October 31, 2014

Mathematics is very much like poetry .... What makes a good poem-a great poem-is that there is a large amount of thought expressed in a very few words. In this sense formula like

$$
e^{\pi i}+1=0 \quad \text { or } \quad \int_{-\infty}^{\infty} e^{-x^{2}} d x=\sqrt{\pi}
$$

are poems. ${ }^{1}$

- Lipman Bers

[^0]
## 1 Sets and Numbers

### 1.1 The concept of a set

The origin of the modern theory of sets can be traced back to the German mathematician Georg Cantor (1854-1918).
$\mathfrak{D e f i n i t i o n ~ ( S e t ) : ~}$
A set is any collection of well-defined and distinguishable objects. These objects are called elements, or members, of the set. Thus if $x$ is an element of a set $A$, then this fact is denoted by writing $x \in A$. If, however, $x$ is not an element of $A$, then we write $x \notin A$.

Curly brackets are usually used to describe the contents of a set. For example, if a set $A$ consists of the element $x_{1}, x_{2}, \ldots, x_{n}$, then it can be represented as $A=\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$. If the event membership in set is determined by satisfaction of a certain property or a relationship, then the description of the satisfaction can be given with the curly bracket. For example, if $A$ consists of all real numbers $x$ such that $x^{2}>1$, then it can be expressed as $A=\left\{x \mid x^{2}>1\right\}$, where the bar $\mid$ is used simply to mean "such that".
$\mathfrak{D e f i n i t i o n ~ ( E m p t y ~ S e t ) : ~}$
The set that contains no element is called the empty set and is denoted by $\varnothing$.
$\mathfrak{D e f i n i t i o n ~ ( S u b s e t ) : ~}$
A set $A$ is a subset of another set $B$, written symbolically as $A \subset B$, if every element of $A$ is an element of $B$. If $B$ contains at least one element that is not in $A$, then $A$ is said to be a proper subset of $B$.

## Definition:

A set $A$ and a set $B$ are equal if $A \subset B$ and $B \subset A$. Thus, every element of $A$ is an element of $B$ and vice versa.
$\mathfrak{D e f i n i t i o n ~ ( U n i v e r s a l ~ S e t ) : ~}$

The set that contains all set under consideration in a certain study is called the universal set and is denoted by $\Omega$.

### 1.2 Set Operations

There are two basic operations for set that produce new sets from existing ones. They are the operations of union and intersection.
$\mathfrak{D e f i n i t i o n ~ ( U n i o n ) : ~}$
The union of two sets $A$ and $B$, denoted by $A \cup B$, is the set of elements that belong to either $A$ or $B$, that is

$$
C=A \cup B=\{x \mid x \in A \text { or } x \in B\} .
$$

The definition can be extended to more than two sets. For example, if $A_{1}, A_{2}, \ldots, A_{n}$ are $n$ given sets, then their union, denoted by $\bigcup_{i=1}^{n} A_{i}$, is a set such that $x$ is an element of it if and only if $x$ belongs to at least one of the $A_{i}, i=1,2, \ldots, n$.
$\mathfrak{D e f i n i t i o n ~ ( I n t e r s e c t i o n ) : ~}$
The intersection of two sets $A$ and $B$, denoted by $A \cap B$, is the set of elements that belong to both $A$ and $B$, that is

$$
C=A \cap B=\{x \mid x \in A \text { and } x \in B\} .
$$

The definition can be extended to more than two sets. As before, if $A_{1}, A_{2}, \ldots, A_{n}$ are $n$ given sets, then their intersection, denoted by $\bigcap_{i=1}^{n} A_{i}$, is the set consisting all elements that belong to all the $A_{i}, i=1,2, \ldots, n$.

## $\mathfrak{D e f i n i t i o n ~ ( D i s j o i n t ) : ~}$

Two sets $A$ and $B$ are disjoint if their intersection is the empty set, that is $A \cap B=\varnothing$.
$\mathfrak{D e f i n i t i o n ~ ( C o m p l e m e n t ) : ~}$

The complement of a set $A$, denoted by $A^{c}$ (or $\bar{A}$ ), is the set consisting of all elements in the universal set that do not belong to $A$. In other words, $x \in A^{c}$ if and only if $x \notin A$.
$\mathfrak{D e f i n i t i o n ~ ( R e l a t i v e ~ C o m p l e m e n t ) : ~}$
The complement of $A$ with respect to s set $B$ is the set $B-A$ which consists of the element of $B$ that do not belong to $A$. This complement is called the relative complement of $A$ with respect to $B$.

From the definition above, the following results can be concluded:
$\mathfrak{R e s u l t s ~ 1 : ~ T h e ~ e m p t y ~ s e t ~} \varnothing$ is a subset of every set.
$\mathfrak{R e s u l t s ~ 2 : ~ T h e ~ e m p t y ~ s e t ~} \varnothing$ is unique.
$\mathfrak{R e s u l t s ~ 3 : ~ T h e ~ c o m p l e m e n t ~ o f ~} \varnothing$ is $\Omega$. Vice versa, the complement of $\Omega$ is $\varnothing$.
$\mathfrak{R e s u l t s ~ 4 : ~ T h e ~ c o m p l e m e n t ~ o f ~} A^{c}$ is $A$.
$\mathfrak{R e s u l t s ~ 5 : ~ F o r ~ a n y ~ s e t ~} A, A \cup A^{c}=\Omega$ and $A \cap A^{c}=\varnothing$.

Results 6: $A-B=A-(A \cap B)$.

Results 7: $A \cup(B \cup C)=(A \cup B) \cup C)$.
$\mathfrak{R e s u l t s ~ 8 : ~} A \cap(B \cap C)=(A \cap B) \cap C$.

Results 9: $A \cup(B \cap C)=(A \cup B) \cap(A \cup C)$.

Results 10: $A \cap(B \cup C)=(A \cap B) \cup(A \cap C)$.

Results ll: $(A \cup B)^{c}=A^{c} \cap B^{c} .{ }^{2}$ More generally, $\left(\bigcup_{i=1}^{n} A_{i}\right)^{c}=\bigcap_{i=1}^{n} A_{i}^{c}$.

[^1]Results 12: $(A \cap B)^{c}=A^{c} \cup B^{c} .{ }^{3}$ More generally, $\left(\bigcap_{i=1}^{n} A_{i}\right)^{c}=\bigcup_{i=1}^{n} A_{i}^{c}$.

Another useful set operation is the Cartesian product defined below.
$\mathfrak{D e f i n i t i o n ~ ( C a r t e s i a n ~ P r o d u c t ) : ~}$
Let $A$ and $B$ be two sets. Their Cartesian product, denoted by $A \times B$, is the set of all ordered pairs $(a, b)$ such that $a \in A$ and $b \in B$, that is,

$$
A \times B=\{(a, b) \mid a \in A \text { and } b \in B\} .
$$

## Example:

Let $A=\{1,2\}$ and $B=\{\alpha, \beta, \gamma\}$. Then

$$
A \times B=\{(1, \alpha),(1, \beta),(1, \gamma),(2, \alpha),(2, \beta),(2, \gamma)\}
$$

The preceding definition can be extended to more than two sets. If $A_{1}, A_{2}, \ldots, A_{n}$ are $n$ given sets, then their Cartesian product is denoted by

$$
\times_{i=1}^{n} A_{i}=\left\{\left(a_{1}, a_{2}, \ldots, a_{n}\right) \mid a_{i} \in A_{i}, i=1,2, \ldots, n\right\} .
$$

In particular, if the $A_{i}$ are equal to $A \forall i=1,2, \ldots, n$, the one writes $A^{n}$ for $\times{ }_{i=1}^{n} A$.

### 1.3 Class of Subsets

## $\mathfrak{D e f i n i t i o n ~ ( P o w e r ~ S e t ) : ~}$

The set of all the subsets of $A$ is called the power set of $A$, denoted $2^{A}$. The power set of a set with $n$ elements has $2^{n}$ elements, which accounts for its name.

When studying the subsets of a given set, particularly their measure-theoretic properties, the power set is often too big for anything very interesting or useful to said about it. The idea behind the following definitions is to specify subset of $2^{A}$ that are large enough to be interesting, but whose characteristics may be more tractable. We typically do this by choosing a base collection of sets with known properties, and

[^2]then specifying certain operations for creating new sets from existing ones. These operations permit an interesting diversity of class members to be generated, but important properties of the sets may be deduced from those of the base collection.
$\mathfrak{D e f i n i t i o n ~ ( S i g m a ~ F i e l d s ) : ~}$
A $\sigma$-field ( $\sigma$-algebra) $\mathcal{F}$ is a class of subsets of $A$ satisfying
(a). $A$ and $\varnothing \in \mathcal{F}$.
(b). If $E \in \mathcal{F}$ then $E^{c}$ (the complements of $E$ in $A$ ) $\in \mathcal{F}$.
(c). If $\left\{E_{n}, n \in \mathbb{N}\right\}$ is a sequence of $\mathcal{F}$-sets, then $\bigcup_{n=1}^{\infty} E_{n} \in \mathcal{F}$.
$\mathfrak{D e f i n i t i o n ~ ( B o r e l ~ F i e l d ) : ~}$
The Borel field of $\mathbb{R}$, denoted by, $\mathcal{B}$, is the smallest $\sigma$-field of $\mathbb{R}$ that contains all the half-lines, the set of the form $(-\infty, x], \forall x \in \mathbb{R}$.

Applying the definition of a $\sigma$-field and the rule of set algebra, consider what sorts of sets $\mathcal{B}$ contains. By the de Morgan's laws it must contain the half-open intervals $\left(x_{1}, x_{2}\right.$ ] for $x_{1}<x_{2}$ (intersections of half-lines). It also contains the open half lines, set of the form $(-\infty, x)$ for $x \in \mathbb{R}$ because

$$
(-\infty, x)=\bigcup_{n=1}^{\infty}\left(-\infty, x-\frac{1}{n}\right]
$$

where all the sets of the countable union are in $\mathcal{B}$. It therefore contains all open intervals, and also the singleton sets

$$
(-\infty,] \cap(-\infty, x)^{c}=\{x\}
$$

for any $x \in \mathbb{R}$. Also, any sets that can be formed from finite or countable infinity union, intersection and complements of theses sets. This is a rich enough collection for our needs to assign probability.

### 1.4 The Real Numbers

Let us first look at the particularly interested set-the real number set. The first questions may be: How the real numbers are created? They are from the following procedures:

Natural Number $(\mathbb{N}){ }^{\text {by addition, subtraction, multiplication, division }}$ Rationales,
Rationales + Irrationals ${ }^{\text {by Dedekind's }} \xrightarrow{\text { Theorem(Completeness) }}$ Real Number $(\mathbb{R})$.

### 1.4.1 Geometry and the Number System

Geometrical language, with its highly suggestive power, can be very useful and can be given arithmetical meaning. We proceed to define a number of geometrical terms.
$\mathfrak{D e f i n i t i o n ~ ( O n e - D i m e n s i o n a l ~ E u c l i d e a n ~ S p a c e ) : ~}$
We speak of the real-number system as one-dimensional space, and of course we visualize it as a line. We denote it as $\mathbb{R}$. In one-dimensional space we use the word point to mean "number". A point set is a collection of points.
$\mathfrak{D e f i n i t i o n ~ ( T w o - D i m e n s i o n a l ~ E u c l i d e a n ~ S p a c e ) : ~}$
Just as one-dimensional space is the collection of all real numbers, the two-dimensional space is the Cartesian product of two one-dimensional space, $\mathbb{R} \times \mathbb{R}=\mathbb{R}^{2}$, i.e. the ordered pairs $(x, y)$ of real numbers, which is also known as the 2 dimensional Euclidean space. These ordered pairs are called points in the two-dimensional space.
$\mathfrak{D e f i n i t i o n ~ ( n - D i m e n s i o n a l ~ E u c l i d e a n ~ S p a c e ) : ~}$
The Cartesian product $\times_{i=1}^{n} \mathbb{R}$ is denoted by $\mathbb{R}^{n}$, which is known as the $n$-dimensional Euclidean space.

### 1.5 Relations and Functions

Let $A \times B$ be the Cartesian product of two sets, $A$ and $B$.
$\mathfrak{D e f i n i t i o n ~ ( R e l a t i o n s ) : ~}$
A relations $\rho$ from $A$ to $B$ is a subset of $A \times B$, that is $\rho$ consists of ordered pairs $(a, b)$ such that $a \in A$ and $b \in B$. In particular, if $A=B$, Then $\rho$ is said to be a relation in $A$. Whenever $\rho$ is a relation and $(x, y) \in \rho$, then $x$ and $y$ are said to be
$\rho$-related. This is denoted by writing $x \rho y$.

## Example:

If $A=\{7,8,9\}$ and $B=\{7,8,9,10\}$, the $\rho=\{(a, b) \mid a<b, a \in A, b \in B\}$ is a relation from $A$ to $B$ that consists of the six ordered pairs $(7,8),(7,9),(7,10),(8,9)$, $(8,10),(9,10)$.

Definition: (Functions):
Let $\rho$ be a relation from $A$ to $B$. Suppose that $\rho$ has the property that for all $x$ in $A$, if $x \rho y$ and $x \rho z$, where $y$ and $z$ are elements in $B$, then $y=z$. Such a relation is called a function.

Thus a function is a relation $\rho$ such that any two elements in $B$ that are $\rho$-related to the same $x$ in $A$ must be identical. In other words, to each element $x$ in $A$, there corresponds only one element $y$ in $B$. We call $y$ the value of the function at $x$ and denote it by writing $y=f(x)$. The set $A$ is called the domain of the function $f$, and the set of all values of $f(x)$ for $x$ in $A$ is called the range of $f$, or the image of $A$ under $f$, and is denoted by $f(A)$. In this case, we say that $f$ is a function, or a mapping, from $A$ to $B$. We express this fact by writing

$$
f: A \mapsto B
$$

Note that $f(A)$ is a subset of $B(f(A) \subset B)$. In particular, if $B=f(A)$, then $f$ is said to be a function from $A$ onto $B$. In this case, every element $b$ in $B$ has a corresponding element $a$ in $A$ such that $b=f(a)$.
$\mathfrak{D e f i n i t i o n ~ ( o n e - t o - o n e ~ f u n c t i o n ) : ~}$
A function $f$ defined on a set $A$ is said to be a one-to-one function if whenever $f\left(x_{1}\right)=f\left(x_{2}\right)$ for $x_{1}, x_{2}$ in $A$, one has $x_{1}=x_{2}$. Equivalently, $f$ is a one-to-one function if whenever $x_{1} \neq x_{2}$, one has $f\left(x_{1}\right) \neq f\left(x_{2}\right) .{ }^{4}$ Thus a function $f: A \rightarrow B$ is one-to-one if to each $y$ in $f(A)$, there corresponds only one element $x$ in $A$ such that $y=f(x)$.
$\mathfrak{D e f i n i t i o n ~ ( o n e - t o - o n e ~ c o r r e s p o n d e n c e ) : ~}$
If $f$ is a one-to-one and onto function, then it is said to provide a one-to-one corre-

[^3]spondence between $A$ and $B$. In this case, the set $A$ and $B$ are said to be equivalent, denoted by $A \mapsto B$.
$\mathfrak{D e f i n i t i o n : ~ ( I n v e r s e ~ F u n c t i o n ) : ~}$
Whenever $A \mapsto B$, there is a function $g: B \rightarrow A$ such that if $y=f(x)$, then $x=g(y)$. The function $g$ is called the inverse function of $f$ and is denoted by $f^{-1}$.
$\mathfrak{D e f i n i t i o n ~ ( C o m p o s i t e ~ F u n c t i o n ) : ~}$
Let $f: A \rightarrow B$ and $h: B \rightarrow C$ be one-to-one and onto functions. Then, the composite function $h \circ f=h[f(x)]$, defines a one-to-one correspondence between $A$ and $C$.

## Example:

The relation $a \rho b$, where $a$ and $b$ are real numbers such that $a=b^{2}$, is not a function. This is true because both pairs $(a, b)$ and $(a,-b)$ belong to $\rho .{ }^{5}$

## Example:

The relation $a \rho b$, where $a$ and $b$ are real numbers such that $a^{2}=b$, is a function. Since for each $a$, there is only one $b$ that is $\rho$-related to $a .^{6}$ However, this is not a one-to-one function since there are two elements in $A$, (i.e. $a$ and $-a$ ) that are $\rho$-related to a given $b$.

### 1.6 Finite, Countable and Uncountable Set

Let $J_{n}=\{1,2, \ldots, n\}$ be s set consisting of the first $n$ positive integers, and let $J^{+}$ denote the set of all positive integers.
$\mathfrak{D e f i n i t i o n ~ ( F i n i t e , ~ C o u n t a b l e ~ a n d ~ U n c o u n t a b l e ~ S e t ) : ~}$
A set $A$ is said to be:
(a). Finite if $A \mapsto J_{n}$ for some positive integer $n$.

[^4](b). Countable if $A \mapsto J^{+}$. In this case, the set $J^{+}$can be used as an index set for $A$, that is, the elements of $A$ are assigned distinct indices (subscripts) that belong to $J^{+}$. Hence, $A$ can be represented as $A=\left\{a_{1}, a_{2}, \ldots, a_{n}, \ldots\right\}$.
(c). Uncountable if $A$ is neither finite nor countable. In this case, the elements of $A$ cannot be indexed by $J_{n}$ for any $n$, or by $J^{+}$.

## Example:

Let $A=\left\{1,4,9, \ldots, n^{2}, \ldots\right\}$. The set is countable, since the function $f: J^{+} \rightarrow A$ defined by $f(n)=n^{2}$ is one-to-one and onto. Hence, $A \mapsto J^{+}$.

## $\mathfrak{E x a m p l e}$ :

Let $A=J$ be the set of all integers. Then $A$ is countable. To show this, consider the function $f: J^{+} \rightarrow A$ defined by

$$
f(n)=\left\{\begin{array}{lc}
(n+1) / 2, & n \text { odd } \\
(2-n) / 2, & n \text { even }
\end{array}\right.
$$

It can be verified that $f$ is one-to-one and onto. Hence, $A \mapsto J^{+}$.

## Example:

Let $A=\{x \mid 0 \leq x \leq 1\}$. The set is uncountable.

This result implies that any subset of $\mathbb{R}$, the set of real numbers, must be uncountable. In particular, $\mathbb{R}$ is uncountable.

## Theorem:

The set $Q$ of all rational number is countable.

### 1.7 Bounded Set

Let us consider the set $\mathbb{R}$ of real numbers.
$\mathfrak{D e f i n i t i o n ~ ( B o u n d e d ~ S e t ) : ~}$
A set $A \subset \mathbb{R}$ is said to be:
(a). Bounded from above if there exists a number $q$ such that $x \leq q, \forall x \in A$. This number is called an upper bound of $A$.
(b). Bounded from below if there exists a number $p$ such that $x \geq q, \forall x \in A$. This number is called an lower bound of $A$.
(c). Bounded if $A$ has an upper bounded $q$ and a lower bounded $p$. In this case, there exist a nonnegative number $r$ such that $-r \leq x \leq r, \forall x \in A$. This number is equal to $\max (|p|,|q|)$.

## $\mathfrak{D e f i n i t i o n ~ ( S u p r e m u m ) : ~}$

Let $A \subset \mathbb{R}$ be a set bounded from above. If there exists a number $l$ that is an upper bounded of $A$ and is less than or equal to any other upper bounded of $A$, then $l$ is called the least upper bound of $A$ and is denoted by $l u b(A)$. Another name for $\operatorname{lub}(A)$ is the supremum of $A$ and is denoted by $\sup _{A} x .^{7}$
$\mathfrak{D e f i n i t i o n ~ ( I n f i m u m ) : ~}$
Let $A \subset \mathbb{R}$ be a set bounded from below. If there exists a number $g$ that is an lower bounded of $A$ and is greater than or equal to any other lower bounded of $A$, then $g$ is called the greatest lower bound of $A$ and is denoted by $g l b(A)$. Another name for $g l b(A)$ is the infimum of $A$ and is denoted by $\inf _{A} x$.

## Theorem:

Let $A \subset \mathbb{R}$ be a non-empty set.
(a). If $A$ is bounded from above, then $\sup _{A} x$ exists.
(b). If $A$ is bounded from below, then $\inf _{A} x$ exists.

### 1.8 The Topology of the Real Line

The purpose of this section is to treat rigorously the idea of 'nearness', as it applies to points of the line. The key intergradient of the theory is the distance between a pair of points $x, y \in \mathbb{R}$.
$\mathfrak{D e f i n i t i o n ~ ( E u c l i d e a n ~ D i s t a n c e ) : ~}$

[^5]The Euclidean distance of $\mathbb{R}$ is $|x-y|$, for $x, y \in \mathbb{R}$.
$\mathfrak{D e f i n i t i o n ~ ( ~} \epsilon$-Neighborhood):
An $\epsilon$-neighborhood of a point $x \in \mathbb{R}$ is a set $N_{\epsilon}(x)=\{y:|x-y|<\epsilon\}$, for some $\epsilon>0$.
$\mathfrak{D e f i n i t i o n ~ ( D e l e t e d ~} \epsilon$-Neighborhood):
An deleted $\epsilon$-neighborhood of a point $x \in \mathbb{R}$ is a set $N_{\epsilon}^{d}(x)=\{y:|x-y|<\epsilon\}$, for some $\epsilon>0$, but $y \neq x$.
$\mathfrak{D e f i n i t i o n}$ (Open Set):
An open set is a set $A \subseteq \mathbb{R}$ such that for each $x \in A$, there exists for some $\epsilon>0$ an $\epsilon$-neighborhood which is a subset of $A$.
$\mathfrak{D e f i n i t i o n ~ ( C l o s e d ~ S e t ) : ~}$
The complement of an open set in $\mathbb{R}$ is a closed set.
$\mathfrak{D e f i n i t i o n ~ ( L i m i t ~ P o i n t ) : ~}$
A limit point of a set $A$, denoted by $\operatorname{Lim}(A)$ is a point $p \in \mathbb{R}$ such that, for every $\epsilon>0$, the set $A \cap N_{\epsilon}^{d}(p)$ is not empty. The limit points of $A$ are not necessarily elements of $A$, open set being a case in point.

Limit points can be used to describe closed sets, as can be seen from the following theorem.

## $\mathfrak{T h e o r e m}$ :

A set $B$ is closed if and only if every limit point of $B$ belongs to $B$.

## $\mathfrak{E x a m p l e}$ :

$A=\{x \mid 0<x<1\}$ is an open subset of $\mathbb{R}$, but is not closed, since both 0 and 1 are limit points of $B$, but do not belong to it.

## Example:

$A=\{x \mid 0 \leq x \leq 1\}$ is closed, but is not open, since any neighborhood of 0 or 1 is not contained in $B$.

## Example:

$A=\{x \mid 0<x \leq 1\}$ is not open, because any neighborhood of 1 is not contained in $B$. It is also not closed, because 0 is a limit point that does not belong to $B$.
$\mathfrak{D e f i n i t i o n ~ ( I n t e r i o r ) : ~}$
A point $x_{0}$ in $\mathbb{R}$ is an interior of a set $A \subset \mathbb{R}$ if there exists an $r>0$ such that $N_{r}\left(x_{0}\right) \subset A$ and is denoted by $\operatorname{Int}(A)$. Thus $A$ is open if it contains entirely of interior points.
$\mathfrak{D e f i n i t i o n ~ ( B o u n d a r y ~ P o i n t ) : ~}$
A point $p \in \mathbb{R}$ is a boundary point of a set $A \subset \mathbb{R}$ if every neighborhood of $p$ contains points of $A$ as well as points of $A^{c}$, the complement of $A$ with respect to $\mathbb{R}$. The set of all boundary points of $A$ is called its boundary and is denoted by $\operatorname{Br}(A)$. Thus it is easy to see that $\operatorname{Br}(A)=\operatorname{Lim}(A)-\operatorname{Int}(A)$.

## Example:

For an open interval $(a, b)$. Every points of $(a, b)$ is a limit point, and $a$ and $b$ are also limit points not belong to $(a, b)$. They are boundary points of both $(a, b)$ and $[a, b]$.
$\mathfrak{D e f i n i t i o n ~ ( C o v e r i n g ) : ~}$
A collection of set $\left\{B_{\alpha}\right\}$ is said to be a covering of a set $A$ if the union $\cup_{\alpha} B_{\alpha}$ contains $A$. If each $B_{\alpha}$ is an open set, then $\left\{B_{\alpha}\right\}$ is called an open covering.

## $\mathfrak{D e f i n i t i o n ~ ( C o m p a c t ) : ~}$

A set $A$ is compact if each open covering $\left\{B_{\alpha}\right\}$ of $A$ has a finite sub-covering, that is, there is a finite sub-collection $B_{\alpha_{1}}, B_{\alpha_{2}}, \ldots, B_{\alpha_{n}}$ of $\left\{B_{\alpha}\right\}$ such that $A \subset \cup_{i=1}^{n} B_{\alpha_{i}}$.

The concept of compactness is motivated by the classical Heine-Borel theorem, which characterizes compact sets in $\mathbb{R}$ as closed and bounded sets.

Theorem (Heine-Borel):
A set $B \subset \mathbb{R}$ is compact if and only if it is closed and bounded.

Thus, according to the Heine-Borel theorem, every closed and bounded interval $[a, b]$ is compact.

## $\mathfrak{E x a m p l e}$ :

Let us examine the open interval $(0,1)$. Consider the collection $B_{\alpha}=\left\{\left(\frac{1}{i}, 1\right), i=\right.$ $1,2, \ldots\}$ and observe that

$$
(0,1)=\left(\frac{1}{2}, 1\right) \cup\left(\frac{1}{3}, 1\right) \cup \cdots,
$$

that is, $B_{\alpha}$ us an open cover of $(0,1)$. Does $B_{\alpha}$ have a finite subset that covers $(0,1)$ ? No, because the greatest lower bound of any finite subset of $B_{\alpha}$ is bounded away from 0 , so no such subset can possibly cover $(0,1)$ entirely. There, we conclude that $(0,1)$ is not a compact subset of $\mathbb{R}$.

## 2 Measure

A measure is a set function, a mapping which associates a real number with a set. Commonplace examples of measure include the lengths, areas, and volumes of geometrical figures, but wholly abstract sets can be 'measured' in an analogous way. Formally, we have the following definition.

## Definition:

Given a class $\mathcal{F}$ of subsets of a set $\Omega$, a measure

$$
\mu: \mathcal{F} \mapsto \mathbb{R}
$$

is a function having the following properties:
(a). $\mu(A) \geq 0, \forall A \in \mathcal{F}$.
(b). $\mu(\varnothing)=0$.
(c). For a countable collection $\left\{A_{j} \in \mathcal{F}, j \in \mathbb{N}\right\}$ with $A_{i} \cap A_{j}=\varnothing$ for $i \neq j$ and $\bigcup_{j} A_{j} \in \mathcal{F}$,

$$
\mu\left(\bigcup_{j} A_{j}\right)=\sum_{j} \mu\left(A_{j}\right) . \quad \text { (Countable Additivity) }
$$

The particular cases at issue in this course are of course the probabilities of random events is a sample space $\Omega$.
$\mathfrak{D e f i n i t i o n ~ ( M e a s u r a b l e ~ S p a c e ) : ~}$
A measurable space is a pair $(\Omega, \mathcal{F})$ where $\Omega$ is any collection of objects, and $\mathcal{F}$ is a $\sigma$-field of subset of $\Omega$.

## $\mathfrak{D e f i n i t i o n ~ ( M e a s u r e ~ S p a c e ) : ~}$

When $(\Omega, \mathcal{F})$ is a measurable space, the triple $(\Omega, \mathcal{F}, \mu)$ is called a measure space. More than one measure can be associated with the measurable space $(\Omega, \mathcal{F})$, hence the distinction between measure space and measurable space is important.

## Example:

The case closest to everyday intuition if Lebesgue measure, $m$, on the measurable space $(\mathbb{R}, \mathcal{B})$, where $\mathcal{B}$ is the Borel field on $\mathbb{R}$. Generalizing the notion of length in geometry, Lebesgue measure assigns $m((a, b])=b-a$ to an interval ( $a, b]$. Additivity is an intuitively plausible property if we think of measuring the total length of a collection of disjoint intervals.

Some additional properties may be deduced from the definition.

## Theorem:

For arbitrary $\mathcal{F}$-set $A, B$, and $\left\{A_{j}, j \in \mathbb{N}\right\}$,
(a). If $A \subseteq B$ then $\mu(A) \leq \mu(B) \quad$ (monotonicity).
(b). $\mu(A \cup B)+\mu(A \cap B)=\mu(A)+\mu(B)$.
(c). $\mu\left(\cup_{j} A_{j}\right) \leq \sum_{j} \mu\left(A_{j}\right) . \quad$ (countable subadditivity).

## 3 Metric Space

Central to the properties of $\mathbb{R}$ was the concept of distance. For any real number $x$ and $y$, the Euclidean distance between them is the number $d_{E}(x, y)=|x-y| \in \backslash^{+}$. Generalizing this idea, a set (otherwise arbitrary) having a distance measure, or metric, defined for each pair of elements is called a metric space. Let $\mathbb{S}$ denote such a set.
$\mathfrak{D e g f i n i t i o n ~ ( M e t r i c ~ a n d ~ M e t r i c ~ S p a c e ) : ~}$
A metric is a mapping $d: \mathbb{S} \times \mathbb{S} \mapsto \mathbb{R}^{+}$having the properties
(a). $d(x, y)=d(y, x)$,
(b). $d(x, y)=0$ iff $x=y$,
(c). $d(x, y)+d(y, z) \geq d(x, z) \quad$ (triangle inequality),

A metric space $(\mathbb{S}, d)$ is a set $\mathbb{S}$ paired with metric $d$, such that conditions (a)-(c) hold for each pair of elements of $\mathbb{S}$.

While the Euclidean metric on $\mathbb{R}$ is the familiar case, and the proof that $d_{E}$ satisfies (a)-(c) is elementary, $d_{E}$ is not the only possible metric on $\mathbb{R}$.

In the space $\mathbb{R}^{2}$ a larger variety of metric is found.

## Example:

The Euclidean distance on $\mathbb{R}^{2}$ is

$$
d_{E}(x, y)=\|x-y\|=\left[\left(x_{1}-y_{1}\right)^{2}+\left(x_{2}-y_{2}\right)^{2}\right]^{1 / 2}
$$

and $\left(\mathbb{R}^{2}, d_{E}\right)$ is the Euclidean plane. An alternative is the 'taxicab' (honest) metric,

$$
d_{T}=\left|x_{1}-y_{1}\right|+\left|x_{2}-y_{2}\right| .
$$

$d_{E}$ is the shortest distance between two address in a city, but $d_{T}$ is the shortest distance by taxi.

In metric space theory, the properties of $\mathbb{R}$ are revealed as a special case.
$\mathfrak{D e f i n i t i o n ~ ( ~} \epsilon$-Neighborhood):
The concept of an open neighborhood in a metric space $(\mathbb{S}, d)$ is the set $N_{\epsilon}(x, d)=$
$\{y: y \in \mathbb{S}, d(x, y)<\epsilon\}$, where $x \in \mathbb{S}$ and $\epsilon>0$.
$\mathfrak{D e f i n i t i o n ~ ( O p e n ~ S e t ) : ~}$
An open set of $(\mathbb{S}, d)$ is a set $A \subseteq \mathbb{S}$ such that for each $x \in A, \exists \delta>0$ such that $N_{\delta}(x, d)$ is a subset of $A$.

## 4 Limit and Continuity of Real Functions

The notation of limits and continuity of functions lie at the kernel of calculus. In this section, we review the concepts of limits and continuity of real-valued functions, and study some of their properties. The domain of the functions will be subsets of R. A

### 4.1 Limits of a Functions

### 4.1.1 What is mean $x \rightarrow a$

Before defining the notation of a limit of a function, let us understand what is meant by the natation $x \rightarrow a$, where $a$ and $x$ are elements in $\mathbb{R}$.
$\mathfrak{D e f i n i t i o n ~}(x \rightarrow a$ when $a$ is finite):
If $a$ is finite, then $x \rightarrow a$ means that $x$ can have value that belong to a neighborhood $N_{r}(a)$ of $a$ for any $r>0$, but $x \neq a$, that is, $0<|x-a|<r$. Such a neighborhood is called a deleted neighborhood of $a$, that is, a neighborhood from which the point $a$ has been removed.
$\mathfrak{D e f i n i t i o n ~}(x \rightarrow a$ when $a$ is infinite):
If $a$ is infinite $(+\infty$ or $-\infty)$, then $x \rightarrow a$ indicates that $|x|$ can get larger and larger without any constraint on the extent of its increase.

Let us now study the behavior of a function $f(x)$ as $x \rightarrow a$.
$\mathfrak{D e f i n i t i o n ~}(f(x)$ is finite and $a$ is finite):
Suppose that the function $f(x)$ is defined in a deleted neighborhood of a point $a \in \mathbb{R}$. Then $f(x)$ is said to have a limit $L$ as $x \rightarrow a$ if for every $\epsilon>0$ there exists a $\delta>0$ such that

$$
|f(x)-L|<\epsilon
$$

for all $x$ for which

$$
0<|x-a|<\delta
$$

In this case, we write $f(x) \rightarrow L$ as $x \rightarrow a$, which is equivalent to saying that $\lim _{x \rightarrow a} f(x)=L$. Less formally, we say that $f(x) \rightarrow L$ as $x \rightarrow a$ if, however small the positive number $\epsilon$ might be, $f(x)$ differs from $L$ by less than $\epsilon$ for value of $x$ sufficiently close to $a .^{8}$

## $\mathfrak{E x a m p l e}$ :

Use the $\epsilon-\delta$ definition of limit to prove that

$$
\lim _{x \rightarrow 2}(3 x-2)=4
$$

## Solution:

You must show that for each $\epsilon>0$, there exists a $\delta>0$ such that $|(3 x-2)-4|<\epsilon$ whenever $0<|x-2|<\delta$. Because your choice of $\delta$ depends on $\epsilon$, you need to establish a connection between the absolute value $|(3 x-2)-4|$ and $|x-2|$.

$$
|(3 x-2)-4|=|3 x-6|=3|x-2| .
$$

So, for a given $\epsilon>0$ you can choose $\delta=\epsilon / 3$. This choice works because

$$
0<|x-2|<\delta=\frac{\epsilon}{3}
$$

implies that

$$
|(3 x-2)-4|=3|x-2|<3\left(\frac{\epsilon}{3}\right)=\epsilon
$$

[^6]$\mathfrak{D e f i n i t i o n ~}(f(x)$ is infinite and $a$ is finite):
For every positive number $M$ there exists a $\delta>0$ such that
$$
|f(x)|>M
$$
for all $x$ for which
$$
0<|x-a|<\delta
$$

## Example:

Discuss the limit of $\lim _{x \rightarrow 0} \frac{1}{x^{2}}$.

## Solition:

Let $f(x)=1 / x^{2}$, we can see that as $x$ approach 0 from either the right or the left, $f(x)$ increase without bound. For instance,

$$
0<|x-0|<\frac{1}{10} \Longrightarrow f(x)=\frac{1}{x^{2}}>100
$$

Similarly, you can force $f(x)$ to be larger than 1000000 as follows.

$$
0<|x-0|<\frac{1}{1000} \Longrightarrow f(x)=\frac{1}{x^{2}}>100000
$$

$\mathfrak{D e f i n i t i o n ~}\left(f(x)\right.$ is finite and $a$ is infinite): ${ }^{9}$ If $a$ is infinite and $L$ is finite, then $f(x) \rightarrow L$ as $x \rightarrow a$ if for any $\epsilon>0$ there exists a positive number $N$ such that

$$
|f(x)-L|<\epsilon
$$

for all $x$ for which

$$
|x|>N .
$$

[^7]$\mathfrak{D e f i n i t i o n ~}(f(x)$ is infinite and $a$ is infinite):
If both $a$ and $L$ are infinite, then $f(x) \rightarrow L$ as $x \rightarrow a$ if for any $B>0$ there exists a positive number $A$ such that
$$
|f(x)|>B
$$
if
$$
|x|>A .
$$

### 4.1.2 One-Sided Limit

The limit of $f(x)$ as described above is actually called a two-sided limit. This is because $x$ can approach $a$ from either side. There are, however, cases where $f(x)$ can have a limit only when $x$ approach $a$ from one side. Such a limit is called a one-sided limit.
$\mathfrak{D e f i n i t i o n ~ ( L e f t - s i d e d ~ l i m i t ) : ~}$
If $f(x)$ has a limit as $x$ approach $a$ from the left, symbolically written as $x \rightarrow a^{-}$, then $f(x)$ has a left-sided limit, which we denote by $L^{-}$. In this case we write

$$
\lim _{x \rightarrow a^{-}} f(x)=L^{-}
$$

It follows that $f(x)$ has a left-sided limit $L^{-}$as $x \rightarrow a^{-}$if for every $\epsilon>0$ there exists a $\delta>0$ such that

$$
\left|f(x)-L^{-}\right|<\epsilon
$$

for all $x$ for which $-\delta<x-a<0$.
$\mathfrak{D e f i n i t i o n ~ ( R i g h t - s i d e d ~ l i m i t ) : ~}$
If $f(x)$ has a limit as $x$ approach $a$ from the right, symbolically written as $x \rightarrow a^{+}$, then $f(x)$ has a right-sided limit, which we denote by $L^{+}$. In this case we write

$$
\lim _{x \rightarrow a^{+}} f(x)=L^{+} .
$$

It follows that $f(x)$ has a right-sided limit $L^{+}$as $x \rightarrow a^{+}$if for every $\epsilon>0$ there exists a $\delta>0$ such that

$$
\left|f(x)-L^{+}\right|<\epsilon
$$

for all $x$ for which $0<x-a<\delta$.
$\mathfrak{T h e o r e m ~ ( T h e ~ E x i s t e n c e ~ o f ~ a ~ L i m i t ) : ~}$
A necessary and sufficient condition that $\lim _{x \rightarrow a} f(x)$ exist is that both $\lim _{x \rightarrow a^{-}} f(x)$ and $\lim _{x \rightarrow a^{+}} f(x)$ exist and be equal. ${ }^{10}$

## Example:

Find the limit of the smallest integer function $f(x)=[x]$ as $x$ approach 0 from the left and from the right.

## Solution:

The limit as $x$ approach 0 from the left is given by

$$
\lim _{x \rightarrow 0^{-}}[x]=-1
$$

[^8]Similarly, there exist a $\delta_{2}$ for which

$$
|f(x)-L|<\epsilon \quad \text { if } a-\delta_{2}<x<a
$$

If we let $\delta=\min \left(\delta_{1}, \delta_{2}\right)$, then

$$
|f(x)-L|<\epsilon \quad \text { if } a-\delta<x<a+\delta
$$

which means that $0<|x-a|<\delta$.
To prove the necessity of the condition we show that if $\lim _{x \rightarrow a} f(x)=L$, then $\lim _{x \rightarrow a^{+}} f(x)=L$.
A similar proof will verify the condition for left-hand limits.
By definition, given $\epsilon>0$ we must find a $\delta>0$ for which

$$
|f(x)-L|<\epsilon \quad \text { if } a<x<a+\delta
$$

For the given $\epsilon$, there is a $\delta_{1}$ for which

$$
|f(x)-L|<\epsilon \quad \text { if } a-\delta_{1}<x<a+\delta_{1}
$$

Hence, merely let $\delta=\delta_{1}$.
and the limit as $x$ approach 0 from the right is given by

$$
\lim _{x \rightarrow 0^{+}}[x]=0
$$

The smallest integer function $\lim _{x \rightarrow 0}[x]$ does not exist. By similar reasoning, you can see that the smallest integer function do not have a limit at any integer $n$.

### 4.2 Some Properties Associated with Limits of Functions

The following theorems give some fundamental properties associated with function limits.

## Theorem:

Let $f(x)$ and $g(x)$ be real-valued functions defined on $D \subset \mathbb{R}$. Suppose that $\lim _{x \rightarrow a} f(x)=L$ and $\lim _{x \rightarrow a} g(x)=M$. Then
(a). $\lim _{x \rightarrow a}[f(x)+g(x)]=L+M$,
(b). $\lim _{x \rightarrow a}[f(x) g(x)]=L M$,
(c). $\lim _{x \rightarrow a}[1 / g(x)]=1 / M$ if $M \neq 0$,
(d). $\lim _{x \rightarrow a}[f(x) / g(x)]=L / M$ if $M \neq 0 .{ }^{11}$

## Theorem:

If $f(x) \leq g(x), \forall x \in D \subset \mathbb{R}$, then $\lim _{x \rightarrow a} f(x) \leq \lim _{x \rightarrow a} g(x)$.

### 4.3 The $o$ and $O$ Notation

These symbols provide a convenient way to describe the limiting behavior of a function $f(x)$ as $x$ tends to a certain limit.

[^9]if $0<|x-a|<\lambda$, where $\lambda=\min \left(\lambda_{1}, \lambda_{2}\right)$.
$\mathfrak{D e f i n i t i o n}(\operatorname{Big} O)$ :
Let $f(x)$ and $g(x)$ be two functions defined on $D \subset \mathbb{R}$. The function $g(x)$ is positive and usually has a simple form such as $1, x$ or $1 / x$. Suppose there exist a positive number $K$ such that
$$
\frac{|f(x)|}{g(x)} \leq K
$$
for all $x \in E$, where $E \subset D$. Then, $f(x)$ is said to be of an order of magnitude not exceeding that of $g(x)$. This fact is denoted by writing
$$
f(x)=O(g(x))
$$
for all $x \in E$. In particular, if $g(x)=1$, then $f(x)$ is necessarily a bounded function on $E$.

## Example:

$$
\begin{aligned}
\cos (x) & =O(1) \quad \text { for all } x \\
x & =O\left(x^{2}\right) \quad \text { for large values of } x \\
x^{2}+x & =O\left(x^{2}\right) \quad \text { for all } x
\end{aligned}
$$

$\mathfrak{D e f i n i t i o n ~ ( S m a l l ~} o$, as $x \rightarrow a$ (finite)):
Suppose that the relationship between $f(x)$ and $g(x)$ is such that

$$
\lim _{x \rightarrow a} \frac{f(x)}{g(x)}=0
$$

Then we say that $f(x)$ is of smaller order of magnitude than $g(x)$ in a deleted neighborhood of $a$. This fact is denoted by writing

$$
f(x)=o(g(x)) \quad \text { as } x \rightarrow a
$$

which is equivalent to saying that $f(x)$ tends to zero more rapidly than $g(x)$ as $x \rightarrow a$.
$\mathfrak{D e f i n i t i o n ~ ( S m a l l ~} o$, as $x \rightarrow \infty$ ):
The o symbol can also be used when $x$ tends to infinity. In this case we write

$$
f(x)=o(g(x)) \quad \text { for } x>A,
$$

where $A$ is some positive number.

## $\mathfrak{E x a m p l e}$ :

$$
\begin{aligned}
x^{2} & =o(x) \quad \text { as } x \rightarrow 0 \\
\sqrt{x} & =o(x) \quad \text { as } x \rightarrow \infty
\end{aligned}
$$

$\mathfrak{D e f i n i t i o n ~ ( A s y m p t o t i c a l l y ~ E q u a l ) : ~}$
If $f(x)$ and $g(x)$ be any two functions such that ${ }^{13}$

$$
\lim _{x \rightarrow a} \frac{f(x)}{g(x)}=1
$$

then $f(x)$ and $g(x)$ are said to be asymptotically equal, written symbolically $f(x) \sim$ $g(x)$, as $x \rightarrow a$.

## $\mathfrak{E x a m p l e}$ :

$$
\begin{aligned}
x^{2} & \sim x^{2}+3 x+1 \text { as } x \rightarrow \infty \\
\sin x & \sim x \text { as } x \rightarrow 0
\end{aligned}
$$

One the basis of the above definitions, the following properties can be deduced:
(a). $O(f(x)+g(x))=O(f(x))+O(g(x)) \cdot{ }^{14}$
${ }^{13}$ Or write as $\frac{f(x)}{g(x)} \rightarrow 1$ as $x \rightarrow a$.
${ }^{14} \mathfrak{P r o o f}$ :
Let $h_{1}(x)=O(f(x))$, then $\left|h_{1}(x)\right| \leq K f(x)$. Similarly $\left|h_{2}(x)\right| \leq K g(x)$. Therefore

$$
\left|h_{1}(x)+h_{2}(x)\right| \leq\left|h_{1}(x)\right|+\left|h_{2}(x)\right| \leq K(f(x)+g(x)),
$$

i.e. $h_{1}(x)+h_{2}(x)$ is $O((f(x)+g(x))$.
(b). $O(f(x) g(x))=O(f(x)) O(g(x)) \cdot{ }^{15}$
(c). $o(f(x) g(x))=O(f(x)) o(g(x))$.
(d). If $f(x) \sim g(x)$ as $x \rightarrow a$, then $f(x)=g(x)+o(g(x))$ as $x \rightarrow a$.

### 4.4 Continuous Functions

A function $f(x)$ may has a limit $L$ as $x \rightarrow a$. This limit, may or may not be equal to the value of the function at $x=a$. In fact, the function may not even be defined at this point. If $f(x)$ is defined at $x=a$ and $L=f(a)$, then $f(x)$ is said to be continuous at $x=a$.
$\mathfrak{D e f i n i t i o n ~ ( C o n t i n u i t y ) : ~}$
Let $f: D \rightarrow \mathbb{R}$, where $D \subset \mathbb{R}$, and let $a \in D .{ }^{16}$ Then $f(x)$ is continuous at $x=a$ if for every $\epsilon>0$ there exists a $\delta>0$ such that

$$
|f(x)-f(a)|<\epsilon
$$

for all $x \in D$ for which $|x-a|<\delta$.

## Note:

To show the continuity of $f(x)$ at $x=a$, the following conditions must be verified:
(a). $f(x)$ is defined at all points inside a neighborhood of the point $a$.
(b). $f(x)$ has a limit from the left and a limit from the right as $x \rightarrow a$, and that these two limits are equal to $L$.
(c). The value of $f(x)$ at $x=a$ is equal to $L$.
$\mathfrak{D e f i n i t i o n ~ ( D i s c o n t i n u i t y ~ o f ~ t h e ~ F i r s t ~ K i n d ) : ~}$
A function $f: D \rightarrow \mathbb{R}$ has a discontinuity of the first kind at $x=a$ if $f\left(a^{+}\right)$and $f\left(a^{-}\right)$exist, but at least one of them is different from $f(a)$.

[^10]$\mathfrak{D e f i n i t i o n ~ ( D i s c o n t i n u i t y ~ o f ~ t h e ~ S e c o n d ~ K i n d ) : ~}$
A function $f: D \rightarrow \mathbb{R}$ has a discontinuity of the second kind at $x=a$ if at least one of $f\left(a^{+}\right)$and $f\left(a^{-}\right)$does not exist.
$\mathfrak{D e f i n i t i o n : ~ ( C o n t i n u i t y ~ o n ~ a ~ P o i n t ~ S e t ) : ~}$
A function $f: D \rightarrow \mathbb{R}$ is continuous on $E \subset D$ if it is continuous at every point of $E$.
$\mathfrak{D e f i n i t i o n ~ ( O n e - s i d e d ~ c o n t i n u i t y ) : ~}$
A function $f: D \rightarrow \mathbb{R}$ is left-continuous at $x=a$ if $\lim _{x \rightarrow a^{-}} f(x)=f(a)$. It is right-continuous at $x=a$ if $\lim _{x \rightarrow a^{+}} f(x)=f(a) .{ }^{17}$
$\mathfrak{D e f i n i t i o n ~ ( C o n t i n u i t y ~ o n ~ a ~ C l o s e d ~ I n t e r v a l ) : ~}$
A function $f$ is continuous on the closed interval $[a, b]$ if it is continuous on the open interval ( $a, b$ ) and
$$
\lim _{x \rightarrow a^{+}} f(x)=f(a) \quad \text { and } \quad \lim _{x \rightarrow b^{-}} f(x)=f(b) .
$$

The function is continuous from the right at $a$ and continuous from the left at $b$.
$\mathfrak{D e f i n i t i o n ~ ( U n i f o r m l y ~ C o n t i n u o u s ) : ~}$
The function $f: D \rightarrow \mathbb{R}$ is uniformly continuous on $E \subset D$ if for every $\epsilon>0$ there exists a $\delta>0$ such that

$$
\left|f\left(x_{1}\right)-f\left(x_{2}\right)\right|<\epsilon
$$

for all $x_{1}, x_{2} \in E$ for which $\left|x_{1}-x_{2}\right|<\delta$.

$$
\begin{aligned}
& { }^{17} \text { For example, the function } \\
& \qquad f(x)=\left\{\begin{array}{cc}
x-1, & x \leq 0, \\
1, & x>0
\end{array}\right.
\end{aligned}
$$

is left-continuous at $x=0$, since $f\left(0^{-}\right)=-1=f(0)$. If $f(x)$ were defined so that $f(x)=x-1$ for $x<0$ and $f(x)=1$ for $x \geq 0$, then it would be right-continuous at $x=0$.

## Example:

Show that $f(x)=x^{2}$ is continuous at $x_{0}$. We Choose any $h$ we please, say $h=1$. We consider $f$, then, in the interval $I:\left\{x_{0}-1<x<x_{0}+1\right\}$. Here, $\left|x-x_{0}\right|<1$. From this we see that

$$
|x|<\left|x_{0}\right|+1
$$

Now for $x$ in $I$ we consider

$$
\begin{aligned}
\left|f(x)-f\left(x_{0}\right)\right| & =\left|x^{2}-x_{0}^{2}\right|=\left|x-x_{0}\right|\left|x+x_{0}\right| \\
& \leq\left|x-x_{0}\right|\left(|x|+\left|x_{0}\right|\right) \leq\left|x-x_{0}\right|\left(2\left|x_{0}\right|+1\right)
\end{aligned}
$$

Consequently, we see that choosing $\left|x-x_{0}\right|<\epsilon /\left(2\left|x_{0}\right|+1\right)$ gives us

$$
\left|x^{2}-x_{0}^{2}\right|<\epsilon .
$$

To achieve this equality, we have imposed two conditions on $x$ :

$$
\begin{array}{r}
\left|x-x_{0}\right|<1 \quad \text { that } i s, x \in I, \\
\left|x-x_{0}\right|<\epsilon /\left(2\left|x_{0}\right|+1\right) .
\end{array}
$$

Thus, if we choose $\delta=\min \left[1, \epsilon /\left(2\left|x_{0}\right|+1\right)\right]$, our definition is satisfied. Note that the dependence of $\delta$ on $x_{0}$ is quite explicit.

## Example:

Show that the function defined by $f(x)=1 / x$ is continuous(uniformly) in the set $D:\{|x| \geq 1 / 2\}$. We have

$$
\left|f(x)-f\left(x_{0}\right)\right|=\left|\frac{1}{x}-\frac{1}{x_{0}}\right|=\frac{\left|x-x_{0}\right|}{\left|x x_{0}\right|}
$$

Hence, if we take $x$ and $x_{0}$ in $D$, and choose $\delta=\left|x_{0}\right| \epsilon / 2,{ }^{18}$ we get

$$
\left|f(x)-f\left(x_{0}\right)\right| \leq \frac{2\left|x-x_{0}\right|}{\left|x_{0}\right|}<\frac{2 \epsilon\left|x_{0}\right|}{2\left|x_{0}\right|}=\epsilon \quad \text { if }\left|x-x_{0}\right|<\delta .
$$

Now, as in any limit, if $\delta$ satisfies the conditions of the definition, so will any smaller number, say $\delta^{\prime} \leq \delta$. In this example we have $\delta=\left|x_{0}\right| \epsilon / 2$, where $x_{0} \geq 1 / 2$. Hence $\delta^{\prime}=\epsilon / 4$ is no larger than $\delta$, and therefore

$$
\left|f(x)-f\left(x_{0}\right)\right|<\epsilon \quad \text { if }\left|x-x_{0}\right|<\delta^{\prime}=\epsilon / 4
$$

[^11]for all $x, x_{0}$ in $D$.
The significance of our new choice, $\delta^{\prime}$, is that it will work for any arbitrary choice of $x_{0}$ in $D$; that is, $\delta^{\prime}$ is independent of $x_{0}$ in $D$. It depends only upon the set $D$ itself, but not on the individual point $x_{0}$ in $D$. Such a $\delta^{\prime}$ is called uniform, and we say that $f(x)$ is uniformly continuous in $D$.

### 4.4.1 Some Properties of Continuous Functions

## $\mathfrak{T h e o r e m}$ :

Let $f(x)$ and $g(x)$ be two continuous functions defined on a set $D \subset \mathbb{R}$. Then:
(a). $f(x)+g(x)$ and $f(x) g(x)$ are continuous on $D$.
(b). $\alpha f(x)$ is continuous on $D$, where $\alpha$ is a constant.
(c). $f(x) / g(x)$ is continuous on $D$ provided that $g(x) \neq 0$ on $D$.

## $\mathfrak{T h e o r e m}$ :

Suppose that $f: D \rightarrow \mathbb{R}$ is continuous on $D$, and $g: f(D) \rightarrow \mathbb{R}$ is continuous on $f(D)$, the image of $D$ under $f$. Then the composite function $h: D \rightarrow \mathbb{R}$ defined as $h(x)=g[f(x)]$ is continuous on $D$.

## $\mathfrak{P r o o f}:$

Let $\epsilon>$ be given, and let $a \in D$. Since $g$ is continuous at $f(a)$, there exist a $\delta^{\prime}$ such that $|g[f(x)]-g[f(a)]|<\epsilon$ if $|f(x)-f(a)|<\delta^{\prime}$. Since $f(x)$ is continuous at $x=a$, there exists a $\delta>0$ such that $|f(x)-f(a)|<\delta^{\prime}$ if $|x-a|<\delta$. It follows that by takeing $|x-a|<\delta$ we must have $|h(x)-h(a)|<\epsilon$.

## Theorem:

If $f(x)$ is continuous at $x=a$ and $f(a)>0$, then these exists a neighborhood $N_{\delta}(a)$ in which $f(x)>0$.

## $\mathfrak{P r o o f}:$

Since $f(x)$ is continuous at $x=a$, there exists a $\delta>0$ such that

$$
|f(x)-f(a)|<\frac{1}{2} f(a)
$$

if $|x-a|<\delta$. This implies that

$$
f(x)>\frac{1}{2} f(a)>0
$$

$\forall x \in N_{\delta}(a)$.
$\mathfrak{T h e o r e m}$ (The Intermediate-Value Theorem):
Let $f: D \rightarrow \mathbb{R}$, and let $[a, b]$ be a closed interval contained in $D$. Suppose that $f(a)>f(b)$. If $\lambda$ is a number such that $f(a)>\lambda>f(b)$, then there exists a point $c$, where $a<c<b$, such that $\lambda=f(c) .{ }^{19}$

In the following we want to establish some properties of continuous functions, defined on closed bounded intervals, that is intervals that not only are bounded above and below, but also include both end points.

## $\mathfrak{T h e o r e m}:$

Suppose that $f: D \rightarrow \mathbb{R}$ is continuous and that $D$ is bounded and closed. Then $f(x)$ is bounded in $D$.

## Theorem:

If $f: D \rightarrow \mathbb{R}$ is continuous, where $D$ is closed and bounded, then $f(x)$ achieves its infimum and supremum at least once in $D$, that is, there exists $\xi, \eta \in D$ such that

$$
\begin{aligned}
& f(\xi) \leq f(x) \quad \forall x \in D \\
& f(\eta) \geq f(x) \quad \forall x \in D
\end{aligned}
$$

Equivalently,

$$
\begin{aligned}
& f(\xi)=\inf _{x \in D} f(x), \\
& f(\eta)=\sup _{x \in D} f(x) .
\end{aligned}
$$

[^12]
## $\mathfrak{T h e o r e m}$ :

Let $f: D \mapsto \mathbb{R}$ be continuous on $D$. If $D$ is closed and bounded, then $f$ is uniformly continuous on $D$.

### 4.4.2 Lipschitz Continuous Functions

Lipschitz continuity is a specialized form of uniform continuity.

## Definition:

The function $f: D \mapsto \mathbb{R}$ is said to satisfy the Lipschitz condition on a set $E \subset D$ if there exist constants, $K$ and $\alpha$, where $K>0$ and $0<\alpha \leq 1$ such that

$$
\left|f\left(x_{1}\right)-f\left(x_{2}\right)\right| \leq K\left|x_{1}-x_{2}\right|^{\alpha}, \forall x_{1}, x_{2} \in E
$$

Notationally, whenever $f(x)$ satisfies the Lipschitz condition with constant $K$ and $\alpha$ on a set $E$, we say that it is $\operatorname{Lip}(K, \alpha) .{ }^{20}$

### 4.5 Convex Functions

Convex functions are frequently used in operations research. They also happen to be continuous. The natural domains for such functions are convex sets.

[^13]Thus,

$$
\left|\sqrt{x_{1}}-\sqrt{x_{2}}\right| \leq\left|x_{1}-x_{2}\right|^{1 / 2} .
$$

$\mathfrak{D e f i n i t i o n ~ ( C o n v e x ~ S e t ) : ~}$
A set $D \subset \mathbb{R}^{1}$ is convex if $\lambda x_{1}+(1-\lambda) x_{2} \in D$ whenever $x_{1}, x_{2}$ belong to $D$ and $0 \leq \lambda \leq 1$. Geometrically, a convex set contains the line segment connecting any two of its points. The same definition actually applies to convex sets in $\mathbb{R}^{n}$.

## Definition:

A function $f: D \rightarrow \mathbb{R}$ is convex if

$$
f\left[\lambda x_{1}+(1-\lambda) x_{2}\right] \leq \lambda f\left(x_{1}\right)+(1-\lambda) f\left(x_{2}\right), \forall x_{1}, x_{2} \in D,
$$

and any $\lambda$ such that $0 \leq \lambda \leq 1$.

Geometrically, a convex function in $\mathbb{R}$ means that if $P$ and $Q$ are any two points on the graph of $y=f(x)$, then the portion of the graph between $P$ and $Q$ lies between the chord $P Q$. Example of convex functions include $f(x)=x^{2}$ on $\mathbb{R}$, $f(x)=e^{x}$ on $\mathbb{R}$, to name just a few.
$\mathfrak{D e f i n i t i o n : ~}$
A function $f: D \rightarrow \mathbb{R}$ is concave if $-f$ is convex.

## $\mathfrak{L e m m a : ~}$

If $f:[a, b] \rightarrow \mathbb{R}$ is convex and the value of $f$ at $a$ and $b$ are finite, then $f(x)$ is bounded from above on $[a, b]$ by $M=\max \{f(a), f(b)\}$, and $f(x)$ is also bounded from below.
$\mathfrak{P r o o f}$ :
(a). Because $x \in[a, b]$, then $x=\lambda a+(1-\lambda) b$ for some $\lambda \in[0,1]$, since $[a, b]$ is a convex set. Hence,

$$
\begin{aligned}
f(x) & \leq \lambda f(a)+(1-\lambda) f(b) \\
& \leq \lambda M+(1-\lambda) M=M
\end{aligned}
$$

(b). We First note that any $x \in[a, b]$ can be written

$$
x=\frac{a+b}{2}+t,
$$

where

$$
a-\frac{a+b}{2} \leq t \leq b-\frac{a+b}{2}
$$

Now, if $(a+b) / 2+t=x$ belong to $[a, b]$, so does $(a+b) / 2-t$, then since

$$
\frac{1}{2}\left[\frac{a+b}{2}+t\right]+\frac{1}{2}\left[\frac{a+b}{2}-t\right]=\frac{a+b}{2}
$$

so

$$
f\left(\frac{a+b}{2}\right) \leq \frac{1}{2} f\left(\frac{a+b}{2}+t\right)+\frac{1}{2} f\left(\frac{a+b}{2}-t\right),
$$

or

$$
f\left(\frac{a+b}{2}+t\right) \geq 2 f\left(\frac{a+b}{2}\right)-f\left(\frac{a+b}{2}-t\right) .
$$

Since

$$
f\left(\frac{a+b}{2}-t\right) \leq M
$$

then

$$
f\left(\frac{a+b}{2}+t\right) \geq 2 f\left(\frac{a+b}{2}\right)-M
$$

That is, $f(x) \geq m \forall x \in[a, b]$, where $m=f\left(\frac{a+b}{2}\right)-M$.

## $\mathfrak{T h e o r e m}:$

Let $f: D \rightarrow \mathbb{R}$ be a convex function, where $D$ is an open interval. Then $f$ is $\operatorname{Lip}(K, 1)$ on any closed interval $[a, b]$ contained in $D$, that is,

$$
\mid f\left(x_{1}\right)-f\left(x_{2}|\leq K| x_{1}-x_{2} \mid, \quad \forall x_{1}, x_{2} \in[a, b] .\right.
$$

## $\mathfrak{P r o o f}:$

Consider the closed interval $[a-\epsilon, b+\epsilon]$, where $\epsilon>0$ is chosen that this interval is contained in $D$. Let $m^{\prime}$ and $M^{\prime}$ be, respectively, the lower and upper bounds of $f$ on $[a-\epsilon, b+\epsilon]$. Let $x_{1}, x_{2}$ be any two distinct points in $[a, b]$. Define $z_{1}$ and $\lambda$ as

$$
\begin{aligned}
z_{1} & =x_{2}+\frac{\epsilon\left(x_{2}-x_{1}\right)}{\left|x_{1}-x_{2}\right|} \\
\lambda & =\frac{\left|x_{1}-x_{2}\right|}{\epsilon+\left|x_{1}-x_{2}\right|} .
\end{aligned}
$$

Then $z_{1} \in[a-\epsilon, b+\epsilon]$. This is true because $\left(x_{2}-x_{1}\right) /\left|x_{1}-x_{2}\right|$ is either equal to 1 or to -1 . Since $x_{2} \in[a, b]$, then

$$
a-\epsilon \leq x_{2}-\epsilon \leq x_{2}+\frac{\epsilon\left(x_{2}-x_{1}\right)}{\left|x_{1}-x_{2}\right|} \leq x_{2}+\epsilon \leq b+\epsilon
$$

Furthermore, it can be verified that

$$
x_{2}=\lambda z_{1}+(1-\lambda) x_{1} .
$$

We then have

$$
f\left(x_{2}\right) \leq \lambda f\left(z_{1}\right)+(1-\lambda) f\left(x_{1}\right)=\lambda\left[f\left(z_{1}\right)-f\left(x_{1}\right)\right]+f\left(x_{1}\right) .
$$

Thus,

$$
\begin{aligned}
f\left(x_{2}\right)-f\left(x_{1}\right) & \leq \lambda\left[f\left(z_{1}\right)-f\left(x_{1}\right)\right] \\
& \leq \lambda\left[M^{\prime}-m^{\prime}\right] \\
& \leq \frac{\left|x_{1}-x_{2}\right|}{\epsilon}\left(M^{\prime}-m^{\prime}\right)=K\left|x_{1}-x_{2}\right|
\end{aligned}
$$

## Corollary:

Let $f: D \rightarrow \mathbb{R}$ be a convex function, where $D$ is an open interval. If $[a, b]$ is any closed interval contained in $D$, then $f(x)$ is uniformly continuous on $[a, b]$ and is therefore continuous on $D$.

## 5 Differentiation

Differentiation originated in connection with the problem of drawing tangents to curves and of finding maxima and minima of functions.

### 5.1 The Derivative of a Function

The notation of differentiation was motivated by the need to find the tangent to a curve at a given point. Fetmat's approach to this problem was inspired by a geometric reasoning. His method uses the idea of a tangent as the limiting position of a secant when two of its points of intersection with the curve tend to coincide.
$\mathfrak{D e f i n i t i o n ~ ( D e r i v a t i v e ) ~ : ~}$
Let $f(x)$ be a function defined in a neighborhood $N_{r}\left(x_{0}\right)$ of a point $x_{0}$. Consider the ration

$$
\phi(h)=\frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h},
$$

where $h$ is a nonzero increment of $x_{0}$ such that $-r<h<r$. If $\phi(h)$ has a limit as $h \rightarrow 0$, then the limit is called the derivative of $f(x)$ at $x_{0}$ and is denoted by $f^{\prime}\left(x_{0}\right)$. It is also common to sue the notation

$$
\left.\frac{d f(x)}{d x}\right|_{x=x_{0}}=f^{\prime}\left(x_{0}\right)
$$

We thus have

$$
\begin{equation*}
f^{\prime}\left(x_{0}\right)=\lim _{h \rightarrow 0} \frac{f\left(x_{0}+h\right)-f\left(x_{0}\right)}{h} . \tag{1}
\end{equation*}
$$

By putting $x=x_{0}+h$,(1) can be written as

$$
f^{\prime}\left(x_{0}\right)=\lim _{x \rightarrow x_{0}} \frac{f(x)-f\left(x_{0}\right)}{x-x_{0}}
$$

If $f^{\prime}\left(x_{0}\right)$ exists, then $f(x)$ is said to be differentiable at $x=x_{0}$. Geometrically, $f^{\prime}\left(x_{0}\right)$ is the slope of the tangent to the graph of the function $y=f(x)$ at the point
$\left(x_{0}, y_{0}\right)$, where $y_{0}=f\left(x_{0}\right) .{ }^{21}$
$\mathfrak{D e f i n i t i o n ~ ( D i f f e r e n t i a b l e ~ o n ~ a ~ s e t ) : ~}$
If $f(x)$ has a derivative at every point of a set $D$, then $f(x)$ is said to be differentiable on $D$.

## $\mathfrak{D e f i n i t i o n ~ ( S e c o n d ~ a n d ~ H i g h e r ~ D e r i v a t i v e ) : ~}$

If $f(x)$ is differentiable on a set $D$, then $f^{\prime}(x)$ is a function itself defined on $D$. In the event $f^{\prime}(x)$ itself is differentiable on $D$, then its derivative is called the second derivative of $f(x)$ and is denoted by $f^{\prime \prime}(x)$. It is also common to use the notation

$$
\frac{d\left(\frac{d f(x)}{d x}\right)}{d x}=\frac{d^{2} f(x)}{d x^{2}}=f^{\prime \prime}(x)
$$

By the same token, we can define the $n$th $(n \geq 2)$ derivative of $f(x)$ as the derivative of the $(n-1)$ st derivative of $f(x)$. We denote this derivative by

$$
\frac{d^{n} f(x)}{d x^{n}}=f^{(n)}(x), \quad n=2,3, \ldots \ldots
$$

Theorem (Differentiability implies Continuity):
Let $f(x)$ be defined at in a neighborhood of a point $x_{0}$. If $f(x)$ has derivative at $x_{0}$, then it must be continuous at $x_{0}$.

## $\mathfrak{P r o o f}:$

From definition above we can write

$$
f\left(x_{0}+h\right)-f\left(x_{0}\right)=h \phi(h)
$$

If the derivative of $f(x)$ exists at $x_{0}$, then $\phi(h) \rightarrow f^{\prime}\left(x_{0}\right)$ as $h \rightarrow 0$. It follows that

$$
f\left(x_{0}+h\right)-f\left(x_{0}\right) \rightarrow 0 \cdot f^{\prime}\left(x_{0}\right)=0
$$

[^14]as $h \rightarrow 0$. Thus for a given $\epsilon>0$ there exists a $\delta>0$ such that
$$
\left|f\left(x_{0}+h\right)-f\left(x_{0}\right)\right|<\epsilon
$$
if $|h|<\delta$, i.e.
$$
\left|f(x)-f\left(x_{0}\right)\right|<\epsilon
$$
if $\left|x-x_{0}\right|<\delta$. This indicates that $f(x)$ is continuous at $x_{0}$.

### 5.1.1 Rules Pertaining to Differentiation

$\mathfrak{T h e o r e m}$ (Power Rule):
If $n$ is a rational number, then the function $f(x)=x^{n}$ is differentiable and

$$
\frac{d f(x)}{d x}=\frac{d x^{n}}{d x}=n x^{n-1}
$$

## Theorem:

Let $f(x)$ and $g(x)$ be defined and differentiable on a set $D$. Then
(a). $[\alpha f(x)+\beta g(x)]^{\prime}=\alpha f^{\prime}(x)+\beta g^{\prime}(x)$.
(b). $[f(x) g(x)]^{\prime}=f^{\prime}(x) g(x)+f(x) g^{\prime}(x)$.
(c). $[f(x) / g(x)]^{\prime}=\left[f^{\prime}(x) g(x)-f(x) g^{\prime}(x)\right] / g^{2}(x)$ if $g(x) \neq 0$.

## $\mathfrak{P r o o f}:$

To prove (b) we write

$$
\begin{aligned}
& \lim _{h \rightarrow 0} \frac{f(x+h) g(x+h)-f(x) g(x)}{h} \\
& =\lim _{h \rightarrow 0} \frac{[f(x+h)-f(x)] g(x+h)+f(x)[g(x+h)-g(x)]}{h} \\
& =\lim _{h \rightarrow 0} \frac{f(x+h)-f(x)] g(x+h)}{h}+\lim _{h \rightarrow 0} \frac{f(x)[g(x+h)-g(x)]}{h} \\
& =\lim _{h \rightarrow 0} g(x+h) \lim _{h \rightarrow 0} \frac{f(x+h)-f(x)}{h}+f(x) \lim _{h \rightarrow 0} \frac{g(x+h)-g(x)}{h} .
\end{aligned}
$$

However, $\lim _{h \rightarrow 0} g(x+h)=g(x)$, since $g(x)$ is continuous (because it is differentiable). Hence,

$$
\lim _{h \rightarrow 0} \frac{f(x+h) g(x+h)-f(x) g(x)}{h}=g(x) f^{\prime}(x)+f(x) g^{\prime}(x) .
$$

Now, to prove (c) we write

$$
\begin{aligned}
& \lim h \rightarrow 0 \frac{f(x+h) / g(x+h)-f(x) / g(x)}{h} \\
& =\lim _{h \rightarrow 0} \frac{g(x) f(x+h)-f(x) g(x+h)}{h g(x) g(x+h)} \\
& =\lim _{h \rightarrow 0} \frac{g(x)[f(x+h)-f(x)]-f(x)[g(x+h)-g(x)]}{h g(x) g(x+h)} \\
& =\frac{\lim _{h \rightarrow 0}\{g(x)[f(x+h)-f(x)] / h-f(x)[g(x+h)-g(x)] / h\}}{g(x) \lim _{h \rightarrow 0} g(x+h)} \\
& =\frac{g(x) f^{\prime}(x)-f(x) g^{\prime}(x)}{g^{2}(x)} .
\end{aligned}
$$

$\mathfrak{T h e o r e m}$ (The Chain Rule):
Let $f: D_{1} \rightarrow \mathbb{R}$ and $g: D_{2} \rightarrow \mathbb{R}$ be two functions. Suppose that $f\left(D_{1}\right) \subset D_{2}$. If $f(x)$ is differentiable on $D_{1}$ and $g(x)$ is differentiable on $D_{2}$, then the composite function $h(x)=g[f(x)]$ is differentiable on $D_{1}$ and

$$
\frac{d g[f(x)]}{d x}=\frac{d g[f(x)]}{d f(x)} \frac{d f(x)}{d x} .
$$

## $\mathfrak{P r o o f}:$

Let $z=f(x)$ and $t=f(x+h)$. By the fact that $g(z)$ is differentiable we can write

$$
\begin{aligned}
g[f(x+h)]-g[f(x)] & =g(t)-g(z) \\
& =(t-z) g^{\prime}(z)+o(t-z) .
\end{aligned}
$$

We then have

$$
\frac{g[f(x+h)]-g[f(x)]}{h}=\frac{t-z}{h} g^{\prime}(z)+\frac{o(t-z)}{t-z} \frac{t-z}{h} .
$$

Since $f$ is differentiable, then it is continuous, so as $h \rightarrow 0, f(x+h) \rightarrow f(x)$, i.e. $t \rightarrow z$. Hence,

$$
\lim _{h \rightarrow 0} \frac{t-z}{h}=\lim _{h \rightarrow 0} \frac{f(x+h)-f(x)}{h}=\frac{d f(x)}{d x} .
$$

Noting that

$$
\lim h \rightarrow 0 \frac{o(t-z)}{t-z}=\lim _{t \rightarrow z} \frac{o(t-z)}{t-z}=0
$$

and we conclude that

$$
\frac{d g[f(x)]}{d x}=\frac{d g[f(x)]}{d f(x)} \frac{d f(x)}{d x} .
$$

## Theorem:

Let $f: D \mapsto \mathbb{R}$, where $D$ is an open set. Suppose that $f^{\prime}(x)$ is positive at a point $x_{0} \in D$. Then there is a neighborhood $N_{\delta}\left(x_{0}\right) \in D$ such that for each $x$ in this neighborhood, $f(x)>f\left(x_{0}\right)$ if $x>x_{0}$, and $f(x)<f\left(x_{0}\right)$ if $x<x_{0} .{ }^{22}$

### 5.2 The Mean Value Theorem

This is one of the most important theorems in differential calculus. It is also known as the theorem of the mean.

Theorem (Rolle's Theorem):
Let $f(x)$ be continuous on the closed interval $[a, b]$ and differentiable on the open interval $(a, b)$. If $f(a)=f(b)$, then there exists a point $c, a<c<b$, such that $f^{\prime}(c)=0$.
$\mathfrak{T h e o r e m}$ (Mean Value Theorem):
If $f(x)$ is continuous on the closed interval $[a, b]$ and differentiable on the open interval $(a, b)$, then there exists a point $c, a<c<b$, such that

$$
f^{\prime}(c)=\frac{f(b)-f(a)}{b-a} .
$$

$\mathfrak{P r o o f}:$
Consider the function

$$
F(x)=f(x)-f(a)-\frac{f(b)-f(a)}{b-a}(x-a) .
$$

The function $F(x)$ is continuous on $[a, b]$ and is differentiable on $(a, b)$, since $F^{\prime}(x)=$ $f^{\prime}(x)-\frac{f(b)-f(a)}{b-a}$. Furthermore, $F(a)=F(b)=0$ it follows from Rolle's theorem that

[^15]there exists a point $c, a<c<b$, such that $F^{\prime}(c)=0$. Thus
$$
f^{\prime}(c)=\frac{f(b)-f(a)}{b-a}
$$
$\mathfrak{T h e o r e m}$ (Cauchy's Mean Value Theorem):
If $f(x)$ and $g(x)$ are continuous on the closed interval $[a, b]$ and differentiable on the open interval ( $a, b$ ), then there exists a point $c, a<c<b$, such that
$$
[f(b)-f(a)] g^{\prime}(c)=[g(b)-g(a)] f^{\prime}(c) .
$$

## $\mathfrak{P r o o f}:$

See Fulks, p. 112 for a good interpretation. When the mean values theorem is a special case of Cauchy's mean values theorem when $g(x)=x$.

An immediate application of the Cauchy's mean theorem is a very popular method in calculating the limit of certain ratios of functions. This method is known as l'Hospital's rule. It deal with the limit of the ratio $f(x) / g(x)$ as $x \rightarrow a$ when both the numerator and the denominator tend simultaneously to zero or to infinity as $x \rightarrow a$. In either case, we have what is called an indeterminate ratio caused by having $0 / 0$ or $\infty / \infty$ as $x \rightarrow a$.

## $\mathfrak{T h e o r e m}$ (l'Hospital's Rule):

Let $f(x)$ and $g(x)$ be continuous on the closed interval $[a, b]$ and differentiable on the open interval $(a, b)$. Suppose that we have the following:
(a). $g(x)$ and $g^{\prime}(x)$ are not zero at any point inside $(a, b)$.
(b). $\lim _{x \rightarrow a^{+}} f^{\prime}(x) / g^{\prime}(x)$ exists.
(c). $f(x) \rightarrow 0$ and $g(x) \rightarrow 0$ as $x \rightarrow a^{+}$, or $f(x) \rightarrow \infty$ and $g(x) \rightarrow \infty$ as $x \rightarrow a^{+}$. Then,

$$
\lim _{x \rightarrow a^{+}} \frac{f(x)}{g(x)}=\lim _{x \rightarrow a^{+}} \frac{f^{\prime}(x)}{g^{\prime}(x)} .
$$

## Example:

$\lim _{x \rightarrow \infty} x \log \left(\frac{x+1}{x-1}\right)$.
This is of the form $\infty \times 0$ as $x \rightarrow \infty$, which is indeterminate. But

$$
\lim _{x \rightarrow \infty} x \log \left(\frac{x+1}{x-1}\right)=\frac{\log \left(\frac{x+1}{x-1}\right)}{1 / x}
$$

is of the form $0 / 0$ as $x \rightarrow \infty$. Hence,

$$
\begin{aligned}
\lim _{x \rightarrow \infty} x \log \left(\frac{x+1}{x-1}\right) & =\lim _{x \rightarrow \infty} \frac{\frac{-2}{(x+1)(x-1)}}{-1 / x^{2}} \\
& =\lim _{x \rightarrow \infty} \frac{2}{(1-1 / x)(1+1 / x)} \\
& =2
\end{aligned}
$$

### 5.3 Taylor's Theorem

This theorem is also known as the general mean value theorem, since it is considered as an extension of the mean value theorem. It was used to expand functions into infinite series.

Theorem (Taylor's Theorem):
If the $(n-1)$ st $(n \geq 1)$ derivative of $f(x)$, namely $f^{(n-1)}(x)$, is continuous on the closed interval $[a, b]$ and the $n$th derivative $f^{(n)}$ exists on the open interval $(a, b)$, then for each $x \in[a, b]$ we have ${ }^{23}$
$f(x)=f(a)+(x-a) f^{\prime}(a)+\frac{(x-a)^{2}}{2!} f^{\prime \prime}(a)+\cdots+\frac{(x-a)^{n-1}}{(n-1)!} f^{(n-1)}(a)+\frac{(x-a)^{n}}{n!} f^{(n)}(\xi)$,
where $a<\xi<x$.

This is known as Taylor's formula. It can also be expressed as

$$
f(a+h)=f(a)+h f^{\prime}(a)+\frac{h^{2}}{2!} f^{\prime \prime}(a)+\cdots+\frac{h^{n-1}}{(n-1)!} f^{(n-1)}(a)+\frac{h^{n}}{n!}\left(a+\theta_{n} h\right)
$$

where $h=x-a$ and $0<\theta_{n}<1 .{ }^{24}$

### 5.4 Maxima and Minima of a Function

[^16]$\mathfrak{D e f i n i t i o n : ~}$
A function $f: D \rightarrow \mathbb{R}$ has a local maximum at a point $x_{0} \in D$ if there exists a $\delta>0$ such that $f(x) \leq f\left(x_{0}\right) \forall x \in N_{\delta}\left(x_{0}\right) \cap D$. The function $f$ has a local minimum at a point $x_{0} \in D$ if there exists a $\delta>0$ such that $f(x) \geq f\left(x_{0}\right) \forall x \in N_{\delta}\left(x_{0}\right) \cap D$.

## Definition:

A function $f: D \rightarrow \mathbb{R}$ has an absolute maximum (minimum) over $D$ if there exist a point $x^{*} \in D$ such that $f(x) \leq f\left(x^{*}\right)\left(f(x) \geq f\left(x^{*}\right)\right) \forall x \in D$.

## $\mathfrak{T h e o r e m}$ :

Let $f(x)$ be differentiable on the open interval $(a, b)$. If $f(x)$ has a local maximum, or a local minimum, at a point $x_{0}$ in $(a, b)$, then $f^{\prime}\left(x_{0}\right)=0$.

## $\mathfrak{P r o o f}:$

Suppose that $f(x)$ has a local maximum at $x_{0}$. Then $f(x) \leq f\left(x_{0}\right), \forall x \in N_{\delta}\left(x_{0}\right) \subset$ $(a, b)$. It follows that

$$
\frac{f(x)-f\left(x_{0}\right)}{x-x_{0}}\left\{\begin{array}{l}
\leq 0 \quad \text { if } x>x_{0}  \tag{2}\\
\geq 0 \\
\text { if } x<x_{0}
\end{array}\right.
$$

$\forall x \in N_{\delta}\left(x_{0}\right)$. As $x \rightarrow x_{0}^{+}$, the ratio in (2) have a non-positive limit, and if $x \rightarrow x_{0}^{-}$, the ration will have a nonnegative limit. Since $f^{\prime}\left(x_{0}\right)$ exists, these two limits must be equal and equal to $f^{\prime}\left(x_{0}\right)$ as $x \rightarrow x_{0}$. We therefore conclude that $f^{\prime}\left(x_{0}\right)=0$.

## Note:

It is important to note that $f^{\prime}\left(x_{0}\right)=0$ is a necessary condition for a differentiable function to have a local optimum at $x_{0}$. It is not, however, a sufficient condition. That is, if $f^{\prime}\left(x_{0}\right)=0$, then it is not necessary true that $x_{0}$ is a point of local optimum. We shall in the next subsection make use of Taylor's expansion to come up with a condition for $f(x)$ to have local optimum at $x=x_{0}$.

## $\mathfrak{N o t e}:$

We recall from last section that if $f(x)$ is continuous on $[a, b]$, then it must achieve its absolute optima at some points inside $[a, b]$. These points can be interior points, that is, points that belong to the open interval $(a, b)$, or they can be end (boundary) points. In particular, if $f^{\prime}(x)$ exists on $(a, b)$, to determine the locations of the
absolute optima we must solve the equation $f^{\prime}(x)=0$ and then compare the value of $f(x)$ at the roots of this equation with $f(a)$ and $f(b)$. The largest of these values is the absolute maximum. In the event $f^{\prime}(x) \neq 0$ on $(a, b)$, then $f(x)$ must achieve its absolute optimum at an end point.

### 5.4.1 A Sufficient Condition for a Local Optimum

Suppose that $f(x)$ has $n$ derivatives in a neighborhood $N_{\delta}\left(x_{0}\right)$ such that $f^{\prime}\left(x_{0}\right)=$ $f^{\prime \prime}\left(x_{0}\right)=\cdots=f^{(n-1)}\left(x_{0}\right)=0$, but $f^{(n)}\left(x_{0}\right) \neq 0$. Then by Taylor's theorem we have

$$
f(x)=f\left(x_{0}\right)+\frac{h^{n}}{n!} f^{(n)}\left(x_{0}+\theta_{n} h\right)
$$

for any $x$ in $N_{\delta}\left(x_{0}\right)$, where $h=x-x_{0}$ and $0<\theta_{n}<1$. Furthermore, if we assume that $f^{(n)}(x)$ is continuous at $x_{0}$, then

$$
f^{(n)}\left(x_{0}+\theta_{n} h\right)=f^{(n)}\left(x_{0}\right)+o(1)
$$

where $o(1) \rightarrow 0$ as $h \rightarrow 0$. We can therefore write ${ }^{25}$

$$
\begin{equation*}
f(x)-f\left(x_{0}\right)=\frac{h^{n}}{n!} f^{(n)}\left(x_{0}\right)+o\left(h^{n}\right) . \tag{3}
\end{equation*}
$$

In order for $f(x)$ to have a local optimum at $x_{0}, f(x)-f\left(x_{0}\right)$ must have the same sign (positive or negative) for small values of $h$ inside a neighborhood of 0 . But from (3), the sign of $f(x)-f\left(x_{0}\right)$ is determined by the sign of $h^{n} f^{(n)}\left(x_{0}\right)$. We can then conclude that:
(a). If $n$ is even, then a local optimum is achieved at $x_{0}$. In this case, a local maximum occurs at $x_{0}$ if $f^{(n)}<0$, whereas $f^{(n)}>0$ indicates that $x_{0}$ is a point of local minimum.
(b). If $n$ is odd, then $x_{0}$ is not a point of local optimum, since $f(x)-f\left(x_{0}\right)$ changes sign around $x_{0}$.

In particular, if $f^{\prime}\left(x_{0}\right)=0$ and $f^{\prime \prime}\left(x_{0}\right) \neq 0$, then $x_{0}$ is a point of local optimum. When $f^{\prime \prime}\left(x_{0}\right)<0, f(x)$ has a local maximum at $x_{0}$, and when $f^{\prime \prime}\left(x_{0}\right)>0, f(x)$ has a local minimum at $x_{0}$.

[^17]
## 6 Infinite Sequence and Series

The study of the theory of infinite sequence and series is an integral part of advanced calculus. All limiting processes, such as differentiation and integration, can be investigated on the basis of this theory. In this chapter we shall study the theory of infinite sequences and series, and investigate their convergence. Unless otherwise stated, the terms of all sequences and series considered in this chapter are realvalued.

### 6.1 Infinite Sequences

$\mathfrak{D e f i n i t i o n ~ ( I n f i n i t e ~ S e q u e n c e s ) : ~}$
An infinite sequence is a particular function $f: J^{+} \rightarrow \mathbb{R}$ defined on the set of all positive integers. For a given $n \in J^{+}$, the value of this function, namely $f(n)$, is called the $n$th term of the infinite sequence and is denoted by $a_{n}$. The sequence itself is denoted by the symbol $\left\{a_{n}\right\}_{n=1}^{\infty}$.

### 6.1.1 Bound, Convergence and Divergence

Since a sequence is a function, then in particular, the sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ can have the following properties:
(a). It is bounded if there exists a constant $K>0$ such that $\left|a_{n}\right| \leq K, \forall n$.
(b). It is monotone increasing if $a_{n} \leq a_{n+1} \forall n$, and is monotone decreasing if $a_{n} \geq a_{n+1} \forall n$.
(c). It converges to a finite number $c$ if $\lim _{n \rightarrow \infty} a_{n}=c$, that is, for a given $\epsilon>0$ there exists an integer $N$ such that

$$
\left|a_{n}-c\right|<\epsilon \quad \text { if } n>N .
$$

In this case, $c$ is called the limit of the sequence and this fact is denoted by writing

$$
a_{n} \rightarrow c \text { as } n \rightarrow \infty .
$$

If the sequence does not converge to a finite limit, then it is said to be divergent.
(d). It is said to oscillate if it does not converge to a finite limit, nor to $+\infty$ or $-\infty$
as $n \rightarrow \infty$.

## $\mathfrak{E x a m p l e}$ :

let $a_{n}=\left(n^{2}+2 n\right) /\left(2 n^{2}+3\right)$. Then $a_{n} \rightarrow \frac{1}{2}$ as $n \rightarrow \infty$, since

$$
\lim _{n \rightarrow \infty}=\lim _{n \rightarrow \infty} \frac{1+2 / n}{2+3 / n^{2}}=\frac{1}{2}
$$

$\mathfrak{T h e o r e m}$ (Convergence implies Boundedness):
Every convergent sequence is bounded.

## $\mathfrak{P r o o f}:$

Let $c=\lim a_{n}$. Then there is an $N$ such that

$$
\left|a_{n}-c\right|<1 \quad \text { if } n>N .
$$

Then

$$
\left|a_{n}\right|-|c| \leq\left|a_{n}-c\right|<1 .
$$

Thus

$$
\begin{equation*}
\left|a_{n}\right|<|c|+1 \equiv K_{1} \quad \text { if } n>N . \tag{4}
\end{equation*}
$$

Now we look among the number $\left|a_{1},\left|a_{2}\right|, \ldots,\left|a_{N}\right|\right.$, and choose the largest, calling it $K_{2}$. Then

$$
\begin{equation*}
\left|a_{n}\right| \leq K_{2} \quad \text { if } n \leq N . \tag{5}
\end{equation*}
$$

It is clear that if we take $K$ to be the larger of $K_{1}$ and $K_{2}$, then by (4) and (5),

$$
\left|a_{n}\right|<K \quad \forall n .
$$

The converse of this theorem is not necessarily true. That is, if a sequence is bounded, then it does not have to be convergent. ${ }^{26}$ To guarantee converge of a bounded sequence we obviously need an additional condition.

[^18]
## $\mathfrak{T h e o r e m}$ :

Every bounded monotone sequence converges.

## $\mathfrak{T h e o r e m}$ :

(a). If $\left\{a_{n}\right\}_{n=1}^{\infty}$ is bounded from above and is monotone increasing, then $\left\{a_{n}\right\}_{n=1}^{\infty}$ converges to $c=\sup _{n \geq 1} a_{n}$.
(b). If $\left\{a_{n}\right\}_{n=1}^{\infty}$ is bounded from below and is monotone decreasing, then $\left\{a_{n}\right\}_{n=1}^{\infty}$ converges to $d=\inf _{n \geq 1} a_{n}$
$\mathfrak{D e f i n i t i o n ~ ( S u b s e q u e n c e ) : ~}$
A sequence $\left\{a_{n}\right\}$ is a function on the integers; that is

$$
a_{n}=f(n) \quad \text { for all } n \geq n_{0} .
$$

Now suppose that we consider the function restricted to a subset of the integers. Let us choose an integer greater than or equal to $n_{0}$ and denote it by $n_{1}$, another greater than $n_{1}$ and denote it by $n_{2}$, another greater than $n_{2}$ and denote it by $n_{3}$, and so forth. Then the new sequence, defined by

$$
b_{k}=a_{n_{k}}=f\left(n_{k}\right), \quad k=1,2, \ldots
$$

we call a subsequence of $\left\{a_{n}\right\}$. It is clear that there are many subsequences of a given sequence. If we assume the $n_{0} \geq 1$, then $n_{k} \geq k$. The theorem we want to prove is the following.

## $\mathfrak{T h e o r e m}$ :

Suppose that $\left\{a_{n}\right\}$ converges; then any subsequence $\left\{a_{n_{k}}\right\}$ also converges and has the same limit.

## $\mathfrak{P r o o f}:$

Let $A$ be the limit of the sequence $\left\{a_{n}\right\}$. We know that for each $\varepsilon>0$ there is an $N$ for which

$$
\left|a_{n_{k}}-A\right|<\varepsilon \quad \text { if } \quad n_{k}>N .
$$

But $n_{k} \geq k$; hence

$$
n_{k}>N \quad \text { if } k>N
$$

Thus

$$
\left|b_{k}-A\right|=\left|a_{n_{k}}-A\right|<\varepsilon \quad \text { if } \quad k>N .
$$

$\mathfrak{N o t e}$ :
It should be noted that if a sequence diverges, then it does not necessarily follows that every one of its subsequences must diverge. A sequence may fail to converge, yet several of its subsequence converges. We have noted that a bounded sequence may not converge. It is possible, however, that one of its subsequences is convergent.

## Theorem:

Every bounded sequence has a convergent subsequence.
$\mathfrak{D e f i n i t i o n ~ ( U p p e r ~ a n d ~ L o w e r ~ L i m i t ~ o f ~ a ~ S e q u e n c e ) : ~}$
Let $\left\{a_{n}\right\}_{n=1}^{\infty}$ be a bounded sequence, ${ }^{27}$ and let $E$ be the set of all its subsequential limits. Then the least upper bound of $E$ is called the upper limit of $\left\{a_{n}\right\}_{n=1}^{\infty}$ and is denoted by $\lim \sup _{n \rightarrow \infty} a_{n}$. Similarly, the greatest lower bound of $E$ is called the lower limit of $\left\{a_{n}\right\}_{n=1}^{\infty}$ and is denoted by $\liminf _{n \rightarrow \infty} a_{n}$.

## $\mathfrak{E x a m p l e}$ :

The sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$, where $a_{n}=(-1)^{n}[1+(1 / n)]$, has two subsequential limits, namely -1 and +1 . Thus $E=\{-1,1\}$, and $\limsup _{n \rightarrow \infty}=1$ and $\liminf _{n \rightarrow \infty}=-1$.

Theorem:
The sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ converges to $c$ if and only if

$$
\lim \sup _{n \rightarrow \infty}=\lim \inf _{n \rightarrow \infty}=c
$$

[^19]
### 6.1.2 The Cauchy Criterion

We have seen earlier that the definition of convergence of a sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ requires finding the limit of $a_{n}$ as $n \rightarrow \infty$. In some cases, such a limit may be difficult to figure out. ${ }^{28}$ Fortunately, however, there is another convergence criterion for sequence, known as the Cauchy Criterion after Augustin-Louis Cauchy.
$\mathfrak{T h} \mathfrak{e r e m}$ (The Cauchy Criterion):
The sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ converges if and only if it satisfies the following condition, known as the $\epsilon$-condition: For each $\epsilon>0$ there is an integer $N$ such that

$$
\left|a_{m}-a_{n}\right|<\epsilon \quad \text { for all } m>N, n>N .
$$

## Definition:

A sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ that satisfies the $\epsilon$-condition of the Cauchy criterion is said to be a Cauchy sequence.

### 6.2 Infinite Series

Let $\left\{a_{n}\right\}_{n=1}^{\infty}$ be a given sequence. Consider the symbolic expression

$$
\sum_{n=1}^{\infty} a_{n}=a_{1}+a_{2}+\cdots+a_{n}+\cdots
$$

By definition, this expression is called an infinite series, or just a series for simplicity, and $a_{n}$ is referred as the $n$th term of the series. The finite sum

$$
s_{n}=\sum_{i=1}^{n} a_{i}, \quad n=1,2, \ldots
$$

[^20]is called the $n$th partial sum of the series.
$\mathfrak{D e f i n i t i o n ~ ( C o n v e r g e n c e ~ o f ~ a n ~ I n f i n i t e ~ S e r i e s ) : ~}$
Consider the series $\sum_{n=1}^{\infty} a_{n}$. Let $s_{n}$ be its $n$th partial sum $(n=1,2, \ldots)$.
(a). The series is said to be convergent if the sequence $\left\{s_{n}\right\}_{n=1}^{\infty}$ converges. In this case, if $\lim _{n \rightarrow \infty} s_{n}=s$, where $s$ is finite, then we say that the series converges to $s$, or that $s$ is the sum of the series. This expressed by writing
$$
s=\sum_{n=1}^{\infty} a_{n} .
$$
(b). If $s_{n}$ does not tend to a finite limit, then the series is said to be divergent.

## Theorem:

The series $\sum_{n=1}^{\infty} a_{n}$, converges if and only if for a given $\epsilon>0$ there is an integer $N$ such that

$$
\begin{equation*}
\left|\sum_{i=m+1}^{n} a_{i}\right|<\epsilon \quad \text { for all } n>m>N . \tag{6}
\end{equation*}
$$

Eq.(6) follows from applying Cauchy Criterion to the sequence $\left\{s_{n}\right\}_{n=1}^{\infty}$ of partial sums of the series and noting that

$$
\left|s_{n}-s_{m}\right|=\left|\sum_{i=m+1}^{n} a_{i}\right| \quad \text { for } n>m
$$

In particular, if $n=m+1$, the (6) becomes

$$
\left|a_{m+1}\right|<\epsilon
$$

for all $m>N$. This implies that $\lim _{m \rightarrow \infty} a_{m+1}=0$ and hence $\lim _{n \rightarrow \infty} a_{n}=0$. We therefore conclude the following result:

## Result:

If $\sum_{n=1}^{\infty} a_{n}$ is a convergent series, then $\lim _{n \rightarrow \infty} a_{n}=0$.

It is important here to note that the convergence of the $n$th term of series to zero as $n \rightarrow \infty$ is a necessary condition for the convergence of the series. It is not,
however, a sufficient condition, that is, if $\lim _{n \rightarrow \infty} a_{n}=0$, then it do not follow that $\sum_{n=1}^{\infty} a_{n}$ converges. It is true, however, that if $\lim _{n \rightarrow \infty} a_{n} \neq 0$, then $\sum_{n=1}^{\infty} a_{n}$ is divergent. This follows from applying the law of contraposition to the necessary condition of convergence. We conclude the following:
(a). If $a_{n} \rightarrow 0$ as $n \rightarrow \infty$, then no conclusion can be reached regarding the convergence or divergence of $\sum_{n=1}^{\infty} a_{n}$.
(b). If $a_{n} \nrightarrow 0$ as $n \rightarrow \infty$, then $\sum_{n=1}^{\infty} a_{n}$ is divergent.

## Example:

One of the simplest series is the geometric series, $\sum_{n=1}^{\infty} a^{n}$. This series is divergent if $|a| \geq 1$, since $\lim _{n \rightarrow \infty} a^{n} \neq 0$. It is convergent if $|a|<1$ by the Cauchy criterion: Let $n>M$. Then

$$
\begin{equation*}
s_{n}-s_{m}=a^{m+1}+a^{m+2}+\cdots+a^{n} . \tag{7}
\end{equation*}
$$

By multiplying the two sides of (7) by $a$, we get

$$
\begin{equation*}
a\left(s_{n}-s_{m}\right)=a^{m+2}+a^{m+3}+\cdots+a^{n+1} . \tag{8}
\end{equation*}
$$

If we now subtract (8) from (7), we obtain

$$
\begin{equation*}
s_{n}-s_{m}=\frac{a^{m+1}-a^{n+1}}{1-a} \tag{9}
\end{equation*}
$$

Since $|a|<1$, we can find an integer $N$ such that for $m>N, n>N$,

$$
\begin{aligned}
|a|^{m+1} & <\frac{\epsilon(1-a)}{2} \\
|a|^{n+1} & <\frac{\epsilon(1-a)}{2} .
\end{aligned}
$$

Hence, for a given $\epsilon>0$,

$$
\left|s_{n}-s_{m}\right|<\epsilon \quad \text { if } n>m>N
$$

Eq. (9) can actually be used to find the sum of the geometric series when $|a|<1$. Let $m=1$. By taking the limits of both sides of (9) as $n \rightarrow \infty$ we get

$$
\begin{aligned}
\lim _{n \rightarrow \infty} s_{n} & =s_{1}+\frac{a^{2}}{1-a} \quad \text { since } \lim _{n \rightarrow \infty} a^{n+1}=0 \\
& =a+\frac{a^{2}}{1-a} \\
& =\frac{a}{1-a}
\end{aligned}
$$

## Theorem:

If $\sum_{n=1}^{\infty} a_{n}$ and $\sum_{n=1}^{\infty}$ are two convergent series, and if $c$ is a constant, then the following series are also convergent:
(a). $\sum_{n=1}^{\infty}\left(c a_{n}\right)=c \sum_{n=1}^{\infty} a_{n}$.
(b). $\sum_{n=1}^{\infty}\left(a_{n}+b_{n}\right)=\sum_{n=1}^{\infty} a_{n}+\sum_{n=1}^{\infty} b_{n}$.
$\mathfrak{D e f i n i t i o n ~ ( A b s o l u t e l y ~ C o n v e r g e n t ) : ~}$
The series $\sum_{n=1}^{\infty} a_{n}$ is absolutely convergent if $\sum_{n=1}^{\infty}\left|a_{n}\right|$ is convergent.

Theorem:
Every absolutely convergent series is convergent.
$\mathfrak{D e g i n i t i o n ~ ( C o n d i t i o n a l ~ C o n v e r g e n t ) : ~}$
In the case that $\sum_{n=1}^{\infty} a_{n}$ is convergent while $\sum_{n=1}^{\infty}\left|a_{n}\right|$ is divergent, then the series $\sum_{n=1}^{\infty} a_{n}$ is said to be conditionally convergent.

### 6.2.1 Multiplication of Series

There are several ways to define the product of two series. We shall consider the so-called Cauchy's product.
$\mathfrak{D e f i n i t i o n ~ ( C a u c h y ' s ~ P r o d u c t ) : ~}$
Let $\sum_{n=0}^{\infty} a_{n}$ and $\sum_{n=0}^{\infty}$ be two series in which the summation index start at zero instead of one. Cauchy's product of these two series is the series $\sum_{n=0}^{\infty} c_{n}$, where

$$
c_{n}=\sum_{k=0}^{n} a_{k} b_{n-k}, \quad n=0,1,2, \ldots
$$

that is,

$$
\sum_{n=0}^{\infty} c_{n}=a_{0} b_{0}+\left(a_{0} b_{1}+a_{1} b_{0}\right)+\left(a_{0} b_{2}+a_{1} b_{1}+a_{2} b_{0}\right)+\cdots
$$

## Theorem:

Let $\sum_{n=0}^{\infty} c_{n}$ be the Cauchy's product of $\sum_{n=0}^{\infty} a_{n}$ and $\sum_{n=0}^{\infty} b_{n}$. Suppose that these two series are convergent and have sums equal to $s$ and $t$, respectively.
(a). If at least one of $\sum_{n=0}^{\infty} a_{n}$ and $\sum_{n=0}^{\infty} b_{n}$ converges absolutely, then $\sum_{n=0}^{\infty} c_{n}$ converges and its sum equal to $s t$.
(b). If both series are absolutely convergent, then $\sum_{n=0}^{\infty} c_{n}$ converges absolutely to the product st.

### 6.3 Sequences and Series of Functions

All the sequences and series considered thus far in this chapter had constant terms. We now extend our study to sequence and series whose terms are function of $x$.

## $\mathfrak{D e f i n i t i o n : ~}$

Let $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ be a sequence of functions defined on a set $D \subset \mathbb{R}$.
(a). If there exists a function $f(x)$ defined on $D$ such that for every $x$ in $D$,

$$
\lim _{n \rightarrow \infty} f_{n}(x)=f(x)
$$

then the sequence $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ is said to converge to $f(x)$ on $D$. Thus for a given $\epsilon>0$ there exists an inter $N$ such that $\left|f_{n}(x)-f(x)\right|<\epsilon$ if $n>N$. In general, $N$ depends on $\epsilon$ as well on $x$. In particular, if $N$ depends on $\epsilon$ but not on $x \in D$, then the sequence $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ is said to converge uniformly to $f(x)$ on $D$.
(b). If $\sum_{n=1}^{\infty} f_{n}(x)$ converges for every $x$ in $D$ to $s(x)$, then $s(x)$ is said to be the sum of the series. In this case, for a given $\epsilon>0$ there exists an inter $N$ such that

$$
\left|s_{n}(x)-s(x)\right|<\epsilon \quad \text { if } n>N,
$$

where $s_{n}(x)$ is the $n$th partial sum of the series $\sum_{n=1}^{\infty} f_{n}(x)$. The integer $N$ depends on $\epsilon$ and in general, on $x$ also. In particular, if $N$ depends on $\epsilon$ but not on $x \in D$, then the series $\sum_{n=1}^{\infty} f_{n}(x)$ is said to converge uniformly to $s(x)$ on $D$.

## Theorem:

Let $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ be a sequence of functions depends on $D \subset \mathbb{R}$ and converging to $f(x)$. Define the number $\lambda_{n}$ as

$$
\lambda_{n}=\sup _{x \in D}\left|f_{n}(x)-f(x)\right| .
$$

Then the sequence converges uniformly to $f(x)$ on $D$ iff $\lambda_{n} \rightarrow 0$ as $n \rightarrow \infty$.

## $\mathfrak{P r o o f}:$

Sufficiency: Suppose that $\lambda_{n} \rightarrow 0$ as $n \rightarrow \infty$. Then there exists an inter $N(\epsilon)$ such that

$$
\left|\lambda_{n}-0\right|<\epsilon \quad \text { if } n>N(\epsilon),
$$

i.e., for $n>N, \lambda_{n}<\epsilon$. Hence for such value of $n$,

$$
\left|f_{n}(x)-f(x)\right| \leq \lambda_{n}<\epsilon
$$

for all $x \in D$. Since $N$ depends only on $\epsilon$, the sequence $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ converges uniformly to $f(x)$ on $D$.

Necessity: Suppose that $f_{n}(x) \rightarrow f(x)$ uniformly on $D$. Then for each $\epsilon>0$ there is an $N(\epsilon)$, independent of $x$, for which

$$
\left|f_{n}(x)-f(x)\right|<\epsilon \quad \text { if } n>N .
$$

This $\epsilon$ is an upper bound for the number $\left|f_{n}(x)-f(x)\right|$. Hence the least upper bound $\lambda_{n}$ is also less equal to $\epsilon$. That is,

$$
0 \leq \lambda_{n} \leq \epsilon \quad \text { if } n>N
$$

Thus $\lambda_{n} \rightarrow 0$.

### 6.3.1 Properties of Uniformly Convergent Sequences and Series

Sequence and series of functions that are uniformly convergent have several interesting properties. We shall study some

[^21]
## Theorem:

Let $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ be uniformly convergent to $f(x)$ on a set $D$. If for each $n, f_{n}(x)$ has a limit $\tau_{n}$ as $x \rightarrow x_{0}$, where $x_{0}$ is a limit point of $D$, then the sequence $\left\{\tau_{n}\right\}_{n=1}^{\infty}$ converges to $\tau_{0}=\lim _{x \rightarrow x_{0}} f(x)$. This is equivalent to stating that

$$
\lim _{n \rightarrow \infty}\left[\lim _{x \rightarrow x_{0}} f_{n}(x)\right]=\lim _{x \rightarrow x_{0}}\left[\lim _{n \rightarrow \infty} f_{n}(x)\right] .
$$

## Corollary:

Let $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$ be a sequence of continuous functions that converges uniformly to $f(x)$ on a set $D$. Then $f(x)$ is continuous on $D$.

## $\mathfrak{P r o o f}:$

We must show that, for each $\epsilon>0$, there is a $\delta(\epsilon)$ for which

$$
\left|f(x)-f\left(x_{0}\right)\right|<\epsilon \quad \text { if }\left|x-x_{0}\right|<\delta
$$

For any $n$,

$$
\left|f(x)-f\left(x_{0}\right)\right|=\left|f(x)-f_{n}(x)+f_{n}(x)-f_{n}\left(x_{0}\right)+f_{n}\left(x_{0}\right)-f\left(x_{0}\right)\right|
$$

so

$$
\begin{equation*}
\left|f(x)-f\left(x_{0}\right)\right| \leq\left|f(x)-f_{n}(x)\right|+\left|f_{n}(x)-f_{n}\left(x_{0}\right)\right|+\left|f_{n}\left(x_{0}\right)-f\left(x_{0}\right)\right| \tag{10}
\end{equation*}
$$

By uniform convergence, there is an $N(\epsilon)$, independent of $x$, for which

$$
\begin{equation*}
\left|f_{n}(x)-f(x)\right|<\epsilon / 3, \quad \text { if } n>N, \forall x \in D \tag{11}
\end{equation*}
$$

From (10) and (11), we then have

$$
\begin{equation*}
\left|f(x)-f\left(x_{0}\right)\right| \leq \epsilon / 3+\left|f_{n}(x)-f_{n}\left(x_{0}\right)\right|+\epsilon / 3 \tag{12}
\end{equation*}
$$

But $f_{n}(x)$ is a continuous function. Hence there is a $\delta(\epsilon)$ for which

$$
\left|f_{n}(x)-f_{n}\left(x_{0}\right)\right|<\epsilon / 3 \quad \text { if }\left|x-x_{0}\right|<\delta .
$$

Putting all this together, we have

$$
\left|f(x)-f\left(x_{0}\right)\right|<\epsilon \quad \text { if }\left|x-x_{0}\right|<\delta
$$

## 7 Integration

The origin of integral calculus can be tracked back to the ancient Greeks. They are motivated by the need to measure the length of a curve, the area of a surface, or the volume of a solid. In the present chapter we shall study integration of real-valued functions of a single variable $x$ according to the concepts put forth by the German mathematician Georg Friedrich Riemann (1826-1866). He was the first to establish a rigorous analytical foundation for integration, based on the older geometric approach.

### 7.1 Some Basic Definitions

Let $f(x)$ be a function defined and bounded on a finite interval $[a, b]$.Suppose that this interval is partitioned into a finite number of subintervals by a set of points $P=\left\{x_{0}, x_{1}, \ldots, x_{n}\right\}$ such that $a=x_{0}<x_{1}<x_{2}<\cdots<x_{n}=b$. This set is called a partition of $[a, b]$. Let $\triangle x_{i}=x_{i}-x_{i-1}(i=1,2, \ldots, n)$, and $\triangle_{p}$ be the largest of $\triangle x_{1}, \triangle x_{2}, \ldots, \triangle x_{n}$. This value is called the norm of $P$ and is denoted by $\|P\|$. Consider the sum

$$
S(P, f)=\sum_{i=1}^{n} f\left(t_{i}\right) \triangle x_{i}
$$

where $t_{i}$ is a point in the subinterval $\left[x_{i-1}, x_{i}\right], i=1,2, \ldots, n$.

Definition: (Riemann Integrable):
The function $f(x)$ is said to be Riemann integrable on $[a, b]$ if a number $A$ exists with the following property: For any given $\epsilon>0$ there exists a number $\delta>0$ such that

$$
|S(P, f)-A|<\epsilon
$$

for any partition $P$ of $[a, b]$ with a norm $\|P\|<\delta$, and for any choice of the points $t_{i}$ in $\left[x_{i-1}, x_{i}\right], i=1,2, \ldots, n$. This is expressed as

$$
\lim _{\|P\| \rightarrow 0} S(P, f)=A .
$$

The number $A$ is called the Riemann integral of $f(x)$ on $[a, b]$ and is denoted by $\int_{a}^{b} f(x) d x$.

In order to investigate the existence of the Riemann integral, we shall need the following theorem:
$\mathfrak{T h e o r e m}$ (The Existence of the Riemann Integral):
Let $f(x)$ be a bounded function on a finite interval $[a, b]$. For every partition $P=$ $\left\{x_{0}, x_{1}, \ldots, x_{n}\right\}$ of $[a, b]$, Let $m_{i}$ and $M_{i}$ be, respectively, the infimum and supremum of $f(x)$ on $\left[x_{i-1}, x_{i}\right], i=1,2, \ldots, n$.
(a). If, for a given $\epsilon>0$ there exist a $\delta>0$ such that ${ }^{29}$

$$
U S_{P}(f)-L S_{P}(f)<\epsilon
$$

whenever $\|P\|<\delta$ and

$$
\begin{aligned}
& L S_{P}(f)=\sum_{i=1}^{n} m_{i} \triangle x_{i} \\
& U S_{P}(f)=\sum_{i=1}^{n} M_{i} \triangle x_{i},
\end{aligned}
$$

then $f(x)$ is Riemann integrable on $[a, b]$. Conversely (b). If $f(x)$ is Riemann integrable, then

$$
U S_{P}(f)-L S_{P}(f)<\epsilon
$$

for any partition $P$ such that $\|P\|<\delta$.

## Example:

Let $f(x):[0,1] \rightarrow \mathbb{R}$ be the function $f(x)=x^{2}$. Then, $f(x)$ is Riemann integrable on $[0,1]$. To show this, let $P=\left\{x_{0}, x_{1}, \ldots, x_{n}\right\}$ be any partition of $[0,1]$, where $x_{0}=0, x_{n}=1$. Then

$$
\begin{aligned}
L S_{P}(f) & =\sum_{i=1}^{n} x_{i-1}^{2} \triangle x_{i}, \\
U S_{P}(f) & =\sum_{i=1}^{n} x_{i}^{2} \triangle x_{i} .
\end{aligned}
$$

Hence,

$$
\begin{aligned}
U S_{P}(f)-L S_{P}(f) & =\sum_{i=1}^{n}\left(x_{i}^{2}-x_{i-1}^{2}\right) \triangle x_{i} \\
& \leq\|P\| \sum_{i=1}^{n}\left(x_{i}^{2}-x_{i-1}^{2}\right) .
\end{aligned}
$$

[^22]But

$$
\sum_{i=1}^{n}\left(x_{i}^{2}-x_{i-1}^{2}\right)=x_{n}^{2}-x_{0}^{2}=1
$$

Thus,

$$
U S_{P}(f)-L S_{P}(f) \leq\|P\| .
$$

It follows that for a given $\epsilon>0$ we can choose $\delta=\epsilon$ such that for any partition $P$ whose norm $\|P\|<\delta$,

$$
U S_{P}(f)-L S_{P}(f)<\epsilon
$$

Thus, $f(x)=x^{2}$ is Riemann integrable on $[0,1]$.

## $\mathfrak{E x a m p l e}$ :

Consider the function $f(x):[0,1] \mapsto \mathbb{R}$ such that $f(x)=0$ if $x$ is a rational number and $f(x)=1$ if $x$ is irrational. Since every (countable) subinterval of $[0,1]$ contains both rational and irrational numbers, then for any partition $P=\left\{x_{0}, x_{1}, \cdots, x_{n}\right\}$ of $[0,1]$ we have

$$
\begin{aligned}
U S_{P}(f) & =\sum_{i=1}^{n} M_{i} \triangle x_{i}=1, \\
L S_{P}(f) & =\sum_{i=1}^{n} m_{i} \triangle x_{i}=0 .
\end{aligned}
$$

It follows that

$$
\inf _{P} U S_{P}(f)=1 \quad \text { and } \quad \sup _{P} L S_{P}(f)=0
$$

Therefore $f(x)$ is not Riemann integrable on $[a, b]$.

### 7.2 Some Classes of Functions That Are Riemann Integrable

There are certain classes of functions that are Riemann integrable. Identifying a given function as a member of such a class can facilitate the determination of its Riemann integrability. Some of there classes of functions include: (a) continuous functions; (b) mon0tone function; (c) functions of bounded variation.

## Theorem:

If $f(x)$ is continuous on $[a, b]$, then it is Riemann integrable there.

## $\mathfrak{P r o o f}:$

Since $f(x)$ is continuous on a closed and bounded interval, then it must be uniformly continuous on $[a, b]$. Consequently, for a given $\epsilon>0$ there exists a $\delta>0$ that depends only on $\epsilon$ such that for any $x_{1}, x_{2}$ in $[a, b]$ we have

$$
\left|f\left(x_{1}\right)-f\left(x_{2}\right)\right|<\frac{\epsilon}{b-a}
$$

if $\left|x_{1}-x_{2}\right|<\delta$. Let $P=\left\{x_{0}, x_{1}, \ldots, x_{n}\right\}$ be a partition of $P$ with a norm $\|P\|<\delta$. Then

$$
U S_{P}(f)-L S_{P}(f)=\sum_{i=1}^{n}\left(M_{i}-m_{i}\right) \triangle x_{i},
$$

where $m_{i}$ and $M_{i}$ are, respectively, the infimum and supremum of $f(x)$ on $\left[x_{i-1}, x_{i}\right], i=$ $1,2, \ldots, n$. By Corollary 3.4. there exist points $\xi_{i}, \eta_{i}$ in $\left[x_{i-1}, x_{i}\right]$ such that $m_{i}=$ $f\left(\xi_{i}\right), M_{i}=f\left(\eta_{i}\right), i=1,2, \ldots, n$. Since $\left|\eta_{i}-\xi_{i}\right| \leq\|P\|<\delta$ for $i=1,2, \ldots, n$, then

$$
\begin{aligned}
U S_{P}(f)-L S_{P}(f) & =\sum_{i=1}^{n}\left[f\left(\eta_{i}\right)-f\left(\xi_{i}\right)\right] \triangle x_{i} \\
& <\frac{\epsilon}{b-a} \sum_{i=1}^{n} \triangle x_{i}=\epsilon .
\end{aligned}
$$

By Theorem of the existence of Riemann Integral, we conclude that $f(x)$ is Riemann integrable on $[a, b]$.

### 7.3 Properties of the Riemann Integral

The Riemann integral has several properties that are useful at both the theoretical and practical levels. The definition of the Riemann integral of $f$ on the interval $[a, b]$ specifies that $a<b$. Now, however, it is convenient to extend the definition to cover cases in which $a=b$ or $a>b$. Geometrically, the following two definitions seem reasonable. For instance, it makes sense to define the area of a region of zero width and finite height to be 0 .

## $\mathfrak{D e f i n i t i o n : ~}$

(a). If $f$ is defined at $x=a$, then $\int_{a}^{a} f(x) d x=0$.
(b). If $f$ is integrable on $[a, b]$, then $\int_{b}^{a} f(x) d x=-\int_{a}^{b} f(x) d x$.

## Theorem:

If $f(x)$ and $g(x)$ are Riemann integrable on $[a, b]$ and if $c_{1}$ and $c_{2}$ are constants, then $c_{1} f(x)+c_{2} g(x)$ is Riemann integrable on $[a, b]$, and

$$
\int_{a}^{b}\left[c_{1} f(x)+c_{2} g(x)\right] d x=c_{1} \int_{a}^{b} f(x) d x+c_{2} \int_{a}^{b} g(x) d x
$$

## Theorem:

If $f(x)$ is Riemann integrable on $[a, b]$ and if $a<c<b$, then

$$
\int_{a}^{b} f(x) d x=\int_{a}^{c} f(x) d x+\int_{c}^{b} f(x) d x
$$

Threoem:
If $f(x)$ and $g(x)$ are Riemann integrable on $[a, b]$, then so is their product $f(x) g(x)$.

Theorem:
If $f(x)$ is Riemann integrable on $[a, b]$, and $m \leq f(x) \leq M$ for all $x$ in $[a, b]$, then

$$
m(b-a) \leq \int_{a}^{b} f(x) d x \leq M(b-a)
$$

Theorem (The Mean Value Theorem for Integrals):
If $f(x)$ is continuous on $[a, b]$, then there exist a point $c \in[a, b]$ such that

$$
\int_{a}^{b} f(x) d x=(b-a) f(c) .
$$

## $\mathfrak{P r o o f}:$

By Theorem above we have

$$
m \leq \frac{1}{b-a} \int_{a}^{b} f(x) d x \leq M
$$

where $m$ and $M$ are, respectively, the infimum and supremum of $f(x)$ on $[a, b]$. By the intermediate- value theorem, there is a point $c \in[a, b]$ such that

$$
f(c)=\frac{1}{b-a} \int_{a}^{b} f(x) d x
$$

$\mathfrak{D e f i n i t i o n ~ ( I n d e f i n i t e ~ I n t e g r a l ) : ~}$
Let $f(x)$ be Riemann integrable on $[a, b]$. The function

$$
F(x)=\int_{a}^{x} f(t) d t, \quad a \leq x \leq b
$$

is called an indefinite integral of $f(x)$.

## Theorem:

If $f(x)$ is Riemann integrable on $[a, b]$, then $F(x)=\int_{a}^{x} f(t) d t$ is uniformly continuous on $[a, b]$.

## $\mathfrak{P r o o f}:$

Let $x_{1}, x_{2}$ be in $[a, b], x_{1}<x_{2}$. Then

$$
\begin{aligned}
\left|F\left(x_{2}\right)-F\left(x_{1}\right)\right| & =\left|\int_{a}^{x_{2}} f(t) d t-\int_{a}^{x_{1}} f(t) d t\right| \\
& =\left|\int_{x_{1}}^{x_{2}} f(t) d t\right| \\
& \leq \int_{x_{1}}^{x_{2}} \mid f(t) d t \\
& \leq M^{\prime}\left(x_{2}-x_{1}\right)
\end{aligned}
$$

where $M^{\prime}$ is the supremum of $|f(x)|$ on $[a, b]$. Thus if $\epsilon>0$ is given, then $\mid F\left(x_{2}\right)-$ $F\left(x_{1}\right) \mid<\epsilon$ provided that $\left|x_{2}-x_{1}\right|<\epsilon / M^{\prime}$. This provides uniform continuity of $F(x)$ on $[a, b]$.
$\mathfrak{T h e o r e m}$ (Fundamental Theorem of Calculus):
Suppose that $f(x)$ is continuous on $[a, b]$. Let $F(x)=\int_{a}^{x} f(t) d t$. Then we have the following:
(a). $\frac{d F(x)}{d x}=f(x), a \leq x \leq b$.
(b). $\int_{a}^{b} f(x) d x=G(b)-G(a)=F(b)-F(a)$, where $G(x)=F(x)+c$, and $c$ is an arbitrary constant.

## $\mathfrak{P r o o f}:$

We have

$$
\begin{aligned}
\frac{d F(x)}{d x} & =\frac{d}{d x} \int_{a}^{x} f(t) d t=\lim _{h \rightarrow 0} \frac{1}{h}\left[\int_{a}^{x+h} f(t) d t-\int_{a}^{x} f(t) d t\right] \\
& =\lim _{h \rightarrow 0} \frac{1}{h} \int_{x}^{x+h} f(t) d t \\
& =\lim _{h \rightarrow 0} f(x+\theta h), \quad \text { (by mean value theorem for integral) }
\end{aligned}
$$

where $0 \leq \theta \leq 1$. Hence,

$$
\frac{d F(x)}{d x}=\lim _{h \rightarrow 0} f(x+\theta h)=f(x)
$$

by the continuity of $f(x)$.

To prove the second part of the theorem, let $G(x)$ be defined on $[a, b]$ as

$$
G(x)=F(x)+c=\int_{a}^{x} f(t) d t+c,
$$

that is, $G(x)$ is an indefinite integral of $f(x)$. If $x=a$, then $G(a)=c$ since $F(a)=0$. Also, if $x=b$, then $G(b)=F(b)+c=\int_{a}^{b} f(t) d t+G(a)$. It follows that

$$
\int_{a}^{b} f(t) d t=G(b)-G(a)=F(b)+c-F(a)-c=F(b)-F(a) .
$$

## Example:

Evaluate $\int_{0}^{2}\left(2 x^{2}-3 x+2\right) d x$.

## Solution:

$$
\begin{aligned}
\int_{0}^{2}\left(2 x^{2}-3 x+2\right) d x & =\left[\frac{2 x^{3}}{3}-\frac{3 x^{2}}{2}+2 x\right]_{0}^{2} \\
& =\left(\frac{16}{3}-6+4\right)-0 \\
& =\frac{10}{3}
\end{aligned}
$$

## Example:

Evaluate $\frac{d}{d x}\left[\int_{0}^{x} \sqrt{t^{2}+1} d t\right]$.

## Solution:

Note that $f(t)=\sqrt{t^{2}+1}$ is continuous on the entire real line. So using the fundamental theorem of calculus, we can write

$$
\frac{d}{d x}\left[\int_{0}^{x} \sqrt{t^{2}+1} d t\right]=\sqrt{x^{2}+1}
$$

### 7.4 Change of Variables in Riemann Integration

There are situations in which the variable $x$ in a Riemann integral is a function of some other variable, sat $u$. In this case, it may be of interest to determine how the integral can be expressed and evaluated under the given transformation.
$\mathfrak{T h e o r e m}$ (Integration by Substitution):
Let $f(x)$ be continuous on $[\alpha, \beta]$, and let $x=g(u)$ be a function whose derivative $g^{\prime}(u)$ exists and is continuous on $[c, d]$. Suppose that the range of $g$ is contained inside $[\alpha, \beta]$. If $a, b$ are points in $[\alpha, \beta]$ such that $a=g(c)$ and $b=g(d)$, then

$$
\int_{a}^{b} f(x) d x=\int_{c}^{d} f[g(u)] g^{\prime}(u) d u
$$

$\mathfrak{P r o o f}:$

Let $F(x)=\int_{a}^{x} f(t) d t$, then $F^{\prime}(x)=f(x)$. According to the chain rule

$$
\begin{aligned}
\frac{d F(g(u))}{d u} & =\frac{d F(g(u))}{d g(u)} \frac{d g(u)}{d u} \\
& =f(g(u)) g^{\prime}(u)
\end{aligned}
$$

By the fundamental theorem,

$$
\begin{aligned}
\int_{c}^{d} f(g(u)) g^{\prime}(u) d u & =\left.F(g(u))\right|_{c} ^{d} \\
& =F(g(d))-F(g(c)) \\
& =F(b)-F(a) \\
& =\int_{a}^{b} f(x) d x
\end{aligned}
$$

## $\mathfrak{E x a m p l e}$ :

Evaluate $\int_{0}^{1} u\left(u^{2}+1\right)^{3} d u$.

## Solution:

To evaluating this integral, let $x=g(u)=\left(u^{2}+1\right), g^{\prime}(u)=2 u$ and $f(x)=x^{3}$. So $g(0)=1$ and $g(1)=2$. Then

$$
\begin{aligned}
\int_{0}^{1} u\left(u^{2}+1\right)^{3} d u & =\frac{1}{2} \int_{0}^{1} g^{\prime}(u) f(g(u)) d u \\
& =\frac{1}{2} \int_{1}^{2} x^{3} d x \\
& =\frac{15}{8}
\end{aligned}
$$

Theorem (Integration by Parts):
If $f$ and $g$ are differentiable on $[a, b]$, and if $f^{\prime}$ and $g^{\prime}$ are integrable there, then

$$
\int_{a}^{b} f(x) g^{\prime}(x) d x=\left.f(x) g(x)\right|_{a} ^{b}-\int_{a}^{b} g(x) f^{\prime}(x) d x
$$

## $\mathfrak{P r o o f}:$

We have the familiar formula

$$
(f g)^{\prime}=f g^{\prime}+f^{\prime} g .
$$

Integrating this formula

$$
\left.f(x) g(x)\right|_{a} ^{b}=\int_{a}^{b}[f(x) g(x)]^{\prime} d x=\int_{a}^{b} f(x) g^{\prime}(x) d x+\int_{a}^{b} f^{\prime}(x) g(x) d x
$$

## Example:

Evaluate $\int x^{2} \ln x d x$.

## Solution:

In this case, $x^{2}$ is more easily integrated than $\ln x$, and the derivative of $\ln x$ is simpler than $\ln x$. There we set $g(x)=\ln x, g^{\prime}(x)=\frac{1}{x}$ and $f^{\prime}(x)=x^{2}, f(x)=1 / 3 x^{3}$. So

$$
\begin{aligned}
\int x^{2} \ln x d x=\int f^{\prime}(x) g(x) d x & =f(x) g(x)-\int f(x) g^{\prime}(x) d x \\
& =\frac{1}{3} x^{3} \ln x-\frac{1}{3} \int x^{3} \frac{1}{x} d x \\
& =\frac{1}{3} x^{3} \ln x-\frac{1}{3} \int x^{2} d x \\
& =\frac{1}{3} x^{3} \ln x-\frac{x^{3}}{9}+C .
\end{aligned}
$$

### 7.5 Improper Riemann Integrals

In our study of the Riemann integral we have only considered integrals of functions that are bounded (talking about $f(x)$ ) on a finite interval $[a, b]$. We now extend the scope of Riemann integration to include situations where the integrand can became unbounded at one or more points inside the range of integration, which can also be infinite. In such situations, the Riemann integral is called an improper integral.

There are two kinds of improper integrals. If $f(x)$ is Riemann integrable on $[a, b]$ for any $b>a$, then $\int_{a}^{\infty} f(x) d x$ is called an improper integral of the first kind, where
the range of integration if infinite. If, however, $f(x)$ becomes infinite at a finite number of points inside the range of integration, then the integral $\int_{a}^{b} f(x) d x$ is said to be improper of the second kind.

### 7.5.1 Improper Riemann Integrals of the First Kind

Definition:
Let $F(z)=\int_{a}^{z} f(x) d x$. Suppose that $F(z)$ exists for any value of $z$ greater than $a$. If $F(z)$ has a finite limit $L$ as $z \rightarrow \infty$, then the improper integral $\int_{a}^{\infty} f(x) d x$ is said to converge to $L$. In this case, $L$ represent the Riemann integral of $f(x)$ on $[a, \infty]$ and we write

$$
L=\int_{a}^{\infty} f(x) d x
$$

On the other hand, if $L= \pm \infty$, then the improper integral $\int_{a}^{\infty} f(x) d x$ is said to diverge. By the same token, we can define the integral $\int_{-\infty}^{a} f(x) d x$ as the limit, if it exists, of $\int_{-z}^{a} f(x) d x$ as $z \rightarrow \infty$. Also, $\int_{-\infty}^{\infty} f(x) d x$ is defined as

$$
\int_{-\infty}^{\infty} f(x) d x=\lim _{u \rightarrow \infty} \int_{-u}^{a} f(x) d x+\lim _{z \rightarrow \infty} \int_{a}^{z} f(x) d x
$$

where $a$ is any finite number, provided that both limit exist.
$\mathfrak{E x a m p l e}$ (An improper integral that diverges):
Evaluate $\int_{1}^{\infty} \frac{1}{x} d x$.

## Solution:

$$
\begin{aligned}
\int_{1}^{\infty} \frac{1}{x} d x & =\lim _{z \rightarrow \infty} \int_{1}^{z} \frac{1}{x} d x \\
& =\lim _{z \rightarrow \infty}[\ln x]_{1}^{z} \\
& =\lim _{z \rightarrow \infty}(\ln z-0) \\
& =\infty
\end{aligned}
$$

Example (An improper integral that converges):
Evaluate $\int_{0}^{\infty} e^{-x} d x$.

## Solution:

$$
\begin{aligned}
\int_{0}^{\infty} e^{-x} d x & =\lim _{z \rightarrow \infty} \int_{0}^{z} e^{-x} d x \\
& =\lim _{z \rightarrow \infty}\left[-e^{-x}\right]_{0}^{z} \\
& =\lim _{z \rightarrow \infty}\left(-e^{-z}+1\right) \\
& =1
\end{aligned}
$$

## Exercise:

Evaluate $\int_{1}^{\infty}(1-x) e^{-x} d x$.

### 7.5.2 Improper Riemann Integrals of the Second Kind

Let us now consider integrals of the form $\int_{a}^{b} f(x) d x$ where $[a, b]$ is a finite interval and the integrand become infinite at a finite number points inside $[a, b]$. Such integrals are called improper integrals of the second kind.
$\mathfrak{D e f i n i t i o n ~ ( E n d ~ p o i n t ) : ~}$
Suppose that $f(x) \rightarrow \infty$ as $x \rightarrow a^{+}$, then $\int_{a}^{b} f(x) d x$ is said to be converge if the limit

$$
\lim _{\epsilon \rightarrow 0^{+}} \int_{a+\epsilon}^{b} f(x) d x
$$

exists and is finite.
$\mathfrak{D e f i n i t i o n ~ ( E n d ~ p o i n t ) : ~}$
Suppose that $f(x) \rightarrow \infty$ as $x \rightarrow b^{-}$, then $\int_{a}^{b} f(x) d x$ is said to be converge if the limit

$$
\lim _{\epsilon \rightarrow 0^{+}} \int_{a}^{b-\epsilon} f(x) d x
$$

exists and is finite.
$\mathfrak{D e f i n i t i o n ~ ( I n t e r i o r ~ p o i n t ) : ~}$
Suppose that $f(x) \rightarrow \infty$ as $x \rightarrow c$, where $a<c<b$, then $\int_{a}^{b} f(x) d x$ is the sum of $\int_{a}^{c} f(x) d x$ and $\int_{c}^{b} f(x) d x$ provided that both integrals converges.

## $\mathfrak{D e f i n i t i o n : ~}$

If $f(x) \rightarrow \infty$ as $x \rightarrow x_{0}$, where $x_{0} \in[a, b]$, then $x_{0}$ is said to be a singularity of $f(x)$.

## Example:

Evaluate $\int_{0}^{1} \frac{1}{\sqrt[3]{x}} d x$.

## Solution:

The integral is improper of the second since $\frac{1}{\sqrt[3]{x}} \rightarrow \infty$ as $x \rightarrow 0^{+}$. Thus, $x=0$ is a singularity of the integrand. We can evaluate this integral as follows.

$$
\begin{aligned}
\int_{0}^{1} x^{-1 / 3} d x & =\lim _{a \rightarrow 0^{+}}\left[\frac{x^{2 / 3}}{2 / 3}\right]_{a}^{1} \\
& =\lim _{a \rightarrow 0^{+}} \frac{3}{2}\left(1-a^{2 / 3}\right) \\
& =\frac{3}{2}
\end{aligned}
$$

## $\mathfrak{E x a m p l e}$ :

Evaluate $\int_{0}^{2} \frac{1}{x^{3}} d x$.

## Solution:

The integral is improper of the second since $\frac{1}{x^{3}} \rightarrow \infty$ as $x \rightarrow 0^{+}$. Thus, $x=0$ is a singularity of the integrand. We can evaluate this integral as follows.

$$
\begin{aligned}
\int_{0}^{2} \frac{1}{x^{3}} d x & =\lim _{a \rightarrow 0^{+}}\left[-\frac{1}{2 x^{2}}\right]_{a}^{2} \\
& =\lim _{a \rightarrow 0^{+}}\left(\frac{1}{8}+\frac{1}{2 a^{2}}\right) \\
& =\infty
\end{aligned}
$$

## $\mathfrak{E x a m p l e}$ :

Evaluate $\int_{-1}^{2} \frac{1}{x^{3}} d x$.

## Solution:

The integral is improper of the second since $\frac{1}{x^{3}} \rightarrow \infty$ as $x \rightarrow 0^{+}$. Thus, $x=0$ is a singularity of the integrand. We can evaluate this integral as follows.

$$
\int_{-1}^{2} \frac{1}{x^{3}} d x=\int_{-1}^{0} \frac{1}{x^{3}} d x+\int_{0}^{2} \frac{1}{x^{3}} d x
$$

From the example above we know that the second integral diverge. Therefore, the original improper integral also diverges.

## Note:

Remember to check for singularity point at interior points as well as endpoints when determining whether an integral is improper. For instance, if you do not recognized that the integral in the above example $\int_{-1}^{2} \frac{1}{x^{3}} d x$ was improper, you would have obtained the incorrect results

$$
\int_{-1}^{2} \frac{1}{x^{3}} d x=\left[-\frac{1}{2 x^{2}}\right]_{-1}^{2}=-\frac{1}{8}+\frac{1}{2}=\frac{3}{8} . \quad \text { (incorrect evaluation) }
$$

## $\mathfrak{E x a m p l e}$ :

Consider the integral $\int_{0}^{2}\left(x^{2}-3 x+1\right) /\left[x(x-1)^{2}\right] d x$. Here, the integrand has two singularities, namely $x=0$ and $x=1$, inside $[0,2]$. We can therefore write

$$
\begin{aligned}
\int_{0}^{2} \frac{x^{2}-3 x+1}{x(x-1)^{2}} d x= & \lim _{t \rightarrow 0^{+}} \int_{t}^{1 / 2} \frac{x^{2}-3 x+1}{x(x-1)^{2}} d x \\
& +\lim _{u \rightarrow 1^{-}} \int_{1 / 2}^{u} \frac{x^{2}-3 x+1}{x(x-1)^{2}} d x \\
& +\lim _{v \rightarrow 1^{+}} \int_{v}^{2} \frac{x^{2}-3 x+1}{x(x-1)^{2}} d x
\end{aligned}
$$

We note that

$$
\frac{x^{2}-3 x+1}{x(x-1)^{2}}=\frac{1}{x}-\frac{1}{(x-1)^{2}} .
$$

Hence,

$$
\begin{aligned}
\int_{0}^{2} \frac{x^{2}-3 x+1}{x(x-1)^{2}} d x= & \lim _{t \rightarrow 0^{+}}\left[\log x+\frac{1}{x-1}\right]_{t}^{1 / 2} \\
& +\lim _{u \rightarrow 1^{-}}\left[\log x+\frac{1}{x-1}\right]_{1 / 2}^{u} \\
& +\lim _{v \rightarrow 1^{+}}\left[\log x+\frac{1}{x-1}\right]_{v}^{2}
\end{aligned}
$$

None of the above limits exists as a finite number. The integral is therefore divergent.

### 7.6 Convergence of a Sequence of Riemann Integrals

In this section we confine our attention to the limiting behavior of the integrals of a sequence of functions $\left\{f_{n}(x)\right\}_{n=1}^{\infty}$.

## Theorem:

Suppose that $f_{n}(x)$ is Riemann integrable on $[a, b]$ for $n \geq 1$. If $f_{n}(x)$ converges uniformly to $f(x)$ on $[a, b]$ as $n \rightarrow \infty$, then $f(x)$ is Riemann integrable on $[a, b]$ and

$$
\lim _{n \rightarrow \infty} \int_{a}^{b} f_{n}(x) d x=\int_{a}^{b} f(x) d x
$$

$\mathfrak{P r o o f}:$
Let us now show the part of

$$
\lim _{n \rightarrow \infty} \int_{a}^{b} f_{n}(x) d x=\int_{a}^{b} f(x) d x
$$

Let $\epsilon>0$ be given. Since $f_{n}(x)$ converges uniformly to $f(x)$, then there exists an integer $n_{0}$ that depends only on $\epsilon$ such that

$$
\left|f_{n}(x)-f(x)\right|<\frac{\epsilon}{b-a}
$$

From the basic properties of Riemann integral, we have

$$
\begin{aligned}
& \left|\int_{a}^{b}\left[f_{n}(x)-f(x)\right] d x\right|=\left|\int_{a}^{b} f_{n}(x) d x-\int_{a}^{b} f(x) d x\right| \\
& \leq \int_{a}^{b}\left|f_{n}(x)-f(x)\right| d x \\
& \leq \int_{a}^{b} \frac{\epsilon}{b-a}=\epsilon
\end{aligned}
$$

and the result follows, since $\epsilon$ is an arbitrary positive number.

### 7.7 Riemann-Stieltjes Integral

The Riemann-Stieltjes integral involves two functions $f(x)$ and $g(x)$, both defined on the interval $[a, b]$, and is denoted by $\int_{a}^{b} f(x) d g(x)$. In particular, if $g(x)=x$ we obtain the Riemann integral $\operatorname{int}_{a}^{b} f(x) d x$. Thus the Riemann integral is the special case of the Riemann-Stieltjes integral
$\mathfrak{D e f i n i t i o n ~ ( T h e ~ R i e m a n n - S t i e l t j e s ~ i n t e g r a l ) : ~}$ If $f(x)$ is bounded on $[a, b]$, if $g(x)$ is monotone increasing on $[a, b]$, and if $P=$ $\left\{x_{0}, x_{1}, \ldots, x_{n}\right\}$ is a partition of $[a, b]$, we define the sums

$$
\begin{aligned}
L S_{P}(f, g) & =\sum_{i=1}^{n} m_{i} \triangle g_{i} \\
U S_{P}(f, g) & =\sum_{i=1}^{n} M_{i} \triangle g_{i}
\end{aligned}
$$

where $m_{i}$ and $M_{i}$ are, respectively, the infimum and supremum of $f(x)$ on $\left[x_{i-1}, x_{i}\right]$, $\triangle g_{i}=g\left(x_{i}\right)-g\left(x_{i-1}\right), i=1,2, \ldots, n$. If for a given $\epsilon>0$ there exists a $\delta>0$ such that

$$
U S_{P}(f, g)-L S_{P}(f, g)<\epsilon
$$

whenever $\triangle_{P}<\delta$, where $\triangle_{P}$ is the norm of $P$, then $f(x)$ is said to be RiemannStieltjes integrable with respect to $g(x)$ on $[a, b]$. In this case,

$$
\int_{a}^{b} f(x) d g(x)=\inf _{P} U S_{P}(f, g)=\sup _{P} L S_{P}(f, g)
$$

Equivalently, suppose that for a given partition $P=\left\{x_{0}, x_{1}, \cdots, x_{n}\right\}$ we define the sum

$$
S(P, f, g)=\sum_{i=1}^{n} f\left(t_{i}\right) \triangle g_{i}
$$

where $t_{i}$ is a point in the interval $\left[x_{i-1}, x_{i}\right], i=1,2, . ., n$. Then $f(x)$ is RiemannStieltjes integrable with respect to $g(x)$ on $[a, b]$ if for any $\epsilon>0$ there exists a $\delta>0$ such that

$$
\left|S(P, f, g)-\int_{a}^{b} f(x) d g(x)\right|<\epsilon
$$

for any partition $P$ on $[a, b]$ with a norm $\triangle_{p}<\delta$, and for any choice of the point $t_{i}$ in $\left[x_{i-1}, x_{i}\right], i=1,2, . ., n$.

The next theorem shows that under certain conditions, the Riemann-Stieltjes integral reduces to the Riemann integral.

## Theorem:

Suppose that $f(x)$ is Riemann-Stieltjes integrable with respect to $g(x)$ on $[a, b]$, where $g(x)$ has continuous derivative $g^{\prime}(x)$ on $[a, b]$. Then

$$
\int_{a}^{b} f(x) d g(x)=\int_{a}^{b} f(x) g^{\prime}(x) d x
$$

## $\mathfrak{E x a m p l e}$ :

Under this theorem, it is easy to see that if $f(x)=1$ and $g(x)=x^{2},{ }^{30}$ then $\int_{a}^{b} f(x) d g(x)=\int_{a}^{b} f(x) g^{\prime}(x) d x=\int_{a}^{b} 2 x d x=b^{2}-a^{2}$.

It is possible, however, for the Riemann-Stieltjes integral to exist even if $g(x)$ is a discontinuous function.

## $\mathfrak{T h e o r e m}$ :

Let $g(x)$ be a step function defined on $[a, b]$ with jump discontinuities at $x=$

[^23]$c_{1}, c_{2}, \ldots, c_{n}$ and $a, c_{1}<c_{2}<\ldots<c_{n}=b$, such that
\[

g(x)=\left\{$$
\begin{array}{cc}
\lambda_{1}, & a \leq x<c_{1} \\
\lambda_{2}, & c_{1} \leq x<c_{2} \\
\cdot & \cdot \\
\cdot & \cdot \\
, & \cdot \\
\lambda_{n}, & c_{n-1} \leq x<c_{n} \\
\lambda_{n+1}, & x=c_{n}
\end{array}
$$\right.
\]

If $f(x)$ is bounded on $[a, b]$ and continuous at $x=c_{1}, c_{2}, \ldots, c_{n}$, then

$$
\int_{a}^{b} f(x) d g(x)=\sum_{i=1}^{n}\left(\lambda_{i+1}-\lambda_{i}\right) f\left(c_{i}\right)
$$

## $\mathfrak{E x a m p l e}$ :

The mathematical expectation of a random variable $X$ is defined by

$$
E(X)=\int x d F(x)
$$

where $F(x)$ is the c.d.f $\mathrm{pf} X$ and the integral is the Riemann-Stieltjes integrals. When $X$ has a probability density function $f(x)$, we see that

$$
E(x)=\int x d F(x)=\int x f(x) d x
$$

reduce to a Riemann integral. For a discrete probability distribution with $F(x)$ as a step function,

$$
E(X)=\int x d F(x)=\sum p_{r} x_{r}
$$

where $x_{r}$ is a discontinuity point of $F$ and $p_{r}$ is the saltus at $x_{r}$. It is interesting to note that in all the three cases the same symbol $\int x d F(x)$ could be employed.

## 8 Multivariate Calculus

In this section we extend the notions of limits, continuity, differentiation, and integration to multivariate functions, that us functions of several variables. These functions can be real-valued or possibly vector- valued. More specifically, if $\mathbb{R}^{n}$ denotes the $n$-dimensional Euclidean space, $n \geq 1$, then we shall in general consider functions defined on a set $D \subset \mathbb{R}^{n}$ and have values in $\mathbb{R}^{m}, m \geq 1$. Such functions are represented symbolically as $\mathbf{f}: D \rightarrow \mathbb{R}^{m}$, where for $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime} \in D$,

$$
\mathbf{f}(\mathbf{x})=\left[f_{1}(\mathbf{x}), f_{2}(\mathbf{x}), \ldots, f_{m}(\mathbf{x})\right]^{\prime}
$$

and $f_{i}(\mathbf{x})$ is a real-valued function of $x_{1}, x_{2}, \ldots, x_{n}(i=1,2, \ldots, m)$.

### 8.1 Some Basic Definition

In this section we extend these concept of one-dimensional Euclidean space to higher dimensional Euclidean spaces.
$\mathfrak{D e f i n i t i o n ~ ( E u c l i d e a n ~ N o r m ) : ~}$
Any point $\mathbf{x}$ in $\mathbb{R}^{n}$ can be represented as a column vector of the form $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime}$, where $x_{i}$ is the $i$ th element of $\mathbf{x}(i=1,2, \ldots, n)$. The Euclidean norm of $\mathbf{x}$ was defined as $\|\mathbf{x}\|=\left(\sum_{i=1}^{n} x_{i}^{2}\right)^{1 / 2}$.
$\mathfrak{D e f i n i t i o n ~ ( N e i g h b o r h o o d ~ o f ~ a ~ p o i n t ~} \mathbf{x}_{0}$ in $\mathbb{R}^{n}$ ):
Let $\mathbf{x}_{0} \in \mathbb{R}^{n}$. A neighborhood $N_{r}\left(\mathbf{x}_{0}\right)$ of $\mathbf{x}_{0}$ is a set of points in $\mathbb{R}^{n}$ that lie within some distance, say $r$, from $\mathbf{x}_{0}$, that is

$$
N_{r}\left(\mathbf{x}_{0}\right)=\left\{\mathbf{x} \in \mathbb{R}^{n} \mid\left\|\mathbf{x}-\mathbf{x}_{0}\right\|<r\right\} .
$$

Thus a neighborhood of $\mathbf{x}_{0}$ is inside a circle centered at $\mathbf{x}_{0} .{ }^{31}$ If $\mathbf{x}_{0}$ is deleted from $N_{r}\left(\mathbf{x}_{0}\right)$, we obtain the so-called deleted neighborhood of $\mathbf{x}_{0}$, which we denote by $N_{r}^{d}\left(\mathbf{x}_{0}\right)$.

[^24]$\mathfrak{D e f i n i t i o n ~ ( L i m i t ~ p o i n t ~ o f ~ a ~ s e t ) : ~}$
A point $\mathbf{x}_{0}$ in $\mathbb{R}^{n}$ is a limit point of a set $A \subset \mathbb{R}^{n}$ if every neighborhood of $\mathbf{x}_{0}$ contains an element $\mathbf{x}$ of $A$ such that $\mathbf{x} \neq \mathbf{x}_{0}$, that is, every deleted neighborhood of $\mathbf{x}_{0}$ contains points of $A$.
$\mathfrak{D e f i n i t i o n ~ ( C l o s e d ~ S e t ) : ~}$
A set $A \subset \mathbb{R}^{n}$ is closed if every limit point of $A$ belongs to $A$.
$\mathfrak{D e f i n i t i o n ~ ( I n t e r i o r ) : ~}$
A point $\mathbf{x}_{0}$ is $\mathbb{R}^{n}$ is an interior of a set $A \subset \mathbb{R}^{n}$ if there exists an $r>0$ such that $N_{r}\left(\mathbf{x}_{0}\right) \subset A$.
$\mathfrak{D e f i n i t i o n ~ ( O p e n ~ S e t ) : ~}$
A set $A \subset \mathbb{R}^{n}$ is open if for every point $\mathbf{x}$ in $A$ there exists a neighborhood $N_{r}(\mathbf{x})$ that is contained in $A$. Thus, $A$ is open if it contains entirely of interior points.
$\mathfrak{D e f i n i t i o n ~ ( B o u n d a r y ~ P o i n t ) : ~}$
A point $p \in \mathbb{R}^{n}$ is a boundary point of a set $A \subset \mathbb{R}^{n}$ if every neighborhood of $p$ contains points of $A$ as well as points of $A^{c}$, the complement of $A$ with respect to $\mathbb{R}^{n}$. The set of all boundary points of $A$ is called its boundary and is denoted by $B r(A)$.
$\mathfrak{D e f i n i t i o n ~ ( B o u n d e d ) : ~}$
A set $A \subset \mathbb{R}^{n}$ is bounded if there exists an $r>0$ such that $\|\mathbf{x}\| \leq r \forall \mathbf{x} \in A$.

## Definition:

Let $\left\{\mathbf{a}_{i}\right\}_{i=1}^{\infty}$ represent a sequence of point in $\mathbb{R}^{n}$. Then $\left\{\mathbf{a}_{i}\right\}_{i=1}^{\infty}$ converges to a point $\mathbf{c} \in \mathbb{R}^{n}$ if for a given $\epsilon>0$ there exists an integer $N$ such that $\left\|\mathbf{a}_{i}-\mathbf{c}\right\|<\epsilon$ whenever $i>N$. This is written as $\lim _{i \rightarrow \infty} \mathbf{a}_{i}=\mathbf{c}$, or $\mathbf{a}_{i} \rightarrow \mathbf{c}$ as $i \rightarrow \infty$.
$\mathfrak{D e f i n i t i o n : ~}$
A sequence $\left\{\mathbf{a}_{i}\right\}_{i=1}^{\infty}$ is bounded if there exist a number $K>0$ such that $\left\|\mathbf{a}_{i}\right\| \leq K$ $\forall i$.

### 8.2 Limits of a Multivariate Function

For a function of several variables, say $x_{1}, x_{2}, \ldots, x_{n}$, its limit at a point $\mathbf{a}=\left(a_{1}, a_{2}, \ldots, a_{n}\right)^{\prime}$ is considered when $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime}$ approach $\mathbf{a}$ in any possible way. Thus when $n>1$ there are infinitely many ways in which $\mathbf{x}$ can approach $\mathbf{a}$.
$\mathfrak{D e f i n i t i o n : ~}$
Let $\mathbf{f}: D \rightarrow \mathbb{R}^{n}$. The $\mathbf{f}(\mathbf{x})$ is said to have a limit $\mathbf{L}=\left(L_{1}, L_{2}, \ldots, L_{m}\right)^{\prime}$ as $\mathbf{x}$ approach $\mathbf{a}$, written symbolically as $\mathbf{x} \rightarrow \mathbf{a}$, where $\mathbf{a}$ is a limit point of $D$, if for a given $\epsilon>0$, there exists a $\delta>0$ such that $\|\mathbf{f}(\mathbf{x})-\mathbf{L}\|<\epsilon$ for all $\mathbf{x}$ in $D \cap N_{\delta}^{d}(\mathbf{a})$, where $N_{\delta}^{d}(\mathbf{a})$ is a deleted neighborhood of $\mathbf{a}$ of radius $\delta$. If it exist, this limit is written symbolically as $\lim _{\mathbf{x} \rightarrow \mathbf{a}} \mathbf{f}(\mathbf{x})=\mathbf{L}$.

Note that whether a limit $\mathbf{f}(\mathbf{x})$ exists, its value must be the same no matter how $\mathbf{x}$ approach $\mathbf{a}$. It is important here to understand the meaning of " $\mathbf{x}$ approach $\mathbf{a}$ ". By this we do not necessarily mean that $\mathbf{x}$ moves along a straight line leading to a. Rather we mean that $\mathbf{x}$ moves closer and closer to $\mathbf{a}$ along any curve that goes through a.

## $\mathfrak{E x a m p l e}:$

Let $f(\mathbf{x})=x^{2} y^{2} /\left(x^{2}+y^{2}\right)$. The domain of this function is the whole plane except the origin. We take $\mathbf{a}=(0,0)^{\prime}$ and investigate the limit of $f(\mathbf{x})$ as $\mathbf{x} \rightarrow \mathbf{a}$.

Let $\epsilon>0$ be given. Then we want to show that there is a $\delta$ such that

$$
\begin{equation*}
|f(\mathbf{x})-0|<\epsilon \quad \text { whenever } \quad\|\mathbf{x}-\mathbf{a}\|=\sqrt{(x-0)^{2}+(y-0)^{2}}<\delta \tag{13}
\end{equation*}
$$

Now clearly

$$
\begin{aligned}
x^{2} & \leq x^{2}+y^{2} \text { and } \\
y^{2} & \leq x^{2}+y^{2}
\end{aligned}
$$

so that

$$
|f(\mathbf{x})|=\frac{x^{2} y^{2}}{x^{2}+y^{2}} \leq \frac{\left(x^{2}+y^{2}\right)^{2}}{x^{2}+y^{2}}=\left(x^{2}+y^{2}\right)<\varepsilon
$$

if

$$
\|\mathbf{x}-\mathbf{a}\|=\sqrt{(x-0)^{2}+(y-0)^{2}}<\sqrt{\epsilon}
$$

Thus we have satisfied (13) with $\delta=\sqrt{\epsilon}$, and we see that

$$
\lim _{x \rightarrow 0} f(x)=0 .
$$

## Example:

Consider the behavior of function

$$
f\left(x_{1}, x_{2}\right)=\frac{x_{1} x_{2}}{x_{1}^{2}+x_{2}^{2}}
$$

as $\mathbf{x}=\left(x_{1}, x_{2}\right)^{\prime} \rightarrow \mathbf{0}$, where $\mathbf{0}=(0,0)^{\prime}$. This function is defined everywhere in $\mathbb{R}^{2}$ except at $\mathbf{0}$. It is convenient here to represent the point $\mathbf{x}$ using polar coordinates, $r$ and $\theta$, such that $x_{1}=r \cos \theta$ and $x_{2}=r \sin \theta, r>0,0 \leq \theta \leq 2 \pi$. We then have

$$
f\left(x_{1}, x_{2}\right)=\cos \theta \sin \theta,
$$

which depends on $\theta$, but not on $r$. Since $\theta$ can have infinitely many values, $f\left(x_{1}, x_{2}\right)$ cannot be made close to any one constant $L$ no matter how small $r$ is. Thus the limit of this function does not exist as $\mathbf{x} \rightarrow \mathbf{0}$.

### 8.3 Continuity of a Multivariate Function

The notation of continuity for a function of several variables is much the same as that for a function of a single variable.

## $\mathfrak{D e f i n i t i o n : ~}$

Let $\mathbf{f}: D \rightarrow \mathbb{R}^{m}$, where $D \subset \mathbb{R}^{n}$, and let $\mathbf{a} \in D$. Then $\mathbf{f}(\mathbf{x})$ is continuous at a if

$$
\lim _{\mathbf{x} \rightarrow \mathbf{a}} \mathbf{f}(\mathbf{x})=\mathbf{f}(\mathbf{a})
$$

where $\mathbf{x}$ remains in $D$ as it approaches a. This is equivalent to stating that for a given $\epsilon>0$ there exist a $\delta>0$ such that

$$
\|\mathbf{f}(\mathbf{x})-\mathbf{f}(\mathbf{a})\|<\epsilon
$$

for all $\mathbf{x} \in D \cap N_{\delta}(\mathbf{a})$. If $\mathbf{f}(\mathbf{x})$ is continuous at every point $\mathbf{x}$ in $D$, then it is said to be continuous in $D$. In particular, if $\mathbf{f}(\mathbf{x})$ is continuous in $D$ and if $\delta$ depends only on $\epsilon$, then $\mathbf{f}(\mathbf{x})$ is said to be uniformly continuous in $D$.

## $\mathfrak{L e m m a : ~}$

Suppose that $f, g: D \mapsto \mathbb{R}$ are real-valued continuous functions, where $D \subset \mathbb{R}^{n}$. Then we have the following:
(a). $f+g, f-g$ and $f g$ are continuous in $D$.
(b). $|f|$ is continuous in $D$.
(c). $1 / f$ is continuous in $D$ provided that $f(\mathbf{x}) \neq 0 \forall \mathbf{x} \in D$.

## $\mathfrak{L e m m a}$ :

Suppose that $\mathbf{f}: D \rightarrow \mathbb{R}^{m}$ is continuous, where $D \subset \mathbb{R}^{n}$, and that $\mathbf{g}: G \rightarrow \mathbb{R}^{v}$ is also continuous, where $G \subset \mathbb{R}^{m}$ is the image of $D$ under $\mathbf{f}$. Then the composite function $\mathbf{g} \circ \mathbf{f}(\mathbf{x}): D \rightarrow \mathbb{R}^{v}$, defined as $\mathbf{g} \circ \mathbf{f}(\mathbf{x})=\mathbf{g}[\mathbf{f}(\mathbf{x})]$, is also continuous in $D$.

## $\mathfrak{T h e o r e m}$ :

Suppose that $f: D \mapsto \mathbb{R}$ be a real-valued continuous function defined on a closed and bounded set $D \subset \mathbb{R}^{n}$. Then there exist points $\mathbf{p}$ and $\mathbf{q}$ in $D$ for which

$$
\begin{aligned}
f(\mathbf{p}) & =\sup _{\mathbf{x} \in D} f(\mathbf{x}) \\
f(\mathbf{q}) & =\inf _{\mathbf{x} \in D} f(\mathbf{x})
\end{aligned}
$$

Thus, $f(\mathbf{x})$ attains each of its infimum and supremum at least once in $D$.

## Theorem:

Suppose that $D$ is a closed and bounded set (or called compact set) in $\mathbb{R}^{n}$. If $\mathbf{f}: D \rightarrow \mathbb{R}^{m}$ is continuous, then it is uniformly continuous in $D$.

### 8.4 Derivatives of A Multivariate Function

In this section we generalize the concept of differentiation given in section 4 to a multivariate function $\mathbf{f}: D \mapsto \mathbb{R}^{m}$, where $D \subset \mathbb{R}^{n}$.

### 8.4.1 Partial Derivatives

$\mathfrak{D e f i n i t i o n ~ ( P a r t i a l ~ D e r i v a t i v e ) : ~}$
Let $\mathbf{a}=\left(a_{1}, a_{2}, \ldots, a_{n}\right)^{\prime}$ be an interior point of $D$. Suppose that the limit

$$
\lim _{h_{i} \rightarrow 0} \frac{\mathbf{f}\left(a_{1}, a_{2}, \ldots, a_{i}+h_{i}, \ldots, a_{n}\right)-\mathbf{f}\left(a_{1}, a_{2}, \ldots, a_{i}, \ldots, a_{n}\right)}{h_{i}}
$$

exists; then $\mathbf{f}$ is said to have partial derivative with respect to $x_{i}$ at a. This derivative is denoted by

$$
\frac{\partial \mathbf{f}(\mathbf{a})}{\partial x_{i}}, \text { or } \quad \mathbf{f}_{x_{i}}(\mathbf{a}), i=1,2, \ldots, n
$$

Hence, partial differentiation with respect to $x_{i}$ is done in the usual way while treating all the remaining variables as constant.

Higher-order partial derivative of $\mathbf{f}$ are defined similarly. For example, the second-order partial derivative of $\mathbf{f}$ with respect to $x_{i}$ at $\mathbf{a}$ is defined as

$$
\lim _{h_{i} \rightarrow 0} \frac{\mathbf{f}_{x_{i}}\left(a_{1}, a_{2}, \ldots, a_{i}+h_{i}, \ldots, a_{n}\right)-\mathbf{f}_{x_{i}}\left(a_{1}, a_{2}, \ldots, a_{i}, \ldots, a_{n}\right)}{h_{i}}
$$

and is denoted by

$$
\frac{\partial^{2} \mathbf{f}(\mathbf{a})}{\partial x_{i}^{2}}, \text { or } \quad \mathbf{f}_{x_{i} x_{i}}(\mathbf{a}), i=1,2, \ldots, n
$$

Also, the second-order partial derivative of $\mathbf{f}$ with respect to $x_{i}$ and $x_{i}, i \neq j$ at $\mathbf{a}$ is defined as

$$
\lim _{h_{j} \rightarrow 0} \frac{\mathbf{f}_{x_{i}}\left(a_{1}, a_{2}, \ldots, a_{j}+h_{j}, \ldots, a_{n}\right)-\mathbf{f}_{x_{i}}\left(a_{1}, a_{2}, \ldots, a_{j}, \ldots, a_{n}\right)}{h_{j}}
$$

and is denoted by

$$
\frac{\partial^{2} \mathbf{f}(\mathbf{a})}{\partial x_{i} \partial x_{j}}, \text { or } \quad \mathbf{f}_{x_{i} x_{j}}(\mathbf{a}), \forall i \neq j
$$

## Example:

Suppose that there is a function $f$ that satisfies the equation

$$
f(x, y)=x^{3}+5 x y-y^{2} .
$$

The first partial derivatives of this function are

$$
f_{x}=3 x^{2}+5 y \quad \text { and } \quad f_{y}=5 x-2 y
$$

Therefore, upon further differentiation, we get

$$
f_{x x}=6 x, \quad f_{y x}=5, \quad f_{x y}=5, \quad f_{y y}=-2
$$

$\mathfrak{D e f i n i t i o n ~ ( J a c o b i a n ~ M a t r i x ) : ~}$
In general if $f_{j}$ is the $j$ th element of $\mathbf{f}(j=1,2, . ., m)$, then the term $\frac{\partial f_{j}(\mathbf{x})}{\partial x_{i}}$ for $i=1,2, \ldots, n ; j=1,2, \ldots, m$, constitute an $m \times n$ matrix called the Jacobian matrix of $\mathbf{f}$ at $\mathbf{x}$ and is denoted by $\mathbf{J}_{\mathbf{f}}(\mathbf{x})$. If $m=n$, the determinant of $\mathbf{J}_{\mathbf{f}}(\mathbf{x})$ is called the Jacobian determinant.; it is some times represented as

$$
\operatorname{det}\left[\mathbf{J}_{\mathbf{f}}(\mathbf{x})\right]=\frac{\partial\left(f_{1}, f_{2}, \ldots, f_{n}\right)}{\partial\left(x_{1}, x_{2}, \ldots, x_{n}\right)}
$$

## Theorem:

Let $f, f_{x}, f_{y}$, and $f_{x y}$ exist and be continuous on a neighborhood of a point $\mathbf{a}=\left(x_{0}, y_{0}\right)$. Then $f_{y x}(\mathbf{a})$ exists and $f_{x y}(\mathbf{a})=f_{y x}(\mathbf{a})$.

## $\mathfrak{P r o o f}:$

Let $\phi(x)=f\left(x, y_{0}+k\right)-f\left(x, y_{0}\right)$, where $k$ and $y$ are held fixed. Then for $x$ sufficiently near $x_{0}$ and $k$ small, $\phi$ is a function of the single variable $x$ near $x_{0}$. To this function we apply the mean value theorem for function of one variable between $x_{0}$ and $x_{0}+h$ :

$$
\phi\left(x_{0}+h\right)-\phi\left(x_{0}\right)=h \phi^{\prime}\left(x_{0}+\theta_{1} h\right),
$$

where the prime (') denote differentiation with respect to $x$ and where $0<\theta_{1}<1$ and $\phi^{\prime}\left(x_{0}+\theta_{1} h\right)=f_{x}\left(x_{0}+\theta_{1} h, y_{0}+k\right)-f_{x}\left(x_{0}+\theta_{1} h, y_{0}\right)$. Thus (for fixed $y$ and $k$ )

$$
\phi\left(x_{0}+h\right)-\phi\left(x_{0}\right)=h\left[f_{x}\left(x_{0}+\theta_{1} h, y_{0}+k\right)-f_{x}\left(x_{0}+\theta_{1} h, y_{0}\right)\right] .
$$

Now for each (fixed) $h$ we apply the mean value theorem for functions of one variable $f_{x}\left(x_{0}+\theta_{1} h, y_{0}+k\right)-f_{x}\left(x_{0}+\theta_{1} h, y_{0}\right)$ between $y_{0}$ and $y_{0}+k$ to obtain

$$
f_{x}\left(x_{0}+\theta_{1} h, y_{0}+k\right)-f_{x}\left(x_{0}+\theta_{1} h, y_{0}\right)=k \cdot f_{x y}\left(x_{0}+\theta_{1} h, y_{0}+\theta_{2} k\right)
$$

where $\theta_{2}<1$. Hence

$$
\begin{aligned}
\phi\left(x_{0}+h\right)-\phi\left(x_{0}\right) & =h\left[f_{x}\left(x_{0}+\theta_{1} h, y_{0}+k\right)-f_{x}\left(x_{0}+\theta_{1} h, y_{0}\right)\right] \\
& =h k\left[f_{x y}\left(x_{0}+\theta_{1} h, y_{0}+\theta_{2} k\right)\right] .
\end{aligned}
$$

Recalling the meaning of $\phi$, we can rewrite this as

$$
\begin{align*}
{\left[f\left(x_{0}+h, y_{0}+k\right)-f\left(x_{0}+h, y_{0}\right)\right]-} & {\left[f\left(x_{0}, y_{0}+k\right)-f\left(x_{0}, y_{0}\right)\right] } \\
& =h k\left[f_{x y}\left(x_{0}+\theta_{1} h, y_{0}+\theta_{2} k\right)\right] . \tag{14}
\end{align*}
$$

Dividing (14) by $k$ and letting $k \rightarrow 0$ we get
$\frac{f\left(x_{0}+h, y_{0}+k\right)-f\left(x_{0}+h, y_{0}\right)}{k}-\frac{f\left(x_{0}, y_{0}+k\right)-f\left(x_{0}, y_{0}\right)}{k}=h\left[f_{x y}\left(x_{0}+\theta_{1} h, y_{0}\right)\right]$,
that is

$$
\begin{equation*}
f_{y}\left(x_{0}+h, y_{0}\right)-f_{y}\left(x_{0}, y_{0}\right)=h\left[f_{x y}\left(x_{0}+\theta_{1} h, y_{0}\right)\right] . \tag{15}
\end{equation*}
$$

Dividing (15) by $h$ and letting $h \rightarrow 0$ we get

$$
\frac{f_{y}\left(x_{0}+h, y_{0}\right)-f_{y}\left(x_{0}, y_{0}\right)}{h}=f_{x y}\left(x_{0}, y_{0}\right)
$$

i.e.

$$
f_{y x}\left(x_{0}, y_{0}\right)=f_{x y}\left(x_{0}, y_{0}\right),
$$

the results desired.

### 8.4.2 The Total Derivatives

Let $f(\mathbf{x})$ be a real valued function defined on a set $D \subset \mathbb{R}^{n}$. Suppose that $x_{1}, x_{2}, \ldots, x_{n}$ are functions of a single variable $t$. Then $f$ is a function of $t$. The ordinary derivative of $4 f$ with respect to $t$, namely, $d f / d t$, is called the total derivative of $t$.

## $\mathfrak{L e m m a : ~}$

Assume that for the value of $t$ under consideration $d x_{i} / d t$ exist for $i=1,2, \ldots, n$ and that $\partial f(\mathbf{x}) / \partial x_{i}$ exists and is continuous in the interior of $D$ for $i=1,2, \ldots, n$. Under theses considerations, the total derivatives of $f$ is given by

$$
\frac{d f}{d t}=\sum_{i=1}^{n} \frac{\partial f(\mathbf{x})}{\partial x_{i}} \frac{d x_{i}}{d t}
$$

## $\mathfrak{P r o o f}:$

Let $n=2$ and $\triangle x_{1}, \triangle x_{2}$ be increment of $x_{1}, x_{2}$ that correspond to an increment of $\Delta t$ of $t$. In turn, $f$ will have the increment $\triangle f$. We then have

$$
\begin{aligned}
\Delta f & =f\left(x_{1}+\triangle x_{1}, x_{2}+\triangle x_{2}\right)-f\left(x_{1}, x_{2}\right) \\
& =\left[f\left(x_{1}+\triangle x_{1}, x_{2}+\triangle x_{2}\right)-f\left(x_{1}, x_{2}+\triangle x_{2}\right)\right]+\left[f\left(x_{1}, x_{2}+\Delta x_{2}\right)-f\left(x_{1}, x_{2}\right)\right]
\end{aligned}
$$

By mean-value theorem,

$$
\begin{aligned}
& \left(x_{1}+\Delta x_{1}, x_{2}+\Delta x_{2}\right)-f\left(x_{1}, x_{2}+\Delta x_{2}\right) \\
& =\left(x_{1}+\triangle x_{1}-x_{1}\right) f^{\prime}\left(x_{1}+\theta \cdot \Delta x_{1}, x_{2}+\triangle x_{2}\right) \\
& =\triangle x_{1} \frac{\partial f\left(x_{1}+\theta \cdot \Delta x_{1}, x_{2}+\triangle x_{2}\right)}{\partial x_{1}}, \quad \theta<1 .
\end{aligned}
$$

we then have

$$
\begin{aligned}
\Delta f= & \Delta x_{1} \frac{\partial f\left(x_{1}+\theta_{1} \Delta x_{1}, x_{2}+\Delta x_{2}\right)}{\partial x_{1}} \\
& +\triangle x_{2} \frac{\partial f\left(x_{1}, x_{2}+\theta_{2} \triangle x_{2}\right)}{\partial x_{2}}
\end{aligned}
$$

Hence

$$
\begin{aligned}
\frac{\Delta f}{\Delta t}= & \frac{\Delta x_{1}}{\Delta t} \frac{\partial f\left(x_{1}+\theta_{1} \Delta x_{1}, x_{2}+\Delta x_{2}\right)}{\partial x_{1}} \\
& +\frac{\triangle x_{2}}{\Delta t} \frac{\partial f\left(x_{1}, x_{2}+\theta_{2} \triangle x_{2}\right)}{\partial x_{2}} .
\end{aligned}
$$

As $\Delta t \rightarrow 0, \theta_{i} \rightarrow 0, \triangle x_{i} / \triangle t \rightarrow d x_{i} / d t$, and by the continuity of $\partial f(\mathbf{x}) / \partial x_{i}$ we also have

$$
\begin{aligned}
\frac{\partial f\left(x_{1}+\theta_{1} \triangle x_{1}, x_{2}+\Delta x_{2}\right)}{\partial x_{1}} & =\frac{\partial f\left(x_{1}, x_{2}\right)}{\partial x_{1}} \\
\frac{\partial f\left(x_{1}, x_{2}+\theta_{2} \triangle x_{2}\right)}{\partial x_{2}} & =\frac{\partial f\left(x_{1}, x_{2}\right)}{\partial x_{2}}
\end{aligned}
$$

The desired results thus obtained.

In general, the expression

$$
d f=\sum_{i=1}^{n} \frac{\partial f(\mathbf{x})}{\partial x_{i}} d x_{i}
$$

is called the total differentials of $f$ at $\mathbf{x}$.

### 8.4.3 Directional Derivatives

Let $\mathbf{f}: D \mapsto \mathbb{R}^{m}$, where $D \subset \mathbb{R}^{n}$, and let $\mathbf{v}$ be a unit vector in $\mathbb{R}^{n}$, which represents a certain direction in the $n$-dimensional Euclidean space. By definition, the directional derivative of $\mathbf{f}$ at a point $\mathbf{x}$ is the interior of $D$ in the direction of $\mathbf{v}$ is given by the limit

$$
\lim _{h \rightarrow 0} \frac{\mathbf{f}(\mathbf{x}+h \mathbf{v})-\mathbf{f}(\mathbf{x})}{h},
$$

if it exists. In particular, if $\mathbf{v}=\mathbf{e}_{i}$, the unit vector in the direction of the $i$ th coordinate axis, then the directional derivative of $\mathbf{f}$ in the direction of $\mathbf{v}$ is just the partial derivatives of $\mathbf{f}$ with respect to $x_{i}(i=1,2, . ., n)$.

## $\mathfrak{L e m m a}$ :

Let $\mathbf{f}: D \mapsto \mathbb{R}^{m}$, where $D \subset \mathbb{R}^{n}$. If the partial derivative $\partial f_{j} / \partial x_{i}$ exist at a point $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime}$ in the interior of $D$ for $i=1,2, \ldots, n ; j=1,2, . ., m$, then the directional derivative of $\mathbf{f}$ at $\mathbf{x}$ in the direction of a unit vector $\mathbf{v}$ exists and is equal to $\mathbf{J}_{\mathbf{f}}(\mathbf{x}) \mathbf{v}$, where $\mathbf{J}_{\mathbf{f}}(\mathbf{x})$ is the $m \times n$ Jacobian of $\mathbf{f}$ at $\mathbf{x}$.

## $\mathfrak{P r o o f}:$

Let $m=1, n=2$. The increment of $f$ is

$$
\begin{aligned}
\Delta f & =f\left(x_{1}+h v_{1}, x_{2}+h v_{2}\right)-f\left(x_{1}, x_{2}\right) \\
& =\left[f\left(x_{1}+h v_{1}, x_{2}+h v_{2}\right)-f\left(x_{1}, x_{2}+h v_{2}\right)\right]+\left[f\left(x_{1}, x_{2}+h v_{2}\right)-f\left(x_{1}, x_{2}\right)\right]
\end{aligned}
$$

By mean-value theorem,

$$
\begin{aligned}
& f\left(x_{1}+h v_{1}, x_{2}+h v_{2}\right)-f\left(x_{1}, x_{2}+h v_{2}\right) \\
& =\left(x_{1}+h v_{1}-x_{1}\right) f^{\prime}\left(x_{1}+\theta \cdot h v_{1}, x_{2}+h v_{2}\right) \\
& =h v_{1} \frac{\partial f\left(x_{1}+\theta \cdot h v_{1}, x_{2}+h v_{2}\right)}{\partial x_{1}}, \quad \theta<1 .
\end{aligned}
$$

we then have

$$
\begin{aligned}
\Delta f= & h v_{1} \frac{\partial f\left(x_{1}+\theta_{1} h v_{1}, x_{2}+h v_{2}\right)}{\partial x_{1}} \\
& +h v_{2} \frac{\partial f\left(x_{1}, x_{2}+\theta_{2} h v_{2}\right)}{\partial x_{2}}
\end{aligned}
$$

Hence

$$
\begin{aligned}
\frac{\Delta f}{h} & =\frac{h v_{1}}{h} \frac{\partial f\left(x_{1}+\theta_{1} h v_{1}, x_{2}+h v_{2}\right)}{\partial x_{1}}+\frac{h v_{2}}{h} \frac{\partial f\left(x_{1}, x_{2}+\theta_{2} h v_{2}\right)}{\partial x_{2}} \\
& =v_{1} \frac{\partial f\left(x_{1}+\theta_{1} h v_{1}, x_{2}+h v_{2}\right)}{\partial x_{1}}+v_{2} \frac{\partial f\left(x_{1}, x_{2}+\theta_{2} h v_{2}\right)}{\partial x_{2}} .
\end{aligned}
$$

As $h \rightarrow 0$, and by the continuity of $\partial f(\mathbf{x}) / \partial x_{i}$ we also have

$$
\begin{aligned}
\frac{\partial f\left(x_{1}+\theta_{1} h v_{1}, x_{2}+h v_{2}\right)}{\partial x_{1}} & =\frac{\partial f\left(x_{1}, x_{2}\right)}{\partial x_{1}} \\
\frac{\partial f\left(x_{1}, x_{2}+\theta_{2} h v_{2}\right)}{\partial x_{2}} & =\frac{\partial f\left(x_{1}, x_{2}\right)}{\partial x_{2}}
\end{aligned}
$$

The desired results thus obtained.

## $\mathfrak{E x a m p l e}$ :

Let $\mathbf{f}: \mathbb{R}^{3} \mapsto \mathbb{R}^{2}$ be defined as

$$
\mathbf{f}\left(x_{1}, x_{2}, x_{3}\right)=\left[\begin{array}{c}
x_{1}^{2}+x_{2}^{2}+x_{3}^{2} \\
x_{1}^{2}-x_{1} x_{2}+x_{3}^{2}
\end{array}\right]
$$

The directional derivative of $\mathbf{f}$ at $\mathbf{a}=(1,2,1)^{\prime}$ in the direction of $\mathbf{v}=(1 / \sqrt{2},-1 / \sqrt{2}, 0)^{\prime}$ is

$$
\begin{aligned}
\mathbf{J}_{\mathbf{f}}(\mathbf{x}) \mathbf{v} & =\left[\begin{array}{ccc}
2 x_{1} & 2 x_{2} & 2 x_{3} \\
2 x_{1}-x_{2} & -x_{1} & 2 x_{3}
\end{array}\right]_{(1,2,1)}\left[\begin{array}{c}
\frac{1}{\sqrt{2}} \\
-\frac{1}{\sqrt{2}} \\
0
\end{array}\right] \\
& =\left[\begin{array}{ccc}
2 & 4 & 2 \\
0 & -1 & 2
\end{array}\right]\left[\begin{array}{c}
\frac{1}{\sqrt{2}} \\
-\frac{1}{\sqrt{2}} \\
0
\end{array}\right]=\left[\begin{array}{c}
-\frac{2}{\sqrt{2}} \\
\frac{1}{\sqrt{2}}
\end{array}\right]
\end{aligned}
$$

$\mathfrak{D e f i n i t i o n ~ ( T h e ~ G r a d i e n t ~ o f ~} f$ ):
Let $f: D \mapsto \mathbb{R}$, where $D \subset \mathbb{R}^{n}$. If the partial derivatives $\partial f / \partial x_{i} i=1,2, \ldots, n$ exists at a point $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime}$ in the interior of $D$, then the vector $\left(\partial f / \partial x_{1}, \partial f / \partial x_{2}, \ldots, \partial f / \partial x_{n}\right)^{\prime}$ is called the gradient of $f$ at $\mathbf{x}$ and is denoted by $\nabla f(\mathbf{x})$.
$\mathfrak{D e f i n i t i o n ~ ( T h e ~ H e s s i a n ~ M a t r i x ) : ~}$
Let $f: D \mapsto \mathbb{R}$, where $D \subset \mathbb{R}^{n}$. Then $\boldsymbol{\nabla} f: D \mapsto \mathbb{R}^{n}$. The Jacobian matrix of $\boldsymbol{\nabla} f(\mathbf{x})$ is called the Hessian matrix of $f$ and is denoted by $\mathbf{H}_{f}(\mathbf{x})$. Thus $\mathbf{H}_{f}(\mathbf{x})=\mathbf{J}_{\boldsymbol{\nabla} f}(\mathbf{x})$, that is,

$$
\mathbf{H}_{f}(\mathbf{x})=\left[\begin{array}{cccccc}
\frac{\partial^{2} f(\mathbf{x})}{\partial x_{1} \partial x_{1}} & \frac{\partial^{2} f(\mathbf{x})}{\partial x_{2} \partial x_{1}} & . & . & \cdot & \frac{\partial^{2} f(\mathbf{x})}{\partial x_{n} \partial x_{1}} \\
\cdot & \cdot & . & . & . & \cdot \\
. & \cdot & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot \\
\frac{\partial^{2} f(\mathbf{x})}{\partial x_{1} \partial x_{n}} & \frac{\partial^{2} f(\mathbf{x})}{\partial x_{2} \partial x_{n}} & . & . & \cdot & \cdot \\
\frac{\partial^{2} f(\mathbf{x})}{\partial x_{n} \partial x_{n}}
\end{array}\right]
$$

The determinant of $\mathbf{H}_{f}(\mathbf{x})$ is called the Hessian determinant.

### 8.5 Taylor's Theorem for a Multivariate Function

$\mathfrak{N o t a t i o n}$ (Del Operator):
Let us first introduce the following notation: Let $\mathbf{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\prime}$. Then $\mathbf{x}^{\prime} \boldsymbol{\nabla}$ denotes a first-order differential operator of the form

$$
\mathbf{x}^{\prime} \boldsymbol{\nabla}=\sum_{i=1}^{n} x_{i} \frac{\partial}{\partial x_{i}}
$$

The symbol, $\boldsymbol{\nabla}=\left(\partial / \partial x_{1}, \ldots, \partial / \partial x_{n}\right)^{\prime}$, is called the del operator. If $m$ is a positive integer, then $\left(\mathbf{x}^{\prime} \nabla\right)^{m}$ denote an $m$ th order differential operator. For example $n=2$,

$$
\begin{aligned}
\left(\mathbf{x}^{\prime} \nabla\right)^{2} & =\left(x_{1} \frac{\partial}{\partial x_{1}}+x_{2} \frac{\partial}{\partial x_{2}}\right)^{2} \\
& =x_{1}^{2} \frac{\partial^{2}}{\partial x_{1}^{2}}+2 x_{1} x_{2} \frac{\partial^{2}}{\partial x_{1} \partial x_{2}}+x_{2}^{2} \frac{\partial^{2}}{\partial x_{2}^{2}} .
\end{aligned}
$$

Thus $\left(\mathbf{x}^{\prime} \nabla\right)^{2}$ is obtained by squaring $x_{1} \frac{\partial}{\partial x_{1}}+x_{2} \frac{\partial}{\partial x_{2}}$ in the usual fashion, except that the squares of $\frac{\partial}{\partial x_{1}}$ and $\frac{\partial}{\partial x_{2}}$ are replaced by $\frac{\partial^{2}}{\partial x_{1}^{2}}$ and $\frac{\partial^{2}}{\partial x_{2}^{2}}$, and the product of $\frac{\partial}{\partial x_{1}}$ and $\frac{\partial}{\partial x_{2}}$ is replaced by $\frac{\partial^{2}}{\partial x_{1} \partial x_{2}}$.

The notation $\left(\mathbf{x}^{\prime} \nabla\right)^{m} f\left(\mathbf{x}_{0}\right)$ indicate that $\left(\mathbf{x}^{\prime} \nabla\right)^{m} f(\mathbf{x})$ is evaluated at $\mathbf{x}_{0}$.
$\mathfrak{T h e o r e m}$ (Multivariate Taylor's Theorem):
Let $f: D \mapsto \mathbb{R}$, where $D \subset \mathbb{R}^{n}$, and let $N_{\delta}\left(\mathbf{x}_{0}\right)$ be a neighborhood of $\mathbf{x}_{0} \in D$ such that $N_{\delta}\left(\mathbf{x}_{0}\right) \subset D$. If $f$ and all its partial derivatives of order $\leq r$ exists and continuous in $N_{\delta}\left(\mathbf{x}_{0}\right)$. Then for $\mathbf{x}$ in $N_{\delta}\left(\mathbf{x}_{0}\right)$,

$$
f(\mathbf{x})=f\left(\mathbf{x}_{0}\right)+\sum_{i=1}^{r-1} \frac{\left[\left(\mathbf{x}-\mathbf{x}_{0}\right)^{\prime} \nabla\right]^{i} f\left(\mathbf{x}_{0}\right)}{i!}+\frac{\left[\left(\mathbf{x}-\mathbf{x}_{0}\right)^{\prime} \nabla\right]^{i} f\left(\mathbf{z}_{0}\right)}{r!},
$$

where $\mathbf{z}_{0}$ is a point on the line segment from $\mathbf{x}_{0}$ to $\mathbf{x}$.

## $\mathfrak{E x a m p l e}$ :

In two dimensions, if we set $\mathbf{x}_{0}=(a, b)^{\prime}$ and $\mathbf{x}=(x, y)^{\prime}$, we get

$$
\begin{aligned}
f(x, y)= & f(a, b)+\left[(x-a) \frac{\partial}{\partial x}+(y-b) \frac{\partial}{\partial y}\right] f(a, b) \\
& +\frac{1}{2!}\left[(x-a) \frac{\partial}{\partial x}+(y-b) \frac{\partial}{\partial y}\right]^{2} f(a, b) \\
& +\frac{1}{3!}\left[(x-a) \frac{\partial}{\partial x}+(y-b) \frac{\partial}{\partial y}\right]^{3} f(a, b) \\
& +\ldots+\frac{1}{r!}\left[(x-a) \frac{\partial}{\partial x}+(y-b) \frac{\partial}{\partial y}\right]^{r} f\left(a_{0}, b_{0}\right) .
\end{aligned}
$$

Expanding, we get

$$
\begin{aligned}
f(x, y)= & f(a, b)+(x-a) \frac{\partial f}{\partial x}(a, b)+(y-b) \frac{\partial f}{\partial y}(a, b) \\
& +\frac{1}{2!}\left[(x-a)^{2} \frac{\partial^{2} f(a, b)}{\partial x^{2}}+2(x-a)(y-b) \frac{\partial^{2} f(a, b)}{\partial x \partial y}+(y-b)^{2} \frac{\partial^{2} f(a, b)}{\partial y^{2}}\right] \\
& +\frac{1}{3!}\left[(x-a)^{3} \frac{\partial^{3} f(a, b)}{\partial x^{3}}+3(x-a)^{2}(y-b) \frac{\partial^{3} f(a, b)}{\partial^{2} x \partial y}\right. \\
& \left.+3(x-a)(y-b)^{2} \frac{\partial^{3} f(a, b)}{\partial x \partial^{2} y}+(y-b)^{3} \frac{\partial^{3} f(a, b)}{\partial y^{3}}\right] \\
& +\frac{1}{4!}[\cdots]+\cdots .
\end{aligned}
$$

In the case of one dimensional, it reduce to

$$
\begin{aligned}
f(x)= & f(a)+(x-a) \frac{\partial f}{\partial x}(a) \\
& +\frac{1}{2!}\left[(x-a)^{2} \frac{\partial^{2} f(a)}{\partial x^{2}}\right] \\
& +\frac{1}{3!}\left[(x-a)^{3} \frac{\partial^{3} f(a)}{\partial x^{3}}\right] \\
& +\frac{1}{4!}[\cdots]+\cdots .
\end{aligned}
$$

### 8.6 Optimum of A Multivariate Function

In general, any point at which $\partial f / \partial x_{i}=0$ for $i=1,2, . ., n$ is called a stationary point. The following theorem gives the conditions needed to have a local optimum at s stationary point.
$\mathfrak{T h e o r e m}:$
Let $f: D \rightarrow \mathbb{R}$, where $D \subset \mathbb{R}^{n}$. Suppose that $f$ has continuous second-order partial derivatives in $D$. If $\mathbf{x}_{0}$ is a stationary point of $f$, then at $\mathbf{x}_{0}$ has the following:
(a). A local minimum if $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)>0$ for all $\mathbf{h}=\left(h_{1}, h_{2}, \ldots, h_{n}\right)^{\prime}$ in a neighborhood of $\mathbf{0}$, where the element of $\mathbf{h}$ are not all equal to zero.
(b). A local maximum if $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)<0$ for all $\mathbf{h}=\left(h_{1}, h_{2}, \ldots, h_{n}\right)^{\prime}$ in a neighborhood of $\mathbf{0}$, where the element of $\mathbf{h}$ are not all equal to zero.
(c). A saddle point if $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)$ changes sign for value of $\mathbf{h}=\left(h_{1}, h_{2}, \ldots, h_{n}\right)^{\prime}$ in a neighborhood of $\mathbf{0}$.

## $\mathfrak{P r o o f}:$

By applying Taylor's theorem to $f\left(\mathbf{x}_{0}\right)$ in a neighborhood of $\mathbf{x}_{0}$ we obtain

$$
f\left(\mathbf{x}_{0}+\mathbf{h}\right)=f\left(\mathbf{x}_{0}\right)+\left(\mathbf{h}^{\prime} \nabla\right) f\left(\mathbf{x}_{0}\right)+\frac{1}{2!}\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{z}_{0}\right),
$$

where $\mathbf{h}$ is a nonzero vector in a neighborhood of $\mathbf{0}$ and $\mathbf{z}_{0}$ is a point on the line segment from $\mathbf{x}_{0}$ to $\mathbf{x}_{0}+\mathbf{h}$. Since $\mathbf{x}_{0}$ is a stationary point, then $\left(\mathbf{h}^{\prime} \nabla\right) f\left(\mathbf{x}_{0}\right)=0$. Hence

$$
f\left(\mathbf{x}_{0}+\mathbf{h}\right)-f\left(\mathbf{x}_{0}\right)=\frac{1}{2!}\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{z}_{0}\right)
$$

Also, since the second-order partial derivative of $f$ is continuous at $\mathbf{x}_{0}$, then we can write

$$
f\left(\mathbf{x}_{0}+\mathbf{h}\right)-f\left(\mathbf{x}_{0}\right)=\frac{1}{2!}\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)+o(\|\mathbf{h}\|)
$$

where $\|\mathbf{h}\|=\left(\mathbf{h}^{\prime} \mathbf{h}\right)^{1 / 2}$ and $o(\|\mathbf{h}\|) \rightarrow 0$ as $\mathbf{h} \rightarrow \mathbf{0}$. We note that for small values of $\|\mathbf{h}\|$, the sign of $f\left(\mathbf{x}_{0}+\mathbf{h}\right)-f\left(\mathbf{x}_{0}\right)$ depends on the value of $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)$. It follows that if
(a). $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)>0$, then $f\left(\mathbf{x}_{0}+\mathbf{h}\right)>f\left(\mathbf{x}_{0}\right)$ for all nonzero value of $\mathbf{h}$ in some neighborhood of $\mathbf{0}$. Thus $\mathbf{x}_{0}$ is a local minimum of $f$.
(b). $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)<0$, then $f\left(\mathbf{x}_{0}+\mathbf{h}\right)<f\left(\mathbf{x}_{0}\right)$ for all nonzero value of $\mathbf{h}$ in some neighborhood of $\mathbf{0}$. In this case, $\mathbf{x}_{0}$ is a local maximum of $f$.
(c). $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)$ changes sign inside a neighborhood of $\mathbf{0}$, then $\mathbf{x}_{0}$ is neither a point of local maximum nor a point of local minimum. In this case, $\mathbf{x}_{0}$ is a saddle point.

## Note:

We note that $\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right)$ can be written as a quadratic form of the form $\mathbf{h}^{\prime} \mathbf{H}_{f}\left(\mathbf{x}_{0}\right) \mathbf{h}$, where $\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)$ is the $n \times n$ Hessian matrix evaluated at $\mathbf{x}_{0}$, that is,

$$
\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)=\left[\begin{array}{cccccc}
f_{11} & f_{12} & \cdot & \cdot & \cdot & f_{1 n} \\
f_{21} & f_{22} & \cdot & \cdot & \cdot & f_{2 n} \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot \\
f_{n 1} & f_{n 2} & \cdot & \cdot & \cdot & f_{n n}
\end{array}\right]_{\mathbf{x}_{0}}
$$

## Example:

For the case of $n=2$,

$$
\begin{aligned}
\left(\mathbf{h}^{\prime} \nabla\right)^{2} f\left(\mathbf{x}_{0}\right) & =h_{1}^{2} \frac{\partial^{2} f\left(\mathbf{x}_{0}\right)}{\partial x_{1}^{2}}+2 h_{1} h_{2} \frac{\partial^{2} f\left(\mathbf{x}_{0}\right)}{\partial x_{1} \partial x_{2}}+h_{2}^{2} \frac{\partial^{2} f\left(\mathbf{x}_{0}\right)}{\partial x_{2}^{2}} \\
& =h_{1}^{2} f_{11}+2 h_{1} h_{2} f_{12}+h_{2}^{2} f_{22} \\
& =\left[h_{1} h_{2}\right]\left[\begin{array}{ll}
f_{11} & f_{12} \\
f_{21} & f_{22}
\end{array}\right]_{\mathbf{x}_{0}}\left[\begin{array}{c}
h_{1} \\
h_{2}
\end{array}\right] \\
& =\mathbf{h}^{\prime} \mathbf{H}_{f}\left(\mathbf{x}_{0}\right) \mathbf{h}, \quad \forall \mathbf{h} \neq \mathbf{0} .
\end{aligned}
$$

## Corollary:

Let $f$ be the same function as in the previous theorem, and let $\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)$ be the Hessian matrix. If $\mathbf{x}_{0}$ is a stationary point of $f$, then at $\mathbf{x}_{0} f$ has the following:
(a). A local minimum if $\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)$ is positive definite.
(b). A local maximum if $\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)$ is negative definite.
(c). A saddle point if $\mathbf{H}_{f}\left(\mathbf{x}_{0}\right)$ is neither positive definite nor negative definite.

### 8.7 The Method of Lagrange Multipliers

The method, which is due to Joseph Louis de Lagrange (1936-1813), is used to optimize a real-valued function $f\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, where $x_{1}, x_{2}, \ldots, x_{n}$ are subject to $m(<n)$ equality constrains of the form

$$
\begin{align*}
g_{1}\left(x_{1}, x_{2}, \ldots, x_{n}\right)= & 0 \\
g_{2}\left(x_{1}, x_{2}, \ldots, x_{n}\right)= & 0 \\
& \cdot  \tag{16}\\
& \cdot \\
g_{m}\left(x_{1}, x_{2}, \ldots, x_{n}\right) & =0
\end{align*}
$$

where $g_{1}, g_{2}, \ldots, g_{m}$ are differentiable functions.
The determination of stationary points in this constrained optimization problem is done by first considering the function

$$
\begin{equation*}
F(\mathbf{x})=f(\mathbf{x})+\sum_{j=1}^{m} \lambda_{j} g_{j}(\mathbf{x}) \tag{17}
\end{equation*}
$$

where $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}$ are scalars called Lagrange multipliers. By differentiating () with respect to $x_{1}, x_{2}, \ldots, x_{n}$ and equating the partial derivatives to zero we obtain

$$
\begin{equation*}
\frac{\partial F}{\partial x_{i}}=\frac{\partial f}{\partial x_{i}}+\sum_{j=1}^{m} \frac{\partial g_{j}}{\partial x_{i}}=0, \quad i=1,2, \ldots, n \tag{18}
\end{equation*}
$$

Equations (16) and (18) consists of $m+n$ equations in $m+n$ unknowns, namely, $x_{1}, x_{2}, \ldots, x_{n} ; \lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}$. The solutions for $x_{1}, x_{2}, \ldots, x_{n}$ determine the locations of stationary points.

### 8.8 The Rinmann Integral of a Multivariate Function

In this section we extend the concept of Riemann integration to real-valued function of $n$ variables, $x_{1}, x_{2}, \ldots, x_{n}$.
$\mathfrak{D e f i n i t i o n ~ ( C e l l ) : ~}$
The set of points in $\mathbb{R}^{n}$ whose coordinates satisfy the inequalities

$$
a_{i} \leq x_{i} \leq b_{i}, \quad i=1,2, \ldots, n
$$

where $a_{i}<b_{i}, i=1,2, . ., n$, form an $n$-dimensional cell denoted by $c_{n}(a, b)$. The content of this cell is $\prod_{i=1}^{n}\left(b_{i}-a_{i}\right)$ and is denoted by $\mu\left[c_{n}(a, b)\right]$.

## $\mathfrak{D e f i n i t i o n ~ ( S u b - C e l l ) : ~}$

Suppose that $P_{i}$ is a partition of the interval $\left[a_{i}, b_{i}\right], i=1,2, \ldots, n$. The Cartesian product $P=\times_{i=1}^{n} P_{i}$ is a partition of $c_{n}(a, b)$ and consists of $n$-dimensional subcells of $c_{n}(a, b)$. We denote these subcells by $S_{1}, S_{2}, \ldots, S_{v}$. The content of $S_{i}$ is denoted by $\mu\left(S_{i}\right), i=1,2, \ldots, v$, where $v$ is the number of subcells.

### 8.8.1 The Riemann Integral on Cells

We first define the Riemann integral of a real-valued function $f(\mathbf{x})$ on an $n$-dimensional cell.
$\mathfrak{D e f i n i t i o n ~ ( L o w e r ~ a n d ~ U p p e r ~ S u m ) : ~}$
Let $f: D \mapsto \mathbb{R}$, where $D \subset \mathbb{R}^{n}$. Suppose that $c_{n}(a, b)$ is an $n$-dimensional cell contained in $D$ and that $f$ is bounded on $c_{n}(a, b)$. Let $P$ be a partition of $c_{n}(a, b)$ consisting of the subcells $S_{1}, S_{2}, \ldots, S_{v}$. Let $m_{i}$ and $M_{i}$ be, respectively the infimum and supremum of $f$ on $S_{i}, i=1,2, \ldots, v$. Consider the sums

$$
\begin{aligned}
& L S_{P}(f)=\sum_{i=1}^{v} m_{i} \mu\left(S_{i}\right), \\
& U S_{P}(f)=\sum_{i=1}^{v} M_{i} \mu\left(S_{i}\right) .
\end{aligned}
$$

We refer to $L S_{P}(f)$ and $U S_{P}(f)$ as the lower and upper sums, respectively, of $f$ respect to the partition $P$.

## $\mathfrak{T h e o r e m}$ :

Let $f: D \mapsto \mathbb{R}$, where $D \subset \mathbb{R}^{n}$. Suppose that $f$ is bounded on $c_{n}(a, b) \subset D$. Then $f$ is Riemann integrable on $c_{n}(a, b)$ if and only if for every $\epsilon>0$ there exists a partition $P$ of $c_{n}(a, b)$ such that

$$
U S_{P}(f)-L S_{P}(f)<\epsilon
$$

## Definition:

Let $f: c_{n}(a, b) \mapsto \mathbb{R}$ be a bounded function. Then $f$ is Riemann integrable on $c_{n}(a, b)$ if and only if

$$
\sup _{P} L S_{P}(f)=\inf _{P} U S_{P}(f) .
$$

Their common value is called the Riemann integral of $f$ on $c_{n}(a, b)$ and is denoted by

$$
\int_{c_{n}(a, b)} f(\mathbf{x}) d \mathbf{x} \equiv \int_{a_{1}}^{b_{1}} \int_{a_{2}}^{b_{2}} \cdots \int_{a_{n}}^{b_{n}} f\left(x_{1}, x_{2}, \ldots, x_{n}\right) d x_{1} d x_{2} \cdots d x_{n}
$$

$\mathfrak{T h e o r e m}$ :
If $f$ is continuous on an $n$-dimensional cell $c_{n}(a, b)$, then it is Riemann integrable there.

### 8.8.2 Iterated Riemann Integrals on Cells

The definition of the $n$-tuple integral does not provide a practicable way to evaluate it. We now show that the evaluation of the integral can be obtained by performing $n$ Riemann integrals each of which is carried out with respect to one variable. Let us first consider a double integral.

## $\mathfrak{L e m m a}$ :

Suppose that $f$ is real-valued and continuous on $c_{2}(a, b)$. Define the function $g\left(x_{2}\right)$ as

$$
g\left(x_{2}\right)=\int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1} .
$$

Then $g\left(x_{2}\right)$ is continuous on $\left[a_{2}, b_{2}\right]$. Therefore $g\left(x_{2}\right)$ is Riemann integrable on $\left[a_{2}, b_{2}\right]$, that is, $\int_{a_{2}}^{b_{2}} g\left(x_{2}\right) d x_{2}$ exists. We call the integral

$$
\int_{a_{2}}^{b_{2}} g\left(x_{2}\right) d x_{2}=\int_{a_{2}}^{b_{2}}\left[\int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1}\right] d x_{2}
$$

an iterated integral of order 2 .

## Theorem:

If $f$ is continuous on $c_{2}(a, b)$, then

$$
\int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1} d x_{2}=\int_{a_{2}}^{b_{2}}\left[\int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1}\right] d x_{2}
$$

That is, the double integral is equal to the iterated integral of order 2.
$\mathfrak{T h e o r e m}$ (Fubini's Theorem):
If $f$ is continuous on $c_{2}(a, b)$, then

$$
\begin{aligned}
\int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1} d x_{2} & =\int_{a_{2}}^{b_{2}}\left[\int_{a_{1}}^{b_{1}} f\left(x_{1}, x_{2}\right) d x_{1}\right] d x_{2} \\
& =\int_{a_{1}}^{b_{1}}\left[\int_{a_{2}}^{b_{2}} f\left(x_{1}, x_{2}\right) d x_{2}\right] d x_{1}
\end{aligned}
$$

### 8.8.3 Integration Over General Set

We now consider $n$-tuple Riemann integration over regions in $\mathbb{R}^{n}$ that are not necessarily cell shaped.

Let $f: D \mapsto \mathbb{R}$ be a bounded and continuous function, where $D$ is a bounded region in $\mathbb{R}^{n}$. There exists an $n$-dimensional cell $c_{n}(a, b)$ such that $D \subset c_{n}(a, b)$. Let
$g: c_{n}(a, b) \mapsto \mathbb{R}$ be defined as

$$
g(\mathbf{x})=\left\{\begin{array}{cl}
f(\mathbf{x}), & \mathbf{x} \in D \\
0, & \mathbf{x} \notin D
\end{array}\right.
$$

Then

$$
\int_{c_{n}(a, b)} g(\mathbf{x}) d \mathbf{x}=\int_{D} f(\mathbf{x}) d \mathbf{x}
$$

In practice, it is not always necessary to make reference to $c_{n}(a, b)$ that encloses $D$ in order to evaluate the integral on $D$. Rather, we only need to recognize that the limits of integration in the iterated Riemann integral depend in general on variables that have not yet been integrated out as the following examples.

## $\mathfrak{E x a m p l e :}$

Let $f\left(x_{1}, x_{2}\right)=x_{1} x_{2}$ and $D$ be the region

$$
D=\left\{\left(x_{1}, x_{2}\right) \mid x_{1}^{2}+x_{2}^{2} \leq 1, x_{1} \geq 0, x_{2} \geq 0\right\}
$$

It is easy to see that $D$ is contained inside the two-dimensional cell

$$
c_{2}(0,1)=\left\{\left(x_{1}, x_{2}\right) \mid 0 \leq x_{1} \leq 1,0 \leq x_{2} \geq 1\right\} .
$$

Then

$$
\begin{aligned}
\iint_{D} x_{1} x_{2} d x_{1} d x_{2} & =\int_{0}^{1}\left[\int_{0}^{\left(1-x_{1}^{2}\right)^{1 / 2}} x_{1} x_{2} d x_{2}\right] d x_{1} \\
& =\int_{0}^{1} x_{1}\left[\int_{0}^{\left(1-x_{1}^{2}\right)^{1 / 2}} x_{2} d x_{2}\right] d x_{1} \\
& =\frac{1}{2} \int_{0}^{1} x_{1}\left(1-x_{1}^{2}\right) d x_{1} \\
& =\frac{1}{8}
\end{aligned}
$$

### 8.9 Differentiation Under the Integral Sign


[^0]:    ${ }^{1}$ From Ok, E.F., 2007, Real Analysis with Economic Application, Princeton University Press.

[^1]:    ${ }^{2}$ Read as "not either $A$ or $B$ " $=$ neither $A$ nor $B=$ both not $A$ and $B$.

[^2]:    ${ }^{3}$ Read as "not both $A$ and $B$ " $=$ either not $A$ or not $B$.

[^3]:    ${ }^{4}$ The function $y=f(x)=3$ is not a one-to-one function since for $x_{1}=2, x_{2}=4, f(2)=f(4)=$ 3.
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[^4]:    ${ }^{5}$ Think of $y= \pm \sqrt{x}$.
    ${ }^{6}$ Think of $y=x^{2}$.

[^5]:    ${ }^{7}$ The $\sup _{A} x$ if it exist, is unique, but it may or may not belong to $A$. For example, let $A=\{x \mid x<0\}$, then $\sup _{A} x=0$, which does not belong to $A$.

[^6]:    ${ }^{8}$ If $f(x)$ has a limit $L$ as $x \rightarrow a$, then $L$ must be unique. To show this, suppose that $L_{1}$ and $L_{2}$ are two limits of $f(x)$ as $x \rightarrow a$. Then, for any $\epsilon>0$ there exist $\delta_{1}>0, \delta_{2}>0$ such that

    $$
    \begin{aligned}
    & \left|f(x)-L_{1}\right|<\frac{\epsilon}{2}, \quad \text { if } 0<|x-a|<\delta_{1} \\
    & \left|f(x)-L_{2}\right|<\frac{\epsilon}{2}, \quad \text { if } 0<|x-a|<\delta_{2}
    \end{aligned}
    $$

    Hence, if $\delta=\min \left(\delta_{1}, \delta_{2}\right)$, then

    $$
    \begin{aligned}
    \left|L_{1}-L_{2}\right| & =\left|L_{1}-f(x)+f(x)-L_{2}\right| \\
    & \leq\left|f(x)-L_{1}\right|+\left|f(x)-L_{2}\right| \\
    & <\epsilon
    \end{aligned}
    $$

    for all $x$ for which $0<|x-a|<\delta$. Since $\left|L_{1}-L_{2}\right|$ is smaller than $\epsilon$, which is an arbitrary positive number, we must have $L_{1}=L_{2}$.

    Note:
    To see the last results, suppose that $L_{1}>L_{2}$, then $L_{1}-L_{2}>0$. Let $L_{1}-L_{2}=2 \epsilon>0$. However this is not the case that $\left|L_{1}-L_{2}\right|<\epsilon$. This is a contradiction.

[^7]:    ${ }^{9}$ This is the case for limit theorem in statistics

[^8]:    ${ }^{10} \mathfrak{P r o o f}:$
    We first prove the sufficiency of the condition. We assume that $\lim _{x \rightarrow a^{+}} f(x)=\lim _{x \rightarrow a^{-}} f(x)=L$, and show $\lim _{x \rightarrow a} f(x)=L$.

    For a given $\epsilon$, there exist a $\delta_{1}$ for which

    $$
    |f(x)-L|<\epsilon \quad \text { if } a<x<a+\delta_{1} .
    $$

[^9]:    ${ }^{11} \mathfrak{P r o o f}$ (for part (c)):
    Let $\epsilon>0$ be given. If $M \neq 0$, then there exists a $\lambda_{1}>0$ such that $|g(x)|>|M| / 2$ if $0<|x-a|<$ $\lambda_{1} .{ }^{12}$ Also, there exists a $\lambda_{2}$ such that $g(x)-M \mid<\epsilon M^{2} / 2$ if $0<|x-a|<\lambda_{2}$. Then,

    $$
    \begin{aligned}
    \left|\frac{1}{g(x)}-\frac{1}{M}\right| & =\frac{|g(x)-M|}{|g(x)||M|} \\
    & <\frac{2|g(x)-M|}{|M|^{2}} \\
    & <\epsilon
    \end{aligned}
    $$

[^10]:    ${ }^{15} \mathfrak{P r o o f}:$
    Let $h_{1}(x)=O(f(x))$, then $\left|h_{1}(x)\right| \leq K f(x)$. Similarly $\left|h_{2}(x)\right| \leq K g(x)$. Therefore

    $$
    \left|h_{1}(x) \cdot h_{2}(x)\right| \leq\left|h_{1}(x)\right| \cdot\left|h_{2}(x)\right| \leq K^{2}(f(x) g(x))
    $$

    i.e. $h_{1}(x) h_{2}(x)$ is $O((f(x) g(x))$.
    ${ }^{16} \mathrm{So}$, the point of interest in defined in the domain of $f$.

[^11]:    ${ }^{18}$ Here, it look like that $\delta$ still depends on $x_{0}$.

[^12]:    ${ }^{19}$ The direct implication of the intermediate-value theorem is that a continuous function possesses the properties of assuming at least once every value between any two distinct value taken inside its domain.

[^13]:    ${ }^{20}$ As an example of Lipschitz continuous function, consider $f(x)=\sqrt{x}, x \geq 0$. We claim that $\sqrt{x}$ is $\operatorname{Lip}(1,1 / 2)$ on its domain. To see this, we first write

    $$
    \left|\sqrt{x_{1}}-\sqrt{x_{2}}\right| \leq \sqrt{x_{1}}+\sqrt{x_{2}} .
    $$

    Hence,

    $$
    \left|\sqrt{x_{1}}-\sqrt{x_{2}}\right|^{2} \leq\left|x_{1}-x_{2}\right| .
    $$

[^14]:    ${ }^{21}$ It is important to note that in order for $f^{\prime}\left(x_{0}\right)$ to exist, the left-sided and right-sided limits of $\phi(h)$ must exist and be equal as $h \rightarrow 0$, or as $x$ approaches $x_{0}$ from either side. It is possible to consider only one-sided derivatives at $x=x_{0}$. These occur when $\phi(h)$ has just a one-sided limit as $h \rightarrow 0$. We shall not, however, concern ourselves with such derivatives in this courses.

[^15]:    ${ }^{22}$ That is, the sign of $f^{\prime}(x)$ provides information about the behavior of $f(x)$ in a neighborhood of $x$.

[^16]:    ${ }^{23}$ Set $x=b$ and $n=1$, we have $f(b)=f(a)+(b-a) f^{\prime}(c)$, the mean value theorem.
    ${ }^{24}$ To make $a<\xi<x$, we need $a<\xi=a+\theta_{n} h<a+1 \cdot h=x$.

[^17]:    ${ }^{25}$ Think of $o(1) \cdot h^{n}=o\left(h^{n}\right)$ since by definition $o\left(h^{n}\right)=\frac{h^{n} o(1)}{h^{n}}=o(1)$ as $h \rightarrow 0$.

[^18]:    ${ }^{26}$ As a example, consider the sequence $a_{n}=(-1)^{n}$. The sequence is bounded, but is not convergent.

[^19]:    ${ }^{27}$ Therefore it has a convergent subsequence.

[^20]:    ${ }^{28}$ Consider the sequence whose $n$th term is $a_{n}=1-\frac{1}{3}+\frac{1}{5}-\frac{1}{7}+\cdots+\frac{(-1)^{n-1}}{2 n-1}, n=1,2, \ldots$. It is not easy to calculate the limit of $a_{n}$ in order to find out if the sequence converges.

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[^22]:    ${ }^{29}$ That is. the lower sum $L S_{P}$ and the upper sum $U S_{P}$ is equal.

[^23]:    ${ }^{30}$ It is noted here to check whether $g(x)$ is monotone increasing on $[a, b]$.

[^24]:    ${ }^{31}$ See the picture on page 250 of Fulks.

