# Mathematics Preparatory Course - MSc in EEBL Lecture Notes 

Simone De Vito

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## General Information

- Instructor: Simone De Vito
- Email: simone.devito@students.uniroma2.eu
- References: The main reference is "Mathematics for Economists" by C. Simon and L. Blume, Norton \& Company
- Objective: The Mathematics preparatory course aims to review the basic concepts of calculus and linear algebra and to provide students with the necessary tools to understand the notions of Economics where a quantitative approach is needed.
- Course content: The course consists in five lectures which will cover the key concepts of the following subjects:
- One variable calculus
- Linear Algebra
- Function of several variable
- Optimization


## 1 Preliminary notions

### 1.1 Set theory

A set is a well-defined collection of distinct objects, considered as an object in its own right. In order to characterize a set we can directly specify the elements that compose it

Example 1.1. $A=\{a, b, c\}$
or it can be identified listing the features of the elements in the set
Example 1.2. $A=\{n \mid(\mathrm{n}$ is odd $) \wedge(\mathrm{n}$ is positive $)\}$

Once we have the defined what a set is, and we have seen how to characterize it, we can define the notion of subset

Definition 1.1. Let A and B be two distinct sets. If every element in A is also in B we can state that A is a subset of B , denoted by $A \subseteq B$. More formally:

$$
\begin{equation*}
\forall a \in A: a \in A \Rightarrow a \in B \Leftrightarrow A \subseteq B \tag{1}
\end{equation*}
$$



Figure 1: Example of $B \subseteq A$ via Euler diagram

If $A \subseteq B$ and $B \subseteq A$, then A and B coincide $\Rightarrow A=B$.
It is called empty set the unique set having no elements, $\varnothing$, and, conventionally, given every set A, the empty set is subset of A

## Definition 1.2.

$$
\begin{equation*}
\varnothing \subseteq A \forall A \tag{2}
\end{equation*}
$$

### 1.2 Set operations

Let $X$ be a set, and $A$ and $B$ subsets of $X$.

Definition 1.3. Union, denoted by $A \cup B$, is the set containing the elements of A or the elements of B

$$
\begin{equation*}
A \cup B=\{x \in X \mid(x \in A) \vee(x \in B)\} \tag{3}
\end{equation*}
$$



Figure 2: Example of $A \cup B$ via Euler diagram

Example 1.3. Let $A=\{a, b, g\}$ and $B=\{a, c, g\} \Rightarrow A \cup B=\{a, b, c, g\}$
Definition 1.4. Intersection, denoted by $A \cap B$, is the set containing the elements either in A and in B

$$
\begin{equation*}
A \cap B=\{x \in X \mid(x \in A) \wedge(x \in B)\} \tag{4}
\end{equation*}
$$



Figure 3: Example of $A \cap B$ via Euler diagram

Example 1.4. Let $A=\{a, b, g\}$ and $B=\{a, c, g\} \Rightarrow A \cap B=\{a, g\}$
Moreover,
Definition 1.5. Two sets are disjoint if

$$
\begin{equation*}
A \cap B=\varnothing \tag{5}
\end{equation*}
$$

Definition 1.6. Difference, denoted by $A \backslash B$, it is given by the set composed by the elements of $A$ minus the elements which are also in $B \Rightarrow A \backslash B=A-(A \cap B)$

Example 1.5. Let $A=\{a, b, g\}$ and $B=\{a, c, g\} \Rightarrow A \backslash B=\{b\}$



Figure 4: Example of two disjoint sets via Euler diagram


Figure 5: Example of $A \backslash B$ via Euler diagram

### 1.3 Sets of numbers

Throughout this course we will use the following sets of numbers:

- Set of Natural numbers: $\mathbb{N}=\{0,1,2,4, \ldots\}$
- Set of Integers: $\mathbb{Z}=\{0, \pm 1, \pm 2, \pm 4, \ldots\}$
- Set of Rational numbers: $\mathbb{Q}=\{a / b \mid a, b \in \mathbf{Z}, b \neq 0\}$
- Set of Real number $\mathbb{R}$. This set includes all rational numbers, together with all irrational numbers.


### 1.4 Cartesian product

The Cartesian product of two sets $A$ and $B$, denoted $A \times B$, is the set of all the ordered pairs $(a, b)$ where $a \in A$ and $b \in B$.

Example 1.6. Let $A=\{0,1\}$ and $B=\{0,2\}$, we have that

$$
\begin{equation*}
A \times B=\{(0,0),(0,2),(1,0),(1,2)\} \tag{6}
\end{equation*}
$$

while

$$
\begin{equation*}
B \times A=\{(0,0),(0,1),(2,0),(2,1)\} \tag{7}
\end{equation*}
$$

## 2 One Variable Calculus

### 2.1 Introduction

One of the main goals of Economics is to understand mechanisms, interactions and relationships between different variables. Mathematically these relations are described by functions.

Definition 2.1. Let $X$ and $Y$ be two subsets of $\mathbb{R}$. A function $f$ defined on $X$ with values in $Y$ is a correspondence associating to each element $x \in X$ at most one element $y \in Y$.

$$
\begin{equation*}
f: x \rightarrow f(x) \tag{8}
\end{equation*}
$$



Figure 6: Graphic representation of (8)
the variable $x \in X$ is called independent, while the variable $y \in Y$ is called dependent.
Roughly speaking a function will take an input $(x)$ and through some rules will give you back an output (y).

Example 2.1. $f(x)=x^{2}+3 \Rightarrow f(5)=5^{2}+3=28$
The set of elements $x \in X$ to which $f$ assigns an element in $Y$ is called the domain, while the elements $y \in Y$ associated to $x$ are called images. The set of all the images is called range.

Useful rules:

$$
\begin{array}{ll}
\text { Function } & \text { Domain } \\
& \\
\ln g(x) & g(x)>0 \\
\sqrt[2 n]{g(x)} & g(x) \geq 0 \\
\sqrt[2 n+1]{g(x)} & \forall x \in \mathbb{R} \\
\frac{k(x)}{g(x)} & g(x) \neq 0
\end{array}
$$

Example 2.2. Let $f(x)=\frac{1}{x-2} \Rightarrow$ the domain of $f$ is given by all the elements in $\mathbb{R}$ except +2
Some of the most important information of a function are contained in the graph.
The graph of $f$ is the subset of $\Gamma(f)$ of the Cartesian product $X \times Y$ made of pairs ( $x, f(x)$ ) when $x$ varies in the domain of $f$, i.e.

$$
\begin{equation*}
\Gamma(f)=\{(x, f(x)) \in X \times Y: x \in \operatorname{dom} f\} \tag{9}
\end{equation*}
$$

A function may be even or odd.
Let be the function

$$
\begin{equation*}
f: x \rightarrow Y \tag{10}
\end{equation*}
$$

Suppose that if $x \in X$ it is also true that $-x \in X$ (It means that the function is symmetric with respect to the origin) then

Definition 2.2. The function f is even if

$$
\begin{equation*}
f(-x)=f(x), \forall x \in X \tag{11}
\end{equation*}
$$

Example 2.3. $f(x)=x^{2}-3$, since $(-x)^{2}=x^{2}$, then $f(x)=f(-x)$
Definition 2.3. The function $f$ is odd if

$$
\begin{equation*}
f(-x)=-f(x), \forall x \in X \tag{12}
\end{equation*}
$$

Example 2.4. $f(x)=3 x^{3}$, since $(-x)^{3}=-\left(x^{3}\right)$ then $f(-x)=-f(x)$

The basic geometric properties of a function are whether it is increasing or decreasing and the location of its local/global maxima/minima.

Definition 2.4. The function $f$ is increasing if

$$
\begin{equation*}
\forall x_{1}, x_{2} \in \operatorname{dom}(f) \quad x_{1}<x_{2} \Rightarrow f\left(x_{1}\right) \leq f\left(x_{2}\right) \tag{13}
\end{equation*}
$$

Conversely, it is decreasing if

$$
\begin{equation*}
\forall x_{1}, x_{2} \in \operatorname{dom}(f) \quad x_{1}>x_{2} \Rightarrow f\left(x_{1}\right) \geq f\left(x_{2}\right) \tag{14}
\end{equation*}
$$




Figure 7: Example of increasing function - Example of decreasing function

The point where the function turns from decreasing to increasing is a minimum for the function and the point where the function turns from decreasing to increasing is a maximum for the function. If there is no greater (smaller) value of the function in its range from that maximum (minimum) then the maximum (minimum) is called global maximum (minimum).

### 2.2 Function types

- A Polynomial is a map of the form $P(x)=\underbrace{a_{n} x^{n}}_{\text {monomial }}+\ldots+a_{1} x+a_{0}$ where $n$ is the degree of polynomial.
- A rational function is of the kind $R(x)=\frac{P(x)}{Q(x)}$ where P and Q are polynomials.
- An exponential function is defined $f(x)=a^{x}$, the domain it is all $\mathbb{R}$ and it satisfies $y(0)=a^{0}=1$.
- The inverse $y=\log _{a} x$ is called logarithm and it is defined in $(0,+\infty)$.


### 2.3 Linear Function

Polynomials of degree 1 are interesting functions, they are also called Linear function

$$
\begin{equation*}
f(x)=a x+b \tag{15}
\end{equation*}
$$



Figure 8: Example of linear functions
the graph of a linear function is a straight line. In order to draw this function knowing two points in the Cartesian plan is enough.

One of the main features which distinguishes two different lines is the slope (steepness) that is given by $a$. This function is increasing as $a>0$ and decreasing if $a<0$; if $a=0$ the function degenerates to the constant function $f(x)=b$. The slope is given by the ratio of the growth in $y\left(y_{2}-y_{1}\right)$ and the growth of $x\left(x_{2}-x_{1}\right)$, that is $a=\frac{y_{2}-y_{1}}{x_{2}-x_{1}}$.

### 2.4 Quadratic Function

A quadratic function is a polynomial of degree 2. The graph of a univariate quadratic function is a parabola whose axis of symmetry is parallel to the $y$-axis

$$
\begin{equation*}
f(x)=a x^{2}+b x+c \tag{16}
\end{equation*}
$$



Figure 9: Example of quadratic functions

### 2.5 Exponential Function

An exponential function is a function of the form

$$
\begin{equation*}
f(x)=a^{x}, a>0 \tag{17}
\end{equation*}
$$

### 2.6 Logarithmic Function

The logarithm is the inverse of the exponential function.

$$
\begin{equation*}
f(x)=\log _{a} x \tag{18}
\end{equation*}
$$

The logarithm of a given number $x$ is the exponent to which another fixed number, the base $a$, must be raised, to produce that number $x$. A particular case of the logarithm is the natural logarithm which has the number $e$ (that is $e \approx 2.718$ ) as the base.

## 3 Limit and continuity

Limit is the value that a function (or sequence) approaches as the input (or index) approaches some value.

Let $X, Y \subseteq \mathbb{R}$ and $f: X \rightarrow Y$ be a function of real variables. If

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}} f(x)=l \tag{19}
\end{equation*}
$$

means that $f(x)$ can be made to be as close to $l$ as desired, by making $x$ sufficiently close to $x_{0}$. Or equivalently, $f(x)$ goes to $l$ as $x$ approaches $x_{0}$

In the most simple cases, if $f(x)$ is a real valued function which goes to $l$ as $x$ goes to $x_{0}$, we can find $l$ (the limit) substituting $x_{0}$ in the function.

## Example 3.1.

$$
\begin{equation*}
\lim _{x \rightarrow 2}\left(3 x-x^{2}\right)=3 \cdot 2-2^{2}=2 \tag{20}
\end{equation*}
$$

Example 3.2.

$$
\begin{equation*}
\lim _{x \rightarrow 2} \ln (x-1)=\ln (2-1)=\ln (1)=0 \tag{21}
\end{equation*}
$$

The limit of the function can be also $+\infty /-\infty$. In order to find these limits, it is important to know the graph of the function.

## Example 3.3.

$$
\begin{equation*}
\lim _{x \rightarrow+\infty} e^{x}=+\infty \tag{22}
\end{equation*}
$$

In calculus, a one-sided limit is either of the two limits of a function $f(x)$ of a real variable $x$ as $x$ approaches a specified point either from the left or from the right

## Definition 3.1. Right limit

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}^{+}} f(x)=l \tag{23}
\end{equation*}
$$

## Definition 3.2. Left limit

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}^{-}} f(x)=l \tag{24}
\end{equation*}
$$

The $\lim _{x \rightarrow x_{0}} f(x)=l$ exists only if the right limit and the left limit exist and they are equal.
Example 3.4. Suppose

$$
\begin{equation*}
f(x)=\frac{1}{x-1} \tag{25}
\end{equation*}
$$

We want to study

$$
\begin{equation*}
\lim _{x \rightarrow 1} \frac{1}{x-1} \tag{26}
\end{equation*}
$$

When $x \rightarrow 1$ the function will assume infinite values; more precisely if $x \rightarrow 1$ with $x>1$ the function will tend to $+\infty$, but if $x \rightarrow 1$ with $x<1$ the function will tend to $-\infty$. In this case,

$$
\begin{equation*}
\lim _{x \rightarrow 1} \frac{1}{x-1} \tag{27}
\end{equation*}
$$

does not exist since

$$
\begin{equation*}
\lim _{x \rightarrow 1^{+}} f(x)=+\infty \quad \text { and } \quad \lim _{x \rightarrow 1^{-}} f(x)=-\infty \tag{28}
\end{equation*}
$$

### 3.1 Mathematical operations

$$
\begin{align*}
\lim _{x \rightarrow x_{0}}[f(x) \pm g(x)] & =\lim _{x \rightarrow x_{0}} f(x) \pm \lim _{x \rightarrow x_{0}} g(x)  \tag{29}\\
\lim _{x \rightarrow x_{0}} f(x) g(x) & =\lim _{x \rightarrow x_{0}} f(x) \lim _{x \rightarrow x_{0}} g(x)  \tag{30}\\
\lim _{x \rightarrow x_{0}} \frac{f(x)}{g(x)} & =\frac{\lim _{x \rightarrow x_{0} f(x)}}{\lim _{x \rightarrow x_{0}} g(x)} \tag{31}
\end{align*}
$$

We may end up in a indeterminate form

$$
\begin{equation*}
\frac{0}{0}, \frac{\infty}{\infty}, 0 \times \infty, \infty-\infty \tag{32}
\end{equation*}
$$

## Example 3.5.

$$
\begin{equation*}
\lim _{x \rightarrow 0} \frac{2 x^{2}-3 x^{5}}{x+2 x^{2}}=\lim _{x \rightarrow 0} \frac{2 x^{2}\left(1-3 x^{3} / 2\right)}{x(1+2 x)}=\lim _{x \rightarrow 0} \frac{2 x\left(1-3 x^{3} / 2\right)}{1+2 x}=0 \tag{33}
\end{equation*}
$$

## Example 3.6.

$$
\begin{equation*}
\lim _{x \rightarrow+\infty} \frac{2 x^{3}+5 x^{2}-x+7}{4 x^{3}-x^{2}+x-3}=\frac{2 x^{3}\left(1+\frac{5}{2 x}-\frac{1}{2 x^{2}}+\frac{7}{2 x^{3}}\right)}{4 x^{3}\left(1-\frac{1}{4 x}+\frac{1}{4 x^{2}}-\frac{3}{4 x^{3}}\right)}=\frac{1}{2} \tag{34}
\end{equation*}
$$

## Example 3.7.

$$
\begin{equation*}
\lim _{x \rightarrow+\infty} e^{x}-\sqrt{x}=+\infty \tag{35}
\end{equation*}
$$

### 3.2 Continuity

We could say if a function is continuous just looking at the graph. A function is continuous if its graph has no breaks (no jumps for the same point) in its domain. However, more rigorously, let $f: X \rightarrow Y$ and $x_{0} \in X . f(x)$ is continuous in $x_{0}$ if

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}} f(x)=f\left(x_{0}\right) \tag{36}
\end{equation*}
$$

### 3.3 Differential

Steepness is a key concept in Economics. We are often interested in evaluating what is the effect an increase of the independent variable to the dependent variable. As we have already seen, regarding linear function the steepness is easily given by the slope coefficient. However,
most of the economic functions are non-linear. How do we measure the marginal effects of these nonlinear functions?

Let $f: X \rightarrow Y$ and $x_{0} \in X$. If $x$ increases by $\Delta x$, e.g. from $x_{0}$ to $x_{0}+\Delta x$, also the function $f(x)$ will increase as follows

$$
\begin{equation*}
\Delta f(x)=f\left(x_{0}+\Delta x\right)-f\left(x_{0}\right) \tag{37}
\end{equation*}
$$



Figure 10: Graphic representation of (37)

However, mostly, it is more interesting studying Difference quotient

$$
\begin{equation*}
\frac{\Delta f(x)}{\Delta x}=\frac{f\left(x_{0}+\Delta x\right)-f\left(x_{0}\right)}{\Delta x} \tag{38}
\end{equation*}
$$

particularly, when $\Delta x$ is very small

$$
\begin{equation*}
\lim _{\Delta x \rightarrow 0} \frac{f\left(x_{0}+\Delta x\right)-f\left(x_{0}\right)}{\Delta x} \tag{39}
\end{equation*}
$$

the limit of the difference quotient as h approaches zero, if it exists, should represent the slope of the tangent line to $\left(x_{0}, f\left(x_{0}\right)\right)$. This limit is defined to be the derivative of the function $f$ at $x_{0}$ :

$$
\begin{equation*}
\lim _{\Delta x \rightarrow 0} \frac{f\left(x_{0}+\Delta x\right)-f\left(x_{0}\right)}{\Delta x}=f^{\prime}\left(x_{0}\right) \tag{40}
\end{equation*}
$$

if the limit of the difference quotient does not exist $f(x)$ is said to be not not differential. If a function is differentiable then it continous the opposite it is not true.

Some derivatives

| Function | Derivative | Function | Derivative |
| :--- | :--- | :--- | :--- |
| $f(x)=k$ | $f^{\prime}(x)=0$ | $f(x)=\sin (x)$ | $f^{\prime}(x)=\cos (x)$ |
| $f(x)=x$ | $f^{\prime}(x)=1$ | $f(x)=\cos (x)$ | $f^{\prime}(x)=-\sin (x)$ |
| $f(x)=x^{\alpha}, \alpha \in \mathbb{R}$ | $f^{\prime}(x)=\alpha x^{\alpha-1}$ | $f(x)=\cot (x)$ | $f^{\prime}(x)=-\frac{1}{\sin ^{2}(x)}$ |
| $f(x)=\alpha^{x}$ | $f^{\prime}(x)=\alpha^{x} \ln (\alpha)$ | $f(x)=\tan (x)$ | $f^{\prime}(x)=\frac{1}{\cos ^{2}(x)}$ |
| $f(x)=e^{x}$ | $f^{\prime}(x)=e^{x}$ | $f(x)=\arcsin (x)$ | $f^{\prime}(x)=\frac{1}{\sqrt{1-x^{2}}}$ |
| $f(x)=\log _{\alpha}(x)$ | $f^{\prime}(x)=\frac{1}{x \ln (\alpha)}$ | $f(x)=\arccos (x)$ | $f^{\prime}(x)=-\frac{1}{\sqrt{1-x^{2}}}$ |
| $f(x)=\ln (x)$ | $f^{\prime}(x)=\frac{1}{x}$ | $f(x)=\arctan (x)$ | $f^{\prime}(x)=\frac{1}{1+x^{2}}$ |
| $f(x)=\|x\|$ | $f^{\prime}(x)=\frac{\|x\|}{x}$ | $f(x)=\operatorname{arccot}(x)$ | $f^{\prime}(x)=-\frac{1}{1+x^{2}}$ |

### 3.4 Derivative of elementary function

- $f(x)=k$, then $f^{\prime}(x)=0$, proof;

$$
\lim _{\Delta x \rightarrow 0} \frac{f(x+\Delta x)-f(x)}{\Delta x}=\lim _{\Delta x \rightarrow 0} \frac{k-k}{\Delta x}=0
$$

- $f(x)=x$, then $f^{\prime}(x)=1$, proof;

$$
\lim _{\Delta x \rightarrow 0} \frac{f(x+\Delta x)-f(x)}{\Delta x}=\lim _{\Delta x \rightarrow 0} \frac{x+\Delta x-x}{\Delta x}=1
$$

- $D[\alpha f(x)]=\alpha D f(x), \alpha \in \mathbb{R}$,
- $D[f(x)+g(x)]=D f(x)+D g(x)$
- $D[f(x) g(x)]=[D f(x)] g(x)+f(x[D g(x)])$
- if $g(x) \neq 0, D\left[\frac{f(x)}{g(x)}\right]=\frac{f^{\prime}(x) g(x)-f(x) g^{\prime}(x)}{[g(x)]^{2}}$
- The chain rule is a fundamental theorem in calculus used to compute the derivative of the composition of two or more functions. If we have two functions $f$ and $g$ such that $y=f(u)$ and $u=g(x)$, then the composition of these functions is $y=f(g(x))$. The chain rule states that the derivative of $y$ with respect to $x$ is given by:

$$
\frac{d y}{d x}=\frac{d y}{d u} \cdot \frac{d u}{d x}
$$

This can be written as:

$$
y^{\prime}(x)=f^{\prime}(g(x)) \cdot g^{\prime}(x)
$$

For a function to be differentiable, it must be continuous but not vice versa.

Example 3.8. $f(x)=|x|$ is continuous but not differentiable $\forall x \in \mathbf{D}$

$$
\begin{equation*}
f^{\prime}(0)=\lim _{h \rightarrow 0} \frac{f(0+h)-f(0)}{h}=\lim _{h \rightarrow 0} \frac{|h|}{h} \tag{41}
\end{equation*}
$$

Therefore, $f(x)=|x|$ is not differentiable in 0 .

### 3.5 Higher order derivative

If $f(x)$ be a differentiable function, and $f^{\prime}(x)$ its derivative. If $f^{\prime}(x)$ is again differentiable, $D\left[f^{\prime}(x)\right]=f^{\prime \prime}(x)$ is its second derivative. If we can do the same with the second derivative we will have the third derivatie, and so on and so forth. Generally, we can talk about higher order derivative, defined as $f^{(n)}(x)$.

### 3.6 Max and min

Let $f: X \rightarrow \mathbb{R}$ be a function, and $x_{0}$ a point in the domain in X of $f(x)$. We can state that

- $x_{0}$ is a minimum if $\exists I_{x_{0}}$ s.t.

$$
\begin{equation*}
f(x)>f\left(x_{0}\right) \forall x \in I_{x_{0}} \backslash\left\{x_{0}\right\} \tag{42}
\end{equation*}
$$

- $x_{0}$ is a maximum if $\exists I_{x_{0}}$ s.t.

$$
\begin{equation*}
f(x)<f\left(x_{0}\right) \forall x \in I_{x_{0}} \backslash\left\{x_{0}\right\} \tag{43}
\end{equation*}
$$

- If $x_{0}$ is a $\max / \min$, then $f^{\prime}\left(x_{0}\right)=0$.
- If $f^{\prime}(x)>0 \forall x \in(a, b)$, then $f(x)$ is increasing in $(a, b)$
- If $f^{\prime}(x)<0 \forall x \in(a, b)$, then $f(x)$ is decreasing in $(a, b)$
- While the first derivative $f^{\prime}(x)$ provides the slope of the tangent line to the function at any given point, the second derivative $f^{\prime \prime}(x)$ gives information about the curvature or concavity of the function. If $f^{\prime \prime}(x)>0$, the function is convex (shaped like a cup) at that point, indicating that the slope of the tangent line is increasing. If $f^{\prime \prime}(x)<0$, the function is concave (shaped like a cap) at that point, indicating that the slope is decreasing.



Figure 11: Example of Minimimum - Example of Maximum

## 4 Integral

The indefinite integral represents a class of functions (antiderivative) whose derivative is the integrand.

Definition 4.1. Let $f(x)$ be a function in $(a, b)$. If a function $F(x)$ exists and it is continous in $[a, b]$ and it is also differentiable in $(a, b)$ such that

$$
\begin{equation*}
F^{\prime}(x)=f(x) \quad \forall x \in(a, b) \tag{44}
\end{equation*}
$$

$F(x)$ is said to be the antiderivative of $f(x)$.
If $F(x)$ is an antiderivative function of $f(x)$, also $G(x)=F(x)+c, c \in \mathbb{R}$ is an antiderivative of $f(x)$, given that

$$
\begin{equation*}
G^{\prime}(x)=D[F(x)+c]=F^{\prime}(x)=f(x) \tag{45}
\end{equation*}
$$

Definition 4.2. Let $f(x)$ be a function in $(a, b)$ and let it have antiderivatives. All the antiderivatives of $f(x)$ are defined as

$$
\begin{equation*}
\int f(x) d x \tag{46}
\end{equation*}
$$

$f(x)$ is said to be the integrand.
As the derivative, the integral is a linear operator, therefore

$$
\begin{align*}
\int[\alpha f(x)+\beta g(x)] d x & =\alpha \int f(x) d x+\beta \int g(x) d x \\
& =\alpha\left[F(x)+c_{1}\right]+\beta\left[G(x)+c_{2}\right]  \tag{47}\\
& =\alpha F(x)+\beta G(x)+c, \quad c \in \mathbb{R}
\end{align*}
$$

|  |  |
| :--- | :--- |
| $\int x^{\alpha} d x$ | $\frac{x^{\alpha+1}}{\alpha+1}, c \in \mathbb{R}$ if $\alpha \neq-1$ |
| $\int \frac{1}{\alpha}$ | $\ln \|x\|+c, c \in \mathbb{R}$ |
| $\int[f(x)]^{\alpha} f^{\prime}(x) d x$ | $\frac{[f(x)]^{\alpha+1}}{\alpha+1}+c, c \in \mathbb{R}$ if $\alpha \neq-1$ |
| $\int \frac{f^{\prime}(x)}{f(x)} d x$ | $\ln \|f(x)\|+c, c \in \mathbb{R}$ |
| $\int e^{f(x)} f^{\prime}(x) d x$ | $e^{f(x)}+c, c \in \mathbb{R}$ |
| $\int e^{x} d x$ | $e^{x}+c, c \in \mathbb{R}$ |
| $\int \cos [f(x)] f^{\prime}(x) d x$ | $\sin [f(x)]+c, c \in \mathbb{R}$ |
| $\int \sin [f(x)] f^{\prime}(x) d x$ | $-\cos [f(x)]+c, c \in \mathbb{R}$ |

The integration by part rule is presented here:

$$
\begin{equation*}
\int f^{\prime}(x) g(x) d x=f(x) g(x)-\int f(x) g^{\prime}(x) d x+c, c \in \mathbb{R} \tag{48}
\end{equation*}
$$

Sometimes, in order to evaluate

$$
\begin{equation*}
\int f(x) d x \tag{49}
\end{equation*}
$$

it is useful to make a change of variables. For instance, let $g(t)$ be a derivable function, and assume

$$
\begin{equation*}
x=g(t) \tag{50}
\end{equation*}
$$

then, $d x$ is changed with

$$
\begin{equation*}
d x=g^{\prime}(t) d t \tag{51}
\end{equation*}
$$

Therefore, we can evaluate Equation (48) as

$$
\begin{equation*}
\int f(g(t)) g^{\prime}(t) d t=G(t)+c, c \in \mathbb{R} \tag{52}
\end{equation*}
$$

and then solving

$$
\begin{equation*}
t=g^{-1}(x) \tag{53}
\end{equation*}
$$

Sometimes, we are interested in evaluating the definite integral; it can be interpreted informally as the signed area of the region in the $x y$-plane that is bounded by the graph of $f(x)$, the x -axis and the vertical lines $x=a$ and $x=b$. If $f(x)$ is a continuous real-valued function defined on a closed interval $[a, b]$, then, once an antiderivative $F(x)$ of $f(x)$ is known, the definite integral of $f(x)$ over that interval is given by

$$
\begin{equation*}
\int_{a}^{b} f(x) d x=[F(x)]_{a}^{b}=F(b)-F(a) \tag{54}
\end{equation*}
$$

## 5 Matrix algebra

The analysis of many economic models reduces to the study of systems of equations. We will, look at the simplest possible system if equations - linear system.

Definition 5.1. Generally, an equation is said to be linear if it has the form

$$
\begin{equation*}
a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{n} x_{n}=b \tag{55}
\end{equation*}
$$

where $a_{1}, \ldots, a_{n}$ are parameters and $x_{1}, \ldots, x_{n}$ are variables.
We can solve linear system of equations using matrices.

## 6 Matrix algebra

We can generate $\mathbb{R}^{n}$ by $n$-ary Cartesian product of $\mathbb{R}$. Therefore, we can represent $\mathbb{R}^{2}$ and $\mathbb{R}^{3}$ though arrows which connect the origin of the cartesian plan to the coordinates given by the pair $\left(x_{1}, x_{2}\right)$ or $\left(x_{1}, x_{2}, x_{3}\right)$

Definition 6.1. Let $\underline{v}_{1}=\left(x_{1}, \ldots, x_{n}\right)$ and $\underline{v}_{2}=\left(y_{1}, \ldots, y_{n}\right)$ be two vectors, their are said to be equal if $x_{i}=y_{i}, \forall i=1,2, . ., N$.


Figure 12: Graphic representation of the vector in $\mathbb{R}^{2} \underline{v}=(2,3)$


Figure 13: Graphic representation of the vector in $\mathbb{R}^{3}$

Definition 6.2. Let $\underline{v}_{1}=\left(x_{1}, \ldots, x_{n}\right)$ and $\underline{v}_{2}=\left(y_{1}, \ldots, y_{n}\right)$ be two vectors, $\underline{v}_{3}=\left(x_{1}+y_{1}, \ldots, x_{n}+\right.$ $\left.y_{n}\right)$ is the sum of former vectors

We can multiply a vector by a scalar, and we will obtain $c \underline{v}_{1}=\left(c x_{1}, \ldots, c x_{n}\right)$ Given $m$ vectors $\underline{v}_{1}, \underline{v}_{2}, \ldots, \underline{v}_{m}$ of $\mathbb{R}^{n}$ and $m$ scalars $c_{1}, c_{2}, \ldots, c_{m}$ in $\mathbb{R}$, a linear combination of the vectors $\underline{v}_{1}, \underline{v}_{2}, \ldots, \underline{v}_{m}$ with coefficients $c_{1}, c_{2}, \ldots, c_{m}$ the vector given by

$$
\begin{equation*}
\underline{x}=\sum_{k=1}^{m} c_{k} \underline{v}_{k} \tag{56}
\end{equation*}
$$

Definition 6.3. A sequence of vectors $\left(\underline{v_{1}}, \underline{v_{2}}, \ldots, \underline{v_{k}}\right)$ from a vector space $V$ is said to be linearly dependent, if there exist scalars $a_{1}, a_{2}, \ldots, a_{k}$, not all zero, such that

$$
a_{1} \underline{v_{1}}+a_{2} \underline{v_{2}}+\cdots+a_{k} \underline{v_{k}}=0
$$

On the other hand, two vectors are said to be linearly independent if there is no linear combination which gives as a result the null vector expect for the one with null coefficients.

### 6.1 Matrix algebra

A matrix is a rectangular array of numbers. The size of a matrix is indicated by the number of its rows and number of its columns. A matrix with $k$ rows and $n$ columns is called a $k \times n$ matrix. The element in row $i$ and column $j$ is called the $(i, j)$ th entry, and it is often written as $a_{i j}$. A matrix with the number of columns equal to the number of rows is called square matrix.

$$
A_{k, n}=\left(\begin{array}{cccc}
a_{1,1} & a_{1,2} & \cdots & a_{1, n}  \tag{57}\\
a_{2,1} & a_{2,2} & \cdots & a_{2, n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{k, 1} & a_{k, 2} & \cdots & a_{k, n}
\end{array}\right)
$$

### 6.2 Operations

Definition 6.4. Let A and B two matrices kxn, their sum is the matrix $C$ whose elements are $c_{i j}=a_{i j}+b_{i j} \forall i, j$

$$
\begin{align*}
&\left(\begin{array}{ccc}
a_{11} & \cdots & a_{1 n} \\
\vdots & a_{i j} & \vdots \\
a_{k 1} & \cdots & a_{k n}
\end{array}\right)+\left(\begin{array}{ccc}
b_{11} & \cdots & b_{1 n} \\
\vdots & b_{i j} & \vdots \\
b_{k 1} & \cdots & b_{k n}
\end{array}\right) \\
&=\left(\begin{array}{ccc}
a_{11}+b_{11} & \cdots & a_{1 n}+b_{1 n} \\
\vdots & a_{i j}+b_{i j} & \vdots \\
a_{k 1}+b_{k 1} & \cdots & a_{k n}+b_{k n}
\end{array}\right) \tag{58}
\end{align*}
$$

Matrices may be multiplied by scalars. This operation is called scalar multiplication. More generally:

Definition 6.5. The product of the matrix $A$ and the number $\alpha$, denoted $\alpha A$, is the matrix whose elements are $\alpha a_{i j} \forall i, j$

$$
\alpha\left(\begin{array}{ccc}
a_{11} & \cdots & a_{1 n}  \tag{59}\\
\vdots & a_{i j} & \vdots \\
a_{k 1} & \cdots & a_{k n}
\end{array}\right)=\left(\begin{array}{ccc}
\alpha a_{11} & \cdots & \alpha a_{1 n} \\
\vdots & \alpha a_{i j} & \vdots \\
\alpha a_{k 1} & \cdots & \alpha a_{k n}
\end{array}\right)
$$

We can define the matrix product $A B$ iff:

$$
\begin{equation*}
\text { number of columns of } A=\text { number of rows of } B \tag{60}
\end{equation*}
$$

To obtain the $(i, j)$ th entry of $A B$, multiply the $i$ th row of $A$ and the $j$ th column of $B$ as follows

$$
\left(\begin{array}{cccc}
a_{i 1} & a_{i 2} & \cdots & a_{i m}
\end{array}\right) \cdot\left(\begin{array}{c}
b_{1 j}  \tag{61}\\
b_{2 j} \\
\vdots \\
b_{n i}
\end{array}\right)=a_{i 1} b_{1 j}+a_{i 2} b_{2 j}+\cdots+a_{i m} b_{m j}
$$

In other words, the $(i, j)$ th entry of the product $A B$ is defined to be

$$
\sum_{n-1}^{m} a_{i n} b_{h j}
$$

If $A$ is a $k \times m$ and $B$ is $m \times n$, then the product $C=A B$ will be $k \times n$.
Usually, $A B \neq B A$.
The $n \times n$ matrix

$$
I=\left(\begin{array}{cccc}
1 & 0 & \cdots & 0  \tag{62}\\
0 & 1 & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & 1
\end{array}\right)
$$

with $a_{i i}=1, \forall i$ and $a_{i j}=0, \forall i \neq j$, has the following property

$$
\begin{equation*}
A I=A \tag{63}
\end{equation*}
$$

for any $m x n$ matrix $A . I$ is called identity matrix.

Definition 6.6. The transpose of a $k x n$ matrix $A$ is the $n x k$ matrix obtained by interchanging the rows and the columns.

$$
\left(\begin{array}{ccc}
a_{11} & a_{12} & a_{13}  \tag{64}\\
a_{21} & a_{22} & a_{23}
\end{array}\right)^{T}=\left(\begin{array}{ll}
a_{11} & a_{21} \\
a_{12} & a_{22} \\
a_{13} & a_{23}
\end{array}\right)
$$

Definition 6.7. A Triangular matrix is a square matrix containing element different from zero only above/below the main diagonal.

## Example 6.1. Upper Triangular

$$
A=\left(\begin{array}{ccc}
1 & -3 & -1  \tag{65}\\
0 & 1 & 1 \\
0 & 0 & -2
\end{array}\right)
$$

## Example 6.2. Lower Triangular

$$
A=\left(\begin{array}{ccc}
1 & 0 & 0  \tag{66}\\
-2 & 1 & 0 \\
3 & 5 & -2
\end{array}\right)
$$

Definition 6.8. To any square matrix is associated a determinant, $|A|$. From a geometric point of view, it represents the area (volume) of the parallelogram generated by the vectors of the matrix.

Let $A$ be a $2 x 2$ matrix,

$$
A=\left(\begin{array}{ll}
a_{11} & a_{12}  \tag{67}\\
a_{21} & a_{22}
\end{array}\right)
$$

the determinant is $|A|=a_{11} a_{22}-a_{12} a_{21}$

In case of a $n \times n$ matrix with $n>3$, in order to find the determinant, we can use the minor of the matrix.

Definition 6.9. A minor of a matrix $A$ is the determinant of some smaller square matrix, cut down from $A$ by removing one or more of its rows and columns. Minors obtained by removing just one row and one column from square matrices (first minors) are required for calculating matrix cofactors, which in turn are useful for computing both the determinant and inverse of square matrices.

The cofactor is:

$$
\begin{equation*}
A_{i k}=(-1)^{i+k} M_{i k} \tag{68}
\end{equation*}
$$

where $M_{i k}$ is a minor of the matrix.
Therefore, using the Laplace theorem we can obtain the determinant of a matrix $n \times n$ as the sum of the product, of any row or column, by their cofactor.

The $\operatorname{rank}$ of $A$ is the largest order of any non-zero minor in $A$.

Definition 6.10. Let $A$ be a $n \times n$ matrix. The $n \times n$ matrix $A^{-1}$ is an inverse for $A$ if $A A^{-1}=A^{-1} A=I_{n}$

A matrix can have at most one inverse
$\Rightarrow$ Nor every matrix are invertible. In order to be invertible, $|A| \neq 0$
Definition 6.11. Any symmetric $A$ is:

- Positive semidefinite if $x^{\prime} A x>0, \forall x \neq 0$
- Positive definite if $x^{\prime} A x \geq 0, \forall x \neq 0$
- Negative definite if $x^{\prime} A x<0, \forall x \neq 0$
- Negative semidefinite if $x^{\prime} A x \leq 0, \forall x \neq 0$


### 6.3 Linear system of equations

As we have stated above, we can use matrix to solve linear system of equations.
Recall: Generally, an equation is said to be linear if it has the form

$$
\begin{equation*}
a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{n} x_{n}=b \tag{69}
\end{equation*}
$$

where $a_{1}, \ldots, a_{n}$ are parameters and $x_{1}, \ldots, x_{n}$ are variables.
The solution of the linear equation is given by $\left(x_{1}, \ldots, x_{n}\right)$ which substituted into the equation solving it.

If there are several linear equations which have to be true all toghether we talk about Linear

## System of equations

## Example 6.3.

$$
\left\{\begin{array}{l}
a_{11} x_{1}+a_{12} x_{2}+\ldots a_{1 n} x_{n}=b_{1}  \tag{70}\\
a_{21} x_{1}+a_{22} x_{2}+\ldots a_{2 n} x_{n}=b_{2} \\
\vdots \\
a_{m 1} x_{1}+a_{m 2} x_{2}+\ldots a_{m n} x_{n}=b_{m}
\end{array}\right.
$$

The solution of the system is given by $\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ which solves all the equations contemporaneously.

The system above, can be expressed much more compactly using matrix notations. Let $A$ denote the coefficient matrix of the system:

$$
A=\left(\begin{array}{ccc}
a_{11} & \cdots & a_{1 n}  \tag{71}\\
\vdots & a_{i j} & \vdots \\
a_{m 1} & \cdots & a_{m n}
\end{array}\right)
$$

Also, let

$$
\mathbf{x}=\left(\begin{array}{c}
x_{1}  \tag{72}\\
\vdots \\
x_{n}
\end{array}\right) \quad \text { and } \quad \mathbf{b}=\left(\begin{array}{c}
b_{1} \\
\vdots \\
b_{m}
\end{array}\right)
$$

Then, the system of equations can be written as

$$
\left(\begin{array}{ccc}
a_{11} & \cdots & a_{1 n}  \tag{73}\\
\vdots & a_{i j} & \vdots \\
a_{m 1} & \cdots & a_{m n}
\end{array}\right) \cdot\left(\begin{array}{c}
x_{1} \\
\vdots \\
x_{n}
\end{array}\right)=\left(\begin{array}{c}
b_{1} \\
\vdots \\
b_{m}
\end{array}\right)
$$

or simply as

$$
\begin{equation*}
A \mathbf{x}=\mathbf{b} \tag{74}
\end{equation*}
$$

Then, if $A$ is nonsingular $(|A| \neq 0)$, we can solve the system as $\mathbf{x}=A^{-1} \mathbf{b}$. To solve a linear system of simultaneous equations we can use also the Cramer's rule. If the matrix A is nonsingular, the linear system of system of $n$ linear equations and $n$ unknowns. Then the theorem states that in this case the system has a unique solution, whose individual values for the unknowns are given by:

$$
\begin{equation*}
x_{i}=\frac{\operatorname{det}\left(A_{i}\right)}{\operatorname{det}(A)} \quad i=1, \ldots, n \tag{75}
\end{equation*}
$$

where $A_{i}$ is the matrix formed by replacing the $i$ th column of $A$ by the column vector $b$.

Example 6.4. Solve the following system

$$
\left\{\begin{array}{l}
x-2 y+z=1  \tag{76}\\
3 x+y-7 z=0 \\
x-z=1
\end{array}\right.
$$

Since $|A| \neq 0$, we can use Cramer's rule

$$
x=\frac{\operatorname{det}\left(\begin{array}{ccc}
1 & -2 & 1 \\
0 & 1 & -7  \tag{77}\\
1 & 0 & -1
\end{array}\right)}{x^{6}\left(\begin{array}{ccc}
1 & 1 & 1 \\
3 & 0 & -7 \\
1 & 1 & -1
\end{array}\right)}=2
$$

In Economics, we may be interested in system of the following form

$$
\begin{equation*}
A x=\lambda x \tag{78}
\end{equation*}
$$

where $A$ is a square matrix.
It is equal to write:

$$
\begin{equation*}
(A-\lambda I) x=0 \tag{79}
\end{equation*}
$$

Definition 6.12. The values $\lambda$ that solve $\operatorname{det}(A-\lambda I)=0$ are called eigenvalues.

Definition 6.13. While the non-trivial vectors, $\mathbf{x}$, obtained as the solution of $(A-\lambda I) x=0$ is called eigenvector.

From a geometric point of view, the eigenvector $(x)$ is the vector which is only scaled by a value $\lambda$ when we apply to it a transformation $A$.

## 7 Function of several variables

As we have seen before, in Economics is interested looking what is the effect of a change in one variable with respect to another one. However, in most real cases, variables depend on several variables. We may be interested in $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$.

We can graph function up to $\mathbb{R}^{3}$.


Figure 14: The graph of $f(x, y)=x^{2}+y^{2}$


Figure 15: The graph of $f(x, y)=x^{2}-y^{2}$

There is another way to visualize function from $\mathbb{R}^{2}$ to $\mathbb{R}^{1}$ which requires only two dimensional sketching: the level curves. They are given by

$$
\begin{equation*}
L_{c}=(x, y): f(x, y)=c \tag{80}
\end{equation*}
$$

It is like to slice the function in many pieces.

### 7.1 Calculus of several variable

When we deal with function of several variable, we are often interested in the partial variation - the variation brought about by the change in only one variable. We will talk, therefore, about partial derivative.

Definition 7.1. Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$. Then for each variable $x_{i}$ at each point $x^{0}$ in the domain of $f$

$$
\begin{equation*}
\frac{\partial f}{\partial x_{i}}\left(x_{1}^{0}, \ldots, x_{n}^{0}\right)=\lim _{h \rightarrow 0} \frac{f\left(x_{1}^{(0)}, \ldots, x_{i}^{0}+h, \ldots, x_{n}^{0}\right)-f\left(x_{1}^{0}, \ldots, x_{i}^{0}, \ldots, x_{n}^{0}\right)}{h} \tag{81}
\end{equation*}
$$

From a practical point of view, we can apply the same rules we apply to one variable function; we take the derivative with respect one variable treating all the other variables as constant.

If we are interested in the behavior of a function $F(x, y)$ of two variables in the neighborhood of a given point $\left(x^{\star}, y^{\star}\right)$, we can, therefore, look at the effect of the change from $x^{\star}\left(y^{\star}\right)$ to $x^{\star}+\Delta x$ $\left(y^{\star}+\Delta y\right)$ respectively:

$$
\begin{align*}
& F\left(x^{*}+\Delta x, y^{*}\right)-F\left(x^{*}, y^{*}\right) \approx \frac{\partial F}{\partial x}\left(x^{*}, y^{*}\right) \Delta x  \tag{82}\\
& F\left(x^{*}, y^{*}+\Delta y\right)-F\left(x^{*}, y^{*}\right) \approx \frac{\partial F}{\partial y}\left(x^{*}, y^{*}\right) \Delta y \tag{83}
\end{align*}
$$

or we can look at the effect of a change of $x$ and $y$ simultaneously as follows

$$
\begin{equation*}
F\left(x^{*}+\Delta x, y^{*}+\Delta y\right)-F\left(x^{*}, y^{*}\right) \approx \frac{\partial F}{\partial x}\left(x^{*}, y^{*}\right) \Delta x+\frac{\partial F}{\partial y}\left(x^{*}, y^{*}\right) \Delta y \tag{84}
\end{equation*}
$$

### 7.2 Higher-order derivatives

The partial derivative $\partial f / \partial x_{i}$ is itself a function - as we have seen in the case of the second derivative. When we take the first derivative with respect to a variable and the second one with respect to another variable we talk about mixed partial derivatives and it is usually written as

$$
\begin{equation*}
\frac{\partial^{2} f}{\partial x_{j} \partial x_{i}} \tag{85}
\end{equation*}
$$

with $i \neq j$.
If a function has $n$ variables, then, it will have $n^{2}$ second order partial derivatives. It is common to arrange these $n^{2}$ partial derivatives into an $n \times n$ matrix whose $(i, j)$ th entry is the $\left(\partial^{2} f / \partial x_{j} \partial x_{i}\right)\left(x^{*}\right)$. This matrix is called Hessian matrix:

$$
D^{2} f_{\mathbf{x}}=\left(\begin{array}{cccc}
\frac{\partial^{2} f}{\partial x_{1}^{2}} & \frac{\partial^{2} f}{\partial x_{2} \partial x_{1}} & \cdots & \frac{\partial^{2} f}{\partial x_{n} \partial x_{1}}  \tag{86}\\
\frac{\partial^{2} f}{\partial x_{1} \partial x_{2}} & \frac{\partial^{2} f}{\partial x_{2}^{2}} & \cdots & \frac{\partial^{2} f}{\partial x_{n} \partial x_{2}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial^{2} f}{\partial x_{1} \partial x_{n}} & \frac{\partial^{2} f}{\partial x_{2} \partial x_{n}} & \cdots & \frac{\partial^{2} f}{\partial x_{n}^{2}}
\end{array}\right)
$$

Since $\frac{\partial^{2} f}{\partial x_{j} \partial x_{i}}=\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}$, the Hessian matrix is often symmetric.

## 8 Optimization

We want to find the point $(x, y)$ s.t. $f^{\prime}(x)=0$, this point is called stationary point (critical point).

A critical point may correspond to:

- Maximum

- Minimum



## - Saddle point

Once we have found the stationary point with $f^{\prime}(x)=0$ (first order condition), we can check what kind of stationary point is looking at the second derivative (second order conditions)


- if $f^{\prime \prime}(x)<0 \Rightarrow \max$
- if $f^{\prime \prime}(x)>0 \Rightarrow$ min
- if $f^{\prime \prime}(x)=0 \Rightarrow$ No clue. It can be max, min or saddle point in $x^{\star}$ we have a local max if $f\left(x^{\star}\right) \geq f(x), \forall x$ in the of $x^{\star}$


### 8.1 Optimization with several variables

The gradient vector is the vector whose components are the partial derivatives of f .

$$
\nabla f(p)=\left[\begin{array}{c}
\frac{\partial f}{\partial x_{1}}(p)  \tag{87}\\
\vdots \\
\frac{\partial f}{\partial x_{n}}(p)
\end{array}\right]
$$

The Hessian Matrix is the $n \times n$ matrix whose $(i, j)$ th entry is the $\left(\partial^{2} f / \partial x_{j} \partial x_{i}\right)\left(\mathbf{x}^{*}\right)$.

$$
D^{2} f_{\mathbf{x}}=\left(\begin{array}{cccc}
\frac{\partial^{2} f}{\partial x_{1}^{2}} & \frac{\partial^{2} f}{\partial x_{2} \partial x_{1}} & \cdots & \frac{\partial^{2} f}{\partial x_{n} \partial x_{1}}  \tag{88}\\
\frac{\partial^{2} f}{\partial x_{1} \partial x_{2}} & \frac{\partial^{2} f}{\partial x_{2}^{2}} & \cdots & \frac{\partial^{2} f}{\partial x_{n} \partial x_{2}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial^{2} f}{\partial x_{1} \partial x_{n}} & \frac{\partial^{2} f}{\partial x_{2} \partial x_{n}} & \cdots & \frac{\partial^{2} f}{\partial x_{n}^{2}}
\end{array}\right)
$$

In order to have a stationay point $\nabla f(p)=0$
In order to attain the second order conditions:

- If the Hessian is positive definite (equivalently, has all eigenvalues positive) at $a$, then $f$ attains a local minimum at $a$
- If the Hessian is negative definite (equivalently, has all eigenvalues negative) at $a$, then $f$ attains a local maximum at $a$
- If the Hessian has both positive and negative eigenvalues then $a$ is a saddle point for $f$ (and in fact this is true even if $a$ is degenerate).

In those cases not listed above, the test is inconclusive.

### 8.2 Constrained optimization

Let $f$ and $h$ be $\mathcal{C}^{1}$ functions of two variables. Suppose that $\mathbf{x}^{\star}=\left(x_{1}^{\star}, x_{2}^{\star}\right)$ is a solution of the problem

$$
\begin{array}{ll}
\operatorname{maximize} & f\left(x_{1}, x_{2}\right)  \tag{89}\\
\text { subject to } & h\left(x_{1}, x_{2}\right)=c
\end{array}
$$

Suppose that $\left(x_{1}^{\star}, x_{2}^{\star}\right)$ is not a critical point of $h$. Then, there is a real number $\mu^{\star}$ s.t. $\left(x_{1}^{\star}, x_{2}^{\star}, \mu^{\star}\right)$ is a critical point of the Lagrangian function

$$
\begin{equation*}
L\left(x_{1}, x_{2}, \mu\right) \equiv f\left(x_{1}, x_{2}\right)-\mu\left[h\left(x_{1}, x_{2}\right)-c\right] \tag{90}
\end{equation*}
$$

In other words at $\left(x_{1}^{*}, x_{2}^{*}, \mu^{*}\right)$

$$
\begin{equation*}
\frac{\partial L}{\partial x_{1}}=0, \quad \frac{\partial L}{\partial x_{2}}=0, \quad \text { and } \quad \frac{\partial L}{\partial \mu}=0 \tag{91}
\end{equation*}
$$

## Reference

Simon, Carl P., and Lawrence Blume. Mathematics for economists. Vol. 7. New York: Norton, 1994

