Lecture notes for Math 417-517
Multivariable Calculus

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1 multivariable calculus

1.1 vectors

We start with some definitions. A real number $x$ is positive, zero, or negative and is rational or irrational. We denote

$$\mathbb{R} = \text{set of all real numbers } x$$  \hspace{1cm} (1)

The real numbers label the points on a line once we pick an origin and a unit of length. Real numbers are also called *scalars*.

Next define

$$\mathbb{R}^2 = \text{all pairs of real numbers } x = (x_1, x_2)$$  \hspace{1cm} (2)

The elements of $\mathbb{R}^2$ label points in the plane once we pick an origin and a pair of orthogonal axes. Elements of $\mathbb{R}^2$ are also called (2-dimensional) *vectors* and can be represented by arrows from the origin to the point represented.

Next define

$$\mathbb{R}^3 = \text{all triples of real numbers } x = (x_1, x_2, x_3)$$  \hspace{1cm} (3)

The elements of $\mathbb{R}^3$ label points in space once we pick an origin and three orthogonal axes. Elements of $\mathbb{R}^3$ are (3-dimensional) vectors. Especially for $\mathbb{R}^3$ one might emphasize that $x$ is a vector by writing it in bold face $\mathbf{x} = (x_1, x_2, x_3)$ or with an arrow $\vec{x} = (x_1, x_2, x_3)$ but we refrain from doing this for the time being.

Generalizing still further we define

$$\mathbb{R}^n = \text{all } n\text{-tuples of real numbers } x = (x_1, x_2, \ldots, x_n)$$  \hspace{1cm} (4)

The elements of $\mathbb{R}^n$ are the points in $n$-dimensional space and are also called ($n$-dimensional) vectors.

Given a vector $x = (x_1, \ldots, x_n)$ in $\mathbb{R}^n$ and a scalar $\alpha \in \mathbb{R}$ the product is the vector

$$\alpha x = (\alpha x_1, \ldots, \alpha x_n)$$  \hspace{1cm} (5)

Another vector $y = (y_1, \ldots, y_n)$ can to added to $x$ to give a vector

$$x + y = (x_1 + y_1, \ldots, x_n + y_n)$$  \hspace{1cm} (6)

Because elements of $\mathbb{R}^n$ can be multiplied by a scalar and added it is called a *vector space*. We can also subtract vectors defining $x - y = x + (-1)y$ and then

$$x - y = (x_1 - y_1, \ldots, x_n - y_n)$$  \hspace{1cm} (7)

For two or three dimensional vectors these operations have a geometric interpretation. $\alpha x$ is a vector in the same direction as $x$ (opposite direction if $\alpha < 0$) with length
increased by $|\alpha|$. The vector $x + y$ can be found by completing a parallelogram with sides $x, y$ and taking the diagonal, or by putting the tail of $y$ on the head of $x$ and drawing the arrow from the tail of $x$ to the head of $y$. The vector $x - y$ is found by drawing $x + (-1)y$. Alternatively if the tail of $x - y$ put a the head of $y$ then the arrow goes from the head of $y$ to the head of $x$. See figure 1.

A vector $x = (x_1, \ldots, x_n)$ has a length which is

$$|x| = \text{length of } x = \sqrt{x_1^2 + \cdots + x_n^2}$$

(8)

Since $x - y$ goes from the point $y$ to the point $x$, the length of this vector is the distance between the points:

$$|x - y| = \text{distance between } x \text{ and } y = \sqrt{(x_1 - y_1)^2 + \cdots + (x_n - y_n)^2}$$

(9)

One can also form the dot product of vectors $x, y$ in $\mathbb{R}^n$. The result is a scalar given by

$$x \cdot y = x_1y_1 + x_2y_2 + \cdots + x_ny_n$$

(10)

Then we have

$$x \cdot x = |x|^2$$

(11)
## 1.2 functions of several variables

We are interested in functions $f$ from $\mathbb{R}^n$ to $\mathbb{R}^m$ (or more generally from a subset $D \subset \mathbb{R}^n$ to $\mathbb{R}^m$ called the domain of the function). A function $f$ assigns to each $x \in \mathbb{R}^n$ a point $y \in \mathbb{R}^m$ and we write

$$y = f(x)$$

(12)

The set of all such points $y$ is the range of the function.

Each component of $y = (y_1, \ldots, y_m)$ is real-valued function of $x \in \mathbb{R}^n$ written $y_i = f_i(x)$ and the function can also be written as the collection of $n$ functions

$$y_1 = f_1(x), \ldots, y_m = f_m(x)$$

(13)

If we also write out the components of $x = (x_1, \ldots, x_n)$, then a function can be written as $m$ functions of $n$ variables each:

$$y_1 = f_1(x_1, \ldots, x_n)$$

$$y_2 = f_2(x_1, \ldots, x_n)$$

$$\ldots$$

$$y_m = f_m(x_1, \ldots, x_n)$$

(14)

The graph of the function is all pairs $(x, y)$ with $y = f(x)$. It is a subset of $\mathbb{R}^{n+m}$.

### special cases:

1. $n = 1, m = 2$ (or $m = 3$). The function has the form

$$y_1 = f_1(x) \quad y_2 = f_2(x)$$

(15)

In this case the range of the function is a curve in $\mathbb{R}^2$.

2. $n = 2, m = 1$. Then function has the form

$$y = f(x_1, x_2)$$

(16)

The graph of the function is a surface in $\mathbb{R}^3$.

3. $n = 2, m = 3$. The function has the form

$$y_1 = f_1(x_1, x_2)$$

$$y_2 = f_2(x_1, x_2)$$

$$y_3 = f_3(x_1, x_2)$$

(17)

The range of the function is a surface in $\mathbb{R}^3$.

4. $n = 3, m = 3$. The function has the form

$$y_1 = f_1(x_1, x_2, x_3)$$

$$y_2 = f_2(x_1, x_2, x_3)$$

$$y_3 = f_3(x_1, x_2, x_3)$$

(18)

The function assigns a vector to each point in space and is called a vector field.
1.3 limits

Consider a function \( y = f(x) \) from \( \mathbb{R}^n \) to \( \mathbb{R}^m \) (or possibly a subset of \( \mathbb{R}^n \)). Let \( x^0 = (x_0^1, \ldots, x_0^n) \) be a point in \( \mathbb{R}^n \) and let \( y^0 = (y_0^1, \ldots, y_0^m) \) be a point in \( \mathbb{R}^m \). We say that \( y^0 \) is the limit of \( f \) as \( x \) goes to \( x^0 \), written

\[
\lim_{x \to x^0} f(x) = y^0
\]

if for every \( \epsilon > 0 \) there exists a \( \delta > 0 \) so that if \( |x - x^0| < \delta \) then \( |f(x) - y^0| < \epsilon \). The function is continuous at \( x^0 \) if

\[
\lim_{x \to x^0} f(x) = f(x^0)
\]

The function is continuous if it is continuous at every point in its domain.

If \( f, g \) are continuous at \( x^0 \) then so are \( f \pm g \). If \( f, g \) are scalars (i.e. if \( m = 1 \)) then the the product \( fg \) is defined and continuous at \( x^0 \). If \( f, g \) are scalars and \( g(x^0) \neq 0 \) then \( f/g \) is defined near \( x^0 \) and and continuous at \( x^0 \).

1.4 partial derivatives

At first suppose \( f \) is a function from \( \mathbb{R}^2 \) to \( \mathbb{R} \) written

\[
z = f(x, y)
\]

We define the partial derivative of \( f \) with respect to \( x \) at \((x_0, y_0)\) to be

\[
f_x(x_0, y_0) = \lim_{h \to 0} \frac{f(x_0 + h, y_0) - f(x_0, y_0)}{h}
\]

if the limit exists. It is the same as the ordinary derivative with \( y \) fixed at \( y_0 \), i.e

\[
\left[ \frac{d}{dx} f(x, y_0) \right]_{x=x_0}
\]

We also define the partial derivative of \( f \) with respect to \( y \) at \((x_0, y_0)\) to be

\[
f_y(x_0, y_0) = \lim_{h \to 0} \frac{f(x_0, y_0 + h) - f(x_0, y_0)}{h}
\]

if the limit exists. It is the same as the ordinary derivative with \( x \) fixed at \( x_0 \), i.e

\[
\left[ \frac{d}{dy} f(x_0, y) \right]_{y=y_0}
\]
We also use the notation

\[ f_x = \frac{\partial z}{\partial x} \quad \text{(or } \frac{\partial f}{\partial x} \text{ or } z_x) \]

\[ f_y = \frac{\partial z}{\partial y} \quad \text{(or } \frac{\partial f}{\partial y} \text{ or } z_y) \]  

(26)

If we let \((x_0, y_0)\) vary the partial derivatives are also functions and we can take second partial derivatives like

\[ (f_x)_x \equiv f_{xx} \quad \text{also written} \quad \frac{\partial}{\partial x} \left( \frac{\partial z}{\partial x} \right) = \frac{\partial^2 z}{\partial x^2} \]  

(27)

The four second partial derivatives are

\[ f_{xx} = \frac{\partial^2 z}{\partial x^2} \quad f_{xy} = \frac{\partial^2 z}{\partial y \partial x} \]

\[ f_{yx} = \frac{\partial^2 z}{\partial x \partial y} \quad f_{yy} = \frac{\partial^2 z}{\partial y^2} \]  

(28)

Usually \(f_{xy} = f_{yx}\) for we have

**Theorem 1** If \(f_x, f_y, f_{xy}, f_{yx}\) exist and are continuous near \((x_0, y_0)\) (i.e. in a little disc centered on \((x_0, y_0)\)) then

\[ f_{xy}(x_0, y_0) = f_{yx}(x_0, y_0) \]  

(29)

**Example:** Consider \(f(x, y) = 3x^2y + 4xy^3\). Then

\[ f_x = 6xy + 4y^3 \quad f_y = 3x^2 + 12xy^2 \]

\[ f_{xy} = 6x + 12y^2 \quad f_{yx} = 6x + 12y^2 \]  

(30)

We also have partial derivatives for a function \(f\) from \(\mathbb{R}^n\) to \(\mathbb{R}\) written \(y = f(x_1, \ldots x_n)\). The partial derivative with respect to \(x_i\) at \((x_1^0, \ldots x_n^0)\) is

\[ f_{x_i}(x_1^0, \ldots x_n^0) = \lim_{h \to 0} \frac{f(x_1^0, \ldots, x_i^0 + h, \ldots, x_n^0) - f(x_1^0, \ldots, x_n^0)}{h} \]  

(31)

It is also written

\[ f_{x_i} = \frac{\partial y}{\partial x_i} \]  

(32)
1.5 derivatives

A function \( z = f(x, y) \) is said to be differentiable at \((x_0, y_0)\) if it can be well-approximated by a linear function near that point. This means there should be constants \( a, b \) such that

\[
f(x, y) = f(x_0, y_0) + a(x - x_0) + b(y - y_0) + \epsilon(x, y)
\]

where the error term \( \epsilon(x, y) \) is continuous at \((x_0, y_0)\) and \( \epsilon(x, y) \to 0 \) as \((x, y) \to (x_0, y_0)\) faster than the distance between the points:

\[
\lim_{(x, y) \to (x_0, y_0)} \frac{\epsilon(x, y)}{|(x, y) - (x_0, y_0)|} = 0
\]

Note that differentiable implies continuous.

Suppose it is true and take \((x, y) = (x_0 + h, y_0)\). Then

\[
f(x_0 + h, y_0) = f(x_0, y_0) + ah + \epsilon(x_0 + h, y_0)
\]

and so

\[
\frac{f(x_0 + h, y_0) - f(x_0, y_0)}{h} = a + \frac{\epsilon(x_0 + h, y_0)}{h}
\]

Taking the limit \( h \to 0 \) we see that \( f_x(x_0, y_0) \) exists and equals \( a \). Similarly if we take \((x, y) = (x_0, y_0 + h)\) we find that \( f_y(x_0, y_0) \) exists and equals \( b \).

Thus if \( f \) is differentiable then

\[
f(x, y) = f(x_0, y_0) + f_x(x_0, y_0)(x - x_0) + f_y(x_0, y_0)(y - y_0) + \epsilon(x, y)
\]

where \( \epsilon \) satisfies the above condition. The linear approximation is the function

\[
z = f(x_0, y_0) + f_x(x_0, y_0)(x - x_0) + f_y(x_0, y_0)(y - y_0)
\]

The graph of the linear approximation is a plane called the tangent plane. See figure 2.

It is possible that the partial derivatives \( f_x(x_0, y_0) \) and \( f_y(x_0, y_0) \) exist but still the function is not differentiable as the following example shows.

**example.** Define a function by

\[
f(x, y) = \begin{cases} 1 & \text{if } x = 0 \text{ or } y = 0 \\ 0 & \text{otherwise} \end{cases}
\]

Then

\[
f_x(0, 0) = 0 \quad f_y(0, 0) = 0
\]

But the function cannot be differentiable at \((0, 0)\) since it is not continuous there. It is not continuous since for example

\[
\lim_{t \to 0} f(t, t) = 0 \quad f(0, 0) = 1
\]

However the following is true:
Theorem 2 If the partial derivatives $f_x, f_y$ exist and are continuous near $(x_0, y_0)$ (i.e. in a little disc centered on $(x_0, y_0)$) then $f$ is differentiable at $(x_0, y_0)$.

**Problem:** Show the function $f(x, y) = y^3 + 3x^2y^2$ is differentiable at any point and find the linear approximation (tangent plane) at $(1, 1)$.

**Solution:** This has partial derivatives

$$f_x(x, y) = 6xy^2 \quad f_y(x, y) = 3y^2 + 6x^2y$$

at any point and they are continuous. Thus the function is differentiable. The tangent plane at $(1, 1)$ is

$$z = f(1, 1) + f_x(1, 1)(x - 1) + f_y(1, 1)(y - 1)$$

$$= 1 + 6(x - 1) + 9(y - 1)$$

$$= 6x + 9y - 11$$

Next consider a function from $\mathbb{R}^n$ to $\mathbb{R}$ written $y = f(x) = f(x_1, \ldots, x_n)$. We say $f$ is differentiable at $x^0 = (x_1^0, \ldots, x_n^0)$ if there is a vector $a = (a_1, \ldots, a_n)$ such that

$$f(x) = f(x^0) + a \cdot (x - x^0) + \epsilon(x)$$

(44)
where as before

\[
\lim_{x \to x^0} \frac{\epsilon(x)}{|x - x^0|} = 0 \tag{45}
\]

If it is true then we find as before that the vector must be

\[ a = (f_{x_1}(x^0), \ldots, f_{x_n}(x^0)) \equiv (\nabla f)(x^0) \tag{46} \]

also called the gradient of \( f \) at \( x^0 \). Thus we have

\[ f(x) = f(x^0) + (\nabla f)(x^0) \cdot (x - x^0) + \epsilon(x) \tag{47} \]

Finally consider a function \( y = f(x) \) from \( \mathbb{R}^n \) to \( \mathbb{R}^m \). We write the points and the function as column vectors:

\[
y = \begin{pmatrix}
y_1 \\
\vdots \\
y_m
\end{pmatrix} = \begin{pmatrix}
f_1(x) \\
\vdots \\
f_m(x)
\end{pmatrix} \quad x = \begin{pmatrix}
x_1 \\
\vdots \\
x_n
\end{pmatrix} \tag{48}
\]

The function is differentiable at \( x^0 \) is there is an \( m \times n \) matrix (m rows, n columns) \( A \) such that

\[ f(x) = f(x^0) + A(x - x^0) + \epsilon(x) \tag{49} \]

where the error \( \epsilon(x) \in \mathbb{R}^m \) satisfies

\[
\lim_{x \to x^0} \frac{|\epsilon(x)|}{|x - x^0|} = 0 \tag{50}
\]

By considering each component separately we find that the \( i^{th} \) row of \( A \) must be the the gradient of \( f_i \) at \( x^0 \). Thus

\[ A = Df(x^0) \equiv \begin{pmatrix}
f_{1,x_1}(x^0) & \cdots & f_{1,x_n}(x^0) \\
\vdots \\
f_{m,x_1}(x^0) & \cdots & f_{m,x_n}(x^0)
\end{pmatrix} \tag{51}
\]

This matrix \( Df(x^0) \) made up of all the partial derivatives of \( f \) at \( x^0 \) is the derivative of \( f \) at \( x^0 \). Thus we have

\[ f(x) = f(x_0) + Df(x^0)(x - x^0) + \epsilon(x) \tag{52} \]

The derivative is also written

\[ Df \equiv \begin{pmatrix}
\frac{\partial y_1}{\partial x_1} & \cdots & \frac{\partial y_1}{\partial x_n} \\
\vdots \\
\frac{\partial y_m}{\partial x_1} & \cdots & \frac{\partial y_m}{\partial x_n}
\end{pmatrix} \tag{53} \]
**problem:** Consider the function from $\mathbb{R}^2$ to $\mathbb{R}^2$ given by

$$\begin{pmatrix} y_1 \\ y_2 \end{pmatrix} = \begin{pmatrix} f_1(x) \\ f_2(x) \end{pmatrix} = \begin{pmatrix} (x_1 + x_2 + 1)^2 \\ x_1(x_2 + 3) \end{pmatrix}$$

(54)

Find the linear approximation at $(0,0)$

**solution:** The derivative is

$$Df = \begin{pmatrix} \frac{\partial y_1}{\partial x_1} & \frac{\partial y_1}{\partial x_2} \\ \frac{\partial y_2}{\partial x_1} & \frac{\partial y_2}{\partial x_2} \end{pmatrix} = \begin{pmatrix} 2(x_1 + x_2 + 1) & 2(x_1 + x_2 + 1) \\ x_2 + 3 & x_1 \end{pmatrix}$$

(55)

The linear approximation is

$$y = f(0) + Df(0)(x - 0)$$

or

$$\begin{pmatrix} y_1 \\ y_2 \end{pmatrix} = \begin{pmatrix} 1 \\ 0 \end{pmatrix} + \begin{pmatrix} 2 & 2 \\ 3 & 0 \end{pmatrix} \begin{pmatrix} x_1 \\ x_2 \end{pmatrix} = \begin{pmatrix} 2x_1 + 2x_2 + 1 \\ 3x_1 \end{pmatrix}$$

(57)

1.6 **the chain rule**

If $f : \mathbb{R}^n \to \mathbb{R}^m$ (i.e. $f$ is a function from $\mathbb{R}^n$ to $\mathbb{R}^m$) and $p : \mathbb{R}^k \to \mathbb{R}^n$, then there is a composite function $h : \mathbb{R}^k \to \mathbb{R}^m$ defined by

$$h(x) = f(p(x))$$

(58)

and we also write $h = f \circ p$. We can represent the situation by the diagram:

$$\begin{array}{c}
\mathbb{R}^k \\
\uparrow p \\
\mathbb{R}^n \\
\downarrow h \\
\mathbb{R}^m \\
\downarrow f \\
\mathbb{R}^m
\end{array}$$

The chain gives a formula for the derivatives of the composite function $h$ in terms of the derivatives of $f$ and $p$. We start with some special cases.

**k=1,n=2,m=1** In this case the functions have the form

$$u = f(x, y)$$

$$x = p_1(t) \quad y = p_2(t)$$

(59)

and the composite is

$$u = h(t) = f(p_1(t), p_2(t))$$

(60)
Theorem 3  (chain rule) If $f$ and $p$ are differentiable, then so is the composite $h = f \circ p$ and the derivative is

$$h'(t) = f_x(p_1(t), p_2(t))p_1'(t) + f_y(p_1(t), p_2(t))p_2'(t)$$  \tag{61}

The idea of the proof is as follows. Since $f$ is differentiable

$$h(t + \Delta t) - h(t) = f((p_1(t + \Delta t), p_2(t + \Delta t)) - f(p_1(t), p_2(t))$$

$$= f_x(p_1(t), p_2(t))\left(p_1(t + \Delta t) - p_2(t)\right) + f_y(p_1(t), p_2(t))\left(p_2(t + \Delta t) - p_2(t)\right)$$

$$+ \epsilon(p_1(t + \Delta t), p_2(t + \Delta t))$$  \tag{62}

Now divide by $\Delta t$ and let $\Delta t \to 0$. One has to show that the error term goes to zero and the result follows.

This form of the chain rule can also be written in the concise form

$$\frac{du}{dt} = \frac{\partial u}{\partial x} \frac{dx}{dt} + \frac{\partial u}{\partial y} \frac{dy}{dt}$$  \tag{63}

But one must keep in mind that $\partial u/\partial x$ and $\partial u/\partial y$ are to be evaluated at $x = p_1(t), y = p_2(t)$. Also note that on the left $u$ stands for the function $u = h(t)$ while on the right it stands for the function $u = f(x, y)$.

example: Suppose $u = x^2 + y^2$ and $x = \cos t, y = \sin t$. Find $du/dt$. We have

$$\frac{du}{dt} = \frac{\partial u}{\partial x} \frac{dx}{dt} + \frac{\partial u}{\partial y} \frac{dy}{dt}$$

$$= 2x(-\sin t) + 2y \cos t$$

$$= 2 \cos t(-\sin t) + 2 \sin t \cos t$$

$$= 0$$  \tag{64}

(This is to be expected since the composite is $u = 1$).

example: Suppose $u = \sqrt{2 + x^2 + y^2}$ and $x = e^t, y = e^{2t}$. Find $du/dt$ at $t = 0$. At $t = 0$ we have $x = 1, y = 1$ and so

$$\frac{du}{dt} = \frac{\partial u}{\partial x} \frac{dx}{dt} + \frac{\partial u}{\partial y} \frac{dy}{dt}$$

$$= \frac{x}{\sqrt{2 + x^2 + y^2}} e^t + \frac{y}{\sqrt{2 + x^2 + y^2}} 2e^{2t}$$

$$= \frac{1}{2} \cdot 1 + \frac{1}{2} \cdot 2 = \frac{3}{2}$$  \tag{65}
In this case the functions have the form

\[ u = f(x_1, \ldots, x_n) \]
\[ x_1 = p_1(t), \ldots, x_n = p_n(t) \]  \hspace{1cm} (66)

and the composite is

\[ u = h(t) = f(p_1(t), \cdots, p_n(t)) \]  \hspace{1cm} (67)

The chain rule says

\[ h'(t) = f_{x_1}(p_1(t), \cdots, p_n(t))p_1'(t) + \cdots + f_{x_n}(p_1(t), \cdots, p_n(t))p_n'(t) \]  \hspace{1cm} (68)

It can also be written

\[ \frac{du}{dt} = \frac{\partial u}{\partial x_1} \frac{dx_1}{dt} + \cdots + \frac{\partial u}{\partial x_n} \frac{dx_n}{dt} \]  \hspace{1cm} (69)

In this case the functions have the form

\[ u = f_1(x, y) \quad v = f_2(x, y) \]
\[ x = p_1(s, t) \quad y = p_2(s, t) \]  \hspace{1cm} (70)

and the composite is

\[ u = h_1(s, t) = f_1(p_1(s, t), p_2(s, t)) \quad v = h_2(s, t) = f_2(p_1(s, t), p_2(s, t)) \]  \hspace{1cm} (71)

Taking partial derivatives with respect to \( s, t \) we can use the formula from the case \( k=1,n=2,m=1 \) to obtain

\[ \frac{\partial u}{\partial s} = \frac{\partial u}{\partial x} \frac{\partial x}{\partial s} + \frac{\partial u}{\partial y} \frac{\partial y}{\partial s} \]
\[ \frac{\partial u}{\partial t} = \frac{\partial u}{\partial x} \frac{\partial x}{\partial t} + \frac{\partial u}{\partial y} \frac{\partial y}{\partial t} \]
\[ \frac{\partial v}{\partial s} = \frac{\partial v}{\partial x} \frac{\partial x}{\partial s} + \frac{\partial v}{\partial y} \frac{\partial y}{\partial s} \]
\[ \frac{\partial v}{\partial t} = \frac{\partial v}{\partial x} \frac{\partial x}{\partial t} + \frac{\partial v}{\partial y} \frac{\partial y}{\partial t} \]  \hspace{1cm} (72)

Here derivatives with respect to \( x, y \) are to be evaluated at \( x = p_1(s, t), y = p_2(s, t) \). This can also be written as a matrix product:

\[
\begin{pmatrix}
\frac{\partial u}{\partial s} & \frac{\partial u}{\partial t} \\
\frac{\partial v}{\partial s} & \frac{\partial v}{\partial t}
\end{pmatrix} = \begin{pmatrix}
\frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\
\frac{\partial v}{\partial x} & \frac{\partial v}{\partial y}
\end{pmatrix} \begin{pmatrix}
\frac{\partial x}{\partial s} & \frac{\partial x}{\partial t} \\
\frac{\partial y}{\partial s} & \frac{\partial y}{\partial t}
\end{pmatrix} \hspace{1cm} (73)
These matrices are just the derivatives of the various functions and the last equation can be written

\[(Dh)(s,t) = (Df)(p(s,t))(Dp)(s,t)\]  \( (74) \)

The last form holds in the general case. If \( p : \mathbb{R}^k \to \mathbb{R}^n \) and \( f : \mathbb{R}^n \to \mathbb{R}^m \) are differentiable then the composite function \( h = f \circ p : \mathbb{R}^k \to \mathbb{R}^m \) is differentiable and

\[
\begin{align*}
(Dh)(x) &= (Df)(p(x)) (Dp)(x)
\end{align*}
\]  \( (75) \)

In other words the derivative of a composite is the matrix product of the derivatives of the two elements. All forms of the chain rule are special cases of this equation.

### 1.7 implicit function theorem -I

Suppose we have an equation of the form

\[f(x, y, z) = 0\]  \( (76) \)

Can we solve it for \( z \)? More precisely, is it the case that for each \( x, y \) in some domain there is a unique \( z \) so that \( f(x, y, z) = 0 \)? If so one can define an implicit function by \( z = \phi(x, y) \). Geometrically points satisfying \( f(x, y, z) = 0 \) are a surface and we are asking whether the surface is the graph of a function.

Suppose there is an implicit function. Then

\[f(x, y, \phi(x, y)) = 0\]  \( (77) \)

Assuming everything is differentiable we can take partial derivatives of this equation. By the chain rule we have

\[
0 = \frac{\partial}{\partial x} [f(x, y, \phi(x, y))] = f_x(x, y, \phi(x, y)) + f_z(x, y, \phi(x, y)) \phi_x(x, y)
\]

\[
0 = \frac{\partial}{\partial y} [f(x, y, \phi(x, y))] = f_y(x, y, \phi(x, y)) + f_z(x, y, \phi(x, y)) \phi_y(x, y)
\]  \( (78) \)

Solve this for \( \phi_x(x, y) \) and \( \phi_y(x, y) \) and get

\[
\phi_x(x, y) = -\frac{f_x(x, y, \phi(x, y))}{f_z(x, y, \phi(x, y))}
\]

\[
\phi_y(x, y) = -\frac{f_y(x, y, \phi(x, y))}{f_z(x, y, \phi(x, y))}
\]  \( (79) \)
This can also be written in the form

\[
\frac{\partial z}{\partial x} = -\frac{\partial f / \partial x}{\partial f / \partial z} \\
\frac{\partial z}{\partial y} = -\frac{\partial f / \partial y}{\partial f / \partial z}
\]  

(80)

keeping in mind that the right side is evaluated at \( z = \phi(x, y) \).

For this to work we need that \( \partial f / \partial z \neq 0 \). If this holds and if we restrict attention to a small region then there always is an implicit function. This is the content of the following

**Theorem 4** (*implicit function theorem*) Let \( f(x, y, z) \) have continuous partial derivatives near some point \((x_0, y_0, z_0)\). If

\[
\begin{align*}
f(x_0, y_0, z_0) &= 0 \\
f_z(x_0, y_0, z_0) &\neq 0
\end{align*}
\]  

(81)

Then for every \((x, y)\) near \((x_0, y_0)\) there is a unique \( z \) near \( z_0 \) such that \( f(x, y, z) = 0 \). The implicit function \( z = \phi(x, y) \) satisfies \( z_0 = \phi(x_0, y_0) \) and has continuous partial derivatives near \((x_0, y_0)\) which satisfy the equations (79).

The theorem does not tell you what \( \phi \) is or how to find it, only that it exists. However if we take the equations (79) at the special point \((x_0, y_0)\) we find

\[
\begin{align*}
\phi_x(x_0, y_0) &= -\frac{f_x(x_0, y_0, z_0)}{f_z(x_0, y_0, z_0)} \\
\phi_y(x_0, y_0) &= -\frac{f_y(x_0, y_0, z_0)}{f_z(x_0, y_0, z_0)}
\end{align*}
\]  

(82)

and these quantities can be computed.

**example:** Suppose the equation is

\[
f(x, y, z) = x^2 + y^2 + z^2 - 1 = 0
\]  

(83)

which describes the surface of a sphere. Is there an implicit function \( z = \phi(x, y) \) near a particular point \((x_0, y_0, z_0)\) on the sphere?

We have the derivatives

\[
\begin{align*}
f_x(x, y, z) &= 2x \\
f_y(x, y, z) &= 2y \\
f_z(x, y, z) &= 2z
\end{align*}
\]  

(84)
Then \( f_z(x_0, y_0, z_0) = 2z_0 \) is not zero if \( z_0 \neq 0 \). So by the theorem there is an implicit function \( z = \phi(x, y) \) near a point with \( z_0 \neq 0 \) and

\[
\begin{align*}
\frac{\partial z}{\partial x} &= -\frac{f_x}{f_z} = -\frac{x}{z} \\
\frac{\partial z}{\partial y} &= -\frac{f_y}{f_z} = -\frac{y}{z}
\end{align*}
\] (85)

This example is special in that we can find the implicit function exactly and so check these results. The implicit function is

\[
z = \pm \sqrt{1 - x^2 - y^2}
\] (86)

with the plus sign if \( z_0 > 0 \) and the minus sign if \( z_0 < 0 \). In either case we have the expected result

\[
\begin{align*}
\frac{\partial z}{\partial x} &= \pm \frac{1}{2} \left(1 - x^2 - y^2\right)^{-1/2}(-2x) = -\frac{x}{z} \\
\frac{\partial z}{\partial y} &= \pm \frac{1}{2} \left(1 - x^2 - y^2\right)^{-1/2}(-2y) = -\frac{y}{z}
\end{align*}
\] (87)

If \( z_0 = 0 \) there is no implicit function near the point as figure 3 illustrates.

**problem:** Consider the equation

\[
f(x, y, z) = x e^z + yz + 1 = 0
\] (88)

Note that \((x, y, z) = (0, 1, -1)\) is one solution. Is there an implicit function \( z = \phi(x, y) \) near this point? What are the derivatives \( \partial z/\partial x, \partial z/\partial y \) at \((x, y) = (0, 1)\)?

**solution** We have the derivatives

\[
\begin{align*}
f_x &= e^z & f_y &= z & f_z &= xe^z + y
\end{align*}
\] (89)

Then \( f_z(0, 1, -1) = 1 \) is not zero so by the theorem there is an implicit function. The derivatives at \((0,1)\) are

\[
\begin{align*}
\frac{\partial z}{\partial x} &= -\frac{f_x}{f_z} = -\frac{e^z}{xe^z + y} = -e^{-1} \\
\frac{\partial z}{\partial y} &= -\frac{f_y}{f_z} = -\frac{z}{xe^z + y} = 1
\end{align*}
\] (90)

**alternate solution:** Once the existence of the implicit function is established we can argue as follows. Take partial derivatives of \( xe^z + yz + 1 = 0 \) assuming \( z = \phi(x, y) \) and obtain

\[
\begin{align*}
\frac{\partial}{\partial x} : & \quad e^z + xe^z \frac{\partial z}{\partial x} + y \frac{\partial z}{\partial x} = 0 \\
\frac{\partial}{\partial y} : & \quad xe^z \frac{\partial z}{\partial y} + z + y \frac{\partial z}{\partial y} = 0
\end{align*}
\] (91)
Figure 3: There is an implicit function near points with $z \neq 0$, but not near points with $z = 0$.

Now put in the point $(x, y, z) = (0, 1, -1)$ and get

\[
\begin{align*}
e^{-1} + \frac{\partial z}{\partial x} &= 0 \\
-1 + \frac{\partial z}{\partial y} &= 0
\end{align*}
\]  

(92)

Solving for the derivatives gives the same result.

**Some remarks:**

1. Why is the implicit function theorem true? Instead of solving $f(x, y, z) = 0$ for $z$ near $(x_0, y_0, z_0)$ one could make a linear approximation and try to solve that. Taking into account that $f(x_0, y_0, z_0) = 0$ the linear approximation is

\[
f_z(x_0, y_0, z_0)(x - x_0) + f_y(x_0, y_0, z_0)(y - y_0) + f_z(x_0, y_0, z_0)(z - z_0) = 0
\]  

(93)

If $f_z(x_0, y_0, z_0) \neq 0$ this can be solved by

\[
z = z_0 - \frac{f_x(x_0, y_0, z_0)}{f_z(x_0, y_0, z_0)}(x - x_0) - \frac{f_y(x_0, y_0, z_0)}{f_z(x_0, y_0, z_0)}(y - y_0)
\]  

(94)
and this has the expected derivatives. To prove the theorem one has to argue that since the actual function is near the linear approximation there is an implicit function in that case as well.

2. Here are some variations of the implicit function theorem

(a) If \( f(x, y) = 0 \) one can solve for \( y = \phi(x) \) near any point where \( f_y \neq 0 \) and \( \frac{dy}{dx} = -\frac{f_x}{f_y} \).

(b) If \( f(x, y, z) = 0 \) one can solve for \( x = \phi(y, z) \) near any point where \( f_x \neq 0 \) and \( \frac{\partial x}{\partial y} = -\frac{f_y}{f_x} \) and \( \frac{\partial x}{\partial z} = -\frac{f_z}{f_x} \).

(c) If \( f(x, y, z, w) = 0 \) one can solve for \( w = \phi(x, y, z) \) near any point where \( f_w \neq 0 \) and \( \frac{\partial w}{\partial x} = -\frac{f_x}{f_w} \), etc.

3. One can also find higher derivatives of the implicit function. If \( f(x, y, z) = 0 \) defines \( z = \phi(x, y) \) and

\[
\phi_x(x, y) = -\frac{f_x(x, y, \phi(x, y))}{f_z(x, y, \phi(x, y))}
\]

then one can find \( \phi_{xx}, \phi_{xy} \) by further differentiation.

**example:** \( f(x, y, z) = x^2 + y^2 + z^2 - 1 = 0 \) defines \( z = \phi(x, y) \) which satisfies

\[
\frac{\partial z}{\partial x} = -\frac{x}{z} \tag{96}
\]

Keeping in mind that \( z \) is a function of \( x, y \) further differentiation yields

\[
\frac{\partial^2 z}{\partial x^2} = -\left( \frac{z - x \frac{\partial z}{\partial x}}{z^2} \right) = -\left( \frac{z - x(-x/z)}{z^2} \right) = -\frac{z^2 - x^2}{z^3} \tag{97}
\]

Since \( x^2 + y^2 + z^2 = 1 \) this can also be written as \( (y^2 - 1)/z^3 \). Similarly

\[
\frac{\partial^2 z}{\partial x \partial y} = -\frac{xy}{z^3} \tag{98}
\]

1.8 implicit function theorem -II

We give another version of the implicit function theorem. Suppose we have a pair of equations of the form

\[
F(x, y, u, v) = 0 \\
G(x, y, u, v) = 0
\]

(99)
Can we solve them for $u, v$ as functions of $x, y$? More precisely is it the case that for each $x, y$ in some domain there is a unique $u, v$ so that the equations are satisfied. If it is so we can define implicit functions

$$u = f(x, y)$$
$$v = g(x, y)$$

(100)

If the implicit functions exist we have

$$F(x, y, f(x, y), g(x, y)) = 0$$
$$G(x, y, f(x, y), g(x, y)) = 0$$

(101)

Take partial derivatives with respect to $x$ and $y$ and get

$$F_x + F_u f_x + F_v g_x = 0$$
$$G_x + G_u f_x + G_v g_x = 0$$
$$F_y + F_u f_y + F_v g_y = 0$$
$$G_y + G_u f_y + G_v g_y = 0$$

(102)

These equations can be written as the matrix equations

$$
\begin{pmatrix}
F_u & F_v \\
G_u & G_v
\end{pmatrix}
\begin{pmatrix}
f_x \\
g_x
\end{pmatrix} =
\begin{pmatrix}
-F_x \\
-G_x
\end{pmatrix}
$$

$$
\begin{pmatrix}
F_u & F_v \\
G_u & G_v
\end{pmatrix}
\begin{pmatrix}
f_y \\
g_y
\end{pmatrix} =
\begin{pmatrix}
-F_y \\
-G_y
\end{pmatrix}
$$

(103)

Note that the matrix on the left is the derivative of the function

$$(u, v) \to (F(x, y, u, v), G(x, y, u, v))$$

(104)

for fixed $(x, y)$. If the determinant of this matrix is not zero we can solve for the partial derivatives $f_x, g_x, f_y, g_y$. One finds for example

$$f_x = \frac{\det\left(\begin{array}{ccc}
-F_x & F_v \\
-G_x & G_v
\end{array}\right)}{\det\left(\begin{array}{ccc}
F_u & F_v \\
G_u & G_v
\end{array}\right)}$$

(105)

The determinant of the matrix is called the Jacobian determinant and is given a special symbol

$$\frac{\partial(F, G)}{\partial(u, v)} = \det\left(\begin{array}{ccc}
F_u & F_v \\
G_u & G_v
\end{array}\right) \equiv F_u G_v - F_v G_u$$

(106)
With this notation the equations for the four partial derivatives can be written

\[
\begin{align*}
\frac{\partial u}{\partial x} &= -\frac{\partial (F,G)}{\partial (x,v)} \frac{\partial (F,G)}{\partial (u,v)} \\
\frac{\partial v}{\partial x} &= -\frac{\partial (F,G)}{\partial (u,x)} \frac{\partial (F,G)}{\partial (u,v)} \\
\frac{\partial u}{\partial y} &= -\frac{\partial (F,G)}{\partial (y,v)} \frac{\partial (F,G)}{\partial (u,v)} \\
\frac{\partial v}{\partial y} &= -\frac{\partial (F,G)}{\partial (u,y)} \frac{\partial (F,G)}{\partial (u,v)}
\end{align*}
\] (107)

where \( u = f(x,y), v = g(x,y) \) on the right. This holds provided \( \frac{\partial (F,G)}{\partial (u,v)} \neq 0 \) and this is the key condition in the following theorem which guarantees the existence of the implicit functions.

**Theorem 5 (implicit function theorem)** Let \( F, G \) have continuous partial derivatives near some point \((x_0, y_0, u_0, v_0)\). If

\[
F(x_0, y_0, u_0, v_0) = 0 \\
G(x_0, y_0, u_0, v_0) = 0
\] (108)

and

\[
\left[ \frac{\partial (F,G)}{\partial (u,v)} \right] (x_0, y_0, u_0, v_0) \neq 0
\] (109)

Then for every \((x,y)\) near \((x_0, y_0)\) there is a unique \((u,v)\) near \((u_0, v_0)\) such that \( F(x,y,u,v) = 0 \) and \( G(x,y,u,v) = 0 \). The implicit functions \( u = f(x,y), v = g(x,y) \) satisfy \( u_0 = f(x_0,y_0), v_0 = g(x_0, y_0) \) and have continuous partial derivative which satisfy the equations (107).

**problem:** Consider the equations

\[
F(x, y, u, v) = x^2 - y^2 + 2uv - 2 = 0 \\
G(x, y, u, v) = 3x + 2xy + u^2 - v^2 = 0
\] (110)

Note that \((x,y,u,v) = (0,0,1,1)\) is a solution. Are there implicit functions \( u = f(x,y), v = g(x,y) \) near \((0,0)\)? What are the derivatives \( \partial u/\partial x, \partial v/\partial x \) at \((x,y) = (0,0)\)?

**solution:** First compute

\[
\frac{\partial (F,G)}{\partial (u,v)} = \text{det} \left( \begin{array}{cc}
F_u & F_v \\
G_u & G_v
\end{array} \right) = \text{det} \left( \begin{array}{cc}
2v & 2u \\
2u & -2v
\end{array} \right) = -4(u^2 + v^2) = -8
\] (111)
Since this is not zero the implicit functions exist by the theorem. We also compute
\[
\frac{\partial(F, G)}{\partial(x, v)} = \det \begin{pmatrix} F_x & F_v \\ G_x & G_v \end{pmatrix} = \det \begin{pmatrix} 2x & 2u \\ 3 + 2y & -2v \end{pmatrix} = -4xv - 2u(3 + 2y) = -6
\]
(112)
and
\[
\frac{\partial(F, G)}{\partial(u, x)} = \det \begin{pmatrix} F_u & F_x \\ G_u & G_x \end{pmatrix} = \det \begin{pmatrix} 2v & 2x \\ 2u & 3 + 2y \end{pmatrix} = 2v(3 + 2y) - 4ux = 6
\]
(113)
Then
\[
\frac{\partial u}{\partial x} = \frac{-6}{-8} = \frac{3}{4}, \quad \frac{\partial v}{\partial x} = \frac{-6}{-8} = \frac{3}{4}
\]
(114)

**alternate solution:** Assuming the implicit functions exist differentiate the equations \( F = 0, G = 0 \) with respect to \( x \) assuming \( u = f(x, y), v = g(x, y) \). This gives
\[
\begin{align*}
2x + 2v \frac{\partial u}{\partial x} + 2u \frac{\partial v}{\partial x} &= 0 \\
3 + 2y + 2u \frac{\partial u}{\partial x} - 2v \frac{\partial v}{\partial x} &= 0
\end{align*}
\]
(115)
Now put in the point \((0, 0, 1, 1)\) and get
\[
\begin{align*}
2 \frac{\partial u}{\partial x} + 2 \frac{\partial v}{\partial x} &= 0 \\
3 + 2 \frac{\partial u}{\partial x} - 2 \frac{\partial v}{\partial x} &= 0
\end{align*}
\]
(116)
This again has the solution \( \frac{\partial u}{\partial x} = -3/4, \frac{\partial v}{\partial x} = 3/4 \).

### 1.9 inverse functions

Suppose we have a function \( f \) from \( U \subseteq \mathbb{R}^2 \) to \( V \subseteq \mathbb{R}^2 \) which we write in the form
\[
\begin{align*}
x &= f_1(u, v) \\
y &= f_2(u, v)
\end{align*}
\]
(117)
Suppose further for every \((x, y)\) in \( V \) there is a unique \((u, v)\) in \( U \) such that \( x = f_1(u, v), y = f_2(u, v) \). (One says the function is one-to-one and onto.) Then there is a an inverse function \( g \) from \( V \) to \( U \) defined by
\[
\begin{align*}
u &= g_1(x, y) \\
v &= g_2(x, y)
\end{align*}
\]
(118)
So $g$ sends every point back where it came from. See figure 4. We have

$$(g \circ f)(u, v) = (u, v)$$

$$(f \circ g)(x, y) = (x, y)$$

(119)

We write $g = f^{-1}$.

**example:** Suppose the function is

$$x = au + bv$$

$$y = cu + dv$$

(120)

defined on all of $\mathbb{R}^2$. This can also be written in matrix form

$$\begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} a & b \\ c & d \end{pmatrix} \begin{pmatrix} u \\ v \end{pmatrix}$$

(121)

This is invertible if we can solve the equation for $(u, v)$ which is possible if and only if

$$\det \begin{pmatrix} a & b \\ c & d \end{pmatrix} = ad - bc \neq 0$$

(122)
The inverse function is given by the inverse matrix

\[
\begin{pmatrix} u \\ v \end{pmatrix} = \begin{pmatrix} a & b \\ c & d \end{pmatrix}^{-1} \begin{pmatrix} x \\ y \end{pmatrix}
\]  

(123)

where

\[
\begin{pmatrix} a & b \\ c & d \end{pmatrix}^{-1} = \frac{1}{ad - bc} \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}
\]  

(124)

**example:** Consider the function

\[
x = r \cos \theta \\
y = r \sin \theta
\]  

(125)

from

\[U = \{(r, \theta) : r > 0, 0 \leq \theta < 2\pi\}\]  

(126)

to

\[V = \{(x, y) : (x, y) \neq 0\}\]  

(127)

This function sends \((r, \theta)\) to the point with polar coordinates \((r, \theta)\). See figure 5. The function is invertible since every point \((x, y)\) in \(V\) has unique polar coordinates \((r, \theta)\) in \(U\). (It would not be invertible if we took \(U = \mathbb{R}^2\) since \((r, \theta)\) and \((r, \theta + 2\pi)\) are sent to the same point ). For \((x, y)\) in the first quadrant the inverse is

\[
r = \sqrt{x^2 + y^2} \\
\theta = \tan^{-1} \left( \frac{y}{x} \right)
\]  

(128)

### 1.10 inverse function theorem

Continuing the discussion of the previous section suppose that \(f\) has an inverse function \(g\) and that both are differentiable. Then differentiating \((f \circ g)(x, y) = (x, y)\) we find by the chain rule

\[
(Df)(g(x, y)) \cdot (Dg)(x, y) = I \quad \quad \quad I = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}
\]  

(129)

and so the derivative of the inverse function is the matrix inverse

\[
(Dg)(x, y) = [(Df)(g(x, y))]^{-1}
\]  

(130)

It is not always easy to tell whether an inverse function exists. The following theorem can be helpful.
Theorem 6 (inverse function theorem) Let \( x = f_1(u, v) \), \( y = f_2(u, v) \) have continuous partial derivatives and suppose
\[
x_0 = f_1(u_0, v_0) \\
y_0 = f_2(u_0, v_0)
\] (131)
and
\[
\left( \frac{\partial(x, y)}{\partial(u, v)} \right)_{(u_0, v_0)} \neq 0
\] (132)
Then there is an inverse function \( u = g_1(x, y) \), \( v = g_2(x, y) \) defined near \((x_0, y_0)\) which satisfies
\[
u_0 = g_1(x_0, y_0) \\
v_0 = g_2(x_0, y_0)
\] (133)
and has a continuous derivative which satisfies (130). In particular \( Dg(x_0, y_0) = [Df(u_0, v_0)]^{-1} \) or
\[
\begin{pmatrix}
\frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\
\frac{\partial v}{\partial x} & \frac{\partial v}{\partial y}
\end{pmatrix}_{(x_0, y_0)} = \left( \begin{pmatrix}
\frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\
\frac{\partial y}{\partial u} & \frac{\partial y}{\partial v}
\end{pmatrix}_{(u_0, v_0)} \right)^{-1}
\] (134)
Proof. The inverse exists if for each \((x, y)\) near \((x_0, y_0)\) there is a unique \((u, v)\) near \((u_0, v_0)\) such that

\[
F(x, y, u, v) \equiv f_1(u, v) - x = 0 \\
G(x, y, u, v) \equiv f_2(u, v) - y = 0
\]  (135)

This follows by the implicit functions theorem since \((x_0, y_0, u_0, v_0)\) is one solution and at this point

\[
\frac{\partial(F, G)}{\partial(u, v)} = \det \begin{pmatrix} F_u & F_v \\ G_u & G_v \end{pmatrix} = \det \begin{pmatrix} f_{1,u} & f_{1,v} \\ f_{2,u} & f_{2,v} \end{pmatrix} = \frac{\partial(x, y)}{\partial(u, v)} \neq 0
\]

The differentiability of the inverse also follows from the implicit function theorem.

Problem: Consider the function

\[
x = u + v^2 \\
y = u^2 + v
\]  (137)

which sends \((u, v) = (1, 2)\) to \((x, y) = (5, 3)\). Show that the function is invertible near \((1, 2)\) and find the partial derivatives of the inverse function at \((5, 3)\).

Solution: We have

\[
\begin{pmatrix}
\frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\
\frac{\partial y}{\partial u} & \frac{\partial y}{\partial v}
\end{pmatrix} = \begin{pmatrix} 1 & 2v \\ 2u & 1 \end{pmatrix} = \begin{pmatrix} 1 & 4 \\ 2 & 1 \end{pmatrix}
\]

at \((1, 2)\)  (138)

Therefore

\[
\frac{\partial(x, y)}{\partial(u, v)} = \det \begin{pmatrix} 1 & 4 \\ 2 & 1 \end{pmatrix} = -7
\]

at \((1, 2)\)  (139)

This is not zero so the inverse exists by the theorem and sends \((5, 3)\) to \((1, 2)\). We have for the derivatives

\[
\begin{pmatrix}
\frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\
\frac{\partial v}{\partial x} & \frac{\partial v}{\partial y}
\end{pmatrix}^{(5,3)} = \begin{pmatrix}
\frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\
\frac{\partial y}{\partial u} & \frac{\partial y}{\partial v}
\end{pmatrix}^{-1}^{(1,2)} = \begin{pmatrix}
1 & 4 \\ 2 & 1
\end{pmatrix}^{-1} = \begin{pmatrix}
-1/7 & 4/7 \\ 2/7 & -1/7
\end{pmatrix}
\]

Alternate solution: Differentiate \(x = u + v^2, y = u^2 + v\) assuming \(u, v\) are functions of \(x, y\), then put in the point.
1.11 maxima and minima
Let \( f(x) = f(x_1, \ldots, x_n) \) be a function from \( \mathbb{R}^n \) to \( \mathbb{R} \) and let \( x^0 = (x_1^0, \ldots, x_n^0) \) be a point in \( \mathbb{R}^n \). We say that \( f \) has a local maximum at \( x^0 \) if \( f(x) \leq f(x^0) \) for all \( x \) near \( x^0 \). We say that \( f \) has a local minimum at \( x^0 \) if \( f(x) \geq f(x^0) \) for all \( x \) near \( x^0 \).

**Theorem 7** If \( f \) is differentiable at \( x_0 \) and \( f \) has a local maximum or minimum at \( x^0 \) then all partial derivatives vanish at the point, i.e.

\[
f_{x_1}(x^0) = \cdots = f_{x_n}(x^0) = 0
\]

**Proof.** For any \( h = (h_1, \ldots, h_n) \in \mathbb{R}^n \) consider the function \( F(t) = f(x^0 + th) \) of the single variable \( t \). If \( f \) has a local maximum or minimum at \( x^0 \) then \( F \) has a local maximum or minimum at \( t = 0 \). By the chain rule \( F(t) \) is differentiable and it is a result of elementary calculus that \( F'(0) = 0 \). But the chain rule says

\[
F'(t) = \sum_{i=1}^{n} f_{x_i}(x^0 + th) \frac{d(x_i^0 + th_i)}{dt} = \sum_{i=1}^{n} f_{x_i}(x^0 + th)h_i
\]

Thus

\[
0 = F'(0) = \sum_{i=1}^{n} f_{x_i}(x^0)h_i
\]

Since this is true for any \( h \) it must be that \( f_{x_i}(x^0) = 0 \).

A point \( x^0 \) with \( f_{x_i}(x^0) = 0 \) is called a critical point for \( f \). We have seen that if \( f \) is has a local maximum or minimum at \( x^0 \) then it is a critical point. We are interested whether the converse is true. If \( x^0 \) is a critical point is it a local maximum or minimum for \( f \)? Which is it?

To answer this question consider again \( F(t) = f(x^0 + th) \) and suppose \( f \) is many times differentiable. By Taylor’s theorem for one variable we have

\[
F(1) = F(0) + F'(0) + \frac{1}{2}F''(0) + \frac{1}{6}F'''(s)
\]

for some \( s \) between 0 and 1. But \( F'(t) \) is computed above, and similarly we have

\[
F''(t) = \sum_{i=1}^{n} \sum_{j=1}^{n} f_{x_ix_j}(x^0 + th)h_ih_j
\]

\[
F'''(t) = \sum_{i=1}^{n} \sum_{j=1}^{n} \sum_{k=1}^{n} f_{x_ix_jx_k}(x^0 + th)h_ih_jh_k
\]
Then the Taylor’s expansion becomes
\[ f(x^0 + h) = f(x^0) + \sum_{i=1}^{n} f_{x_i}(x^0)h_i + \frac{1}{2} \sum_{i=1}^{n} \sum_{j=1}^{n} f_{x_i,x_j}(x^0)h_ih_j + R(h) \] (146)

This is an example of a multivariable Taylor’s theorem with remainder. The remainder \( R(h) = F'''(s)/6 \) is small if \( h \) is small and one can show that there is a constant \( C \) such that for \( h \) small \( |R(h)| \leq C|h|^3 \).

Now suppose \( x^0 \) is a critical point so the first derivatives vanish. Also define a matrix of second derivatives
\[
A = \begin{pmatrix}
  f_{x_1,x_1}(x^0) & \cdots & f_{x_1,x_n}(x^0) \\
  \vdots & \ddots & \vdots \\
  f_{x_n,x_1}(x^0) & \cdots & f_{x_n,x_n}(x^0)
\end{pmatrix}
\] (147)
called the Hessian of \( f \) at \( x^0 \). Then we can write our expansion as
\[ f(x^0 + h) = f(x^0) + \frac{1}{2} h \cdot Ah + R(h) \] (148)

We are interested in whether \( f(x^0 + h) \) is greater than or less than \( f(x^0) \) for \( |h| \) small. Then idea is that since \( R(h) \) is much smaller than \( \frac{1}{2} h \cdot Ah \) it is the latter term which determines the answer.

A \( n \times n \) matrix \( A \) is called positive definite if there is a constant \( M \) so \( h \cdot Ah > M|h|^2 \) for all \( h \neq 0 \). It is called negative definite if \( h \cdot Ah < -M|h|^2 \) for all \( h \neq 0 \).

**Theorem 8** Let \( x^0 \) be a critical point for \( f \) and suppose the Hessian \( A \) has \( \det A \neq 0 \).

1. If \( A \) is positive definite then \( f \) has a local minimum at \( x^0 \).

2. If \( A \) is negative definite then \( f \) has a local maximum at \( x^0 \).

3. Otherwise \( f \) has a saddle point at \( x^0 \), i.e \( f \) increases in some directions and decreases in other directions as you move away from \( x^0 \).

**Proof.** We prove the first statement. We have \( h \cdot Ah > M|h|^2 \). If also \( |h| \leq M/4C \), then
\[ |R(h)| \leq C|h|^3 \leq \frac{1}{4} M|h|^2 \] (149)

Therefore
\[ f(x^0 + h) \geq f(x^0) + \frac{1}{2} M|h|^2 - \frac{1}{4} M|h|^2 \] (150)
or
\[ f(x^0 + h) \geq f(x^0) + \frac{1}{4} M|h|^2 \] (151)
Thus \( f(x^0 + h) > f(x^0) \) for \( 0 < |h| \leq M/4C \) which means we have a strict local maximum at \( x^0 \).

Thus our problem is to decide whether or not \( A \) is positive or negative definite. For any symmetric matrix \( A \) one can show that \( A \) is positive definite if and only if all eigenvalues are positive and \( A \) is negative definite if and only if all eigenvalues are negative. Recall that \( \lambda \) is an eigenvalue if there is a vector \( v \neq 0 \) such that \( Av = \lambda v \). One can find the eigenvalues by solving the equation

\[
\det(A - \lambda I) = 0
\]

where \( I \) is the identity matrix.

**Example:** Consider the function

\[
f(x, y) = \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right)
\]

The critical points are solutions of

\[
\begin{align*}
f_x &= (-x + 1) \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right) = 0 \\
f_y &= (-y + 1) \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right) = 0
\end{align*}
\]

Thus the only critical point is \((x, y) = (1, 1)\). The second derivatives at this point are

\[
\begin{align*}
f_{xx} &= (x^2 - 2x) \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right) = -e \\
f_{xy} &= (-x + 1)(-y + 1) \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right) = 0 \\
f_{yy} &= (y^2 - 2y) \exp\left(-\frac{x^2}{2} - \frac{y^2}{2} + x + y\right) = -e
\end{align*}
\]

Thus the Hessian is

\[
A = \begin{pmatrix} f_{xx} & f_{xy} \\ f_{yx} & f_{yy} \end{pmatrix} = \begin{pmatrix} -e & 0 \\ 0 & -e \end{pmatrix}
\]

It has eigenvalues \(-e, -e\) which are negative. Hence \( A \) is negative definite and the function has a local maximum at \((1, 1)\).

**Example:** Consider the function

\[
f(x, y) = x \sin y
\]
The critical points are solutions of
\begin{align}
  f_x &= \sin y = 0 \\
  f_y &= x \cos y = 0
\end{align}

These are points with \( x = 0 \) and \( y = 0, \pm \pi, \pm 2\pi, \ldots \). The Hessian is
\begin{equation}
  A = \begin{pmatrix}
    f_{xx} & f_{xy} \\
    f_{yx} & f_{yy}
  \end{pmatrix} = \begin{pmatrix}
    0 & \cos y \\
    \cos y & -x \sin y
  \end{pmatrix}
\end{equation}

At points \( (0, \pm \pi), (0, \pm 3\pi), \ldots \) this is
\begin{equation}
  A = \begin{pmatrix}
    0 & -1 \\
    -1 & 0
  \end{pmatrix}
\end{equation}

At points \( (0, 0), (0, \pm 2\pi), \ldots \) this is
\begin{equation}
  A = \begin{pmatrix}
    0 & 1 \\
    1 & 0
  \end{pmatrix}
\end{equation}

In either case the eigenvalues are solutions of
\begin{equation}
  \det(A - \lambda I) = \det\begin{pmatrix}
    -\lambda & \pm 1 \\
    \pm 1 & -\lambda
  \end{pmatrix} = \lambda^2 - 1 = 0
\end{equation}

Thus the eigenvalues are \( \lambda = \pm 1 \). Since they have both signs every critical point is a saddle point.

1.12 differentiation under the integral sign

**Theorem 9** If \( f(t, x), (\partial f/\partial t)(t, x) \) exist and are continuous
\begin{equation}
  \frac{d}{dt} \left[ \int_a^b f(t, x) \, dx \right] = \int_a^b \frac{\partial f}{\partial t}(t, x) \, dx
\end{equation}

**Proof.** Form the difference quotient
\begin{equation}
  \frac{\int_a^b f(t, x + h) - \int_a^b f(t, x)}{h} = \int_a^b \frac{f(t, x + h) - f(t, x)}{h}
\end{equation}

and take the limit \( h \to 0 \). The only issue is whether we can take the limit under the integral sign on the right. This can be justified under the hypotheses of the theorem.
example:

\[
\frac{d}{dt} \left[ \int_0^1 \log(x^2 + t^2) \, dx \right] = \int_0^1 \frac{2t}{x^2 + t^2} \, dx \\
= \left[ 2 \tan^{-1} \left( \frac{x}{t} \right) \right]_{x=0}^{x=1} \\
= 2 \tan^{-1} \left( \frac{1}{t} \right)
\]  

(165)

problem: Find

\[
\int_0^1 \frac{\sqrt{x} - 1}{\log x} \, dx
\]  

(166)

solution: We solve a more general problem which is to evaluate

\[
\phi(t) = \int_0^1 \frac{x^t - 1}{\log x} \, dx
\]  

(167)

Then \(\phi(1/2)\) is the answer to the original problem. Differentiating under the integral sign and taking account that \(d(x^t)/dt = x^t \log x\) we have

\[
\phi'(t) = \int_0^1 x^t \, dx = \left[ \frac{x^{t+1}}{t+1} \right]_{x=0}^{x=1} = \frac{1}{t+1}
\]  

(168)

Hence

\[
\phi(t) = \log(t + 1) + C
\]  

(169)

for some constant \(C\). But we know \(\phi(0) = 0\) so we must have \(C = 0\). Thus \(\phi(t) = \log(t + 1)\) and the answer is \(\phi(1/2) = \log(3/2)\).

1.13 Leibniz’ rule

The following is a generalization of the previous result where we allow the endpoints to be functions of \(t\).

Theorem 10

\[
\frac{d}{dt} \left[ \int_{a(t)}^{b(t)} f(t, x) \, dx \right] = f\left(t, b(t)\right)b'(t) - f\left(t, a(t)\right)a'(t) + \int_{a(t)}^{b(t)} \frac{\partial f}{\partial t}(t, x) \, dx
\]  

(170)
Proof. Let

$$F(b,a,t) = \int_a^b f(t,x) \, dx$$

(171)

What we want is the derivative of $F(b(t),a(t),t)$ and by the chain rule this is

$$\frac{d}{dt}[F(b(t),a(t),t)] = F_b(b(t),a(t),t)b'(t) + F_a(b(t),a(t),t)a'(t) + F_t(b(t),a(t),t)$$

(172)

But

$$F_b(b,a,t) = f(t,b) \quad F_a(b,a,t) = -f(t,a) \quad F_t(b,a,t) = \int_a^b \frac{\partial f}{\partial t}(t,x) \, dx$$

(173)

which gives the result.

example:

$$\frac{d}{dt}\left[ \int_t^{t^2} \sin(t^2 + x^2) \, dx \right] = \sin(t^2 + x^2)|_{x=t^2} \frac{d(t^2)}{dt} - \sin(t^2 + x^2)|_{x=t} \frac{d(t)}{dt} + \int_t^{t^2} \frac{\partial}{\partial t} \sin(t^2 + x^2) \, dx$$

(174)

$$= \sin(t^2 + t^4)2t - \sin(2t^2) + 2t \int_t^{t^2} \cos(t^2 + x^2) \, dx$$

example: (forced harmonic oscillator). Suppose we want to find a function $x(t)$ whose derivatives $x'(t), x''(t)$ solve the ordinary differential equation

$$mx'' + kx = f(t)$$

(175)

with the initial conditions

$$x(0) = 0 \quad x'(0) = 0$$

(176)

Here $k, m$ are positive constants and $f(t)$ is an arbitrary function.

We claim that a solution is

$$x(t) = \frac{1}{m \omega} \int_0^t \sin(\omega(t - \tau)) f(\tau) \, d\tau \quad \omega = \sqrt{\frac{k}{m}}$$

(177)

To check this note first that $x(0) = 0$. Then take the first derivative using the Leibniz rule and find

$$x'(t) = \frac{1}{m \omega}[\sin(\omega(t - \tau)) f(\tau)]_{\tau=t} + \frac{1}{m \omega} \int_0^t \omega \cos(\omega(t - \tau)) f(\tau) \, d\tau$$

(178)
The first term is zero and we also note that \( x'(0) = 0 \). Taking one more derivative yields
\[
x''(t) = \frac{1}{m} [\cos(\omega(t - \tau))] f(\tau)_{\tau=t} + \frac{1}{m\omega} \int_{0}^{t} (\omega^{2}) \sin(\omega(t - \tau)) f(\tau) d\tau
\]
which is the same as
\[
x''(t) = \frac{1}{m} f(t) - \omega^{2} x(t)
\]
Now multiply by \( m \), use \( m\omega^{2} = k \) and we recover the differential equation.

### 1.14 calculus of variations

We consider the problem of finding maxima and minima for a functional - i.e. for a function of functions.

As an example consider the problem of finding the curve between two points \((x_0, y_0)\) and \((x_1, y_1)\) which has the shortest length. We assume that there is such a curve and that it is the graph of a function \( y = y(x) \) with \( y(x_0) = y_0 \) and \( y(x_1) = y_1 \). The length of such a curve is
\[
I(y) = \int_{x_0}^{x_1} \sqrt{1 + (y'(x))^2} \, dx
\]
and we seek to find the function \( y = y(x) \) which gives the minimum value.

More generally suppose we have a function \( F(x, y, y') \) of three real variables \((x, y, y')\). For any differentiable function \( y = y(x) \) satisfying \( y(x_0) = y_0 \) and \( y(x_1) = y_1 \), form the integral
\[
I(y) = \int_{x_0}^{x_1} F(x, y(x), y'(x)) \, dx
\]
The question is which function \( y = y(x) \) minimizes (or maximizes) the functional \( I(y) \). We are looking for a local minimum (or maximum), that is we want to find a function \( y \) such that \( I(\tilde{y}) \geq I(y) \) for all functions \( \tilde{y} \) near to \( y \) in the sense that \( \tilde{y}(x) \) is close to \( y(x) \) for all \( x_0 \leq x \leq x_1 \).

**Theorem 11** If \( y = y(x) \) is a local maximum or minimum for \( I(y) \) with fixed endpoints then it satisfies
\[
F_y(x, y(x), y'(x)) - \frac{d}{dx} (F_{y'}(x, y(x), y'(x))) = 0
\]

**Remarks.**

1. The equation is called *Euler’s equation* and is abbreviated as
\[
F_y - \frac{d}{dx} F_{y'} = 0
\]
2. Note that \( y, y' \) mean two different things in Euler’s equation. First one evaluates the partial derivatives \( F_y, F_{y'} \) treating \( y, y' \) as independent variables. Then in computing \( d/dx \) one treats \( y, y' \) as a function and its derivative.

3. The converse may not be true, i.e solutions of Euler’s equation are not necessarily maxima or minima. They are just candidates and one should decide by other criteria.

**proof:** Suppose \( y = y(x) \) is a local minimum. Pick any differentiable function \( \eta(x) \) defined for \( x_0 \leq x \leq x_1 \) and satisfying \( \eta(x_0) = \eta(x_1) = 0 \). Then for any real number \( t \) the function

\[
y_t(x) = y(x) + t\eta(x)
\]

is differentiable with respect to \( x \) and satisfies \( y_t(x_0) = y_t(x_1) = 0 \). Consider the function

\[
J(t) = I(y_t) = \int_{x_0}^{x_1} F(x, y_t(x), y'_t(x)) dx
\]

Since \( y_t \) is near \( y_0 = y \) for \( t \) small, and since \( y_0 \) is a local minimum for \( I \), we have

\[
J(t) = I(y_t) \geq I(y_0) = J(0)
\]

Thus \( t = 0 \) is a local minimum for \( J(t) \) and it follows that \( J'(0) = 0 \).

To see what this says we differentiate under the integral sign and compute

\[
J'(t) = \int_{x_0}^{x_1} \frac{\partial}{\partial t} F(x, y_t(x), y'_t(x)) dx
\]

\[
= \int_{x_0}^{x_1} \left( F_y(x, y_t(x), y'_t(x)) \frac{\partial}{\partial t}(y_t(x)) + F_{y'}(x, y_t(x), y'_t(x)) \frac{\partial}{\partial t}(y'_t(x)) \right) dx
\]

\[
= \int_{x_0}^{x_1} (F_y(x, y_t(x), y'_t(x)) \eta(x) + F_{y'}(x, y_t(x), y'_t(x)) \eta'(x)) dx
\]

Here in the second step we have used the chain rule. Now in the second term integrate by parts taking the derivative off \( \eta'(x) = (d/dx)\eta \) and putting it on \( F_{y'}(x, y_t(x), y'_t(x)) \). The term involving the endpoints vanishes because \( \eta(x_0) = \eta(x_1) = 0 \). Then we have

\[
J'(t) = \int_{x_0}^{x_1} \left( F_y(x, y_t(x), y'_t(x)) - \frac{d}{dx} F_{y'}(x, y_t(x), y'_t(x)) \right) \eta(x) dx
\]

Now put \( t = 0 \) and get

\[
0 = J'(0) = \int_{x_0}^{x_1} \left( F_y(x, y(x), y'(x)) - \frac{d}{dx} F_{y'}(x, y(x), y'(x)) \right) \eta(x) dx
\]
Since this is true for an arbitrary function \( \eta \) it follows that the expression in parentheses must be zero which is our result. \(^1\)

**example:** We return to the problem of finding the curve between two points with the shortest length. That is we seek to minimize

\[
I(y) = \int_{x_0}^{x_1} \sqrt{1 + y'(x)^2} \, dx
\]  

(191)

This has the form (182) with

\[
F(x, y, y') = \sqrt{1 + (y')^2}
\]  

(192)

The minimizer must satisfy Euler’s equation. Since \( F_y = 0 \) and \( F_{yy'} = \frac{y'}{\sqrt{1 + (y')^2}} \) this says

\[
F_y - \frac{d}{dx} F_{yy'} = \frac{d}{dx} \left[ \frac{y'}{\sqrt{1 + (y')^2}} \right] = 0
\]  

(193)

Evaluating the derivatives yields

\[
\frac{\sqrt{1 + (y')^2}y'' - (y')^2 y'' / \sqrt{1 + (y')^2}}{1 + (y')^2} = 0
\]  

(194)

Now multiple by \((1 + (y')^2)^{3/2}\) and get

\[
(1 + (y')^2)y'' - (y')^2 y'' = 0
\]  

(195)

which is the same as \( y'' = 0 \). Thus the minimizer must have the form \( y = ax + b \) for some constants \( a, b \). So the shortest distance between two points is along a straight line as expected.

**example:** The problem is to find the function \( y = y(x) \) with \( y(0) = 0 \) and \( y(1) = 1 \) which minimizes the integral

\[
I(y) = \frac{1}{2} \int_0^1 (y(x)^2 + (y'(x))^2) \, dx
\]  

(196)

and again we assume there is such a minimizer. The integral has the form (182) with

\[
F(x, y, y') = \frac{1}{2} (y^2 + (y')^2)
\]  

(197)

---

\(^1\)In general if \( f \) is a continuous function \( \int_a^b f(x) \, dx = 0 \) does not imply that \( f \equiv 0 \). However it is true if \( f(x) \geq 0 \). If \( \int_a^b f(x) \eta(x) \, dx = 0 \) for any continuous function \( \eta \) then we can take \( \eta(x) = f(x) \) and get \( \int_a^b f(x)^2 \, dx = 0 \), hence \( f(x)^2 = 0 \), hence \( f(x) = 0 \). This is not quite the situation above since we also restricted \( \eta \) to vanish at the endpoints. But the conclusion is still valid.
Then Euler’s equation says

\[ F_y - \frac{d}{dx}(F_y') = y - \frac{d}{dx}y' = y - y'' = 0 \tag{198} \]

Thus we must solve the second order equation \( y - y'' = 0 \). Since the equation has constant coefficients one can find solutions by trying \( y = e^{rx} \). One finds that \( r^2 = 1 \) and so \( y = e^{\pm x} \) are solutions. The general solution is

\[ y(x) = c_1 e^x + c_2 e^{-x} \tag{199} \]

The constants \( c_1, c_2 \) are fixed by the condition \( y(0) = 0 \) and \( y(1) = 1 \) and one finds

\[ y(x) = \frac{e}{e^2 - 1}(e^x - e^{-x}) = \frac{2e}{e^2 - 1} \sinh x \tag{200} \]

**example**: Suppose that an object moves on a line and its position at time \( t \) is given by a function \( x(t) \). Suppose also we know that \( x(t_0) = x_0 \) and \( x(t_1) = x_1 \) and that it is moving in a force field \( F(x) = -dV/dx \) determined by some potential function \( V(x) \). What is the trajectory \( x(t) \)?

One way to proceed is to form a function called the Lagrangian by taking the difference of the kinetic and potential energy:

\[ L(x, x') = \frac{1}{2}m(x')^2 - V(x) \tag{201} \]

For any trajectory \( x = x(t) \) one forms the action

\[ A(x) = \int_{t_0}^{t_1} L(x(t), x'(t))dt \tag{202} \]

According to D’Alembert’s principle the actual trajectory is the one that minimizes the action. This is also called the principle of least action.

To see what it says we observe that this problem has the form (182) with new names for the variables. Euler’s equation says

\[ L_x - \frac{d}{dt}L_{x'} = 0 \tag{203} \]

But \( L_x = -dV/dx = F \) and \( L_{x'} = mx' \) so this is

\[ F - mx'' = 0 \tag{204} \]

which is Newton’s second law. Thus the principle of least action is an alternative to Newton’s second law. This turns out to be true for many dynamical problems.
2 vector calculus

2.1 vectors

We continue our survey of multivariable calculus but now put special emphasis on $\mathbb{R}^3$ which is a model for physical space.

Vectors in $\mathbb{R}^3$ will now be indicated by arrows or bold face type as in $\mathbf{u} = (u_1, u_2, u_3)$. Any such vector can be written

$$\mathbf{u} = (u_1, u_2, u_3) = u_1(1, 0, 0) + u_2(0, 1, 0) + u_3(0, 0, 1)$$

where we have defined

$$\mathbf{i} = (1, 0, 0) \quad \mathbf{j} = (0, 1, 0) \quad \mathbf{k} = (0, 0, 1)$$

Any vector can be written as a linear combination of the independent vectors $\mathbf{i}, \mathbf{j}, \mathbf{k}$ so these form a basis for $\mathbb{R}^3$ called the standard basis.

We consider several products of vectors:

**dot product:** The dot product is defined either by

$$\mathbf{u} \cdot \mathbf{v} = u_1v_1 + u_2v_2 + u_3v_3$$

or by

$$\mathbf{u} \cdot \mathbf{v} = |\mathbf{u}||\mathbf{v}|\cos \theta$$

where $\theta$ is the angle between $\mathbf{u}$ and $\mathbf{v}$. Note that $\mathbf{u} \cdot \mathbf{u} = |\mathbf{u}|^2$. Also note that $\mathbf{u}$ is orthogonal (perpendicular) to $\mathbf{v}$ if and only if $\mathbf{u} \cdot \mathbf{v} = 0$.

The dot product has the properties

$$\mathbf{u} \cdot \mathbf{v} = \mathbf{v} \cdot \mathbf{u}$$

$$(\alpha \mathbf{u}) \cdot \mathbf{v} = \alpha (\mathbf{u} \cdot \mathbf{v}) = \mathbf{u} \cdot (\alpha \mathbf{v})$$

$$(\mathbf{u}_1 + \mathbf{u}_2) \cdot \mathbf{v} = \mathbf{u}_1 \cdot \mathbf{v} + \mathbf{u}_2 \cdot \mathbf{v}$$

Examples are

$$\mathbf{i} \cdot \mathbf{i} = 1 \quad \mathbf{j} \cdot \mathbf{j} = 1 \quad \mathbf{j} \cdot \mathbf{j} = 1$$

$$\mathbf{i} \cdot \mathbf{j} = 0 \quad \mathbf{j} \cdot \mathbf{k} = 0 \quad \mathbf{k} \cdot \mathbf{i} = 0$$

This says that $\mathbf{i}, \mathbf{j}, \mathbf{k}$ are orthogonal unit vectors. They are an example of an orthonormal basis for $\mathbb{R}^3$. 

cross product: (only in \( \mathbb{R}^3 \)) The cross product of \( \mathbf{u} \) and \( \mathbf{v} \) is a vector \( \mathbf{u} \times \mathbf{v} \) which has length

\[ |\mathbf{u} \times \mathbf{v}| = |\mathbf{u}| |\mathbf{v}| \sin \theta \]  

Here \( \theta \) is the positive angle between the vectors. The direction of \( \mathbf{u} \times \mathbf{v} \) is specified by requiring that it be perpendicular to \( \mathbf{u} \) and \( \mathbf{v} \) in such a way that \( \mathbf{u}, \mathbf{v}, \mathbf{u} \times \mathbf{v} \) form a right-handed system. (See figure 6)

The length \( |\mathbf{u} \times \mathbf{v}| \) is interpreted as the area of the parallelogram spanned by \( \mathbf{u}, \mathbf{v} \). This follows since the parallelogram has base \( |\mathbf{u}| \) and height \( |\mathbf{v}| \sin \theta \) (See figure 6) and so

\[ \text{area} = \text{base} \times \text{height} \]
\[ = |\mathbf{u}| |\mathbf{v}| \sin \theta \]  

\[ = |\mathbf{u} \times \mathbf{v}| \]  

An alternate definition of the cross-product uses determinants. Recall that

\[
\det \begin{pmatrix} a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \\ c_1 & c_2 & c_3 \end{pmatrix} = a_1 \det \begin{pmatrix} b_2 & b_3 \\ c_2 & c_3 \end{pmatrix} - a_2 \det \begin{pmatrix} b_1 & b_3 \\ c_1 & c_3 \end{pmatrix} + a_3 \det \begin{pmatrix} b_1 & b_2 \\ c_1 & c_2 \end{pmatrix}
\]

\[ = a_1(b_2c_3 - b_3c_2) + a_2(b_3c_1 - b_1c_3) + a_3(b_1c_2 - b_2c_1) \]  

(213)
The other definition is
\[ \mathbf{u} \times \mathbf{v} = \det \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ u_1 & u_2 & u_3 \\ v_1 & v_2 & v_3 \end{pmatrix} = \mathbf{i}(u_2v_3 - u_3v_2) + \mathbf{j}(u_3v_1 - u_1v_3) + \mathbf{k}(u_1v_2 - u_2v_1) \] (214)

The cross product has the following properties
\[ \mathbf{u} \times \mathbf{u} = 0 \]
\[ \mathbf{u} \times \mathbf{v} = -\mathbf{v} \times \mathbf{u} \]
\[ (\alpha \mathbf{u}) \times \mathbf{v} = \alpha (\mathbf{u} \times \mathbf{v}) = \mathbf{u} \times (\alpha \mathbf{v}) \]
\[ (\mathbf{u}_1 + \mathbf{u}_2) \times \mathbf{v} = \mathbf{u}_1 \times \mathbf{v} + \mathbf{u}_2 \times \mathbf{v} \]

Examples are
\[ \mathbf{i} \times \mathbf{j} = \mathbf{k} \quad \mathbf{j} \times \mathbf{k} = \mathbf{i} \quad \mathbf{k} \times \mathbf{i} = \mathbf{j} \] (216)

**triple product:** The triple product of vectors \( \mathbf{w}, \mathbf{u}, \mathbf{v} \) is defined by
\[ \mathbf{w} \cdot (\mathbf{u} \times \mathbf{v}) = w_1(u_2v_3 - u_3v_2) + w_2(u_3v_1 - u_1v_3) + w_3(u_1v_2 - u_2v_1) \]
\[ = \det \begin{pmatrix} w_1 & w_2 & w_3 \\ u_1 & u_2 & u_3 \\ v_1 & v_2 & v_3 \end{pmatrix} \] (217)

The absolute value \( |\mathbf{w} \cdot (\mathbf{u} \times \mathbf{v})| \) is the volume of the parallelipiped spanned by \( \mathbf{u}, \mathbf{v}, \mathbf{w} \) (see figure 16). This is so because if \( \phi \) is the angle between \( \mathbf{w} \) and \( \mathbf{u} \times \mathbf{v} \) then
\[ \text{volume} = (\text{area of base}) \times \text{height} \]
\[ = (|\mathbf{u} \times \mathbf{v}|)(|\mathbf{w}| \cos \phi) \]
\[ = |\mathbf{w} \cdot (\mathbf{u} \times \mathbf{v})| \] (218)

**problem:** Find the volume of the parallelipiped spanned by \( \mathbf{i} + \mathbf{j}, \mathbf{j}, \mathbf{i} + \mathbf{j} + \mathbf{k} \).

**solution:**
\[ \det \begin{pmatrix} 1 & 1 & 0 \\ 0 & 1 & 0 \\ 1 & 1 & 1 \end{pmatrix} = 1 \] (219)
Planes: A plane is determined by a particular point \( \mathbf{R}_0 = x_0\mathbf{i} + y_0\mathbf{j} + z_0\mathbf{k} \) in the plane and a vector \( \mathbf{N} = N_1\mathbf{i} + N_2\mathbf{j} + N_3\mathbf{k} \) perpendicular to the plane, called a normal vector. If \( \mathbf{R} = xi + yj + zk \) is any other point in the plane, then \( \mathbf{R} - \mathbf{R}_0 \) lies in the plane, hence it is orthogonal to \( \mathbf{N} \) and hence (see figure 8)

\[
\mathbf{N} \cdot (\mathbf{R} - \mathbf{R}_0) = 0 \tag{220}
\]

In fact this is the equation of the plane. That is a point \( \mathbf{R} \) lies on the plane if and only if it satisfies the equation. Written out it says

\[
N_1(x - x_0) + N_2(y - y_0) + N_3(z - z_0) = 0 \tag{221}
\]

Problem: Find the plane determined by the three points \( \mathbf{a} = \mathbf{i}, \mathbf{b} = 2\mathbf{j}, \mathbf{c} = 3\mathbf{k} \).

Solution: \( \mathbf{b} - \mathbf{a} = -\mathbf{i} + 2\mathbf{j} \) and \( \mathbf{c} - \mathbf{a} = -\mathbf{i} + 3\mathbf{k} \) both lie in the plane. Their cross product is orthogonal to both, hence to the plane, and can be take as a normal vector:

\[
\mathbf{N} = (\mathbf{b} - \mathbf{a}) \times (\mathbf{c} - \mathbf{a}) = \det \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ -1 & 2 & 0 \\ -1 & 0 & 3 \end{pmatrix} = 6\mathbf{i} + 3\mathbf{j} + 2\mathbf{k} \tag{222}
\]

For the particular point take \( \mathbf{R}_0 = \mathbf{a} = \mathbf{i} \). Then the equation of the plane is

\[
\mathbf{N} \cdot (\mathbf{R} - \mathbf{R}_0) = 6(x - 1) + 3y + 2z = 0 \tag{223}
\]
2.2 vector-valued functions

A vector-valued function is a function from $\mathbb{R}$ to $\mathbb{R}^3$ (or more generally $\mathbb{R}^n$). It is written

$$\mathbf{R}(t) = (x(t), y(t), z(t)) = x(t)\mathbf{i} + y(t)\mathbf{j} + z(t)\mathbf{k}$$  \hspace{1cm} (224)

To say $\mathbf{R}(t)$ has limit

$$\lim_{t \to t_0} \mathbf{R}(t) = \mathbf{R}_0$$  \hspace{1cm} (225)

means that

$$\lim_{t \to t_0} |\mathbf{R}(t) - \mathbf{R}_0| = 0$$  \hspace{1cm} (226)

If $\mathbf{R}_0 = x_0\mathbf{i} + y_0\mathbf{j} + z_0\mathbf{j}$ then

$$|\mathbf{R}(t) - \mathbf{R}_0| = \sqrt{(x(t) - x_0)^2 + (y(t) - y_0)^2 + (z(t) - z_0)^2}$$  \hspace{1cm} (227)

Hence $\lim_{t \to t_0} \mathbf{R}(t) = \mathbf{R}_0$ is the same as the three limits

$$\lim_{t \to t_0} x(t) = x_0 \quad \lim_{t \to t_0} y(t) = y_0 \quad \lim_{t \to t_0} z(t) = z_0$$  \hspace{1cm} (228)

The function $\mathbf{R}(t)$ is continuous if

$$\lim_{t \to t_0} \mathbf{R}(t) = \mathbf{R}_0$$  \hspace{1cm} (229)
This is the same as saying that the components \( x(t), y(t), z(t) \) are all continuous. The function \( \mathbf{R}(t) \) is differentiable if

\[
\mathbf{R}'(t) = \frac{d\mathbf{R}}{dt} = \lim_{h \to 0} \frac{\mathbf{R}(t+h) - \mathbf{R}(t)}{h}
\]

(230)

exists. Since

\[
\frac{\mathbf{R}(t+h) - \mathbf{R}(t)}{h} = \frac{x(t+h) - x(t)}{h} \mathbf{i} + \frac{y(t+h) - y(t)}{h} \mathbf{j} + \frac{z(t+h) - z(t)}{h} \mathbf{k}
\]

(231)

this is the same as saying that the components \( x(t), y(t), z(t) \) are all differentiable, in which case the derivative is

\[
\mathbf{R}'(t) = x'(t) \mathbf{i} + y'(t) \mathbf{j} + z'(t) \mathbf{k}
\]

(232)
In other words you find the derivative by differentiating the components

The range of a continuous function \( R(t) \) is a curve. The derivative \( R'(t) \) has the interpretation of being a tangent to the curve as figure 9 shows.

A common application is that \( R(t) \) is the location of some object at time \( t \). Then \( R'(t) \) is the velocity and the magnitude of the velocity \( |R'(t)| \) is the speed. The second derivative \( R''(t) \) is the acceleration.

Here are some examples illustrated in figure 10

**example:** straight line. Consider

\[
R(t) = at + b
\]  

(233)
Then \( \mathbf{R}(0) = \mathbf{b} \) and \( \mathbf{R}'(t) = \mathbf{a} \) so it is a straight line through \( \mathbf{b} \) in the direction \( \mathbf{a} \).

**example**: circle. Let \( \mathbf{R}(t) \) be the point on a circle of radius \( a \) with polar angle \( t \). As \( t \) increases it travels around the circle at uniform speed. The point has Cartesian coordinates \( x = a \cos t, y = a \sin t \) so

\[
\mathbf{R}(t) = a \cos t \, \mathbf{i} + a \sin t \, \mathbf{j}
\]  
(234)

The velocity is

\[
\mathbf{R}'(t) = -a \sin t \, \mathbf{i} + a \cos t \, \mathbf{j}
\]  
(235)

and the speed is \( |\mathbf{R}'| = a \).

**example**: helix. To the previous example we add a constant velocity in the \( z \)-direction

\[
\mathbf{R}(t) = a \cos t \, \mathbf{i} + a \sin t \, \mathbf{j} + bt \mathbf{k}
\]  
(236)

This describes a helix and we have

\[
\mathbf{R}'(t) = -a \sin t \, \mathbf{i} + a \cos t \, \mathbf{j} + b \mathbf{k}
\]  
(237)

### 2.3 other coordinate systems

We next describe vector valued functions using other coordinate systems.

**A. Polar**: First some general remarks about vectors and polar coordinates in \( \mathbb{R}^2 \). Let \( \mathbf{R}(r, \theta) \) be the point with polar coordinates \( r, \theta \). This has Cartesian coordinates \( x = r \cos \theta, y = r \sin \theta \) so

\[
\mathbf{R}(r, \theta) = r \cos \theta \, \mathbf{i} + r \sin \theta \, \mathbf{j}
\]  
(238)

If we vary \( r \) with \( \theta \) fixed we get an "r-line". The tangent vector to this line is

\[
\frac{\partial \mathbf{R}}{\partial r}(r, \theta) = \cos \theta \, \mathbf{i} + \sin \theta \, \mathbf{j}
\]  
(239)

If we vary \( \theta \) with \( r \) fixed we get an "\( \theta \)-line". The tangent vector to this line is

\[
\frac{\partial \mathbf{R}}{\partial \theta}(r, \theta) = -r \sin \theta \, \mathbf{i} + r \cos \theta \, \mathbf{j}
\]  
(240)

We also consider unit tangent vectors to these coordinate lines:

\[
\mathbf{e}_r(\theta) = \frac{\partial \mathbf{R}/\partial r}{|\partial \mathbf{R}/\partial r|} = \cos \theta \, \mathbf{i} + \sin \theta \, \mathbf{j}
\]
\[
\mathbf{e}_\theta(\theta) = \frac{\partial \mathbf{R}/\partial \theta}{|\partial \mathbf{R}/\partial \theta|} = -\sin \theta \, \mathbf{i} + \cos \theta \, \mathbf{j}
\]  
(241)

See figure 11.
Note that
\[ e_r(\theta) \cdot e_\theta(\theta) = \cos \theta (-\sin \theta) + \sin \theta \cos \theta = 0 \] (242)
Thus for any \( \theta \) the vectors \( e_r(\theta), e_\theta(\theta) \) form an orthonormal set of vectors in \( \mathbb{R}^2 \) and hence an orthonormal basis. Also note for future reference that
\[
\frac{de_r}{d\theta} = e_\theta \\
\frac{de_\theta}{d\theta} = -e_r \] (243)

Now a curve is specified in polar coordinates by a pair of functions \( r(t), \theta(t) \). The Cartesian coordinates are \( x(t) = r(t) \cos \theta(t) \) and \( y(t) = r(t) \sin \theta(t) \). So the curve is
described in polar coordinates with polar basis vectors by
\[ \mathbf{R}(t) = x(t)\mathbf{i} + y(t)\mathbf{j} \]
\[ = r(t)(\cos \theta(t)\mathbf{i} + \sin \theta(t)\mathbf{j}) \]
\[ = r(t)e_r(\theta(t)) \]  
(244)

The velocity is
\[ \mathbf{R}'(t) = r'(t)e_r(\theta(t)) + r(t)\frac{d}{dt}e_r(\theta(t)) \]  
(245)
But
\[ \frac{d}{dt}e_r(\theta(t)) = \frac{de_r}{d\theta}(\theta(t))\frac{d\theta}{dt} = \theta'(t)e_\theta(\theta(t)) \]  
(246)
and so
\[ \mathbf{R}'(t) = r'(t)e_r(\theta(t)) + r(t)\theta'(t)e_\theta(\theta(t)) \]  
(247)
By differentiating this we get a formula for the acceleration \( \mathbf{R}''(t) \):
\[ \mathbf{R}''(t) = \left( r''(t) - r'(t)(\theta'(t)^2)e_r(\theta(t)) + \left( r(t)\theta''(t) + 2r'(t)\theta'(t) \right)e_\theta(\theta(t)) \right) \]  
(248)
We summarize in an abbreviated notation
\[ \mathbf{R} = r e_r \]
\[ \mathbf{R}' = r' e_r + r\theta' e_\theta \]  
(249)
\[ \mathbf{R}'' = (r'' - r(\theta')^2)e_r + (r\theta'' + 2r'\theta')e_\theta \]

**Example:** Consider the spiral described in polar coordinates by \( r = at \) and \( \theta = bt \). Then \( r' = a, \theta' = b \) and \( r'' = 0, \theta'' = 0 \) and so
\[ \mathbf{R} = at e_r \]
\[ \mathbf{R}' = a e_r + abt e_\theta \]  
(250)
\[ \mathbf{R}'' = -ab^2t e_r + 2ab e_\theta \]
In these formulas \( e_r = e_r(bt) = \cos(bt)\mathbf{i} + \sin(bt)\mathbf{j} \) and \( e_\theta = e_\theta(bt) = -\sin(bt)\mathbf{i} + \cos(bt)\mathbf{j} \).

**B. Cylindrical:** Cylindrical coordinates in \( \mathbb{R}^3 \) replace \( x, y \) by polar coordinates \( r, \theta \) and leave \( z \) alone. Thus
\[ x = r \cos \theta \]
\[ y = r \sin \theta \]  
(251)
\[ z = z \]
A point with cylindrical coordinates \( r, \theta, z \) is
\[ \mathbf{R}(r, \theta, z) = r \cos \theta \mathbf{i} + r \sin \theta \mathbf{j} + z \mathbf{k} \]  
(252)
Unit tangent vectors to the coordinate lines are $e_r, e_\theta$ as before and $e_z = k$. See figure 11.

A curve in cylindrical coordinate is given by $r(t), \theta(t), z(t)$ and we have

$$\mathbf{R}(t) = r(t)e_r(\theta(t)) + z(t)e_z$$

(253)

As before:

$$\mathbf{R} = re_r + ze_z$$
$$\mathbf{R}' = r'e_r + r\theta'e_\theta + z'e_z$$
$$\mathbf{R}'' = (r'' - r(\theta')^2)e_r + (r\theta'' + 2r'\theta')e_\theta + z''e_z$$

(254)

C. **spherical:** Spherical coordinates $\rho, \phi, \theta$ label a point by its distance to the origin, the angle with the $z$-axis, and the polar angle when projected into the $x, y$ plane. The corresponding Cartesian coordinates are

$$x = \rho \sin \phi \cos \theta$$
$$y = \rho \sin \phi \sin \theta$$
$$z = \rho \cos \phi$$

(255)

The point with spherical coordinates $\rho, \phi, \theta$ is

$$\mathbf{R}(\rho, \phi, \theta) = \rho \sin \phi \cos \theta \mathbf{i} + \rho \sin \phi \sin \theta \mathbf{j} + \rho \cos \phi \mathbf{k}$$

(256)

Tangent vectors to the coordinate lines are

$$\frac{\partial \mathbf{R}}{\partial \rho} = \sin \phi \cos \theta \mathbf{i} + \sin \phi \sin \theta \mathbf{j} + \cos \phi \mathbf{k}$$
$$\frac{\partial \mathbf{R}}{\partial \phi} = \rho \cos \phi \cos \theta \mathbf{i} + \rho \cos \phi \sin \theta \mathbf{j} - \rho \sin \phi \mathbf{k}$$
$$\frac{\partial \mathbf{R}}{\partial \theta} = -\rho \sin \phi \sin \theta \mathbf{i} + \rho \sin \phi \cos \theta \mathbf{j}$$

(257)

Divide by the length and get unit tangent vectors to the coordinate lines: (see figure 12)

$$\mathbf{e}_\rho(\phi, \theta) = \sin \phi \cos \theta \mathbf{i} + \sin \phi \sin \theta \mathbf{j} + \cos \phi \mathbf{k}$$
$$\mathbf{e}_\phi(\phi, \theta) = \cos \phi \cos \theta \mathbf{i} + \cos \phi \sin \theta \mathbf{j} - \sin \phi \mathbf{k}$$
$$\mathbf{e}_\theta(\phi, \theta) = -\sin \theta \mathbf{i} + \cos \theta \mathbf{j}$$

(258)

A curve is spherical coordinates is specified by three functions $r(t), \phi(t), \theta(t)$. The corresponding vector-valued function is

$$\mathbf{R}(t) = x(t)\mathbf{i} + y(t)\mathbf{j} + z(t)\mathbf{k}$$
$$= \rho(t) \left( \sin \phi(t) \cos \theta(t)\mathbf{i} + \sin \phi(t) \sin \theta(t)\mathbf{j} + \cos \phi(t)\mathbf{k} \right)$$
$$= \rho(t) \mathbf{e}_\rho(\phi(t), \theta(t))$$

(259)
Now we can take derivatives and we find

\[
\begin{align*}
\mathbf{R} &= \rho \mathbf{e}_\rho \\
\mathbf{R}' &= \rho' \mathbf{e}_\rho + \rho \phi' \mathbf{e}_\phi + \rho \theta' \sin \phi \mathbf{e}_\theta \\
\mathbf{R}'' &= \left( \rho'' - \rho (\phi')^2 - \rho (\theta')^2 \sin^2 \phi \right) \mathbf{e}_\rho \\
&\quad + \left( \rho \phi'' + 2 \rho' \phi' - \rho (\theta')^2 \sin \phi \cos \phi \right) \mathbf{e}_\phi \\
&\quad + \left( \rho \theta'' \sin \phi + 2 \rho' \phi' \sin \phi + 2 \rho \theta' \phi' \cos \phi \right) \mathbf{e}_\theta
\end{align*}
\]  

(260)

**example:** Suppose that \( \rho = 1, \phi = at, \theta = bt \) with \( b \) much greater than \( a \). This represents a point on a sphere sprialing down from the north pole. Then \( \rho' = 0, \phi' = \)
\[ a, \theta' = b \text{ and } \rho'' = 0, \phi'' = 0, \theta'' = 0 \text{ and we have with } e_\rho = e_\rho(at, bt), \text{ etc} \]

\[
R = e_\rho,
R' = a e_\phi + b \sin at \, e_\theta
\]

\[
R'' = \left( -a^2 - b^2 \sin^2 at \right) e_\rho + \left( -b^2 \sin at \cos at \right) e_\phi + \left( 2ab \cos at \right) e_\theta
\]

### 2.4 Line Integrals

We want to define the length of a curve \( C \) in \( \mathbb{R}^3 \). Suppose the curve is the range of a vector valued function \( R(t) = x(t)i + y(t)j + z(t)k, \ a \leq t \leq b \). We say that \( R(t) \) is a parametrization of \( C \). There will be many parametrizations, but we pick one. We divide up the interval \([a, b]\) by picking points \( a = t_0 < t_1 < t_2 < \cdots < t_n = b \) (262)

This gives a sequence of points on the curve \( R(t_0), R(t_1), \ldots R(t_n) \). (see figure 13). If \( \Delta t_i = t_{i+1} - t_i \) is small then for any \( t^*_i \) in the interval \([t_i, t_{i+1}]\)

\[
R(t_{i+1}) - R(t_i) = (x(t_{i+1}) - x(t_i))i + (y(t_{i+1}) - y(t_i))j + (z(t_{i+1}) - z(t_i))k
\approx x'(t^*_i)\Delta t_i i + y'(t^*_i)\Delta t_i j + z'(t^*_i)\Delta t_i k
= R'(t^*_i)\Delta t_i
\]

(The mean value theorem says there is a point \( t^*_i \) so that \( (x(t_{i+1}) - x(t_i)) = x'(t^*_i)\Delta t_i \). Changing to an arbitrary point in the interval is second order small and negligible).

Let \( \Delta s_i \) be the length of the straight line from \( R(t_i) \) to \( R(t_{i+1}) \). Then

\[
\Delta s_i = |R(t_{i+1}) - R(t_i)| \approx |R'(t^*_i)|\Delta t_i
\]

(264)

Then we have

\[
\text{length of } C \approx \sum_{i=0}^{n-1} \Delta s_i \approx \sum_{i=0}^{n-1} |R'(t^*_i)|\Delta t_i
\]

(265)

This is a Riemann sum and as the division becomes increasingly fine, i.e as \( \max_i \Delta t_i \) tends to 0, this converges a Riemann integral which we take as the definition

\[
\text{length of } C = \int_a^b |R'(t)|dt
\]

(266)

One can show that this depends only on \( C \) and not on the particular parametrization.
More generally we want to define the integral of a function $f(\mathbf{R}) = f(x, y, z)$ over the curve $C$. An approximation to what we want is

$$\sum_{i=0}^{n-1} f(\mathbf{R}(t_i^*)) \Delta s_i \approx \sum_{i=0}^{n-1} f(\mathbf{R}(t_i^*)) |\mathbf{R}'(t_i^*)| \Delta t_i \quad (267)$$

As the division becomes fine this converges to a Riemann integral which we take as the definition of the integral of $f$ over $C$. It is denoted $\int_C f(\mathbf{R}) ds$ and is given by

$$\int_C f(\mathbf{R}) ds = \int_a^b f(\mathbf{R}(t)) |\mathbf{R}'(t)| dt \quad (268)$$

This is also independent of parametrization. A short way to remember it is to replace $\mathbf{R}$ by its parametrization $\mathbf{R}(t)$, replace $C$ by the interval $[a, b]$ and replace the formal
symbol $ds$ by

$$ds = \left| \frac{d\mathbf{R}}{dt} \right| dt$$  \hspace{1cm} (269)$$

Note also that if $f(\mathbf{R}) = 1$ we have

$$\int_{C} ds = \text{length of } C$$  \hspace{1cm} (270)$$

The same formulas hold in $\mathbb{R}^2$ but now $\mathbf{R}(t) = x(t)i + y(t)j$.

**example:** Consider the helix parametrized by

$$\mathbf{R}(t) = a \cos t \, i + a \sin t \, j + bt \, k$$  \hspace{1cm} (271)$$

with $0 \leq t \leq 2\pi$. Then

$$\frac{d\mathbf{R}}{dt} = -a \sin t \, i + a \cos t \, j + b \, k$$  \hspace{1cm} (272)$$

and

$$ds = \left| \frac{d\mathbf{R}}{dt} \right| dt = \sqrt{a^2 + b^2} \, dt$$  \hspace{1cm} (273)$$

The length is then

$$\int_{C} ds = \int_{0}^{2\pi} \sqrt{a^2 + b^2} \, dt = 2\pi \sqrt{a^2 + b^2}$$  \hspace{1cm} (274)$$

**example:** Suppose we have a thin semi-circular wire of radius $a$ with uniform linear density $\rho$ (mass per unit of length). We want to find the $y$-component of the center of mass. This is defined by dividing the wire up into segments of length $\Delta s_i$ and mass $\Delta m_i = \rho \Delta s_i$ and computing

$$\bar{y} = \frac{\sum y_i \Delta m_i}{\sum \Delta m_i} = \frac{\sum y_i \Delta s_i}{\sum \Delta s_i}$$  \hspace{1cm} (275)$$

where $y_i$ is the $y$-coordinate of the $i^{th}$ segment. As the division becomes fine this is expressed as a ratio of line integrals over the semi-circle $C$

$$\bar{y} = \frac{\int_{C} y ds}{\int_{C} ds}$$  \hspace{1cm} (276)$$

To compute it parametrize the semi-circle by

$$\mathbf{R}(\theta) = a \cos \theta i + a \sin \theta j$$  \hspace{1cm} (277)$$
with $0 \leq \theta \leq \pi$. Then
\[
\frac{d\mathbf{R}}{d\theta} = -a \sin \theta \, \mathbf{i} + a \cos \theta \, \mathbf{j}
\]
and
\[
ds = \left| \frac{d\mathbf{R}}{d\theta} \right| \, d\theta = a \, d\theta
\]
Then since $y = a \sin \theta$
\[
\int_C y \, ds = \int_0^\pi a \sin \theta \, a \, d\theta = 2a^2
\]
and
\[
\int_C ds = \int_0^\pi a \, d\theta = a \pi
\]
Thus
\[
\bar{y} = \frac{2a^2}{a \pi} = \frac{2a}{\pi}
\]

2.5 double integrals

Let $\mathcal{R}$ be a region in $\mathbb{R}^2$ and let $f(x, y)$ be a function defined on $\mathcal{R}$. We want to define the integral of $f$ over $\mathcal{R}$ denoted by
\[
\int_\mathcal{R} f(x, y) \, dA \quad \text{or} \quad \int_\mathcal{R} f(x, y) \, dx \, dy \quad \text{or} \quad \iint_\mathcal{R} f(x, y) \, dx \, dy
\]
It is supposed to be the sum of the values of the function weighted by area.

To define it put a rectangular grid over the region (see figure 14) and suppose the rectangles are enumerated by some index $i$. The $i^{th}$ rectangle will have some dimensions $\Delta x_i, \Delta y_i$. Let $\Delta A_i = \Delta x_i \Delta y_i$ be the area of the $i^{th}$ rectangle. Also let $(x_i^*, y_i^*)$ be any point in the $i^{th}$ rectangle. An approximation to what we want is the Riemann sum
\[
\sum_i f(x_i^*, y_i^*) \Delta A_i
\]
If these expressions approach a definite number as the grid becomes fine then this is the integral we want. The fineness of the grid can be measured by
\[
h = \max_i \{\Delta x_i, \Delta y_i\}
\]
Here is an exact definition of the integral.

**definition:** If there is a number $I$ so that for any $\epsilon > 0$ there is a $\delta > 0$ such that for any grid over $\mathcal{R}$ with $h < \delta$ and any choice of points $(x_i^*, y_i^*)$ in the grid we have
\[
| \sum_i f(x_i^*, y_i^*) \Delta A_i - I | < \epsilon
\]
Figure 14: double integral

then $f$ is integrable over $\mathcal{R}$ and we define

$$\int_{\mathcal{R}} f(x, y) dA = I$$  \hfill (287)

For short one can write this as

$$\int_{\mathcal{R}} f(x, y) dA = \lim_{h \to 0} \sum_i f(x_i^*, y_i^*) \Delta A_i$$  \hfill (288)

although it is not an ordinary limit since the right side is not a function of $h$.

One can show:

**Theorem 12** *Continuous functions are integrable.*

Here are some applications of double integrals:

1. With $f = 1$

$$\int_{\mathcal{R}} dA = \text{area of } \mathcal{R}$$  \hfill (289)
2. If $\mathcal{R}$ represents a thin plate and $f(x, y)$ is the density of the plate (mass per unit area) then $f(x_i^*, y_i^*) \Delta A_i$ is the approximate mass of the $i^{th}$ rectangle and so
\[
\int_{\mathcal{R}} f(x, y) dA = \text{total mass of plate} \tag{290}
\]

3. If $f(x, y) \geq 0$ then $f(x_i^*, y_i^*) \Delta A_i$ is the approximate volume of the column above the $i^{th}$ rectangle and under the graph and so
\[
\int_{\mathcal{R}} f(x, y) dA = \text{volume under the graph of } z = f(x, y) \text{ above } \mathcal{R} \tag{291}
\]

Here are some properties of double integrals:

1. For any two functions $f_1, f_2$ on $\mathcal{R}$
\[
\int_{\mathcal{R}} (f_1 + f_2) dA = \int_{\mathcal{R}} f_1 dA + \int_{\mathcal{R}} f_2 dA
\]

2. If $\alpha$ is a constant
\[
\int_{\mathcal{R}} \alpha f dA = \alpha \int_{\mathcal{R}} f dA
\]

3. If $\mathcal{R}$ can be written as a disjoint union $\mathcal{R} = \mathcal{R}_1 \cup \mathcal{R}_2$ then
\[
\int_{\mathcal{R}} f dA = \int_{\mathcal{R}_1} f dA + \int_{\mathcal{R}_2} f dA
\]

To compute double integrals one writes them as iterated integrals in one variable and then uses the fundamental theorem of calculus.

**Theorem 13** Suppose the region $\mathcal{R}$ has the form
\[
\mathcal{R} = \{(x, y) : a \leq x \leq b, p(x) \leq y \leq q(x)\} \tag{292}
\]
for some functions $p, q$. (See figure 15). Then
\[
\int_{\mathcal{R}} f(x, y) dA = \int_{a}^{b} \left( \int_{p(x)}^{q(x)} f(x, y) dy \right) dx \tag{293}
\]

This says fix $x$ and integrate over the $y$ values in the region for this value of $x$. This gives you a function of $x$ which you integrate over the $x$ values for the region.
Figure 15: iterated integral

**Example:** Suppose the region $\mathcal{R}$ below the graph of $y = -x^2 + 1$ in the first quadrant. Thus $\mathcal{R}$ is defined by $0 \leq x \leq 1$ and $0 \leq y \leq -x^2 + 1$. We compute

\[
\int_{\mathcal{R}} x \, dA = \int_0^1 \left( \int_0^{-x^2+1} x \, dy \right) \, dx = \int_0^1 [xy]_{y=0}^{y=-x^2+1} \, dx \\
= \int_0^1 (-x^3 + x) \, dx = \frac{-1}{4} + \frac{1}{2} = \frac{1}{4}
\]

(294)

Alternatively $\mathcal{R}$ can be regarded as the region $0 \leq y \leq 1$ and $0 \leq x \leq \sqrt{1-y}$ (draw a picture). Then we can do the $x$ integral first:

\[
\int_{\mathcal{R}} x \, dA = \int_0^1 \left( \int_0^{\sqrt{1-y}} x \, dx \right) \, dy = \int_0^1 \left[ \frac{x^2}{2} \right]_{x=0}^{x=\sqrt{1-y}} \, dy \\
= \int_0^1 \frac{1-y}{2} \, dy = \frac{1}{2} - \frac{1}{4} = \frac{1}{4}
\]

(295)
2.6 triple integrals

Now \( \mathcal{R} \) be a region in \( \mathbb{R}^3 \) and let \( f \) be a function defined on \( \mathcal{R} \). We want to define the integral of \( f \) over \( \mathcal{R} \) which will be denoted

\[
\int_{\mathcal{R}} f(x, y, z) dV \quad \text{or} \quad \int_{\mathcal{R}} f(x, y, z) dx dy dz \quad (296)
\]

To defined it divide the region \( \mathcal{R} \) up into many small rectangular boxes. Suppose the \( i^{th} \) box has dimensions \( \Delta x_i, \Delta y_i, \Delta z_i \) and volume \( \Delta V_i = \Delta x_i \Delta y_i \Delta z_i \) and let \((x_i^*, y_i^*, z_i^*)\) be an arbitrary point in the \( i^{th} \) box. Also let

\[
h = \max_i (\Delta x_i, \Delta y_i, \Delta z_i) \quad (297)
\]

be the large dimension in the grid. Then we define

\[
\int_{\mathcal{R}} f(x, y, z) dV = \lim_{h \to 0} \sum_i f(x_i^*, y_i^*, z_i^*) \Delta V_i \quad (298)
\]

If \( f = 1 \) then \( \int_{\mathcal{R}} dV \) is interpreted as the volume of \( \mathcal{R} \). Another application is that \( \mathcal{R} \) could represent a solid object. If \( f(x, y, z) \) is the density at the point \((x, y, z)\) (mass per unit volume) then \( \int_{\mathcal{R}} f(x, y, z) dV \) is the total mass of the object.

Suppose the region \( \mathcal{R} \) is the region between the graphs of \( z = \phi(x, y) \) and \( z = \psi(x, y) \) with \((x, y)\) restricted to some plane region \( E \) (see figure 16). Then we can write the triple integral as a single integral followed by a double integral:

\[
\int_{\mathcal{R}} f(x, y, z) dV = \int_E \left( \int_{\phi(x, y)}^{\psi(x, y)} f(x, y, z) dz \right) dA \quad (299)
\]

If in addition the plane region \( E \) is the region between two curves \( y = p(x) \) and \( y = q(x) \) with \( a \leq x \leq b \) then the double integral can be written as an iterated in integral and we have

\[
\int_{\mathcal{R}} f(x, y, z) dV = \int_a^b \left( \int_{p(x)}^{q(x)} \left( \int_{\phi(x, y)}^{\psi(x, y)} f(x, y, z) dz \right) dy \right) dx \quad (300)
\]

d example: Suppose we are given the problem of finding the volume between the paraboloid \( z = 2 - x^2 - y^2 \) and the plane \( z = 1 \).

These surfaces intersect when \( x^2 + y^2 = 1 \). The problem must be refering to the region with \( x^2 + y^2 \leq 1 \) since the region with \( x^2 + y^2 \geq 1 \) is infinite. Thus we want to find the volume of the region \( \mathcal{R} \) below \( z = 2 - x^2 - y^2 \) and above \( z = 1 \) with \( x^2 + y^2 \leq 1 \).
It is

\[ \int_{\mathcal{R}} dV = \int_{x^2+y^2 \leq 1} \left( \int_1^{1-x^2-y^2} dz \right) dA \]
\[ = \int_{x^2+y^2 \leq 1} (1 - x^2 - y^2) dA \]
\[ = \int_{-1}^1 \left( \int_{-\sqrt{1-x^2}}^{\sqrt{1-x^2}} (1 - x^2 - y^2) dy \right) dx \]
\[ = \int_{-1}^1 \frac{4}{3} (1 - x^2)^{3/2} dx \]
\[ = \frac{\pi}{2} \] (301)

The last integral is left as an exercise. (An alternative is to evaluate the integral \( \int_{x^2+y^2 \leq 1} (1 - x^2 - y^2) dA \) in polar coordinates, a topic we take up later.)

**example:** Let \( \mathcal{R} \) be the region bounded by the planes \( x = 0, y = 0, z = 0 \) and \( x + y + z = 1 \) and suppose we want to write \( \int_{\mathcal{R}} x \ dV \) as an iterated integral.

The intersection of \( \mathcal{R} \) with the plane \( z = 0 \) is the region \( E \) bounded by the lines \( x = 0, y = 0, x + y = 1 \). In fact the region \( \mathcal{R} \) lies between \( z = 1 - x - y \) and \( z = 0 \) and
above \( E \). (draw a picture). Thus we have

\[
\int_{\mathcal{R}} x \, dV = \int_{E} \left( \int_{0}^{1-x-y} x \, dz \right) \tag{302}
\]

But since \( E \) lies between \( y = 0 \) and \( y = 1 - x \) with \( 0 \leq x \leq 1 \) this can be expressed as

\[
\int_{\mathcal{R}} x \, dV = \int_{0}^{1} \left( \int_{0}^{1-x} \left( \int_{0}^{1-x-y} x \, dz \right) \, dy \right) \, dx \tag{303}
\]

The evaluation is left as an exercise.

### 2.7 parametrized surfaces

Consider a function from \( \mathcal{R} \subset \mathbb{R}^2 \) to \( \mathbb{R}^3 \) which we write as

\[
x = x(u, v) \quad y = y(u, v) \quad z = z(u, v) \tag{304}
\]

The range of this function is a surface \( S \) and the function is called a parametrization of the surface. (A surface has many possible parametrizations, but we pick one). The function can also be written

\[
\mathbf{R}(u, v) = x(u, v)\mathbf{i} + y(u, v)\mathbf{j} + z(u, v)\mathbf{k} \tag{305}
\]

**example:** Consider the function

\[
x = a \sin \phi \cos \theta \\
y = a \sin \phi \sin \theta \\
z = a \cos \phi \tag{306}
\]

with \( 0 < \phi < \pi \) and \( 0 < \theta < 2\pi \). Then \( S \) is the surface of a sphere of radius \( a \), and it is parametrized by spherical coordinates.

**example:** Suppose \( S \) is the graph of a function \( z = \phi(x, y) \) with \( (x, y) \in \mathcal{R} \). Then \( S \) can be parametrized by

\[
x = u \quad y = v \quad z = \phi(u, v) \tag{307}
\]

with \( (u, v) \in \mathcal{R} \). This can also be written

\[
x = x \quad y = y \quad z = \phi(x, y) \tag{308}
\]

with \( (x, y) \in \mathcal{R} \).
Now suppose $S$ is a surface parametrized by $R(u,v)$. At $(u_0,v_0)$ we fix $v$ and vary $u$ we get a $u$-line in $S$. Then

$$\frac{\partial R}{\partial u}(u_0,v_0) = \text{tangent vector to } u\text{-line through } R(u_0,v_0)$$

$$\frac{\partial R}{\partial v}(u_0,v_0) = \text{tangent vector to } v\text{-line through } R(u_0,v_0)$$

Together these two tangent vectors determine the tangent plane to the surface at $R(u_0,v_0)$. A normal to this tangent plane is (see figure 17)

$$N = \frac{\partial R}{\partial u}(u_0,v_0) \times \frac{\partial R}{\partial v}(u_0,v_0)$$

(309)

With this $N$ the equation of the tangent plane to the surface $S$ at $R_0 = R(u_0,v_0)$ is

$$N \cdot (R - R_0) = 0$$

(310)

**example:** Suppose we want to find the tangent plane to the surface

$$x = u + v \quad y = u - v \quad z = 2uv$$

(311)
at the point \( u = 1, v = 1 \). In this case
\[
\mathbf{R}(u, v) = (u + v)\mathbf{i} + (u - v)\mathbf{j} + 2uv\mathbf{k} \tag{312}
\]
and the point is
\[
\mathbf{R}_0 = \mathbf{R}(1, 1) = 2\mathbf{i} + 2\mathbf{k} \tag{313}
\]
At this point
\[
\frac{\partial \mathbf{R}}{\partial u} = \mathbf{i} + \mathbf{j} + 2v\mathbf{k} = \mathbf{i} + \mathbf{j} + 2\mathbf{k} \tag{314}
\]
\[
\frac{\partial \mathbf{R}}{\partial v} = \mathbf{i} - \mathbf{j} + 2u\mathbf{k} = \mathbf{i} - \mathbf{j} + 2\mathbf{k}
\]
so the normal is
\[
\mathbf{N} = \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} = \mathbf{det} \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ 1 & 1 & 2 \\ 1 & -1 & 2 \end{pmatrix} = 4\mathbf{i} - 2\mathbf{k} \tag{315}
\]
The equation of the tangent plane at this point is \( \mathbf{N} \cdot (\mathbf{R} - \mathbf{R}_0) = 0 \). Since \( \mathbf{R} - \mathbf{R}_0 = (x - 2)\mathbf{i} + y\mathbf{j} + (z - 2)\mathbf{k} \) this says
\[
4(x - 2) - 2(z - 2) = 0 \tag{316}
\]
which can also be written \( z = 2x - 2 \).

**example:** Suppose our surface is the graph of a function \( z = \phi(x, y) \). Then it can be parametrized by
\[
\mathbf{R}(x, y) = x\mathbf{i} + y\mathbf{j} + \phi(x, y)\mathbf{k} \tag{317}
\]
We want to find the normal and the tangent plane to the graph at \((x_0, y_0)\) This is the point
\[
\mathbf{R}_0 = \mathbf{R}(x_0, y_0) = x_0\mathbf{i} + y_0\mathbf{j} + \phi(x_0, y_0)\mathbf{k} \tag{318}
\]
The derivatives at \((x_0, y_0)\) are
\[
\frac{\partial \mathbf{R}}{\partial x} = \mathbf{i} + \mathbf{j} + \phi_x(x_0, y_0)\mathbf{k} \tag{319}
\]
\[
\frac{\partial \mathbf{R}}{\partial y} = \mathbf{j} + \mathbf{k} + \phi_y(x_0, y_0)\mathbf{k}
\]
and so the normal is
\[
\mathbf{N} = \frac{\partial \mathbf{R}}{\partial x} \times \frac{\partial \mathbf{R}}{\partial y} = \mathbf{det} \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ 1 & 0 & \phi_x(x_0, y_0) \\ 0 & 1 & \phi_y(x_0, y_0) \end{pmatrix} \tag{320}
\]
which says
\[ \mathbf{N} = -\phi_x(x_0, y_0)\mathbf{i} - \phi_y(x_0, y_0)\mathbf{j} + \mathbf{k} \] (321)

The equation of the tangent plane \( \mathbf{N} \cdot (\mathbf{R} - \mathbf{R}_0) = 0 \) is then
\[ -\phi_x(x_0, y_0)(x - x_0) - \phi_y(x_0, y_0)(y - y_0) + (z - \phi(x_0, y_0)) = 0 \] (322)

This can also be written
\[ z = \phi(x_0, y_0) + \phi_x(x_0, y_0)(x - x_0) + \phi_y(x_0, y_0)(y - y_0) \] (323)

which agrees with our earlier definition of the tangent plane.

2.8 surface area

Let \( \mathcal{S} \) be a surface parametrized by a function \( \mathbf{R}(u, v) \) with \((u, v) \in \mathcal{R}\). We assume the function is one-to-one so it only covers \( \mathcal{S} \) once. We want to define the area of \( \mathcal{S} \).

To do so we divide up \( \mathcal{R} \) into a fine rectangular grid (see figure 18). The lines of the grid are mapped to lines in the surface and this divides up the surface into little pieces (no longer rectangles). Suppose the \( i^{th} \) rectangle has lower left corner \((u_i, v_i)\) and dimensions \( \Delta u_i, \Delta v_i \). The image of this rectangle is a patch with corners
\[ \mathbf{R}(u_i, v_i), \quad \mathbf{R}(u_i + \Delta u_i, v_i), \quad \mathbf{R}(u_i, v_i + \Delta v_i), \quad \mathbf{R}(u_i + \Delta u_i, v_i + \Delta v_i) \]

The area of this patch is approximated as the area of the parallelogram spanned by
\[ \mathbf{a}_i = \mathbf{R}(u_i + \Delta u_i, v_i) - \mathbf{R}(u_i, v_i) \]
\[ \mathbf{b}_i = \mathbf{R}(u_i, v_i + \Delta v_i) - \mathbf{R}(u_i, v_i) \] (324)

This area is
\[ \Delta \sigma_i = |\mathbf{a}_i \times \mathbf{b}_i| \] (325)

However since \( \Delta u_i \) and \( \Delta v_i \) are assumed small we have the approximations
\[ \mathbf{a}_i \approx \frac{\partial \mathbf{R}}{\partial u}(u_i, v_i)\Delta u_i \]
\[ \mathbf{b}_i \approx \frac{\partial \mathbf{R}}{\partial v}(u_i, v_i)\Delta v_i \] (326)

Hence
\[ \Delta \sigma_i \approx \left| \frac{\partial \mathbf{R}}{\partial u}(u_i, v_i) \times \frac{\partial \mathbf{R}}{\partial v}(u_i, v_i) \right| \Delta u_i \Delta v_i \] (327)

Now if \( h = \max\{\Delta u_i, \Delta v_i\} \) is the maximum dimension in the grid we define
\[
\text{Area of } \mathcal{S} = \lim_{h \to 0} \sum_i \Delta \sigma_i \\
= \lim_{h \to 0} \sum_i \left| \frac{\partial \mathbf{R}}{\partial u}(u_i, v_i) \times \frac{\partial \mathbf{R}}{\partial v}(u_i, v_i) \right| \Delta u_i \Delta v_i \\
= \int_{\mathcal{R}} \left| \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} \right| \, du \, dv
\]
This definition turns out to be independent of the particular parametrization we have chosen.

**Example:** Find the area of a spherical cap of radius $a$ and angle $\alpha$. In spherical coordinates this is the surface described

$$r = a \quad 0 \leq \phi \leq \alpha \quad 0 \leq \theta \leq 2\pi$$

We parametrize with spherical coordinates and take

$$\mathbf{R}(\phi, \theta) = a \sin \phi \cos \theta \mathbf{i} + a \sin \phi \sin \theta \mathbf{j} + a \cos \phi \mathbf{k} \quad (329)$$
with \(0 \leq \phi \leq \alpha, 0 \leq \theta \leq 2\pi\). Then we compute

\[
\frac{\partial \mathbf{R}}{\partial \phi} \times \frac{\partial \mathbf{R}}{\partial \theta} = \text{det} \begin{pmatrix}
i & j & k \\
a \cos \phi \cos \theta & a \cos \phi \sin \theta & -a \sin \phi \\
-a \sin \phi \sin \theta & a \sin \phi \cos \theta & 0
\end{pmatrix}
\]

\[= a^2 \sin^2 \phi \cos \theta \mathbf{i} + a^2 \sin^2 \phi \sin \theta \mathbf{j} + a^2 \cos \phi \sin \phi \mathbf{k}
\]  

Equation (330)

Then

\[
\left| \frac{\partial \mathbf{R}}{\partial \phi} \times \frac{\partial \mathbf{R}}{\partial \theta} \right| = \sqrt{a^4 \sin^4 \phi (\cos^2 \theta + \sin^2 \theta) + a^4 \cos^2 \phi \sin^2 \phi}
\]

\[= a^2 \sin \phi \sqrt{\sin^2 \phi + \cos^2 \phi}
\]

\[= a^2 \sin \phi
\]  

Equation (331)

and the area is

\[
\text{Area} = \int_0^{2\pi} \int_0^\alpha \left| \frac{\partial \mathbf{R}}{\partial \phi} \times \frac{\partial \mathbf{R}}{\partial \theta} \right| d\phi \, d\theta
\]

\[= \int_0^{2\pi} \int_0^\alpha a^2 \sin \phi \, d\phi \, d\theta
\]

\[= a^2 \left( \int_0^{2\pi} d\theta \right) \left( \int_0^\alpha \sin \phi \, d\phi \right)
\]

\[= 2\pi a^2 (1 - \cos \alpha)
\]  

Equation (332)

Note that if \(\alpha = \pi\) the area is \(4\pi a^2\) which is what we expect for the whole sphere.

### 2.9 Surface Integrals

We continue to consider a surface \(S\) parametrized by a function \(\mathbf{R}(u, v)\) with \((u, v) \in \mathcal{R}\). Also let \(f(\mathbf{R}) = f(x, y, z)\) be a function defined on \(S\) (and possibly elsewhere in \(\mathbb{R}^3\)). We want to define the integral of \(f\) over \(S\) which will be written \(\int_S f(\mathbf{R})d\sigma\).

To define it we again divide up the parameter space into a rectangular grid. We also let \((u_i, v_i)\) be the corner point in the \(i^{th}\) rectangle. (see figure 18 again). Then we define

\[
\int_S f(\mathbf{R})d\sigma = \lim_{h \to 0} \sum_i f(\mathbf{R}(u_i, v_i)) \Delta \sigma_i
\]

\[= \lim_{h \to 0} \sum_i f(\mathbf{R}(u_i, v_i)) \left| \frac{\partial \mathbf{R}}{\partial u}(u_i, v_i) \times \frac{\partial \mathbf{R}}{\partial v}(u_i, v_i) \right| \Delta u_i \Delta v_i
\]  

Equation (333)

\[= \int_{\mathcal{R}} f(\mathbf{R}(u, v)) \left| \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} \right| du \, dv
\]
For short one just has to remember

\[ d\sigma = \left| \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} \right| du \, dv \]  

(334)

A special case is with \( f(\mathbf{R}) = 1 \) which gives

\[ \int_S d\sigma = \int_\mathcal{R} \left| \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} \right| du \, dv = \text{Area of } \mathcal{S} \]  

(335)

**example:** Let \( \mathcal{S} \) represent a thin hemispherical shell of uniform density which has radius \( a \) and is centered on the origin. We want to find the \( z \)-component of the center of mass. Since it is a thin shell it is reasonable to represent in terms of surface integrals and we take the definition

\[ \bar{z} = \frac{\int_S z \, d\sigma}{\int_S d\sigma} \]  

(336)

The hemisphere is parametrized as before by \( x = a \sin \phi \cos \theta, \ y = a \sin \phi \sin \theta, \) and \( z = \cos \phi \) with \( 0 \leq \phi \leq \pi/2 \) and \( 0 \leq \theta \leq 2\pi \). We also have as before

\[ d\sigma = \left| \frac{\partial \mathbf{R}}{\partial \phi} \times \frac{\partial \mathbf{R}}{\partial \theta} \right| d\phi \, d\theta = a^2 \sin \phi \, d\phi \, d\theta \]  

(337)

Then we can compute

\[ \int_S z \, d\sigma = \int_0^{2\pi} \int_0^{\pi/2} (a \cos \phi) a^2 \sin \phi \, d\phi \, d\theta \]

\[ = a^3 \left( \int_0^{2\pi} d\theta \right) \left( \int_0^{\pi/2} \cos \phi \sin \phi \, d\phi \right) \]

\[ = a^3 \cdot 2\pi \cdot \frac{1}{2} \]

\[ = \pi a^3 \]  

(338)

From our earlier calculation of area we have

\[ \int_S d\sigma = 2\pi a^2 \]  

(339)

Thus

\[ \bar{z} = \frac{\pi a^3}{2\pi a^2} = \frac{a}{2} \]  

(340)

**example:** Suppose that the surface \( \mathcal{S} \) is the graph of a function \( z = \phi(x, y) \) with \( (x, y) \in \mathcal{R} \). As noted previously we can parametrize \( \mathcal{S} \) by

\[ \mathbf{R}(x, y) = xi + yj + \phi(x, y)k \quad (x, y) \in \mathcal{R} \]  

(341)
We also computed earlier
\[
\frac{\partial \mathbf{R}}{\partial x} \times \frac{\partial \mathbf{R}}{\partial y} = -\phi_x \mathbf{i} - \phi_y \mathbf{j} + \mathbf{k}
\] (342)

Therefore
\[
d\sigma = \left| \frac{\partial \mathbf{R}}{\partial x} \times \frac{\partial \mathbf{R}}{\partial y} \right| dxdy = \sqrt{1 + \phi_x^2 + \phi_y^2} dxdy
\] (343)

and so
\[
\int_S f(x, y, z)d\sigma = \int_R f(x, y, \phi(x, y))\sqrt{1 + \phi_x^2 + \phi_y^2} dxdy
\] (344)

In particular
\[
\text{Area of } S = \int_S d\sigma = \int_R \sqrt{1 + \phi_x^2 + \phi_y^2} dxdy
\] (345)

### 2.10 change of variables in \( \mathbb{R}^2 \)

Consider a special case of the surface integral in which the surface \( S \) lies in the \( xy \) plane. Then the parametrization has the form \( x = x(u, v), y = y(u, v), z = 0 \) for \( (u, v) \in \mathcal{R} \).

In vector form
\[
\mathbf{R}(u, v) = x(u, v)\mathbf{i} + y(u, v)\mathbf{j}
\] (346)

In this case
\[
\frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} = \det \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ x_u & y_u & 0 \\ x_v & y_v & 0 \end{pmatrix} = \det \begin{pmatrix} x_u & y_u \\ x_v & y_v \end{pmatrix} \mathbf{k}
\] (347)

and so
\[
d\sigma = \left| \frac{\partial \mathbf{R}}{\partial u} \times \frac{\partial \mathbf{R}}{\partial v} \right| dudv = \left| \frac{\partial(x, y)}{\partial(u, v)} \right| dudv
\] (348)

Then our surface integral is evaluated as
\[
\int_S f(x, y)d\sigma = \int_R f(x(u, v), y(u, v)) \left| \frac{\partial(x, y)}{\partial(u, v)} \right| dudv
\] (349)

However since \( S \) is flat the surface integral \( \int_S f(x, y)d\sigma \) is the same as the area integral \( \int_S f(x, y)dA \). It has been divided up in an irregular fashion but the result is the same.

Thus we have demonstrated the following change of variables formula:

**Theorem 14** Let a region \( S \subset \mathbb{R}^2 \) be the image of a region \( \mathcal{R} \subset \mathbb{R}^2 \) under a differentiable invertible function \( x = x(u, v), y = y(u, v) \). Then
\[
\int_S f(x, y)dA = \int_\mathcal{R} f(x(u, v), y(u, v)) \left| \frac{\partial(x, y)}{\partial(u, v)} \right| dudv
\] (350)
One can think of \((u, v)\) as new coordinates for the region \(S\). Then a short version of the theorem is
\[
dA = \frac{\partial(x, y)}{\partial(u, v)} \, dudv
\]
expressing the area element in the new coordinates.

As a special case consider polar coordinates \(x = r \cos \theta, y = r \sin \theta\). Then
\[
\frac{\partial(x, y)}{\partial(r, \theta)} = \det \begin{pmatrix} x_r & x_\theta \\ y_r & y_\theta \end{pmatrix} = \det \begin{pmatrix} \cos \theta & -r \sin \theta \\ \sin \theta & r \cos \theta \end{pmatrix} = r
\]
and so \(dA = rdrd\theta\). The change of variables formula is
\[
\int_S f(x, y) \, dA = \int_{\mathcal{R}} f(r \cos \theta, r \sin \theta) \, rdrd\theta
\]
Here \(\mathcal{R}\) is all points \((r, \theta)\) such that \((x, y) = (r \cos \theta, r \sin \theta)\) is in \(S\). Thus \(\mathcal{R}\) is just \(S\) described in polar coordinates.

**example:** Suppose \(S\) is the half-disc \(x^2 + y^2 \leq 4, y > 0\) and we want to evaluate
\[
\int_S (3 - x^2 - y^2) \, dA
\]
In polar coordinates \(S\) becomes the region \(\mathcal{R}\) defined by \(0 < r < 2, 0 < \theta < \pi\). Thus
\[
\int_S (3 - x^2 - y^2) \, dA = \int_{\mathcal{R}} (3 - r^2) \, rdrd\theta
\]
\[
= \int_0^\pi \int_0^2 (3r - r^3) \, drd\theta
\]
\[
= \pi \left[ \frac{3r^2}{2} - \frac{r^4}{4} \right]_0^2
\]
\[
= 2\pi
\]

**example:** Let \(S\) be the region bounded by the lines \(x + y = -1, x + y = 3, 2x - y = 0, 2x - y = 4\). (see figure 19) We want to evaluate the integral
\[
\int_S (x + y) \, dA
\]
We make a change of variables suggested by the boundary lines and set
\[
u = x + y \quad \text{and} \quad v = 2x - y
\]
The lines $x + y = c$ are sent to the line $u = c$ and the lines $2x - y = c$ are sent to the lines $v = c$. Thus the region $S$ is sent to the region $R$ bounded by the lines $u = -1, u = 3, v = 0, v = 4$.

For the change of variables formula we need the inverse function

$$x = \frac{u + v}{3}, \quad y = \frac{2u - v}{3},$$  \hspace{1cm} (358)

which sends $R$ back to $S$.

We compute

$$\frac{\partial(x, y)}{\partial(u, v)} = \det \begin{pmatrix} 1/3 & 1/3 \\ 2/3 & -1/3 \end{pmatrix} = -\frac{1}{3}$$  \hspace{1cm} (359)
and then

\[ \int_S (x+y)dA = \int_R u \left| \frac{\partial(x,y)}{\partial(u,v)} \right| dudv \]
\[ = \int_0^4 \int_{-1}^1 u \cdot \frac{1}{3} dudv \]
\[ = 4 \left[ \frac{u^2}{6} \right]_{-1}^1 = \frac{16}{3} \] (360)

### 2.11 change of variables in \( \mathbb{R}^3 \)

In \( \mathbb{R}^3 \) the change of variables formula is the following:

**Theorem 15** Let a region \( V \subset \mathbb{R}^3 \) be the image of a region \( R \subset \mathbb{R}^3 \) under a differentiable invertible function \( x = x(u,v,w), y = y(u,v,w), z = z(u,v,w) \). Then

\[ \int_V f(x,y,z)dV = \int_R f(x(u,v,w),y(u,v,w),z(u,v,w)) \left| \frac{\partial(x,y,z)}{\partial(u,v,w)} \right| dudvdw \] (361)

For short we can write

\[ dV = \left| \frac{\partial(x,y,z)}{\partial(u,v,w)} \right| dudvdw \] (362)

Also note that in the special case \( f(x,y,z) = 1 \) we have

\[ \text{Volume of } V = \int_V dV = \int_R \left| \frac{\partial(x,y,z)}{\partial(u,v,w)} \right| dudvdw \] (363)

**Proof.** Divide up the region \( R \) into a grid of small boxes. The \( i^{th} \) box will have a corner \( (u_i, v_i, w_i) \) and dimensions \( (\Delta u_i, \Delta v_i, \Delta w_i) \). The volume of this box is \( \Delta V_i = \Delta u_i \Delta v_i \Delta w_i \).

Let \( \Delta V_i' \) be the volume of image of the \( i^{th} \) box. The image has corners \( x_i = x(u_i, v_i, w_i), y_i = y(u_i, v_i, w_i), z_i = z(u_i, v_i, w_i) \) and curved sides. The volume is approximately the volume of a parallelopiped spanned by vectors \( a_i, b_i, c_i \) joining the corners. (See figure 20). Thus

\[ \Delta V_i' \approx |a_i \cdot (b_i \times c_i)| \] (364)

If we write the the function as

\[ R(u,v,w) = x(u,v,w)i + y(u,v,w)j + z(u,v,w)k \] (365)

67
then

\[ a_i = R(u_i + \Delta u_i, v_i, w_i) - R(u_i, v_i, w_i) \approx \frac{\partial R}{\partial u_i}(u_i, v_i, w_i) \Delta u_i \]  

(366)

and similarly

\[ b_i \approx \frac{\partial R}{\partial v}(u_i, v_i, w_i) \Delta v_i \]

\[ c_i \approx \frac{\partial R}{\partial w}(u_i, v_i, w_i) \Delta w_i \]

Therefore

\[
\Delta V'_i \approx \left| \frac{\partial R}{\partial u} \cdot \left( \frac{\partial R}{\partial v} \times \frac{\partial R}{\partial w} \right) \right| \Delta u_i \Delta v_i \Delta w_i \\
= \left| \det \begin{pmatrix}
x_u & y_u & z_u \\
x_v & y_v & z_v \\
x_w & y_w & z_w 
\end{pmatrix} \right| \Delta V_i = \left| \det \begin{pmatrix}
x_u & x_v & x_w \\
y_u & y_v & y_w \\
z_u & z_v & z_w 
\end{pmatrix} \right| \Delta V_i 
\]

(368)

So the Jacobian determinant tells how volumes increase.
Then we have

\[
\sum_i f(x_i, y_i, z_i) \Delta V_i'
= \sum_i f(x(u_i, v_i, w_i), y(u_i, v_i, w_i), z(u_i, v_i, w_i)) \left| \frac{\partial(x, y, z)}{\partial(u, v, w)}(u_i, v_i, w_i) \right| \Delta V_i
\]  

(369)

Now taking the limit as the grid size goes to zero we obtain the result (although the expression on the left is not the standard Riemann integral).

**special cases:**

1. **cylindrical coordinates:**

\[
\begin{align*}
x &= r \cos \theta \\
y &= r \sin \theta \\
z &= z
\end{align*}
\]  

(370)

In this case

\[
dV = \left| \frac{\partial(x, y, z)}{\partial(r, \theta, z)} \right| drd\theta dz = r drd\theta dz
\]  

(371)

and

\[
\int_V f(x, y, z) dV = \int_\mathcal{R} f(r \cos \theta, r \sin \theta, z) r drd\theta dz
\]  

(372)

where \( \mathcal{R} \) is the region \( V \) described in cylindrical coordinates.

2. **spherical coordinates:**

\[
\begin{align*}
x &= \rho \sin \phi \cos \theta \\
y &= \rho \sin \phi \sin \theta \\
z &= \rho \cos \phi
\end{align*}
\]  

(373)

In this case (check it!)

\[
dV = \left| \frac{\partial(x, y, z)}{\partial(\rho, \phi, \theta)} \right| d\rho d\phi d\theta = \rho^2 \sin \phi d\rho d\phi d\theta
\]  

(374)

and

\[
\int_V f(x, y, z) dV = \int_\mathcal{R} f(\rho \sin \phi \cos \theta, \rho \sin \phi \sin \theta, \rho \cos \phi) \rho^2 \sin \phi d\rho d\phi d\theta
\]  

(375)

where \( \mathcal{R} \) is the region \( V \) described in spherical coordinates.


**Example:** Suppose we want to find the volume of the quarter-cone in a sphere $V$ which is described in spherical coordinates by $0 < \rho < 1$, $0 < \phi < \pi/4$, $0 < \theta < \pi/2$. We compute

$$
\text{Volume} = \int_V dV = \int_0^{\pi/2} \int_0^{\pi/4} \int_0^1 \rho^2 \sin \phi \ d\rho \ d\phi \ d\theta
= \left( \int_0^1 \rho^2 d\rho \right) \left( \int_0^{\pi/4} \sin \phi \ d\phi \right) \left( \int_0^{\pi/2} d\theta \right)
= \frac{1}{3} \cdot \left[ -\cos \phi \right]_0^{\pi/4} \cdot \pi/2
= \frac{\pi}{6} \left( 1 - \frac{1}{\sqrt{2}} \right)
$$

(376)

### 2.12 derivatives in $\mathbb{R}^3$

In $\mathbb{R}^3$ we continue to use the notation

$$
\mathbf{R} = x \mathbf{i} + y \mathbf{j} + z \mathbf{k} = (x, y, z)
$$

(377)

A *scalar* is a function from $\mathbb{R}^3$ to $\mathbb{R}$ and has the form

$$
\mathbf{R} = u = u(\mathbf{R}) = u(x, y, z)
$$

(378)

A scalar can be drawn (not very well) by shading points in $\mathbb{R}^3$ proportional to the value of $u$ at that point. Examples of quantities that can be represented by scalars are density and temperature.

A *vector field* is a function from $\mathbb{R}^3$ to $\mathbb{R}^3$ and has the form

$$
\mathbf{v} = \mathbf{v}(\mathbf{R}) = v_1(\mathbf{R}) \mathbf{i} + v_2(\mathbf{R}) \mathbf{j} + v_3(\mathbf{R}) \mathbf{k} = v_1(x, y, z) \mathbf{i} + v_2(x, y, z) \mathbf{j} + v_3(x, y, z) \mathbf{k}
$$

(379)

A vector field can be represented by drawing a vector $\mathbf{v}(\mathbf{R})$ at the point $\mathbf{R}$ for some representative points $\mathbf{R}$. Examples of quantities that can be represented by scalars are forces and the velocity of a fluid.

We want to define various derivatives of scalars and vector fields. These are specified with the operator

$$
\nabla = \mathbf{i} \frac{\partial}{\partial x} + \mathbf{j} \frac{\partial}{\partial y} + \mathbf{k} \frac{\partial}{\partial z}
$$

(380)

which is called *del* or *nabla*. If $u$ is a scalar we define a vector field

$$
\text{gradient of } u = \nabla u = \mathbf{i} \frac{\partial u}{\partial x} + \mathbf{j} \frac{\partial u}{\partial y} + \mathbf{k} \frac{\partial u}{\partial z}
$$

(381)
If \( \mathbf{v} = v_1 \mathbf{i} + v_2 \mathbf{j} + v_3 \mathbf{k} \) is a vector field we define a scalar

\[
\text{divergence of } \mathbf{v} = \nabla \cdot \mathbf{v} = \frac{\partial v_1}{\partial x} + \frac{\partial v_2}{\partial y} + \frac{\partial v_3}{\partial z}
\]

and also a vector field

\[
\text{curl of } \mathbf{v} = \nabla \times \mathbf{v} = \text{det} \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ \partial/\partial x & \partial/\partial y & \partial/\partial z \\ v_1 & v_2 & v_3 \end{pmatrix}
\]

\[
= \mathbf{i} \left( \frac{\partial v_3}{\partial y} - \frac{\partial v_2}{\partial z} \right) + \mathbf{j} \left( \frac{\partial v_1}{\partial z} - \frac{\partial v_3}{\partial x} \right) + \mathbf{k} \left( \frac{\partial v_2}{\partial x} - \frac{\partial v_1}{\partial y} \right)
\]

Finally if \( u \) is a scalar we define

\[
\text{Laplacian of } u = \Delta u = \nabla \cdot \nabla u = \frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial y^2} + \frac{\partial^2 u}{\partial z^2}
\]

**example:** If \( u = x^2 + xy + y^2 + yz + z^2 + zx \) then

\[
\nabla u = (2x + y + z)\mathbf{i} + (2y + x + z)\mathbf{j} + (2z + z + y)\mathbf{k}
\]

\[
\Delta u = \nabla \cdot \nabla u = 2 + 2 + 2 = 6
\]

**example:** If \( \mathbf{v} = yz^2\mathbf{i} + xz^2\mathbf{j} + (2xyz + z)\mathbf{k} \) then

\[
\nabla \cdot \mathbf{v} = 2xy + 1
\]

\[
\nabla \times \mathbf{v} = \text{det} \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ \partial/\partial x & \partial/\partial y & \partial/\partial z \\ yz^2 & xz^2 & (2xyz + z) \end{pmatrix}
\]

\[
= (2xz - 2xz)\mathbf{i} + (2zy - 2zy)\mathbf{j} + (z^2 - z^2)\mathbf{k} = 0
\]

The derivatives satisfy various identities some of which we list. These hold for any scalar \( u \) and any vector field \( \mathbf{v} \), assuming only they are twice continuously differentiable.

1. \( \nabla \times \nabla u = 0 \)
2. \( \nabla \cdot (\nabla \times \mathbf{v}) = 0 \)
3. \( \nabla \cdot (u \mathbf{v}) = \nabla u \cdot \mathbf{v} + u(\nabla \cdot \mathbf{v}) \)
4. \( \nabla \cdot (\mathbf{v} \times \mathbf{w}) = (\nabla \times \mathbf{v}) \cdot \mathbf{w} - \mathbf{v} \cdot (\nabla \times \mathbf{w}) \)
5. $\nabla \times (uv) = u(\nabla \times v) + \nabla u \times v$

6. $\nabla \times (\nabla \times v) = \nabla (\nabla \cdot v) - \Delta v$ where $\Delta v = \Delta v_1 i + \Delta v_2 j + \Delta v_3 k$

**Proof of (1.)** We compute

$$\nabla \times \nabla u = \text{det} \begin{pmatrix} i & j & k \\ \frac{\partial}{\partial x} & \frac{\partial}{\partial y} & \frac{\partial}{\partial z} \\ \frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} & \frac{\partial u}{\partial z} \end{pmatrix}$$

$$= \left( \frac{\partial^2 u}{\partial y \partial z} - \frac{\partial^2 u}{\partial z \partial y} \right) i + \left( \frac{\partial^2 u}{\partial z \partial x} - \frac{\partial^2 u}{\partial x \partial z} \right) j + \left( \frac{\partial^2 u}{\partial x \partial y} - \frac{\partial^2 u}{\partial y \partial x} \right) k$$

$$= 0$$

We also note that the chain rule

$$\frac{d}{dt} u(x(t), y(t), z(t))$$

$$= \frac{\partial u}{\partial x} (x(t), y(t), z(t)) \frac{dx}{dt} + \frac{\partial u}{\partial y} (x(t), y(t), z(t)) \frac{dy}{dt} + \frac{\partial u}{\partial z} (x(t), y(t), z(t)) \frac{dz}{dt}$$

can be written in a vector notation as

$$\frac{d}{dt} u(R(t)) = \nabla u(R(t)) \cdot \frac{dR}{dt}$$

### 2.13 Gradient

One of our goals is to interpret the gradient, divergence, and curl. Here we give three interpretations of the gradient.

1. $(\nabla u)(R_0)$ is normal to the level surface $u = \text{constant}$ through $R_0$, i.e. the surface $u(R) = u(R_0)$.

To see this let $R(t)$ be any curve in the surface with $R(0) = R_0$. Thus

$$u(R(t)) = u(R_0)$$

Taking the derivative with respect to $t$ and using the chain rule gives

$$(\nabla u)(R(t)) \cdot R'(t) = 0$$

At $t = 0$ this says

$$(\nabla u)(R_0) \cdot R'(0) = 0$$

Any tangent vector to the surface at $R_0$ has the form $R'(0)$ for some curve through $R_0$. Thus $(\nabla u)(R_0)$ is normal to any tangent vector at $R_0$ and hence is normal to the surface. See figure 21.
2. If \( \mathbf{n} \) is a unit vector, then \((\nabla u)(\mathbf{R}_0) \cdot \mathbf{n}\) is the rate of change of \( u \) at \( \mathbf{R}_0 \) in the direction \( \mathbf{n} \), also called the directional derivative.

To see this use the chain rule to calculate the rate of change as

\[
\frac{d}{dt} u(\mathbf{R}_0 + t\mathbf{n})\bigg|_{t=0} = (\nabla u)(\mathbf{R}_0 + t\mathbf{n}) \cdot \frac{d}{dt}(\mathbf{R}_0 + t\mathbf{n})\bigg|_{t=0} = (\nabla u)(\mathbf{R}_0) \cdot \mathbf{n} \quad (393)
\]

3. \((\nabla u)(\mathbf{R}_0)\) is the direction of greatest increase for \( u \) at \( \mathbf{R}_0 \).

To see this consider that the direction of greatest increase is the unit vector \( \mathbf{n} \) which maximizes the directional derivative \((\nabla u)(\mathbf{R}_0) \cdot \mathbf{n}\). This occurs when \( \mathbf{n} \) is parallel to \((\nabla u)(\mathbf{R}_0)\).

**problem** Find the normal to the surface which is the graph of \( z = f(x, y) \)

**solution** The surface is

\[
u(x, y, z) = -f(x, y) + z = 0 \quad (394)
\]
A normal is
\[ \nabla u = u_x \mathbf{i} + u_y \mathbf{j} + u_z \mathbf{k} = -f_x \mathbf{i} - f_y \mathbf{j} + \mathbf{k} \] (395)
as before.

**problem** Find the direction of greatest change for \( u = x^2 + 2xy + y^2 + 3z^2 \) at the point \((1,1,1)\), i.e. find a unit vector.

**solution** The gradient at this point is
\[ \nabla u = (2x + 2y) \mathbf{i} + (2x + 2y) \mathbf{j} + 6z \mathbf{k} = 4 \mathbf{i} + 4 \mathbf{j} + 6 \mathbf{k} \] (396)
A unit vector in this direction is
\[ \mathbf{n} = \frac{\nabla u}{|\nabla u|} = \frac{4 \mathbf{i} + 4 \mathbf{j} + 6 \mathbf{k}}{\sqrt{68}} \] (397)

### 2.14 divergence theorem

The divergence is important because it appears in the following theorem.

**Theorem 16** (Divergence Theorem). Let \( \mathcal{R} \) be a solid region in \( \mathbb{R}^3 \) with boundary surface \( \mathcal{S} \). Let \( \mathbf{n} \) be the unit outward normal on \( \mathcal{S} \). Then for any continuously differentiable vector field \( \mathbf{v} \) on \( \mathcal{R} \)
\[ \int_{\mathcal{R}} \nabla \cdot \mathbf{v} \, dV = \int_{\mathcal{S}} \mathbf{v} \cdot \mathbf{n} \, d\sigma \] (398)

**Proof.** Suppose \( \mathbf{n} = n_1 \mathbf{i} + n_2 \mathbf{j} + n_3 \mathbf{k} \). If suffices to show that
\[ \int_{\mathcal{R}} \frac{\partial v_1}{\partial x} \, dV = \int_{\mathcal{S}} v_1 n_1 \, d\sigma \]
\[ \int_{\mathcal{R}} \frac{\partial v_2}{\partial y} \, dV = \int_{\mathcal{S}} v_2 n_2 \, d\sigma \]
\[ \int_{\mathcal{R}} \frac{\partial v_3}{\partial z} \, dV = \int_{\mathcal{S}} v_3 n_3 \, d\sigma \] (399)
Then adding them together gives the result.

We prove the last, the others are similar. To prove it suppose that \( \mathcal{R} \) if the region between the graphs of two functions. It is defined by \( \phi(x,y) \leq z \leq \psi(x,y) \) with with \((x,y)\) in some region \( E \) in the plane. Then we have
\[ \int_{\mathcal{R}} \frac{\partial v_3}{\partial z} \, dV = \int_{E} \left( \int_{\phi(x,y)}^{\psi(x,y)} \frac{\partial v_3}{\partial z} \right) \, dy \, dx \]
\[ = \int_{E} (v_3(x,y,\psi(x,y)) - v_3(x,y,\phi(x,y))) \, dy \, dx \] (400)
We need to show that \( \int_S v_3 n_3 \, d\sigma \) has the same expression. Now \( S \) has three parts \( S_1, S_2, S_3 \), see figure 22, and

\[
\int_S v_3 n_3 \, d\sigma = \int_{S_1} v_3 n_3 \, d\sigma + \int_{S_2} v_3 n_3 \, d\sigma + \int_{S_3} v_3 n_3 \, d\sigma \tag{401}
\]

The surface \( S_1 \) is the graph of \( z = \psi(x, y) \) and a normal vector is \( N = -\psi_x \mathbf{i} - \psi_y \mathbf{j} + \mathbf{k} \). This is an upward normal since the third component is positive. On this surface upward is outward and so the unit outward normal is

\[
\mathbf{n} = \frac{N}{|N|} = \frac{-\psi_x \mathbf{i} - \psi_y \mathbf{j} + \mathbf{k}}{\sqrt{1 + \psi_x^2 + \psi_y^2}} \tag{402}
\]

Also on \( S_1 \) we have

\[
d\sigma = \sqrt{1 + \psi_x^2 + \psi_y^2} \, dx \, dy \tag{403}
\]
Therefore
\[
\int_{S_1} v_3 n_3 \, d\sigma = \int_E v_3(x, y, \psi(x, y)) \frac{1}{\sqrt{1 + \psi_x^2 + \psi_y^2}} \sqrt{1 + \psi_x^2 + \psi_y^2} \, dxdy
\]
\[= \int_E v_3(x, y, \psi(x, y)) \, dxdy \tag{404}\]

The surface \(S_2\) is the graph of \(z = \phi(x, y)\) and an upward normal vector is \(N = -\phi_x i - \phi_y j + k\). On this surface upward is inward and so the unit outward normal is
\[
n = -\frac{N}{|N|} = \frac{\phi_x i + \phi_y j - k}{\sqrt{\phi_x^2 + \phi_y^2}} \tag{405}\]

Also on \(S_2\) we have
\[
d\sigma = \sqrt{1 + \phi_x^2 + \phi_y^2} \, dxdy \tag{406}\]

Therefore
\[
\int_{S_2} v_3 n_3 \, d\sigma = \int_E v_3(x, y, \psi(x, y)) \frac{-1}{\sqrt{1 + \phi_x^2 + \phi_y^2}} \sqrt{1 + \phi_x^2 + \phi_y^2} \, dxdy
\]
\[= -\int_E v_3(x, y, \phi(x, y)) \, dxdy \tag{407}\]

On the surface \(S_3\) we have \(n_3 = 0\) and so \(\int_{S_3} v_3 n_3 \, d\sigma = 0\)

Adding the contributions from the three surfaces gives the desired result
\[
\int_S v_3 n_3 \, d\sigma = \int_E (v_3(x, y, \psi(x, y)) - v_3(x, y, \phi(x, y))) \, dxdy \tag{408}\]

The divergence theorem can be used in various ways. Here we just offer a couple of examples illustrating what it says.

**Example** Let \(\mathcal{R}\) be a sphere of radius \(a\) with surface \(\mathcal{S}\). Check the divergence theorem for this region and the vector field \(\mathbf{v}(\mathbf{R}) = \mathbf{R} = xi + yj + zk\).

For the volume integral note that \(\nabla \cdot \mathbf{R} = 3\) so
\[
\int_{\mathcal{R}} \nabla \cdot \mathbf{v} \, dV = 3 \int_{\mathcal{R}} dV = 3 \times \text{volume of } \mathcal{R} = 3 \left(\frac{4}{3} \pi a^3\right) = 4\pi a^3 \tag{409}\]

On the other hand for a sphere the unit normal to the surface \(\mathcal{S}\) at \(\mathbf{R}\) is \(\mathbf{n} = \mathbf{R}/|\mathbf{R}|\). Hence on \(\mathcal{S}\) we have \(\mathbf{v} \cdot \mathbf{n} = \mathbf{R} \cdot \mathbf{R}/|\mathbf{R}| = |\mathbf{R}| = a\). Therefore
\[
\int_{\mathcal{S}} \mathbf{v} \cdot \mathbf{n} d\sigma = a \int_{\mathcal{S}} d\sigma = a \text{ area of } \mathcal{S} = a \left(4\pi a^2\right) = 4\pi a^3 \tag{410}\]

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which agrees with the volume integral.

**example** Let \( R \) be the region defined by \( 0 \leq z \leq 1 - x^2 - y^2 \) with \( x^2 + y^2 \leq 1 \). Let \( S \) be the boundary of \( R \). Check the divergence theorem for this region and the vector field \( \mathbf{v} = \frac{1}{2}(x^2 + y^2)\mathbf{k} \).

The surface has two pieces. The top piece called \( S_1 \) is the graph of the function \( z = 1 - x^2 - y^2 \) above the disc \( x^2 + y^2 \leq 1 \). A normal on this surface is

\[
\mathbf{N} = -\frac{\partial z}{\partial x} \mathbf{i} - \frac{\partial z}{\partial y} \mathbf{j} + \mathbf{k} = 2xi + 2yj + k \tag{411}
\]

This points upward which is outward for this surface. Thus the unit outward normal is

\[
\mathbf{n} = \frac{\mathbf{N}}{|\mathbf{N}|} = \frac{2xi + 2yj + k}{\sqrt{1 + 4x^2 + 4y^2}} \tag{412}
\]

and on the surface we have

\[
\mathbf{v} \cdot \mathbf{n} = \frac{1}{2} \frac{x^2 + y^2}{\sqrt{1 + 4x^2 + 4y^2}} \tag{413}
\]

Furthermore we have for this surface

\[
d\sigma = \sqrt{1 + 4x^2 + 4y^2}dxdy \tag{414}
\]

Combining the above we have

\[
\int_{S_1} \mathbf{v} \cdot \mathbf{n} \, d\sigma = \int_{x^2+y^2\leq1} \frac{1}{2} \frac{x^2 + y^2}{\sqrt{1 + 4x^2 + 4y^2}} \sqrt{1 + 4x^2 + 4y^2}dxdy
\]

\[
= \int_{x^2+y^2\leq1} \frac{1}{2} (x^2 + y^2)dxdy
\]

\[
= \int_{r\leq1} \frac{1}{2} r^2 \cdot rdrd\theta
\]

\[
= \frac{1}{2} \int_{0}^{1} r^3 \, dr \int_{0}^{2\pi} \, d\theta
\]

\[
= \frac{1}{2} \cdot \frac{1}{4} \cdot (2\pi) = \frac{\pi}{4} \tag{415}
\]

The bottom piece is called \( S_2 \). It is the disc \( x^2 + y^2 \leq 1, z = 0 \). The unit outward normal is \( \mathbf{n} = -\mathbf{k} \), so we have

\[
\mathbf{v} \cdot \mathbf{n} = -\frac{1}{2}(x^2 + y^2) \tag{416}
\]

The surface is flat so \( d\sigma = dxdy \) and we have

\[
\int_{S_2} \mathbf{v} \cdot \mathbf{n} \, d\sigma = -\int_{x^2+y^2\leq1} \frac{1}{2}(x^2 + y^2)dxdy = -\frac{\pi}{4} \tag{417}
\]
Combining the two pieces we have for the surface integral

\[ \int_S \mathbf{v} \cdot \mathbf{n} \, d\sigma = \int_{S_1} \mathbf{v} \cdot \mathbf{n} \, d\sigma + \int_{S_2} \mathbf{v} \cdot \mathbf{n} \, d\sigma = \frac{\pi}{4} - \frac{\pi}{4} = 0 \quad (418) \]

For the volume integral we note that \( \nabla \cdot \mathbf{v} = \partial / \partial z \left( \left( x^2 + y^2 \right) / 2 \right) = 0 \) and so

\[ \int_\mathcal{R} \nabla \cdot \mathbf{v} \, dV = 0 \quad (419) \]

which agrees with the surface integral.

**Remark:** In evaluating surface integrals we are frequently canceling awkward square roots. We can avoid this as follows. If \( S \) is a surface with unit normal vector \( \mathbf{n} \) define a formal symbol \( d\vec{\sigma} = \mathbf{n} d\sigma \) and then

\[ \int_S \mathbf{v} \cdot \mathbf{n} \, d\sigma = \int_S \mathbf{v} \cdot d\vec{\sigma} \quad (420) \]

If \( S \) is parametrized by a function \( \mathbf{R}(u, v) \) then

\[ \mathbf{n} = \pm \left( \frac{\mathbf{R}_u \times \mathbf{R}_v}{|\mathbf{R}_u \times \mathbf{R}_v|} \right) \]

\[ d\sigma = |\mathbf{R}_u \times \mathbf{R}_v| \, dudv \]

\[ d\vec{\sigma} = \pm (\mathbf{R}_u \times \mathbf{R}_v) \, dudv \quad (421) \]

If \( S \) is the graph of a function \( z = f(x, y) \) then

\[ \mathbf{n} = \pm \left( -f_x \mathbf{i} - f_y \mathbf{j} + \mathbf{k} / \sqrt{1 + f_x^2 + f_y^2} \right) \]

\[ d\sigma = \sqrt{1 + f_x^2 + f_y^2} \, dxdy \]

\[ d\vec{\sigma} = \pm \left( -f_x \mathbf{i} - f_y \mathbf{j} + \mathbf{k} \right) \, dxdy \quad (422) \]

In either case the square roots are gone from \( d\vec{\sigma} \). The only difficulty is that one must still think about which normal one wants to determine whether to take the plus sign or the minus sign.

### 2.15 Applications

We give some applications of the divergence theorem. In the first we answer the question "what is divergence?". The others are derivations of some basic partial differential equations.
A. steady fluid flow  The steady (i.e. time independent) flow of a fluid is described by the following quantities:

\[ \mathbf{v}(\mathbf{R}) = \text{velocity of the fluid at } \mathbf{R} \quad (\text{cm/sec}) \]

\[ \rho(\mathbf{R}) = \text{density of the fluid at } \mathbf{R} \quad (\text{gr/cm}^3) \]

\[ \mathbf{u}(\mathbf{R}) = \rho(\mathbf{R}) \mathbf{v}(\mathbf{R}) = \text{mass flow density} \quad (\text{gr/cm}^2 \cdot \text{sec}) \]

Also if \( S \) is a surface with unit normal \( \mathbf{n} \) we define

\[ \int_S \mathbf{u} \cdot \mathbf{n} \, d\sigma = \text{flux of } \mathbf{u} \text{ over } S \text{ in direction } \mathbf{n} \]

We want to interpret this flux and also the divergence of \( \mathbf{u} \). (They are related).

Pick a point \( \mathbf{R} \) on the the surface \( S \), let \( \Delta \sigma \) be a small piece of surface around \( \mathbf{R} \). Then (see figure 23)

\[
\text{mass through } \Delta \sigma \text{ in time } \Delta t \quad (\text{gr}) \\
\approx \text{density at } \mathbf{R} \times \text{volume through } \Delta \sigma \text{ in time } \Delta t \\
\approx \rho(\mathbf{R}) \times \Delta \sigma \times (\mathbf{v}(\mathbf{R})\Delta t) \cdot \mathbf{n} \\
= \mathbf{u}(\mathbf{R}) \cdot \mathbf{n} \Delta \sigma \Delta t
\]
If we divide by $\Delta t$ we get

$$\text{rate of mass flow thru } \Delta \sigma = \mathbf{u}(\mathbf{R}) \cdot \mathbf{n} \Delta \sigma \quad (\text{gr/sec})$$

Now sum over pieces $\Delta \sigma_i$ covering the surface $\mathcal{S}$ and take the limit as the partition becomes fine. This gives

$$\text{rate of mass flow thru } \mathcal{S} = \lim_{\Delta \sigma \to 0} \sum_i \mathbf{u}(\mathbf{R}_i) \cdot \mathbf{n}_i \Delta \sigma_i = \int_{\mathcal{S}} \mathbf{u}(\mathbf{R}) \cdot \mathbf{n} \, d\sigma \quad (\text{gr/sec})$$

Thus we have an interpretation of the flux of $\mathbf{u}$ over $\mathcal{S}$. It is the rate of mass flow through $\mathcal{S}$.

Now for the divergence of $\mathbf{u}$ at any point $\mathbf{R}$ let $D_{\epsilon}$ be a sphere of radius $\epsilon$ around $\mathbf{R}$. let $S_{\epsilon}$ be the surface of that sphere with outward normal $\mathbf{n}$. Then by the divergence theorem

$$\nabla \cdot \mathbf{u}(\mathbf{R}) = \lim_{\epsilon \to 0} \frac{1}{\text{Vol } D_{\epsilon}} \int_{D_{\epsilon}} \nabla \cdot \mathbf{u} \, dV = \frac{1}{\text{Vol } D_{\epsilon}} \left( \int_{S_{\epsilon}} \mathbf{u} \cdot \mathbf{n} \, d\sigma \right) \quad (424)$$

This is the flux divided by the volume. Thus the divergence is the rate of outward mass flow per unit volume $(\text{gr/cm}^3 \cdot \text{sec})$. A large divergence at a point means a lot of fluid is entering the system at that point.

**B. fluid dynamics**

Again we consider fluid flow, but now the velocity $\mathbf{v}(\mathbf{R}, t)$, the density $\rho(\mathbf{R}, t)$, and the mass flow density $\mathbf{u}(\mathbf{R}, t) = \rho(\mathbf{R}, t) \mathbf{v}(\mathbf{R}, t)$ all depend on the time $t$.

In our fluid consider any region $\mathcal{R}$ with surface $\mathcal{S}$ and outward normal $\mathbf{n}$. Conservation of mass says that

$$\text{rate of mass flow out of } \mathcal{S} = \text{rate of decrease of mass in } \mathcal{R} \quad (425)$$

which means that

$$\int_{\mathcal{S}} \mathbf{u} \cdot \mathbf{n} \, d\sigma = -\frac{d}{dt} \left( \int_{\mathcal{R}} \rho(\mathbf{R}, t) \, dV \right) \quad (426)$$

Use the divergence theorem on the left, and differentiate under the integral sign on the right to obtain

$$\int_{\mathcal{R}} \nabla \cdot \mathbf{u} \, dV = -\int_{\mathcal{R}} \frac{\partial \rho}{\partial t} \, dV \quad (427)$$

which is the same as

$$\int_{\mathcal{R}} \left( \nabla \cdot \mathbf{u} + \frac{\partial \rho}{\partial t} \right) \, dV = 0 \quad (428)$$
Since this holds for an arbitrary region $\mathcal{R}$ it must be that
\[ \nabla \cdot \mathbf{u} + \frac{\partial \rho}{\partial t} = 0 \] (429)
which is also written
\[ \nabla \cdot (\rho \mathbf{v}) + \frac{\partial \rho}{\partial t} = 0 \] (430)
This is the continuity equation which must be satisfied by any flow.

We can rewrite this as
\[ \rho \nabla \cdot \mathbf{v} + \nabla \rho \cdot \mathbf{v} + \frac{\partial \rho}{\partial t} = 0 \] (431)
Now define the total derivative of $\rho$ to be
\[ \frac{D\rho}{Dt} = \frac{\partial \rho}{\partial t} + \nabla \rho \cdot \mathbf{v} \] (432)
Then the continuity equation can be written as
\[ \frac{D\rho}{Dt} + \rho \nabla \cdot \mathbf{v} = 0 \] (433)

The total derivative has the interpretation
\[ \frac{D\rho}{Dt} = \text{rate of change of the density at a test particle moving in the fluid} \] (434)
To see that this is true let $\mathbf{R}(t)$ be the trajectory of the test particle. Moving with the fluid means that \( \frac{d\mathbf{R}}{dt} = \mathbf{v}(\mathbf{R}(t), t) \). The rate of change of the density at the test particle is by the chain rule
\[ \frac{d}{dt}\rho(\mathbf{R}(t), t) = \nabla \rho(\mathbf{R}(t), t) \cdot \frac{d\mathbf{R}}{dt} + \frac{\partial \rho}{\partial t}(\mathbf{R}(t), t) \]
\[ = \nabla \rho(\mathbf{R}(t), t) \cdot \mathbf{v}(\mathbf{R}(t), t) + \frac{\partial \rho}{\partial t}(\mathbf{R}(t), t) \] (435)
\[ = \frac{D\rho}{Dt}(\mathbf{R}(t), t) \]
as claimed.

The field is incompressible if the test particle sees no change in the density, that is
\[ \frac{D\rho}{Dt} = 0 \] (436)
By the second form of the continuity equation this is the same as
\[ \nabla \cdot \mathbf{v} = 0 \] (437)
C. heat equation  For any object let $T(R, t)$ be the temperature of the object at position $R$ and time $t$. We want to derive an equation which describes how the temperature evolves in time. For this we also need to consider the heat (energy) in the system. The heat density (calories/ cm$^3$) at $R, t$ is proportional to the temperature there and has the form $c\rho T(R, t)$ where the specific heat $c$ is a constant depending on the material and $\rho$ is the mass density, assumed constant. We also need the heat flow $u(R, t)$ at $R, t$. This is analogous to the mass flow in the previous examples but is now the flow of energy (calories/ cm$^2$· sec).

In our object consider any region $\mathcal{R}$ with surface $\mathcal{S}$ and outward normal $n$. Conservation of energy says that

$$\text{heat flow out of } \mathcal{S} = \text{rate of decrease of heat in } \mathcal{R}$$

which means that

$$\int_{\mathcal{S}} u \cdot n \, d\sigma = -\frac{d}{dt} \left( \int_{\mathcal{R}} c\rho T(R, t) \, dV \right)$$

Again use the divergence theorem on the left, and differentiate under the integral sign on the right to obtain

$$\int_{\mathcal{R}} \nabla \cdot u \, dV = -\int_{\mathcal{R}} c\rho \frac{\partial T}{\partial t} \, dV$$

which is the same as

$$\int_{\mathcal{R}} \left( \nabla \cdot u + c\rho \frac{\partial T}{\partial t} \right) \, dV = 0$$

Since this holds for any region $\mathcal{R}$ it must be that

$$\nabla \cdot u + c\rho \frac{\partial T}{\partial t} = 0$$

Now we need another fact. This is the thermal conduction law which says that the heat flow is proportional to the negative gradient of the temperature:

$$u = -k \nabla T$$

Here $k$ is a positive constant called the thermal conductivity. The law says that heat flows in the direction of greatest temperature decrease. Inserting this in the above equation gives

$$c\rho \frac{\partial T}{\partial t} - k \Delta T = 0$$

This is called the heat equation. If the temperature is independent of time then this becomes

$$\Delta T = 0$$

which is known as Laplace’s equation.
2.16 more line integrals

We define a line integral of vector fields. Let $C$ be a directed curve parametrized by $R(t), a \leq t \leq b$ and let $v(R)$ be a vector field defined on $C$. We define

$$\int_C v \cdot dR = \int_a^b v(R(t)) \cdot \frac{dR}{dt} \, dt$$  \hspace{1cm} (446)

Thus we replace the curve by the parameter and interpret $dR = (dR/dt)dt$. This definition turns out to be independent of parametrization as long as we respect the direction of the curve.

Let $T = [dR/dt]/|dR/dt|$  \hspace{1cm} (447)

be a unit tangent vector to the curve. Then the the vector line integral is related to a scalar line integral by

$$\int_C v \cdot dR = \int_C v \cdot T \, ds$$  \hspace{1cm} (448)

Thus we are integrating the tangential component of $v$ along the curve. To see that it is true compute

$$\int_C v \cdot T \, ds = \int_a^b v(R(t)) \cdot \frac{dR/dt}{|dR/dt|} \cdot |dR/dt| \, dt$$

$$= \int_a^b v(R(t)) \cdot dR/ dt \, dt$$  \hspace{1cm} (449)

$$= \int_C v \cdot dR$$

**Problem** Let $C$ be a straight line from $a = i + k$ to $b = 2i + j + 3k$. Let $v(R) = xi + yz k$. Evaluate $\int_C v \cdot dR$.

**Solution:** $C$ can be parametrized by $R(t) = (1 - t)a + tb$ with $0 \leq t \leq 1$ which is the same as

$$R(t) = (t + 1)i + tj + (2t + 1)k$$  \hspace{1cm} (450)

Then

$$dR = \frac{dR}{dt} \, dt = (i + j + 2k)dt$$  \hspace{1cm} (451)

Also

$$v(R(t)) = (t + 1)i + t(2t + 1)k$$  \hspace{1cm} (452)

Then we compute

$$\int_C v \cdot dR = \int_0^1 ((t + 1) + 2t(2t + 1)) \, dt = \int_0^1 (1 + 3t + 4t^2) \, dt = \frac{23}{6}$$  \hspace{1cm} (453)
We list some properties of line integrals

- For any vector fields $\mathbf{v}, \mathbf{w}$
  \[ \int_C (\mathbf{v} + \mathbf{w}) \cdot d\mathbf{R} = \int_C \mathbf{v} \cdot d\mathbf{R} + \int_C \mathbf{w} \cdot d\mathbf{R} \]  
  \[ (454) \]

- For a constant $\alpha$
  \[ \int_C \alpha \mathbf{v} \cdot d\mathbf{R} = \alpha \int_C \mathbf{v} \cdot d\mathbf{R} \]  
  \[ (455) \]

- Let $C_2$ start where $C_1$ finishes and let $C_1 + C_2$ be the curve which first traverses $C_1$ and then traverses $C_2$. Then
  \[ \int_{C_1+C_2} \mathbf{v} \cdot d\mathbf{R} = \int_{C_1} \mathbf{v} \cdot d\mathbf{R} + \int_{C_2} \mathbf{v} \cdot d\mathbf{R} \]  
  \[ (456) \]

- Let $-C$ be the curve $C$ traversed in the opposite direction. Then
  \[ \int_{-C} \mathbf{v} \cdot d\mathbf{R} = -\int_C \mathbf{v} \cdot d\mathbf{R} \]  
  \[ (457) \]

(For scalar integrals on the other hand we have $\int_{-C} f ds = \int_C f ds$)

**application:** Let $\mathbf{F}(\mathbf{R})$ be the force applied to an object at position $\mathbf{R}$. Then the line integral $\int_C \mathbf{F}(\mathbf{R}) \cdot d\mathbf{R}$ is interpreted as the work done (energy expended) in moving the object along $C$.

**another notation** If $\mathbf{v} = v_1 \mathbf{i} + v_2 \mathbf{j} + v_3 \mathbf{k}$ and $d\mathbf{R} = dx \mathbf{i} + dy \mathbf{j} + dz \mathbf{k}$ then
\[ \mathbf{v} \cdot d\mathbf{R} = v_1 dx + v_2 dy + v_3 dz \]  
\[ (458) \]

This is called a *differential form*. For us it is just a formal symbol whose integral has a meaning, but it can be given a separate precise meaning in higher mathematics. In this notation our definition of the line integral of $\mathbf{v}$ along a curve $C$ parametrized by $\mathbf{R}(t) = x(t) \mathbf{i} + y(t) \mathbf{j} + z(t) \mathbf{k}$ with $a \leq t \leq b$ is
\[ \int_C v_1 dx + v_2 dy + v_3 dz 
= \int_a^b \left( v_1(x(t), y(t), z(t)) \frac{dx}{dt} + v_2(x(t), y(t), z(t)) \frac{dy}{dt} + v_3(x(t), y(t), z(t)) \frac{dz}{dt} \right) dt \]  
\[ (459) \]

In other words we interpret $dx$ as $(dx/dt)dt$ and so forth.
This notation is also used for $\mathbb{R}^2$. If the curve $C$ is parametrized by $x = x(t), y = y(t)$ with $a \leq t \leq b$ and $M(x,y)$ and $N(x,y)$ are functions on $C$, then

$$\int_C Mdx + Ndy = \int_a^b \left( M(x(t), y(t)) \frac{dx}{dt} + N(x(t), y(t)) \frac{dy}{dt} \right) dt$$  \hspace{1cm} (460)

**definitions** A curve $C$ is *simple* if it does not intersect itself (except possibly at the endpoints). A curve is *closed* if two endpoints are the same point. See figure 24 for examples.

**Theorem 17** (Green’s Theorem). Let $C$ be a simple closed curve in the plane traversed counterclockwise with interior $R$. If $M, N$ are continuously differentiable everywhere in $R$ then

$$\int_C Mdx + Ndy = \int_R \left( \frac{\partial N}{\partial x} - \frac{\partial M}{\partial y} \right) dxdy$$  \hspace{1cm} (461)

**Proof.** Suppose the region $R$ lies between the graphs of two functions $y = p(x)$ and $y = q(x)$ with $a \leq x \leq b$ and $p(x) \leq q(x)$. Call these two curves $C_1$ and $C_2$ both
traversed in the direction of increasing $x$. Then $C = C_1 - C_2$. (see figure 25). We compute

$$
\int_{\mathcal{R}} \frac{-\partial M}{\partial y} \, dxdy = -\int_a^b \left[ \int_{p(x)}^{q(x)} \frac{\partial M}{\partial y}(x,y) \, dy \right] \, dx
$$

$$
= -\int_a^b (M(x,q(x)) - M(x,p(x))) \, dx 
$$

(462)

Similarly one shows that

$$
\int_{\mathcal{R}} \frac{-\partial N}{\partial x} \, dxdy = \int_{\mathcal{C}} Ndy
$$

(463)

Here we have used that $C_2$ can be parametrized by $x = x, y = q(x), a \leq x \leq b$ and that $C_1$ can be parametrized by $x = x, y = p(x), a \leq x \leq b$.

Adding the two equations gives the result.
**Corollary:** If $\frac{\partial N}{\partial x} = \frac{\partial M}{\partial y}$ everywhere inside a simple closed curve $C$ then

$$\int_C M\,dx + N\,dy = 0 \quad (464)$$

**Corollary:** With $M = -y/2$ and $N = x/2$

$$\frac{1}{2} \int_C -y\,dx + x\,dy = \int_\mathcal{R} dxdy = \text{area of } \mathcal{R} \quad (465)$$

**Problem** Find the area inside the ellipse

$$\frac{x^2}{a^2} + \frac{y^2}{b^2} = 1 \quad (466)$$

**Solution:** The ellipse can be parametrized by $x = a\cos t$, $y = b\sin t$ with $0 \leq t \leq 2\pi$. Then we have

$$dx = -a\sin t\,dt \quad dy = b\cos t\,dt \quad (467)$$

Hence

$$\text{Area} = \frac{1}{2} \int_C -y\,dx + x\,dy$$

$$= \frac{1}{2} \int_0^{2\pi} (ab\sin^2 t + ab\cos^2 t)\,dt$$

$$= \frac{1}{2} \int_0^{2\pi} ab\,dt = \pi ab \quad (468)$$

### 2.17 Stoke’s theorem

A surface $S$ is **orientable** if there is a continuous family of unit normal vectors $\mathbf{n}$. All the surfaces we encounter will be orientable. An example of a surface that is not orientable is the Mobius strip.

**Theorem 18** (Stoke’s theorem). Let $S$ be an orientable surface with continuous unit normal $\mathbf{n}$ and boundary $C$ which is a simple closed curve. Then for any continuously differentiable vector field $\mathbf{v}$ on $S$

$$\int_S (\nabla \times \mathbf{v}) \cdot \mathbf{n} \,d\sigma = \int_C \mathbf{v} \cdot d\mathbf{R} \quad (469)$$

provided $C$ is traversed in a right-handed sense relative to $\mathbf{n}$. (see figure 26)
Proof. We have

\[ (\nabla \times \mathbf{v}) \cdot \mathbf{n} = \det \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ \frac{\partial v_1}{\partial x} & \frac{\partial v_2}{\partial y} & \frac{\partial v_3}{\partial z} \\ v_1 & v_2 & v_3 \end{pmatrix} \cdot \mathbf{n} \]

\[ = \left( \frac{\partial v_3}{\partial y} - \frac{\partial v_2}{\partial z} \right) n_1 + \left( \frac{\partial v_1}{\partial z} - \frac{\partial v_3}{\partial x} \right) n_2 + \left( \frac{\partial v_2}{\partial x} - \frac{\partial v_1}{\partial y} \right) n_3 \]  

(470)

We pick out the the \( v_1 \) part of this and show

\[ \int_S \left( \frac{\partial v_1}{\partial z} n_2 - \frac{\partial v_1}{\partial y} n_3 \right) \, d\sigma = \int_C v_1 \, dx \]  

(471)

There will be similar equations for the \( v_2 \) part and the \( v_3 \) and when we add them together we get the result.

Suppose that \( S \) is given as the graph of a function \( z = \phi(x, y) \) with \((x, y)\) in \( E \) and upward normal \( \mathbf{n} \). (see figure 27). Then

\[ \mathbf{n} d\sigma = d\tilde{\sigma} = (-\phi_x \mathbf{i} - \phi_y \mathbf{j} + \mathbf{k}) \, dxdy \]  

(472)
Here $C'$ is a boundary curve for $E$ and the second to last step follows by Green's theorem. The last step follows since if $(x,y)$ traverses $C'$ then $(x,y,\phi(x,y))$ traverses $C$.

**example:** Let $S$ be the graph of the paraboloid $z = 1 - x^2 - y^2$ which lies above the disc $x^2 + y^2 \leq 1$ and let $n$ be the the upward normal. Also let $v = yi + zj + xk$. We check Stokes theorem in this case.
First for the surface integral we have

\[
\mathbf{n} d\sigma = d\vec{\sigma} = (-\frac{\partial z}{\partial x} \mathbf{i} - \frac{\partial z}{\partial y} \mathbf{j} + \mathbf{k}) \, dx\, dy = (2x \mathbf{i} + 2y \mathbf{j} + \mathbf{k}) \, dx\, dy
\]  

(474)

Also

\[
\nabla \times \mathbf{v} = \det \begin{pmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ \partial/\partial x & \partial/\partial y & \partial/\partial z \\ y & z & x \end{pmatrix} = -\mathbf{i} - \mathbf{j} - \mathbf{k}
\]

(475)

Therefore

\[
\int_S (\nabla \times \mathbf{v}) \cdot d\vec{\sigma} = \int_{x^2+y^2 \leq 1} (-2x - 2y - 1) \, dx\, dy
\]

\[
= \int_0^{2\pi} \int_0^1 (-2r \cos \theta - 2r \sin \theta - 1) \, rdrd\theta
\]

\[
= 2\pi \int_0^1 (-r) \, dr = -\pi
\]  

(476)

For the line line integral note that \( C \), the boundary of \( S \), is the circle \( x^2 + y^2 = 1, z = 0 \). We want to go around it counterclockwise so we parametrize by

\[
\mathbf{R} = \cos t \, \mathbf{i} + \sin t \, \mathbf{j} \quad 0 \leq t \leq 2\pi
\]  

(477)

Then

\[
d\mathbf{R} = (\sin t \, \mathbf{i} + \cos t \, \mathbf{j}) \, dt
\]

\[
\mathbf{v} = \sin t \, \mathbf{i} + \cos t \, \mathbf{k}
\]  

(478)

and so

\[
\int_C \mathbf{v} \cdot d\mathbf{R} = \int_0^{2\pi} (-\sin^2 t) \, dt = -\pi
\]

(479)

as expected.

\textbf{circulation:} Suppose \( \mathbf{v}(\mathbf{R}) \) describes the velocity of a fluid at \( \mathbf{R} \), and \( C \) is a directed simple closed curve. We define

\[
\text{circulation of } \mathbf{v} \text{ around } C = \int_C \mathbf{v} \cdot d\mathbf{R} = \int_C (\mathbf{v} \cdot \mathbf{T}) ds
\]

(480)

This tells how much the fluid is circulating around the curve. See figure 28.

\textbf{what is curl?} Now we can answer this question, still thinking of \( \mathbf{v}(\mathbf{R}) \) as the velocity of a fluid at \( \mathbf{R} \). Given \( \mathbf{R} \) and a unit vector \( \mathbf{n} \), let \( S_\mathbf{n} \) be the disc centered on \( \mathbf{R} \) with
radius $\epsilon$ and normal $\mathbf{n}$, and let $\mathcal{C}_\epsilon$ be the circle which is the boundary of $\mathcal{S}_\epsilon$. Then we have by Stoke’s theorem

$$(\nabla \times \mathbf{v}(\mathbf{R}) \cdot \mathbf{n} = \lim_{\epsilon \to 0} \frac{1}{\text{area of } \mathcal{S}_\epsilon} \int_{\mathcal{S}_\epsilon} (\nabla \times \mathbf{v}) \cdot \mathbf{n} \, d\sigma$$

$$= \lim_{\epsilon \to 0} \frac{1}{\text{area of } \mathcal{S}_\epsilon} \int_{\mathcal{C}_\epsilon} \mathbf{v} \cdot d\mathbf{R}$$

$$= \text{circulation density of } \mathbf{v} \text{ around } \mathbf{n} \text{ at } \mathbf{R}. \quad (481)$$

**general remark:** In thinking about the fundamental theorem of calculus, Green’s theorem, Stoke’s theorem, and the divergence theorem note that they all have the form

$$\int_{\text{region}} (\text{derivative of function}) = \int_{\text{boundary of region}} (\text{function}) \quad (482)$$

This may help in remembering them. It also suggests that there is a more general theorem of this form which holds in any dimension. This is true. It is also called Stoke’s theorem and uses a general theory of differential forms.
2.18 still more line integrals

We investigate when line integrals are independent of the path.

**Theorem 19** If \( \mathbf{v} = \nabla u \) and \( C \) is any curve from \( R_0 \) to \( R_1 \) then

\[
\int_C \mathbf{v} \cdot d\mathbf{R} = \int_C \nabla u \cdot d\mathbf{R} = u(R_1) - u(R_0) \tag{483}
\]

So in this case the integral is independent of the path taken from \( R_0 \) to \( R_1 \). We can write without ambiguity.

\[
\int_{R_0}^{R_1} \mathbf{v} \cdot d\mathbf{R} \tag{484}
\]

Another way to write this is to define the differential \( du \) by

\[
du = \frac{\partial u}{\partial x} \, dx + \frac{\partial u}{\partial y} \, dy + \frac{\partial u}{\partial z} \, dz = \nabla u \cdot d\mathbf{R} \tag{485}
\]

Then the theorem says that for any path \( C \) from \( R_0 \) to \( R_1 \)

\[
\int_C du = u(R_1) - u(R_0) \tag{486}
\]

**Proof.** Suppose \( C \) is parametrized by \( \mathbf{R}(t) \) with \( a \leq t \leq b \). Then \( \mathbf{R}(a) = R_0 \) and \( \mathbf{R}(b) = R_1 \) and we have by the chain rule

\[
\int_C \nabla u \cdot d\mathbf{R} = \int_a^b \nabla u(\mathbf{R}(t)) \cdot \frac{d\mathbf{R}}{dt} \, dt
\]

\[
= \int_a^b \frac{d}{dt}(u(\mathbf{R}(t))) \, dt
\]

\[
= u(\mathbf{R}(b)) - u(\mathbf{R}(a))
\]

\[
= u(R_1) - u(R_0) \tag{487}
\]

application: Suppose \( \mathbf{F}(\mathbf{R}) \) is the force applied to an object at \( \mathbf{R} \). As noted before \( \int_C \mathbf{F} \cdot d\mathbf{R} \) is the work done in moving the object along \( C \). If \( C \) goes from \( R_0 \) to \( R_1 \) and \( \mathbf{F} = \nabla u \) then the work is \( u(R_1) - u(R_0) \) independent of \( C \). In this case one says that the force is \textit{conservative} and \( u \) is called the \textit{potential energy}.

For example the electrostatic force around a charge at the origin has the form

\[
\mathbf{F}(\mathbf{R}) = -C \frac{\mathbf{R}}{|\mathbf{R}|^3} \tag{488}
\]
This is a conservative force with potential (check it)

\[ u(\mathbf{R}) = \frac{C}{|\mathbf{R}|} \]  

**Definition:** A region \( \mathcal{R} \) in \( \mathbb{R}^3 \) is *connected* if any two points in \( \mathcal{R} \) can be joined by a continuous curve in \( \mathcal{R} \).

**Theorem 20** Let \( \mathbf{v} \) be a vector field in a connected region \( \mathcal{R} \) in \( \mathbb{R}^3 \). Then the following statements are equivalent (i.e. either they are all true or all false)

1. \( \int_{\mathcal{C}} \mathbf{v} \cdot d\mathbf{R} = 0 \) for any closed curve \( \mathcal{C} \) in \( \mathcal{R} \).

2. \( \int_{\mathcal{C}} \mathbf{v} \cdot d\mathbf{R} = 0 \) is independent of the path \( \mathcal{C} \) in \( \mathcal{R} \) (\( \mathcal{C} \) with fixed endpoints in \( \mathcal{R} \)).

3. \( \mathbf{v} = \nabla u \) for some function \( u \) in \( \mathcal{R} \).

**Proof.** First we show that (1.) implies (2.). If (1.) is true and \( \mathcal{C}_1 \) and \( \mathcal{C}_2 \) are any two paths from \( \mathbf{R}_0 \) to \( \mathbf{R}_1 \) then \( \mathcal{C}_1 - \mathcal{C}_2 \) is a closed curve and so

\[ \int_{\mathcal{C}_1} \mathbf{v} \cdot d\mathbf{R} - \int_{\mathcal{C}_2} \mathbf{v} \cdot d\mathbf{R} = \int_{\mathcal{C}_1 - \mathcal{C}_2} \mathbf{v} \cdot d\mathbf{R} = 0 \]  

Thus (2.) is true.

By a similar argument one can show that (2.) implies (1.). Thus (1.) and (2.) are equivalent.

We already know that (3.) implies (2.). Thus we need only show that (2.) implies (3.). Assuming (2.) we define for any point \( \mathbf{R}_0 \) in \( \mathcal{R} \)

\[ u(\mathbf{R}) = \int_{\mathbf{R}_0}^{\mathbf{R}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' \]  

(491)

Because \( \mathcal{R} \) is connected there are paths from \( \mathbf{R}_0 \) to \( \mathbf{R} \) and because integrals are independent of path we do not have to specify which one.
We have to show that the gradient of $u$ is $v$. We compute

$$\frac{\partial u}{\partial x}(\mathbf{R}) = \lim_{h \to 0} \frac{1}{h} (u(\mathbf{R} + h\mathbf{i}) - u(\mathbf{R}))$$

$$= \lim_{h \to 0} \frac{1}{h} \left[ \int_{\mathbf{R} + h\mathbf{i}}^{\mathbf{R}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' - \int_{\mathbf{R}}^{\mathbf{R} + h\mathbf{i}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' \right]$$

$$= \lim_{h \to 0} \frac{1}{h} \left[ \int_{\mathbf{R}}^{\mathbf{R} + h\mathbf{i}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' + \int_{\mathbf{R}}^{\mathbf{R} + h\mathbf{i}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' \right]$$

$$= \lim_{h \to 0} \frac{1}{h} \left[ \int_{\mathbf{R}}^{\mathbf{R} + h\mathbf{i}} \mathbf{v}(\mathbf{R}') \cdot d\mathbf{R}' \right]$$

$$= \lim_{h \to 0} \frac{1}{h} \left[ \int_{0}^{h} \mathbf{v}_1(\mathbf{R} + t\mathbf{i}) dt \right]$$

$$= v_1(\mathbf{R})$$

(492)

Here in the second to last step we have chosen a particular path from $\mathbf{R}$ to $\mathbf{R} + h\mathbf{i}$, namely $\mathbf{R}' = \mathbf{R} + t\mathbf{i}$ with $0 \leq t \leq h$ and $d\mathbf{R}' = dt$.

Similarly one shows that $(\partial u/\partial y)(\mathbf{R}) = v_2(\mathbf{R})$ and $(\partial u/\partial z)(\mathbf{R}) = v_3(\mathbf{R})$ and hence $\nabla u(\mathbf{R}) = \mathbf{v}(\mathbf{R})$. This completes the proof.

**example:** Consider again the vector field

$$\mathbf{v} = yz^2 \mathbf{i} + xz^2 \mathbf{j} + (2xyz + z) \mathbf{k}$$

(493)

Are integrals $\int_C \mathbf{v} \cdot d\mathbf{R}$ independent of path in $\mathbb{R}^3$? According to the theorem this is the same as asking whether $\mathbf{v} = \nabla u$. If there is such a $u$ it must satisfy

$$\frac{\partial u}{\partial x} = yz^2$$

$$\frac{\partial u}{\partial y} = xz^2$$

$$\frac{\partial u}{\partial z} = 2xyz + z$$

(494)

The first equation says that $u = xyz^2 + h(y, z)$ for some function $h$. Taking derivatives of this we have

$$\frac{\partial u}{\partial y} = xz^2 + \frac{\partial h}{\partial y}$$

$$\frac{\partial u}{\partial z} = 2xyz + \frac{\partial h}{\partial z}$$

(495)

Comparing this with the second and third equations above yields

$$\frac{\partial h}{\partial y} = 0$$

$$\frac{\partial h}{\partial z} = z$$

(496)
The second equation says \( h(y, z) = \frac{1}{2}z^2 + g(y) \) for some function \( g \). But then the first equation implies that \( g(y) = C \) for some constant \( C \). Thus \( h(y, z) = \frac{1}{2}z^2 + C \) for some constant \( C \). Therefore

\[
u = xyz^2 + \frac{1}{2}z^2 + C
\]

is the function we are looking for and the answer is yes.

The theorem actually holds in any dimension. We state the theorem for \( \mathbb{R}^2 \) in the language of differential forms.

**Theorem 21** Let \( M, N \) be functions in a connected region \( \mathcal{R} \) in \( \mathbb{R}^2 \). Then the following statements are equivalent

1. \( \int_C Mdx + Ndy = 0 \) for any closed curve \( C \) in \( \mathcal{R} \).
2. \( \int_C Mdx + Ndy = 0 \) is independent of the path \( C \) in \( \mathcal{R} \).
3. \( M = \partial u/\partial x \) and \( N = \partial u/\partial y \) for some function \( u \) in \( \mathcal{R} \), i.e. \( Mdx + Ndy = du \).

**Definition:** A region \( \mathcal{R} \) is *simply connected* if it is connected and every simple closed curve can be continuously shrunk to a point without leaving \( \mathcal{R} \). Figure 29 gives some examples.

**Theorem 22** If \( \mathbf{v} \) is a continuously differentiable vector field in a simply connected region \( \mathcal{R} \) then the conditions (1.), (2.), (3.) of the last two theorems are equivalent to

\[
\begin{aligned}
\nabla \times \mathbf{v} &= 0 \\
\partial N/\partial x - \partial M/\partial y &= 0
\end{aligned}
\]

in \( \mathcal{R} \subset \mathbb{R}^2 \) (498)

**Proof.** We give the proof for \( \mathbb{R}^3 \). (3.) says that \( \mathbf{v} = \nabla u \) and we know that this implies \( \nabla \times \mathbf{v} = 0 \) which is (4.).

On the other hand suppose that (4.) is true. Since \( \mathcal{R} \) is simply connected every simple closed curve \( C \) in \( \mathcal{R} \) is the boundary of a surface \( S \) in \( \mathcal{R} \) - the curves shrinking down \( C \) sweep out the surface. Then by Stokes theorem

\[
\int_C \mathbf{v} \cdot d\mathbf{R} = \int_S (\nabla \times \mathbf{v}) \cdot \mathbf{n} \, d\sigma = 0
\]

Thus (1.) is true for simple closed curves. But closed curve that is not simple can be broken up into pieces that are simple. Thus (1.) is true in general.
example Let $v = yi + 2xj + zk$. Are line integrals $\int_C v \cdot dR$ independent of path in $\mathbb{R}^3$? Since $\mathbb{R}^3$ is simply connected we just have to check whether the curl is zero. We have

$$\nabla \times v = \begin{pmatrix} \frac{\partial}{\partial x} & \frac{\partial}{\partial y} & \frac{\partial}{\partial z} \\ i & j & k \end{pmatrix} = k \quad (500)$$

Since this is not zero the answer is no.

example Let $M = -y/(x^2 + y^2)$ and $N = x/(x^2 + y^2)$. Are line integrals $\int_C Mdx + Ndy$ independent of path

1. in $\mathbb{R}^2$? Actually we cannot ask the question in $\mathbb{R}^2$ since $M, N$ are not defined at the origin.
2. in \( \mathbb{R}^2 \) with the origin deleted? This is not simply connected so we cannot use the test \( \partial N/\partial x - \partial M/\partial y = 0 \). No conclusion.

3. in \( \mathbb{R}^2 \) with the negative \( x \)-axis deleted? This is simply connected so we can apply the test and if we do we find that \( \partial N/\partial x - \partial M/\partial y = 0 \) (check it). So the answer is yes.

We can say more. Since we do have path independence there must be a \( u \) so that \( M = \partial u/\partial x \) and \( N = \partial u/\partial y \). By solving these equations one finds

\[
 u(x, y) = \tan^{-1} \left( \frac{y}{x} \right) = \text{polar angle of } (x, y) \tag{501}
\]

This works in the region (3.), but not in the region (2.) since the function is not continuous across the negative \( x \)-axis. It takes the value \( \pi \) from above and \( -\pi \) from below. In fact there is no function in the entire region (2.) and the answer to the question is no.

**definitions** A vector field \( \mathbf{v} \) is *irrotational* if \( \nabla \times \mathbf{v} = 0 \). We have seen that in a simply connected region this occurs if and only if \( \mathbf{v} = \nabla u \) for some scalar \( u \). A vector field is *solenoidal* is \( \nabla \cdot \mathbf{v} = 0 \). One can show that in a simply connected region this occurs if and only if \( \mathbf{v} = \nabla \times \mathbf{w} \) for some vector field \( \mathbf{w} \).

It is a theorem that in a simply connected region every vector field \( \mathbf{v} \) can be written in the form \( \mathbf{v} = \mathbf{v}_1 + \mathbf{v}_2 \) where \( \mathbf{v}_1 \) is solenoidal and \( \mathbf{v}_2 \) is irrotational.

### 2.19 more applications

We discuss some applications to electromagnetism. The players in our drama are an electric field \( \mathbf{E}(\mathbf{R}, t) \), a magnetic field \( \mathbf{B}(\mathbf{R}, t) \), a current density \( \mathbf{j}(\mathbf{R}, t) \) and a charge density \( \rho(\mathbf{R}, t) \). These are time dependent vector fields, except the charge density which is a time dependent scalar.

The charge density and the current obey the charge conservation equation

\[
 \nabla \cdot \mathbf{j} + \frac{d \rho}{dt} = 0 \tag{502}
\]

This can be derived just as we derived the continuity equation for fluid flow. The fields obey *Maxwell’s equations*:

\[
 \begin{align*}
 \nabla \cdot \mathbf{E} &= 4\pi \rho \\
 \nabla \cdot \mathbf{B} &= 0 \\
 \nabla \times \mathbf{E} &= -\frac{1}{c} \frac{\partial \mathbf{B}}{\partial t} \\
 \nabla \times \mathbf{B} &= \frac{1}{c} \frac{\partial \mathbf{E}}{\partial t} + \frac{4\pi}{c} \mathbf{j}
\end{align*} \tag{503}
\]
Here $c$ is a constant of nature ($3 \times 10^{10} \text{cm/sec}$).

We work out some special cases, assuming we are in a simply connected region.

1. Suppose $\mathbf{B}$ is constant in the time $t$. Then third equation says that $\nabla \times \mathbf{E} = 0$. This means that $\mathbf{E}$ is the gradient of a scalar. We write

$$\mathbf{E} = -\nabla \Phi$$

(504)

and $\Phi$ is called the electromagnetic potential. Inserting this equation into the first gives

$$-\Delta \Phi = 4\pi \rho$$

(505)

This known as Poisson’s equation. It is easier to solve than the full Maxwell system and a great deal of mathematics is devoted to its solution. Especially important is the case $\rho = 0$ in which case we again have Laplace’s equation:

$$\Delta \Phi = 0$$

(506)

2. The second equation $\nabla \cdot \mathbf{B} = 0$ implies that

$$\mathbf{B} = \nabla \times \mathbf{A}$$

(507)

for some vector field $\mathbf{A}$ known as the magnetic potential. The potential can be written $\mathbf{A} = \mathbf{A}_1 + \mathbf{A}_2$ where $\nabla \cdot \mathbf{A}_1 = 0$ and $\nabla \times \mathbf{A}_2 = 0$. However $\mathbf{A}_2$ does not contribute to $\mathbf{B}$ so we can take $\mathbf{A} = \mathbf{A}_1$ and still have $\mathbf{B} = \nabla \times \mathbf{A}$. Then $\nabla \cdot \mathbf{A} = \nabla \cdot \mathbf{A}_1 = 0$.

If $\mathbf{E}$ is constant in time, then the last Maxwell equation becomes

$$\nabla \times (\nabla \times \mathbf{A}) = \frac{4\pi}{c} \mathbf{j}$$

(508)

But by one of our vector identities

$$\nabla \times (\nabla \times \mathbf{A}) = \nabla (\nabla \cdot \mathbf{A}) - \Delta \mathbf{A} = -\Delta \mathbf{A}$$

(509)

thus our equation becomes

$$-\Delta \mathbf{A} = \frac{4\pi}{c} \mathbf{j}$$

(510)

Thus each component of $\mathbf{A}$ satisfies Poisson’s equation and so is amenable to solution.

3. Now suppose both $\rho$ and $\mathbf{j}$ are zero. Then the equations become

$$\nabla \cdot \mathbf{E} = 0$$

$$\nabla \cdot \mathbf{B} = 0$$

$$\nabla \times \mathbf{E} = -\frac{1}{c} \frac{\partial \mathbf{B}}{\partial t}$$

(511)

$$\nabla \times \mathbf{B} = \frac{1}{c} \frac{\partial \mathbf{E}}{\partial t}$$
Taking the curl of the last equation, and using the third yields

$$\nabla \times (\nabla \times \mathbf{B}) = \frac{1}{c} \frac{\partial}{\partial t} (\nabla \times \mathbf{E}) = -\frac{1}{c^2} \frac{\partial^2 \mathbf{B}}{\partial t^2} \quad (512)$$

But on the other hand

$$\nabla \times (\nabla \times \mathbf{B}) = \nabla (\nabla \cdot \mathbf{B}) - \nabla^2 \mathbf{B} = -\nabla^2 \mathbf{B} \quad (513)$$

Comparing the last two gives

$$\frac{1}{c^2} \frac{\partial^2 \mathbf{B}}{\partial t^2} - \nabla^2 \mathbf{B} = 0 \quad (514)$$

This is called the wave equation and has also been studied at length. By a similar argument we can also show the \( \mathbf{E} \) obeys the same equation:

$$\frac{1}{c^2} \frac{\partial^2 \mathbf{E}}{\partial t^2} - \nabla^2 \mathbf{E} = 0 \quad (515)$$

A characteristic feature of solutions of these equations is that disturbances propagate with speed \( c \). But \( c \) is the speed of light. This makes it possible to interpret light as waves in the electromagnetic field. This was one of the triumphs of Maxwell’s equations.

### 2.20 general coordinate systems

Consider a differentiable invertible function

$$x = x(u_1, u_2, u_3)$$
$$y = y(u_1, u_2, u_3)$$
$$z = z(u_1, u_2, u_3) \quad (516)$$

from some region \( \mathcal{R'} \subset \mathbb{R}^3 \) onto \( \mathcal{R} \subset \mathbb{R}^3 \). Then \( \mathbf{u} = (u_1, u_2, u_3) \) in \( \mathcal{R'} \) can be considered as new coordinates for \( \mathcal{R} \). These were considered previously in our change of variables formula for \( \mathbb{R}^3 \). In a vector notation we have

$$\mathbf{R}(\mathbf{u}) = x(u_1, u_2, u_3) \mathbf{i} + y(u_1, u_2, u_3) \mathbf{j} + z(u_1, u_2, u_3) \mathbf{k} \quad (517)$$

As before we let \( \partial \mathbf{R}/\partial u_i \) be the tangent vector to a \( u_i \)-line. The length of these vectors are called the scale factors:

$$h_i = h_i(\mathbf{u}) = |\partial \mathbf{R}/\partial u_i| \quad (518)$$

We also consider

$$\mathbf{e}_i = \mathbf{e}_i(\mathbf{u}) = \frac{1}{h_i} \frac{\partial \mathbf{R}}{\partial u_i} \quad (519)$$
which are unit tangent vectors to the \( u_i \)-lines.

Our assumption of invertibility implies that
\[
\frac{\partial (x, y, z)}{\partial (u_1, u_2, u_3)} \neq 0
\]

(This is a converse to the inverse function theorem). This is a determinant with columns \( \partial R/\partial u_i \). It follows that \( \partial R/\partial u_1, \partial R/\partial u_2, \partial R/\partial u_3 \) are linearly independent. Hence \( e_1, e_2, e_3 \) are linearly independent and so form a basis for \( \mathbb{R}^3 \). Any vector \( v \) in \( \mathbb{R}^3 \) can be uniquely written in the form
\[
v = v_1 e_1 + v_2 e_2 + v_3 e_3 = \sum_{i=1}^{3} v_i e_i
\]

The new coordinates are said to be **orthogonal** if
\[
\frac{\partial R}{\partial u_i} \cdot \frac{\partial R}{\partial u_j} = 0 \quad \text{for } i \neq j
\]

It is equivalent to say that
\[
e_i \cdot e_j = \begin{cases} 0 & i \neq j \\ 1 & i = j \end{cases}
\]

and the \( e_i \) form an orthonormal basis. Then \( v \cdot e_i = v_i \) and we can write
\[
v = \sum_{i=1}^{3} (v \cdot e_i) e_i
\]

For the rest of this section we assume we have an orthogonal coordinate system.

We give some examples of orthogonal coordinate systems.

1. (cylindrical coordinates) These are defined by
\[
R = r \cos \theta i + r \sin \theta j + zk
\]

We computed \( \partial R/\partial r, \partial R/\partial \theta, \partial R/\partial z \), saw that they were orthogonal, and computed \( e_r, e_\theta, e_z \). The scale factors are
\[
h_r = |\partial R/\partial r| = 1 \\
h_\theta = |\partial R/\partial \theta| = r \\
h_z = |\partial R/\partial z| = 1
\]
2. (spherical coordinates) These are defined by

\[ \mathbf{R} = \rho \sin \phi \cos \theta \mathbf{i} + \rho \sin \phi \sin \theta \mathbf{j} + \rho \cos \phi \mathbf{k} \quad (527) \]

We computed \( \partial \mathbf{R} / \partial \rho \), \( \partial \mathbf{R} / \partial \phi \), \( \partial \mathbf{R} / \partial \theta \). One can check that they are orthogonal. Also we computed \( e_\rho, e_\phi, e_\theta \). The scale factors are

\[ h_\rho = |\partial \mathbf{R} / \partial \rho| = 1 \]
\[ h_\phi = |\partial \mathbf{R} / \partial \phi| = \rho \]
\[ h_\theta = |\partial \mathbf{R} / \partial \theta| = \rho \sin \phi \quad (528) \]

3. (parabolic coordinates) These are defined by

\[ \mathbf{R} = (u_1 u_2 \cos u_3) \mathbf{i} + (u_1 u_2 \sin u_3) \mathbf{j} + \frac{1}{2}(u_1^2 - u_2^2) \mathbf{k} \quad (529) \]

We compute

\[ \frac{\partial \mathbf{R}}{\partial u_1} = u_2 \cos u_3 \mathbf{i} + u_2 \sin u_3 \mathbf{j} + u_1 \mathbf{k} \]
\[ \frac{\partial \mathbf{R}}{\partial u_2} = u_1 \cos u_3 \mathbf{i} + u_1 \sin u_3 \mathbf{j} - u_2 \mathbf{k} \]
\[ \frac{\partial \mathbf{R}}{\partial u_3} = u_1 u_2 (-\sin u_3) \mathbf{i} + u_1 u_2 (\cos u_3) \mathbf{j} \quad (530) \]

These are orthogonal since for example

\[ \frac{\partial \mathbf{R}}{\partial u_1} \cdot \frac{\partial \mathbf{R}}{\partial u_2} = u_1 u_2 \cos^2 u_3 + u_1 u_2 \sin^2 u_3 - u_1 u_2 = 0 \quad (531) \]

The other pairs are similar. The scale factors are

\[ h_1 = |\partial \mathbf{R} / \partial u_1| = \sqrt{u_1^2 + u_2^2} \]
\[ h_2 = |\partial \mathbf{R} / \partial u_2| = \sqrt{u_1^2 + u_2^2} \]
\[ h_3 = |\partial \mathbf{R} / \partial u_3| = u_1 u_2 \quad (532) \]

From the above one can compute \( e_1, e_2, e_3 \).

Orthogonal coordinate systems are special; there are only a finite number of them. A vector field \( \mathbf{v}(\mathbf{R}) \) in new coordinates \( \mathbf{R}(\mathbf{u}) \) is

\[ \mathbf{\hat{v}}(\mathbf{u}) = \mathbf{v}(\mathbf{R}(\mathbf{u})) \quad (533) \]

The vector field in new coordinates and new basis \( \mathbf{e}_i(\mathbf{u}) \) is

\[ \mathbf{\hat{v}}(\mathbf{u}) = \sum_{i=1}^{3} v_i(\mathbf{u}) \mathbf{e}_i(\mathbf{u}) \quad v_i(\mathbf{u}) = \mathbf{\hat{v}}(\mathbf{u}) \cdot \mathbf{e}_i(\mathbf{u}) \quad (534) \]
**example:** Consider the vector field
\[ \mathbf{v} = yi + xj + zk \]  
(535)

In cylindrical coordinates this is
\[ \hat{\mathbf{v}} = r \sin \theta \mathbf{i} + r \cos \theta \mathbf{j} + zk \]  
(536)

We want to express it in the cylindrical basis
\[ \mathbf{e}_r = \cos \theta \mathbf{i} + \sin \theta \mathbf{j} \]
\[ \mathbf{e}_\theta = -\sin \theta \mathbf{i} + \cos \theta \mathbf{j} \]
\[ \mathbf{e}_z = k \]  
(537)

It will have the form
\[ \hat{\mathbf{v}} = v_r \mathbf{e}_r + v_\theta \mathbf{e}_\theta + v_z \mathbf{e}_z \]  
(538)

where
\[ v_r = \hat{\mathbf{v}} \cdot \mathbf{e}_r = 2r \cos \theta \sin \theta = r \sin(2\theta) \]
\[ v_\theta = \hat{\mathbf{v}} \cdot \mathbf{e}_\theta = -r \sin^2 \theta + r \cos^2 \theta = r \cos(2\theta) \]  
(539)
\[ v_z = \hat{\mathbf{v}} \cdot \mathbf{e}_z = z \]

Thus
\[ \hat{\mathbf{v}} = r \sin(2\theta) \mathbf{e}_r + r \cos(2\theta) \mathbf{e}_\theta + z \mathbf{e}_z \]  
(540)

Now suppose we are given a scalar function \( f(\mathbf{R}) \) and a vector field \( \mathbf{v}(\mathbf{R}) \). We consider the problem of expressing the derivatives \( \nabla f, \nabla \cdot \mathbf{v}, \nabla \times \mathbf{v}, \Delta f \) in a general orthogonal coordinate system \( \mathbf{R} = \mathbf{R}(\mathbf{u}) \) with orthonormal basis \( \mathbf{e}_i(\mathbf{u}) \).

We start with the gradient. In new coordinates \( \hat{f}(\mathbf{u}) = f(\mathbf{R}(\mathbf{u})) \) The new gradient is defined as
\[ (\text{grad} \hat{f})(\mathbf{u}) \equiv (\nabla f)(\mathbf{R}(\mathbf{u})) = \sum_{i=1}^{3} [(\nabla f)(\mathbf{R}(\mathbf{u})) \cdot \mathbf{e}_i(\mathbf{u})] \cdot \mathbf{e}_i(\mathbf{u}) \]  
(541)

But the derivatives are still in Cartesian coordinates. By the chain rule
\[ \frac{\partial \hat{f}}{\partial u_i}(\mathbf{u}) = (\nabla f)(\mathbf{R}(\mathbf{u})) \cdot \frac{\partial \mathbf{R}}{\partial u_i} = (\nabla f)(\mathbf{R}(\mathbf{u})) \cdot \mathbf{e}_i(\mathbf{u}) h_i(\mathbf{u}) \]  
(542)

Therefore
\[ (\nabla f)(\mathbf{R}(\mathbf{u})) \cdot \mathbf{e}_i(\mathbf{u}) = \frac{1}{h_i(\mathbf{u})} \frac{\partial \hat{f}}{\partial u_i}(\mathbf{u}) \]  
(543)

102
Substituting this into the expression for $\text{grad } \hat{f}$ we find that
\[
(\text{grad } \hat{f})(\mathbf{u}) = \sum_{i=1}^{3} \frac{1}{h_i(\mathbf{u})} \frac{\partial \hat{f}}{\partial u_i}(\mathbf{u}) \mathbf{e}_i(\mathbf{u})
\] (544)

**examples:** In cylindrical coordinates
\[
\text{grad } \hat{f} = \frac{1}{h_r} \frac{\partial \hat{f}}{\partial r} \mathbf{e}_r + \frac{1}{h_\theta} \frac{\partial \hat{f}}{\partial \theta} \mathbf{e}_\theta + \frac{1}{h_z} \frac{\partial \hat{f}}{\partial z} \mathbf{e}_z
\]
\[
= \frac{\partial \hat{f}}{\partial r} \mathbf{e}_r + \frac{1}{r} \frac{\partial \hat{f}}{\partial \theta} \mathbf{e}_\theta + \frac{\partial \hat{f}}{\partial z} \mathbf{e}_z
\] (545)

In spherical coordinates
\[
\text{grad } \hat{f} = \frac{1}{h_\rho} \frac{\partial \hat{f}}{\partial \rho} \mathbf{e}_\rho + \frac{1}{h_\phi} \frac{\partial \hat{f}}{\partial \phi} \mathbf{e}_\phi + \frac{1}{h_\theta} \frac{\partial \hat{f}}{\partial \theta} \mathbf{e}_\theta
\]
\[
= \frac{\partial \hat{f}}{\partial \rho} \mathbf{e}_\rho + \frac{1}{\rho} \frac{\partial \hat{f}}{\partial \phi} \mathbf{e}_\phi + \frac{1}{\rho \sin \phi} \frac{\partial \hat{f}}{\partial \theta} \mathbf{e}_\theta
\] (546)

**example**. Consider the function
\[
f(\mathbf{R}) = \frac{1}{|\mathbf{R}|} = \frac{1}{\sqrt{x^2 + y^2 + z^2}}
\] (547)

We want to find the gradient. It is easiest to change to spherical coordinates by
\[
\hat{f}(\rho, \phi, \theta) = \frac{1}{\rho}
\] (548)

Then since this function does not depend on $\phi$ or $\theta$
\[
\text{grad } \hat{f} = \frac{\partial \hat{f}}{\partial \rho} \mathbf{e}_\rho = -\frac{1}{\rho^2} \mathbf{e}_\rho
\] (549)

(And since $\rho = |\mathbf{R}|$ and $\mathbf{e}_\rho = \mathbf{R}/|\mathbf{R}|$ this is $-\mathbf{R}/|\mathbf{R}|^3$ back in Cartesian coordinates.)

Next consider the divergence. Recall that a vector field $\mathbf{v}(\mathbf{R})$ is expressed in new coordinates $\mathbf{R}(\mathbf{u})$ by $\mathbf{\hat{v}}(\mathbf{u}) = \mathbf{v}(\mathbf{R}(\mathbf{u}))$ and that in the new basis it has the form $\mathbf{\hat{v}} = \sum_i v_i \mathbf{e}_i$ with $v_i = \mathbf{\hat{v}} \cdot \mathbf{e}_i$. The divergence of the vector field is defined by $(\text{div } \mathbf{\hat{v}})(\mathbf{u}) = (\nabla \cdot \mathbf{v})(\mathbf{R}(\mathbf{u}))$. Then one can show that (after a somewhat lengthy calculation)
\[
\text{div } \mathbf{\hat{v}} = \frac{1}{h_1 h_2 h_3} \left( \frac{\partial (h_2 h_3 v_1)}{\partial u_1} + \frac{\partial (h_1 h_3 v_2)}{\partial u_2} + \frac{\partial (h_1 h_2 v_3)}{\partial u_3} \right)
\] (550)
For example consider spherical coordinates with \( h_\rho = 1, h_\phi = \rho, h_\theta = \rho \sin \phi \). Then we have
\[
\hat{v} = v_\rho e_\rho + v_\phi e_\phi + v_\theta e_\theta
\]
and
\[
\text{div } \hat{v} = \frac{1}{\rho^2 \sin \phi} \left( \frac{\partial (\rho^2 \sin \phi v_\rho)}{\partial \rho} + \frac{\partial (\rho \sin \phi v_\phi)}{\partial \phi} + \frac{\partial (\rho v_\theta)}{\partial \theta} \right)
\]
\[
= \frac{1}{\rho^2} \frac{\partial (\rho^2 v_\rho)}{\partial \rho} + \frac{1}{\rho \sin \phi} \left( \frac{\partial (\rho \sin \phi v_\phi)}{\partial \phi} + \frac{\partial \rho v_\theta}{\partial \theta} \right)
\]
(552)

For a general coordinate system the curl is \( (\text{curl } \hat{v})(\mathbf{u}) = (\nabla \times \mathbf{v})(R(\mathbf{u})) \). One can show that
\[
\text{curl } \hat{v} = \frac{1}{h_1 h_2 h_3} \det \begin{pmatrix}
h_1 e_1 & h_2 e_2 & h_2 e_3 \\
\partial/\partial u_1 & \partial/\partial u_2 & \partial/\partial u_3 \\
h_1 v_1 & h_2 v_2 & h_3 v_3
\end{pmatrix}
\]
(553)

Finally for a scalar if \( \hat{f}(\mathbf{u}) = f(R(\mathbf{u})) \) and \( (\hat{\Delta} \hat{f})(\mathbf{u}) = (\Delta f)(R(\mathbf{u})) \) then one can show
\[
\hat{\Delta} \hat{f} = \frac{1}{h_1 h_2 h_3} \left( \frac{\partial}{\partial u_1} \left( h_2 h_3 \frac{\partial \hat{f}}{\partial u_1} \right) + \frac{\partial}{\partial u_2} \left( h_1 h_3 \frac{\partial \hat{f}}{\partial u_2} \right) + \frac{\partial}{\partial u_3} \left( h_1 h_2 \frac{\partial \hat{f}}{\partial u_3} \right) \right)
\]
(554)

For example in spherical coordinates
\[
\hat{\Delta} \hat{f} = \frac{1}{\rho^2 \sin \phi} \left( \frac{\partial}{\partial \rho} \left( \rho^2 \sin \phi \frac{\partial \hat{f}}{\partial \rho} \right) + \frac{\partial}{\partial \phi} \left( \sin \phi \frac{\partial \hat{f}}{\partial \phi} \right) + \frac{\partial}{\partial \theta} \left( \frac{1}{\sin \phi} \frac{\partial \hat{f}}{\partial \theta} \right) \right)
\]
(555)

**example** Find all spherically symmetric solutions of \( \Delta f = 0 \). Spherically symmetric means that \( f(R) \) depends only on \( |R| \). In spherical coordinates it means that \( \hat{f}(\rho, \phi, \theta) = \hat{f}(\rho) \). Then the equation is
\[
\hat{\Delta} \hat{f} = \frac{1}{\rho^2} \frac{\partial}{\partial \rho} \left( \rho^2 \frac{\partial \hat{f}}{\partial \rho} \right) = 0
\]
(556)

This implies
\[
\rho^2 \frac{\partial \hat{f}}{\partial \rho} = c_1
\]
(557)
for some constant \( c_1 \). Then
\[
\frac{\partial \hat{f}}{\partial \rho} = \frac{c_1}{\rho^2}
\]
(558)

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which has the general solution

\[ \hat{f}(\rho) = \frac{-c_1}{\rho} + c_2 \]  \hspace{1cm} (559)

Back in Cartesian coordinates it says that

\[ f(\mathbf{R}) = \frac{-c_1}{|\mathbf{R}|} + c_2 \]  \hspace{1cm} (560)
3 \hspace{4pt} \textbf{complex variables}

3.1 \hspace{4pt} \textbf{complex numbers}

A complex number is a pair of real numbers, hence a vector in $\mathbb{R}^2$. It is written

$$z = (x, y)$$

(561)

Define addition and multiplication by a scalar just as for vectors:

$$z_1 + z_2 = (x_1 + x_2, y_1 + y_2)$$
$$\alpha z = (\alpha x, \alpha y) \quad \alpha \in \mathbb{R}$$

(562)

We also define multiplication by

$$z_1 z_2 = (x_1 x_2 - y_1 y_2, x_1 y_2 + y_1 x_2)$$

(563)

Then, as one can check, the ordinary rules of arithmetic apply:

$$z_1 (z_2 z_3) = (z_1 z_2) z_3$$
$$z_1 z_2 = z_2 z_1$$
$$z_1 (z_2 + z_3) = z_1 z_2 + z_1 z_3$$

Two special complex numbers are

$$1 = (1, 0) \quad i = (0, 1)$$

(565)

Then any complex $z$ can be written

$$z = (x, y) = x(1, 0) + y(0, 1) = x1 + yi$$

(566)

Consider complex numbers of the form $x1$. We have

$$(x_11)(x_21) = (x_1, 0)(x_2, 0) = (x_1 x_2, 0) = x_1 x_2 1$$
$$x_1 1 + x_2 1 = (x_1, 0) + (x_2, 0) = (x_1 + x_2, 0) = (x_1 + x_2) 1$$
$$(x1)(x', y') = (x, 0)(x', y') = (xx', xy') = x(x', y')$$

(567)

These behave just like the real numbers. Hence we can identify the complex number $x1 = (x, 0)$ with the real number $x$

Now consider complex numbers of the form $yi$. We have

$$(y_1 i)(y_2 i) = (0, y_1)(0, y_2) = (-y_1 y_2, 0) = -y_1 y_2 1 = -y_1 y_2$$

(568)

In particular

$$(yi)^2 = -y^2 \quad i^2 = -1$$

(569)
These complex numbers have a square which is negative. They are called imaginary
numbers.

A general complex number \( z = (x, y) = x + yi \) can now written

\[
z = x + iy
\]

Points in the plane can be labeled in this form. For example the point \( (3, 2) \) could be
labeled \( 3 + 2i \). Also in this form the multiplication law need not be remembered since
it follows from the relation \( i^2 = -1 \). Indeed we have

\[
z_1z_2 = (x_1 + iy_1)(x_2 + iy_2) \\
= x_1x_2 + ix_1y_2 + iy_1x_2 + i^2y_1y_2 \\
= x_1x_2 - y_1y_2 + i(x_1y_2 + y_1x_2)
\]

### 3.2 definitions and properties

1. (definitions) For a complex number \( z = x + iy \) define

\[
x = \text{Re} \, z = \text{real part of } z \\
y = \text{Im} \, z = \text{imaginary part of } z
\]

We also define

\[
|z| = \sqrt{x^2 + y^2}
\]

called the ”length of \( z \)” or the ”absolute value of \( z \)” or the ”modulus of \( z \)”, and

\[
\arg z = \text{polar angle for } z
\]

and

\[
\bar{z} = x - iy
\]

called the ”complex conjugate” of \( z \). See figure 30 for the associated geometric picture.

The complex conjugate has the properties

\[
\overline{z_1 + z_2} = \overline{z}_1 + \overline{z}_2 \\
\overline{z_1z_2} = \overline{z}_1\overline{z}_2 \\
|\overline{z}| = |z|
\]

2. (inverses) Note that

\[
z\overline{z} = (x + iy)(x - iy) = x^2 + y^2 = |z|^2
\]

If \( z \neq 0 \) this can be written

\[
z(\overline{z}/|z|^2) = 1
\]
This says that $\overline{z}/|z|^2$ is an inverse for $z$. We write
\[ z^{-1} = \frac{\overline{z}}{|z|^2} \quad \text{or} \quad (x + iy)^{-1} = \frac{x - iy}{x^2 + y^2} \] (579)

Then we can divide by any $z \neq 0$ by defining
\[ \frac{w}{z} = wz^{-1} \] (580)

**Examples:** By the formula $i^{-1} = -i$ and
\[ (3 + 4i)^{-1} = \frac{3 - 4i}{25} = \frac{3}{25} - \frac{4}{25}i \] (581)

However, one does not have to remember the formula. Instead, just multiply in the numerator and denominator by the complex conjugate. For example
\[ \frac{1}{3 + 4i} = \frac{1}{3 + 4i} \cdot \frac{3 - 4i}{3 - 4i} = \frac{3 - 4i}{25} \]

3. (Distances) Just as for vectors
\[ |z_1 - z_2| = \text{distance from } z_1 \text{ to } z_2 \] (582)
The triangle inequality holds:

\[ |z_1 - z_2| \leq |z_1 \pm z_2| \leq |z_1| + |z_2| \] (582)

4. (exponentials) We want to define

\[ e^z = e^{x+iy} = e^x e^{iy} \] (582)

We know what \( e^x \) means and we know that it can be expressed as a convergent series

\[ e^x = 1 + x + \frac{1}{2!}x^2 + \frac{1}{3!}x^3 + \cdots \] (582)

Suppose we try to define \( e^{iy} \) in the same way ignoring questions of convergence. We would have

\[ e^{iy} = 1 + iy + \frac{1}{2!}(iy)^2 + \frac{1}{3!}(iy)^3 + \cdots = \left(1 - \frac{1}{2!}y^2 + \frac{1}{4!}y^4 + \cdots \right) + i \left(y - \frac{1}{3!}y^3 + \frac{1}{5!}y^5 + \cdots \right) = \cos y + i \sin y \]

We take this as the definition, that is \( e^{iy} \equiv \cos y + i \sin y \) and more generally

\[ e^z = e^{x+iy} = e^x (\cos y + i \sin y) \] (582)

This does obey the law of exponents:

\[ e^{iy_1} e^{iy_2} = (\cos y_1 + i \sin y_1)(\cos y_2 + i \sin y_2) \]
\[ = (\cos y_1 \cos y_2 - \sin y_1 \sin y_2) + i(\cos y_1 \sin y_2 + \sin y_1 \cos y_2) \]
\[ = \cos(y_1 + y_2) + i \sin(y_1 + y_2) = e^{i(y_1 + y_2)} \] (582)

and in general \( e^{z_1} e^{z_2} = e^{z_1 + z_2} \). Also note

\[ e^{iy} = \cos y + i \sin y = \cos (-y) + i \sin (-y) = e^{-iy} \]
\[ |e^{iy}| = \cos^2 y + \sin^2 y = 1 \]
\[ (e^{iy})^{-1} = \frac{e^{iy}}{|e^{iy}|^2} = e^{-iy} \] (582)

The last could also be deduced from the law of exponents.

examples:

\[ e^{i0} = 1 \quad e^{i\pi/2} = i \quad e^{i\pi} = -1 \quad e^{i3\pi/2} = -i \quad e^{i2\pi} = 1 \] (582)

\[ e^{3 + i\pi/4} = e^3 (\cos \pi/4 + i \sin \pi/4) = e^3 \left( \frac{1}{\sqrt{2}} + i \frac{1}{\sqrt{2}} \right) \] (582)
3.3 polar form

Points in the plane have polar coordinates \((r, \theta)\) related to Cartesian coordinate by \(x = r \cos \theta, y = r \sin \theta\). A complex number can then be written in polar form by

\[
z = x + iy = r \cos \theta + ir \sin \theta = re^{i\theta}
\]

Every point \(z \neq 0\) has a unique polar representation \(z = re^{i\theta}\) with \(r > 0\) and \(0 \leq \theta < 2\pi\). Now we have

\[
|z| = r \quad \arg z = \theta
\]

If \(z_1 = r_1e^{i\theta_1}\) and \(z_2 = r_2e^{i\theta_2}\) then the product is

\[
z_1z_2 = r_1r_2e^{i(\theta_1 + \theta_2)}
\]

This says that in complex multiplication we multiply lengths and add angles. Another way to write it is

\[
|z_1z_2| = |z_1||z_2| \quad \arg(z_1z_2) = \arg(z_1) + \arg(z_2)
\]

The second equation only holds if both sides are in the interval \([0, 2\pi)\). For example if \(z_1 = z_2 = -i\), then \(\arg(z_1) + \arg(z_2) = 3\pi/2 + 3\pi/2 = 3\pi\) but \(\arg(z_1z_2) = \arg(-1) = \pi\) and the identity fails.

**examples**

1. To find the polar form of \(z = 1 + i\sqrt{3}\) note that the length is 2 and the angle is \(\tan^{-1}(\sqrt{3}) = \pi/3\) hence \(z = 2e^{i\pi/3}\).

2. To find the polar form of \(z = 1 + i\) note that the length is \(\sqrt{2}\) and the angle is \(\tan^{-1}(1) = \pi/4\). Thus \(z = \sqrt{2}e^{i\pi/4}\).

3. To find \((1 + i)^8\) write it in polar form and compute

\[
(1 + i)^8 = (\sqrt{2}e^{i\pi/4})^8 = 2^4e^{2\pi i} = 16
\]

**problem:** Find all complex \(z\) such that \(z^3 = 1\)

**solution:** Try \(z = re^{i\theta}\) with \(r > 0\) and \(0 \leq \theta < 2\pi\). This is a solution if

\[
r^3e^{i3\theta} = 1
\]

Comparing lengths gives \(r^3 = 1\) and so \(r = 1\). Then \(\theta\) must satisfy \(e^{i3\theta} = 1\) or \(\cos(3\theta) = 1, \sin(3\theta) = 0\). The solutions are

\[
3\theta = 0, \pm2\pi, \pm4\pi, \pm6\pi, \ldots
\]
or
\[ \theta = 0, \pm 2\pi/3, \pm 4\pi/3, \pm 6\pi/3, \ldots \]  \hspace{1cm} (582)

But the only solutions with \(0 \leq \theta < 2\pi\) are \(\theta = 0, 2\pi/3, 4\pi/3\). Thus the answer is
\[
z = e^{i0}, \ e^{i2\pi/3}, \ e^{i4\pi/3} = 1, \ -\frac{1}{2} + \frac{\sqrt{3}}{2}i, \ -\frac{1}{2} - \frac{\sqrt{3}}{2}i \hspace{1cm} (582)
\]

One of the uses of complex numbers is that every polynomial has complex roots. In the problem we found the solutions of \(z^3 - 1 = 0\). Here are some more examples:

**examples:**

1. The equation \(z^n = a\) (a real) has \(n\) solutions:
\[
z = a^{1/n}, \ a^{1/n}e^{i(2\pi/n)}, \ a^{1/n}e^{i2(2\pi/n)}, \ldots, a^{1/n}e^{i(n-1)(2\pi/n)} \hspace{1cm} (582)
\]

2. The equation \(az^2 + bz + c = 0\) (a, b, c real) has solutions
\[
z = \begin{cases} 
-\frac{b \pm \sqrt{b^2 - 4ac}}{2a} & \text{if } b^2 - 4ac > 0 \\
-\frac{b}{2a} & \text{if } b^2 - 4ac = 0 \\
-\frac{b \pm i\sqrt{4ac - b^2}}{2a} & \text{if } b^2 - 4ac < 0 
\end{cases} \hspace{1cm} (582)
\]

### 3.4 functions

Let \(\mathbb{C}\) stand for the set of all complex numbers. Thus \(\mathbb{C}\) is \(\mathbb{R}^2\) with a special multiplication. We are interested in functions from \(\mathbb{C}\) (or a subset) to \(\mathbb{C}\) written \(w = f(z)\) with both \(w, z\) complex. For example \(w = z^2\) or \(w = 1/(3z + 2)\) or \(w = |z|\).

If \(w = u + iv\) and
\[
f(z) = f(x + iy) = u(x, y) + iv(x, y) \hspace{1cm} (582)
\]

then the equation \(w = f(z)\) can be written as the pair of equations
\[
u = u(x, y) \hspace{1cm} v = v(x, y) \hspace{1cm} (582)
\]

The function \(u(x, y)\) is called the real part and the function \(v(x, y)\) is called the imaginary part.

**examples**

1. If \(w = e^z = e^x (\cos y + i \sin y)\) then
\[
u = e^x \cos y \hspace{1cm} v = e^x \sin y \hspace{1cm} (582)
\]
To represent a function \( w = f(z) \) geometrically we show how it affects straight lines.

**Example:** Consider the function \( w = z^2 \) which takes the first quadrant to the upper half plane. We have \( w = z^2 = (x + iy)^2 = (x^2 - y^2) + i2xy \) so it is equivalent to the pair of functions

\[
\begin{align*}
    u &= x^2 - y^2 \\
    v &= 2xy
\end{align*}
\]  

(582)

It maps the line \( x = c, y > 0 \) to the line \( u = c^2 - y^2, v = 2cy \) which is the half parabola

\[
\begin{align*}
    u &= c^2 - \frac{v^2}{4c^2}, \quad v > 0
\end{align*}
\]  

(582)

It maps the line \( y = c, x > 0 \) to the line \( u = x^2 - c^2, v = 2xc \) which is the half parabola

\[
\begin{align*}
    u &= \frac{v^2}{4c^2} - c^2, \quad v > 0
\end{align*}
\]  

(582)

This is represented in figure 31.
3.5 special functions

(A.) trigonometric functions

If \( y \) is real so \( e^{iy} = \cos y + i \sin y \) and \( e^{-iy} = \cos y - i \sin y \) then

\[
\frac{1}{2}(e^{iy} + e^{-iy}) = \cos y \quad \frac{1}{2}(e^{iy} - e^{-iy}) = i \sin y
\]

Accordingly we define for complex \( z \)

\[
\cos z = \frac{1}{2}(e^{iz} + e^{-iz}) \quad \sin z = \frac{1}{2i}(e^{iz} - e^{-iz})
\]

Other trig functions can be defined from these, for example

\[
\tan z = \frac{\sin z}{\cos z}
\]

If \( z = i\alpha \) is purely imaginary then

\[
\cos(i\alpha) = \frac{1}{2}(e^{-\alpha} + e^{\alpha}) = \cosh \alpha
\]

\[
\sin(i\alpha) = \frac{1}{2i}(e^{-\alpha} - e^{\alpha}) = i\frac{1}{2}(e^{\alpha} - e^{-\alpha}) = i \sinh \alpha
\]

Thus the complex trig functions include both the usual trig functions and the hyperbolic trig functions as special cases.

(B.) logarithm

We want to define the natural logarithm \( \log z \) of a complex number \( z \). If \( z \neq 0 \) and such a log has the same properties as the log of positive numbers then we expect

\[
\log z = \log re^{i\theta} = \log r + \log e^{i\theta} = \log r + i\theta
\]

Accordingly we take as the definition for \( z \neq 0 \)

\[
\log z = \log |z| + i \arg z
\]

where \( 0 \leq \arg z < 2\pi \). However sometimes we may want to make another choice of the polar angle \( \arg z \). For example we could take \(-\pi \leq \arg z < \pi \) or more generally

\[
c \leq \arg z < c + 2\pi
\]

Each choice of \( c \) gives a different log functions, called a \textit{branch of the logarithm}. An alternative is to take all possible values of \( \arg z \). If \( \theta \) is one possible value then these would be the infinite sequence

\[
\ldots, \theta - 2\pi, \theta, \theta + 2\pi, \theta + 4\pi, \ldots
\]
In this case we have a multi-valued function.

**Example:** Take the branch of \( \log z \) with \( 0 \leq \arg z < 2\pi \). Then

\[
\log i = \log |i| + i \arg i = i \frac{\pi}{2}
\]
\[
\log(-2) = \log |-2| + i \arg(-2) = \log 2 + i\pi
\]
\[
\log(-3i) = \log |-3i| + i \arg(-3i) = \log 3 + i\frac{3\pi}{2}
\]
\[
\log(1 + i) = \log |1 + i| + i \arg(1 + i) = \log \sqrt{2} + i\frac{\pi}{4}
\]

Is the logarithm the inverse of the exponential? We have the following result

**Theorem 23**

1. For any branch of the logarithm

\[
e^{\log z} = z
\]

2. If \( \log z \) is defined with \( c \leq \arg z < c + 2\pi \) and if \( z = x + iy \) with \( c \leq y < c + 2\pi \) then

\[
\log(e^z) = z
\]

**Proof.** For the first if \( z = r e^{i\theta} \), then \( \log z = \log r + i(\theta + 2\pi k) \) for some \( k \). Then

\[
e^{\log z} = e^{\log r e^{i\theta} e^{i2\pi k}} = r e^{i\theta} = z
\]

For the second since the restriction on \( y \) matches the definition of \( \arg \)

\[
\log(e^z) = \log(e^x e^{iy}) = \log e^x + i \arg e^{iy} = x + iy = z
\]

(C.) **complex exponentials**

Since \( z = e^{\log z} \) we can define for complex \( z, w \)

\[
z^w = e^{w \log z}
\]

Here we have to specify the branch of the logarithm we are using. An alternative is to take all possible values of \( \log z \) and get a multi-valued function \( z^w \).

If \( w \) is a positive integer \( n \) then

\[
z^n \text{ (new)} = e^{n \log z} = (e^{\log z})^n = z^n \text{ (old)}
\]
is the same for any branch of the logarithm. But
\[ z^{1/n} = e^{\log z / n} \] (582)
will have \( n \) different values depending on the branch of the log. Each value does satisfy
\[(z^{1/n})^n = e^{\log z} = z.\]

**Example** Since \( \log 1 = i \arg 1 \)
\[1^{1/3} = e^{\frac{1}{3} \log 1} = e^{\frac{1}{3} \arg 1} = \begin{cases} 1 & \text{if } \arg 1 = 0 \\ e^{i2\pi/3} & \text{if } \arg 1 = 2\pi \\ e^{i4\pi/3} & \text{if } \arg 1 = 4\pi \end{cases} (582)\]
We find the solutions of \( z^3 = 1 \) as \( z = 1^{1/3} \) with the different branches for the cube root.

**Example** Since \( \log i = i \arg i \)
\[i^{1/3} = e^{\frac{1}{3} \log i} = e^{\frac{1}{3} \arg i} = \begin{cases} e^{i\pi/6} & \text{if } \arg i = \frac{\pi}{2} \\ e^{i5\pi/6} & \text{if } \arg i = \frac{5\pi}{2} \\ e^{i9\pi/6} & \text{if } \arg i = \frac{9\pi}{2} \end{cases} (582)\]

**Example** Since \( \log i = i \arg i \)
\[i^i = e^{i \log i} = e^{-\arg i} (582)\]
This takes infinitely many values
\[\ldots, e^{3\pi/2}, e^{-\pi/2}, e^{-5\pi/2}, e^{-9\pi/2}, \ldots (582)\]

### 3.6 derivatives

(A.) First consider functions from \( \mathbb{R} \) to \( \mathbb{C} \) written \( z = f(t) \) where \( z \) is complex and \( t \) is real. For example \( z = e^{it} \) or \( z = t^2 + i \log t \). Such a function can always be written
\[ z = f(t) = x(t) + iy(t) \] (582)
and \( x(t) \) is called the real part of the function and \( y(t) \) is called the imaginary part of the function. The derivative is defined just as for a vector valued function:
\[ \frac{dz}{dt} = f'(t) = x'(t) + iy'(t) \] (582)
example: If
\[ z = e^{iat} = \cos(at) + i \sin(at) \quad a \text{ real} \] then
\[ \frac{dz}{dt} = -a \sin(at) + ia \cos(at) = ia\cos(at) + i \sin(at) = iae^{iat} \] (582)

(B.) Now consider functions from $\mathbb{C}$ to $\mathbb{C}$ written $w = f(z)$. We start with the definition of a limit.

**definition:** $\lim_{z \to z_0} f(z) = w_0$ means for every $\epsilon > 0$ there is a $\delta > 0$ so if $|z - z_0| < \delta$ then $|f(z) - w_0| < \epsilon$.

This says that $f(z) \to w_0$ as $z \to z_0$ from any direction.

Now write
\[ z = x + iy \quad z_0 = x_0 + iy_0 \]
\[ f(z) = u(x, y) + iv(x, y) \quad w_0 = u_0 + iv_0 \] (582)

Then
\[ |f(z) - w_0| = \sqrt{(u(x, y) - u_0)^2 + (v(x, y) - v_0)^2} \]
\[ |z - z_0| = \sqrt{(x - x_0)^2 + (y - y_0)^2} \] (582)

From this we deduce that $\lim_{z \to z_0} f(z) = w_0$ is the same as
\[ \lim_{(x,y) \to (x_0, y_0)} u(x, y) = u_0 \]
\[ \lim_{(x,y) \to (x_0, y_0)} v(x, y) = v_0 \] (582)

**definition:** $w = f(z)$ is continuous at $z_0$ if $\lim_{z \to z_0} f(z) = f(z_0)$

This is the same as the statement that $u(x, y)$ and $v(x, y)$ are continuous at $(x_0, y_0)$. One can show that if $f, g$ are continuous at $z_0$ then so are $f \pm g, f \cdot g$, and $f/g$, the last provided $g(z_0) \neq 0$. Also if $g$ is continuous at $z_0$ and $f$ is continuous at $g(z_0)$ then the composition $f \circ g$ is continuous at $z_0$.

**definition:** $w = f(z)$ is differentiable at $z_0$ if the derivative
\[ f'(z_0) = \lim_{\Delta z \to 0} \frac{f(z_0 + \Delta z) - f(z_0)}{\Delta z} \quad \text{exists} \] (582)

Varying $z_0$ it is also a function. We also write
\[ \frac{dw}{dz} = f'(z) \] (582)
Theorem 24  If \( f \) is differentiable at \( z_0 \), then it is continuous at \( z_0 \)

Proof.

\[
\lim_{z \to z_0} f(z) - f(z_0) = \lim_{z \to z_0} \left( f(z) - f(z_0) \right) \left( \frac{z - z_0}{z - z_0} \right) = f'(z_0) \cdot 0 = 0 \tag{582}
\]

Theorem 25  If \( f, g \) are differentiable at \( z \) so are \( f \pm g, f \cdot g, f/g \) (provided \( g(z) \neq 0 \)) and

\[
(f \pm g)'(z) = f'(z) \pm g'(z)
\]

\[
(f \cdot g)'(z) = f'(z)g(z) + f(z)g'(z)
\]

\[
\left( \frac{f'}{g} \right)'(z) = \frac{g(z)f'(z) - f(z)g'(z)}{g(z)^2} \tag{582}
\]

If \( g \) is differentiable at \( z \) and \( f \) is differentiable at \( g(z) \) then \( f \circ g \) is differentiable at \( z \) and

\[
(f \circ g)'(z) = f'(g(z))g'(z) \tag{582}
\]

These are proved just as for real variables.

examples:

1. If \( f(z) = c \) then \( f'(z) = 0 \).

2. If \( f(z) = z \) then

\[
f'(z) = \lim_{\Delta z \to 0} \frac{(z + \Delta z) - z}{\Delta z} = \lim_{\Delta z \to 0} 1 = 1 \tag{582}
\]

3. If \( f(z) = z^2 \) then by the product rule

\[
f'(z) = \frac{dz}{dz} z + z \frac{dz}{dz} = 2z \tag{582}
\]

4. If \( f(z) = z^n \) then \( f'(z) = nz^{n-1} \).

5. Any polynomial \( P(z) \) is differentiable.

6. Any rational function \( P(z)/Q(z) \) is differentiable except at points where \( Q(z) = 0 \).
   (where it is not even defined)
Let \( f(z) = \bar{z} \). If it exists the derivative is
\[
\lim_{\Delta z \to 0} \frac{f(z + \Delta z) - f(z)}{\Delta z} = \lim_{\Delta z \to 0} \frac{\Delta \bar{z}}{\Delta z}
\] (582)
The limit must be the same from any direction. But if \( \Delta z = \Delta x \) is real then
\[
\lim_{\Delta z \to 0} \frac{\Delta \bar{z}}{\Delta z} = \lim_{\Delta x \to 0} \frac{\Delta x}{\Delta z} = 1
\] (582)
and \( \Delta z = i\Delta y \) is imaginary then
\[
\lim_{\Delta z \to 0} \frac{\Delta \bar{z}}{\Delta z} = \lim_{\Delta y \to 0} \frac{-i\Delta y}{\Delta y} = -1
\] (582)
The limit depends on the direction which means there is no limit in the complex sense. Thus the derivative does not exist. (Even though there are no kinks or discontinuities in this function).

### 3.7 Cauchy-Riemann equations

Let \( f(z) \) be differentiable so \( \lim_{\Delta z \to 0} (f(z + \Delta z) - f(z)) / \Delta z \) exists. What does this say about the real and imaginary parts \( f(z) = u(x, y) + iv(x, y) \)? First let \( \Delta z = \Delta x \) be real. Then
\[
f'(z) = \lim_{\Delta x \to 0} \frac{f(z + \Delta x) - f(z)}{\Delta x}
= \lim_{\Delta x \to 0} \left( \frac{u(x + \Delta x, y) - u(x, y)}{\Delta x} \right) + i \lim_{\Delta x \to 0} \left( \frac{v(x + \Delta x, y) - v(x, y)}{\Delta x} \right)
= \frac{\partial u}{\partial x} + i \frac{\partial v}{\partial x}
\] (583)
Now let \( \Delta z = i\Delta y \) be imaginary. Then
\[
f'(z) = \lim_{\Delta y \to 0} \frac{f(z + i\Delta y) - f(z)}{i\Delta y}
= \frac{1}{i} \lim_{\Delta y \to 0} \left( \frac{u(x, y + \Delta y) - u(x, y)}{\Delta y} \right) + \lim_{\Delta y \to 0} \left( \frac{v(x, y + \Delta y) - v(x, y)}{\Delta y} \right)
= -i \frac{\partial u}{\partial y} + \frac{\partial v}{\partial y}
\] (584)
These two expressions must agree since the limit must be the same from any direction. Thus we must have
\[
\frac{\partial u}{\partial x} = \frac{\partial v}{\partial y} \quad \frac{\partial u}{\partial y} = -\frac{\partial v}{\partial x}
\] (585)
These are the *Cauchy-Riemann equations* or CR equations.

Thus if \( f(z) \) is differentiable the the real and imaginary parts satisfy the CR equations. The converse is also true if we strengthen the hypotheses a bit.
**Theorem 26** If \( u(x,y), v(x,y) \) have continuous partial derivatives which satisfy the CR equations near a point, then \( f(z) = u(x,y) + iv(x,y) \) is differentiable at that point and

\[
f'(z) = \frac{\partial u}{\partial x} + i \frac{\partial v}{\partial x} = \frac{\partial v}{\partial y} - i \frac{\partial u}{\partial y}
\]

**examples**

1. Let \( f(z) = z^2 \) so that \( u = x^2 - y^2 \) and \( v = 2xy \). We have

\[
\frac{\partial u}{\partial x} = 2x = \frac{\partial v}{\partial y} \quad \frac{\partial u}{\partial y} = -2y = -\frac{\partial v}{\partial x}
\]

Thus the CR equations hold so \( f(z) \) is differentiable (which we already knew) and

\[
\frac{d(z^2)}{dz} = \frac{\partial u}{\partial x} + i \frac{\partial v}{\partial x} = 2x + i2y = 2z
\]

(which we also knew)

2. Let \( f(z) = e^z = e^x \cos y + ie^x \sin y \) so \( u = e^x \cos y, v = e^x \sin y \). Then

\[
\frac{\partial u}{\partial x} = e^x \cos y = \frac{\partial v}{\partial y} \quad \frac{\partial u}{\partial y} = -e^x \sin y = -\frac{\partial v}{\partial x}
\]

Thus the CR equations hold so \( f(z) \) is differentiable and

\[
\frac{d(e^z)}{dz} = \frac{\partial u}{\partial x} + i \frac{\partial v}{\partial x} = e^x \cos y + ie^x \sin y = e^z
\]

We can use this result and the chain rule to find derivatives of trigonometric functions

\[
\frac{d(\cos z)}{dz} = \frac{d}{dz} \left[ \frac{e^{iz} + e^{-iz}}{2} \right] = \frac{ie^{iz} - ie^{-iz}}{2} = -\frac{e^{iz} - e^{-iz}}{2i} = -\sin z
\]

Similarly

\[
\frac{d(\sin z)}{dz} = \cos z
\]

3. Let \( f(z) = \bar{z} = x - iy = u + iv \). Then

\[
\frac{\partial u}{\partial x} = 1 \neq -1 = \frac{\partial v}{\partial y}
\]

Thus the CR equations fail everywhere and so the function is not differentiable anywhere (which we already knew).
4. Let \( f(z) = \log z = \log |z| + i \arg z \). Consider only the right half plane with \(-\frac{1}{2}\pi < \arg z < \frac{1}{2}\pi\). Then \( \arg z = \tan^{-1}(y/x) \) so we have

\[
\begin{align*}
  u &= \log \sqrt{x^2 + y^2} = \frac{1}{2} \log(x^2 + y^2) \\
  v &= \tan^{-1}(y/x)
\end{align*}
\]

Then we compute

\[
\frac{\partial u}{\partial x} = \frac{x}{x^2 + y^2} = \frac{\partial v}{\partial y} \quad \frac{\partial u}{\partial y} = \frac{y}{x^2 + y^2} = -\frac{\partial v}{\partial x}
\]

(595)

Thus the CR equations hold so \( f(z) \) is differentiable and

\[
\frac{d(\log z)}{dz} = \frac{\partial u}{\partial x} + i \frac{\partial v}{\partial x} = \frac{x - iy}{x^2 + y^2} = \frac{\bar{z}}{|z|^2} = z^{-1}
\]

(596)

This result actually holds for any branch of the logarithm.

3.8 analyticity

A **neighborhood** of a point \( z_0 \) is a disc of some radius \( \epsilon \) centered on \( z_0 \). It is written \( \{ z \in \mathbb{C} : |z - z_0| < \epsilon \} \). A set \( D \subset \mathbb{C} \) is defined to be **open** if every point \( z_0 \in D \) has a neighborhood entirely contained in \( D \). For example the disc \( \{ z \in \mathbb{C} : |z - z_0| < r \} \) is open, but the disc \( \{ z \in \mathbb{C} : |z - z_0| \leq r \} \) is not open since any neighborhood of a point with \( |z - z_0| = r \) will have points outside the disc.

A function \( w = f(z) \) is **analytic** on an open set \( D \) if it is defined and differentiable at every point in \( D \). (We restrict to open sets so that if \( z_0 \in D \) then \( z_0 + \Delta z \in D \) for \( \Delta z \) small enough, hence we can form the difference quotient \( (f(z_0 + \Delta z) - f(z_0))/\Delta z \) and test the differentiability.)

Thus analyticity is essentially the same as differentiability, expect that the latter refers to points and the former refers to regions. When we take up integration theory it will be important to identify domains of analyticity for various functions. The following examples give some practice at this.

**Examples:**

1. A polynomial \( P(z) \) is analytic in the entire plane \( \mathbb{C} \).
2. The functions \( e^z, \cos z, \sin z \) are analytic in plane.
3. The function \( 1/(z^2 + 1) \) is analytic in the plane with \( z = \pm i \) deleted. (It is not defined at the deleted points.)
4. The function \( 1/(z^3 - 1) \) is analytic in the plane with \( z = 1, e^{2\pi i/3}, e^{4\pi i/3} \) deleted.
5. A rational function \( P(z)/Q(z) \) is analytic in the plane with the roots of \( Q(z) \) deleted.
6. Consider \( \log z = \log |z| + i \arg z \) with \( z \neq 0 \) and \(-\pi < \arg z < \pi\). This is analytic in the cut plane with the negative real axis deleted. We could define it on the whole plane with only \( z = 0 \) deleted if we specified say \(-\pi \leq \arg z < \pi\). But then it would not be analytic because it would be discontinuous across the negative real axis and hence not differentiable on the negative real axis.

3.9 complex line integrals

(A.) First consider functions \( f \) from (a subset of) \( \mathbb{R} \) to \( \mathbb{C} \) written \( z = f(t) = x(t) + iy(t) \). We define

\[
\int_a^b f(t) dt = \int_a^b x(t) dt + i \int_a^b y(t) dt
\]

(597)

Then we have

\[
\int_a^b f'(t) dt = \int_a^b x'(t) dt + i \int_a^b y'(t) dt
\]

\[
= [x(t)]_a^b + i[y(t)]_a^b
\]

(598)

\[
= [f(t)]_a^b
\]

example:

\[
\int_0^{2\pi} e^{iat} dt = \left[ \frac{e^{iat}}{ia} \right]_0^{2\pi} = \frac{e^{2\pi ia} - 1}{ia}
\]

(599)

(B.) Next consider functions \( f \) from (a subset of ) \( \mathbb{C} \) to \( \mathbb{C} \) written \( w = f(z) \). We could consider integrals over regions in the plane, but the most interesting case turns out to be line integrals.

Let \( \mathcal{C} \) be a directed curve in the plane. Choose a sequence of points \( z_0, z_1, \ldots, z_n \) such that \( z_0 \) is the start point and \( z_n \) is the finishing point. (see figure 32) Let \( z_i^* \) be a point on \( \mathcal{C} \) between \( z_i \) and \( z_{i+1} \). Finally let \( \Delta z_i = z_{i+1} - z_i \) and \( h = \max_i |\Delta z_i| \). We define the line integral of \( f \) over \( \mathcal{C} \) by

\[
\int_{\mathcal{C}} f(z) dz = \lim_{h \to 0} \sum_{i=0}^{n-1} f(z_i^*) \Delta z_i
\]

(600)

Note that this involves complex multiplication, so it is different from the line integrals we have considered previously.

Now suppose that \( \mathcal{C} \) is parametrized by

\[
z(t) = x(t) + iy(t) \quad a \leq t \leq b
\]

(601)
Take points \( a = t_0 \leq t_1 \leq t_2 \leq \cdots \leq t_n = b \). Take \( z_i = z(t_i) \) and for any points \( t_i^* \) in \( [t_i, t_{i+1}] \) let \( z_i^* = z(t_i^*) \). If \( \Delta t_i = t_{i+1} - t_i \) then
\[
\Delta z_i = z(t_{i+1}) - z(t_i) \approx z'(t_i^*)\Delta t_i
\]
(602)

Hence
\[
\int_C f(z)dz = \lim_{h \to 0} \sum_{i=0}^{i-1} f(z(t_i^*))z'(t_i^*)\Delta t_i = \int_a^b f(z(t))z'(t)dt
\]
(603)

This is the basic definition. It can also be written
\[
\int_C f(z)dz = \int_a^b f(z(t))\frac{dz}{dt}dt
\]
(604)

and for short one can just remember \( dz = (dz/dt)dt \). This integral is independent of the choice of parametrization.

**examples**

1. We want to find \( \int_C zdz \) where \( C \) be the unit circle traversed counterclockwise.

The usual parametrization \( x = \cos t, y = \sin t \) with \( 0 \leq t \leq 2\pi \) becomes in the complex form
\[
z = x + iy = \cos t + i\sin t = e^{it} \quad 0 \leq t \leq 2\pi
\]
(605)
Then
\[ dz = ie^{it} dt \]  
and so
\[ \int_C zdz = \int_0^{2\pi} e^{it}ie^{it} dt = \int_0^{2\pi} ie^{2it} dt = \left[ \frac{e^{2it}}{2} \right]_0^{2\pi} = 0 \]  

2. With the same curve \( C \)
\[ \int_C \bar{z}dz = \int_0^{2\pi} e^{-it}ie^{it} dt = \int_0^{2\pi} idt = 2\pi i \]  

3. Let \( C = C_1 + C_2 \) where \( C_1 \) is a straight line from 0 to 1 and \( C_2 \) is a straight line from 1 to 1 + i. We evaluate
\[ \int_C \bar{z}dz = \int_{C_1} \bar{z}dz + \int_{C_2} \bar{z}dz \]  
The line \( C_1 \) is parametrized by \( z = t, \ 0 \leq t \leq 1. \) Hence \( dz = dt \) and \( \bar{z} = t \) and
\[ \int_{C_1} \bar{z}dz = \int_0^1 t dt = \frac{1}{2} \]  
The line \( C_2 \) is parametrized by \( z = 1 + it, \ 0 \leq t \leq 1. \) Hence \( dz = idt \) and \( \bar{z} = 1 - it \) and
\[ \int_{C_1} \bar{z}dz = \int_0^1 (1 - it) idt = \int_0^1 (i + t) dt = \left[ it + \frac{t^2}{2} \right]_0^1 = i + \frac{1}{2} \]  
Thus \( \int_C \bar{z}dz = 1 + i. \)

### 3.10 properties of line integrals

To estimate the size of a line integral we have:

**Theorem 27** Suppose a curve \( C \) has length \( L \) and \( |f(z)| \leq M \) for all \( z \) on \( C. \) Then
\[ \left| \int_C f(z)dz \right| \leq ML \]  

**Proof.** Approximate the integral by a sum \( \sum_i f(z_i^*) \Delta z_i. \) By the triangle inequality
\[ \left| \sum_i f(z_i^*) \Delta z_i \right| \leq \sum_i |f(z_i^*) \Delta z_i| \]
\[ = \sum_i |f(z_i^*)||\Delta z_i| \]
\[ \leq M \sum_i |\Delta z_i| \]  

123
Now take the limit as $h = \max_i |\Delta z_i| \to 0$ and get the result.

**example:** Let $C$ be a semi-circle of radius $R$ in the upper half plane centered on the origin. Suppose we want to estimate the size of $\int_C (z + 1)^3dz$. For $z$ on $C$ we have

$$|(z + 1)^3| = |z + 1|^3 \leq (|z| + 1)^3 = (R + 1)^3$$  \hfill (614)

The length of $C$ is $\pi R$. Hence by the theorem

$$\left|\int_C (z + 1)^3dz\right| \leq (R + 1)^3 \pi R$$  \hfill (615)

Next we want to express our complex line integrals in terms of their real and imaginary parts. Start with the definition

$$\int_C f(z)dz = \int_a^b f(z(t)) \frac{dz}{dt} dt$$  \hfill (616)

where $z(t), a \leq t \leq b$ is a parametrization of $C$. Then insert the expressions

$$f(z) = u(x, y) + iv(x, y)$$

$$z(t) = x(t) + iy(t)$$  \hfill (617)

This gives

\[
\int_C f(z)dz = \int_a^b \left( u(x(t), y(t)) + iv(x(t), y(t)) \right) \left( \frac{dx}{dt} + i \frac{dy}{dt} \right) dt \\
= \int_a^b \left( u(x(t), y(t)) \frac{dx}{dt} - v(x(t), y(t)) \frac{dy}{dt} \right) dt \\
+ i \int_a^b \left( u(x(t), y(t)) \frac{dy}{dt} + v(x(t), y(t)) \frac{dx}{dt} \right) dt \\
= \left( \int_C udx - vdy \right) + i \left( \int_C udy + vdx \right) 
\]  \hfill (618)

For short one can think of making the substitutions $f = u + iv$ and $dz = dx + idy$.

This formula expresses complex line integrals in terms of real line integrals. From this we can deduce that the complex line integrals have all the properties of real line integrals. In particular

$$\int_C f(z)dz = -\int_C f(z)dz$$  \hfill (619)
3.11 Cauchy’s theorem

Now we can prove:

**Theorem 28** *(Cauchy’s theorem)* If \( f(z) \) is analytic everywhere inside a simple closed curve \( C \) then

\[
\int_C f(z)\,dz = 0 \quad (620)
\]

**Proof.** Let \( R \) be the region inside \( C \). Then by Green’s theorem followed by the CR equations:

\[
\int_C f(z)\,dz = \left( \int_C u\,dx - v\,dy \right) + i \left( \int_C v\,dx + u\,dy \right) = \int_R \left(- \frac{\partial v}{\partial x} - \frac{\partial u}{\partial y}\right)\,dxdy + i \int_R \left( \frac{\partial u}{\partial x} - \frac{\partial v}{\partial y}\right)\,dxdy \quad (621)
\]

\[= 0\]

**examples:** Let \( C \) be the unit circle traversed counterclockwise. Then we have

\[
\int_C z^n\,dz = 0 \quad n = 0, 1, 2, \ldots
\]

\[
\int_C e^{z^2+\cos(3z)}\,dz = 0
\]

\[
\int_C \frac{1}{z^2}\,dz = 0
\]

since in each case the integrand is analytic inside the circle. But for integrals like

\[
\int_C e^z\,dz \quad \int_C \frac{1}{z - 1/2}\,dz \quad \int_C \log z\,dz \quad (623)
\]

there is no conclusion since the integrand is not analytic inside the circle

We define a **region** as a connected open set.

**Corollary** If \( f \) is analytic inside a simply connected region \( D \) then \( \int_C f(z)\,dz = 0 \) for any closed curve in \( D \).

**Proof.** Suppose \( C \) is a simple closed curve. Since \( D \) is simply connected the interior of \( C \) is in \( D \), hence \( f \) is analytic inside \( C \) and the result follows by Cauchy’s theorem. If \( C \) is not simple break it up into pieces which are simple.
Corollary If \( f \) is analytic inside a simply connected region \( D \) then integrals \( \int_C f(z)dz \) are independent of path in \( D \).

**Proof.** If \( C_1 \) and \( C_2 \) are paths in \( D \) with the same endpoints then \( C_1 - C_2 \) is a closed curve in \( D \) and so \( \int_{C_1} f - \int_{C_2} f = \int_{C_1 - C_2} f = 0 \) by the previous Corollary.

**Corollary** (deformation theorem) Let \( f \) be analytic in a region \( D \) and let \( C_1, C_2 \) be simple closed curves in \( D \) such that \( C_1 \) can be continuously deformed to \( C_2 \) in \( D \). Then

\[
\int_{C_1} f = \int_{C_2} f
\]

(624)

**Proof.** We prove the result in the case where \( C_1, C_2 \) are simple closed curves and \( C_1 \) is inside \( C_2 \). The hypotheses of the theorem imply that \( f \) is analytic in the region between them.

Let \( C \) be a curve which traverses \( C_1 \), then jumps to \( C_2 \) along a path \( \gamma \), then traverses \(-C_2\), then jumps back to \( C_1 \) along \(-\gamma\). (see figure 33) Then \( C \) is a closed curve and \( f \) is analytic inside it and so by Cauchy’s theorem

\[
\int_{C_1} f - \int_{C_2} f = \int_{C_1} f + \int_{\gamma} f + \int_{-C_2} f + \int_{-\gamma} f = \int_{C} f = 0
\]

(625)
3.12 Cauchy integral formula

We need formulas for the direct evaluation of complex line integrals, i.e. without making a parametrization. The first is a generalization of the fundamental theorem of calculus.

**Theorem 29** Let $f$ be analytic in a region $D$. Then $\int_C f'(z)dz$ is independent of path in $D$ and so can be written $\int_{z_0}^{z_1} f'(z)dz$. We have

$$\int_{z_0}^{z_1} f'(z)dz = f(z_1) - f(z_0) \quad (626)$$

**Proof.** Given points $z_0, z_1$ in $D$ let $C$ be any path from $z_0$ to $z_1$ and let $z(t), a \leq t \leq b$ be any parametrization of $C$ with $z(a) = z_0$ and $z(b) = z_1$. Then by the chain rule

$$\int_C f'(z)dz = \int_a^b f'(z(t))z'(t)dt$$

$$= \int_a^b \frac{d}{dt}[f(z(t))]dt$$

$$= f(z(b)) - f(z(a))$$

$$= z_1 - z_0 \quad (627)$$

**example:**

$$\int_0^{1+i} z^3 dz = \left[ \frac{z^4}{4} \right]_0^{1+i} = \frac{(1 + i)^4}{4} = \frac{\sqrt{2}e^{i\pi/4})^4}{4} = e^{i\pi} = -1 \quad (628)$$

To use the theorem one still has to find an anti-derivative for the integrand. For closed curves there is another formula which evaluates integrals with no work at all.

**Theorem 30** *(Cauchy Integral Formula)* Let $f(z)$ be analytic inside a simple closed curve $C$ traversed counterclockwise. Let $z_0$ be a point inside $C$. Then

$$\int_C \frac{f(z)}{z - z_0} = 2\pi i f(z_0) \quad (629)$$

Note that if $z_0$ is outside $C$ then $f(z)/(z - z_0)$ is analytic inside $C$ and so the integral is zero by Cauchy’s theorem.
Proof. Let $C'$ be a little circle of radius $r$ centered on $z_0$. (see figure 34) Since $f$ is analytic between $C$ and $C'$ we have by the deformation theorem

$$\int_{C'} \frac{f(z)}{z-z_0} \, dz = \int_{C'} \frac{f(z)}{z-z_0} \, dz$$

(630)

Now parametrize $C'$ by

$$z = z_0 + re^{i\theta} \quad 0 \leq \theta \leq 2\pi$$

$$dz = ire^{i\theta} \, d\theta$$

(631)

Then

$$\int_{C'} \frac{f(z)}{z-z_0} \, dz = \int_0^{2\pi} \frac{f(z_0 + re^{i\theta})ire^{i\theta}}{re^{i\theta}} \, d\theta = i \int_0^{2\pi} f(z_0 + re^{i\theta}) \, d\theta$$

(632)

This is true for any small $r > 0$. Thus it is also true in the limit $r \to 0$ which gives

$$i \int_0^{2\pi} f(z_0) \, d\theta = 2\pi if(z_0)$$

(633)
**Examples:** Let $C$ be the circle $|z| = 2$ traversed counterclockwise. Then

\[
\int_C \frac{e^z}{z - 1} \, dz = 2\pi i [e^z]_{z=1} = 2\pi i e
\]

\[
\int_C \frac{e^z}{z - 3} \, dz = 0
\]

\[
\int_C \frac{4z^2 + 7}{z - 1} \, dz = 2\pi i [4z^2 + 7]_{z=1} = 22\pi i
\]

**Example:** Let $C$ be the square with corners $i, -i, 2 - i, 2 + i$ traversed in that order. We want to evaluate

\[
\int_C \frac{e^z}{z^2 - 1} \, dz = \int_C \frac{e^z}{(z + 1)(z - 1)} \, dz
\]

The integrand goes bad at both $z = \pm 1$ but only $z = 1$ is inside $C$. Thus $e^z/(z + 1)$ is analytic inside the curve and we can evaluate the integral as

\[
\int_C \frac{e^z/(z + 1)}{z - 1} \, dz = 2\pi i [e^z/(z + 1)]_{z=1} = \pi i e
\]

**Example:** Let $C$ be the circle $|z| = 2$ traversed counterclockwise. We want to evaluate

\[
\int_C \frac{e^z}{z^2 + 1} \, dz = \int_C \frac{e^z}{(z + i)(z - i)} \, dz
\]

The integrand goes bad at $z = \pm i$ and this time both points are inside $C$. Thus we cannot use the Cauchy integral formula as it stands. We give two methods to modify the problem so that we can use it.

**Solution (1).** Break the denominator up using partial fractions We look for $A, B$ such that

\[
\frac{1}{(z + i)(z - i)} = \frac{A}{z - i} + \frac{B}{z + i}
\]

for all $z$. This is the same as

\[
1 = A(z + i) + B(z - i)
\]

Matching coefficients gives $A + B = 0$ and $Ai - Bi = 1$. The solution is $A = 1/2i$ and $B = -1/2i$. Therefore

\[
\frac{1}{(z + i)(z - i)} = \frac{1}{2i} \left( \frac{1}{z - i} - \frac{1}{z + i} \right)
\]

Inserting this into the integral it becomes

\[
\frac{1}{2i} \int_C \frac{e^z}{z - i} \, dz - \frac{1}{2i} \int_C \frac{e^z}{z + i} \, dz = 2\pi i \left( \frac{e^i - e^{-i}}{2i} \right) = 2\pi i \sin(1)
\]
Solution (2). Deform the curve to a pair of little circles \( C_1 \) around \( z = i \) and \( C_2 \) around \( z = -i \). (see figure 35.) Then the integral is

\[
\int_{C_1} \frac{e^z}{(z+i)(z-i)} \, dz + \int_{C_2} \frac{e^z}{(z+i)(z-i)} \, dz
\]

\[
= 2\pi i \left. \frac{e^z}{z+i} \right|_{z=i} + 2\pi i \left. \frac{e^z}{z-i} \right|_{z=-i}
\]

\[
= 2\pi i \left( e^i - e^{-i} \right) = 2\pi i \sin(1)
\]
3.13 higher derivatives

Let \( f \) be analytic in a region \( D \), let \( z \) be a point in \( D \), and let \( \mathcal{C} \) be a simple closed curve enclosing \( z \) such that the interior is contained in \( D \), for example \( \mathcal{C} \) could be a little circle around \( z \). By the Cauchy integral formula

\[
f(z) = \frac{1}{2\pi i} \int_{\mathcal{C}} \frac{f(\zeta)}{\zeta - z} d\zeta \tag{643}
\]

The integrand here is a differentiable function of \( z \) since

\[
\frac{d}{dz} \left[ \frac{1}{\zeta - z} \right] = \frac{1}{(\zeta - z)^2} \tag{644}
\]

One can differentiate under the integral sign obtain

\[
f'(z) = \frac{1}{2\pi i} \int_{\mathcal{C}} \frac{f(\zeta)}{(\zeta - z)^2} d\zeta \tag{645}
\]

Repeat the argument and conclude that \( f'(z) \) is differentiable and that

\[
f''(z) = \frac{2}{2\pi i} \int_{\mathcal{C}} \frac{f(\zeta)}{(\zeta - z)^3} d\zeta \tag{646}
\]

Repeat the argument and conclude that \( f''(z) \) is differentiable and that

\[
f'''(z) = \frac{3 \cdot 2}{2\pi i} \int_{\mathcal{C}} \frac{f(\zeta)}{(\zeta - z)^4} d\zeta \tag{647}
\]

In fact \( f \) can be differentiated any number of times and

\[
f^{(n)}(z) = \frac{n!}{2\pi i} \int_{\mathcal{C}} \frac{f(\zeta)}{(\zeta - z)^{n+1}} d\zeta \tag{648}
\]

This is a remarkable result: assuming only \( f \) is analytic, i.e. once differentiable, we conclude that it is infinitely differentiable. This can certainly fail for functions of a real variable.

The last formula can also be used to evaluate integrals. We replace \( z \) by \( z_0 \) and \( \zeta \) by \( z \) and state it as follows.

**Theorem 31** Let \( f \) be analytic inside a simple closed curve \( \mathcal{C} \) traversed counterclockwise. Let \( z_0 \) be a point inside \( \mathcal{C} \). Then for \( n = 0, 1, 2, \ldots \)

\[
\int_{\mathcal{C}} \frac{f(z)}{(z - z_0)^{n+1}} dz = \frac{2\pi i}{n!} f^{(n)}(z_0) \tag{649}
\]
**Notation:** If $C$ is a circle of radius $r$ centered on $z_0$ traversed counterclockwise then

$$\int_C f(z)dz$$

can be written

$$\oint_{|z-z_0|=r} f(z)dz \quad (650)$$

**examples**

$$\oint_{|z|=2} \frac{e^{6z}}{z-1} dz = 2\pi i \left[ e^{6z} \right]_{z=1} = 2\pi i e^6$$

$$\oint_{|z|=2} \frac{e^{6z}}{(z-1)^2} dz = \frac{2\pi i}{1!} \left[ \frac{d}{dz} e^{6z} \right]_{z=1} = 12\pi i e^6 \quad (651)$$

$$\oint_{|z|=2} \frac{e^{6z}}{(z-1)^3} dz = \frac{2\pi i}{2!} \left[ \frac{d^2}{dz^2} e^{6z} \right]_{z=1} = 36\pi i e^6$$

**example**

$$\oint_{|z|=1} \frac{\sin(z)\sin(z)}{z^2} dz = 2\pi i \left[ \frac{d}{dz} \sin(z) \right]_{z=0} = 2\pi i [\cos(z)\cos(z)]_{z=0} = 2\pi i \quad (652)$$

### 3.14 Cauchy inequalities

Let $C$ be a circle of radius $R$ around a point $a$. Suppose that $f(z)$ is analytic inside the circle and that on the circle we have the bound

$$|f(z)| \leq M \quad z \in C \quad (653)$$

The Cauchy integral formula says that

$$f(a) = \frac{1}{2\pi i} \oint_C \frac{f(z)}{z-a} dz \quad (654)$$

We use this to estimate $|f(a)|$. For $z$ on $C$ we have

$$\left| \frac{f(z)}{z-a} \right| = \frac{|f(z)|}{|z-a|} = \frac{|f(z)|}{R} \leq \frac{M}{R} \quad (655)$$

Then length of the curve is $2\pi R$. Thus

$$|f(a)| \leq \frac{M}{2\pi R} \cdot 2\pi R = M \quad (656)$$

We could take $M$ to be the maximum of $|f(z)|$ on the circle and then this says that $|f|$ at the center of a circle is less than or equal to the maximum of $|f|$ on the circle.
More generally we have
\[ f^{(n)}(a) = \frac{n!}{2\pi i} \int_C \frac{f(z)}{(z-a)^{n+1}} dz \]  \hspace{1cm} (657)

Now \( z \) on \( C \) we have
\[ \left| \frac{f(z)}{(z-a)^{n+1}} \right| \leq \frac{M}{R^{n+1}} \]  \hspace{1cm} (658)

and so
\[ |f^{(n)}(a)| \leq \frac{n!}{2\pi} \frac{M}{R^{n+1}} 2\pi R = \frac{n!M}{R^n} \]  \hspace{1cm} (659)

We restate the result as a theorem.

**Theorem 32** (*Cauchy inequalities.*) Let \( f \) be analytic inside a circle of radius \( R \) centered on \( a \) and suppose \( |f| \leq M \) on the circle. Then for \( n = 0, 1, 2, \ldots \)
\[ |f^{(n)}(a)| \leq \frac{n!M}{R^n} \]  \hspace{1cm} (660)

This has some interesting consequences.

**Theorem 33** (*Liouville’s theorem*) Suppose \( f \) is analytic in the entire plane and \( f \) is bounded, that is \( |f(z)| \leq M \) for some \( M \) and all \( z \). Then \( f \) is constant.

**Proof.** \( f \) is analytic inside any circle of radius \( R \) centered on any \( a \), and is bounded by \( M \) on the circle. By the \( n = 1 \) Cauchy inequality
\[ |f'(a)| \leq \frac{M}{R} \]  \hspace{1cm} (661)

Taking the limit \( R \to \infty \) gives
\[ |f'(a)| \leq 0 \]  \hspace{1cm} (662)

Hence \( f'(a) = 0 \) for all \( a \) and so \( f \) is a constant.

**Theorem 34** (*Fundamental theorem of algebra*) Every polynomial \( P(z) \) of degree \( n \geq 1 \) has \( n \) roots (not necessarily distinct)

**Proof.** Suppose it is not true and \( P(z) \neq 0 \) for all \( z \). Then \( 1/P(z) \) is analytic everywhere. Furthermore if \( P(z) = a_n z^n + \cdots + a_1 z + a_0 \) with \( a_n \neq 0 \) then
\[ \lim_{z \to \infty} \frac{1}{P(z)} = \lim_{z \to \infty} \frac{1}{a_n z^n + \cdots + a_1 z + a_0} \]
\[ = \lim_{z \to \infty} \frac{z^{-n}}{a_n + a_{n-1} z^{-1} \cdots + a_1 z^{-n+1} + a_0 z^{-n}} \]
\[ = 0 \]
\[ a_n = 0 \]
From this and the fact that \(1/P(z)\) is continuous one can deduce that it is bounded. Then by Liouville’s theorem \(1/P(z)\) is constant and hence \(P(z)\) is constant. This is a contradiction. Thus we must have \(P(z) = 0\) for some \(z\). The polynomial has at least one root.

If we call the root \(a_1\) then \(P(z)\) must have \(z - a_1\) as a factor so

\[
P(z) = (z - a_1)P_1(z) \tag{664}
\]

where \(P_1(z)\) is a polynomial of degree \(n - 1\). Repeating the argument \(P_1\) must have a root \(a_2\) and hence a factor \(z - a_2\). Then

\[
P(z) = (z - a_1)(z - a_2)P_2(z) \tag{665}
\]

where \(P_2(z)\) has degree \(n - 2\). After \(n\) steps we are left with \(n\) factors \(z - a_i\) and a polynomial of degree 0 which is a constant \(c\)

\[
P(z) = (z - a_1)(z - a_2) \cdots (z - a_n)c \tag{666}
\]

This exhibits the \(n\) roots.

### 3.15 real integrals

The Cauchy integral formula is so easy to use that it is worthwhile going to a lot of trouble to mold other integrals into this form. In particular this can be done for certain real integrals.

Consider first integrals of the form \(\oint_{|z|=1} f(z)dz\). Parametrizing the unit circle by \(z = e^{i\theta}, 0 \leq \theta \leq 2\pi\) and \(dz = ie^{i\theta}d\theta\) we have

\[
\oint_{|z|=1} f(z)dz = \int_0^{2\pi} f(e^{i\theta})ie^{i\theta}d\theta \tag{667}
\]

Now we reverse this process. Suppose we an integral of the form \(\int_0^{2\pi} F(\cos \theta, \sin \theta)d\theta\) for some real function \(F\). (Such integrals come up for example when computing Fourier series.) Then we can write it a complex integral over the unit circle by

\[
\int_0^{2\pi} F(\cos \theta, \sin \theta)d\theta \\
= \int_0^{2\pi} F\left(\frac{e^{i\theta} + e^{-i\theta}}{2}, \frac{e^{i\theta} - e^{-i\theta}}{2i}\right)ie^{i\theta}d\theta \\
= \int_{|z|=1} F\left(\frac{z + z^{-1}}{2}, \frac{z - z^{-1}}{2i}\right)\frac{dz}{iz}
\]

Thus the method is to make the substitutions

\[
\cos \theta = \frac{z + z^{-1}}{2} \quad \sin \theta = \frac{z - z^{-1}}{2i} \quad d\theta = \frac{dz}{iz} \tag{669}
\]
and replace the integral over \([0, 2\pi]\) by an integral over \(|z| = 1\).

In evaluating integrals over the unit circle it is sometimes useful to keep in mind that for any integer \(n\)

\[
\oint_{|z|=1} z^n dz = \begin{cases} 
0 & n \geq 0 \\
2\pi i & n = -1 \\
0 & n \leq -2 
\end{cases} \tag{670}
\]

This follows by Cauchy's theorem for \(n > 0\) and by the Cauchy integral formula in the other cases.

**example**

\[
\int_0^{2\pi} \sin^4 \theta \, d\theta = \oint_{|z|=1} \frac{(z - z^{-1})^4 \, dz}{2iz} = \frac{1}{16} \oint_{|z|=1} (z^4 - 4z^2 + 6 - 4z^{-2} + z^{-4}) \frac{dz}{iz}
\]

\[
= \frac{1}{16i} \oint_{|z|=1} (z^3 - 4z + 6z^{-1} - 4z^{-3} + z^{-5}) \, dz = \frac{1}{16i} \cdot 6 \cdot 2\pi i = \frac{3}{4} \pi \tag{671}
\]

**example**

\[
\int_0^{2\pi} \frac{d\theta}{2 + \cos \theta} = \oint_{|z|=1} \frac{1}{2 + \frac{z + z^{-1}}{2}} \frac{dz}{iz} = \frac{1}{i} \oint_{|z|=1} \frac{2dz}{z^2 + 4z + 1}
\]

The denominator vanishes when \(z^2 + 4z + 1 = 0\) which occurs at

\[
z = \frac{-4 \pm \sqrt{16 - 4}}{2} = -2 \pm \sqrt{3} \tag{673}
\]

Hence the polynomial factors as

\[
z^2 + 4z + 1 = (z + 2 + \sqrt{3})(z + 2 - \sqrt{3}) \tag{674}
\]

Put this into the integral and note that the only part not analytic inside the circle is the factor \((z + 2 - \sqrt{3})^{-1}\). Thus the integral can be evaluated as

\[
\frac{1}{i} \oint_{|z|=1} \frac{2dz}{(z + 2 + \sqrt{3})(z + 2 - \sqrt{3})} = \frac{1}{i} \cdot 2\pi i \left[ \frac{2}{z + 2 + \sqrt{3}} \right]_{z=-2+\sqrt{3}} = \frac{2\pi}{\sqrt{3}} \tag{675}
\]
There is another class of real integrals which can be treated by complex variable
techniques. These are integrals over the entire real line. We illustrate with an example.

**example:** Suppose we want to find the integral \( \int_{-\infty}^{\infty} (1 + x^2)^{-1} \, dx \). We can proceed as follows

\[
\int_{-\infty}^{\infty} \frac{dx}{1 + x^2} = \lim_{R \to \infty} \int_{-R}^{R} \frac{dx}{1 + x^2}
\]

\[
= \lim_{R \to \infty} \int_{L_R} \frac{dz}{1 + z^2}
\]

\[
= \lim_{R \to \infty} \left( \int_{L_R + C_R} \frac{dz}{1 + z^2} - \int_{C_R} \frac{dz}{1 + z^2} \right)
\]

Here we treat the real integral as a complex integral over the line \( L_R \) from \(-R \) to \( R \), then we turn it into a closed curve by adding a semi-circle \( C_R \) of radius \( R \) in the upper half plane. (see figure 36) Of course we also have to subtract the contribution of the semi-circle. Now the idea is to evaluate the integral over \( L_R + C_R \) by the Cauchy integral formula and show that the integral over \( C_R \) goes to zero as \( R \to \infty \).

We have for any \( R > 1 \)

\[
\int_{L_R + C_R} \frac{dz}{1 + z^2} = \int_{L_R + C_R} \frac{dz}{(z + i)(z - i)} = 2\pi i \left[ \frac{1}{z + i} \right]_{z=i} = \pi
\]

(677)

Note that only \( z = i \) is inside the curve, not \( z = -i \). The result is independent of \( R \) and so the limit as \( R \to \infty \) is also \( \pi \).

For the other term note that if \( z \) is on \( C_R \) then \( |z^2 + 1| \geq |z|^2 - 1 = R^2 - 1 \) and therefore

\[
\frac{1}{|1 + z^2|} \leq \frac{1}{R^2 - 1}
\]

(678)

Then since \( C_R \) has length \( \pi R \) we have

\[
0 \leq \left| \int_{C_R} \frac{dz}{1 + z^2} \right| \leq \frac{\pi R}{R^2 - 1}
\]

(679)

But \( \pi R/(R^2 - 1) \) goes to zero at \( R \to \infty \). Therefore

\[
\lim_{R \to \infty} \int_{C_R} \frac{dz}{1 + z^2} = 0
\]

(680)

Thus the answer is

\[
\int_{-\infty}^{\infty} \frac{dx}{1 + x^2} = \pi
\]

(681)
3.16 Fourier and Laplace transforms

(A.) If \( f(x) \) is a function defined for \( -\infty < x < \infty \), the Fourier transform of \( f \) is a new function \( \tilde{f} \) defined by

\[
\tilde{f}(k) = \int_{-\infty}^{\infty} e^{ikx} f(x) \, dx
\]  

(682)

The function can be recovered from its Fourier transform by the inversion formula

\[
f(x) = \frac{1}{2\pi} \int_{-\infty}^{\infty} e^{-ikx} \tilde{f}(k) \, dk
\]

(683)

Fourier transforms are useful for solving partial differential equations on the whole line, among other things. Complex integration techniques are useful for computing Fourier transforms.
**Example:** Suppose we want to find the Fourier transform of the function \((1 + x^2)^{-1}\). We proceed as in the last example

\[
\int_{-\infty}^{\infty} \frac{e^{ikx}}{1 + x^2} dx = \lim_{R \to \infty} \int_{-R}^{R} \frac{e^{ikx}}{1 + x^2} dx
\]

\[
= \lim_{R \to \infty} \int_{L_R} \frac{e^{ikz}}{1 + z^2} dz
\]

\[
= \lim_{R \to \infty} \left( \int_{L_R + C_R} \frac{e^{ikz}}{1 + z^2} dz - \int_{C_R} \frac{e^{ikz}}{1 + z^2} dz \right)
\]

Suppose that \(k > 0\). Then for \(z = x + iy\) on the semi-circle \(C_R\) we have \(y > 0\) and so \(|e^{ikz}| = |e^{ikx}|e^{-ky}| = e^{-ky} \leq 1\). From before \(|1 + z^2| \geq R^2 - 1\) and so on \(C_R\)

\[
\left| \frac{e^{ikz}}{1 + z^2} \right| \leq \frac{1}{R^2 - 1}
\]

Then

\[
\left| \int_{C_R} \frac{e^{ikz}}{1 + z^2} dz \right| \leq \frac{\pi R}{R^2 - 1} \to 0 \text{ as } R \to \infty
\]

There is no contribution from the integral over \(C_R\).

On the other hand by the Cauchy integral formula since \(z = i\) is inside the curve and \(z = -i\) is not

\[
\int_{L_R + C_R} \frac{e^{ikz}}{1 + z^2} dz = \int_{L_R + C_R} \frac{e^{ikz}}{(z + i)(z - i)} dz = 2\pi i \left[ \frac{e^{ikz}}{z + i} \right]_{z=i} = \pi e^{-k}
\]

This holds for any \(R > 1\) and hence also in the limit \(R \to \infty\) and so

\[
\int_{-\infty}^{\infty} \frac{e^{ikx}}{1 + x^2} dx = \pi e^{-k} \quad k > 0
\]

This analysis fails if \(k < 0\) since then \(|e^{ikz}| = e^{-ky}\) grows exponentially in \(y\) and we cannot argue that the contribution from \(C_R\) vanishes. Instead we close the curve with a semi-circle \(C'_R\) in the lower half plane. We have

\[
\int_{-\infty}^{\infty} \frac{e^{ikx}}{1 + x^2} dx = \lim_{R \to \infty} \left( \int_{L_R + C'_R} \frac{e^{ikz}}{1 + z^2} dz - \int_{C'_R} \frac{e^{ikz}}{1 + z^2} dz \right)
\]

For \(z = x + iy\) on \(C'_R\) we have \(y < 0\) and so \(|e^{ikz}| = e^{-ky} < 1\) for \(k < 0\). Then again

\[
\left| \int_{C'_R} \frac{e^{ikz}}{1 + z^2} dz \right| \leq \frac{\pi R}{R^2 - 1} \to 0 \text{ as } R \to \infty
\]

and there is no contribution from \(C'_R\).
For the other term it is now the point \( z = -i \) which is inside the closed curve \( L_R + C'_R \). We are going around the curve clockwise so when we use the Cauchy integral formula there is an overall minus sign. We have

\[
\int_{L_R+C'_R} \frac{e^{ikz}}{1+z^2} \, dz = \int_{L_R+C'_R} \frac{e^{ikz}}{(z+i)(z-i)} \, dz = -2\pi i \left[ \frac{e^{ikz}}{z-i} \right]_{z=-i} = \pi e^k \quad (691)
\]

Therefore

\[
\int_{-\infty}^{\infty} \frac{e^{ikx}}{1+x^2} \, dx = \pi e^k \quad k < 0 \quad (692)
\]

The results for the three cases \( k > 0, k = 0, k < 0 \) can be summarized by

\[
\int_{-\infty}^{\infty} \frac{e^{ikx}}{1+x^2} \, dx = \pi e^{-|k|} \quad (693)
\]

(B.) If \( f(t) \) is a function defined for \( 0 \leq t < \infty \) the Laplace transform of \( f \) is a new function defined by

\[
F(s) = \int_0^\infty e^{-st} f(t) \, dt \quad (694)
\]

if the integral converges. The Laplace transform is useful for solving linear ordinary differential equations.

Suppose there are constants \( K, c \) such that

\[
|f(t)| \leq Ke^{ct} \quad (695)
\]

then

\[
|e^{-st} f(t)| \leq Ke^{-(s-c)t} \quad (696)
\]

This is rapidly decreasing if \( s > c \), in which case the integral converges. Thus \( F(s) \) is defined for \( s > c \).

We could also let \( s \) be complex. Then \( |e^{-st}| = e^{-(Re \, s)t} \) and

\[
|e^{-st} f(t)| \leq Ke^{-((Re \, s)-c)t} \quad (697)
\]

Then the integral converges and the Laplace transform is defined in the half plane \( \text{Re} \, s > c \). Furthermore one can show that it is analytic in this region (think of differentiating under the integral sign).

There is also an inversion formula for the Laplace transform. If \( f(s) \) is analytic for \( \text{Re} \, (s) > c \) and \( \gamma > c \), then for \( t > 0 \)

\[
f(t) = \frac{1}{2\pi i} \int_{\gamma-i\infty}^{\gamma+i\infty} e^{st} F(s) \, ds \quad (698)
\]
Here the integral is over the vertical line \( \text{Re} \,(s) = \gamma \). Such lines can be deformed to each other in the region of analyticity, so it does not matter what \( \gamma \) we take.

**problem:** Find the function whose Laplace transform is \((s - 2)^{-3}\)

**solution:** The function is analytic for \( \text{Re}(s) > 2 \) so the function is given by the inversion formula with \( \gamma > 2 \). We evaluate it by closing the curve in the left half plane:

\[
f(t) = \frac{1}{2\pi i} \int_{\gamma - i\infty}^{\gamma + i\infty} \frac{e^{st}}{(s - 2)^3} \, ds
= \lim_{R \to \infty} \frac{1}{2\pi i} \int_{\gamma - iR}^{\gamma + iR} \frac{e^{st}}{(s - 2)^3} \, ds
= \lim_{R \to \infty} \frac{1}{2\pi i} \int_{L_R} e^{st} \, ds
= \lim_{R \to \infty} \left( \frac{1}{2\pi i} \int_{L_R + C_R} e^{st} \, ds - \frac{1}{2\pi i} \int_{C_R} e^{st} \, ds \right)
\]

Here \( L_R \) is the straight line from \( \gamma - iR \) to \( \gamma + iR \) and \( C_R \) is the semi-circle of radius \( R \) centered on \( \gamma \), see figure 37.
Now for $s$ on $C_R$ we have $\text{Re}(s) \leq \gamma$ and so $|e^{st}| = e^{\text{Re}(s)t} \leq e^{\gamma t}$. Also the point on $C_R$ closest to 2 is $\gamma - R$ and so $|s - 2| \geq |2 - (\gamma - R)|$. Therefore

$$\left| \frac{1}{2\pi i} \int_{C_R} \frac{e^{st}}{(s - 2)^3} \, ds \right| \leq \frac{1}{2\pi} \frac{e^{\gamma t}}{|2 - (\gamma - R)|^3} \cdot \pi R \to 0 \text{ as } R \to \infty \quad (700)$$

The integral over $C_R$ does not contribute.

If $R$ is large enough, then the point 2 is inside $L_R + C_R$ and so by the higher derivative Cauchy integral formula

$$\frac{1}{2\pi i} \int_{L_R+C_R} \frac{e^{st}}{(s - 2)^3} \, ds = \frac{1}{2!} \frac{d^2}{ds^2} [e^{st}]_{s=2} = \frac{1}{2} t^2 e^{2t} \quad (701)$$

This also holds in the limit $R \to \infty$ and so this answer is $f(t) = \frac{1}{2} t^2 e^{2t}$. 

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