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**Kunthavai Naachiyaar Govt. Arts College (W) Autonomous,
Thanjavur**

M.Stat

18KPI502

CORE COURSE - II

REAL ANALYSIS AND MATRIX THEORY

Hours:6

Credit:5

Unit – I

Functions – Real valued function, Equivalence, countability, least upper bounds
Sequence of real numbers – Definition, limit of a sequence, convergent sequence, divergent
sequence, bounded sequence, monotone sequence, limit superior and limit inferior, Series of real
numbers – convergence and divergence, series with non-negative terms, alternating series,
conditional convergence and absolute convergence.

Unit – II

Calculus – Sets of Measure zero, Definition of the Riemann integral, existence of the
Riemann integral, Derivatives, Rolle's Theorem, the law of the mean, Fundamental theorems of
calculus, improper integrals.

Unit – III

Rank of a matrix – elementary transformation of a matrix, Equivalent Matrices,
Elementary matrices, Echelon matrix – Hermite Canonical form, Sylvester's law, Frobenius
inequality, certain results on the rank of an Idempotent matrix.

Unit – IV

Eigen values and Eigen Vectors – properties, Cayley – Hamilton theorem, application of
Cayley-Hamilton theorem – simple problems.

Unit – V

Generalized inverse of a matrix: Definition, different classes of generalized inverse,
Properties of G- inverse – classes properties – properties of Moore and Penrose – application
of generalized inverse in the solution of system of linear equations. Quadratic forms – Definition,
classification of the quadratic form, positive semi-definite quadratic form and Canonical
reduction.

Text Books:

1. Goldberg, R. (1963). Method of Real Analysis, Oxford & IBH publishers. New Delhi
(Unit – I: Chapter – 1: Page no.3 – 17, 21-75, Unit – II: Chapter-7 Page no.(156 –
194).Biswas.S(1996) , A Text book of Matrix Algebra, New Age International
Publishers, New Delhi. (Unit – III: Chapter – 5, Unit – IV: Chapter – 7: Page no.185 –
198, 208 – 209, 213 – 227. Unit – V: Chapter – 8: page no.228 – 245 and Chapter – 9:
page no. 267 – 268, 317 – 323).

REAL ANALYSIS AND MATRIX THEORY

UNIT-1

Functions :

" If to each x (in a set S) there corresponds one and only one value of y , then we say that y is a function of x ".

Definition :

If A, B are sets, then the Cartesian Product of A and B (denoted $A \times B$) is the set of all ordered pairs $\langle a, b \rangle$ where $a \in A$ and $b \in B$.

Definition : [THE COMPOSITION OF FUNCTIONS]

If $f: A \rightarrow B$ and $g: B \rightarrow C$, then we define the function $g \circ f$ by,

$$g \circ f(x) = g[f(x)] \quad (x \in A)$$

That is the image of x under $g \circ f$ is defined to be the image of $f(x)$ under g . Then function $g \circ f$ is called the composition of function of f with g .

Real valued Functions :

The range of a given function f is contained in the set of all real numbers. (We henceforth denote the set of all real numbers by \mathbb{R}). If $f: A \rightarrow \mathbb{R}$, we call f a real valued function. If $x \in A$, then $f(x)$ (heretofore called the image of x under f) is also called the value of f at x .

Definition :

If $f: A \rightarrow \mathbb{R}$ and $g: A \rightarrow \mathbb{R}$, We define $f+g$ as the function whose value at $x \in A$ is equal to $f(x) + g(x)$. That is

$$(f+g)(x) = f(x) + g(x) \quad (x \in A)$$

EQUIVALENCE AND COUNTABILITY :

Definition :

If $f: A \Rightarrow B$ and f is 1-1 then f is called a 1-1 correspondence (between A and B). If there exists a 1-1 correspondence between the sets A and B , then A and B are called equivalent.

1. Every Set A is equivalent to itself.
2. If A and B are equivalent, then B and A are equivalent.
3. If A and B are equivalent and B and C are equivalent then A and C are equivalent.

Definition :

The Set A is said to be countable (or denumerable) if A is equivalent to the Set I of Positive Integers. An uncountable Set is an infinite Set which is not countable.

The element of A are then the integers of images $f(1), f(2), \dots$, of the positive integers,

$$A = \{ f(1), f(2), \dots \}$$

Least Upper Bounds:

Definition:

If $A \subset \mathbb{R}$ is bounded above, then N is called an upper bound for A if $x \leq N$ for all $x \in A$. If $A \subset \mathbb{R}$ is bounded below, then M is called a lower bound for A if $M \leq x$ for every $x \in A$.

Definition:

Let the subset A of \mathbb{R} be bounded above. The number L is called the upper bound for A if (1) L is an upper bound for A , and (2) no number smaller than L is an upper bound for A .



Sequence of Real Numbers:

A sequence of real numbers is a function from the set \mathbb{N} of natural numbers to the set \mathbb{R} of the real numbers. If $f: \mathbb{N} \rightarrow \mathbb{R}$ is a sequence, and if $a_n = f(n)$ for $n \in \mathbb{N}$, then we write the sequence f as (a_n) (or) (a_1, a_2, \dots) . A sequence of real numbers is also called a real sequences.

Convergence Sequence and Divergent :

A Sequence (a_n) in \mathbb{R} is said to converge to a real Number a if for every $\epsilon > 0$, there exists Positive integer N (in general depending on ϵ) Such that

$$|a_n - a| < \epsilon \quad \forall n \geq N,$$

and in that case, the number a is called a limit of a Sequence (a_n) and (a_n) is called a Convergent sequences.

A Sequence which does not converge is called a divergent Sequence.

Bounded Sequence :

A Sequence (a_n) is said to be bounded above if there exists a real number M such that $a_n \leq M$ for all $n \in \mathbb{N}$; and the Sequence (a_n) is said to be bounded below if there exists a real number M' such that $a_n \geq M'$ for all the $n \in \mathbb{N}$.

A Sequence which is bound above and bounded below is said to be a bounded Sequence.

Monotonic Sequences :

We can infer the convergence or divergence of a Sequence in certain cases by observing the way the terms of the Sequence varies.

Definition :

Consider a Sequence (a_n)

(i) (a_n) is said to be monotonically increasing, if $a_n \leq a_{n+1}$ for all $n \in \mathbb{N}$.

(ii) (a_n) is said to be monotonically decreasing, if $a_n \geq a_{n+1}$ for all $n \in \mathbb{N}$.

If strict inequality occur in (i) and (ii) then we say that the sequences is strictly increasing and decreasing.

Definition :

Let S be a subset of \mathbb{R} , Then

(i) S is said to be bounded above if there exists $b \in \mathbb{R}$ such that $x \leq b$ for all $x \in S$, and in that case b is called a upper bound of S .

(ii) S is said to be bounded below if there exists $a \in \mathbb{R}$ such that $a \leq x$ for all $x \in S$, and in that case of a is called a lower bound of S .

Limit Superior :

If $\{S_n\}_{n=1}^{\infty}$ is a sequence of real numbers that is not bounded above, we write the,

$$\limsup_{n \rightarrow \infty} S_n = \infty$$

$$\limsup_{n \rightarrow \infty} n = \infty.$$

Limit inferior :

If $\{S_n\}_{n=1}^{\infty}$ is a sequence of real numbers that is not bounded below, we write $\liminf_{n \rightarrow \infty} S_n = -\infty$

Thus $\liminf_{n \rightarrow \infty} (-1)^n = -1$, $\liminf_{n \rightarrow \infty} n = \infty$,

$\liminf_{n \rightarrow \infty} (-n) = -\infty$. The sequence $1, -1, 1, -2, 1, -3, 1, -4, \dots$ has $\liminf = -\infty$.

Real Number:

A Series of Number is an expression of the form

$a_1 + a_2 + a_3 + \dots$, or more compactly as $\sum_{n=1}^{\infty} a_n$; where (a_n) is a sequence of real numbers. The number is called the n -th term of the series and the sequence $S_n = \sum_{i=1}^n a_i$ is called the n th partial sum of the series $\sum_{n=1}^{\infty} a_n$.

Convergence and divergence of Series :

A series $\sum_{n=1}^{\infty} a_n$ is said to converge (to $S \in \mathbb{R}$) if the sequence $\{S_n\}$ of partial sum of the series converge (to $S \in \mathbb{R}$).

If $\sum_{n=1}^{\infty} a_n$ converges to S , then we write $\sum_{n=1}^{\infty} a_n = S$.

A series which does not converge is called a divergent series

Alternating Series:

A Series of the form $\sum_{n=1}^{\infty} (-1)^{n+1} u_n$ where (u_n) is a sequence of positive terms is called an alternating Series.

Absolute convergence and conditional convergence:

A Series $\sum_{n=1}^{\infty} a_n$ is said to be a converge absolutely, if $\sum_{n=1}^{\infty} |a_n|$ converges.

Every absolutely convergent Series converges, let S_n and σ_n be the n -th partial sums of the Series $\sum_{n=1}^{\infty} |a_n|$ respectively, then for $n > m$, we have

$$|S_n - S_m| = \left| \sum_{j=m+1}^n a_j \right| \leq \sum_{j=m+1}^n |a_j| = |\sigma_n - \sigma_m|.$$

Since $\{\sigma_n\}$ converges, it is a Cauchy sequence. Hence from that $\{S_n\}$ is also a Cauchy sequence.

Definition:

A Series $\sum_{n=1}^{\infty} a_n$ is said to be converge conditionally, if $\sum_{n=1}^{\infty} a_n$ converges but not absolutely to the sequence.

Example :

(i) The Series $\sum_{n=1}^{\infty} \frac{(-1)^{n+1}}{n}$ is conditionally convergent
sequences.

(ii) The Series $\sum_{n=1}^{\infty} \frac{(-1)^{n+1}}{n^2}$ is absolutely convergent

sequences.

(iii) The Series $\sum_{n=1}^{\infty} \frac{(-1)^{n+1}}{n!}$ is absolutely convergent.

For any $\alpha \in \mathbb{R}$ the Series $\sum_{n=1}^{\infty} \frac{\sin(n\alpha)}{n^2}$ is absolutely
convergent :

$$\left| \frac{\sin(n\alpha)}{n^2} \right| \leq \frac{1}{n^2} \quad \forall n \in \mathbb{N}.$$

Since $\sum_{n=1}^{\infty} \frac{1}{n^2}$ converges,

By comparison test, Hence the value is

$\sum_{n=1}^{\infty} \left| \frac{\sin(n\alpha)}{n^2} \right|$ also convergent to the

sequences.

Series with non-negative terms:

The easiest series to deal with are those with non-negative terms. For these series all the theory on the convergence and divergence is embodied in the following theorem:

If $\sum_{n=1}^{\infty} a_n$ is a series of non-negative numbers with $S_n = a_1 + \dots + a_n$ ($n \in \mathbb{I}$), then,

(a) $\sum_{n=1}^{\infty} a_n$ converges if the sequence $\{S_n\}_{n=1}^{\infty}$ is bounded;

(b) $\sum_{n=1}^{\infty} a_n$ diverges if $\{S_n\}_{n=1}^{\infty}$ is not a bounded.

Definition:

Let $\sum_{n=1}^{\infty} a_n$ be a series of real numbers,

(a) If $\sum_{n=1}^{\infty} |a_n|$ converges we say that $\sum_{n=1}^{\infty} a_n$ converges absolutely.

(b) If $\sum_{n=1}^{\infty} a_n$ converges but $\sum_{n=1}^{\infty} |a_n|$ diverges, we say that $\sum_{n=1}^{\infty} a_n$ converges conditionally.

Unit-2

Sets of Measure Zero :

Definition :

The Subset E of \mathbb{R} is said to be of measure zero if for each $\epsilon > 0$, there exists a finite or a countable number of open intervals I_1, I_2 such that

$$E \subset \cup I_n \quad \text{and} \quad \sum |I_n| < \epsilon.$$

Definition :

A statement is said to hold at almost every point of $[a, b]$ (or almost everywhere in $[a, b]$) if the set of points of which the statement does not hold is of measure zero.

Definition of the Riemann Integral :

Definition :

Let I be a bounded interval of real numbers and let f be a bounded function defined on I . We define

$M[f; I]$ and $m[f; I]$ by

$$M[f; I] = \sup_{x \in I} f(x)$$

$$m[f; I] = \inf_{x \in I} f(x)$$

Definition : 2

A Subdivision of the closed bounded interval $[a, b]$ we mean a finite subset $\{x_0, x_1, \dots, x_n\}$ of $[a, b]$ such that,

$$a = x_0 < x_1 < x_2 < \dots < x_n = b$$

If σ and τ are two subdivisions of $[a, b]$ we say that τ is a refinement of σ if $\sigma \subset \tau$.

If $\sigma = \{x_0, x_1, \dots, x_n\}$ is a subdivision of $[a, b]$ then the closed intervals.

$$I_1 = [x_0, x_1], I_2 = [x_1, x_2], \dots, I_n = [x_{n-1}, x_n]$$

are called the components of σ .

Definition : 3

Let f be a bounded function on $[a, b]$ and let σ be a subdivision of $[a, b]$ with components I_1, I_2, \dots, I_n . We define $U[f; \sigma]$, called the upper sum for f corresponding to σ , by

$$U[f; \sigma] = \sum_{k=1}^n M[f; I_k] |I_k|.$$

Here $|I_k|$ is the length of I_k which is given by $|I_k| = x_k - x_{k-1}$.

$$L[f; \sigma] = \sum_{k=1}^n m[f; I_k] |I_k|.$$

It is denoted to the lower sum of the corresponding to σ .

Definition : 4

Let f be a bounded function on $[a, b]$ we define the upper integral of f over $[a, b]$ by,

$$\int_a^b f(x) dx = g. l. b U [f; T]$$

We define the lower integral of f over $[a, b]$ by

$$\int_a^b f(x) dx = l. u. b L [f; T]$$

It is common to denote upper integrals and lower integrals

of f ; $\overline{\int_a^b f}$ and $\underline{\int_a^b f}$.

Definition : 5

If f is bounded on $[a, b]$ we say that f is Riemann integrable on $[a, b]$ if

$$\underline{\int_a^b f} \leq \overline{\int_a^b f}.$$

In this case, we define the Riemann integrals,

$$\int_a^b f(x) dx = \int_a^b f = \underline{\int_a^b f} \leq \overline{\int_a^b f}$$

The class of the Riemann integral is denoted by $R [a, b]$.

Existence of the Riemann Integral:

Let f be a bounded function on $[a, b]$ then $f \in \mathcal{R}$ if and only if f is continuous at almost every point of $[a, b]$.

Properties of Riemann Integral:

(i) If $f \in \mathcal{R} [a, b]$ and $a < c < b$, then $f \in \mathcal{R} [a, c]$, $f \in \mathcal{R} [c, b]$ and

$$\int_a^b f = \int_a^c f + \int_c^b f.$$

(ii) If $f \in \mathcal{R} [a, b]$ and λ is any real number then $\lambda f \in \mathcal{R} [a, b]$ and

$$\int_a^b \lambda f = \lambda \int_a^b f.$$

(iii) If $f, g \in \mathcal{R} [a, b]$ then $f+g \in \mathcal{R} [a, b]$ and

$$\int_a^b (f+g) = \int_a^b f + \int_a^b g$$

If $f \in \mathcal{R} [a, b]$ and if $f(x) \geq 0$, almost where in $[a, b]$ then

$$\int_a^b f \geq 0.$$

Derivative:

Definition:

Let f be defined on an interval J . If $c \in J$, we say that f has a derivative at c ,

$$\lim_{x \rightarrow c} \frac{f(x) - f(c)}{x - c}$$

(i) If f has a derivative at c , then f is continuous at c .

(ii) If f and g both have derivatives at c , then do $f+g$, $f-g$, fg , and

$$(f+g)'(c) = f'(c) + g'(c)$$

$$(f-g)'(c) = f'(c) - g'(c)$$

$$(fg)'(c) = f'(c)g(c) + f(c)g'(c)$$

If $g(c) \neq 0$, then f/g has derivative at c ,

$$\left(\frac{f}{g}\right)'(c) = \frac{f'(c)g(c) - f(c)g'(c)}{g(c)^2}$$

(iii) If g has derivative at c and f has derivative at $g(c)$ then $f \circ g$ has derivative at c and,

$$(f \circ g)'(c) = f'(g(c))g'(c)$$

(iv) Let f be one-to-one function and let ϕ be the inverse of f . If f is continuous at c and if ϕ has derivative at $d = f(c)$ with $\phi'(d) \neq 0$, then $f'(c)$ exists and,

$$f'(c) = \frac{1}{\phi'(d)}$$

Rolle's Theorem:

(i) Let f be a continuous function on $[a, b]$. If the maximum value of f is attained at $c \in (a, b)$ and if $f'(c)$ exists then $f'(c) = 0$.

(ii) Let f be a continuous function on $[a, b]$. If the minimum value of f is attained at $c \in (a, b)$ and if $f'(c)$ exists, then $f'(c) = 0$.

(iii) Let f be continuous on $[a, b]$ with $f(a) = f(b) = 0$.

If $f'(x)$ exists for every $x \in (a, b)$ then there exists some point $c \in (a, b)$ such that $f'(c) = 0$.

(iv) If f has derivatives at every point of $[a, b]$ then f' takes value between $f'(a)$ and $f'(b)$.

Mean Value Theorem :

Theorem 1: If f is continuous on $[a, b]$ and differentiable on (a, b) then there exists a point $c \in (a, b)$ such that

$$f'(c) = \frac{f(b) - f(a)}{b - a}$$

Theorem 2: If f is continuous on $[a, b]$ and differentiable on (a, b) and $f'(x) > 0$ for all $x \in (a, b)$ then f is strictly increasing on $[a, b]$.

Theorem 3: If f, g are continuous on $[a, b]$, $f(a) \neq f(b)$, if f, g are both differentiable on (a, b) and if $f'(x)$ and $g'(x)$ are both zero for any $x \in (a, b)$ then there exists a

Points $c \in (a, b)$ such that,

$$\frac{f'(c)}{g'(c)} = \frac{f(b) - f(a)}{g(b) - g(a)}.$$

The Mean Value Theorem:

Theorem 1: If f is continuous on $[a, b]$ and differentiable on (a, b) then there exists a point $c \in (a, b)$ such that,

$$f'(c) = \frac{f(b) - f(a)}{b - a}$$

Theorem 2: If f is continuous on $[a, b]$ and differentiable on (a, b) and if $f'(x) > 0$ for all $x \in (a, b)$ then f is strictly increasing on $[a, b]$.

Theorem 3: If f, g are continuous on $[a, b]$, $g(a) \neq g(b)$ if f, g are both differentiable on (a, b) and if $f'(x)$ and $g'(x)$ are not both zero for any $x \in (a, b)$ then there exists a point $c \in (a, b)$ such that,

$$\frac{f'(c)}{g'(c)} = \frac{f(b) - f(a)}{g(b) - g(a)}.$$

The Fundamental Theorem of Calculus:

First fundamental Theorem:

Theorem:

If f is continuous on $[a, b]$ and if,

$$F(x) = \int_a^x f(t) dx.$$

for each $x \in (a, b)$ then $F'(x) = f(x)$ for all $x \in [a, b]$.

Theorem: 2

If $f'(x) = 0$ for every $x \in (a, b)$ then $f(x) = f(a) \forall x \in (a, b)$

Theorem: 3

If $f'(x) = g'(x)$ for all $x \in (a, b)$ then there exists a constant C such that,

$$f(x) = g(x) + C \quad \forall x \in (a, b).$$

The Second fundamental Theorem of Calculus:

If f is a continuous function on $[a, b]$ and if,

$$\phi(x) = f(x), \quad \forall x \in (a, b)$$

$$\int_a^b f(x) dx = \phi(b) - \phi(a).$$

Theorem: 2

Let ϕ be a function on $[a, b]$ such that ϕ is continuous on $[a, b]$ let $A = \phi(a)$ and $B = \phi(b)$, then if f is continuous on $\phi([a, b])$ we have,

$$\int_A^B f(x) dx = \int_a^b f[\phi(u)] \phi'(u) du.$$

Improper Integrals:

Definition:

Suppose $f \in R[a, s]$ for every $s > a$, we define

$$F(s) = \int_a^s f(x) dx$$

$\int_a^\infty f(x) dx$ is convergent to A ,

$$\text{if } \lim_{s \rightarrow \infty} F(s) = A$$

If $\int_a^\infty f(x) dx$ does not converge, we say it diverges.

Definition:

If $f \in R[a, s]$ for every $s > a$ and if $\int_a^\infty |f(x)| dx$ converges, we say that, $\int_a^\infty f$ converges absolutely

If $\int_a^\infty f$ converges but $\int_a^\infty |f|$ diverges we say that converges conditionally.

Definition:

We define $\int_{-\infty}^a f(x) dx = \lim_{s \rightarrow -\infty} \int_s^a f(x) dx$ provided the limits exists. If the limit exists we say $\int_{-\infty}^a f$ converges and if the limits does not exist we say the $\int_{-\infty}^a f$ diverges.

Improper Integral (continued) :

Definition:

Suppose $f \in R[a+\epsilon, b]$ for every ϵ such that $0 < \epsilon < b-a$ but $f \notin R[a, b]$,

$$\text{Define } f(\epsilon) = \int_{a+\epsilon}^b f(x) dx.$$

We say the Improper integral,

$\int_a^b f(x) dx$ converges to A , if

$$\lim_{\epsilon \rightarrow 0} f(\epsilon) = A.$$

If the limit does not exist we say $\int_a^b f$ diverges.

An improper integral of the form $\int_a^b f$ is called an improper integral of second kind.

Real Analysis And Matrix Theory

Unit - 2

1) Rank of Matrix :

Let, $A = (a_{ij})_{m \times n}$ ($i = 1, 2, \dots, m, j = 1, 2, \dots, n$)

be a matrix of order $m \times n$.

Where $a_{ij} \in \mathbb{R}$, \mathbb{R} is a real field (or complex field) i.e.,

$$A = \begin{bmatrix} a_{11} & a_{12} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & a_{2n} \\ \vdots & & & \\ a_{m1} & a_{m2} & \dots & a_{mn} \end{bmatrix}$$

The whole array can be imagined as a set of m row vectors or as a set of n column vectors, the number of linearly independent rows or columns is called the rank of the matrix.

2) Elementary Transformation of a Matrix :

$E_1 \rightarrow$ Replacing a row/column say r th row of A by a scalar multiple of the same, In notation, $R_r \rightarrow CR_r$ or $C_r \rightarrow CC_r$ where R_r and C_r stands for r th row and r th column respectively.

$E_2 \rightarrow$ Replacing the r th row/column plus λ times the s th row/column, i.e., $R_r \rightarrow R_r + \lambda R_s$ (or) $C_r \rightarrow C_r + \lambda C_s$

$E_3 \rightarrow$ Interchanging r th row/column with s th row/column i.e., $R_r \leftrightarrow R_s$ (or) $C_r \leftrightarrow C_s$

3) EQUIVALENT MATRICES :

Any matrix obtained from a given matrix by performing a series of elementary transformations is called as equivalent to the given matrix.

4) Elementary Matrices :

A Matrix obtained from a performing a single elementary transformation on a unit matrix is called an elementary matrix.

Problem :

$$A = \begin{pmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \\ 5 & 7 & 9 \end{pmatrix}$$

$$e_{23} = \lambda$$

$$r = 2, s = 3$$

$$E_{rs}(A) A = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & \lambda \\ 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \\ 5 & 7 & 9 \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 \\ 4 + 5\lambda & 5 + 7\lambda & 6 + 9\lambda \\ 5 & 7 & 9 \end{pmatrix}$$

which consists of replacing the second row to second row + λ times of the row.

$E_{rs}(A)$ is non-singular

which consists the second column + λ times of the third of third column.

ECHELON MATRIX :

A Matrix A is said to be echelon if elements in a row are either zero or a unity preceded by zeros (but not necessarily succeeded by zeros) and if in a column there is a unity which is the first unity of a row then all the other elements of the column are zeros,

$$E = \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & \delta_1 & 0 & \delta_2 \\ 0 & 0 & 1 & \delta_3 \\ 0 & 0 & 0 & 0 \end{pmatrix}$$

HERMITE CANONICAL FORM :

A Matrix H of order $m \times m$ is said to be in Hermite Canonical form if its principal diagonal consists of only zeros or unities and all elements below the diagonal are zero, such that when a diagonal element is zero, the entire row is zero and when a diagonal element is unity, the rest of the elements in the column are all zeros.

$$H = [\delta_{ij}]_{m \times m}$$

$$\delta_{ij} = 0 \quad \text{and} \quad \delta_{ki} = 0 \quad \forall k < i \quad (i, k = 1, 2, \dots, m)$$

Subject to,

$$(i) \quad \delta_{ii} = 0 \Rightarrow \delta_{ij} = 0 \quad \forall j$$

$$(ii) \quad \delta_{ii} = 1 \Rightarrow \delta_{ij} = 0 \quad \forall i$$

PROBLEM:

$H = \begin{pmatrix} 1 & \delta_{12} \\ 0 & 0 \end{pmatrix}$ is a Hermitian canonical form
which is an idempotent matrix.

$$H^2 = \begin{pmatrix} 1 & \delta_{12} \\ 0 & 0 \end{pmatrix} \begin{pmatrix} 1 & \delta_{12} \\ 0 & 0 \end{pmatrix} = \begin{pmatrix} 1 & \delta_{12} \\ 0 & 0 \end{pmatrix} = H$$

FROBENIUS INEQUALITY:

Theorem:

If A, B, C are any three matrices such that products considered are defined, then.

$$P(AB) + P(BC) \leq P(ABC) + P(B)$$

Proof: $P(BC) \geq P(ABC) \quad \rightarrow (1)$

$$P(B) \geq P(AB) \quad \rightarrow (2)$$

$$P(B) \geq P(BC), \quad P(AB) \geq P(ABC).$$

Sub (1) from (2) we get,

$$P(B) - P(BC) \geq P(AB) - P(ABC)$$

$$\Rightarrow P(B) + P(BC) \geq P(AB) + P(BC)$$

$$\Rightarrow P(AB) + P(BC) \leq P(ABC) + P(B).$$

Hence Proved.

(CERTAIN RESULTS ON THE RANK OF AN IDEMPOTENT MATRIX :

A Matrix $(A)_{n \times n}$ is said to be idempotent if $A^2 = A$.

I. The rank of an idempotent matrix A is equal to its trace (i.e. sum of the diagonal elements).

$P(A)$ be the rank of A

$$P(A) = r$$

\exists two non-singular matrices P and $Q \Rightarrow$

$$PAQ = \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix}$$

$$\Rightarrow A = P^{-1} \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix} Q^{-1}$$

Let us denote $\begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix} = [I_r, 0]_{n \times n}$.

$$\begin{aligned} A^2 = A &= P^{-1} I_{r,0} Q^{-1}, P^{-1} I_{r,0} Q^{-1} = P^{-1} I_{r,0} Q^{-1} \\ &= PP^{-1} I_{r,0} Q^{-1} P^{-1} I_{r,0} Q^{-1} Q = PP^{-1} I_{r,0} Q^{-1} Q = \\ & \quad I_{r,0}. \end{aligned}$$

$$= I_n I_{r,0} Q^{-1} P^{-1} I_{r,0} I_n = I_n I_{r,0} I_n$$

$$= I_{r,0} (Q^{-1} P^{-1} - I_n) I_{r,0} = 0$$

$$= Q^{-1} P^{-1} = I_n$$

$$= Q^{-1} = P$$

$$PAP^{-1} = \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix}$$

$$PAP^{-1} = \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix}$$

$$\begin{aligned} \text{tr}(PAP^{-1}) &= \text{tr}(P^{-1}PA) = \text{tr}(A) \\ &= \text{tr} \begin{pmatrix} I_r & 0 \\ 0 & 0 \end{pmatrix} = r = \rho(A) \end{aligned}$$

Problems:

$$(i) \quad A = \begin{pmatrix} 1 & 2 \\ 3 & 4 \end{pmatrix}$$

E_1 consists of multiplying the 2nd row by 2.

Here, $r=2$; $c=2$.

$$\underbrace{\begin{pmatrix} 1 & 0 \\ 0 & 2 \end{pmatrix}}_{= \Delta} \begin{pmatrix} 1 & 2 \\ 3 & 4 \end{pmatrix} = \begin{pmatrix} 1 & 2 \\ 6 & 8 \end{pmatrix}$$

$\Delta = [\delta_{ij}]$ is non-singular.

(ii) Suppose we want to replace the second column by 3 times the second column.

$$C=3 ; r=2$$

$$\begin{pmatrix} 1 & 2 \\ 3 & 4 \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & 3 \end{pmatrix} = \begin{pmatrix} 1 & 6 \\ 3 & 12 \end{pmatrix}$$

E_2 : Replacing any row (column) say r th row (column) by r th row/column any other row/column say s th row/column.

In this case we require to pre or post multiply A by a non-replacement.

$E_{rs}(A)$ where $E_{rs}(A)$ is defined by,

$$E_{rs}(A) = [e_{ij}]$$

$$e_{ii} = 1 \quad \forall i$$

$$e_{ij} = 0 \quad i \neq j$$

$$E_{rs} = A$$

which consists of replacing of the rows and columns by the singular matrix.

Example :

$$A = \begin{pmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \end{pmatrix} \quad r=1, s=2.$$

$$E_{sr}(-1) = \begin{pmatrix} 1 & 0 \\ -1 & 1 \end{pmatrix}$$

$$E_{rs}(1) = \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}$$

$$E_{sr}(-1) = \begin{pmatrix} 1 & 0 \\ -1 & 1 \end{pmatrix}$$

$$\Delta = \begin{pmatrix} -1 & 0 \\ 0 & 1 \end{pmatrix}$$

$$E_{sr}(-1) E_{sr}(1) E_{sr}(-1) \Delta$$

$$= \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix} \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix} \begin{pmatrix} 1 & 0 \\ -1 & 1 \end{pmatrix} \begin{pmatrix} -1 & 0 \\ 0 & 1 \end{pmatrix}$$

$$= \begin{pmatrix} +0 & 1 \\ 1 & 0 \end{pmatrix}$$

$$E_{sr}(-1) E_{rs}(1) E_{sr}(-1) \Delta$$

$$= \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \end{pmatrix}$$

$$= \begin{pmatrix} 4 & 5 & 6 \\ 1 & 2 & 3 \end{pmatrix}.$$

Unit - 4

Eigen Values and Eigen Vectors :

Let,

$$(A)_n \cdot n(x)_{n-1} = \lambda (x)_{n-1}$$

Then λ is called the characteristic root (or eigen value or the latent root) of A and x is called the characteristic vector (or eigen vector) of A or the right characteristic vector of A .

If the A is non-symmetric (or non-Hermitian) then $A^{-1} \neq A$ and if it satisfies from some scalar λ .

$$A^{-1}x = \lambda x$$

Then x will be called the left characteristic vector corresponding to λ . If A is a symmetric matrix then left and right characteristic vectors are the same.

$$Ax = \lambda x$$

$$(A - \lambda I)x = 0.$$

is a system of n linear and homogenous equation.

$|A - \lambda I|$ is called a characteristic polynomial of A . If $A =$

$(a_{ij})_{n \times n}$ then,

$$|A - \lambda I| = \begin{vmatrix} a_{11} - \lambda & a_{12} & \dots & a_{1n} \\ a_{21} & a_{22} - \lambda & \dots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{n1} & a_{n2} & \dots & a_{nn} - \lambda \end{vmatrix}$$

$$= (-\lambda)^n + (-\lambda)^{n-1} \sigma_1 + (-\lambda)^{n-2} \sigma_2 + \dots + a_n$$

Where, $\sigma_0 = 1$, $\sigma_n = |A|$

σ_i = Sum of the principal minor of order i in A .

$i = 1, 2, \dots, n-1$.

Problem :

Show that a skew symmetric matrix has no real characteristic roots other than zero.

A is a skew symmetric matrix

$$AX = \lambda X \text{ holds}$$

$$iAX = (i\lambda)X.$$

If A is a skew symmetric then iA is Hermitian whose characteristic roots are real $i\lambda$ is real λ is imaginary unless $\lambda = 0$.

the characteristic root of A is same as the characteristic root of A' .

$$|A - \lambda I| = 0 = |A' - \lambda I| = 0$$

$$= -A - \lambda I = 0 \text{ if } A \text{ is Skew Symmetric}$$

$$= |A + \lambda I| = 0$$

$$= \lambda^n + \lambda^{n-1}\sigma_1 + \lambda^{n-2}\sigma_2 + \dots + \underbrace{|A|}_{=\sigma_n} = 0$$

$$(-\lambda)^n + (-\lambda)^{n-1}\sigma_1 + (-\lambda)^{n-2}\sigma_2 + \dots + \underbrace{(-1)^{n-n}|A|}_{=|A|} = 0$$

This shows that if $\lambda = \lambda_0$ is a root then $\lambda = -\lambda_0$ is also a root of the characteristic equation which implies that either λ_0 is imaginary or zero.

CAYLEY HAMILTON THEOREM:

Every Matrix Satisfies its characteristic equation.

Proof: We have,

$$[A - \lambda I] = \frac{\text{adj}[A - \lambda I]}{[A - \lambda I]} = I_n$$

$$= [A - \lambda I] \text{adj}[A - \lambda I] = [A - \lambda I]$$

$$\text{adj}[A - \lambda I] = C_0 + C_1\lambda + \dots + C_{n-1}\lambda^{n-1}$$

C_i 's are matrices,

$$[A - \lambda I] = a_0 + a_1 \lambda + a_2 \lambda^2 + \dots + a_n \lambda^n$$

a_i 's are scalars,

$$\begin{aligned} [A - \lambda I] (C_0 + C_1 \lambda + \dots + C_{n-1} \lambda^{n-1}) \\ = a_0 + a_1 \lambda + \dots + a_n \lambda^n \end{aligned}$$

Equating Co-efficients of $\lambda, \lambda^2, \dots$ etc on both sides we have,

$$A C_0 = a_0$$

$$A C_1 - C_0 I = a_1 = A^2 C_1 - A C_0 I = A a_1$$

$$= A^{n-1} C_{n-2} - A^{n-2} C_{n-3} I = A^{n-2} a_{n-2}$$

$$= A C_{n-1} - C_{n-2} I = a_{n-1}$$

$$= A^n C_{n-1} - A^{n-1} C_{n-2} I = A^{n-1} a_{n-1}$$

$$= - C_{n-1} I = a_n$$

$$= - A^n C_{n-1} = A^n a_n$$

Adding up, we have,

$$a_0 + a_1 A + a_2 A^2 + \dots + a_{n-1} A^{n-1} + a_n A^n = 0.$$

Application of Cayley Hamilton Theorem :

We have if,

$$|A - \lambda I| = a_0 = a_1 \lambda + a_2 \lambda^2 + \dots + a_n \lambda^n = 0$$

$$a_0 + a_1 A + a_2 A^2 + \dots + a_n A^n = 0$$

For a non-singular matrix, Pre multiplying both sides by A^{-1} ,

$$a_0 A^{-1} + a_1 I + a_2 A + \dots + a_n A^{n-1} = 0.$$

The equation can be employed to obtain A^{-1} when the Powers of A are known.

Example :

$$A = \begin{bmatrix} 1 & 2 \\ 1 & 3 \end{bmatrix}; \text{ ch equation}$$

$$\lambda^2 - (1+3)\lambda + 1 = 0$$

$$\lambda^2 - 4\lambda + 1 = 0$$

$$\Rightarrow A^2 - 4A + I = 0$$

$$= A - 4I + A^{-1} = 0$$

$$A - 4I + A^{-1} = 0$$

$$\begin{bmatrix} 1 & 2 \\ 1 & 3 \end{bmatrix} - 4 \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix} = -A^{-1}$$

$$= \begin{bmatrix} 1 & -4 & 2 & -0 \\ 1 & -0 & 3 & -0 \end{bmatrix} = -A^{-1}$$

$$= \begin{bmatrix} 3 & -2 \\ -1 & -3 \end{bmatrix} = A^{-1}.$$

Example Problems:

[S is a real -skew symmetric matrix $\Rightarrow (I-S)$ is non-singular and $A = (I+S)(I-S)^{-1}$ is orthogonal].

Proof: $|I-S| \neq 0$

$$\begin{aligned} A^{-1} &= [(I+S)(I-S)^{-1}]^{-1} \\ &= [(I-S)^{-1}]^{-1} (I+S)^{-1} \\ &= (I+S) (I-S) \end{aligned}$$

$$Q^{-1}BP^{-1} = \left[\begin{array}{c|c} C_{11} & C_{12} \\ \hline C_{21} & C_{22} \end{array} \right]$$

Partitioned in the same way as $\left[\begin{array}{c|c} I_r & 0 \\ \hline 0 & 0 \end{array} \right]$

$$PAQ Q^{-1}BP^{-1} = \left[\begin{array}{c|c} I_r & 0 \\ \hline 0 & 0 \end{array} \right] \left[\begin{array}{c|c} C_{11} & C_{12} \\ \hline C_{21} & C_{22} \end{array} \right]$$

$$\left[\begin{array}{c|c} C_{11} & C_{12} \\ \hline 0 & 0 \end{array} \right]$$

It follows that the characteristic roots of AB are either those of C_{11} or zeros.

$Q^{-1}BAQ$ where Q is non-singular

$$= \underbrace{Q^{-1}BP^{-1}} \underbrace{PAQ}$$

$$\left[\begin{array}{c|c} C_{11} & C_{12} \\ \hline C_{21} & C_{22} \end{array} \right] \left[\begin{array}{c|c} I_r & 0 \\ \hline 0 & 0 \end{array} \right]$$

$$= \left[\begin{array}{c|c} C_{11} & 0 \\ \hline C_{21} & 0 \end{array} \right]$$

Again the characteristic roots of BA are either those of C_{11} or zeros. Hence the characteristic roots of AB and BA are identical.

$$\begin{aligned}
A^T A &= (I+S)^{-1} (I-S) (I+S) (I-S)^{-1} \\
&= (I+S)^{-1} (I^2 - S^2) (I-S)^{-1} \\
&= (I+S)^{-1} (I+S) (I-S) (I-S)^{-1} \\
&= I \quad (\because A \text{ is orthogonal}).
\end{aligned}$$

Problem :

Show that AB and BA have same characteristic roots :

Proof : A and its transform $P^{-1}AP$ (P is a non-singular)

$$\begin{aligned}
AX &= \lambda X \\
&= P^{-1}AX = \lambda P^{-1}X \\
&= P^{-1}APP^{-1}X = \lambda P^{-1}X \\
&= X(P^{-1}X).
\end{aligned}$$

which shows that $P^{-1}AP$ and A have the same characteristic roots.

$$P(A) = 0, \quad PAQ = \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix}$$

$P(AB)P^{-1}$ which has the characteristic roots.

$$P(AB)P^{-1} = PAQQ^{-1}BP^{-1}.$$

Q is non-singular matrix of the same order such that,

Generalised Inverse of Matrix :

$$\text{Let, } (A)_{m \times n} (X)_{n \times 1} = (Y)_{m \times 1}$$

be a consistent system of equation. A generalised inverse (or a g-inverse) of A is a $n \times m$ matrix denoted by A^- (or A^g)

Such that,

$X = A^- Y$ (or $X - A^- Y$) is a solution of the equation.

$$Y \in R(A).$$

Example : 1

$$(A)_{m \times n} \text{ and } P(A) = r.$$

\exists a non-singular sub matrix $(B)_{r \times r} \ni P(B) = r$

By a suitable interchange of rows and columns of A ,

$$A = \begin{pmatrix} B & C \\ D & E \end{pmatrix}$$

$$A^- = \begin{pmatrix} B^{-1} & 0 \\ 0 & 0 \end{pmatrix} \text{ then one may show}$$

$$\left. \begin{array}{l} AA^-A = A \\ A^-AA^- = A^- \end{array} \right\} P(A^-) = P(A)$$

A^- is a reflexive generalised inverse. It is also easy to see $E = DB^-C$.

Example : 2

Obtain a g-inverse of the matrix given by.

$$A = \begin{bmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \\ 7 & 8 & 9 \end{bmatrix}$$

Sol : It may be seen that $\rho(A) = 2$

$$|A| = 0$$

$\rho(A) = 2$ since \exists one

$$\text{Minor of } A \text{ viz } \begin{vmatrix} 1 & 2 \\ 4 & 5 \end{vmatrix} \neq 0$$

A generalised inverse of A is obtained by arbitrary values

of U, V and W .

$$\begin{bmatrix} \begin{pmatrix} 1 & 2 \\ 4 & 5 \end{pmatrix}^{-1} (U)_{2 \times 1} \\ V' \quad W' \end{bmatrix}$$

$$U = 0, V' = 0, W = 0$$

We have, $A^{\#} = \begin{bmatrix} -5/3 & 2/3 & 0 \\ 4/3 & -1/3 & 0 \\ 0 & 0 & 0 \end{bmatrix}$

DIFFERENT CLASSES OF GENERALISED INVERSE :

We define matrices over a complex number field clearly analogous results are obtainable when the matrices are defined over a real field.

$$AXA = A$$

$$XAX = X$$

$$(XA)^* = XA$$

$$(AX)^* = AX$$

* = Conjugate transpose.

Example :

Show that,

$$A^+ = (A^*A)^{-1} A^* \text{ if } (A^*A)^{-1} \text{ exists.}$$

$$B^+ = B^* (BB^*)^{-1} \text{ if } (BB^*)^{-1} \text{ exists.}$$

Sol: Put $B = I$ in solution,

$$(AI)^+ = I^* (II^*)^{-1} (A^*A)^{-1} A^*$$

$$\Rightarrow A^+ = (A^*A)^{-1} A^*$$

$$A = I.$$

$$(IB^+) = B^* (BB^*)^{-1} (I^*I)^{-1} I^*$$

$$\Rightarrow B^+ = B^* (BB^*)^{-1}.$$

PROPERTIES OF G -INVERSE ($A\#$)

- (a) $P(A\#) \supseteq P(A)$
- (b) One choice of $(A^*)\#$ is $(A\#)^*$
- (c) $A(A^*A)\#A^*A = A$
- (d) $A^*A(A^*A)\#A^* = A^*$
- (e) If A is Idempotent Matrix
 $A\# = A$ is a choice of $A\#$
- (f) $A(A^*A)\#A^*$ is a Hermitian idempotent matrix
Projecting vectors onto the subspace spanned by the
columns of A .
- (g) $A\#A$ and $AA\#$ are Hermitian.

PROPERTIES OF G -INVERSE ($A\#$)

- (a) All the Properties from (b) to (g) for generalised -
inverse
- (b) $A\#$ is reflective g -inverse if and only if
 $P(A\#) = P(A)$.

PROPERTIES OF LEFT PSEUDO-INVERSE ($A^{\#}$):

- (a) Properties (c), (d), (f) and of replacing A with A^*
- (b) One choice for $A^{\#}$ is $A^*(AA^*)^{\#}$ and for $(A^*)^{\#}$ is $(A^{\#})^*$ $a \neq 0$ $(A^*A)^{\#}$ is $A^{\#}(A^*)^{\#}$
- (c) One choice for $(UAV)^{\#}$ is $V^{\#}A^{\#}U^{\#}$
Where U and V are unitary matrices.
- (d) $A + A^{\#}A = A^{\#}A$
for any choice of $A^{\#}$.
- (e) If B_1 and B_2 are any two choices of $A^{\#}$ then
 $(B_1 - B_2)A = 0$
- (f) If $\mathcal{P}(A) = n = \min(m, n)$
Then, $\forall A^{\#}$ is $A^{\#}$.
- (g) If A is Hermitian and idempotent of the matrix.

PROPERTIES OF MOORE PENROSE PSEUDO INVERSE:

The following properties are true for A^{\dagger} of A (A^{\dagger}).

- (a) Properties (c), (d), (f), (g) of A^{\dagger} in Replacing of A^{\dagger}
- (b) One choice for A^{\dagger} is $(A^*A)^{\dagger} A^*$
- (c) One choice for $(AA^*)^{\dagger}$ is $(A^*)^{\dagger} A^{\dagger}$.
- (d) One choice for $(UAV)^{\dagger}$ is $V^*A^{\dagger}U^*$ where U and V are unitary matrices.
- (e) $AA^{\dagger} = AA^{\dagger}$ for any choice of A^{\dagger} .
- (f) If B_1 and B_2 are any two choices of A^{\dagger} then $A(B_1 - B_2) = 0$.
- (g) If $r(A) = m = \min(m, n)$ then every A^{\dagger} is A^{\dagger} .
- (h) If A is hermitian and idempotent one choice for A^{\dagger} is A .

APPLICATIONS OF GENERALISED INVERSES IN THE SOLUTION OF SYSTEMS OF LINEAR EQUATIONS :

Theorem :

A Necessary and Sufficient condition for the equation
 $(A)XB = C$

to have a solution is

$A(A^{\dagger}CB^{\dagger})B = C$ in which case the general solution is

$$X = A^{\dagger}CB^{\dagger} + Y - A^{\dagger}AYBB^{\dagger}$$

Y is arbitrary.

Proof : $AXB = C$

$$C = AXB = AA^{\dagger}A X (BB^{\dagger}B)$$

$$\Rightarrow C = AA^{\dagger}CB^{\dagger}B \quad (\because AXB = C)$$

$$AXB = 0$$

$$AXB = A(Y - A^{\dagger}AYBB^{\dagger})B$$

$$= AYB - AA^{\dagger}AYBB^{\dagger}B$$

$$= AYB - AYB = 0.$$

The general solution of the non-homogenous equation -

$$AXB = C$$

\Rightarrow general Solution of $AXB = 0$ + a particular solution of $AXB = C$

Hence the general of $AXB = C$ is given by

$$X = Y - A^+ A Y B B^+ + A^+ C B^+ \text{ (Q.E.D.)}$$

Theorem 6

$$\text{If } P = \begin{pmatrix} P_1 \\ P_2 \end{pmatrix} P^{-1} = [X, Y]$$

$$Q = [Q_1, Q_2] Q^{-1} = \begin{bmatrix} T_1 \\ T_2 \end{bmatrix}$$

$$(a) P_1 X = I, P_2 X = 0.$$

(b) $X Q_1 = I, X Q_2 = 0$ have a unique common solution.

$$X = (Q_1 - Q_2 Q_2^+ Q_1) + (I - Q_2 Q_2^+).$$

Proof :

$$P^{-1} = [X, Y]$$

$$P P^{-1} = \begin{bmatrix} P_1 \\ P_2 \end{bmatrix} [X, Y] = \begin{bmatrix} P_1 X & P_1 Y \\ P_2 X & P_2 Y \end{bmatrix}$$

$$P_1 X = I \text{ and } P_2 X = 0$$

$$= I \text{ (identically)}$$

$$P_2 X = (P_2 - P_2 P_2^+ P_2) (P_1 - P_2^+ P_2) +$$

$$= 0.$$

Similarly it can easily be verified that the common solution of $XQ_1 = I$ and $XQ_2 = 0$.

is given by

$$X = (Q_1 - Q_2 Q_1 + Q_1) + (I - Q_2 Q_1)$$

Quadratic forms :

If X is a column vector of order $n \times 1$ and there is a symmetric matrix $A = [a_{ij}]_{n \times n}$ i.e, $a_{ij} = a_{ji}$

Then $Q = (X)'_{1 \times n} (A)_{n \times n} (X)_{n \times 1}$ is called a Quadratic form.

Classification of the Quadratic form :

If we apply linear transformation say $X = PY$ then

$$Q = Y' (P'AP) Y$$

$$P'AP = I_n$$

where P is a non-singular matrix

$$P(A) = n, \text{ then } Q = \sum_{i=1}^n y_i^2$$

$Q > 0$, unless all y 's are zero and y 's cannot be zero unless all x 's are zero.

Positive Semi-Definite Quadratic form:

If $P(A) = r < n$ and if we apply the same linear transformation $X = PY$, P being a non-singular matrix

Then

$$\begin{aligned} Q &= X^T A X \\ &= Y^T P^T A P Y \\ &= Y^T I_r Y \quad (\because P(A) = r) \\ &= \sum_{i=1}^r Y_i^2 \geq 0 \end{aligned}$$

is called a Positive Semi-definite Quadratic form.

For eg: Correlation matrix A is Positive definite or Semi definite matrix in case of multinomial distribution

is positively Semi-definite.

When Q transformed in the form,

$$Q = - \sum_{i=1}^r Y_i^2 < 0.$$

then Q is called Negative Semi-definite Quadratic form.

The Number of Positive terms S is called the index of the Quadratic forms.

Compute a g -inverse:

$$\begin{pmatrix} 18 & 2 & 46 \\ 2 & 1 & 2 \\ 46 & 2 & 130 \end{pmatrix}$$

Solution: Note $P(A) = 2$

$$\begin{bmatrix} 18 & 2 \\ 2 & 1 \end{bmatrix}^{-1} = \frac{1}{14} \begin{bmatrix} 1 & -2 \\ -2 & 18 \end{bmatrix}$$

$$\therefore A^g = \frac{1}{14} \begin{pmatrix} 1 & -2 & 0 \\ -2 & 18 & 0 \\ 0 & 0 & 0 \end{pmatrix}$$

2) Compute the right weak g -inverse (A^n) of A

given by, $A = \begin{bmatrix} 1 & 2 \\ 2 & 1 \\ 1 & 1 \\ 1 & -1 \end{bmatrix}_{4 \times 2}$

Solution: $P(A) = 2$, $\min(m, n) = 2$

A right weak generalised inverse = $(A^*A)^g A^*$

$$= \left(\begin{bmatrix} 1 & 2 & 1 & 1 \\ 2 & 1 & 1 & -1 \end{bmatrix}_{2 \times 4} \begin{bmatrix} 1 & 2 \\ 2 & 1 \\ 1 & 1 \\ 1 & -1 \end{bmatrix}_{4 \times 2} \right)$$

$$= \begin{pmatrix} 7 & 4 \\ 4 & 7 \end{pmatrix}^{-1} \begin{pmatrix} 1 & 2 & 1 & 1 \\ 2 & 1 & 1 & -1 \end{pmatrix}$$

$$= \frac{1}{33} \begin{bmatrix} 7 & -4 \\ -4 & 7 \end{bmatrix} \begin{bmatrix} 1 & 2 & 1 & 1 \\ 2 & 1 & 1 & -1 \end{bmatrix}$$

$$= \frac{1}{33} \begin{bmatrix} -1 & 10 & 3 & 11 \\ 10 & -1 & 3 & -11 \end{bmatrix}$$

$$\therefore A^{-1} = \frac{1}{33} \begin{bmatrix} -1 & 10 & 3 & 11 \\ 10 & -1 & 3 & -11 \end{bmatrix}.$$