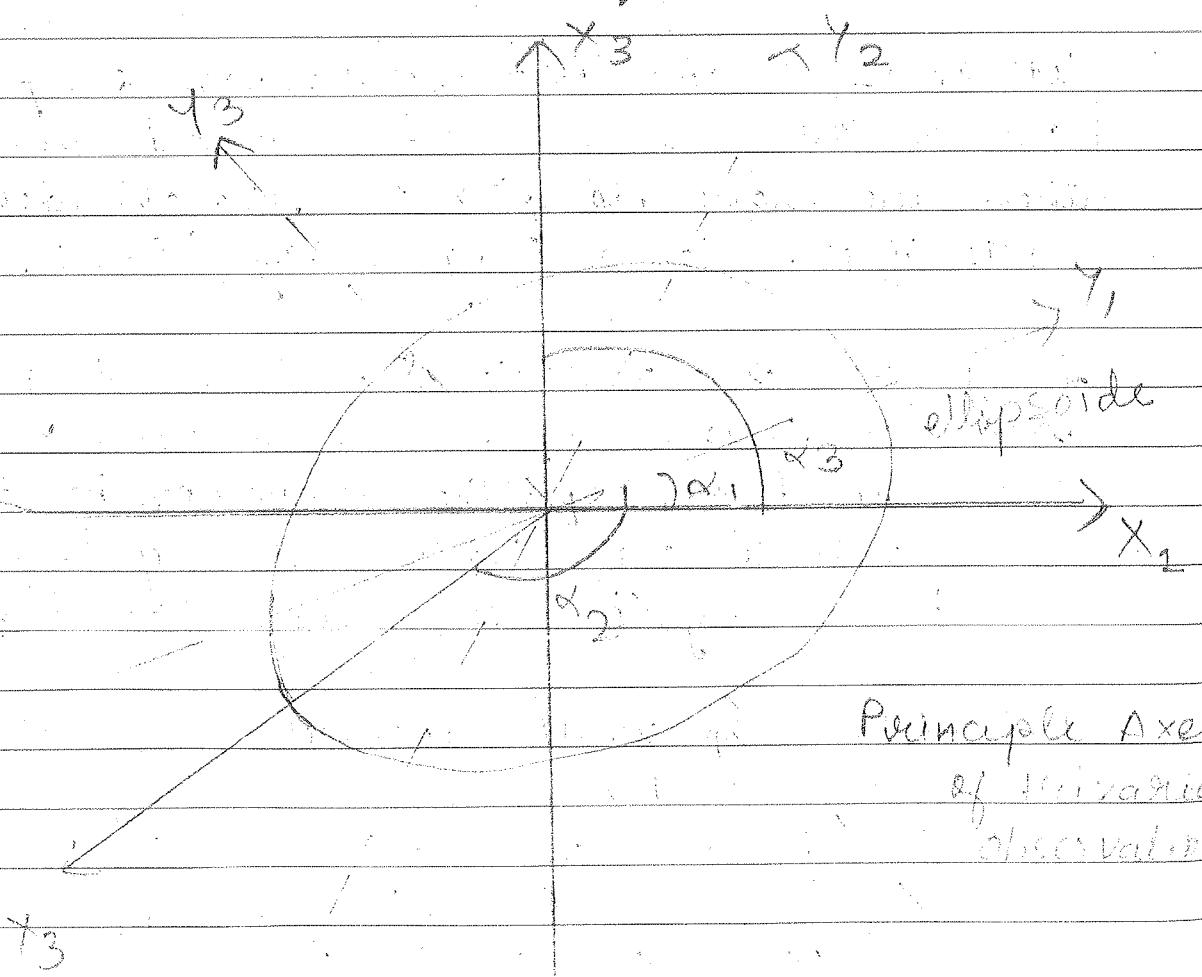
i.e. $S = (b_1, b_2, \dots, b_p) \begin{pmatrix} l_1 & 0 & \dots & 0 \\ 0 & l_2 & \dots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \dots & l_p \end{pmatrix} \begin{pmatrix} b_1 \\ b_2 \\ \vdots \\ b_p \end{pmatrix}$

ie $[S = l_1 b_1 b_1' + l_2 b_2 b_2' + \dots + l_p b_p b_p']$

if the rank of $S = r$! ($r < p$) then

$$[S = b_1 b_1' + \dots + b_r b_r b_r']$$

II Geometrical Meaning Of Principle Component (Trivariate set of observations)



so I imagine that a sample of N bivariate obs
would have the scatter plot shown as follows where
the origin of the response axes has been taken
at the sample mean.

The swarm of points seems to have generally
an ellipsoidal shape with MAJOR axis y_1 and
less well defined minor axes y_2 and y_3 .

Let us denote the angles with the
original response axis as: x_1, x_2, x_3, y_1 passes
through the sample mean point, its obs
orientation is completely determined by the
direction cosines.

$$\begin{aligned} a_{11} &= \cos \alpha_1 \\ a_{21} &= \cos \alpha_2 \\ a_{31} &= \cos \alpha_3 \end{aligned}$$

where $a_{11}^2 + a_{21}^2 + a_{31}^2 = 1$

It is known from the analytic geometry
that, the value of the obs $[x_{11}, x_{12}, x_{13}]$
on the new coordinate axis y_1 will be
given by

$$y_{11} = a_{11}(x_{11} - \bar{x}_1) + a_{21}(x_{12} - \bar{x}_2) + a_{31}(x_{13} - \bar{x}_3) \quad (2)$$

Now the mean of y_1 is

$$\bar{y}_1 = \frac{1}{N} \sum_{l=1}^N y_{1l} = \frac{1}{N} \sum_{l=1}^N \sum_{j=1}^{N-3} (x_{lj} - \bar{x}_j) a_{j1} \quad (3)$$

Let us suppose the notion of major axis as being referring that axis which passes through the direction of maximum variance in the points. In the present case of three responses that variance is :-

$$\frac{1}{(N-1)} \sum_{l=1}^N Y_l^2 = \frac{1}{(N-1)} \sum_{l=1}^N \left[\sum_{j=1}^3 a_{lj} (x_{lj} - \bar{x}_j)^2 \right] \quad (4)$$

the angles of Y_l could be found by differentiating this expression w.r.t a_{lj} and solving for this a_{lj} which makes the derivatives zero.

The solution will be the characteristic vectors of the greatest root of the sample VCM and Y_l would be the centre of the continuum of the 1st P.C. of the system.
(continuum - 3rd image which contains 1st p.c)

Definition:- The P.C. of the sample of N , p-dimensional obs are the new variates specified by the axes of a rigid rotation of the original response coordinate system into an orientation corresponding to the directions of maximum variance in the sample scatter configuration.

The direction cosine of the new axes are normalised characteristic vectors corresponding to the successively a smaller characteristic roots of the sample covariance matrix. If in a given situation two or more roots were equal the directions of the associated axes are not unique and it

may also therefore choose in several positions.

Therefore, it was suggested that if the components are computed using the correlation matrix instead of VCM, the same geometrical interpretation holds and it would be easy to find out the Principle Component without overlapping.

Some Patterned Matrices And their P.C :-

(1) Equi-correlation Matrix :-

In this case the (p x p) covariance matrix has the form

$$\Sigma = \sigma^2 = \begin{bmatrix} 1 & \frac{1}{p} & \dots & \frac{1}{p} \\ \frac{1}{p} & 1 & \dots & \frac{1}{p} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{1}{p} & \frac{1}{p} & \dots & 1 \end{bmatrix}$$

We shall require only first column

The greatest eigen root of the matrix is

$$\lambda_1 = \sigma^2 [1 + (p-1)\frac{1}{p}] \quad \text{--- (2)}$$

and its normalized eigenvector is -

$$\underline{a}_1 = \left[\frac{1}{\sqrt{p}}, \dots, \frac{1}{\sqrt{p}} \right] \text{ as } \underline{a}_1 \underline{a}_1 = \begin{bmatrix} \frac{1}{\sqrt{p}} \\ \vdots \\ \frac{1}{\sqrt{p}} \end{bmatrix} \left[\frac{1}{\sqrt{p}} \dots \frac{1}{\sqrt{p}} \right]$$

$$= \frac{1}{p} + \dots + \frac{1}{p} = 1$$

--- (3)

∴ the 1st P.C. will be

$$Y_1 = \frac{1}{\sqrt{P}} \left(\sum_{j=1}^P X_j \right) \quad (4)$$

i.e. it is merely proportional to the MEAN of the original p responses and it accounts for

$$\frac{100 E + (P-1) S}{P}$$

total variance i.e. the 1st PC

The remaining $(P-1)$ ch. roots are all equal

$$\text{to } \sigma^2 \text{ or } \sigma^2 (1-g) \quad (5)$$

and their corresponding vector are any of the $(P-1)$ linearly independent soln of the eqn $\sigma^2 g (a_{12} + a_{22} + \dots + a_{P2}) = 0$ — (6)

Note that this eqn amounts to the requirement that each of the $(P-1)$ ch. vector be orthogonal to the 1st vector i.e.

$$a_1 a_2 = 0 \quad (7)$$

Note: In many real life situations eqi correlation matrix corresponds to the biological variables such as size of the living things

Eg consider the matrix $P = \begin{bmatrix} 1 & 0.6 & 0.6 & 0.6 \\ 0.6 & 1 & 0.6 & 0.6 \\ 0.6 & 0.6 & 1 & 0.6 \\ 0.6 & 0.6 & 0.6 & 1 \end{bmatrix}$
assume that $\sigma^2 = 1$

Here $\lambda_1 = [1 + (4-1)8] = 1 + (4-1) \times 0.6 = 2.8$

Since $100 \times 2.8 = 70.1\%$ of the variation will be explained by 1st P.C since the 1st eigenvector is $a_1 = \left(\frac{1}{2}, \frac{1}{2}, \frac{1}{2}, \frac{1}{2}\right)$.

Thus $Y_1 = \frac{1}{2}(X_1 + X_2 + X_3 + X_4)$.

The remaining 30% of the variance is **EQUALLY** attributable to the three new variates symmetrically distributed about Y_1 axis.

Principal Component Regression Analysis (PCRA)

Consider the regression model $Y = X\beta + \epsilon$ (1)

where $Y \in \mathbb{R}^n \rightarrow$ dependent variable

$X \in (\mathbb{R}^{n \times p})$ matrix of p -explanatory variables, β

vector parameter of p.m.s and ϵ residual vector of random errors. It is assumed

that $E(\epsilon) = 0$ and $E(\epsilon \epsilon') = \sigma^2 I_n$. Then the OLS

estimate of β is $\hat{\beta} = (X'X)^{-1}X'Y$

This estimator exists subjected to the condition that the rank of X matrix is ' p ' but this is violated if the explanatory variable are correlated. The perfect correlation among the predictors gives the $X'X = 0$ and in this case $(X'X)^{-1}$ does not exist of course g-inverse will exist and in such situation the OLS method of estimation break down.

Situation where the predictors are correlated is known as Multicollinearity problem where

the variables x_1, x_2, \dots, x_p may be not perfectly collinear but there may be some kind of relationship exists among them.

The problem of VIF can be tackled if the columns of X matrix were 'orthogonalized' or the variables which cause the VIF problem are dropped from the analysis but dropping of a variable may not be a wise decision because the variable being dropped may be important to study the variation in the dependent variable.

The principal component technique is used to orthogonalized the X matrix. In PCA the component are so selected that the selected component can explain maximum variation in data. But in case of regression analysis the variables are selected in such a way that the correlation between y and selected variables (or) component is maximum. Little, Spamer and Bartlett (1977) have discussed the method of selection of principal component in regression analysis.

* METHOD:

Without loss of generality let us consider that $\bar{y} = \bar{x}_{ij} = 0$ and $\text{var}(y) = \text{var}(\bar{x}_j) = 1$

and let $U \sim N_n(0, \sigma^2 H)$ where $H = I - \frac{1}{n} 1 1'$ where $1 = (1, 1, \dots, 1)$ and

let $\Gamma = (\gamma_{1j}, \gamma_{2j}, \dots, \gamma_{pj})$ be an orthogonal matrix where γ_j is the j th eigen vector

Corresponding to the eigen value λ_j , let y_{ij} be the i th element of y_j vector. If $\lambda_1 > \lambda_2 > \dots > \lambda_p$ consider the transformation of X matrix to Z matrix s.t.

$Z = \Gamma' X$ where the elements of Z are z_1, z_2, \dots, z_p . Now if we replace X by Z in (2) we have

$$Y = Z \Gamma' B + U = Z B + U \quad \text{--- (2)}$$

where $B = \Gamma' \beta$

∴ the LSE of B is $\hat{B} = (Z'Z)^{-1}Z'y$

This is \hat{B} is unbiased as $E(\hat{B}) = E(Z'Z)^{-1}Z'(ZB+U)$

$$= E(Z'Z)^{-1}Z'ZB + (Z'Z)^{-1}Z'E(U) \\ \text{as } E(U) = 0 \\ = B$$

Also,

$$\text{Var}(\hat{B}) = E(\hat{B} - B)(\hat{B} - B)' \\ = E(Z'Z)^{-1}Z'UU'Z(Z'Z)^{-1} \\ (\text{as } U = \Delta^{-1}\varepsilon^2 \text{ (say)})$$

∴ $\hat{B} \sim N_p(B, \Delta^{-1}\varepsilon^2)$ FIM we can write

$$\hat{B}_i = \frac{1}{n} \sum_{j=1}^n z_{ij} y_j \text{ and}$$

$$\text{Var}(B_i) = \frac{\sigma^2}{n} \quad i=1, 2, \dots, p$$

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From this PCRA we obtain the sum of squares due to regression

$$SS_{\text{(Reg)}} = \hat{B}'Z'y$$

$$\text{and } SS_{\text{(error)}} = y'y - \hat{B}'Z'y$$

$$\text{an estimator of } \sigma^2 \text{ is: } \frac{[y'y - \hat{B}'Z'y]}{(n-p-1)}$$

Now, the test statistic for testing the significance of $H_0: \beta_i = 0$, i.e. if β_i is zero we have

$$S.E(\hat{\beta}_i)$$

estimated $\beta \rightarrow$ the vector of pm's from this the original parameter vector β is estimated $\hat{\beta}$, where

$$\hat{\beta} = \Gamma \hat{B}$$

(3) thus, $E(\hat{\beta}_i) = \sum_{j=1}^p x_{ij} \beta_{ji}$ with

$$\text{var}(\hat{\beta}_i) = (\sigma^2)^2 \sum_{j=1}^p \frac{x_{ij}^2}{n x_{ii}}$$

from sample the estimator $\hat{\beta}_i$ is obtained by

$$\hat{\beta}_i = \sum_{j=1}^p g_{ij} \beta_j \quad \text{and} \quad \text{var}(\hat{\beta}_i) = \sum_{j=1}^p g_{ij}^2 / n l_j$$

where l_j is the j th eigen value of Σ ($n \times n$ matrix)

g_{ij} : i th element of eigen values.

corresponding to l_j ($j = 1, 2, \dots, p$)

UNIT 3

(One-way) MULTIVARIATE ANALYSIS OF VARIANCE: (MANOVA)

In many real life situations often more than two populations need to be compared, for this purpose random samples are calculated from each of the 'g' populations (say) and these were arranged as:-

$$\text{pop } 1 : x_{11}, x_{12}, \dots, x_{1n},$$

$$\text{pop } 2 : x_{21}, x_{22}, \dots, x_{2n},$$

!

$$\text{pop } g : x_{g1}, x_{g2}, \dots, x_{gn}$$

(1)

Assumptions:-

- (1) $X_{11}, X_{12}, \dots, X_{pq}$ is a random sample of size n_1 from a popⁿ having mean μ_l , $l=1, 2, \dots, g$ the obs from different popⁿ are independent.
- (2) All the popⁿ share common covariance matrix.
- (3) Each population is multivariate normal then the hypothesis of interest is :-

$$H_0 : \begin{bmatrix} \mu_{11} \\ \mu_{12} \\ \vdots \\ \mu_{1p_1} \end{bmatrix} = \begin{bmatrix} \mu_{12} \\ \mu_{22} \\ \vdots \\ \mu_{2p_2} \end{bmatrix} = \dots = \begin{bmatrix} \mu_{1g} \\ \mu_{2g} \\ \vdots \\ \mu_{pg} \end{bmatrix}$$

OR

$$H_0: \mu_1 = \mu_2 = \dots = \mu_g$$

THE MODEL:

Analogy to univariate case, the model for comparing g -population mean vector is :

$$X_{ij} = \bar{\mu} + T_l + \epsilon_{ij}; j=1, 2, \dots, n_l, l=1, \dots, g$$

where $\epsilon_{ij} \sim N\ N_p(0, \Sigma)$

$\bar{\mu}$ = Overall mean vector

T_l = l th treatment effect with $\sum_{l=1}^g T_l = 0$

According to model (?), each component of the obs vector X_{ij} satisfies the assumption of univariate model. Now, a vector of observations may be decomposed (as suggested by the model)

$$\text{Thus } \underline{x}_{ej} = \bar{x} + (\bar{x}_e - \bar{x}) + (\underline{x}_{ej} - \bar{x}_e) \quad (3)$$

↑ ↓ ↓ ↗
 (Obs & env) overall sample mean estimated treatment Residual
 ↓ ↓ ↓ ↗
 effect $\rightarrow t_{ej}$ \bar{x}_e

the decomposition in (3) leads to the Multivariate analogue of the univariate SS. break-up.

Note that the cross-product (C.P)

$(\underline{x}_{ej} - \bar{x})(\underline{x}_{ej} - \bar{x})'$ can be written as:-

$$\begin{aligned}
 (\underline{x}_{ej} - \bar{x})(\underline{x}_{ej} - \bar{x})' &= [(\underline{x}_{ej} - \bar{x}_e) + (\bar{x}_e - \bar{x})] \\
 &\quad [(\underline{x}_{ej} - \bar{x}_e) + (\bar{x}_e - \bar{x})]' \\
 &= (\underline{x}_{ej} - \bar{x}_e)(\underline{x}_{ej} - \bar{x}_e)' + (\underline{x}_{ej} - \bar{x}_e)(\bar{x}_e - \bar{x})' \\
 &\quad + (\bar{x}_e - \bar{x})(\underline{x}_{ej} - \bar{x}_e)' + (\bar{x}_e - \bar{x})(\bar{x}_e - \bar{x})'
 \end{aligned}$$

$$\sum_{j=1}^{n_e} (\underline{x}_{ej} - \bar{x}_e) = 0 \rightarrow \text{matrix}$$

$$\text{Next, } \sum_{l=1}^q n_e (\bar{x}_e - \bar{x}_l) = 0 \quad \text{as } \sum_l n_e t_{el} = 0$$

\therefore Summing the C.P over j and l gives:-

$$\sum_{l=1}^q \sum_{j=1}^{n_e} (\underline{x}_{ej} - \bar{x})(\underline{x}_{ej} - \bar{x})'$$

$$\begin{aligned}
 \text{Total SS} &= \sum_{l=1}^q n_e (\bar{x}_e - \bar{x})(\bar{x}_e - \bar{x})' + \sum_{l=1}^q \sum_{j=1}^{n_e} (\underline{x}_{ej} - \bar{x}_e)(\underline{x}_{ej} - \bar{x}_e)' \\
 &\quad \downarrow \qquad \qquad \qquad \downarrow \\
 &\quad \text{treat SS (Between)} \qquad \text{Residual S.S (within)}
 \end{aligned}$$

now, the within SS and C.P matrix can be expressed as:-

$$W = \sum_{l=1}^q \sum_{j=1}^{n_e} (\underline{x}_{ej} - \bar{x}_e)(\underline{x}_{ej} - \bar{x}_e)'$$

$$\Delta (n_1-1)S_1 + (n_2-1)S_2 + \dots + (n_g-1)S_g$$

where S_e is the sample cov. matrix for i^{th} sample W matrix is the generalization of $(n_1+n_2-2)S_p$ matrix. Analogous to univariate, test of the hypothesis of No. treatment effect is $H_0: \pi_1 = \pi_2 = \dots = \pi_g = 0$

MANOVA Table

Source of Variation	Degrees of freedom.	Matrix of SS & CP
Treatments	$(g-1)$	$B = \sum_{l=1}^g n_l (\bar{x}_l - \bar{\bar{x}})(\bar{x}_l - \bar{\bar{x}})'$
Residual (Error)	$\sum_{l=1}^g (n_l - g)$	$W = \sum_{l=1}^g \sum_{j=1}^{n_l} (x_{lj} - \bar{x}_{lj})(x_{lj} - \bar{x}_{lj})'$
Total	$(\sum_{l=1}^g n_l - 1)$	$B + W = \sum_{l=1}^g \sum_{j=1}^{n_l} (x_{lj} - \bar{\bar{x}})(x_{lj} - \bar{\bar{x}})'$

One of the tests to the test H_0 involves GENERALISED VARIANCE.

We reject H_0 if the ratio of the generalized variances $\Lambda^* = |W| / |B|$ is too small. Proposed by WILKS

Wilks Λ^* has relation with likelihood ratio criterion and the exact distribution of Λ^* can be derived for some special cases. However for large sample a modification of Λ^* is due to Bartlett. He has shown that BARTLETT and can be used to test H_0 .

He has shown that $\left[(m-1) - (p+g) \right] \frac{\log \Delta^2}{2} \sim \chi^2_{(g-1)d}$

* MANOVA Two-way fixed effect Model (with Interaction)

→ (Univariate) → Denoting the i th observation at the level ' l ' of factor 1 and level ' k ' of factor 2 by x_{lkr} , we may write a two-way

univariate model as: $x_{lkr} = \mu + \tau_l + \beta_k + \gamma_{lk} + \epsilon_{lkr}$
 $l = 1, \dots, g$ and $k = 1, \dots, b$, $r = 1, 2, \dots, n$

proceeding by analogy, we specify two-way fixed effect for a vector response consisting of p -component as:

$$x_{lkr} = \mu + \tau_l + \beta_k + \gamma_{lk} + \epsilon_{lkr} \quad (1)$$

$$\text{where } \sum_{l=1}^g \tau_l = \sum_{k=1}^b \beta_k = \sum_{l=1}^g \gamma_{lk} = \sum_{k=1}^b \gamma_{lk} = 0$$

the vectors are all of order (px_1) and $\text{Sol}_K \sim N_p(0, \Sigma)$. Thus, the responses consist of p -measurements repeated ' n ' times at each possible combinations of levels of factor 1 and factor 2. We can decompose the observation vector x_{lkr} as:

$$x_{lkr} = \bar{x} + (\bar{x}_{l\cdot} - \bar{x}) + (\bar{x}_{\cdot k} - \bar{x}) + (\bar{x}_{lk} - \bar{x}_{l\cdot} - \bar{x}_{\cdot k} + \bar{x}) + (x_{lkr} - \bar{x}_{lk}) \quad (2)$$

where \bar{X} : overall average of the obser. vectors

\bar{X}_l : overall average of the obser vectors at the l^{th} level of factor 1

\bar{X}_k : Overall average of the observ vectors at the k^{th} level of factor 2

\bar{X}_{lk} : Overall average of the obser vector at l^{th} level of factor 1 and k^{th} level of factor 2

We have :

$$\sum_{l=1}^g \sum_{k=1}^b \sum_{r=1}^n (\bar{X}_{lkr} - \bar{X}) (\bar{X}_{lkr} - \bar{X})'$$

$$= \sum_{l=1}^g b n (\bar{X}_{l\cdot} - \bar{X}) (\bar{X}_{l\cdot} - \bar{X})' + \sum_{k=1}^b g n (\bar{X}_{\cdot k} - \bar{X}) (\bar{X}_{\cdot k} - \bar{X})'$$

$$+ \sum_{l=1}^g \sum_{k=1}^b n (\bar{X}_{l k} - \bar{X}_{l\cdot} - \bar{X}_{\cdot k} + \bar{X}) (\bar{X}_{l k} - \bar{X}_{l\cdot} - \bar{X}_{\cdot k} + \bar{X})'$$

$$+ \sum_{l=1}^g \sum_{k=1}^b \sum_{r=1}^n (\bar{X}_{lkr} - \bar{X}_{lk}) (\bar{X}_{lkr} - \bar{X}_{lk})'$$

— (3)

degree of freedom

$$(gbn-1) = (g-1)f(b-1) + (g-1)(b-1) + gb(n-1)$$

— (4)

MANOVA Table

Source of Variation	Degrees of freedom	Matrix of SS
Factor - 1	(g-1)	S.S. Factor 1 = $\sum_{l=1}^g b_n (\bar{X}_{l.} - \bar{X})(\bar{X}_{l.} - \bar{X})'$
Factor - 2	(b-1)	S.S. Factor 2 = $\sum_{k=1}^b g_n (\bar{X}_{.k} - \bar{X})(\bar{X}_{.k} - \bar{X})'$
Interaction	((g-1)(b-1))	S.S. Interaction = $\sum_{l=1}^g \sum_{k=1}^b n (\bar{X}_{lk} - \bar{X}_{l.} - \bar{X}_{.k} + \bar{X})(\bar{X}_{lk} - \bar{X}_{l.} - \bar{X}_{.k} + \bar{X})'$
Residual	: g b (n-1)	S.S. Residual = $\sum_{l=1}^g \sum_{k=1}^b \sum_{t=1}^n (\bar{X}_{lkt} - \bar{X}_{lk})(\bar{X}_{lkt} - \bar{X}_{lk})'$
Total.	(gbn-1)	S.S. Total = $\sum_{l=1}^g \sum_{k=1}^b \sum_{t=1}^n (\bar{X}_{lkt} - \bar{X})(\bar{X}_{lkt} - \bar{X})'$

Hypothesis:-

A test of $H_0: \gamma_{11} = \gamma_{12} = \dots = \gamma_{gb} = 0$
 (No interaction effect)

or

H_1 : At least one $\gamma_{lk} \neq 0$
 is conducted we reject H_0 for smaller values of

$$\Delta^* = \frac{(\text{S.S. Pres})}{(\text{S.S. Int} + \text{S.S. Pres})} - (5)$$

Note: The LRT procedure requires that $p \leq g_b(n-1)$ so that S.S. residual will be a positive definite matrix.

For large sample Δ^* can be compared with χ^2 -distribution using Bartlett's connection to improper χ^2 -approximation. Therefore we rej H_0 if

$$- [g_b(n-1) + (p+1) - (g-1)(b-1)] \log \Delta^* \sim \chi^2_{(g-1)(b-1), df} \quad (6)$$

Remark

- In two way manova generally the test for interested interaction is carried out before testing the main factor effects.
- As if the interaction effect exists then the factor effects do not have clear interpretation. Therefore from a practical interpretation, therefore, from a practical standpoint it is not advisable to proceed with additional multivariate test.
- In multivariate MANOVA model we proceed for testing of factor-1 and factor-2 as the main effects if no interaction is present.

First we consider the testing of

$$H_0: \tau_1 = \tau_2 = \dots = \tau_g = 0$$

vs H_1 : atleast one $\tau_i \neq 0$

This hypo specifies No factor - 1 effects and Some factors - 1 effects respectively.

The prop test stat is

$$\Delta^* = \frac{|S.S.P \text{ res}|}{|S.S.P \text{ fac 1} + S.S.P \text{ res}|} \quad - (7)$$

→ Again using Bartlett's correction, the L.R for testing of No factor - 1 effects is conducted by rej H_0 at level α if

$$-\left[gb(n-1) - (p+1)(g-1)\right] \log \Delta^* > \chi^2_{(g-1)} \text{ i.d.r.} \quad - (8)$$

In similar manner if factor - 2 effects are tested by considering

$$H_0: \beta_1 = \beta_2 = \dots = \beta_k = 0$$

vs H_1 : atleast one $\beta_j \neq 0$,

then the test stat is :

$$\Delta = \frac{|S.S.P \text{ Res}|}{|S.S.P \text{ fac 2} + S.P \text{ res}|} \quad - (9)$$

Again using Bartlett's correction the L.R. for testing of No factor - 2 effects is conducted by rej. H_0 at level α if

$$-\frac{[gb(n-1) - (p+1)(b-1)]}{2} \log \Delta^* > \chi^2_{(b-1)}; \alpha.$$

— (10)

—x—

