Introductory Econometrics

Basic Statistics

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We begin by recalling some elementary notions from linear algebra needed for introducing random vectors.

A vector \mathbf{x} of dimension n is an ordered collection of n numbers, which are called **components** or elements:

$$\mathbf{x} = (x_1, \ldots, x_n)$$

Example

$$\mathbf{x} = (2,3), \quad \mathbf{y} = (-1,2), \quad \mathbf{z} = (\sqrt{2},0,\pi)$$

Two vectors \mathbf{x} and \mathbf{y} of common dimension n are *added* component by component:

$$\mathbf{x} + \mathbf{y} = (x_1 + y_1, \dots, x_n + y_n)$$

Example

If $\mathbf{x} = (2,3)$ and $\mathbf{y} = (-1,2)$, then

$$\mathbf{x} + \mathbf{y} = (2 + (-1), 3 + 2) = (1, 5),$$

Two vectors \mathbf{x} and \mathbf{y} of common dimension n are *added* component by component:

$$\mathbf{x} + \mathbf{y} = (x_1 + y_1, \dots, x_n + y_n)$$

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Spring, 2025 2 / 54

If a is a number, we will sometimes refer to a as a scalar.

If a is a scalar and **x** is a vector, then the product a**x** of a and **x** is

$$a\mathbf{x} = (ax_1, \ldots, ax_n)$$

If a = 0 we get $a\mathbf{x} = \mathbf{0}$ where $\mathbf{0}$ is the **zero vector**, i.e. the vector of same dimension as \mathbf{x} with all components equal to zero.

A **matrix** can be viewed in different ways: either as an array of numbers ordered into rows and columns, or as a collection of vectors.

Example

$$\mathbf{A} = \left(\begin{array}{rrr} 2 & 2 & -1 \\ 3 & 1 & 0 \end{array}\right)$$

The vectors (2,3), (2,1) and (-1,0) are the columns of A.

Linear Algebra

A matrix has m rows and n columns

$$\mathbf{A} = \left(\begin{array}{ccccc} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{array}\right)$$

To indicate the dimension of a matrix we sometimes write $\mathbf{A}_{m \times n}$. Example

$$\mathbf{B} = \left(\begin{array}{cc} 2 & 2\\ 3 & 1 \end{array}\right)$$

B is a 2×2 matrix with, for example, element $b_{12} = 2$. **B** is a so-called **square matrix**, i.e. a matrix with the same number of rows as columns (m = n).

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Spring, 2025 5 / 54

Matrix addition of two matrices of common dimension is analogous to vector addition, and it can be shown that the previously mentioned properties of vectors also hold for matrices.

Example

$$\begin{pmatrix} 1 & 3 \\ 7 & -1 \end{pmatrix} + \begin{pmatrix} 2 & 5 \\ 8 & 1 \end{pmatrix} = \begin{pmatrix} 3 & 8 \\ 15 & 0 \end{pmatrix}$$
$$\begin{pmatrix} 2 & 3 \\ 1 & 5 \\ 6 & 4 \end{pmatrix} - \begin{pmatrix} 1 & 9 \\ 7 & 4 \\ 2 & 5 \end{pmatrix} = \begin{pmatrix} 1 & -6 \\ -6 & 1 \\ 4 & -1 \end{pmatrix}$$

Linear Algebra

We can also define **matrix multiplication**:

Let **A** be a $m \times l$ matrix, with row *i* and column *k* element a_{ik} , and let **B** be a $l \times n$ matrix, with row *k* and column *j* element b_{kj} . Then c_{ij} , the row *i* and column *j* element of the matrix **C** = **AB**, is given by

$$c_{ij} = \sum_{k=1}^{l} a_{ik} b_{kj}$$

Example

$$\begin{pmatrix} 2 & 1 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} 0 & 4 \\ 2 & 1 \end{pmatrix} = \begin{pmatrix} 2 \times 0 + 1 \times 2 & 2 \times 4 + 1 \times 1 \\ 1 \times 0 + 0 \times 2 & 1 \times 4 + 0 \times 1 \end{pmatrix} = \begin{pmatrix} 2 & 9 \\ 0 & 4 \end{pmatrix}$$

Remark

From the definition it follows that the number of columns of \mathbf{A} must equal the number of rows of \mathbf{B} in order for the product $\mathbf{C} = \mathbf{AB}$ to be well-defined:

$$\mathbf{A}_{m \times l} \mathbf{B}_{l \times n} = \mathbf{C}_{m \times n}$$

Moreover, in general, **AB** is different from **BA**.

Some elementary properties of matrices:

Let a be a scalar, and let **A**, **B** and **C** be matrices of dimensions such that the left-hand side expressions below are well-defined. Then,

A(B + C) = AB + ACA(aB) = a(AB)(AB)C = A(BC)

The **identity matrix**, **I**, of dimension n is the $n \times n$ matrix with elements along its main diagonal equal to 1, and all other elements equal to 0:

$$\mathbf{I}_{n} = \left(\begin{array}{ccccc} 1 & 0 & \cdots & 0 \\ 0 & 1 & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ 0 & \cdots & 0 & 1 \end{array} \right)$$

It is easy to verify that

$$\mathbf{I}_m \mathbf{A} = \mathbf{A} = \mathbf{A} \mathbf{I}_n$$

for all $m \times n$ matrices **A**. Yaohan Chen (AHU) Let **A** be a square matrix. If there exists a matrix \mathbf{A}^{-1} such that $\mathbf{A}^{-1}\mathbf{A} = \mathbf{I}$, where **I** is the identity matrix of same dimension as **A**, then \mathbf{A}^{-1} is called the inverse of **A**. Not all matrices have an inverse.

Example

If
$$A_{1\times 1} = a_{11}$$
 and $a_{11} \neq 0$, then $\mathbf{A}^{-1} = \frac{1}{a_{11}}$
If $\mathbf{B} = \begin{pmatrix} 2 & 1\\ 1 & 0 \end{pmatrix}$, then $\mathbf{B}^{-1} = \begin{pmatrix} 0 & 1\\ 1 & -2 \end{pmatrix}$

The **transpose** \mathbf{A}^{T} of a matrix $\mathbf{A}_{m \times n}$ is the $n \times m$ matrix whose *i*th column is the *i*th row of \mathbf{A} :

$$\mathbf{A} = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{pmatrix}, \quad \mathbf{A}^{\mathsf{T}} = \begin{pmatrix} a_{11} & a_{21} & \cdots & a_{m1} \\ a_{12} & a_{22} & \cdots & a_{m2} \\ \vdots & \vdots & \ddots & \vdots \\ a_{1n} & a_{2n} & \cdots & a_{mn} \end{pmatrix}$$

A square matrix A with A[↑] = A is called a symmetric matrix.

The **determinant** of a 2×2 matrix

$$\mathbf{A} = \left(\begin{array}{cc} a_{11} & a_{12} \\ a_{21} & a_{22} \end{array}\right)$$

is given by the scalar

$$|\mathbf{A}| = a_{11}a_{22} - a_{12}a_{21}$$

The determinant of a general $n \times n$ matrix is more complicated but can be defined recursively.

One case is particularly simple:

The determinant of a diagonal matrix is the product of the elements along its main diagonal:

 $|\mathbf{A}| = a_{11} \times \cdots \times a_{nn}$

Example

Since $|\mathbf{I}_n| = 1 \times \cdots \times 1$, the determinant of the identity matrix is one.

We begin by recalling some elementary notions from set theory needed for introducing concepts from probability theory.

- A set is a collection of objects. Sets are usually denoted by upper-case letters such as A, B or C.
- If an object c belongs to a set C we write $c \in C$ (read, "c in C").
- If c does not belong to C we write c ∉ C (read, "c not in C").
- If c_1, \ldots, c_n are objects, the set consisting of precisely these *n* objects is denoted by $\{c_1, \ldots, c_n\}$.

- The objects belonging to a set A are called its **elements** (or members). The elements can be **anything**, e.g. numbers, outcomes or other sets.
- Let A and B be sets. The **intersection** of A and B is the set whose elements are those objects c such that $c \in A$ and $c \in B$. We write $A \cap B$ (read, "A intersect B").
- The **union** of A and B is the set whose elements are those objects c such that c belongs to **at least** one of the two sets A,B (i.e. either $c \in A$ or $c \in B$, or both). We write $A \cup B$ (read, A union B).

When flipping a coin, 'head' or 'tail' occurs. If $A = \{\text{head}, \text{tail}\}$ and $B = \{\text{head}\}$, then

 $A \cap B = \{\text{head}\} \text{ and } A \cup B = \{\text{head}, \text{tail}\}$



Figure : The union of A and B illustrated using a Venn diagram



Figure : The intersection of A and B

- If all elements of a set A are also elements of a set B, we say that A is a subset of B, and write A ⊂ B.
- **Empty** set is the set that has no members. The empty set is denoted by \emptyset .
- The set \emptyset is a subset of **any** set.

The **difference** of two sets A and B, written in A-B, is the set of all elements that are in A but not in B.

Example

If $A = \{STI, HSI, SSE\}$ and $B = \{HSI\}$, then $A-B = \{STI, SSE\}$.

If A and B have no elements in common we say that the two sets are **disjoint**, $A \cap B = \emptyset$.

Let $A \subset C$. The **complement** of A in C is the set of elements that belong to C but not to A. We write A^c .

Example The possible outcomes of a coin tossing experiment are 'head' and 'tail'. Here $C = \{head, tail\}$. Hence, if $A = \{head\}$, then $A^c = \{tail\}$. Yaohan Chen (AHU) Spring, 2025 21 / 54

Function

A **function** is a rule that associates each member of one set with a member of another set.



Figure : A function f takes an input x and returns an output f(x).

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Spring, 2025 22 / 54

- Assigning a value to each random outcome.
- When tossing a coin we can write '1' for 'head' and '0' for tail. In this way, we get a **random variable** $X(\omega) \mapsto \{0,1\}$, where ω belongs to the **sample space** $\mathcal{F} = \{\text{head, tail}\}$.
- \mathcal{F} is a abstract space collecting **all possible outcomes** of the underlying experiment.
- The random variable $X(\omega)$ is nothing but a **real-valued** function defined on \mathcal{F} (i.e. a numerical summary of a random outcome).

- How to make \mathcal{F} complete ? \mathcal{F} should be like:
 - If $A \in \mathcal{F}$, so is its complement A^c ;
 - If $A, B \in \mathcal{F}$, so are $A \cap B$, $A \cup B$, $A \cup B^c$, $B \cup A^c$, $A \cap B^c$, $B \cap A^c$, etc.
- In some advanced textbooks \mathcal{F} is called σ -filed.

Probability

- Probability is a measure such that for each $A \in \mathcal{F}$, it assigns a number $P(A) \in [0, 1]$.
- Probability should satisfy:
 - For $A, B \in \mathcal{F}$

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

- If A and B are disjoint,

$$\mathbf{P}(A \cup B) = \mathbf{P}(A) + \mathbf{P}(B)$$

- Moreover,

$$P(A^c) = 1 - P(A), \quad P(\mathcal{F}) = 1 \quad \text{and} \quad P(\emptyset) = 0$$

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Spring, 2025 25 / 54

Distribution

• The collection of the probabilities

$$F_X(x) = P(X \le x) = P(\{\omega : X(\omega) \le x\}), \quad x \in \mathbf{R}$$

is the **Cumulative Distribution Function (CDF)** $F_X(x)$ of X. $F_X(x)$ of X gives the probability that X belongs to the interval (a, b] as

$$P(\{\omega : a < X(\omega) \le b\}) = F_X(b) - F_X(a), \quad a < b$$

Continuous distributions have Probability Density Function (PDF) f_X(x):

$$F_X(x) = \int_{-\infty}^x f_X(t)dt, \quad x \in \mathbf{R}, \quad \text{and } \int_{-\infty}^\infty f_X(t)dt = 1$$

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Spring, 2025 26 / 54

In the following lectures we will frequently make use of certain finite-dimensional random structures.

We consider finite-dimensional random vectors:

 $\mathbf{X} = (X_1, \dots, X_n)$ is a *n*-dimensional **random vector** if its components X_1, \dots, X_n are one-dimensional real-valued random variables.

Toss a coin. We consider the pairs

(head,head), (tail,tail), (head,tail), (tail,head)

as outcomes of the experiment. These four pairs form the sample space \mathcal{F} . We can write '1' for 'head' and '0' for 'tail'. In this way, we get two random variables X_1 and X_2 , and $\mathbf{X} = (X_1, X_2)$ is a two-dimensional random vector.

Example (Cont'd)

Note that

$$X$$
 (head, head) = (1, 1), X (tail, tail) = (0, 0)

$$X$$
 (head, tail) = (1,0), X (tail, head) = (0,1)

If the coin is 'fair', we can assign the probability 0.25 to each of the four outcomes, i.e.

$$P(\{\omega : \mathbf{X}(\omega) = (k, i)\}) = 0.25, \quad k, i \in \{0, 1\}$$

• The collection of the probabilities

$$F_{\mathbf{X}}(\mathbf{x}) = P(X_1 \le x_1, \dots, X_n \le x_n)$$

= P({\{\omega : X_1(\omega) \le x_1, \dots, X_n(\omega) \le x_n}\})

where $\mathbf{x} = (x_1, \ldots, x_n) \in \mathbf{R}^n$, is the **joint CDF** $F_{\mathbf{X}}$ of \mathbf{X} .

- $F_{\mathbf{X}}(\mathbf{x})$ is the shorthand for $F_{X_1,\ldots,X_n}(x_1,\ldots,x_n)$.
- Correspondingly, joint **PDF** of $F_{\mathbf{X}}$ is

$$F_{\mathbf{X}}(\mathbf{x}) = F_{\mathbf{X}}(x_1, \dots, x_n) = \int_{-\infty}^{x_1} \cdots \int_{-\infty}^{x_n} f_{\mathbf{X}}(w_1, \dots, w_n) dw_1 \cdots dw_n$$

Random Variable:

Expectation, Variance, Moment

• Expectation

$$\mu_X = \mathbb{E}(X) = \int_{-\infty}^{\infty} x f_X(x) dx$$

• Variance

$$\sigma_X^2 = \operatorname{Var}(X) = \int_{-\infty}^{\infty} (x - \mu_X)^2 f_X(x) dx$$

• *I*-th moment

$$\mathbb{E}(X^{I}) = \int_{-\infty}^{\infty} x^{I} f_{X}(x) dx$$

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Spring, 2025 31 / 54

• **Correlation** between two random variables X_1 and X_2 is defined as

$$\operatorname{Corr}(X_1, X_2) = \frac{\operatorname{Cov}(X_1, X_2)}{\sigma_{X_1} \sigma_{X_2}} = \frac{\mathbb{E}\left[(X_1 - \mu_{X_1}) (X_2 - \mu_{X_2}) \right]}{\sigma_{X_1} \sigma_{X_2}}$$

• $-1 \leq \operatorname{Corr}(X_1, X_2) \leq 1$. Why ?

Independence

Two **events** A_1 and A_2 are **independent** if

$$P(A_1 \cap A_2) = P(A_1) P(A_2)$$

Two random variables X_1 and X_2 are independent if

$$P(X_1 \in B_1, X_2 \in B_2) = P(X_1 \in B_1) P(X_2 \in B_2)$$

for all suitable subsets of B_1 and B_2 of **R**. This means that the events

$$\{c: X_1(\omega) \in B_1\}$$
 and $\{\omega: X_2(\omega) \in B_2\}$

are independent.

• The random variables X_1, \ldots, X_n are mutually independent if and only if their joint CDF can be written as

$$F_{X_1,...,X_n}(x_1,...,x_n) = F_{X_1}(x_1) \times \cdots \times F_{X_n}(x_n), (x_1,...,x_n) \in \mathbf{R}^n$$

• If the random vector $\mathbf{X} = (X_1, \dots, X_n)$ has joint PDF $f_X = f_{X_1,\dots,X_n}$ with marginal pdfs f_{X_1},\dots,f_{X_n} , then X_1,\dots,X_n are mutually independent if and only if

$$f_{X_{1},...,X_{n}}(x_{1},...,x_{n}) = f_{X_{1}}(x_{1}) \times \cdots \times f_{X_{n}}(x_{n}), (x_{1},...,x_{n}) \in \mathbf{R}^{n}$$

- Uncorrelated random variables ⇒ Independent Random Variables ? Any counterexample ?
- If the random variables X_1, \ldots, X_n are mutually independent and have the *same distribution*, we say that they are **independent and identically distributed (iid)**.

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Spring, 2025 35 / 54

Random Vectors:

Expectation, Variance-Covariance Matrix

• Expectation

$$\mu_{\mathbf{X}} = \mathbb{E}(\mathbf{X}) = (\mathbb{E}(X_1), \dots, \mathbb{E}(X_n))$$

• The Variance-covariance matrix of **X** is defined as the matrix $\Sigma_{\mathbf{X}}$ with row *i* column *j* element given by

$$\operatorname{Cov}(X_i, X_j), \quad i, j = 1, \dots, n$$

where

$$\operatorname{Cov} (X_i, X_j) = \mathbb{E} \left[(X_i - \mu_{X_i}) \left(X_j - \mu_{X_j} \right) \right]$$
$$= \mathbb{E} (X_i X_j) - \mu_{X_i} \mu_{X_j}$$

and $\operatorname{Cov}(X_i, X_i) = \sigma_{X_i}^2$.

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Spring, 2025 36 / 54

The most important continuous distribution is the normal or Gaussian distribution:

A random variable X is said to be **normally distributed** or $\mathcal{N}(\mu, \sigma^2)$ with parameters $-\infty < \mu < \infty$ if

$$f_X(x) = \frac{1}{\sqrt{2\pi\sigma}} e^{-\frac{1}{2}\left(\frac{x-\mu}{\sigma}\right)^2}, \quad -\infty < x < \infty$$

and when $\mu = 0$ and $\sigma = 1$, X is called **standard normal**.

The CDF of standard normal distribution has its own notation $\Phi(x)$.

The Multivariate Normal Distribution (*n*dimensional normal) or Gaussian Distribution is given by its joint PDF

$$f_{\mathbf{X}}(\mathbf{x}) = \frac{1}{(2\pi)^{n/2} |\mathbf{\Sigma}|^{1/2}} \exp\left\{-\frac{1}{2}(\mathbf{x}-\boldsymbol{\mu})\mathbf{\Sigma}^{-1}(\mathbf{x}-\boldsymbol{\mu})^{\mathsf{T}}\right\}, \quad \mathbf{x} \in \mathbf{R}^{n}$$

with parameters $\boldsymbol{\mu} \in \mathbf{R}^n$ and $\boldsymbol{\Sigma}$ is symmetric (positive definite) $n \times n$ matrix, $\boldsymbol{\Sigma}^{-1}$ its inverse and $|\boldsymbol{\Sigma}|$ its determinant. Multivariate Normal Distribution is denoted by $\mathcal{N}(\boldsymbol{\mu}, \boldsymbol{\Sigma})$

Suppose $\mathbf{X} = (X_1, X_2)$ is 2-dimensional normal with

$$\boldsymbol{\mu} = (0,0)$$
 and $\boldsymbol{\Sigma} = \mathbf{I}_2 = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}$

then

$$f_{\mathbf{X}}(\mathbf{x}) = f_{\mathbf{X}}(x_1, x_2) = \frac{1}{2\pi} e^{-\frac{1}{2}(x_1^2 + x_2^2)}$$

for $\mathbf{x} \in \mathbf{R}^2$.

If $\mu = 0$ and $\Sigma = I_n$, then the density f_X is simply the product of *n* standard normal densities:

$$f_{\mathrm{X}}(x_1,\ldots,x_n) = \varphi(x_1) \times \cdots \times \varphi(x_n), \quad \varphi(x) = \frac{1}{\sqrt{2\pi}} e^{-\frac{x^2}{2}}$$

Let $\mathbf{X} = (X_1, \dots, X_n)$ have an $\mathcal{N}(\boldsymbol{\mu}, \boldsymbol{\Sigma})$ distribution and \mathbf{A} be an $m \times n$ matrix. Then $\mathbf{A}\mathbf{X}^{\mathsf{T}}$ has an $\mathcal{N}(\mathbf{A}\boldsymbol{\mu}^{\mathsf{T}}, \mathbf{A}\boldsymbol{\Sigma}\mathbf{A}^{\mathsf{T}})$ distribution.

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Let $\mathbf{X} = (X_1, \dots, X_n)$ have an $\mathcal{N}(\boldsymbol{\mu}, \boldsymbol{\Sigma})$ distribution and \mathbf{A} be an $m \times n$ matrix. Then $\mathbf{A}\mathbf{X}^{\mathsf{T}}$ has an $\mathcal{N}(\mathbf{A}\boldsymbol{\mu}^{\mathsf{T}}, \mathbf{A}\boldsymbol{\Sigma}\mathbf{A}^{\mathsf{T}})$ distribution.

In the Gaussian case, uncorrelatedness and independence are equivalent notions. This statement is wrong for non-Gaussian vectors.

• If Z_1, \ldots, Z_k are **iid standard normal** random variables, then

$$\sum_{i=1}^{k} Z_i^2 = Z_1^2 + Z_2^2 + \dots + Z_k^2 \sim \chi^2_{(k)}$$

- A random variable T follows the student t distribution with k degrees of freedom, written as $T \sim t(k)$ if $T = \frac{U}{\sqrt{V/k}}$, where $U \sim \mathcal{N}(0,1), V \sim \chi^2_{(k)}$, and U and V are independent.
- A random variable F follows the F-distribution with (m, n) degrees of freedom, written as F ~ F(m, n), if F = U/m/V/n, where U ~ χ²_(m) and V ~ χ²_(n), and U and V are independent.

- In a typical real-world statistical problem, we have a random variable X of interest, but the PDF f(x) is not known.
- Our lack of knowledge can be classified in on e of two ways:
 - f(x) is completely unknown.
 - The functional form of f(x) is assumed to be known up to a parameter vector θ .

X has a normal distribution $\mathcal{N}(\mu, \sigma^2)$, where $\theta = (\mu, \sigma^2)$.

We often write $f(x; \theta)$, where $\theta \in \Omega$ for a specified set Ω , to emphasize that the PDF is known up to θ .

Example

If X has a normal distribution $\mathcal{N}(\mu, \sigma^2)$, then

$$\Omega = \{\theta = (\mu, \sigma^2) : \mu \in \mathbf{R}, \sigma^2 > 0\}$$

- We call θ a **parameter** of the distribution.
- As θ is unknown, we want to **estimate** it.

In practice, our information about the unknown distribution of X, or the unknown parameters of the distribution of X, comes from a **sample** of X.

The sample observations have the same distribution as X, and we denote them as the random variables X_1, \ldots, X_n .

- n denotes the sample size.
- When the sample is actually drawn, we use lower case letters x_1, \ldots, x_n to denote the values or **realizations** of the sample

We often use functions of a sample to summarize the information in it:

Let X_1, \ldots, X_n denote a sample of a random variable X, and let $T = T(X_1, \ldots, X_n)$ be a function of the sample, then T is called a statistic.

- T is a random variable, and has nothing to do with θ .
- Once is the sample is drawn, t is called a **realization** of T, where $t = T(x_1, \ldots, x_n)$ and x_1, \ldots, x_n is the realization of the sample.

- Let X_1, \ldots, X_n denote a sample of random variable X with a pdf $f(x; \theta)$, where $\theta \in \Omega$ for a specified set Ω .
- Statistic T is called a **point estimator** of θ , usually denoted by $\hat{\theta}$.
- The realization t of T is called an **estimate** of θ .

Point estimator might have various properties: unbiasedness, consistency, and efficiency.

We temporarily focus on **unbiasedness**, but leave the discussion about consistency and efficiency later.

Let X_1, \ldots, X_n denote a sample of a random variable X with pdf $f(x; \theta), \theta \in \Omega$, and let $T = T(X_1, \ldots, X_n)$ be a statistic. Then T is an **unbiased** estimator of θ if $\mathbb{E}(T) = \theta$.

- The mean of an unbiased estimators sampling distribution is located at the true (but unknown) value of the parameter of interest.
- An estimator T of θ with $\mathbb{E}(T) \neq \theta$ is **biased**.

Consider the dataset where the variable of interest X is the number of operating hours until the first failure of airconditioning units for Boeing 720 airplanes.

A sample with size n = 13 was obtained, with realized values:

359, 413, 25, 130, 90, 50, 50, 487, 102, 194, 55, 74, 97

- $\hat{\theta} = \frac{1}{n} \sum_{i=1}^{n} X_i$ is a point **estimator** of θ .
- $\frac{1}{n} \sum_{i=1}^{n} x_i = \frac{1}{13} (359 + 413 + 25 + 130 + 90 + 50 + 50 + 487 + 102 + 194 + 55 + 74 + 97) = 163.5385$ is the corresponding **estimate** of θ .

We use confidence intervals to measure the error of the corresponding estimate:

Let X_1, \ldots, X_n be a sample of a random variable X, where X has pdf $f(x; \theta)$, $\theta \in \Omega$. Let $0 < \alpha < 1$ be specified. Let $L = L(X_1, \ldots, X_n)$ and $U = U(X_1, \ldots, X_n)$ be two statistics. The interval (L, U) is a $(1 - \alpha)100\%$ confidence interval for θ if

 $1 - \alpha = \mathcal{P}_{\theta}[\theta \in (L, U)]$

where P_{θ} refers to the probability when θ is the true parameter. That is, the probability that the random interval includes θ is $1 - \alpha$, which is called the **confidence level** of the interval.

Suppose, for simplicity, that $\hat{\theta}$ denotes the estimator of θ for $\mathcal{N}(\theta, \sigma^2)$ with σ^2 known. Then,

$$P_{\theta}\left(-z_{\alpha/2} < \frac{\hat{\theta} - \theta}{\sigma} < z_{\alpha/2}\right) = P_{\theta}\left(-z_{\alpha/2} < Z < z_{\alpha/2}\right) = 1 - \alpha$$

where we have used that Z is a standard normal random variable.

A $(1-\alpha)100\%$ confidence interval is now obtained by solving for θ in the above equation. This gives

$$\mathrm{P}_{\theta}\left(\hat{\theta}-z_{\alpha/2}\sigma<\theta<\hat{\theta}+z_{\alpha/2}\sigma\right)=1-\alpha$$

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Spring, 2025 50 / 54

Point estimation and confidence intervals are useful statistical inference procedures. Another type of inference that is frequently used concerns tests of hypotheses.

We label these hypotheses as

 $H_0: \theta \in \Omega_0$ against $H_1: \theta \in \Omega_1$

where $\Omega_0 \cap \Omega_1 = \emptyset$ and $\Omega_1 \subset \Omega$, $\Omega_0 \subset \Omega$

 H_0 is called the **null hypothesis**, and H_1 the **alternative** hypothesis.

The decision take H_0 or H_1 is based on a sample X_1, \ldots, X_n from the distribution of X:

Denote the range of the random sample $X = (X_1, \ldots, X_n)$ by \mathcal{D} . A **test** of H_0 against H_1 is based on a subset C of \mathcal{D} . This set C is called the **critical region** and its corresponding decision rule is

Reject H_0 (accept H_1) if $\mathbf{X} \in C$

Accept H_0 (reject H_1) if $\mathbf{X} \in C^c$

Since the decision rule to take H_0 or H_1 is based on a random sample, the decision could be wrong.

	True State of Nature	
Decision	H_0 is true	H_1 is true
Reject H_0	Type I error	Correct Decision
Accept H_0	Correct Decision	Type II error

In Econometrics, we often refer to statistical significance by controlling the probability making Type I error, say by setting a probability threshold α .

- Basic linear algebra operation.
- Basic set theory.
- Random Variables and Random Vectors.
- Some fundamental distributions: Normal Distribution, Multivariate Normal Distribution, χ^2 -Distribution, t-Distribution, F-Distribution.
- Sampling, Estimation (Estimator), and Statistics.
- Confidence Intervals and Hypothesis Testing.