

# **Mathematical Methods for Economic Dynamics (ECON 696)**

## **Lecture 2: Compactness, Continuity and Linear Algebra**

## Boundedness in terms of Sequences on $\mathbb{R}^n$ (Ch. 13.2)

**Definition:** A set  $S$  in  $\mathbb{R}^n$  is **bounded** if there exists a number  $M \in \mathbb{R}_+$  such that  $\|x\| \leq M$  for all  $x \in S$ . A set that is not bounded is called **unbounded**. Here  $\|x\| = d(x, \mathbf{0}) = \sqrt{x_1^2 + \cdots + x_n^2}$ , called the **Euclidean norm**.

Similarly, a sequence  $\{x^k\}$  in  $\mathbb{R}^n$  is **bounded** if the set  $\{x^k | k \in \mathbb{N}\}$  is bounded.

**Lemma:** Any convergent sequence  $\{x^k\}$  in  $\mathbb{R}^n$  is bounded.

**Proof:** We skip the proof. ■

On the other hand, a bounded sequence  $\{x^k\}$  in  $\mathbb{R}^n$  is “not” necessarily convergent. This is the same as sequences in  $\mathbb{R}$ .

Hence, we obtain

**Theorem:** A subset  $S$  of  $\mathbb{R}^n$  is bounded if and only if every sequence of points in  $S$  has a convergent subsequence.

**Proof:** We skip the proof. ■

# Compactness

## Compactness (Ch. 13.2)

**Bolzano-Weierstrass Theorem:** A subset  $S$  of  $\mathbb{R}^n$  is closed and bounded if and only if every sequence of points in  $S$  has a subsequence that converges to a point in  $S$ .

**Proof of Bolzano-Weierstrass's Theorem:** We skip this. ■

**Definition:** A set  $S$  in  $\mathbb{R}^n$  is **compact** if it is closed and bounded.

# Continuity

## Continuous Functions (Ch. 13.3)

Roughly speaking,  $f(\cdot)$  is continuous if small changes in the independent variables cause only small changes in the function value.

In what follows, let  $S$  be a nonempty subset of  $\mathbb{R}^n$ .

**Theorem:** A function  $f$  from  $S$  into  $\mathbb{R}$  is **continuous** at a point  $x^0$  in  $S$  if and only if  $(\forall \{x^k\} \in S)(x^k \rightarrow x^0 \Rightarrow f(x^k) \rightarrow f(x^0))$ . If  $f(\cdot)$  is continuous at every point in a set  $S$ ,  $f(\cdot)$  is continuous on  $S$ .

**Proof:** We skip the proof. ■

**Example:** Let  $f : \mathbb{R}_+ \rightarrow \mathbb{R}$  be given below.

$$f(x) = \begin{cases} 1 & \text{if } x \geq 1 \\ 0 & \text{if } x < 1 \end{cases}$$

Let  $x^0 = 1$ .

If  $x$  converges to 1 “from below,” we have  $\lim_{x \rightarrow 1^-} f(x) = 0$ .

And, if  $x$  converges to 1 “from above,” we have  $\lim_{x \rightarrow 1^+} f(x) = 1$ .

Therefore,  $f(\cdot)$  is **not** a continuous function at  $x^0 = 1$ .



## Continuity of Vector-Valued Functions (Ch. 13.3)

Let  $f = (f^{(1)}, \dots, f^{(m)})$  be a function from a subset  $S$  to  $\mathbb{R}^m$ .

**Theorem:** A function  $f = (f^{(1)}, \dots, f^{(m)})$  from  $S \subseteq \mathbb{R}^n$  to  $\mathbb{R}^m$  is continuous at a point  $x^0$  in  $S$  if and only if each component function  $f^{(j)} : S \rightarrow \mathbb{R}$ ,  $j = 1, \dots, m$ , is continuous at  $x^0$ .

**Proof:** We omit the proof. ■

## Two Preliminary Results for “Weierstrass Theorem” (Ch. 3.1)

**Theorem:** Let  $S \subseteq \mathbb{R}^n$  and let  $f : S \rightarrow \mathbb{R}^m$  be continuous. Then  $f(K) = \{f(x) | x \in K\}$  is compact for every compact subset  $K$  of  $S$ .

**Proof:** We omit the proof. ■

**Theorem:** Let  $S$  be a compact set in  $\mathbb{R}$  and let  $x_*$  be the greatest lower bound of  $S$  and  $x^*$  be the least upper bound of  $S$ . Then,  $x_* \in S$  and  $x^* \in S$ .

**Proof:** We skip the proof. ■

## Weierstrass (Extreme Value) Theorem (Ch. 3.1 and 13.3)

**Weierstrass Theorem:** Let  $f : S \rightarrow \mathbb{R}$  be a continuous real-valued mapping where  $S$  is a nonempty compact subset of  $\mathbb{R}^n$ . Then there exist two vectors  $x^*, x_* \in S$  such that for all  $x \in S$ ,

$$f(x_*) \leq f(x) \leq f(x^*).$$

**Proof:** The first theorem shows that  $f(S)$  is a nonempty compact set. The second theorem shows that there exist  $y^*, y_* \in f(S)$  such that  $y_* \leq y \leq y^*$  for any  $y \in f(S)$ . The rest of the proof is completed by finding  $x_*, x^* \in S$  such that  $y_* = f(x_*)$  and  $y^* = f(x^*)$ . ■

## Why are the compactness and the continuity needed

(1) Let  $S = [0, \infty)$  and  $f(x) = x$ . Then  $f(\cdot)$  cannot attain a maximum because  $S$  is not bounded from above. But  $S$  is closed and  $f(\cdot)$  is continuous.

(2) Let  $S = (0, 1)$  and  $f(x) = x$ . Then  $f(\cdot)$  cannot attain a maximum or minimum because  $S$  is not closed. But  $S$  is bounded and  $f(\cdot)$  is continuous.

(3) Let  $S = [0, 1]$ . Define  $f : S \rightarrow \mathbb{R}$  as follows:

$$f(x) = \begin{cases} 1 - x & \text{if } x \in (0, 1] \\ 0 & \text{if } x = 0 \end{cases}$$

This  $f(\cdot)$  fails to attain a maximum because  $f(\cdot)$  is not continuous at  $x = 0$ . But  $S$  is compact.

# **Linear Algebra (Ch. 1)**

## Basic Concepts of Matrix Algebra in $\mathbb{R}^n$ (Ch. 1.1)

An  $m \times n$  **matrix** is a rectangular array with  $m$  rows and  $n$  columns:

$$A = (a_{ij})_{m \times n} = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{pmatrix}$$

Here  $a_{ij}$  denotes the elements in the  $i$ th row and the  $j$ th column.

A matrix can be considered a generalization of numbers. We want to understand to what extent we can treat matrices like numbers.

## Sum and Subtraction of Matrices

If  $A = (a_{ij})_{m \times n}$ ,  $B = (b_{ij})_{m \times n}$ , and  $\alpha \in \mathbb{R}$  is a scalar, I define

- $A + B = (a_{ij} + b_{ij})_{m \times n}$ ,
- $\alpha A = (\alpha a_{ij})_{m \times n}$ ,
- $A - B = A + (-1)B = (a_{ij} - b_{ij})_{m \times n}$ .



The **zero matrix**  $\mathbf{0}_{m \times n}$  is defined as a matrix where all entries are zero:

$$\mathbf{0}_{m \times n} = \begin{pmatrix} 0 & 0 & \cdots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 0 \end{pmatrix}.$$

## Multiplications of Matrices

$A$ :  $m \times n$  matrix; and  $B$ :  $n \times p$  matrix.

$C = AB$  :  $m \times p$  matrix  $C = (c_{ij})_{m \times p}$  such that for each  $(i, j)$ ,

$$c_{ij} = \sum_{r=1}^n a_{ir}b_{rj} = \underbrace{a_{i1}b_{1j} + a_{i2}b_{2j} + \cdots + a_{ik}b_{kj} + \cdots + a_{in}b_{nj}}_{n \text{ terms}}$$

$AB$  is well-defined only if the number of columns in  $A$  is equal to the number of rows in  $B$ .

### Example of Multiplication of Matrices:

$$\begin{pmatrix} 1 & 0 & 3 \\ 2 & 1 & 5 \end{pmatrix} \begin{pmatrix} 1 & 3 \\ 2 & 5 \\ 6 & 2 \end{pmatrix} = \begin{pmatrix} 1 \cdot 1 + 0 \cdot 2 + 3 \cdot 6 & 1 \cdot 3 + 0 \cdot 5 + 3 \cdot 2 \\ 2 \cdot 1 + 1 \cdot 2 + 5 \cdot 6 & 2 \cdot 3 + 1 \cdot 5 + 5 \cdot 2 \end{pmatrix}$$
$$= \begin{pmatrix} 19 & 9 \\ 34 & 21 \end{pmatrix}.$$

**Theorem:** If  $A$ ,  $B$ , and  $C$  are matrices such that the given operations are well-defined, then

- $(AB)C = A(BC)$  (**associative law**)
- $A(B + C) = AB + AC$  (**left distributive law**)
- $(A + B)C = AC + BC$  (**right distributive law**)

**Proof:** We omit the proof. ■

However, matrix multiplication is **not** commutative. In fact,

(1)  $AB \neq BA$ , except in special cases:

Let

$$A = \begin{pmatrix} 1 & 1 \\ 1 & 0 \end{pmatrix} \text{ and } B = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}.$$

Then,

$$\begin{aligned} AB &= \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix} \\ BA &= \begin{pmatrix} 1 & 0 \\ 1 & 1 \end{pmatrix}. \end{aligned}$$

(2)  $AB = \mathbf{0}$  does not imply that  $A$  or  $B$  is  $\mathbf{0}$

Let

$$A = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \text{ and } B = \begin{pmatrix} 1 & 1 \\ -1 & -1 \end{pmatrix}.$$

but

$$AB = \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix}.$$

(3)  $AB = AC$  and  $A \neq 0$  do not imply that  $B = C$

Let

$$A = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}; \quad B = \begin{pmatrix} 1 & 1 \\ -1 & -1 \end{pmatrix}; \quad \text{and} \quad C = \begin{pmatrix} -1 & -1 \\ 1 & 1 \end{pmatrix}.$$

$$AB = AC = \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix}.$$

A matrix is **square** if it has an equal number of rows and columns.

$$A = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{n1} & a_{n2} & \cdots & a_{nn} \end{pmatrix}$$

The elements  $a_{11}, a_{22}, \dots, a_{nn}$  form the **principal diagonal** of the matrix  $A$ .



The **identity matrix** of order  $n$ , denoted by  $\mathbf{I}_n$ , is the  $n \times n$  matrix having ones along the main diagonal and zeros elsewhere:

$$\mathbf{I}_n = \begin{pmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \end{pmatrix} \text{ (identity matrix)}$$

If  $A$  is any  $m \times n$  matrix, then  $A\mathbf{I}_n = A = \mathbf{I}_m A$ . In particular,

$$A\mathbf{I}_n = \mathbf{I}_n A = A \text{ for every } n \times n \text{ matrix } A$$

Let  $A = (a_{ij})_{m \times n}$ . The **transpose** of  $A$  is defined as  $A^T = (a_{ji})_{n \times m}$  where the subscripts  $i$  and  $j$  are interchanged.

A square matrix is said to be **symmetric** if  $A = A^T$ .

**Theorem:** The following rules apply to matrix transposition:

1.  $(A^T)^T = A$

2.  $(A + B)^T = A^T + B^T$

3.  $(\alpha A)^T = \alpha A^T$

4.  $(AB)^T = B^T A^T$

**Proof:** We omit the proof. ■

# **Determinants (Ch. 1.1)**

Let  $A$  be an  $n \times n$  matrix.

The **determinant** of  $A$  generates a real number that summarizes some information about what matrix  $A$  is.

This is very convenient because we are usually not good at dealing with matrices but we should be OK to handle a single number.

The determinant of such  $A$  is denoted  $|A|$ . If

$$A = \begin{pmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{pmatrix},$$

we have  $|A| = a_{11}a_{22} - a_{12}a_{21}$ .

The **cofactor**  $A_{ij}$  is the determinant of  $(n-1) \times (n-1)$  matrices given by deleting  $i$ th row and  $j$ th columns from the matrix  $A$ :

$$A_{ij} = (-1)^{i+j} \begin{vmatrix} a_{11} & \cdots & a_{1,j-1} & a_{1,j+1} & \cdots & a_{1n} \\ a_{21} & \cdots & a_{2,j-1} & a_{2,j+1} & \cdots & a_{2n} \\ \vdots & \ddots & \vdots & \vdots & \ddots & \vdots \\ a_{i-1,1} & \cdots & a_{i-1,j-1} & a_{i-1,j+1} & \cdots & a_{i-1,n} \\ a_{i+1,1} & \cdots & a_{i+1,j-1} & a_{i+1,j+1} & \cdots & a_{i+1,n} \\ \vdots & \ddots & \vdots & \vdots & \ddots & \vdots \\ a_{n1} & \cdots & a_{n,j-1} & a_{n,j+1} & \cdots & a_{nn} \end{vmatrix}$$

**Definition:** Let  $A$  be an  $n \times n$  matrix. The **determinant** of  $A$  is computed by

$$|A| = a_{1j}A_{1j} + a_{2j}A_{2j} + \cdots + a_{nj}A_{nj}.$$

for each  $j = 1, \dots, n$ . This is called the **expansion of the determinant  $|A|$  with respect to the  $j$ -th column**.

Naturally, we can write a similar formula for any **row** of the determinant  $|A|$ . For example, for the  $i$ -th row,

$$|A| = a_{i1}A_{i1} + a_{i2}A_{i2} + \cdots + a_{ij}A_{ij} + \cdots + a_{in}A_{in}.$$

**The role of determinants in solving a system of linear equations:**

$$\begin{aligned}a_{11}x_1 + a_{12}x_2 &= b_1 \\a_{21}x_1 + a_{22}x_2 &= b_2\end{aligned}$$

Eliminating one of the unknowns in the usual way, one can easily obtain the formulas:

$$x_1 = \frac{b_1a_{22} - b_2a_{12}}{a_{11}a_{22} - a_{21}a_{12}} \text{ and } x_2 = \frac{a_{11}b_2 - a_{21}b_1}{a_{11}a_{22} - a_{21}a_{12}}$$

assuming that  $a_{11}a_{22} - a_{21}a_{12} \neq 0$ .

The numerators and denominators of the ratio can be represented by:

$$\begin{aligned}a_{11}a_{22} - a_{12}a_{21} &= \begin{vmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{vmatrix}, \\b_1a_{22} - b_2a_{12} &= \begin{vmatrix} b_1 & a_{12} \\ b_2 & a_{22} \end{vmatrix}, \\a_{11}b_2 - a_{21}b_1 &= \begin{vmatrix} a_{11} & b_1 \\ a_{21} & b_2 \end{vmatrix}.\end{aligned}$$

A square matrix  $A$  is **nonsingular** if  $|A| \neq 0$  and **singular** if  $|A| = 0$ .



## Some Rules for Manipulating Determinants

- (1) If two rows (or two columns) of  $A$  are interchanged, the determinant changes sign but its absolute value remains unchanged.
- (2) If all the elements in a single row (or column) of  $A$  are multiplied by a number  $c$ , the determinant is multiplied by  $c$ .
- (3) If two of the rows (or columns) of  $A$  are proportional, then  $|A| = 0$ .
- (4) The value of  $|A|$  remains unchanged if a multiple of one row (or one column) is added to another row (or column).

## Cramer's Rule (Ch. 1.1)

A linear system of  $n$  equations and  $n$  unknowns is given:

$$\begin{array}{ccccccccc} a_{11}x_1 + a_{12}x_2 + \cdots + a_{1n}x_n & = & b_1 \\ a_{21}x_1 + a_{22}x_2 + \cdots + a_{2n}x_n & = & b_2 & (*) \\ \vdots & & \vdots & & \vdots & & \vdots & & \vdots \\ a_{n1}x_1 + a_{n2}x_2 + \cdots + a_{nn}x_n & = & b_n \end{array}$$

**Theorem (Cramer's Rule):**  $(*)$  has a unique solution if and only if  $|A| \neq 0$ . The solution is then

$$x_j = \frac{|A_j|}{|A|}, \quad j = 1, \dots, n$$

where the determinant  $|A_j|$  is defined as:

$$|A_j| = \begin{vmatrix} a_{11} & \cdots & a_{1,j-1} & \boxed{b_1} & a_{1,j+1} & \cdots & a_{1n} \\ a_{21} & \cdots & a_{2,j-1} & \boxed{b_2} & a_{2,j+1} & \cdots & a_{2n} \\ \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots \\ a_{n1} & \cdots & a_{n,j-1} & \boxed{b_n} & a_{n,j+1} & \cdots & a_{nn} \end{vmatrix}.$$

Note that  $|A_j|$  is obtained by replacing the  $j$ th column of  $|A|$  by the column whose components are  $b_1, b_2, \dots, b_n$ .

**Proof:** We omit the proof. ■

A linear system of  $n$  equations and  $n$  unknowns  $Ax = b$  is given:

$$\begin{array}{ccccccccc} a_{11}x_1 + a_{12}x_2 + \cdots + a_{1n}x_n & = & b_1 \\ a_{21}x_1 + a_{22}x_2 + \cdots + a_{2n}x_n & = & b_2 & (*) \\ \vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\ a_{n1}x_1 + a_{n2}x_2 + \cdots + a_{nn}x_n & = & b_n \end{array}$$

If  $b_1 = \cdots = b_n = 0$ , i.e.,  $Ax = 0$ , the system  $(*)$  is called **homogeneous**.

A homogeneous system always has the **trivial** solution:  $x_1 = \cdots = x_n = 0$ .

## Matrix Inverse (Ch. 1.1)

Let  $A$  be an  $n \times n$  matrix.

When can we find an  $n \times n$  matrix  $B$  such that  $BA = AB = I$ ?

If such  $B$  exists, it is called the **inverse** of matrix  $A$  and is denoted by  $A^{-1}$ .

If  $A^{-1}$  exists, matrix  $A$  itself is said to be **invertible (nonsingular)**.

**Theorem:** Every nonsingular matrix  $A$  has a unique inverse matrix  $B$  such that

$$AB = BA = I.$$

If  $A = (a_{ij})_{n \times n}$  and  $|A| \neq 0$ , the unique inverse of  $A$  is given by

$$A^{-1} = \frac{1}{|A|} \text{adj}(A), \text{ where } \text{adj}(A) = \begin{pmatrix} A_{11} & A_{21} & \cdots & A_{n1} \\ A_{12} & A_{22} & \cdots & A_{n2} \\ \vdots & \vdots & \ddots & \vdots \\ A_{1n} & A_{2n} & \cdots & A_{nn} \end{pmatrix}$$

with  $A_{ij}$ , the **cofactor** of the element  $a_{ij}$ . Note carefully the order of the indices in  $\text{adj}(A)$  with the column number preceding the row number.

**Proof:** We omit the proof. ■

We provide the results on matrix inverse without proofs.

**Lemma (Rules for Matrix Inverse):**

- $(A^{-1})^{-1} = A,$
- $(AB)^{-1} = B^{-1}A^{-1},$
- $(A^T)^{-1} = (A^{-1})^T,$
- $(\alpha A)^{-1} = \alpha^{-1}A^{-1},$  where  $\alpha \in \mathbb{R}.$

**Proposition:** If  $|A| \neq 0,$  then  $|A^{-1}| = 1/|A|.$

# **Vectors (Ch. 1.1)**



An  $n$ -**vector** is an ordered  $n$ -tuple of (real) numbers.

An  $n$ -vector can be understood either as a  $1 \times n$  matrix  $\mathbf{a} = (a_1, a_2, \dots, a_n)$  (a **row vector**) or as an  $n \times 1$  matrix (a **column vector**)

$$\mathbf{a} = \begin{pmatrix} a_1 \\ a_2 \\ \vdots \\ a_n \end{pmatrix}.$$

The operations of addition, subtraction and “scalar” multiplication of vectors are defined in the obvious way.

## Inner Product (Dot Product)

We introduce the multiplication of two vectors:

The **dot product** (or **inner product**) of the  $n$ -vectors  $\mathbf{a}$  and  $\mathbf{b}$  is defined as

$$\mathbf{a} \cdot \mathbf{b} = a_1b_1 + a_2b_2 + \cdots + a_nb_n = \sum_{i=1}^n a_ib_i.$$

**Proposition (Properties of the Inner Product):** If  $\mathbf{a}$ ,  $\mathbf{b}$ , and  $\mathbf{c}$  are  $n$ -vectors and  $\alpha$  is a scalar, then

1.  $\mathbf{a} \cdot \mathbf{b} = \mathbf{b} \cdot \mathbf{a}$ ,

2.  $\mathbf{a} \cdot (\mathbf{b} + \mathbf{c}) = \mathbf{a} \cdot \mathbf{b} + \mathbf{a} \cdot \mathbf{c}$ ,

3.  $(\alpha \mathbf{a}) \cdot \mathbf{b} = \mathbf{a} \cdot (\alpha \mathbf{b}) = \alpha(\mathbf{a} \cdot \mathbf{b})$ .

4.  $\mathbf{a} \cdot \mathbf{a} = 0 \Rightarrow \mathbf{a} = \mathbf{0}$

5.  $(\mathbf{a} + \mathbf{b}) \cdot (\mathbf{a} + \mathbf{b}) = \mathbf{a} \cdot \mathbf{a} + 2(\mathbf{a} \cdot \mathbf{b}) + \mathbf{b} \cdot \mathbf{b}$ .

**Proof:** We skip the proof. ■

The **Euclidean norm** or **length** of the vector  $\mathbf{a}$  is defined:

$$\|\mathbf{a}\| = \sqrt{\mathbf{a} \cdot \mathbf{a}} = \sqrt{a_1^2 + a_2^2 + \cdots + a_n^2}$$

NOTE:  $\|\alpha \mathbf{a}\| = |\alpha| \|\mathbf{a}\|$  for all scalars and vectors.

**Lemma:** The following useful inequalities hold.

1.  $|\mathbf{a} \cdot \mathbf{b}| \leq \|\mathbf{a}\| \cdot \|\mathbf{b}\|$  (**Cauchy-Schwartz inequality**)
2.  $\|\mathbf{a} + \mathbf{b}\| \leq \|\mathbf{a}\| + \|\mathbf{b}\|$  (**Minkowski inequality**)

**Proof:** We omit the proof. ■

**Remark:** You are referred to Ch.B.1 for trigonometric functions.

Cauchy-Schwartz inequality implies that, for any  $\mathbf{a}, \mathbf{b} \in \mathbb{R}^n$ ,

$$-1 \leq \frac{\mathbf{a} \cdot \mathbf{b}}{\|\mathbf{a}\| \|\mathbf{b}\|} \leq 1.$$

Thus, the **angle**  $\theta$  between nonzero vectors  $\mathbf{a}$  and  $\mathbf{b}$  is defined by

$$\cos \theta = \frac{\mathbf{a} \cdot \mathbf{b}}{\|\mathbf{a}\| \|\mathbf{b}\|}, \quad \theta \in [0, \pi]$$

This definition reveals that  $\cos \theta = 0$  if and only if  $\mathbf{a} \cdot \mathbf{b} = 0$ . Then  $\theta = \pi/2 = 90^\circ$ . In symbols,

$$\mathbf{a} \perp \mathbf{b} \iff \mathbf{a} \cdot \mathbf{b} = 0$$

# Linear Independence (Ch. 1.2)

**Definition:** The  $n$  vectors  $\mathbf{a}_1, \mathbf{a}_2, \dots, \mathbf{a}_n$  of an Euclidean space are **linearly dependent** if there exist numbers  $c_1, c_2, \dots, c_n \in \mathbb{R}$ , not all zero, such that

$$c_1\mathbf{a}_1 + c_2\mathbf{a}_2 + \dots + c_n\mathbf{a}_n = \mathbf{0}$$

If this equation holds only when  $c_1 = c_2 = \dots = c_n = 0$ , then the vectors are **linearly independent**.

The next result is a characterization of linear dependence in a linear space.

**Proposition:** A set of  $n$  vectors  $\mathbf{a}_1, \mathbf{a}_2, \dots, \mathbf{a}_n$  of an Euclidean space is linearly dependent if and only if at least one of them can be written as a linear combination of the others.

**Proof:** Suppose that  $\mathbf{a}_1, \mathbf{a}_2, \dots, \mathbf{a}_n$  are linearly dependent. Then the equation  $c_1\mathbf{a}_1 + \dots + c_n\mathbf{a}_n = 0$  holds with at least **one** of the coefficients  $c_i$  differently from 0. We can, without loss of generality, assume that  $c_1 \neq 0$ . Solving the equation for  $\mathbf{a}_1$  yields

$$\mathbf{a}_1 = -\frac{c_2}{c_1}\mathbf{a}_2 - \dots - \frac{c_n}{c_1}\mathbf{a}_n.$$

Thus,  $\mathbf{a}_1$  is a linear combination of the other vectors. ■



The relation between the determinant and linear dependence:

**Theorem:** Let  $A$  be an  $n \times n$  matrix. Then,  $|A| = 0$  if and only if there is a linear dependence between its columns of  $A$ .

**Proof:** We omit the proof. ■

## The Rank of a Matrix (Ch. 1.3)

**Definition:** The **rank** of a matrix  $A$ , written  $\text{rank}(A)$ , is the maximum number of linearly independent column vectors in  $A$ . If  $A$  is the  $0$  matrix, we put  $\text{rank}(A) = 0$ .

$A = (\mathbf{a}_1, \dots, \mathbf{a}_n)$ :  $n \times n$  matrix

$\text{rank}(A) = n \Rightarrow |A| \neq 0$ ;

$\text{rank}(A) < n \Rightarrow |A| = 0$ .

**Theorem:** Let  $A = (\mathbf{a}_1, \dots, \mathbf{a}_n)$  be an  $n \times n$  matrix. Then,  $|A| \neq 0$  if and only if  $\mathbf{a}_1, \dots, \mathbf{a}_n$  are linearly independent.

**Proof:** We skip the proof. ■

## Example

$$A = \begin{pmatrix} 1 & 0 & 2 & 1 \\ 0 & 2 & 4 & 2 \\ 0 & 2 & 2 & 1 \end{pmatrix}$$

We can easily see that the third and fourth columns are parallel to each other:

$$\begin{pmatrix} 1 \\ 2 \\ 1 \end{pmatrix} = \frac{1}{2} \begin{pmatrix} 2 \\ 4 \\ 2 \end{pmatrix}.$$

So, we now consider the following truncated matrix:

$$B = \begin{pmatrix} 1 & 0 & 2 \\ 0 & 2 & 4 \\ 0 & 2 & 2 \end{pmatrix}$$

We compute the determinant of the truncated matrix:

$$\begin{vmatrix} 1 & 0 & 2 \\ 0 & 2 & 4 \\ 0 & 2 & 2 \end{vmatrix} = 1 \cdot (-1)^{1+1} \begin{vmatrix} 2 & 4 \\ 2 & 2 \end{vmatrix} \\ = 4 - 8 = -4 \neq 0.$$

Since  $|B| \neq 0$ , all the columns vectors in  $B$  are linearly independent.

Thus,  $\text{rank}(A) = 3$ .