Abstract

In this course, we shall study differential vector calculus, which is the branch of mathematics that deals with differentiation and integration of scalar and vector fields. We shall encounter many examples of vector calculus in physics.
Timetable

- Tuesday and Friday 11:10-12:00 Lecture (JCMB Lecture Theatre A)
- Tuesday 14:10-16:00 Tutorial Workshop (JCMB Teaching Studio 3217)
- Thursday 14:10-16:00 Tutorial Workshop (JCMB Room 1206c)

Students should attend both lectures and one tutorial workshop each week. Tutorials start in Week 2.

Genealogy

For historians of pre-Honours courses . . .

This course was known as *Mathematics for Physics 4: Fields* until 2012-13, when it became *Vector Calculus*.

There will be some evolution from last year’s instance of the course, but I’m not planning any major structural changes. There should be some new material on index notation.

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Synopsis

We will cover all the topics below, but not necessarily in that order.

- Introduction, scalar and vector fields in gravitation and electrostatics. Revision of vector algebra and products.
- Fields, potentials, grad, div and curl and their physical interpretation, the Laplacian, vector identities involving grad, div, curl and the Laplacian. Physical examples.
- Lines and surfaces. Line integrals, vector integration, physical applications.
- Surface and volume integrals, divergence and Stokes’ theorems, Green’s theorem and identities, scalar and vector potentials; applications in electromagnetism and fluids.
- Curvilinear coordinates, line, surface, and volume elements; grad, div, curl and the Laplacian in curvilinear coordinates. More examples.

Syllabus

The Contents section of this document is the course syllabus!

Books

The course will not use any particular textbook. The first six listed below are standard texts; Spiegel contains many examples and problems. Matthews works with index notation ($\delta_{ij}$- and $\epsilon_{ijk}$-tensors) which is very useful. All are available from Amazon.co.uk

- KF Riley and MP Hobson, Essential Mathematical Methods for the Physical Sciences (CUP) (also useful for Junior Honours)
  or Foundation Mathematics for the Physical Sciences (CUP) (not so good for JH)
- KF Riley, MP Hobson and SJ Bence, Mathematical Methods for Physics and Engineering, (CUP).
  (This is an older, but more comprehensive version of the books above.)
- PC Matthews, Vector Calculus, (Springer). (Also useful for JH SoCM)
- Any mathematical methods book you’re comfortable with.
Prologue

This course deals with vector calculus and its differential version. In particular we will study the vector (or more generally the tensor tensor) formalism of the three dimensional Euclidian space. Intuitively the latter is the space we live in and it is therefore not surprising that there are many applications. Having our three dimensional space as a basis will allow us to resort to geometrical understanding. One of the most notable applications of the three dimensional vector calculus is Maxwell’s theory of electromagnetism.

There is no problem in extending any of the learnt material to higher dimensional spaces. This becomes relevant when studying Einstein’s theory of special relativity where space and time are united into a four dimensional space for example. Results even find their counterparts in curved spaces which is the basis of differential geometry and Einstein’s theory of general relativity. The underlying structure of the three dimensional vector calculus is the rotational symmetry of the three dimensional space which is described by the rotation group denoted by \(SO(3)\). From a mathematical point of view vectors are objects that transform in a definite way under rotations. The latter subject is known as representation theory and is part of the course on symmetry of quantum mechanics which I also happen to teach. Of course these matters go beyond the scope of this course but the reason for mentioning them is to broaden your horizon and mainly to motivate you for this course. In fact one might say that having a solid basis in vector calculus opens the door to an understanding of these subsequent matters.

1 Fields and why we need them in Physics

1.1 Vectors and scalars

We start by recalling two basic definitions (simple-minded Physics versions) in order to establish our notation (which is similar to that used in Linear Algebra and Several Variable Calculus.)

**Scalar:** a quantity specified by a single number. This quantity is invariant (unchanged) under rotations.

**Vector:** a quantity specified by a number (magnitude) and a direction (two numbers in three dimensions, e.g. two angles). This quantity transform specifically under rotations.

**Examples:** mass is a scalar, velocity is a vector.

**Example:** A position vector is a vector bound to some origin and gives the position of some point \(P\), say, relative to that origin. It is often denoted by \(\mathbf{r}\) (or \(x\) or \(\overrightarrow{OP}\)).

Define an orthonormal right-handed Cartesian basis of unit vectors \(\{\xi_1, \xi_2, \xi_3\}\), such that

\[1\] The word orthonormal means mutually orthogonal (perpendicular) and normalised to have unit length.
In such a basis, we may write the position vector \( \mathbf{r} \) in terms of its Cartesian components \( (x_1, x_2, x_3) \) as follows

\[
\mathbf{r} \equiv x_1 \mathbf{e}_1 + x_2 \mathbf{e}_2 + x_3 \mathbf{e}_3 .
\]

This is an important notation that you will have to remember throughout this course. The length or magnitude of the vector \( \mathbf{r} \) is a scalar and is denoted by \( r \equiv |\mathbf{r}| = \sqrt{x_1^2 + x_2^2 + x_3^2} \).

Any vector \( \mathbf{a} \) may written in this notation as \[\mathbf{a} = a_1 \mathbf{e}_1 + a_2 \mathbf{e}_2 + a_3 \mathbf{e}_3 = \sum_{i=1}^{3} a_i \mathbf{e}_i .\]

1.2 Fields

In physics we have quantities that vary in some region of space, e.g. the temperature \( T(r) \) of the ocean depends on position \( \mathbf{r} \). To study this variation we require the concept of a field. In these lectures we shall develop the calculus of scalar fields and vector fields.

If to each point \( r \) in some region of space there corresponds a scalar \( \phi(x_1, x_2, x_3) \), then \( \phi(r) \) is a scalar field: \( \phi \) is a function of the three Cartesian position coordinates \( (x_1, x_2, x_3) \).

Examples: the temperature distribution in a body \( T(r) \), pressure in the atmosphere \( P(r) \), electric charge (or mass) density \( \rho(r) \), electrostatic potential \( \phi(r) \), the Higgs field \( h(r) \).

Similarly a vector field assigns a vector \( \mathbf{V}(r) \) to each point \( r \) of some region.

Examples: velocity in a fluid \( \mathbf{v}(r) \), electric current density \( \mathbf{j}(r) \), electric field \( \mathbf{E}(r) \), magnetic field \( \mathbf{B}(r) \) (actually a pseudo-vector field).

A vector field in 2-d can be represented graphically, at a carefully selected set of points \( r \), by an arrow whose length and direction is proportional to \( \mathbf{V}(r) \), e.g. wind velocity on a weather forecast chart.

1.3 Examples: Gravitation and Electrostatics

Let us revisit two familiar examples of fundamental fields in Nature.

Gravitation: The foundation of Newtonian Gravity is Newton’s Law of Gravitation, which Newton deduced from observations of the motion of the planets by Tycho Brahe, and their analysis by Kepler.

\[\mathbf{e}_1 \cdot \mathbf{e}_1 = \mathbf{e}_2 \cdot \mathbf{e}_2 = \mathbf{e}_3 \cdot \mathbf{e}_3 = 1 \quad \text{and} \quad \mathbf{e}_1 \cdot \mathbf{e}_2 = \mathbf{e}_1 \cdot \mathbf{e}_3 = \mathbf{e}_2 \cdot \mathbf{e}_3 = 0. \] (1)

---

\[\text{The ‘xyz’ notation is a popular alternative in which the Cartesian components of a vector } \mathbf{a} \text{ are written as } (a_x, a_y, a_z), \text{ and a vector is written in terms of orthonormal basis vectors } \{\mathbf{i}, \mathbf{j}, \mathbf{k}\} \text{ as } \mathbf{a} = a_x \mathbf{i} + a_y \mathbf{j} + a_z \mathbf{k} \quad \text{or perhaps as } \mathbf{a} = a_x \mathbf{e}_x + a_y \mathbf{e}_y + a_z \mathbf{e}_z .\]

There are good reasons for our conventions: the ‘123’ notation is succinct; it’s easier to generalise to an arbitrary number of dimensions; and it avoids possible confusion between the index \( i \) and the unit vector \( \mathbf{i} \) and it more suitable to the tensor calculus.
The force $F$ on a particle of mass $m_1$ at the point with position vector $\mathbf{r}$ due to a particle of mass $m$ situated at the origin is given (in SI units) by

$$F(r) = -Gm m_1 \frac{\mathbf{r}}{r^3}$$

where $G = 6.67259(85) \times 10^{-11} \text{Nm}^2\text{kg}^{-2}$ is Newton’s Gravitational Constant. The magnitude of $F(r)$ is proportional to the length of the vector $|\mathbf{r}|/r^3$, which is just $|\mathbf{r}|/r^3 = r/r^3 = 1/r^2$, so we have the well-known inverse square law.

The gravitational field $\vec{G}(r)$ at $\mathbf{r}$ due to the mass $m$ at the origin is defined by

$$\vec{F}(r) \equiv m_1 \vec{G}(r) \implies \vec{G}(r) = \frac{-Gm}{r^3} \mathbf{r}$$

(2)

where the test mass $m_1$ is so small that its own gravitational field can be ignored. $\vec{G}(r)$ is a vector field.

At this point, we shall simply state that the gravitational potential due to the field $\vec{G}(r)$ is

$$\phi(r) = \frac{Gm}{r}$$

(3)

and the potential energy of the mass $m_1$ in the field is

$$V(r) = m_1 \phi(r).$$

Gravitational potential and potential energy are scalar fields. The distinction (a convention) between potential and potential energy is a common source of confusion. We shall return to these potentials later in the course, so don’t worry if you don’t know how to obtain them.

**Electrostatics:** Coulomb’s Law was also deduced experimentally; it states that the force $\vec{F}(r)$ on a particle of charge $q_1$ situated at $\mathbf{r}$ due to a particle of charge $q$ situated at the origin is given (in SI units) by

$$\vec{F}(\mathbf{r}) = \frac{q_1 q}{4\pi\epsilon_0 \mathbf{r}^3},$$

where $\epsilon_0 \equiv 10^7/(4\pi c^2) = 8.854187817 \cdots \times 10^{-12} \text{C}^2\text{N}^{-1}\text{m}^{-2}$ is called the permittivity of free space. The electric field $\vec{E}(r)$ at $\mathbf{r}$ due to the charge $q$ at the origin is defined by

$$\vec{E}(r) \equiv q_1 \vec{E}(r) \implies \vec{E}(r) = \frac{q}{4\pi\epsilon_0 \mathbf{r}^3} \mathbf{r}$$

(4)

Again the test charge $q_1$ is taken as small, so as not to disturb the electric field from $q$.

The electrostatic potential $\phi(r)$ is then (see later)

$$\phi(r) = \frac{q}{4\pi\epsilon_0 r}$$

(5)

and the potential energy of a charge $q_1$ in the electric field is $V = q_1 \phi$.

Note that electrostatics and gravitation are very similar mathematically, the only real difference being that the gravitational force between two masses is always attractive, whereas like charges repel.\(^3\)

\(^3\)In this course, a “particle” will refer to an idealised point particle, i.e. a particle of negligible size.

\(^4\)From the view of particle physics this originates in the so-called spin of the force carrier. The electromagnetic and gravitational force are transmitted by so-called photons and gravitons which have spin one and two respectively. You may think of a spin one object as an ordinary vector.
1.4 The need for vector calculus

At this point, we may ask several questions:

(i) How are equations (2) through (5) changed when the mass \( m \) or the charge \( q \) moves away from the origin?

(ii) How do the vector fields in equations (2) and (4) change when the mass \( m_1 \) or the charge \( q_1 \) moves a small distance from position \( r \) to position \( r + \delta r \) — where \( \delta r \) is very small but its direction is arbitrary? In the language of mathematics, how do we define derivatives of a vector field with respect to the position vector?

(iii) Similarly, how do the potentials change when \( r \to r + \delta r \)?

(iv) What happens when there are more than two masses or charges, or when the masses and charges have a finite size? For example, when the masses are lumpy asteroids or the charged objects are irregular lumps of metal.

You should be able to answer the first question already (exercise).

In order to address the second and third questions, we need to develop the sub-branch of mathematics known as differential vector calculus, to which we shall soon turn our attention. The answer to the fourth question requires integral vector calculus, which will come later.

In what follows, we will assume some familiarity with several variable calculus at the level of Linear Algebra and Several Variable Calculus (or the specialist Mathematics course Several Variable Calculus and Differential Equations), but these notes should be largely self-contained. We will also assume a working knowledge of vectors and bases, matrices and determinants. We shall develop the mathematics of scalar and vector fields required for third- and fourth-year courses on electromagnetism, quantum mechanics, etc, and for courses on meteorology, fluid mechanics, etc, from other schools.

1.5 Revision of vector algebra

In this section we collect together many of the results on vector algebra that will be assumed in these lectures. The index notation is introduced in subsection 1.6.

A vector \( \underline{a} \) is represented in some basis \( \{ \underline{e}_1, \underline{e}_2, \underline{e}_3 \} \) by its components,

\[
\underline{a} = a_1 \underline{e}_1 + a_2 \underline{e}_2 + a_3 \underline{e}_3
\]  

(6)

An example of a basis is given by the Cartesian basis

\[
\underline{e}_1 = (1,0,0) \quad \underline{e}_2 = (0,1,0) \quad \underline{e}_3 = (0,0,1)
\]  

(7)

which is also an orthonormal (and therefore called orthonormal basis). All quantities in this course will first be defined in Cartesian coordinates before investigating them in other types of coordinates. Often

\( \underline{a} \) is represented by \( (a_1, a_2, a_3) \).
The various ‘products’ of vectors are defined in an orthonormal basis as follows:

The **scalar product** is denoted by $\mathbf{a} \cdot \mathbf{b}$ and is the single number defined as

$$\mathbf{a} \cdot \mathbf{b} \equiv |\mathbf{a}| |\mathbf{b}| \cos \theta_{ab} = a_1 b_1 + a_2 b_2 + a_3 b_3 = \sum_{i=1}^{3} a_i b_i$$

(8)

where $\theta_{ab}$ is the angle between $\mathbf{a}$ and $\mathbf{b}$. The geometric interpretation is given in figure 1 (left).

The vector product or **cross product** is denoted by $\mathbf{a} \times \mathbf{b}$ and is defined in a right-handed orthonormal basis as the vector

$$\mathbf{a} \times \mathbf{b} \equiv |\mathbf{a}| |\mathbf{b}| \sin \theta_{ab} \mathbf{n} = \begin{vmatrix} \mathbf{e}_1 & \mathbf{e}_2 & \mathbf{e}_3 \\ a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \end{vmatrix} = (a_2 b_3 - a_3 b_2) \mathbf{e}_1 + (a_3 b_1 - a_1 b_3) \mathbf{e}_2 + (a_1 b_2 - a_2 b_1) \mathbf{e}_3 ,$$

(9)

where the vertical lines signify the determinant of the matrix, and $\mathbf{n}$ is a unit vector orthogonal to $\mathbf{a}$ and $\mathbf{b}$. The geometric interpretation is given in figure 1 (right).

Eq. (9) gives

$$\mathbf{a} \times \mathbf{b} = -\mathbf{b} \times \mathbf{a} \quad \text{and hence} \quad \mathbf{a} \times \mathbf{a} = 0 .$$

The basis vectors satisfy the cyclic properties

$$\mathbf{e}_1 \times \mathbf{e}_2 = \mathbf{e}_3 , \quad \mathbf{e}_2 \times \mathbf{e}_3 = \mathbf{e}_1 , \quad \mathbf{e}_3 \times \mathbf{e}_1 = \mathbf{e}_2 .$$

The **scalar triple product** is the single number

$$(\mathbf{a}, \mathbf{b}, \mathbf{c}) \equiv \mathbf{a} \cdot (\mathbf{b} \times \mathbf{c}) = \begin{vmatrix} a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \\ c_1 & c_2 & c_3 \end{vmatrix}$$

(10)

5To be more precise, the result of the vector product an axial- or pseudo-vector, but this subtle difference is not needed for this course. See Junior Honours *Symmetries of Classical Mechanics.*
The properties of the determinant imply that \((a, b, c) = (b, c, a) = (c, a, b)\) – cyclic permutation, and \((a, b, c) = -(b, a, c)\), etc – non-cyclic permutation.

If \(a\), \(b\) and \(c\) are three concurrent edges of a parallelepiped, its volume is \((a, b, c)\). This assertion follows from the combined geometric interpretation of the scalar product and cross-product given in figure [1].

### 1.6 Index notation - the use of tensor calculus

The index notation is an alternative way to write vector and algebra calculus. You have seen it, to some extent, in your course on **Linear Algebra and Several Variable Calculus**. It is an efficient notation that we will use for doing proofs. You are welcome and encouraged to use it yourself but you are not required to do so in order to solve exercises or the exam. More globally the formalism is known as tensor formalism and is standard in general relativity (differential geometry), quantum field theory (elementary particle physics), electromagnetism and many other fields of physics and mathematics.

An \(n\)-tensor \(t_{i_1...i_n}\) is an object \(t_{i_1...i_n}\) \((1 \leq i_1,...,i_n \leq 3\) since we are working in the three dimensional space) which takes values in the real numbers for the purpose of this course. The indices indicate how the object transforms under a rotation of the coordinate system. For example a vector is a 1-tensor and a scalar is a 0-tensor. Besides those two objects we will need the two specific 2- and 3-tensors: the Kronecker \(\delta\)-tensor (also called Kronecker symbol at times) and the Levi-Civita \(\varepsilon\)-tensor. It is noteworthy that a matrix can be interpreted as a 2-tensor.

The symbol \(\delta_{ij}\) (in words “delta i j”), where \(i\) and \(j\) can take on the values 1, 2, 3, is defined as

\[
\delta_{ij} = \begin{cases} 
1 & i = j \\
0 & \text{otherwise}
\end{cases},
\]

i.e. \(\delta_{11} = \delta_{22} = \delta_{33} = 1\) and \(\delta_{12} = \delta_{13} = \delta_{23} = \cdots = 0\).

The symbol \(\varepsilon_{ijk}\) \((i, j, k = 1..3)\) is defined as

\[
\varepsilon_{ijk} = \begin{cases} 
+1 & (i, j, k) = (1, 2, 3), (2, 3, 1) \text{ or } (3, 1, 2) \\
-1 & (i, j, k) = (2, 1, 3), (3, 2, 1) \text{ or } (1, 3, 2) \\
0 & \text{otherwise}
\end{cases},
\]

The scalar product \((8)\) is written as

\[
a \cdot b = \sum_{i,j=1}^{3} \delta_{ij}a_ib_j \quad \text{ESC} = \delta_{ij}a_ib_j = a_ib_i,
\]

where ESC stands for the Einstein summation convention. The latter implies that whenever there are two identical indices then a summation over that index is implied. Such indices

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6The word tensor originates from physics, namely from Euler stress tensor which describes the tension in a material. A stress tensor is a 2-tensor with symmetric indices.

7If one works in spaces with curvature such as in general relativity then one needs to distinguish so called co- and contravariant indices. This is far beyond the scope of this course and so is the general notion of a tensor which we just mention in passing by here. As indicated above the notions introduced below will safe us a lot of writing in the demonstrations.
are also known as dummy indices. In using tensor calculus you MUST pay attention never to introduce the same dummy index twice (e.g. \(a \cdot b \cdot c \cdot d\) can be replaced by \(a_ib_jcjd_j\) but \(a_ib_jcjd_i\) is ambiguous and therefore wrong). In order to save ourselves writing we shall use the Einstein summation convention as do many classic textbooks on physics and mathematics.

The vector product (9) can be written as

\[
\mathbf{a} \times \mathbf{b} = \sum_{i,j=1}^{3} \varepsilon_{ijk}a_ib_j \mathbf{e}_k ,
\]

where we made use of the Einstein summation convention. The \(i^{th}\) component of the vector product is then given by

\[
(a \times b)_i = \varepsilon_{ijk}a_jb_k .
\]

The latter form is the one which is most useful for manipulations. The triple scalar product (10) is written in indices notation as

\[
(a, b, c) = \varepsilon_{ijk}a_ib_jc_k .
\]

Beyond the somewhat minimal definition above we intend to give a few useful rules and properties:

- **Contraction with \(\delta_{ij}\):**

  \[
  \delta_{ij}t_{jk} = t_{ik} .
  \]

  e.g. \(\delta_{ij}a_j = a_i\) which we have already assumed in (8).

- **The orthonormality condition may be written in a rather compact form**

  \[
  \varepsilon_i \cdot \varepsilon_j = \delta_{ij} .
  \]

- **The Lagrange identity**

  \[
  \varepsilon_{ijk}\varepsilon_{kab} = \delta_{ai}\delta_{bj} - \delta_{aj}\delta_{bi} ,
  \]

  will be most useful for this course. You can convince yourself of the correctness of this equation by checking it for fixed values of the indices. For what this course is concerned we will just use it without being concerned about the deeper reason why such an identity has to exist.

- **A tensor is said to be (anti)symmetric in two indices if**

  \[
  t_{i_{n+m}} = (-)^n t_{i_n...i_{n+m}} .
  \]

  This notion is best illustrated on the two tensor used throughout this course:

  - The \(\delta_{ij}\) symbol is symmetric in its two indices. I.e.

    \[
    \delta_{ij} = \delta_{ji} .
    \]
The $\epsilon_{ijk}$ symbol is antisymmetric in its three indices. I.e.

$$\epsilon_{ijk} = -\epsilon_{jik} = -\epsilon_{kji} = \ldots \quad (21)$$

- If any two antisymmetric tensor indices take the same value then the entry is zero. For example $\epsilon_{113} = 0$.
- The contraction of two antisymmetric indices with two symmetric indices is zero. Let $s_{ij} = +s_{ji}$ and $a_{ij} = -a_{ji}$ symmetric and antisymmetric indices respectively. Hence

$$s_{ij}a_{ij} = (+s_{ji})(-a_{ji}) = (+s_{ij})(-a_{ij}) \Rightarrow s_{ij}a_{ij} = 0 \quad (22)$$

where in the second step we have relabelled the dummy indices $i$ and $j$. For example $\delta_{ij}\epsilon_{ijk} = 0$.

As one immediate application of the tensor calculus we are going to rederive the vector triple product $\mathbf{a} \times (\mathbf{b} \times \mathbf{c})$

$$\mathbf{a} \times (\mathbf{b} \times \mathbf{c}) = (\mathbf{a} \cdot \mathbf{c}) \mathbf{b} - (\mathbf{a} \cdot \mathbf{b}) \mathbf{c} \quad (23)$$

which was derived in *Linear Algebra and Several Variable Calculus* by writing out each component. Using the tensor formalism and the Lagrange identity one gets

$$(\mathbf{a} \times (\mathbf{b} \times \mathbf{c})), \quad \epsilon_{ijk}a_j(b \times c)_k \epsilon_{ijk}a_j c_{klm} b_l c_m$$

$$\delta_{il}\delta_{jm} - \delta_{jl}\delta_{im} a_j b_l c_m \left((\mathbf{a} \cdot \mathbf{c})b_i - (\mathbf{a} \cdot \mathbf{b})c_i \right) \quad (24)$$

This might appear to be complicated at first view but one gets used to it rather quickly and it is much more elegant and compact than writing out everything in terms of components. A large part of this course consists in generalising identities like (23) to the case where one of the vectors is the differential vector nabla $\nabla$. This mainly involves taking the Leibniz (product) rule of differentiation into account.

## 2 Level surfaces, gradient and directional derivative

### 2.1 Level surfaces/equipotentials of a scalar field

If $\phi(r)$ is a non-constant scalar field, then the equation $\phi(r) = c$ where $c$ is a constant, defines a level surface or equipotential of the field. Different level surfaces do not intersect, or $\phi$ would be multi-valued at the point of intersection.

Familiar examples in two dimensions, where they are level curves rather than level surfaces, are contours of constant height on a geographical map, $h(x_1, x_2) = c$. Similarly, isobars on a weather map are level curves of pressure $P(x_1, x_2) = c$.

Examples in three dimensions:
(i) Suppose that 
\[ \phi(r) = r^2 \equiv x_1^2 + x_2^2 + x_3^2 \]

The level surface \( \phi(r) = c \) is the surface of a sphere of radius \( \sqrt{c} \) centred on the origin. If we vary \( c \), we obtain a family of level surfaces or equipotentials which are concentric spheres.\(^9\)

(ii) The electrostatic potential at \( \underline{r} \) due to a point charge \( q \) situated at the point \( \underline{a} \) is

\[ \phi(\underline{r}) = \frac{q}{4\pi\epsilon_0} \frac{1}{|\underline{r} - \underline{a}|} \]

where \( |\underline{r} - \underline{a}| \) denotes the length of the vector \( \underline{r} - \underline{a} \):

\[ |\underline{r} - \underline{a}| = \sqrt{(\underline{r} - \underline{a}) \cdot (\underline{r} - \underline{a})} = \sqrt{(x_1 - a_1)^2 + (x_2 - a_2)^2 + (x_3 - a_3)^2} \]

The equipotentials or level surfaces are concentric spheres centred on the point \( \underline{a} \), as shown in the figure.

If \( \underline{a} = 0 \), the equipotentials are centered on the origin.

(iii) Let \( \phi(\underline{r}) = \overrightarrow{k} \cdot \underline{r} \).

The level surfaces are planes \( \overrightarrow{k} \cdot \underline{r} = \text{constant} \), with \( \overrightarrow{k} \) normal to the planes.

(iv) Let \( \phi(\underline{r}) = \exp(i\overrightarrow{k} \cdot \underline{r}) \), which is a complex scalar field.

Since \( \overrightarrow{k} \cdot \underline{r} = \text{constant} \) is the equation for a plane, the level surfaces are again planes.

2.2 Gradient of a scalar field

How do we describe mathematically the variation of a scalar field as a function of small changes in position?

As an example, think of a 2-d contour map of the height \( h = h(x_1, x_2) \) of a hill. \( h(x_1, x_2) \) is a scalar field. If we are on the hill and move in the \( x_1-x_2 \) plane then the change in height will depend on the direction in which we move (unless the hill is completely flat!) For example there will be a direction in which the height increases most steeply: ‘straight up the hill.’

We now introduce a formalism to describe how a scalar field \( \phi(\underline{r}) \) changes as a function of \( \underline{r} \). We begin by recalling Taylor’s theorem and the definition of partial derivatives.

\(^9\)It may seem strange to call the surface of a sphere a level surface! The point is that the scalar field is equal to a constant everywhere on the level surface, and it is in this sense that the surface is said to be level.
**Taylor’s theorem:** Recall that if \( f(x) \) is a function of a single variable \( x \), Taylor’s theorem states that \( f(x + \delta x) \) can be expanded in powers of \( \delta x \)

\[
f(x + \delta x) = f(x) + \delta x \frac{df(x)}{dx} + \frac{(\delta x)^2}{2!} \frac{d^2f(x)}{dx^2} + \cdots + \frac{(\delta x)^n}{n!} \frac{d^nf(x)}{dx^n} + \cdots
\]

If \( \delta x \) is very small, we may approximate \( f(x + \delta x) = f(x) + \delta x \frac{df(x)}{dx} + O((\delta x)^2) \)

**Partial derivatives:** If \( f(x_1, x_2, x_3) \) is a function of the three independent variables \( x_1, x_2 \) and \( x_3 \), then the partial derivative

\[
\frac{\partial f(x_1, x_2, x_3)}{\partial x_1}
\]

is obtained by differentiating \( f(x_1, x_2, x_3) \) with respect to \( x_1 \), whilst keeping \( x_2 \) and \( x_3 \) fixed. Similarly for the partial derivatives with respect to \( x_2 \) and \( x_3 \). For partial derivatives we shall introduce the following shorthand

\[
\partial_n f(a, b, ...) \equiv \frac{\partial f(a, b, ...)}{\partial a}, \quad (25)
\]

which will save us a lot of writing in the notes and in particular during the lecture\(^{10}\)

**Mathematical aside:** A scalar field \( \phi(r) = \phi(x_1, x_2, x_3) \) is said to be continuously differentiable in a region \( R \) if its partial derivatives

\[
\frac{\partial x_1 \phi(r)}{\partial x_1}, \quad \frac{\partial x_2 \phi(r)}{\partial x_2} \quad \text{and} \quad \frac{\partial x_3 \phi(r)}{\partial x_3}
\]

exist, and are continuous at every point \( r \in R \). We will generally assume scalar fields are continuously differentiable.

Let \( \phi(r) \) be a scalar field, and consider 2 nearby points: \( P \) with position vector \( r \), and \( Q \) with position vector \( r + \delta r \), where \( \delta r \) has components \((\delta x_1, \delta x_2, \delta x_3)\). Assume \( P \) and \( Q \) lie on different level surfaces as shown:

---

\(^{10}\) For the derivatives with respect to Cartesian coordinates it is common practice to introduce the even more compact notation \( \partial_i f(x_1, x_2, x_3) = \partial_{x_i} f(x_1, x_2, x_3) \equiv \frac{\partial f(x_1, x_2, x_3)}{\partial x_i} \) for \( i = 1..3 \); whenever there is no ambiguity. In this course I shall refrain from using this shorthand and remain at the level of \((25)\).
Now use Taylor’s theorem to first order in each of the 3 variables $x_1$, $x_2$ and $x_3$ to evaluate the change in $\phi$ as we move from $P$ to $Q$:

$$\delta \phi = \phi(r + \delta r) - \phi(r)$$

$$= \phi(x_1 + \delta x_1, x_2 + \delta x_2, x_3 + \delta x_3) - \phi(x_1, x_2, x_3)$$

$$= [\phi(x_1, x_2, x_3) + \partial_{x_1}\phi(r) \delta x_1 + \partial_{x_2}\phi(r) \delta x_2 + \partial_{x_3}\phi(r) \delta x_3 + O(\delta x_a \delta x_b)] - \phi(x_1, x_2, x_3)$$

$$= \partial_{x_i}\phi(r) \delta x_i + O(\delta x_a \delta x_b)$$

where we assumed that the higher order partial derivatives exist. For sufficiently small $\delta r$ we can neglect these higher order terms, and we have

$$\delta \phi = \nabla \phi \cdot \delta r$$

where the 3 quantities ($\partial_{x_1}\phi$, $\partial_{x_2}\phi$, $\partial_{x_3}\phi$) form the Cartesian components of a vector field $\nabla \phi(r)$, which we can write as

$$\nabla \phi(r) \equiv \xi_1 \partial_{x_1}\phi + \xi_2 \partial_{x_2}\phi + \xi_3 \partial_{x_3}\phi = \xi_i \partial_{x_i}\phi$$

The vector field $\nabla \phi(r)$ is called the gradient of $\phi(r)$, and is pronounced ‘grad phi’.

**Example:** Calculate the gradient of the scalar field $\phi(r) = r^2 = x_1^2 + x_2^2 + x_3^2$.

First, recall that the partial derivative $\partial_{x_1}\phi$ is the derivative of $\phi(r) = \phi(x_1, x_2, x_2)$ with respect to $x_1$, keeping $x_2$ and $x_3$ fixed, etc.

So $\partial_{x_1} x_1 = 1$ and $\partial_{x_1} x_2 = 0$. Similarly $\partial_{x_1} x_1^2 = 2x_1$ and $\partial_{x_1} x_2^2 = 0$, etc.

The first component of $\nabla r^2$ is then

$$\partial_{x_1}(x_1^2 + x_2^2 + x_3^2) = 2x_1 + 0 + 0$$

Similarly for the 2nd and 3rd components (exercise), and hence

$$\nabla r^2 = 2x_1 \xi_1 + 2x_2 \xi_2 + 2x_3 \xi_3 = 2r$$

The vector $\nabla r^2 = 2r$ points radially outwards from the origin with magnitude $2r$. The level surfaces, $r^2 = \text{constant}$, are spheres centred on the origin.

**Example:** Calculate the gradient of $\phi(r) = \sin x_1 + 2x_1 x_2^2$

$$\partial_{x_1}\phi = \cos x_1 + 2x_2^2, \quad \partial_{x_2}\phi = 0 + 4x_1 x_2, \quad \partial_{x_3}\phi = 0 + 0$$

$$\Rightarrow \quad \nabla \phi = (\cos x_1 + 2x_2^2) \xi_1 + 4x_1 x_2 \xi_2$$
Example: Calculate $\nabla (a \cdot r)$ where $a$ is a constant vector. The first component is

$$\partial x_i (a_1 x_1 + a_2 x_2 + a_3 x_3) = a_1 + 0 + 0$$

Similarly for the other two components. Hence

$$\nabla (a \cdot r) = a_1 \xi_1 + a_2 \xi_2 + a_3 \xi_3 = a$$

This is a very important result – as we shall see.

Note: A useful shorthand for partial derivatives is

$$\frac{\partial x_i}{\partial x_j} \equiv \partial_{x_j} x_i = \delta_{ij} \text{ which holds for all } i, j = 1, 2, 3.$$  (26)

2.3 Interpretation of the gradient

In deriving the expression for $\delta \phi$ above, we assumed that the points $P$ and $Q$ lie on different level surfaces. Now consider the situation where $P$ and $Q$ are nearby points on the same level surface. In this case, $\delta \phi = 0$ and so

$$\delta \phi = \nabla \phi \cdot \delta r = 0$$

In this case, the infinitesimal vector $\delta r$ lies in the level surface at $r$, and since the above equation holds for all such $\delta r$, we may deduce that

$$\nabla \phi(r) \text{ is normal (i.e. perpendicular) to the level surface at } r$$

To construct a unit normal $n(r)$ to the level surface at $r$, we divide $\nabla \phi$ by its length

$$n(r) = \frac{\nabla \phi(r)}{|\nabla \phi(r)|} \quad \text{(when } |\nabla \phi(r)| \neq 0)$$

2.4 Directional derivative

Consider the change, $\delta \phi$, produced in $\phi(r)$ by moving a distance $\delta s$ in the direction of the unit vector $\hat{s}$, so that $\delta r = \delta s \hat{s}$. Then

$$\delta \phi = \nabla \phi \cdot \delta r = (\nabla \phi) \cdot \hat{s} \delta s$$
Now divide by $\delta s$, and then let $\delta s \to 0$. The rate of change of $\phi$, with respect to distance, $s$, in the direction of $\hat{s}$, is then
\[
\frac{d\phi}{ds} = \hat{s} \cdot \nabla \phi = |\nabla \phi| \cos \theta
\] (27)
where $\theta$ is the angle between $\hat{s}$ and the normal to the level surface at $r$.

$\frac{d\phi}{ds} = \hat{s} \cdot \nabla \phi$ is called the \textit{directional derivative} of the scalar field $\phi$ in the direction of $\hat{s}$.

Equivalently, $d\phi/ds$ is the component of $\nabla \phi$ in the $\hat{s}$ direction.

From equation (27), the directional derivative has its \textit{maximum} value, $|\nabla \phi|$, when $\hat{s}$ is \textit{parallel} to $\nabla \phi$, and is \textit{zero} when $\delta s \hat{s}$ lies \textit{in} the level surface (where $\phi$ is constant.) Therefore

$\nabla \phi$ points in the direction of the \textit{maximum} rate of increase in $\phi$.

Recall that this direction is normal (perpendicular) to the level surface. A familiar example is that of contour lines on a map: the steepest direction is \textit{perpendicular} to the contour lines, \textit{i.e.} straight up the hill.

\textbf{Example:} Find the directional derivative of $\phi(r) = x_1 x_2 (x_1 + x_3)$ at the point $(1, 2, -1)$ in the direction of the unit vector $(e_1 + e_2)/\sqrt{2}$.

$\nabla \phi = (2x_1 x_2 + x_2 x_3) e_1 + x_1 (x_1 + x_3) e_2 + x_1 x_2 e_3 = 2e_1 + 2e_3$ at $(1, 2, -1)$

Thus, at this point, the directional derivative in the direction $(e_1 + e_3)/\sqrt{2}$ is
\[
\frac{1}{\sqrt{2}} (e_1 + e_2) \cdot \nabla \phi = \sqrt{2}
\]

\textbf{Physical example:} Let $T(r)$ be the temperature of the atmosphere at the point $r$. An object flies through the atmosphere with velocity $v \equiv dr/dt$. Obtain an expression for the rate of change of temperature experienced by the object.

As the object moves from $r$ to $r + \delta r$ in time $\delta t$, it experiences a change in temperature
\[
\delta T = \nabla T \cdot \delta r = \left( \nabla T \cdot \frac{\delta r}{\delta t} \right) \delta t
\]

Dividing by $\delta t$ and taking the limit $\delta t \to 0$, we obtain
\[
\frac{dT(r)}{dt} = v \cdot \nabla T(r)
\]
3 More on gradient, the operator $\nabla$

3.1 Examples of the gradient in physical laws

3.1.1 Gravitational force due to the Earth

The potential energy of a particle of mass $m$ at a modest height, $z$, above the Earth’s surface is $V = mgz$. The force due to gravity can be written as

$$
F = -\nabla V = -mg \hat{e}_z
$$

**Exercise:** Show that this last expression is correct.

Note that we choose to put a minus sign in the expression $F = -\nabla V$ so that the force acts **down** the potential energy gradient – as observed in nature!

3.1.2 More examples on $\nabla$

Before looking at Newton’s Universal Law of Gravitation, we consider two straightforward (but important) examples of gradients.

(i) Calculate $\nabla \phi$ for $\phi(r) = r = \sqrt{x_1^2 + x_2^2 + x_3^2}$

Using the chain rule, the first component of $\nabla r$ is

$$
\partial_{x_1}(x_1^2 + x_2^2 + x_3^2)^{1/2} = \frac{1}{2}(x_1^2 + x_2^2 + x_3^2)^{-1/2} 2x_1 = \frac{x_1}{r}
$$

Similarly, the second and third components are:

$$
\partial_{x_2}(x_1^2 + x_2^2 + x_3^2)^{1/2} = \frac{x_2}{r} \quad \text{and} \quad \partial_{x_3}(x_1^2 + x_2^2 + x_3^2)^{1/2} = \frac{x_3}{r}
$$

Putting these together, we get

$$
\nabla r = (x_1 \hat{e}_1 + x_2 \hat{e}_2 + x_3 \hat{e}_3) \frac{1}{r}
$$

Hence

$$
\nabla r = \frac{1}{r} \vec{r} = \hat{r}
$$

We conclude that for $\phi(r) = r = \text{“the length of the position vector at position } r\text{”}$, the gradient $\nabla \phi$ is just the unit vector $\hat{r}$ which points **radially outwards** from the origin. It has the same **magnitude** everywhere, but its **direction** is normal to the level surfaces which are spheres centered on the origin.

(ii) Calculate $\nabla \phi$ for $\phi(r) = \frac{1}{r} = \frac{1}{\sqrt{x_1^2 + x_2^2 + x_3^2}}$

The first component of $\nabla \left(\frac{1}{r}\right)$ is

$$
\partial_{x_1}(x_1^2 + x_2^2 + x_3^2)^{-1/2} = -\frac{1}{2}(x_1^2 + x_2^2 + x_3^2)^{-3/2} 2x_1 = -\frac{x_1}{r^3}
$$

---

11 We use the notations $\hat{r}$ and $\vec{e}_r$ for a unit vector in the direction of the position $r$. 
Similarly for the second and third components:
\[ \frac{\partial}{\partial x_2} (x_1^2 + x_2^2 + x_3^2)^{-1/2} = -\frac{x_2}{r^3} \quad \text{and} \quad \frac{\partial}{\partial x_3} (x_1^2 + x_2^2 + x_3^2)^{-1/2} = -\frac{x_3}{r^3} \]

Hence
\[ \nabla \left( \frac{1}{r} \right) = -\frac{1}{r^3} \hat{r} = -\frac{1}{r^2} \hat{r} \]

For \( \phi(r) = 1/r = \text{“the inverse of the length of the position vector”} \), the gradient \( \nabla \phi \) points radially inwards towards the origin, with magnitude \( 1/r^2 \). The level surfaces are again spheres centered on the origin.

### 3.1.3 Newton’s Law of Gravitation:

Armed with this result, we now return to the force \( \mathbf{F}(r) \) on a mass \( m_1 \) at \( r \) due to a mass \( m \) at the origin:
\[ \mathbf{F}(r) = -\frac{Gm_1m}{r^3} \hat{r} \]

Using the result \( \nabla (1/r) = -r/r^3 \) derived in (ii) above, we may write this as
\[ \mathbf{F}(r) = -\nabla V(r) \]

where the gravitational potential energy is \( V(r) = -Gm_1m/r \).

Similarly, we can write the gravitational field
\[ \mathbf{G}(r) = -\frac{Gm}{r^3} \hat{r} = -\nabla \phi(r) \]

where the gravitational potential is \( \phi(r) = -Gm/r \).

Similar results hold for the electrostatic force and its associated electric field.

We shall show later that there is a very general class of vector fields that can be written as the gradient of a scalar field known as a scalar potential.

### 3.2 Identities for gradients

Thus far, we have calculated gradients in a rather tedious and repetitive fashion – we worked out each example from scratch, and we calculated each of the three components of \( \nabla \phi \) in each example. This was deliberate. . . . We shall now see what we can gain by becoming a little more sophisticated.

We shall derive several identities which hold for the gradient of any scalar field, and which we may use to speed up the evaluation of the gradient of more complicated scalar fields.

If \( \phi(r) \) and \( \psi(r) \) are scalar fields, then:

(i) **Distributive law**
\[ \nabla (\phi + \psi) = \nabla \phi + \nabla \psi \]
**Proof:** It follows from the partial derivatives being distributive. More explicitly for the $i^{\text{th}}$ component,

$$\left( \nabla (\phi + \psi) \right)_i \equiv \partial_{x_i}(\phi + \psi) = \partial_{x_i}(\phi) + \partial_{x_i}(\psi) = (\nabla \phi)_i + (\nabla \psi)_i .$$

If you’re not comfortable with evaluating the $i^{\text{th}}$ component, then put $i = 1$ to recover the expression in the first derivation. Then put $i = 2$, then $i = 3$. After a while, it seems perfectly natural to consider the $i^{\text{th}}$ component from the beginning.

(ii) **Product rule**

$$\nabla (\phi \psi) = (\nabla \phi) \psi + \phi (\nabla \psi)$$

**Proof:** Follows from the Leibniz rule ($\partial (fg) = (\partial f)g + f(\partial g)$) of ordinary calculus. More explicitly for the $i^{\text{th}}$ component,

$$\left( \nabla (\phi \psi) \right)_i = \partial_{x_i}(\phi \psi) = \partial_{x_i}(\phi) \psi + \partial_{x_i}(\psi) \phi = (\nabla \phi)_i \psi + \phi (\nabla \psi)_i .$$

(iii) **Chain rule:** If $F(\phi(r))$ is a scalar field, then

$$\nabla F(\phi) = \partial_{\phi} F(\phi) \nabla \phi$$

**Proof:** Follows from the chain rule of ordinary calculus. The $i^{\text{th}}$ component is

$$\left( \nabla F(\phi) \right)_i = \partial_{x_i} F(\phi) = \partial_{\phi} F(\phi) \partial_{x_i} \phi = \partial_{\phi} F(\phi) (\nabla \phi)_i$$

where we used the ordinary chain rule to get the second-last expression.

**Example of chain rule:** If $\phi(r) = r = \sqrt{x_1^2 + x_2^2 + x_3^2}$ we can use result (i) from section (3.1.2), $\nabla r = r/r$, to get

$$\nabla F(r) = \partial_r F(r) \nabla r = \frac{F'(r)}{r} r$$

where, as usual, $F'(r)$ denotes the derivative of the function $F(r)$ with respect to $r$.

If $F(\phi(r)) = r^n$, we have $\phi(r) = r$ as in the previous example, and we obtain the important result

$$\nabla (r^n) = \partial_r r^n (\nabla r) = \left( n r^{n-1} \right) \frac{1}{r} r = (n r^{n-2}) r$$

**NB:** This is a much quicker way of evaluating $\nabla (r^n)$ than writing out the components – as you were asked to do in Homework Problem 1.8.

Setting $n = -1$ gives

$$\nabla \left( \frac{1}{r} \right) = -\frac{r}{r^3}$$

which reproduces result (ii) from section (3.1.2).
Example: Calculate $\nabla \phi$ when $\phi(r) = r^n (a \cdot r)^m$.

\[
\nabla \{r^n (a \cdot r)^m\} = (\nabla r^n) (a \cdot r)^m + r^n \{\nabla (a \cdot r)^m\} \\
= (\nabla r^n) (a \cdot r)^m + r^n m (a \cdot r)^{m-1} \{\nabla (a \cdot r)\} \\
= n r^{n-2} (a \cdot r)^m r + m r^n (a \cdot r)^{m-1} a
\]

where, in the last line, we used $\nabla r^n = n r^{n-2} r$ and $\nabla (a \cdot r) = a$.

This example demonstrates the “toolkit” approach to evaluating gradients of complicated expressions. Here, we combined the product rule for grad with known gradients of “simple” scalar fields $(a \cdot r)$ and $r^n$, which you should know, and be able to work out from first principles.

3.3 The operator $\text{del}$

We can think of $\nabla$ as a vector operator, called del, which acts on the scalar field $\phi(r)$ to produce the vector field $\nabla \phi(r)$, which is pronounced grad phi. In Cartesian coordinates

\[
\nabla \equiv \hat{\epsilon}_1 \frac{\partial}{\partial x_1} + \hat{\epsilon}_2 \frac{\partial}{\partial x_2} + \hat{\epsilon}_3 \frac{\partial}{\partial x_3} \equiv \hat{\epsilon}_i \partial_{x_i}
\]

where we have used the standard Einstein summation convention and the shorthand (25).

We call $\nabla$ an operator since it operates on something to its right. It is a vector operator because it produces a vector field when it operates on a scalar field. The derivatives transform in the same way as coordinates\(^{12}\) which are vectors and are therefore themselves vectors.

We have seen how $\nabla$ acts on a scalar field to produce a vector field. We can take products of the vector operator $\nabla$ with other vector quantities to produce new operators and fields in the same way we can take scalar and vector products of two vectors.

For example, the directional derivative of $\phi$ in the direction $\hat{s}$, is given by $\hat{s} \cdot \nabla \phi$.

More generally, if $a = a_i \hat{e}_i$ is any vector (or any vector field), we can interpret $a \cdot \nabla = a_i \partial_{x_i}$ as a scalar operator.

The operator $a \cdot \nabla$ acts on a scalar field to its right to produce another scalar field

\[
(a \cdot \nabla) \phi = a_i \partial_{x_i} \phi = a \cdot (\nabla \phi)
\]

We can also act with this operator on a vector field $b(r)$ to get another vector field,

\[
(a \cdot \nabla) b = \hat{e}_1 (a \cdot \nabla) b_1 + \hat{e}_2 (a \cdot \nabla) b_2 + \hat{e}_3 (a \cdot \nabla) b_3
\]

The alternative expression $a \cdot (\nabla b)$ is undefined because $\nabla b$ doesn’t make sense – just like $a \cdot b$ doesn’t make sense.

(For this reason, the parentheses are sometimes omitted, and $a \cdot \nabla b$ is taken to mean $(a \cdot \nabla) b$, but I wouldn’t recommend doing this because it often leads to errors.)

\(^{12}\)On a curved space this statement has to be refined. The derivatives transform contravariant and the coordinates transform covariant like bra and ket vectors in quantum mechanics. This can also be deduced from the orthogonality relation (26). These remarks go beyond the scope of this course and are given to you to encourage to broaden your horizon in further studies.
Attention to ordering:

• We have written the orthonormal vectors $e_i$ to the left of the derivatives. For Cartesian coordinates $\partial_x e_i = 0$ and it does therefore not matter whether we write the orthonormal vectors $e_i$ to the left or to the right of the partial derivatives. For cylindrical and spherical coordinates this is not the case and this is why it is better to be strict, i.e. $e_i$ to the left of the derivatives in (28).

• Great care is required with the order in products since, in general, products involving operators are not commutative. For example, if $a(r)$ is a vector field

$$a \cdot \nabla \neq \nabla \cdot a$$

The quantity $a \cdot \nabla$ is a scalar differential operator, whereas $\nabla \cdot a$ is a scalar field called the divergence of $a$—see later.

Example: If $a(r)$ is a vector field, show that $(a \cdot \nabla) r = a$. This is left as an (important) tutorial exercise for the student.

Examples: In section (2.2) we showed that, for small displacements $\delta r$, we have

$$\phi(r + \delta r) = \phi(r) + (\nabla \phi) \cdot \delta r + O((\delta r)^2) = \phi(r) + \delta r \cdot \nabla \phi + O((\delta r)^2)$$

(29)

If we set $\delta r = a$, where $a$ is an arbitrary (but small) constant vector, we have

$$\phi(r + a) = \phi(r) + a \cdot \nabla \phi(r) + O(a^2)$$

(30)

As we shall see, this expression is is very useful when $\phi(r)$ is the electrostatic potential.

We can expand a vector field about some point $r$ in exactly the same way. For example, let $E(r)$ be the electric field at $r$. If we take the result of equation (29), and simply replace $\phi(r)$ by (say) the first component of the electric field $E_1(r)$, we obtain:

$$E_1(r + \delta r) = E_1(r) + (\delta r \cdot \nabla) E_1(r) + O((\delta r)^2)$$

Doing the same for the second and third components of the electric field, and again setting $\delta r = a$, gives

$$E(r + a) = E(r) + (a \cdot \nabla) E(r) + O(a^2)$$

(31)

Equations (30) and (31) arise in the study of dipoles in electrostatics. Incidentally (31) and (29) answer questions (ii) and (iii) posed in section ??.

3.4 Equations of points, lines and planes

(This section may be revision.)

13 We can write this equation as $\phi(r + a) = \phi(r) + a \frac{d\phi}{dr} + O(a^2)$ where $\frac{d\phi}{dr}$ is the directional derivative in the direction of the unit vector $\hat{a}$, and $a = |a|$, but this is rarely done in practice - possibly because its apparent simplicity leads to errors!
3.4.1 The position vector

The equation for a point is simply \( \mathbf{r} = \mathbf{a} \) where \( \mathbf{a} \) is some constant vector.

3.4.2 The equation of a line

Suppose that the point \( P \) lies on a line which passes through a point \( A \), which has a position vector \( \mathbf{a} \) with respect to an origin \( O \). Let \( P \) have position vector \( \mathbf{r} \) relative to \( O \), and let \( \mathbf{u} \) be a vector through the origin in a direction parallel to the line.

From the figure, \( \mathbf{r} = \mathbf{OA} + \mathbf{AP} \), which, for some \( \lambda \), we may write as

\[
\mathbf{r} = \mathbf{a} + \lambda \mathbf{u}
\]

This is the explicit or parametric equation of the line, i.e. as we vary the parameter \( \lambda \) from \(-\infty\) to \(\infty\), \( \mathbf{r} \) describes all points on the line.

Rearranging and using \( \mathbf{u} \times \mathbf{u} = 0 \), we can also write this as

\[
(r - a) \times \mathbf{u} = 0
\]

or

\[
\boxed{\mathbf{r} \times \mathbf{u} = c}
\]

where \( c = \mathbf{a} \times \mathbf{u} \) is normal to the plane containing the line and the origin.

Notes

(i) \( \mathbf{r} \times \mathbf{u} = c \) is an implicit equation for a line parallel to the vector \( \mathbf{u} \).

(ii) \( \mathbf{r} \times \mathbf{u} = 0 \) is the equation of a line through the origin (\( \mathbf{a} \) and \( \mathbf{u} \) are parallel in this case.)

3.4.3 The equation of a plane

\[
\mathbf{r} \text{ is the position vector of an arbitrary point } P \text{ on the plane;}
\]
\[
\mathbf{a} \text{ is the position vector of a fixed point } A \text{ in the plane;}
\]
\[
\mathbf{u} \text{ and } \mathbf{v} \text{ are any vectors parallel to the plane, but non-collinear: } \mathbf{u} \times \mathbf{v} \neq 0.
\]

We can express the vector \( \mathbf{AP} \) in terms of \( \mathbf{u} \) and \( \mathbf{v} \), so that

\[
\mathbf{r} = \mathbf{a} + \mathbf{AP} = \mathbf{a} + \lambda \mathbf{u} + \mu \mathbf{v}
\]
for some $\lambda$ and $\mu$. This is the parametric equation of the plane.

Now define the unit normal to the plane

$$n = \frac{u \times v}{|u \times v|}.$$ 

Clearly $u \cdot n = v \cdot n = 0$. Rearranging equation (32), and using these results, we find the implicit equation for the plane

$$(r - a) \cdot n = 0.$$ 

Alternatively, we can write this as

$$r \cdot n = p$$

where $p = a \cdot n$ is the perpendicular distance of the plane from the origin.

This is a very important equation which you must be able to recognise.

**Note:** $r \cdot n = 0$ is the equation for a plane through the origin (with unit normal $n$).

**Example:** Recall that $\nabla (a \cdot r) = a$ when $a$ is a constant vector. Thus the level surfaces of the scalar field $a \cdot r$ are planes orthogonal to $a$.

## 4 Div, curl and the Laplacian

We now combine the vector operator $\nabla$ (del) with a vector field to define two new operations div and curl. Then we define the Laplacian.

### 4.1 Divergence

The *divergence* of a vector field $a$ (pronounced ‘div $a$’) is obtained by taking the scalar product of the vector with the del operator as follows

$$\text{div } a(r) = \nabla \cdot a(r) = \partial_x a_1 (\xi_i \cdot \xi_j) = \partial_x a_j \delta_{ij} = \partial_x a_i$$

where we have assumed that $\partial_x \xi_j = 0$ which is true for Cartesian coordinates in particular. This is no longer true for cylindrical and spherical coordinates and complicates the derivation of the divergence, Laplacian. We shall investigate this viewpoint in the very last subsection 15.6 which will complicate matters slightly. For later reference we rewrite the divergence in Cartesian coordinates using the standard Einstein summation convention and the shorthand (25):

$$\nabla \cdot a \equiv \frac{\partial a_1}{\partial x_1} + \frac{\partial a_2}{\partial x_2} + \frac{\partial a_3}{\partial x_3} = \partial_x a_i$$

$\nabla \cdot a$ is a single number at each point $r$, and it’s invariant under a change of coordinates (basis) and therefore a scalar or *scalar field*. 

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Example: If \( \text{a}(r) = r \) then \[ \nabla \cdot r = 3 \] which is a very important and useful result.

Explicitly:
\[
\nabla \cdot r = \partial_r x_i = 1 + 1 + 1 = 3
\]

(35)

Example: Let \( \text{a}(r) = x_1^2 x_3 e_1 - 2 x_2^3 x_3^2 e_2 + x_1 x_2^2 x_3 e_3 \)
\[
\nabla \cdot \text{a} = \partial_{x_1} (x_1^2 x_3) - \partial_{x_2} (2 x_2^3 x_3^2) + \partial_{x_3} (x_1 x_2^2 x_3)
\]
\[
= 2 x_1 x_3 - 6 x_2^2 x_3^2 + x_1 x_2^2
\]

Then, at the point \((1, 1, 1)\) for instance, \( \nabla \cdot \text{a} = 2 - 6 + 1 = -3 \).

4.2 Curl

We define the curl of a vector field \( \text{a}(r) \) by
\[
\text{curl } \text{a}(r) \equiv \nabla \times \text{a}(r)
\]
\( \nabla \times \text{a} \) is a vector field\(^{14} \)

We can write the curl in determinant form, as we did for the ordinary vector product:
\[
\nabla \times \text{a} = \begin{vmatrix}
\xi_1 & \xi_2 & \xi_3 \\
\partial_{x_1} & \partial_{x_2} & \partial_{x_3} \\
a_1 & a_2 & a_3
\end{vmatrix} = \epsilon_{ijk} \epsilon_i \partial_{x_j} a_k
\]

(36)

where we have given the definition using the determinant as well as the one using the \( \epsilon \)-tensor. You may or may not find the former form more convenient for actual computation whereas the latter is certainly more efficient for proofs. More explicitly, the components of \( \nabla \times \text{a} \) are:
\[
(\nabla \times \text{a})_1 = \partial_{x_2} a_3 - \partial_{x_3} a_2 \\
(\nabla \times \text{a})_2 = \partial_{x_3} a_1 - \partial_{x_1} a_3 \\
(\nabla \times \text{a})_3 = \partial_{x_1} a_2 - \partial_{x_2} a_1
\]

Example: If \( \text{a}(r) = r \) then
\[
\nabla \times r = \epsilon_{ijk} \epsilon_i \partial_{x_j} x_k \epsilon_{ijk} \epsilon_j \delta_{jk} 0
\]

(37)

which is another very important and useful result.

\(^{14}\)More precisely, \( \nabla \times \text{a} \) is a pseudo-vector field, if \( \text{a} \) is a vector field.
Example: Compute the curl of \( \mathbf{a} = x_1^2x_2 \mathbf{e}_1 + x_2^2x_1 \mathbf{e}_2 + x_1x_2x_3 \mathbf{e}_3 \):

\[
\nabla \times \mathbf{a} = \epsilon_{ijk} \epsilon_i \partial_j a_k = \begin{vmatrix}
\mathbf{e}_1 & \mathbf{e}_2 & \mathbf{e}_3 \\
\partial_{x_1} & \partial_{x_2} & \partial_{x_3} \\
x_1^2x_2 & x_2^2x_1 & x_1x_2x_3
\end{vmatrix}
= (x_1x_3 - 0) \mathbf{e}_1 - (x_2x_3 - 0) \mathbf{e}_2 + (x_2^2 - x_1^2) \mathbf{e}_3
\]

We have given both forms for the curl and leaving it up to the student to choose which one to use. Maybe you should use both and convince yourself that they indeed give the same result.

4.3 The Laplacian operator \( \nabla^2 \)

Recall that we interpreted the operator \( \nabla \) as a vector differential operator (28). Therefore one can form a scalar differential operator by taking the scalar product.

\[
\nabla \cdot \nabla = \epsilon_i \partial_i \cdot \epsilon_j \partial_j = (\epsilon_i \cdot \epsilon_j) \partial_i \partial_j = \delta_{ij} \partial_i \partial_j = \sum_{i=1}^{3} \partial^2_{x_i} \equiv \nabla^2.
\]

This defines the so-called \textbf{Laplacian operator} (sometimes called the \textbf{Laplacian})

\[
\nabla^2 = \frac{\partial^2}{\partial x_1^2} + \frac{\partial^2}{\partial x_2^2} + \frac{\partial^2}{\partial x_3^2} = \sum_{i=1}^{3} \partial^2_{x_i},
\]

where we have used the shorthand (25) but have not dared to imply the Einstein summation convention in the last equality because of an anticipated degree of ambiguity. In Cartesian coordinates applied to a scalar function \( \phi(r) \):

\[
\nabla^2 \phi = \frac{\partial^2 \phi}{\partial x_1^2} + \frac{\partial^2 \phi}{\partial x_2^2} + \frac{\partial^2 \phi}{\partial x_3^2} = \sum_{i=1}^{3} \partial^2_{x_i} \phi
\]

Remarks

• The Laplace operator (39), operates on functions to its \textit{right}. In (40) the Laplace operator acts on a scalar. The next item clarifies how this generalises.

• This Laplacian, being a scalar, can act on any tensor without changing its degree. For example for a vector one (just) replaces \( \phi \rightarrow \mathbf{v} \) in Eq. (40). In this course it will only act on scalars and vectors.

• From Eq. (38) it ought to be understood that \( \nabla \cdot (\nabla \phi) = \nabla^2 \phi. \)

• In Eq. (38) in the second equality we have assumed that \( \partial_{x_i} \mathbf{e}_j = 0 \). This assumption is valid for Cartesian coordinates but not for cylindrical or spherical coordinates which will complicate matters slightly.
Example: Evaluate $\nabla^2 r^2$ and $\nabla^2 (a \cdot r)$ where $a$ is a constant vector.

\[
\nabla^2 r^2 = \nabla \cdot (\nabla r^2) = \nabla \cdot (2r) = 2 \times 3 = 6.
\]

\[
\nabla^2 (a \cdot r) = \nabla \cdot (\nabla (a \cdot r)) = \nabla \cdot a = 0 \quad \text{(because $a$ is a constant)}
\]

where we used the basic results $\nabla r^2 = 2r$, $\nabla \cdot r = 3$ and $\nabla (a \cdot r) = a$, derived previously. We could have performed these calculations using components, but that would take longer.

Example: Evaluate $\nabla^2 r^n$.

\[
\nabla^2 r^n = \nabla \cdot (\nabla r^n) = \nabla \cdot (nr^{n-2}r) = ?
\]

We could evaluate the last expression by writing out its components, but we shall be patient and work out a better way of doing it in the next section.

5 Vector operator identities

There are many identities involving div, grad, curl and the Laplacian. It is not necessary to know all of these, but you should know and be able to use the product and chain rules for gradients (see Section 3.2), together with the product laws for div and curl given below. Fortunately, most of these are almost almost obvious . . .

You should be familiar with the rest and to be able to derive and use them when necessary.

It is also extremely useful to know and be able to derive the results for elementary quantities such as $\nabla r$, $\nabla r^n$, $\nabla \cdot r$, $\nabla \times r$, $(a \cdot \nabla)r$, $\nabla (a \cdot r)$, $\nabla \times (a \times r)$ and $\nabla \cdot (a \times r)$ where $a$ is a constant vector. This is similar to learning and understanding multiplication tables, or knowing the derivative of elementary functions such as $\sin x$.

We shall use vector identities and these elementary results to enable us to evaluate div, grad and curl of complicated expression with a minimum of effort – which is surely a good thing!

Most importantly you should be at ease with div, grad and curl. This only comes through practice and deriving the various identities gives you just that.

In what follows $\phi(r)$, $a(r)$ and $b(r)$ are all continuously-differentiable scalar and vector fields.

5.1 Distributive laws

A. $\nabla \cdot (a + b) = \nabla \cdot a + \nabla \cdot b$

B. $\nabla \times (a + b) = \nabla \times a + \nabla \times b$

The proofs of these are straightforward using components and they follow from the fact that div and curl are linear operations.
5.2 Product laws: one scalar field and one vector field

The results of taking the \( \text{div} \) or \( \text{curl} \) of products of vector and scalar fields are the most useful. They are predictable but they need a little care.

C. \[ \nabla \cdot (\phi a) = (\nabla \phi) \cdot a + \phi (\nabla \cdot a) \]
D. \[ \nabla \times (\phi a) = (\nabla \phi) \times a + \phi (\nabla \times a) \]

Proof of (C): [Tutorial exercise]
Proof of (D): Follows from Leibniz rule. In terms of the \( i \)-th component
\[
(\nabla \times (\phi a))_i = \varepsilon_{ijk} \partial_{x_j} (\phi a_k) = \varepsilon_{ijk} ((\partial_{x_j} \phi) a_k + \phi \partial_{x_j} a_k) \\
= \varepsilon_{ijk} (\partial_{x_j} \phi) a_k + \phi \varepsilon_{ijk} \partial_{x_j} a_k = ((\nabla \phi) \times a)_i + (\phi (\nabla \times a))_i .
\]

5.3 Product laws: two vector fields

The following identities are extremely useful but less obvious.

E. \[ \nabla \cdot (a \times b) = b \cdot (\nabla \times a) - a \cdot (\nabla \times b) \]
F. \[ \nabla \times (a \times b) = (\nabla \cdot b) a + (b \cdot \nabla) a - (\nabla \cdot a) b - (a \cdot \nabla) b \]
G. \[ \nabla (a \cdot b) = (a \cdot \nabla) b + (b \cdot \nabla) a + a \times (\nabla \times b) + b \times (\nabla \times a) \]

Proof of (E): Follows from the Leibniz rule
\[
\nabla \cdot (a \times b) = \varepsilon_{ijk} \partial_{x_j} a_j b_k = \varepsilon_{ijk} (\partial_{x_j} a_j) b_k + \varepsilon_{ijk} a_j (\partial_{x_i} b_k) = b \cdot (\nabla \times a) - a \cdot (\nabla \times b) .
\]

Proof of (F): Using Leibniz rule and the double cross product formula (6) can be derived.
\[
\nabla \times (a \times b) = \nabla \cdot b a - \nabla \cdot a b = (\nabla \cdot b) a + (b \cdot \nabla) a - (\nabla \cdot a) b - (a \cdot \nabla) b .
\]

After the first equality the nabla operator \( \nabla \) is meant to act on both \( a \) and \( b \) on the left. In this case it did not seem necessary to go to index notation but you could do that yourself as an additional exercise if you wanted.

Proof of (G): [Tutorial exercise]

Identities (E), (F) & (G) are very useful in explicit calculations. You don’t need to memorise them in gory detail, but you should know they exist, and you should understand the derivations. Other results involving one \( \nabla \) can be derived similarly.
Example: Show that $\nabla \cdot (r^{-3}r) = 0$, for $r \neq 0$, where as usual $r = |r|$.

Solution: Using identity (3), we have

$$\nabla \cdot (r^{-3}r) = (\nabla r^{-3}) \cdot r + r^{-3} (\nabla \cdot r)$$

$$= \left( -\frac{3}{r^5} r \right) \cdot r + \frac{3}{r^3} = -\frac{3}{r^5} r^2 + \frac{3}{r^3} = 0 \quad (\text{except at } r = 0)$$

where we used the results $\nabla r^n = nr^{n-2}r$ and $\nabla \cdot r = 3$ that we derived previously.

5.4 Identities involving two $\nabla$s

H. $\nabla \times (\nabla \phi) = 0$ \hspace{1cm} (curl grad $\phi$ is always zero)

I. $\nabla \cdot (\nabla \times a) = 0$ \hspace{1cm} (div curl $a$ is always zero)

J. $\nabla \times (\nabla \times a) = \nabla (\nabla \cdot a) - \nabla^2 a$

Proofs of (H) and (I) are readily obtained directly in index notation. You should know (H) and (I), and knowing (J) is really, really useful!\(^{15}\)

Proof of (H):

$$\nabla \times (\nabla \phi) = \varepsilon_{ijk} \partial_{x_j} \partial_{x_k} \phi \quad \text{(22)}$$

where in the last step we used the result that partial derivatives and therefore $\partial_{x_j} \partial_{x_k}$ is a symmetric tensor.

Proof of (I): This is (very) similar to (H) and is left as a tutorial exercise.

Proof of (J): This is left as an exercise. [tutorial].

Identity (J) is used in curvilinear coordinate systems to define the action of the Laplacian on a vector field as

$$\nabla^2 a \equiv \nabla \left( \nabla \cdot a \right) - \nabla \times \left( \nabla \times a \right)$$

We shall do this at the end the course.

Identities (H), (I) and (J) are extremely important: they’re used a lot in electromagnetism; fluid mechanics; elasticity theory and other field theories.

Example: Scalar fields which depend only on $r$.

In many physical examples the scalar field $\phi$ depends only on the length of the position vector $r = |r|$, and we have

$$\nabla^2 \phi(r) = \phi''(r) + \frac{2\phi'(r)}{r} = \frac{1}{r^2} \left( r^2 \phi'(r) \right)'$$

\(^{15}\) A mnemonic for the Laplacian acting on a vector field is GDMCC – Grad-Div Minus Curl-Curl, pronounced g(i)d(u)mcc!
where the prime denotes differentiation with respect to \( r \). Proof of this relation utilises the chain rule

\[ \nabla \phi(r) = \frac{d\phi(r)}{dr} \nabla r = \phi'(r) \frac{r}{r} \]

and is left to the tutorial.

**Example:** Evaluate \( \nabla \times \{ (c \times r) / r^3 \} \) where \( c \) is a constant vector.

Start with the product rule \( \nabla \times (\phi a) = (\nabla \phi) \times a + \phi (\nabla \times a) \) where \( \phi = 1/r^3 \) and \( a = c \times r \).

\[
\nabla \times \left( \frac{c \times r}{r^3} \right) = \nabla \left( \frac{1}{r^3} \right) \times (c \times r) + \frac{1}{r^3} \nabla \times (c \times r)
\]

\[
= -3 \frac{r}{r^3} \times (c \times r) + \frac{1}{r^3} 2c
\]

\[
= -3 \frac{1}{r^3} (r^2 c - (r \cdot c) r) + \frac{2}{r^3} c = \frac{3}{r^5} \frac{(r \cdot c) r - c}{r^3}
\]

where we used \( \nabla r^n = nr^{n-2}r \) with \( n = -3 \), and \( \nabla \times (c \times r) = 2c \).

**Physical example:** The vector potential of an electric dipole has the form \( \vec{A}(r) = (c \times r) / r^3 \).

### 5.5 Summary

**Elementary results:** We can calculate \( \text{div}, \text{grad}, \text{curl} \) and the \( \text{Laplacian} \) of many of the scalar and vector fields that occur in physics using the following elementary results, which you must know and be able to derive:

\[
\begin{align*}
\nabla r &= \hat{r} \\
\nabla r^n &= nr^{n-2}r \\
\nabla \cdot r &= 3 \\
\nabla \times r &= 0 \\
\n\nabla (a \cdot r) &= a \\
\n(a \cdot \nabla) r &= a \\
\n\nabla \times (a \times r) &= 2a \\
\n\nabla \cdot (a \times r) &= 0
\end{align*}
\]

where \( \vec{r} \) is the position vector, \( r = |\vec{r}| \) is its magnitude (length), and \( \vec{a} \) is a constant vector.

The identity \( (a \cdot \nabla) r = a \) holds also for vector fields \( \vec{a}(r) \), because no derivatives act on \( \vec{a} \).

**Identities for scalar fields:** You must know, and be able to derive, all of the following for scalar fields \( \phi(\vec{r}) \) and \( \psi(\vec{r}) \):

(i) Distributive law: \( \nabla (\phi + \psi) = \nabla \phi + \nabla \psi \)

(ii) Product rule: \( \nabla (\phi \psi) = (\nabla \phi) \psi + \phi (\nabla \psi) \)

(iii) Chain rule: If \( F(\phi(\vec{r})) \) is a scalar field, then \( \nabla F(\phi) = \frac{dF(\phi)}{d\phi} \nabla \phi(\vec{r}) \)

**Important example:** \( \nabla f(r) = (f'(r)/r) r \)
Identities for vector fields: For vector fields $a(r)$ and $b(r)$, and scalar fields $\phi(r)$:

A. $\nabla \cdot (a + b) = \nabla \cdot a + \nabla \cdot b$

B. $\nabla \times (a + b) = \nabla \times a + \nabla \times b$

C. $\nabla \cdot (\phi a) = (\nabla \phi) \cdot a + \phi (\nabla \cdot a)$

D. $\nabla \times (\phi a) = (\nabla \phi) \times a + \phi (\nabla \times a)$

E. $\nabla \cdot (a \times b) = b \cdot (\nabla \times a) - a \cdot (\nabla \times b)$

F. $\nabla \times (a \times b) = (\nabla \cdot b) a + (b \cdot \nabla) a - (\nabla \cdot a) b - (a \cdot \nabla) b$

G. $\nabla \cdot (a \cdot b) = (a \cdot \nabla) b + (b \cdot \nabla) a + a \times (\nabla \times b) + b \times (\nabla \times a)$

H. $\nabla \times (\nabla \phi) = 0$ \hspace{1cm} (curl grad $\phi$ is always zero)

I. $\nabla \cdot (\nabla \times a) = 0$ \hspace{1cm} (div curl $a$ is always zero)

J. $\nabla \times (\nabla \times a) = \nabla (\nabla \cdot a) - \nabla^2 a$

You must know and be able to use & derive identities A-D and H-I. You must be familiar with the others, and be able to use them in calculations.
6 Geometrical/physical interpretation of \( \text{div} \) and \( \text{curl} \)

A full interpretation of the divergence and curl of a vector field is best left until after we have studied the divergence theorem and Stokes’ theorem respectively. However, we can gain some intuitive understanding by looking at simple examples where div and/or curl vanish.

**Example:** Consider the radial field \( \vec{a}(r) = r \). We have just shown that \( \nabla \cdot r = 3 \) and \( \nabla \times r = 0 \). We may sketch the vector field \( \vec{a}(r) \) by drawing vectors of the appropriate direction and magnitude at selected points. These give the tangents of ‘flow lines’. Roughly speaking, in this example the divergence is positive because bigger arrows come out of any point than go into it. So the field ‘diverges’. (Once the concept of flux of a vector field is understood this will make more sense.)

**Angular velocity:** Consider a point in a rigid body rotating with angular velocity \( \vec{\omega} = \omega \hat{\omega} \). The magnitude \( \omega = |\vec{\omega}| \) is the angular speed of rotation measured in radians per second, and the unit vector \( \hat{\omega} \) lies along the axis of rotation. Let the position vector of the point with respect to an origin \( O \) on the axis of rotation be \( \vec{r} \).

You should convince yourself that the velocity of the point is given by \( \vec{v} = \vec{\omega} \times \vec{r} \) by checking that this gives the right direction for \( \vec{v} \); that it is perpendicular to the plane of \( \vec{\omega} \) and \( \vec{r} \); that the magnitude \( |\vec{v}| = \omega r \sin \theta = \omega \rho \), where \( \rho \) is the radius of the circle in which the point is travelling.

**Example:** Consider the field \( \vec{v}(r) = \vec{\omega} \times \vec{r} \) where \( \vec{\omega} \) is a constant vector. One can think of \( \vec{v}(r) \) as the velocity of a point in a rigid rotating body. The sketch shows a cross-section of the field \( \vec{v}(r) \) with \( \vec{\omega} \) chosen to point out of the page.

To evaluate \( \nabla \times (\vec{\omega} \times \vec{r}) \) and \( \nabla \cdot (\vec{\omega} \times \vec{r}) \), we recall that \( \nabla \times r = 0 \) \((37)\) and \( \nabla \cdot r = 3 \) \((35)\) and the identities \((5)\) and \((6)\) from section 5.5. One gets

\[
\nabla \times (\vec{\omega} \times \vec{r}) = \nabla \cdot (\vec{\omega} \times r) = 3 \vec{\omega} - \omega_l \partial_{x_l} (x_j \hat{e}_j) = 3 \vec{\omega} - \omega_i \delta_{ij} \hat{e}_j = 2 \vec{\omega},
\]

and

\[
\nabla \cdot (\vec{\omega} \times \vec{r}) = \nabla \cdot (\vec{\omega} \times r) = \vec{r} \cdot (\nabla \times \vec{\omega}) - \omega \cdot (\nabla \times r) = 0.
\]

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Thus we have two more important and useful results, which hold for all constant vectors $\omega$

$$\nabla \times (\omega \times r) = 2\omega \quad \text{and} \quad \nabla \cdot (\omega \times r) = 0$$

To understand intuitively the non-zero curl of $v(r) = \omega \times r$, imagine that the flow lines are those of a rotating fluid with a ball centred on a flow line of the field. The centre of the ball will follow the flow line. However the effect of the neighbouring flow lines is to make the ball rotate. Therefore the field has non-zero ‘curl’, and the axis of rotation gives the direction that the field rotates or ‘curls’ around.

In the previous example, $a(r) = r$, the ball would just move away from the origin without rotating, therefore the field $r$ has zero curl. Colloquially, the vector field $a(r) = r$, doesn’t “curl around” any point at all.

**Terminology:** For any vector field $a(r)$

(i) If $\nabla \cdot a = 0$ in some region $R$, $a$ is said to be **solenoidal** in $R$.

(ii) If $\nabla \times a = 0$ in some region $R$, $a$ is said to be **irrotational** in $R$.

## 7 Line integrals

Having completed our study of differential vector calculus using Cartesian coordinates, we now embark on integral vector calculus.

We begin by revising the standard definition of the integral of a function of a single variable. We then introduce line integrals, surface integrals and volume integrals. We shall assume familiarity with integrals over plane areas (double integrals) using Cartesian and plane polar coordinates, as covered in *Linear Algebra and Several Variable Calculus* or the School of Mathematics course *Several Variable Calculus and Differential Equations*.

### 7.1 Revision of ordinary integrals

We start with some formal definitions and discuss some limits, but we shall not be rigorous. Let $f(x)$ be a continuous real-valued function defined on the interval $a \leq x \leq b$. Begin by subdividing this interval into $n$ subintervals:

$$a = x_0 < x_1 < x_2 \ldots < x_n = b.$$ 

In interval $j$, pick an arbitrary point $x_j^*$ with $x_j \leq x_j^* \leq x_{j+1}$, as illustrated in the figure.
Define the Riemann sum
\[ S_n = \sum_{j=0}^{n-1} f(x_j^*) (x_{j+1} - x_j) \]

It can be shown that \( S_n \to \text{unique limit} \) as we let \( n \to \infty \) and \( (x_{j+1} - x_j) \to 0 \) for all \( j \), with \( \sum_{j=0}^{n-1} (x_{j+1} - x_j) = (b - a) \) kept fixed.\(^{16}\)

In this limit
\[ S_n \to \int_a^b f(x) \, dx \]

One can then prove the usual properties of integrals. Note that we also define
\[ \int_b^a f(x) \, dx = -\int_a^b f(x) \, dx \]

### 7.2 Motivation and formal definition of line integrals

**Motivation:** As a physical example, consider a particle constrained to move along a wire under the influence of a force \( F(r) \).

Only the component of the force along the wire does work.

Thus the work, \( dW \), done by the force in moving the particle along the wire from \( r \) to \( r + dr \) is
\[ dW = (|F| \cos \theta) \, |dr| = F \cdot dr \]

\(|F| \cos \theta\) is the component of \( F \) along \( dr \), where \( \theta \) is the angle between \( F(r) \) and \( dr \).

The total work done in moving the particle along a wire which follows some curve \( C \) from point \( P \) to point \( Q \) is the sum of the (infinitesimal) \( dW \) factors, which is an integral
\[ W_C = \int_P^Q dW = \int_C F(r) \cdot dr \]

\(^{16}\)We assume that \( f(x) \) is sufficiently well behaved that the limit actually exists.
This is called a line integral along the curve $C$.

**Formal definition:** Let $a(r)$ be a vector field defined in the region $R$, and let $C$ be a curve in $R$ between two points $P$ and $Q$. Let $r$ be the position vector at some point on the curve, and let $dr$ be an infinitesimal vector along the curve at $r$.

The magnitude of $dr$ is the infinitesimal arc length: $ds = \sqrt{dr \cdot dr}$

If $t$ is the unit vector tangent to the curve at $r$ (i.e. $t$ points in the direction of $dr$ at $r$), then

$$dr = t \, ds$$

The line integral is defined formally as a Riemann sum by dividing the curve into $n$ intervals

$$\int_C a \cdot dr = \int_C a \cdot t \, ds \equiv \lim_{\delta s \to 0} \sum_{i=0}^{n-1} \left( a(r^{(i)}) \cdot t^{(i)} \right) \delta s^{(i)}$$

the $i$th interval having length $\delta s^{(i)}$, unit tangent vector $t^{(i)}$, etc. It can be shown that the limit is unique for sufficiently smooth vector fields $a(r)$.

The integrand $a \cdot t$ is the component of $a$ along $dr$ at the point $r$.

In Cartesian coordinates, we have

$$\int_C a \cdot dr = \int_C (a_1 dx_1 + a_2 dx_2 + a_3 dx_3) = \int_C \sum_{i=1}^{3} a_i \, dx_i \quad (43)$$

**NB:** In general, the line integral depends on the path joining $P$ and $Q$. For example, the $a_1$ component is $a_1(x_1, x_2, x_3)$ and all three coordinates will generally change at once along the path. Therefore, you cannot compute $\int a_1 \, dx_1$ as an ordinary integral over $x_1$ holding $x_2$ and $x_3$ constant. That would only be correct if $x_2$ and $x_3$ were constant along the path, i.e. if the path were parallel to the $x_1$ axis. This is a common source of mistakes.

### 7.3 Parametric representation of line integrals

The definition above was rather formal. What follows is much more useful.

A smooth curve in 3D can be parameterised by a single parameter. For example, if the curve is the trajectory of a particle, then a natural parameter is the time $t$. Sometimes the parameter is chosen to be the arc-length $s$ along the curve $C$.

If we parameterise the curve by the parameter $\lambda$ (varying from $\lambda_P$ to $\lambda_Q$), we can write the coordinates $x_i$ as functions of $\lambda$

$$x_1 = x_1(\lambda), \quad x_2 = x_2(\lambda), \quad x_3 = x_3(\lambda), \quad \text{with } \lambda_P \leq \lambda \leq \lambda_Q$$

and

$$\int_C a \cdot dr = \int_{\lambda_P}^{\lambda_Q} \left( a \cdot \frac{dr}{d\lambda} \right) \, d\lambda = \int_{\lambda_P}^{\lambda_Q} \left( a_1 \frac{dx_1}{d\lambda} + a_2 \frac{dx_2}{d\lambda} + a_3 \frac{dx_3}{d\lambda} \right) \, d\lambda$$

If necessary, the curve $C$ may be subdivided into sections, each with a different parameterisation (piecewise smooth curve).
Example: Let \( \mathbf{a}(r) = -ky \mathbf{e}_x + kx \mathbf{e}_y \), where \( k \) is a positive constant.

Evaluate \( \int_C \mathbf{a} \cdot d\mathbf{r} \) between the points with Cartesian coordinates \((1, 0, 0)\) and \((0, 1, 0)\) along the curve \( C: (x = \cos \lambda, y = \sin \lambda, z = 0) \), where \( 0 \leq \lambda \leq \pi/2 \).

Convince yourself that \( C \) is one quarter of a unit circle. Sketch the curve \( C \) and the field \( \mathbf{a}(r) \) on \( C \).

On the curve \( C \), we have
\[
\mathbf{a}(r) = -ky \mathbf{e}_x + kx \mathbf{e}_y = -k \sin \lambda \mathbf{e}_x + k \cos \lambda \mathbf{e}_y
\]
\[
\mathbf{r} = \cos \lambda \mathbf{e}_x + \sin \lambda \mathbf{e}_y
\]
\[
\frac{d\mathbf{r}}{d\lambda} = \left( -\sin \lambda \mathbf{e}_x + \cos \lambda \mathbf{e}_y \right)
\]

Therefore
\[
\int_C \mathbf{a} \cdot d\mathbf{r} = \int_0^{\pi/2} \left( \mathbf{a} \cdot \frac{d\mathbf{r}}{d\lambda} \right) d\lambda = \int_0^{\pi/2} k (\sin^2 \lambda + \cos^2 \lambda) d\lambda = \int_0^{\pi/2} k d\lambda = \frac{k\pi}{2}
\]

In this example, the result is simple because the field \( \mathbf{a} \) is parallel to \( \frac{d\mathbf{r}}{d\lambda} \) and hence to \( d\mathbf{r} \).

Example: Let \( \mathbf{a}(r) = (3x^2 + 6y) \mathbf{e}_x - 14yz \mathbf{e}_y + 20xz^2 \mathbf{e}_z \).

Evaluate \( \int_C \mathbf{a} \cdot d\mathbf{r} \) between the points with Cartesian coordinates \((0, 0, 0)\) and \((1, 1, 1)\), along the two paths \( C \):

(i) \((0, 0, 0) \to (1, 0, 0) \to (1, 1, 0) \to (1, 1, 1)\)

(These are 3 contiguous straight lines parallel to the \( x \), \( y \) & \( z \) axes respectively.)

(ii) \((x = \lambda, y = \lambda^2, z = \lambda^3) \) from \( \lambda = 0 \) to \( \lambda = 1 \).

(i) (a) Along the line from \((0, 0, 0)\) to \((1, 0, 0)\), we have \( y = z = 0 \), so \( dy = dz = 0 \), hence
\[
\frac{d\mathbf{r}}{d\lambda} = \mathbf{e}_x d\lambda \text{ and } \mathbf{a} = 3x^2 \mathbf{e}_x \text{ (here the parameter is just } x \text{ itself), and}
\]
\[
\int_{(0,0,0)}^{(1,0,0)} \mathbf{a} \cdot d\mathbf{r} = \int_{x=0}^{x=1} 3x^2 \, dx = \left[ x^3 \right]_0^1 = 1
\]
(b) Along the line from \((1,0,0)\) to \((1,1,0)\), we have \(x = 1\), \(dx = 0\), \(z = dz = 0\), so 
\[
dr = e_y \, dy \quad \text{(here the parameter is } y)\text{, and}
\]
\[
a = \left(3x^2 + 6y \right) e_x \bigg|_{x=1} \quad = (3+6y)e_x
\]
\[
\Rightarrow \quad \int_{(1,0,0)}^{(1,1,0)} a \cdot dr = \int_{y=0}^{y=1} (3+6y) \, e_x \cdot e_y \, dy = 0
\]

(c) Along the line from \((1,1,0)\) to \((1,1,1)\), we have \(x = y = 1\), \(dx = dy = 0\), and hence \(dr = e_z \, dz\) and \(a = 9e_x - 14ze_y + 20z^2e_z\). Therefore 
\[
\int_{(1,1,0)}^{(1,1,1)} a \cdot dr = \int_{z=0}^{z=1} 20z^2 \, dz = \left[ \frac{20}{3} z^3 \right]_0^1 = \frac{20}{3}
\]

Adding up the 3 contributions, we get
\[
\int_C a \cdot dr = 1 + 0 + \frac{20}{3} = \frac{23}{3} \quad \text{along path (i)}
\]

(ii) To integrate \(a = (3x^2 + 6y)e_x - 14yze_y + 20xz^2e_z\) along path (ii), we parameterise
\[
\vec{r} = \lambda e_x + \lambda^2 e_y + \lambda^3 e_z
\]
\[
\frac{dr}{d\lambda} = e_x + 2\lambda e_y + 3\lambda^2 e_z
\]
\[
a = (3\lambda^2 + 6\lambda^3) e_x - 14\lambda^5 e_y + 20\lambda^7 e_z \quad \text{so that}
\]
\[
\int_C \left( a \cdot \frac{dr}{d\lambda} \right) \, d\lambda = \int_{\lambda=0}^{\lambda=1} \left( 9\lambda^2 - 28\lambda^6 + 60\lambda^9 \right) \, d\lambda = \left[ 3\lambda^3 - 4\lambda^7 + 6\lambda^{10} \right]_0^1 = 5
\]

Hence 
\[
\int_C a \cdot dr = 5 \quad \text{along path (ii)}
\]

In this case, the integral of \(a\) from \((0,0,0)\) to \((1,1,1)\) depends on the path taken.

Notes:

(i) The line integral \(\int_C a \cdot dr\) is a scalar quantity. Another scalar line integral is \(\int_C f \, ds\) 
where \(f(r)\) is a scalar field and \(ds = \sqrt{dr \cdot dr}\) is the infinitesimal arc-length introduced earlier.

(ii) A line integral around a simple (doesn’t intersect itself) closed curve \(C\) is denoted by the symbol \(\oint_C\)

Example: \(\oint_C a \cdot dr\) \(\equiv\) the circulation of \(a\) around \(C\)
Example: Let \( f(r) = ax^2 + by^2 \). Evaluate \( \oint_C f \, ds \) around the unit circle \( C \) centred on the origin in the \( x-y \) plane:

\[
x = \cos \phi, \; y = \sin \phi, \; z = 0; \quad 0 \leq \phi \leq 2\pi.
\]

On the curve \( C \):

\[
\begin{align*}
f(r) &= ax^2 + by^2 = a \cos^2 \phi + b \sin^2 \phi \\
l &= \cos \phi \hat{e}_x + \sin \phi \hat{e}_y \\
dr &= \left( -\sin \phi \hat{e}_x + \cos \phi \hat{e}_y \right) \, d\phi
\end{align*}
\]

\[
\Rightarrow \quad ds = \sqrt{dr \cdot dr} = (\sin^2 \phi + \cos^2 \phi)^{1/2} \, d\phi = d\phi
\]

Therefore, in this example,

\[
\oint_C f \, ds = \int_0^{2\pi} (a \cos^2 \phi + b \sin^2 \phi) \, d\phi = \pi (a + b)
\]

The length \( s \) of a curve \( C \) is given by \( s = \int_C ds \). In this example \( s = \int_0^{2\pi} d\phi = 2\pi \).

We can also define vector line integrals, whose result is a vector:

\[
\begin{align*}
(i) \quad & \int_C a \, ds = \int_C (\hat{e}_1 a_1 + \hat{e}_2 a_2 + \hat{e}_3 a_3) \, ds \quad \text{in Cartesian coordinates = a vector.} \\
(ii) \quad & \int_C a \times dr = \int_C \left[ \hat{e}_1 (a_2 dx_3 - a_3 dx_2) + \hat{e}_2 (a_3 dx_1 - a_1 dx_3) + \hat{e}_3 (a_1 dx_2 - a_2 dx_1) \right] \\
& \quad \quad \text{in Cartesian coordinates. The parametric form is simply} \int_C a \times dr = \int_C \left( a \times \frac{dr}{d\lambda} \right) \, d\lambda \\
(iii) \quad & \int_C f \, dr = \int_C f(\hat{e}_1 dx_1 + \hat{e}_2 dx_2 + \hat{e}_3 dx_3) \quad \text{in Cartesian coordinates. In parametric form, this becomes} \int_C f \, dr = \int_C f \frac{dr}{d\lambda} \, d\lambda
\end{align*}
\]

7.4 Current loop in a magnetic field

Consider an electric current of magnitude \( I \) flowing along a thin wire in the shape of a closed path \( C \).

The magnetic force on an element \( dr \) of the wire at \( r \) due to an external magnetic field \( \mathbf{B}(r) \) is given by the Lorentz force

\[
\mathbf{dF}(r) = I \, dr \times \mathbf{B}(r).
\]

The total force \( F \) on the wire is the vector sum of the forces on the individual elements, which is given by the line integral of \( \mathbf{dF} \) around the closed curve \( C \).

\[
\mathbf{F} = \oint_C \mathbf{dF} = \oint_C I \, dr \times \mathbf{B}(r) = -I \oint_C \mathbf{B}(r) \times dr
\]
**Example:** For the case where the external magnetic field is $B(r) = B_0(x e_x + y e_y)$, evaluate the total force on a circular current loop of radius $a$ which lies in the $x-y$ plane and is centred on the origin.

We parameterise the curve by the angle $\phi$ (as in plane polars), so that on the curve $C$, we have

\[
\begin{align*}
  r &= a \cos \phi e_x + a \sin \phi e_y \\
  dr &= (-\sin \phi e_x + a \cos \phi e_y) \, d\phi \\
  B &= B_0(a \cos \phi e_x + a \sin \phi e_y)
\end{align*}
\]

\[
\Rightarrow \oint_C B \times dr = B_0 \int_0^{2\pi} (a^2 \cos^2 \phi + a^2 \sin^2 \phi) e_z \, d\phi = B_0 e_z \int_0^{2\pi} d\phi = 2\pi a^2 B_0 e_z
\]

So $F = -2\pi a^2 B_0 I e_z$ which is in a vertically downward direction.

### 8 Surface integrals

Let $S$ be a two-sided surface in three-dimensional space as shown. If an infinitesimal element of surface with (scalar) area $dS$ has unit normal $\vec{n}$, then the infinitesimal vector element of area is defined by

\[
d\vec{S} = \vec{n} \, dS
\]

**Example:** If $S$ lies in the $(x, y)$ plane, then in Cartesian coordinates the infinitesimal scalar element of area is $dS = dx \, dy$, and the infinitesimal vector element of area is $d\vec{S} = e_z \, dx \, dy$.

**Geometrical interpretation:** $\vec{m} \cdot d\vec{S}$ gives the projected (scalar) element of area onto the plane with unit normal $\vec{m}$. See later for more details.

For closed surfaces (e.g. a sphere) we always choose $\vec{n}$ to be the outward normal. For open surfaces (e.g. the curved surface of a hemisphere), the sense of $\vec{n}$ is arbitrary – except that it is chosen in the same sense for all elements of the surface.
If \( a(r) \) is a vector field defined on the surface \( S \), we define the (normal) surface integral formally by the Riemann sum

\[
\int_S a \cdot dS = \int_S a \cdot n \, dS = \lim_{m \to \infty} \sum_{i=0}^{m-1} \left( a \left( r^{(i)} \right) \cdot n^{(i)} \right) \delta S^{(i)}
\]

where we divided the surface \( S \) into \( m \) small areas, the \( i \)th area having vector area \( \delta S^{(i)} \).

The quantity \( a \left( r^{(i)} \right) \cdot n^{(i)} \) is the component of \( a \) normal to the surface at the point \( r^{(i)} \). This is similar to the definition of planar integrals given in Linear Algebra and Several Variable Calculus, except that the (scalar) area elements \( \delta S^{(i)} \) are not constrained to lie in a plane.

- We use the notation \( \int_S a \cdot dS \) for both open and closed surfaces. Sometimes the integral over a closed surface is denoted by \( \oint_S a \cdot dS \) (not used here).

- The integral over \( S \) is an ordinary double integral in each case.\(^{17}\)

**Example:** Let \( S \) be the surface of a unit cube. Note that \( S \) is the sum of all six faces of the cube.

On the front face, parallel to the \((y, z)\) plane, at \( x = 1 \),

\[
dS = n \, dS = e_x \, dy \, dz
\]

On the back face at \( x = 0 \) in the \((y, z)\) plane,

\[
dS = n \, dS = -e_x \, dy \, dz
\]

Similarly for the other four faces.

In each case, the unit normal \( n \) is an outward normal because \( S \) is a closed surface.

If \( a(r) \) is a vector field, the integral \( \int_S a \cdot dS \) over the front face, where \( dS = e_x \, dy \, dz \), is

\[
\int_{z=0}^{z=1} \int_{y=0}^{y=1} a \cdot (e_x \, dy \, dz) = \int_{z=0}^{z=1} dz \int_{y=0}^{y=1} a \cdot e_x \, dy = \int_{z=0}^{z=1} dz \int_{y=0}^{y=1} a_x \bigg|_{x=1} \, dy
\]

where \( a_x \bigg|_{x=1} = a_x(1, y, z) \). The integral over \( y \) and \( z \) is an ordinary double integral over a square of side 1.

Similarly, the integral over the back face, where \( dS = -e_x \, dy \, dz \) is

\[
\int_{z=0}^{z=1} \int_{y=0}^{y=1} a \cdot (-e_x \, dy \, dz) = -\int_{z=0}^{z=1} dz \int_{y=0}^{y=1} a \cdot e_x \, dy = -\int_{z=0}^{z=1} dz \int_{y=0}^{y=1} a_x \bigg|_{x=0} \, dy
\]

\(^{17}\)Surface integrals are sometimes denoted by \( \iint_S a \cdot dS \) (not used here).
where \( a_x \bigg|_{x=0} \equiv a_x(0, y, z) \). Again, this is an ordinary double integral over a square of side 1. The total integral over \( S \) is the sum of integrals over all 6 faces of the cube. See Tutorial 6 for an explicit example.

8.1 Parametric form of the surface integral

We now introduce a parametric representation for surface integrals\(^{[18]}\).

Suppose the points on a surface \( S \) can be specified by two real parameters \( u \) and \( v \), so that the position vector on the surface may be written as

\[
\mathbf{r} = \mathbf{r}(u, v) = x(u, v) \mathbf{e}_x + y(u, v) \mathbf{e}_y + z(u, v) \mathbf{e}_z
\]

Then

- the lines \( \mathbf{r}(u, v) \) for fixed \( u \), variable \( v \), and
- the lines \( \mathbf{r}(u, v) \) for fixed \( v \), variable \( u \)

are parametric lines and form a grid on the surface \( S \) as shown.

If we change \( u \) by \( du \), and \( v \) by \( dv \), then \( \mathbf{r} \) changes by \( d\mathbf{r} \), where the infinitesimal vector \( d\mathbf{r} \) lies in the surface, and is given by

\[
d\mathbf{r} = \frac{\partial \mathbf{r}}{\partial u} du + \frac{\partial \mathbf{r}}{\partial v} dv
\]

Along the curves \( v = constant \), we have \( dv = 0 \), so \( d\mathbf{r} \) is just

\[
d\mathbf{r}_u = \frac{\partial \mathbf{r}}{\partial u} du
\]

The vector \( \partial \mathbf{r} / \partial u \) is tangent to the surface at \( \mathbf{r} \), and tangent to the lines \( v = constant \).

Similarly, for \( u = constant \), we have

\[
d\mathbf{r}_v = \frac{\partial \mathbf{r}}{\partial v} dv
\]

so \( \partial \mathbf{r} / \partial v \) is tangent to the surface at \( \mathbf{r} \), and tangent to the lines \( u = constant \), as shown in the figure.

---

\(^{[18]}\)This is the extension to 3D of the representation used for planar integrals in LAeSVC Section (18.4).
The infinitesimal vectors $d\mathbf{r}_u$ and $d\mathbf{r}_v$ lie in the surface at $\mathbf{r}$.

The vectors $\partial_u \mathbf{r}$ and $\partial_v \mathbf{r}$ lie in the tangent plane to the surface at $\mathbf{r}$.

We can construct a unit vector $n$, normal to the surface at $\mathbf{r}$

$$n = \frac{(\partial_u \mathbf{r} \times \partial_v \mathbf{r})}{|\partial_u \mathbf{r} \times \partial_v \mathbf{r}|}$$

Since the vector element of area, $d\mathbf{S}$, has magnitude equal to the area of the infinitesimal parallelogram shown in the figure above, and it points in the direction of $n$, we can write by virtue of the equation above and the geometric interpretation of the vector product (9)

$$d\mathbf{S} = d\mathbf{r}_u \times d\mathbf{r}_v = (\partial_u \mathbf{r} \, du) \times (\partial_v \mathbf{r} \, dv) = (\partial_u \mathbf{r} \times \partial_v \mathbf{r}) \, du \, dv$$

Finally, our integral is parameterised as

$$\int_S a \cdot d\mathbf{S} = \int_u \int_v a \cdot (\partial_u \mathbf{r} \times \partial_v \mathbf{r}) \, du \, dv$$

We use two integral signs when writing surface integrals in terms of explicit parameters $u$ and $v$. The limits for the integrals over $u$ and $v$ must be chosen appropriately for the surface.

This was a little abstract and a little complicated...

Fortunately, in most practical cases we don’t need the detailed form of this expression, since we tend to use orthogonal curvilinear coordinates, where the vectors $\partial \mathbf{r}/\partial u$ and $\partial \mathbf{r}/\partial v$ are perpendicular to each other. It is normally clear from the geometry of the situation whether this is the case, as it is for cylindrical coordinates and spherical polar coordinates - as we shall see.
9 Curvilinear coordinates, flux and surface integrals

9.1 Curvilinear coordinates

It is often convenient to work with coordinate systems other than the Cartesian coordinates \((x_1, x_2, x_3)\) or \((x, y, z)\).

9.1.1 Plane polar coordinates

We begin by revising a familiar example in the \(x-y\) plane.

Define two new variables \(\rho, \phi\) as 

\[
\rho = \sqrt{x^2 + y^2} \quad \text{and} \quad \tan \phi = \frac{y}{x}
\]

with \(0 \leq \rho \leq \infty\) and \(0 \leq \phi < 2\pi\).

Clearly, \(\rho\) is the distance from the origin \(O\) to the point \(P\) which has Cartesian coordinates \((x, y)\), and \(\phi\) is the angle between the \(x\)-axis and the line from the origin to the point \(P\) (measured in an anti-clockwise direction from the positive \(x\)-axis.)

The inverse transformation is clearly

\[
x = \rho \cos \phi \quad \text{and} \quad y = \rho \sin \phi
\]

The variables \(\rho\) and \(\phi\) are plane polar coordinates. Any point in the plane can be specified by its Cartesian coordinates \((x, y)\), or by its plane polar coordinates \((\rho, \phi)\). Clearly, \(\rho = 0\) at the origin, but \(\phi\) is undefined there.

Basis vectors: Thus far, we have used only Cartesian basis vectors \(\{\mathbf{e}_1, \mathbf{e}_2, \mathbf{e}_3\}\), also known as \(\{\mathbf{e}_x, \mathbf{e}_y, \mathbf{e}_z\}\). These point in the same direction everywhere in space.

Restricting ourselves to vectors in the \((x-y)\) plane for now, we can define two new orthonormal basis vectors in the plane.

Define the unit vector \(\mathbf{e}_\rho\) parallel to the position vector \(\mathbf{OP} \equiv \rho\).

Then define the unit vector \(\mathbf{e}_\phi\) orthogonal to \(\mathbf{e}_\rho\) and pointing in the direction of increasing \(\phi\).

Note that \(\mathbf{e}_\rho\) and \(\mathbf{e}_\phi\) point in different directions at different points in the plane, as shown in the figure for two points \(P\) and \(Q\).

Since \(\rho = x \mathbf{e}_x + y \mathbf{e}_y = \rho (\cos \phi \mathbf{e}_x + \sin \phi \mathbf{e}_y)\) then

\[
\mathbf{e}_\rho = \cos \phi \mathbf{e}_x + \sin \phi \mathbf{e}_y \quad \text{and} \quad \mathbf{e}_\phi = -\sin \phi \mathbf{e}_x + \cos \phi \mathbf{e}_y (44)
\]
Components of vectors in plane polars: Any vector \( \mathbf{a} \) which lies in the \( x-y \) plane may be expressed in the original Cartesian basis \( \{ \mathbf{e}_x, \mathbf{e}_y \} \), or in the plane polar basis \( \{ \mathbf{e}_\rho, \mathbf{e}_\phi \} \)

\[
\mathbf{a} = a_x \mathbf{e}_x + a_y \mathbf{e}_y = a_\rho \mathbf{e}_\rho + a_\phi \mathbf{e}_\phi
\]  

(45)

The quantities \( a_\rho \) and \( a_\phi \) are the components of the vector \( \mathbf{a} \) in the basis \( \{ \mathbf{e}_\rho, \mathbf{e}_\phi \} \).

Example: The position vector of a point in the \( x-y \) plane (which we call \( \mathbf{r} \) to distinguish it from the position vector \( \mathbf{r} \) in 3-d) is

\[
\mathbf{r} = x \mathbf{e}_x + y \mathbf{e}_y = \rho \cos \phi \mathbf{e}_x + \rho \sin \phi \mathbf{e}_y = \rho \mathbf{e}_\rho
\]  

(46)

NB: \( \mathbf{r} \neq \rho \mathbf{e}_\rho + \phi \mathbf{e}_\phi \). This is a common mistake and it’s very important not to make it! **By definition, \( \mathbf{r} \) has no component in the direction of \( \mathbf{e}_\phi \), as can be seen in the figure.**

### 9.1.2 Cylindrical coordinates

This is the simplest extension of plane polars to three dimensions. The usual plane polar coordinates \( \rho \) and \( \phi \) replace the \( x \) and \( y \) coordinates, but the Cartesian \( z \) coordinate is retained.

The cylindrical coordinates \( (\rho, \phi, z) \) of a point are therefore related to its Cartesian coordinates \( (x, y, z) \) by

\[
x = \rho \cos \phi, \quad y = \rho \sin \phi, \quad z = z
\]

where \( 0 \leq \rho < \infty \), \( 0 \leq \phi < 2\pi \), and \( -\infty < z < \infty \).

The position vector of the point with Cartesian coordinates \( (x, y, z) \) is

\[
\mathbf{r} = \rho \cos \phi \mathbf{e}_x + \rho \sin \phi \mathbf{e}_y + z \mathbf{e}_z
\]

where \( \{ \mathbf{e}_x, \mathbf{e}_y, \mathbf{e}_z \} \) are the usual Cartesian basis vectors.

The basis vectors for cylindrical coordinates are the plane polar basis vectors, plus the third Cartesian basis vector

\[
\mathbf{e}_\rho = \cos \phi \mathbf{e}_x + \sin \phi \mathbf{e}_y, \quad \mathbf{e}_\phi = -\sin \phi \mathbf{e}_x + \cos \phi \mathbf{e}_y, \quad \mathbf{e}_z = \mathbf{e}_z
\]

In this basis, the position vector for the point with cylindrical coordinates \( (\rho, \phi, z) \) is

\[
\mathbf{r} = \rho \mathbf{e}_\rho + z \mathbf{e}_z
\]

A general vector \( \mathbf{a} \) has components \( (a_x, a_y, a_z) \) in the Cartesian basis, and components \( (a_\rho, a_\phi, a_z) \) in the cylindrical basis, so that

\[
\mathbf{a} = a_x \mathbf{e}_x + a_y \mathbf{e}_y + a_z \mathbf{e}_z = a_\rho \mathbf{e}_\rho + a_\phi \mathbf{e}_\phi + a_z \mathbf{e}_z
\]

Cylindrical coordinates are useful for problems with cylindrical symmetry, and also for problems involving cones and paraboloids.

---

19. Cylindrical coordinates are also known as **circular cylindrical coordinates** or **cylindrical polar coordinates**.

20. Note that \( \mathbf{r} \neq \rho \mathbf{e}_\rho + \phi \mathbf{e}_\phi + z \mathbf{e}_z \).
9.1.3 Spherical polar coordinates

The spherical polar coordinates \((r, \theta, \phi)\) of a point with Cartesian coordinates \((x, y, z)\) are defined by

\[ x = r \sin \theta \cos \phi, \quad y = r \sin \theta \sin \phi, \quad z = r \cos \theta \]

with \(0 \leq r < \infty, \ 0 \leq \theta \leq \pi, \ 0 \leq \phi < 2\pi\).

The azimuthal angle \(\phi\) runs from 0 to \(2\pi\), but since \(\theta = \pi\) describes a point on the negative z axis, the polar angle \(\theta\) runs from 0 to \(\pi\) only, so that \(\theta\) and \(\phi\) are specified uniquely.

Note that \(\phi\) is not defined anywhere on the z axis, and \(\theta\) is not defined at the origin.

We may write the position vector as

\[
\mathbf{r} = r \sin \theta \cos \phi \mathbf{e}_1 + r \sin \theta \sin \phi \mathbf{e}_2 + r \cos \theta \mathbf{e}_3 = r \mathbf{e}_r
\]

where \(\{\mathbf{e}_1, \mathbf{e}_2, \mathbf{e}_3\}\) are the usual Cartesian basis vectors, and \(\mathbf{e}_r\) is a unit vector in the direction of \(\mathbf{r}\),

\[
\mathbf{e}_r = \sin \theta \cos \phi \mathbf{e}_1 + \sin \theta \sin \phi \mathbf{e}_2 + \cos \theta \mathbf{e}_3
\]

Unit vectors \(\mathbf{e}_\theta\) and \(\mathbf{e}_\phi\) will be derived later.

We may invert the expressions for \((x, y, z)\) in spherical polars \((r, \theta, \phi)\) to obtain

\[
r = \sqrt{x^2 + y^2 + z^2}, \quad \theta = \cos^{-1}\left\{\frac{z}{\sqrt{x^2 + y^2 + z^2}}\right\}, \quad \phi = \tan^{-1}\left(\frac{y}{x}\right).
\]

**NB** with our conventions, the angle \(\phi\) is the same in each of plane polar coordinates, cylindrical coordinates and spherical polar coordinates, and \(r\) is always the length of the position vector in three dimensions. Beware of other conventions!

Note that \(r \sin \theta\) in spherical polars is equal to the coordinate \(\rho\) in cylindrical coordinates.
9.2 Flux of a vector field through a surface

Let \( \mathbf{v}(r) \) be the velocity at a point \( r \) in a moving fluid. In a small region, where \( \mathbf{v} \) is approximately constant, the \textit{volume} of fluid crossing the element of vector area \( dS = \hat{n} dS \) in time \( dt \) is

\[
\left( |\mathbf{v}| dt \right) (dS \cos \theta) = (\mathbf{v} \cdot \mathbf{dS}) dt
\]

This is just the distance, \( |\mathbf{v}| dt \), travelled by the fluid, multiplied by the scalar area \textit{normal} to the direction of motion, that it flows through. This scalar area is \( dS \cos \theta = \hat{\mathbf{v}} \cdot \mathbf{dS} \), where \( \hat{\mathbf{v}} \) is a unit vector in the direction of \( \mathbf{v} \).

Therefore

\[
\mathbf{v} \cdot \mathbf{dS} = \text{volume per unit time of fluid crossing } dS
\]

\[
\Rightarrow \int_S \mathbf{v} \cdot \mathbf{dS} = \text{volume per unit time of fluid crossing a finite surface } S
\]

More generally, for a vector field \( \mathbf{a}(\mathbf{r}) \),

\[
\text{The surface integral } \int_S \mathbf{a} \cdot \mathbf{dS} \text{ is called the flux of } \mathbf{a} \text{ through the surface } S
\]

The concept of flux is useful in many different contexts \textit{e.g.} flux of molecules in a flow of gas; electric or magnetic flux through a surface, \textit{etc}.

**Example:** Let \( S \) be the surface of the sphere \( x_1^2 + x_2^2 + x_3^2 = R^2 \), with radius \( R \). Find the unit normal \( \hat{n} \), the vector element of area \( dS \), and evaluate the total flux of the vector field \( \mathbf{a}(\mathbf{r}) = \mathbf{r} / r^3 \) out of the sphere.

An arbitrary point \( \mathbf{r} \) on the surface \( S \) may be parameterised using the spherical polar coordinates \( \theta \) and \( \phi \) as

\[
\mathbf{r} = R \sin \theta \cos \phi \mathbf{e}_1 + R \sin \theta \sin \phi \mathbf{e}_2 + R \cos \theta \mathbf{e}_3 \quad \{0 \leq \theta \leq \pi, \ 0 \leq \phi < 2\pi\}
\]

so

\[
\partial_\theta \mathbf{r} = R \cos \theta \cos \phi \mathbf{e}_1 + R \cos \theta \sin \phi \mathbf{e}_2 - R \sin \theta \mathbf{e}_3
\]

and

\[
\partial_\phi \mathbf{r} = -R \sin \theta \sin \phi \mathbf{e}_1 + R \sin \theta \cos \phi \mathbf{e}_2 + \mathbf{e}_3
\]

Therefore

\[
\partial_\theta \mathbf{r} \times \partial_\phi \mathbf{r} = \begin{vmatrix}
\mathbf{e}_1 & \mathbf{e}_2 & \mathbf{e}_3 \\
R \cos \theta \cos \phi & R \cos \theta \sin \phi & -R \sin \theta \\
-R \sin \theta \sin \phi & R \sin \theta \cos \phi & 0
\end{vmatrix} = R^2 \sin^2 \theta \cos \phi \mathbf{e}_1 + R^2 \sin^2 \theta \sin \phi \mathbf{e}_2 + R^2 \cos \theta \left( \cos^2 \phi + \sin^2 \phi \right) \mathbf{e}_3
\]

\[
= R^2 \sin \theta \left( \sin \phi \mathbf{e}_1 + \cos \phi \mathbf{e}_2 \right)
\]

\[
= R^2 \sin \theta \mathbf{e}_r
\]

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where $\hat{e}_r$ is the unit vector ($\hat{e}_r \cdot \hat{e}_r = 1$) in the direction of $r$ that we introduced previously

$$\hat{e}_r = \sin \theta \cos \phi \hat{e}_1 + \sin \theta \sin \phi \hat{e}_2 + \cos \theta \hat{e}_3$$

Hence the unit normal $\hat{n} = \hat{e}_r$. Note: $\hat{e}_r = \hat{e}_r(\theta, \phi)$ is a function of $\theta$ and $\phi$, but not of $r$.

The vector element of area on the surface of the sphere is then

$$dS = \partial_\theta r \times \partial_\phi r \ d\theta \ d\phi = R^2 \sin \theta \ d\theta \ d\phi \hat{e}_r$$

On $S$, we have $r = R$, so the vector field $\mathbf{a}(r)$ on the surface $S$ is $\mathbf{a} = (R \hat{e}_r)/R^3 = \hat{e}_r/R^2$.

The flux of $\mathbf{a}$ through the closed surface $S$ is then

$$\int_S \mathbf{a} \cdot d\mathbf{S} = \int_S \frac{r}{r^3} \cdot d\mathbf{S} = \int_0^\pi \sin \theta \ d\theta \int_0^{2\pi} \ d\phi \left( \frac{\hat{e}_r}{R^2} \right) \cdot (R^2 \sin \theta \hat{e}_r) = \int_0^\pi \sin \theta \ d\theta \int_0^{2\pi} \ d\phi = 4\pi$$

In this example, the result of the integral is the surface area of a unit sphere. This is important in physics, as we now show...

**Physics example:** The electric field $E(r)$ at $r$ due to a point charge $q$ situated at the origin is

$$E(r) = \frac{q}{4\pi \epsilon_0} \frac{r}{r^3}$$

Evaluate the total flux of $E(r)$ out of the sphere $x_1^2 + x_2^2 + x_3^2 = R^2$.

Noting that the electric field $E(r) = q/(4\pi \epsilon_0) \mathbf{a}(r)$ in the example above, we can use our previous result to obtain

$$\int_S E \cdot d\mathbf{S} = \frac{q}{4\pi \epsilon_0} \int_S \frac{r}{r^3} \cdot d\mathbf{S} = \frac{q}{4\pi \epsilon_0} 4\pi = \frac{q}{\epsilon_0}$$

This is an example of Gauss’ Law of Electromagnetism. We shall show later that the total electric flux through any closed surface is the total charge enclosed by the surface, divided by the constant $\epsilon_0$.

**Basis vectors for spherical polars:**

In spherical polars, the position vector at the point $r$ with spherical polar coordinates $(r, \theta, \phi)$ is

$$r = r \sin \theta \cos \phi \hat{e}_1 + r \sin \theta \sin \phi \hat{e}_2 + r \cos \theta \hat{e}_3$$

Consider the change in $r$ as we let $r \rightarrow r + dr$

$$dr_r = r (r + dr, \theta, \phi) - r (r, \theta, \phi) = \frac{\partial r}{\partial r} dr$$

Clearly $dr_r$ and hence $\partial r/\partial r$ are parallel to $r$. 

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Similarly, we define the infinitesimal vectors \( \mathrm{d}r_\theta \) and \( \mathrm{d}r_\phi \).

Then the normalised vectors

\[
\mathbf{e}_r = \frac{\partial r}{|\partial r|}, \quad \mathbf{e}_\theta = \frac{\partial \theta}{|\partial \theta|}, \quad \mathbf{e}_\phi = \frac{\partial \phi}{|\partial \phi|},
\]

are [important exercise for the student]

\[
\mathbf{e}_r = \sin \theta \cos \phi \mathbf{e}_1 + \sin \theta \sin \phi \mathbf{e}_2 + \cos \theta \mathbf{e}_3, \quad \mathbf{e}_\theta = \cos \theta \cos \phi \mathbf{e}_1 + \cos \theta \sin \phi \mathbf{e}_2 - \sin \theta \mathbf{e}_3, \quad \mathbf{e}_\phi = -\sin \phi \mathbf{e}_1 + \cos \phi \mathbf{e}_2.
\]

These form an orthonormal basis, i.e. \( \mathbf{e}_r \cdot \mathbf{e}_r = \mathbf{e}_\theta \cdot \mathbf{e}_\theta = \mathbf{e}_\phi \cdot \mathbf{e}_\phi = 1 \), and \( \mathbf{e}_r \cdot \mathbf{e}_\theta = \mathbf{e}_\theta \cdot \mathbf{e}_\phi = \mathbf{e}_\phi \cdot \mathbf{e}_r = 0 \) [important exercise]

The unit vector \( \mathbf{e}_\theta \) points in the direction of increasing \( \theta \), with \( r \) and \( \phi \) fixed – as illustrated in the diagram above. Similarly for \( \mathbf{e}_r \) and \( \mathbf{e}_\phi \).

The basis vectors \( \mathbf{e}_r, \mathbf{e}_\theta \) and \( \mathbf{e}_\phi \) depend on \( \theta \) and \( \phi \), but not on \( r \). They point in different directions at different points. We say they form a non-Cartesian or curvilinear basis.

### 9.3 Other surface integrals

If \( f(r) \) is a scalar field, we may define a scalar surface integral

\[
\int_S f \, dS
\]

For example, the (scalar) area of the surface \( S \) is just

\[
S = \int_S dS = \int_S |dS| = \int_u \int_v |\partial_u r \times \partial_v r| \, du \, dv
\]

We may also define vector surface integrals

\[
\int_S f \, dS, \quad \int_S a \, dS, \quad \int_S a \times dS
\]

Each of these is a double integral, and is evaluated in a way similar to the scalar integrals, the result being a vector in each case.

The vector area of a surface is defined as \( S = \int_S d\mathbf{S} \).

**Example:** The vector area \( S \) of the (open) hemisphere, \( x_1^2 + x_2^2 + x_3^2 = R^2, \ (x_3 \geq 0) \), of radius \( R \), is, using spherical polars,

\[
S \equiv \int_S d\mathbf{S} = \int_{\phi=0}^{2\pi} \int_{\theta=0}^{\pi/2} R^2 \sin \theta \, \mathbf{e}_r \, d\theta \, d\phi.
\]
NB Since $\mathbf{e}_r = \sin \theta \cos \phi \mathbf{e}_1 + \sin \theta \sin \phi \mathbf{e}_2 + \cos \theta \mathbf{e}_3$ is not a constant (it depends on $\theta$ and $\phi$), we can’t take it out of the integral. So we have

$$
S = e_1 R^2 \int_0^{\pi/2} \sin^2 \theta \, d\theta \int_0^{2\pi} \cos \phi \, d\phi + e_2 R^2 \int_0^{\pi/2} \sin^2 \theta \, d\theta \int_0^{2\pi} \sin \phi \, d\phi +
$$

$$
e_3 R^2 \int_0^{\pi/2} \sin \theta \cos \theta \, d\theta \int_0^{2\pi} \, d\phi
$$

$$= 0 + 0 + \pi R^2 e_3$$

Notes

The vector surface integral over the full sphere $x_1^2 + x_2^2 + x_3^2 = R^2$ is zero because the contributions from the upper and lower hemispheres cancel. Similarly, the vector area of a closed hemisphere is zero because the vector area of the bottom face, which is a circular disc of radius $R$ in the $x_1 - x_2$ plane, is $-\pi R^2 e_3$. This is just the projection of the sum of infinitesimal vector areas onto the base of the hemisphere.

In fact, for any closed surface,

$$\int_S dS = 0$$

To show this, it is simplest to use the divergence theorem – see later.
Volume integrals

Volume integrals are conceptually simpler than line and surface integrals because the element of volume \( dV \) is a **scalar** quantity. We define integrals over a volume \( V \) in the standard way using the Riemann sum.

### 10.1 Integrals over scalar fields

Let \( f(r) \) be a scalar field defined in a volume \( V \). Divide \( V \) into \( n \) small volumes \( \delta V^{(i)} \). Then

\[
\int_V f(r) \, dV = \lim_{n \to \infty} \sum_{i=0}^{n-1} f(r^{(i)}) \, \delta V^{(i)}
\]

where \( f(r^{(i)}) \) is the value of the function \( f \) at some point \( r^{(i)} \) in the element of volume \( \delta V^{(i)} \).

In Cartesian coordinates

\[
\int_V f(r) \, dV = \int_z \int_y \int_x f(x, y, z) \, dx \, dy \, dz
\]

where the integrals over \( x, y \) and \( z \) have appropriate limits. Note that \( \int_V f(r) \, dV \) is a **triple** integral, but we use three integral signs only when writing it in terms of explicit integration variables, namely \( x, y, z \) in this case.

**Example:** Integrate \( f(x, y, z) = xyz^2 \) over the cuboid: \( \{ 0 \leq x < a, 0 \leq y < b, 0 \leq z < c \} \).

We choose to perform the integral over \( x \) first, keeping \( y \) and \( z \) fixed; then the integral over \( y \), keeping \( z \) fixed; and finally the integral over \( z \):

\[
I = \int_V f(r) \, dV = \int_y \int_x f(x, y, z) \, dx \, dy \, dz = \int_y \int_x x y z^2 \, dx \, dy \, dz
\]

The volume of the cuboid is simply

\[
V = \int_V dV = \int_x \int_y \int_z dx \, dy \, dz = abc
\]

The integrals may be performed in any order, but the limits on the integrals must be chosen appropriately to cover the volume \( V \).

For example, if we choose to perform the \( z \) integral first, its limits may depend on \( x \) and \( y \). The limits on the second integral over \( y \) may then depend on \( x \) (but not on \( z \), because we have already integrated over \( z \)). The limits on the last integral over \( x \) can’t depend on either \( y \) or \( z \) (because we have already integrated over them.)

**Example:**

\[
\int_0^1 dx \int_0^{x^2} dy \int_x^1 2x^2 z \, dz = \frac{28}{165} \text{ (exercise)}
\]
We now illustrate how to find the limits in an explicit example, albeit a rather complicated (and slightly masochistic) one.
**Example:** Use Cartesian coordinates to evaluate
\[ \int_V (x + y + z) \, dV \]
where \( V \) is the positive octant of the unit sphere:
\[ x^2 + y^2 + z^2 \leq 1, \quad x \geq 0, \quad y \geq 0, \quad z \geq 0 \]
If we choose to perform the \( z \) integral first:

(i) For fixed \((x, y)\), we first integrate with respect to \( z \), from the point \((x, y, 0)\) \( (i.e. \ z = 0) \), to the point \((x, y, z)\) with \( z = \sqrt{1 - x^2 - y^2} \). This is the integral up the strip shown.

(ii) Then, for fixed \( x \), we integrate wrt \( y \) from the vertical strip at \( y = 0 \) to the vertical strip at \( y = \sqrt{1 - x^2} \). This sums over all such strips in the planar quadrant shown.

(iii) Finally, we integrate from \( x = 0 \) to \( x = 1 \), which sums over all planes from \( x = 0 \) to \( 1 \).

\[ \Rightarrow \int_V (x + y + z) \, dV = \int_0^1 dx \int_0^{\sqrt{1-x^2}} dy \int_0^{\sqrt{1-x^2-y^2}} (x + y + z) \, dz \]

### 10.2 Parametric form of volume integrals

The example above was very complicated – because we used Cartesians. It’s much easier using a parametric representation. (We would use spherical polars in this example.)

Suppose we can write the position vector \( \vec{r} \) in terms of three real parameters \( u, v \) and \( w \), so that \( \vec{r} = r(u, v, w) \). If we make a small change in each of these parameters, then \( \vec{r} \) changes by

\[ \text{dr} = \partial_u \vec{r} \, du + \partial_v \vec{r} \, dv + \partial_w \vec{r} \, dw \]

Along the curves \( \{v = \text{constant}, \ w = \text{constant}\} \), we have \( dv = dw = 0 \), so \( dr \) is simply

\[ dr_u = \partial_u \vec{r} \, du \]

with \( dr_v \) and \( dr_w \) having similar definitions.

As shown in the figure, the vectors \( dr_u, dr_v \) and \( dr_w \) form the sides of an infinitesimal parallelepiped. The volume of a of a parallelepiped is given by the triple vector product (or the determinant) \( [8] \)

\[ dV = | \text{dr}_u \cdot (\text{dr}_v \times \text{dr}_w) | \]

\[ \Rightarrow dV = | \partial_u \vec{r} \cdot (\partial_v \vec{r} \times \partial_w \vec{r}) | \, du \, dv \, dw \]

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We take the magnitude of the scalar triple product so that the element of volume \( dV \) is always positive.

**Jacobians revisited:** We may write the volume element in terms of the Jacobian \( J \),
\[
dV = |J| \, du \, dv \, dw,
\]
where
\[
J \equiv \begin{vmatrix} \frac{\partial (x, y, z)}{\partial (u, v, w)} \end{vmatrix} = \begin{vmatrix} \partial_u x & \partial_u y & \partial_u z \\ \partial_v x & \partial_v y & \partial_v z \\ \partial_w x & \partial_w y & \partial_w z \end{vmatrix}
\]
This is the three dimensional version of the \( 2 \times 2 \) Jacobian derived in *LA&SVC* Section 18 and (presumably) in *SVC&DE*. It generalises to higher dimensions in a straightforward way.

**Example:** Consider a circular cylinder of radius \( a \), height \( c \). We can parameterise \( r \) using cylindrical coordinates. Within the cylinder, we have
\[
r = \rho \cos \phi \, \mathbf{e}_1 + \rho \sin \phi \, \mathbf{e}_2 + z \, \mathbf{e}_3 \quad \text{with} \quad \{ 0 \leq \rho \leq a, \ 0 \leq \phi \leq 2\pi, \ 0 \leq z \leq c \}
\]
Then
\[
\begin{align*}
\partial_\rho r &= \cos \phi \, \mathbf{e}_1 + \sin \phi \, \mathbf{e}_2 \\
\partial_\phi r &= -\rho \sin \phi \, \mathbf{e}_1 + \rho \cos \phi \, \mathbf{e}_2 \\
\partial_z r &= \mathbf{e}_3
\end{align*}
\]
And hence (easy exercise)
\[
dV = \left| \partial_\rho r \cdot (\partial_\phi r \times \partial_z r) \right| \, d\rho \, d\phi \, dz = \rho \, d\rho \, d\phi \, dz
\]
The volume of the cylinder is
\[
V = \int_V dV = \int_{\theta=0}^{\theta=2\pi} \int_{\phi=0}^{\phi=2\pi} \int_{\rho=0}^{\rho=a} \rho \, d\rho \, d\phi \, dz = \pi a^2 c.
\]

**Cylindrical basis:** the normalised vectors
\[
\mathbf{e}_\rho = \frac{\partial_\rho r}{|\partial_\rho r|} ; \quad \mathbf{e}_\phi = \frac{\partial_\phi r}{|\partial_\phi r|} ; \quad \mathbf{e}_z = \frac{\partial_z r}{|\partial_z r|}
\]
(shown in the figure) are the orthonormal basis vectors that we wrote down earlier for cylindrical coordinates (exercise).

**Exercise:** For spherical polars: \( r = r \sin \theta \cos \phi \, \mathbf{e}_1 + r \sin \theta \sin \phi \, \mathbf{e}_2 + r \cos \theta \, \mathbf{e}_3 \), show that
\[
dV = \left| \partial_r r \cdot (\partial_\theta r \times \partial_\phi r) \right| \, dr \, d\theta \, d\phi = r^2 \sin \theta \, dr \, d\theta \, d\phi
\]
10.3 Integrals over vector fields

The integral of a vector field \( \mathbf{a}(\mathbf{r}) \) over a volume \( V \) is

\[
\int_V \mathbf{a} \cdot d\mathbf{V} = \sum_{i=1}^{3} \varepsilon_i \int_V a_i \, dV = \sum_{i=1}^{3} \varepsilon_i \int_y \int_x a_i \, dx \, dy \, dz \quad \text{in Cartesian coordinates}
\]

The result of the integral is a vector, and we must evaluate 3 triple integrals for each of the 3 components of the vector.

Example: Consider a solid hemisphere of radius \( a \) centered on the \( \varepsilon_3 \) axis, with its bottom face at \( x_3 = 0 \). If the mass density (mass/unit volume), a scalar field, is \( \rho(\mathbf{r}) = \rho_0 / r \) where \( \rho_0 \) is a constant, what is the total mass, \( M \), of the hemisphere?

It is most convenient to use spherical polar coordinates. Then \( dV = r^2 \sin \theta \, dr \, d\theta \, d\phi \) and

\[
M = \int_V \rho(r) \, dV = \int_0^a \rho(r) \, r^2 \, dr \int_0^{\pi/2} \sin \theta \, d\theta \int_0^{2\pi} \, d\phi = 2\pi \rho_0 \int_0^a r \, dr = \pi \rho_0 a^2
\]

Now consider the centre of mass vector \( \mathbf{R} = (X \varepsilon_1 + Y \varepsilon_2 + Z \varepsilon_3) \), defined by

\[
\mathbf{MR} \equiv \int_V \rho(\mathbf{r}) \mathbf{r} \, dV
\]

We integrate each component of the vector field \( \rho(r) \mathbf{r} \) in turn using

\[
\mathbf{r} = r \sin \theta \cos \phi \varepsilon_1 + r \sin \theta \sin \phi \varepsilon_2 + r \cos \theta \varepsilon_3 \quad \text{and} \quad dV = r^2 \sin \theta \, dr \, d\theta \, d\phi
\]

which gives

\[
MX = \int_0^a \rho(r) r^3 \, dr \int_0^{\pi/2} \sin^2 \theta \, d\theta \int_0^{2\pi} \cos \phi \, d\phi = 0 \quad \text{(since the \( \phi \) integral gives 0)}
\]

\[
MY = \int_0^a \rho(r) r^3 \, dr \int_0^{\pi/2} \sin^2 \theta \, d\theta \int_0^{2\pi} \sin \phi \, d\phi = 0 \quad \text{(since the \( \phi \) integral gives 0)}
\]

\[
MZ = \int_0^a \rho(r) r^3 \, dr \int_0^{\pi/2} \sin \theta \cos \theta \, d\theta \int_0^{2\pi} d\phi = 2\pi \rho_0 \int_0^a r^2 \, dr \left[ \frac{1}{2} \sin^2 \theta \right]_0^{\pi/2}
\]

\[
= 2\pi \rho_0 \frac{a^3}{3} \frac{1}{2} = \frac{\pi \rho_0 a^3}{3}
\]

Hence

\[
\mathbf{R} = \frac{\pi \rho_0 a^3}{3} \varepsilon_3 = \frac{a}{3} \varepsilon_3 \quad \text{(independent of \( \rho_0 \))}
\]

Note that the integrals over \( \phi \) in the first two components of \( \mathbf{MR} \) are zero. Watch out for integrals that are zero – spotting them can save you a lot of unnecessary work!
10.4 Summary of polar coordinate systems

To conclude this section, we give a brief summary of polar coordinate systems.

In the figures below, \( dS \) indicates an area element and \( dV \) a volume element. We also sketch geometrical “derivations” of the infinitesimal elements of area and volume.

**Plane polar coordinates:** \((\rho, \phi)\)

\[
\begin{align*}
    x &= \rho \cos \phi \\
    y &= \rho \sin \phi \\
    dS &= \rho \, d\rho \, d\phi
\end{align*}
\]

**Cylindrical coordinates:** \((\rho, \phi, z)\)

\[
\begin{align*}
    x &= \rho \cos \phi \\
    y &= \rho \sin \phi \\
    z &= z \\
    dV &= \rho \, d\rho \, d\phi \, dz
\end{align*}
\]

**Spherical polar coordinates:** \((r, \theta, \phi)\)

\[
\begin{align*}
    x &= r \sin \theta \cos \phi \\
    y &= r \sin \theta \sin \phi \\
    z &= r \cos \theta \\
    dV &= r^2 \sin \theta \, dr \, d\theta \, d\phi
\end{align*}
\]
11 The divergence theorem

11.1 Integral definition of divergence

Let \( \mathbf{a} \) be a vector field in the region \( R \), and let \( P \) be a point in \( R \), then the divergence of \( \mathbf{a} \) at \( P \) may be defined by

\[
\text{div} \mathbf{a} = \lim_{\delta V \to 0} \frac{1}{\delta V} \int_{\delta S} \mathbf{a} \cdot d\mathbf{S}
\]

where \( \delta S \) is a closed surface in \( R \) which encloses the volume \( \delta V \). The limit must be taken so that the point \( P \) is within \( \delta V \). It can be shown that the limit is independent of the shape of \( \delta V \).

This definition of \( \text{div} \mathbf{a} \) is also basis independent, i.e. the result doesn’t depend on whether we evaluate it in Cartesian coordinates, spherical polars, etc.

We now show that our original definition of \( \nabla \cdot \mathbf{a} \) is recovered in Cartesian co-ordinates.

Let \( P \) be a point with Cartesian coordinates \((x_0, y_0, z_0)\) situated at the centre of a small rectangular block of size \( \delta x \times \delta y \times \delta z \), with volume \( \delta V = \delta x \delta y \delta z \).

- On the front face of the block, parallel to the \( y-z \) plane at \( x = x_0 + \delta x / 2 \), we have outward normal \( \mathbf{n} = e_x \) and so \( d\mathbf{S} = e_x dy \, dz \)

- On the back face of the block, parallel to the \( y-z \) plane at \( x = x_0 - \delta x / 2 \), we have outward normal \( \mathbf{n} = -e_x \) and so \( d\mathbf{S} = -e_x dy \, dz \)

Hence \( \mathbf{a} \cdot d\mathbf{S} = \pm a_x \, dy \, dz \) on these two faces. Let us denote the union (sum) of the two faces orthogonal to the \( x \) axis by \( \delta S_x \).

The contribution of these two surfaces to the integral \( \int_S \mathbf{a} \cdot d\mathbf{S} \) is given by

\[
\int_{\delta S_x} \mathbf{a} \cdot d\mathbf{S} = \int_z \int_y \left\{ a_x(x_0 + \delta x / 2, y, z) - a_x(x_0 - \delta x / 2, y, z) \right\} dy \, dz
\]

\[
= \int_z \int_y \left\{ a_x(x_0, y, z) + \frac{\delta x}{2} \left. \frac{\partial}{\partial x} a_x \right|_{(x_0, y, z)} + O(\delta x^2) \right\}
\]

\[
- \left\{ a_x(x_0, y, z) - \frac{\delta x}{2} \left. \frac{\partial}{\partial x} a_x \right|_{(x_0, y, z)} + O(\delta x^2) \right\}
\}

\[
= \int_z \int_y \left. \delta x \frac{\partial}{\partial x} a_x \right|_{(x_0, y, z)} \, dy \, dz
\]

\[
= \int_z \int_y \left. \delta x \frac{\partial}{\partial x} a_x \right|_{(x_0, y, z)} \, dy \, dz
\]

\[
= \int_z \int_y \delta x \frac{\partial}{\partial x} a_x \bigg|_{(x_0, y, z)} \, dy \, dz
\]

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where we Taylor-expanded \( a_x(x, y, z) \) about the point \((x_0, y, z)\), and we have dropped the \( O(\delta^2) \) terms which will vanish when we divide by \( \delta V \) at the end.

So

\[
\frac{1}{\delta V} \int_{\delta S_x} a \cdot dS = \frac{1}{\delta y \delta z} \int_x \int_y \frac{\partial_x a_x}{(x_0, y, z)} dy \, dz
\]

As we take the limit \( \delta_x, \delta_y, \delta_z \to 0 \), the integrand is approximately constant in the volume \( \delta V \), so that

\[
\int_x \int_y \frac{\partial_x a_x}{(x_0, y, z)} dy \, dz \to \delta_y \delta_z \frac{\partial_x a_x}{(x_0, y_0, z_0)}
\]

and hence

\[
\lim_{\delta V \to 0} \frac{1}{\delta V} \int_{\delta S_x} a \cdot dS = \frac{\partial_x a_x}{(x_0, y_0, z_0)}
\]

Adding similar contributions from the other 4 faces, we find

\[
\text{div} \, a = \partial_x a_x + \partial_y a_y + \partial_z a_z = \nabla \cdot a
\]

in agreement with our original definition in Cartesian co-ordinates. Thus we can use the notation \( \nabla \cdot a \) for both the differential and the integral definitions of the divergence of a vector field \( a(r) \).

The integral definition

\[
\text{div} \, a = \nabla \cdot a = \lim_{\delta V \to 0} \frac{1}{\delta V} \int_{\delta S} a \cdot dS
\]

provides a precise intuitive understanding of divergence as the (net) flux per unit volume leaving a small volume around a point \( r \). In pictures, for an infinitesimal volume \( dV \),

\[
\text{div} \, a > 0 \quad \text{(flux out > flux in)}
\]

\[
\text{div} \, a < 0 \quad \text{(flux out < flux in)}
\]

\[
\text{div} \, a = 0 \quad \text{(flux out = flux in)}
\]

11.2 The divergence theorem (Gauss’ theorem)

Let \( a \) be a vector field in a volume \( V \), and let \( S \) be the closed surface bounding \( V \), then

\[
\int_V \nabla \cdot a \, dV = \int_S a \cdot dS
\]
Proof: We derive the divergence theorem using the integral definition of $\nabla \cdot a$

$$\nabla \cdot a = \lim_{\delta V \to 0} \frac{1}{\delta V} \int_{\delta S} a \cdot dS.$$ 

Since this definition of $\nabla \cdot a$ is valid for volumes of arbitrary shape, we can build a smooth surface $S$ from a large number, $N$, of small blocks of volume $\delta V^{(i)}$ and area $\delta S^{(i)}$. For small, but not infinitesimal, $\delta S^{(i)}$ we may write

$$\nabla \cdot a(r^{(i)}) = \frac{1}{\delta V^{(i)}} \int_{\delta S^{(i)}} a \cdot dS + \epsilon^{(i)}$$

where $\epsilon^{(i)} \to 0$ as $\delta V^{(i)} \to 0$. Now multiply both sides by $\delta V^{(i)}$ and sum over all $i$

$$\sum_{i=1}^{N} \nabla \cdot a(r^{(i)}) \delta V^{(i)} = \sum_{i=1}^{N} \int_{\delta S^{(i)}} a \cdot dS + \sum_{i=1}^{N} \epsilon^{(i)} \delta V^{(i)} \quad (47)$$

On the RHS the contributions from surface elements interior to $S$ cancel. This is because where two blocks touch, the outward normals are in opposite directions, implying that the contributions to the respective integrals cancel.

To illustrate this, consider two adjacent blocks, 1 and 2.

The volume of block 1 is $\delta V_1$ and its surface is $\delta S_1$. Similarly for block 2. Denote the shaded surface common to the two blocks by $\delta S$, then

$$\int_{\delta S} a \cdot dS \text{ regarded as part of block 1}$$
$$= - \int_{\delta S} a \cdot dS \text{ regarded as part of block 2}$$

because their outward normals $dS_1$ and $dS_2$ are equal and opposite. Therefore

$$\int_{\delta S_1} a \cdot dS + \int_{\delta S_2} a \cdot dS = \int_{\delta S_{1+2}} a \cdot dS \quad \text{(because contributions from } \delta S \text{ cancel)}$$

where we have denoted the exterior surface of the compound block by $\delta S_{1+2}$.

Thus the contributions from all interior surface elements cancel pairwise. Using this result, and letting $N \to \infty$ in equation (47), the $O(\epsilon^{(i)} \delta V^{(i)})$ terms go to zero, and we get

$$\int_V \nabla \cdot a \, dV = \int_S a \cdot dS.$$ 

For an alternative proof of the divergence theorem, see Bourne & Kendall, Chapter 6.

Note: The divergence theorem is the generalisation to 3-D of

$$\int_a^b \frac{df(x)}{dx} \, dx = f(b) - f(a)$$
11.3 Volume of a body using the divergence theorem

The volume of a body is \( V = \int_V \, dV \)

Recalling that \( \nabla \cdot \mathbf{r} = 3 \), we can write

\[
V = \frac{1}{3} \int_V \nabla \cdot \mathbf{r} \, dV = \frac{1}{3} \int_S \mathbf{r} \cdot d\mathbf{S}
\]

where we used the divergence theorem in the last step.
Example: Consider the hemisphere \( x^2 + y^2 + z^2 \leq R^2 \) centered on \( \mathbf{e}_3 \) with its bottom face in the \( x-y \) plane. Recalling that the divergence theorem holds for a closed surface, the volume of the hemisphere is

\[
V = \frac{1}{3} \left[ \int_{S_C} \mathbf{r} \cdot d\mathbf{S} + \int_{S_B} \mathbf{r} \cdot d\mathbf{S} \right]
\]

where \( S_C \) is the curved surface of the hemisphere and \( S_B \) is its bottom. On \( S_B \), we have \( d\mathbf{S} = -\mathbf{e}_z dS \) and \( z = 0 \), so \( \mathbf{r} \cdot d\mathbf{S} = -z dS = 0 \). Therefore the only contribution comes from the (open) curved surface \( S_C \),

\[
V = \frac{1}{3} \int_{S_C} \mathbf{r} \cdot d\mathbf{S}
\]

We can evaluate this surface integral using spherical polars. For a hemisphere of radius \( R \) we showed previously that \( d\mathbf{S} = R^2 \sin \theta \, d\theta \, d\phi \, \mathbf{e}_r \).

On the hemisphere, \( \mathbf{r} \cdot d\mathbf{S} = R \mathbf{e}_r \cdot d\mathbf{S} = R^3 \sin \theta \, d\theta \, d\phi \), therefore

\[
V = \frac{1}{3} \int_{S_C} \mathbf{r} \cdot d\mathbf{S} = \frac{R^3}{3} \int_0^{\pi/2} \sin \theta \, d\theta \int_0^{2\pi} d\phi = \frac{2\pi R^3}{3}
\]

as anticipated.

11.4 The continuity equation

Consider a fluid with mass density \( \rho(r, t) \) and velocity field \( \mathbf{v}(r, t) \). We have seen previously that the volume flux (volume per unit time) of fluid flowing across a surface \( S \) is \( \int_S \mathbf{v} \cdot d\mathbf{S} \).

The corresponding mass flux (mass per unit time) flowing across the surface is

\[
\int_S (\rho \mathbf{v}) \cdot d\mathbf{S} \equiv \int_S J \cdot d\mathbf{S} \quad (48)
\]

where \( J(r, t) \equiv \rho(r, t) \mathbf{v}(r, t) \) is the mass current density.

Now consider a volume \( V \) bounded by the closed surface \( S \) containing no sources or sinks of fluid. Conservation of mass means that the outward mass flux through the surface \( S \) must be equal to the rate of decrease of mass of fluid contained in the volume \( V \),

\[
\int_S J \cdot d\mathbf{S} = -\partial_t M \quad , \quad (49)
\]

where \( M \) is the total mass in \( V \), which may be written as \( M = \int_V \rho \, dV \). Substituting this into equation \( (49) \), we get

\[
\frac{\partial}{\partial t} \int_V \rho \, dV + \int_S J \cdot d\mathbf{S} = 0 \, .
\]

Using the divergence theorem to rewrite the second term as a volume integral, we obtain

\[
\int_V \left[ \partial_t \rho + \nabla \cdot J \right] \, dV = 0
\]

Since this holds for arbitrary volumes \( V \), we must have
everywhere. This is the *continuity equation*. In this case, it expresses conservation of mass locally at each point \( \mathbf{r} \).
The continuity equation appears in many different contexts because it holds for any conserved quantity. Here we considered mass density $\rho$ and mass current density $J$ in a fluid, but equally it could have been thermal energy density and heat-current density, electric charge density and electric current density, or more abstract quantities such as probability density and probability-current density in quantum mechanics [tutorial].

In the case of fluid flow, the continuity equation tells us that

\[
\begin{align*}
\text{if } \nabla \cdot J &> 0 \text{ then } \partial_t \rho < 0 \text{ and the mass density at } r \text{ decreases} \\
\text{if } \nabla \cdot J &< 0 \text{ then } \partial_t \rho > 0 \text{ and the mass density at } r \text{ increases}
\end{align*}
\]

If the mass density at each point is constant in time, so that $\partial \rho / \partial t = 0$, the continuity equation tells us that for the mass density to be constant in time, and we must have $\nabla \cdot J = 0$, i.e. the mass flux into a point equals the flux out.

**Incompressible flow:** If $\rho(r,t)$ is a constant (i.e. time-independent and independent of $r$), then

\[
\nabla \cdot J = \rho \nabla \cdot \mathbf{v} = 0 \Rightarrow \nabla \cdot \mathbf{v} = 0
\]

Fluid flows with $\rho =$ constant, and hence $\nabla \cdot \mathbf{v} = 0$, are said to be incompressible flows.

### 11.5 Sources and sinks

We can generalise these ideas. For a vector field $\mathbf{a}(r)$, the quantity

\[
\frac{1}{V} \int_S \mathbf{a} \cdot d\mathbf{S}
\]

gives us information about whether there are sources or sinks of the vector field $\mathbf{a}$ within the volume $V$. If $V$ contains

- a net source, then \( \int_S \mathbf{a} \cdot d\mathbf{S} = \int_V \nabla \cdot \mathbf{a} \, dV > 0 \)
- a net sink, then \( \int_S \mathbf{a} \cdot d\mathbf{S} = \int_V \nabla \cdot \mathbf{a} \, dV < 0 \)

If $S$ contains neither sources nor sinks, or sources and sinks in equal measure, then

\[
\int_S \mathbf{a} \cdot d\mathbf{S} = 0.
\]

The quantity

\[
\nabla \cdot \mathbf{a} = \lim_{V \to 0} \frac{1}{V} \int_S \mathbf{a} \cdot d\mathbf{S}
\]

is therefore a measure of the density of sources or sinks,

\[
\nabla \cdot \mathbf{a} = \text{net outward flux per unit volume at } r.
\]

If we have a net source or sink of a vector field $\mathbf{a}$ at the point $r$, then $\nabla \cdot \mathbf{a} \neq 0$ at $r$.

These ideas can be applied to electric and magnetic fields.
11.6 Electrostatics - Gauss’ law and Maxwell’s first equation

As a very important example, we consider electrostatics. The electric field at \( r \) due to a point charge \( q \) at the origin is

\[
E(r) = \frac{q}{4\pi\varepsilon_0} \frac{r}{r^3}
\]

Then, for \( r \neq 0 \),

\[
\nabla \cdot E = \frac{q}{4\pi\varepsilon_0} \nabla \cdot \left( \frac{r}{r^3} \right) = \frac{q}{4\pi\varepsilon_0} \left( \nabla \left( \frac{1}{r^3} \right) \cdot r + \frac{\nabla \cdot r}{r^3} \right) = \frac{q}{4\pi\varepsilon_0} \left( -\frac{3r}{r^5} \cdot r + \frac{3}{r^3} \right) = 0
\]

In section (9.2), we showed that

\[
\int_{\text{sphere}} E \cdot dS = \frac{q}{4\pi\varepsilon_0} \int_{\text{sphere}} \frac{r \cdot dS}{r^3} = \frac{q}{4\pi\varepsilon_0} 4\pi = \frac{q}{\varepsilon_0} \quad \text{(50)}
\]

where the integral is over the surface of a sphere centred on the origin. The key result was \( \int_{\text{sphere}} (r \cdot dS)/r^3 = 4\pi \), independent of the radius of the sphere.

Now consider an arbitrary closed surface \( S \) which encloses the charge at the origin. Define the volume \( V \) to be the region between the surfaces \( S \) and \( S_1 \), where \( S_1 \) is a small sphere, radius \( \delta \), centred on the origin. The volume \( V \) is then bounded by the closed surface \( S + S_1 \).

(Ignore the spheres \( S_2 \) and \( S_3 \) in the figure for now.)

Since the volume \( V \) does not contain the origin, \( \nabla \cdot E = 0 \) everywhere in \( V \), and the divergence theorem tells us that

\[
\int_{S + S_1} E \cdot dS = \int_S E \cdot dS + \int_{S_1} E \cdot dS = \int_V \nabla \cdot E \, dV = 0 \quad \text{(51)}
\]

Since the outward normal on the sphere \( S_1 \) (i.e. outward from within the volume \( V \)) points towards the origin, equation (50) gives

\[
\int_{S_1} E \cdot dS = -\frac{q}{\varepsilon_0}
\]

independent of \( \delta \), and we may safely take the limit \( \delta \to 0 \). Equation (51) then becomes

\[
\int_S E \cdot dS = \frac{q}{\varepsilon_0} \quad \text{(52)}
\]

This holds for any closed surface \( S \) which encloses the charge at the origin.
If, instead of charge $q$ at the origin, we have charge $q_i$ at position $r_i$ inside $S$, we can change integration variable from $r$ to $\rho = r - r_i$ when integrating over the sphere $S_i$ (with outward normal pointing towards its centre), and we get

$$\int_{S_i} \mathbf{E}_i \cdot dS = \frac{q_i}{4\pi\epsilon_0} \int_{S_i} \frac{(r - r_i) \cdot dS}{|r - r_i|^3} = \frac{q_i}{4\pi\epsilon_0} \int_{S_i} \frac{\rho \cdot dS}{\rho^3} = -\frac{q_i}{\epsilon_0} \quad (53)$$

Equations (53) and (51) then give

$$\int_{S} \mathbf{E} \cdot dS = \frac{q_i}{\epsilon_0} \quad (54)$$

Now let’s replace the single charge by a set of $N$ charges $q_i$ at positions $r_i$. Experiment tells us that the total electric field $\mathbf{E}$ is the sum of the electric fields $\mathbf{E}_i$ due to the individual charges. Therefore, using equation (54), we get

$$\int_{S} \mathbf{E} \cdot dS = \int_{S} \left( \sum_{i=1}^{N} \mathbf{E}_i \right) \cdot dS = \sum_{i=1}^{N} \int_{S} \mathbf{E}_i \cdot dS = \sum_{i=1}^{N} \frac{q_i}{\epsilon_0} = \frac{Q}{\epsilon_0}$$

where $Q = \sum_{i=1}^{N} q_i$ is the total charge enclosed by $S$. This is Gauss’ Law of electrostatics.

Generalising further, if we have a charge density $\rho(r)$ (charge/unit volume), then the total charge in a volume $V$ is

$$Q = \int_{V} \rho(r) \, dV$$

Applying the divergence theorem and Gauss’ Law (respectively), we get

$$\int_{V} \nabla \cdot \mathbf{E} \, dV = \int_{S} \mathbf{E} \cdot dS = \frac{Q}{\epsilon_0} = \frac{1}{\epsilon_0} \int_{V} \rho(r) \, dV$$

Since this holds for arbitrary volumes $V$, we must have

$$\nabla \cdot \mathbf{E}(r) = \frac{\rho(r)}{\epsilon_0}$$

which holds for all $r$. This is Maxwell’s first equation of electromagnetism.

Evidently, it states that the divergence of the electric field at any point $r$ is equal to the charge density at that point divided by (in SI units) the constant $\epsilon_0$.

A positive charge is a source of electric field (i.e. it creates a positive flux) and a negative charge is a sink (i.e. it absorbs flux, or, equivalently, creates a negative flux).

### 11.7 Corollaries of the divergence theorem

We may deduce several immediate consequences of the divergence theorem

$$\int_{V} \nabla \cdot \mathbf{a} \, dV = \int_{S} \mathbf{a} \cdot dS$$

\[21\] We have derived it for static electric fields, but it also holds in electrodynamics – see Junior Honours Electromagnetism or Electromagnetism and Relativity.
(i) Let \( \mathbf{a} = \mathbf{c} \) where \( \mathbf{c} \) is an arbitrary constant vector, then \( \nabla \cdot \mathbf{c} = 0 \), and hence
\[
\int_S \mathbf{c} \cdot dS = \mathbf{c} \cdot \int_S dS = 0
\]
Since this holds for all constant vectors \( \mathbf{c} \), we must have
\[
\int_S dS = 0
\]
for any closed surface \( S \) (as claimed previously).

(ii) Let \( \mathbf{a}(\mathbf{r}) = p(\mathbf{r}) \mathbf{c} \) where \( p(\mathbf{r}) \) is a scalar field and \( \mathbf{c} \) is a constant vector, then (tutorial)
\[
\int_V (\nabla p) \, dV = \int_S p \, dS
\]
This is used to derive the condition \( \rho \mathbf{F} = \nabla p \), for hydrostatic equilibrium in a fluid or gas of density \( \rho \), subject to force per unit mass \( \mathbf{F} \) at pressure \( p \) in tutorial problem (4.6).

(iii) Consider the vector field \( \mathbf{a} \times \mathbf{c} \), where \( \mathbf{a}(\mathbf{r}) \) is an arbitrary vector field and \( \mathbf{c} \) is a constant vector. We have, using a standard identity,
\[
\nabla \cdot (\mathbf{a} \times \mathbf{c}) = \mathbf{c} \cdot (\nabla \times \mathbf{a}) - \mathbf{a} \cdot (\nabla \times \mathbf{c}) = \mathbf{c} \cdot (\nabla \times \mathbf{a}) - 0
\]
Now apply the divergence theorem to \( \mathbf{a} \times \mathbf{c} \)
\[
\mathbf{c} \cdot \int_V (\nabla \times \mathbf{a}) \, dV = \int_V \nabla \cdot (\mathbf{a} \times \mathbf{c}) \, dV = \int_S \mathbf{c} \cdot (\nabla \times \mathbf{a}) \, dS = \mathbf{c} \cdot \int_S dS \times \mathbf{a}
\]
This holds for all constant vectors \( \mathbf{c} \), hence
\[
\int_V \nabla \times \mathbf{a} \, dV = \int_S dS \times \mathbf{a}
\]

(iv) Green’s theorem in the plane

Let \( V \) be the volume inside the cylinder \( 0 < z < 1 \), and define the vector field \( \mathbf{a}(\mathbf{r}) \) as
\[
\mathbf{a} = P(x, y) \mathbf{e}_x + Q(x, y) \mathbf{e}_y
\]
Then
\[
\int_V \nabla \cdot \mathbf{a} \, dV = \int_V (\partial_x P + \partial_y Q) \, dx \, dy \, dz = \int_A (\partial_x P + \partial_y Q) \, dx \, dy
\]
where we performed the (trivial) integral over \( z \) to get unity, and defined \( A = S_B \) to be the bottom surface of the cylinder in the \( x-y \) plane.

Let \( S = S_C + S_T + S_B \) be the closed surface bounding \( V \). On the top and bottom surfaces, \( S_T \) and \( S_B \), we have \( dS = \pm dS \mathbf{e}_z \) and therefore \( \mathbf{a} \cdot dS = 0 \). On the curved surface, \( S_C \),
\[
dS = (dx \mathbf{e}_x + dy \mathbf{e}_y) \times \mathbf{e}_z \, dz = (dy \mathbf{e}_x - dx \mathbf{e}_y) \, dz
\]
Hence
\[ \int_S a \cdot dS = \int_{S_C} a \cdot dS = \int \int (P \, dy - Q \, dx) \, dz \]
The (trivial) integral over \( z \) again gives unity. Using the divergence theorem, we get
\[ \int_A (\partial_x P + \partial_y Q) \, dx \, dy = \oint (P \, dy - Q \, dx) \]
which is Green’s theorem in the plane (sometimes called the two dimensional divergence theorem) relating the integral over a planar surface \( A \) to the line integral around the closed curve \( C \) enclosing \( A \). The theorem applies to any surface in the \( x-y \) plane, because the proof above doesn’t rely on the base of the cylinder being circular.

12 Line integral definition of curl, Stokes’ theorem

12.1 Line integral definition of curl

Consider a small planar surface with unit normal \( \mathbf{n} \) and (scalar) area \( \delta S \), bounded by a closed curve \( \delta C \), and containing the point \( P \). Let \( \mathbf{a} \) be a vector field defined on this surface.

The component of \( \nabla \times \mathbf{a} \) in the direction of \( \mathbf{n} \) is defined to be
\[ n \cdot (\nabla \times \mathbf{a}) = \lim_{\delta S \to 0} \frac{1}{\delta S} \oint_{\delta C} \mathbf{a} \cdot d\mathbf{r} \]

NB: The integral around \( \delta C \) is taken in the right-hand sense with respect to the normal \( \mathbf{n} \) to the surface — as shown in the figure above.

This definition of curl is independent of the choice of basis.

Cartesian form of curl: The usual Cartesian form for curl can be recovered from this general definition by considering small rectangles parallel to the \( x-y \), \( y-z \), and \( z-x \) planes respectively.
Let \( P \) be a point with Cartesian coordinates \((x_0, y_0, z_0)\) situated at the centre of a small rectangle \( \delta C = ABCD \) of size \( \delta x \times \delta y \), area \( \delta S = \delta x \delta y \), parallel to the \( x-y \) plane.

The line integral around \( \delta C \) is given by the sum of four terms

\[
\oint_{\delta C} a \cdot dr = \int_A^B a \cdot dr + \int_B^C a \cdot dr + \int_C^D a \cdot dr + \int_D^A a \cdot dr
\]

\[
= \int_A^B a \cdot dr - \int_B^C a \cdot dr - \int_C^D a \cdot dr + \int_D^A a \cdot dr
\]

Since \( r = x\xi_x + y\xi_y + z\xi_z \) we have \( dr = e_x dx \) along \( D \to A \) and \( C \to B \), and \( dr = e_y dy \) along \( A \to B \) and \( D \to C \). Therefore

\[
\oint_{\delta C} a \cdot dr = \int_A^B a_y dy - \int_B^C a_x dx - \int_C^D a_y dy + \int_D^A a_x dx
\]

For small \( \delta x \) and \( \delta y \), we can Taylor expand the integrands about \( y = y_0 \),

\[
\int_D^A a_x dx = \int_D^A a_x(x, y_0 - \delta y/2, z_0) dx
\]

\[
= \int_{x_0 - \delta x/2}^{x_0 + \delta x/2} \left[ a_x(x, y_0, z_0) - \frac{\delta y}{2} \frac{\partial}{\partial y} a_x \bigg|_{(x, y_0, z_0)} + O(\delta y^2) \right] dx
\]

\[
\int_C^D a_x dx = \int_C^D a_x(x, y_0 + \delta y/2, z_0) dx
\]

\[
= \int_{x_0 - \delta x/2}^{x_0 + \delta x/2} \left[ a_x(x, y_0, z_0) + \frac{\delta y}{2} \frac{\partial}{\partial y} a_x \bigg|_{(x, y_0, z_0)} + O(\delta y^2) \right] dx
\]

so

\[
\frac{1}{\delta S} \left[ \int_D^A a \cdot dr + \int_B^C a \cdot dr \right] = \frac{1}{\delta x \delta y} \left[ \int_D^A a_x dx - \int_C^D a_x dx \right]
\]

\[
= \frac{1}{\delta x \delta y} \int_{x_0 - \delta x/2}^{x_0 + \delta x/2} \left[ -\delta_y \frac{\partial}{\partial y} a_x \bigg|_{(x, y_0, z_0)} + O(\delta y^2) \right] dx
\]

\[
\to -\frac{\partial_y}{(x_0, y_0, z_0)} \frac{\partial}{\partial y} a_x \bigg|_{(x_0, y_0, z_0)} \quad \text{as} \quad \delta_x, \delta_y \to 0
\]

In the last step, as we take the limit \( \delta_x \to 0 \), the integrand tends to a constant in the region of integration:

\[
\frac{\partial_y}{(x_0, y_0, z_0)} \frac{\partial}{\partial y} a_x \bigg|_{(x_0, y_0, z_0)} \to \frac{\partial_y}{(x_0, y_0, z_0)} \frac{\partial}{\partial y} a_x \bigg|_{(x_0, y_0, z_0)}
\]
and the integral over $x$ is then trivial.

A similar analysis of the line integrals along $A \to B$ and $C \to D$ gives (exercise)

$$\frac{1}{\delta S} \left[ \int_{A}^{B} a \cdot dr + \int_{C}^{D} a \cdot dr \right] \to \partial_x a_y \bigg|_{(x_0,y_0,z_0)} \text{ as } \delta_x, \delta_y \to 0$$

Adding the results gives, for our line integral definition of curl,

$$\varepsilon_z \cdot (\nabla \times a) = (\nabla \times a)_z = \left[ \partial_x a_y - \partial_y a_x \right]_{(x_0,y_0,z_0)}$$

in agreement with our original definition in Cartesian coordinates.

The other components of $\nabla \times a$ can be obtained from similar rectangles parallel to the $y-z$ and $x-z$ planes, respectively.

It can be shown that $\nabla \times a$, when defined in this way, is independent of the shape of the infinitesimal area $\delta S$. 

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12.2 Physical/geometrical interpretation of curl

Consider a force field \( \mathbf{F}(\mathbf{r}) \), and let \( \delta C \) be a small rectangular contour which encloses an area \( \delta S \) in the \( x-y \) plane – as in the line-integral definition of curl above.

The work done on a (point) test particle in moving it around the closed curve \( \delta C \) is

\[
\delta W = \oint_{\delta C} \mathbf{F} \cdot d\mathbf{r} = \text{circulation of } \mathbf{F}(\mathbf{r}) \text{ about } \delta C
\]

From the integral definition of curl, we know that for small \( \delta S \)

\[
\oint_{\delta C} \mathbf{F} \cdot d\mathbf{r} \approx (\nabla \times \mathbf{F})_z \delta S
\]

Therefore \( (\nabla \times \mathbf{F})_z \neq 0 \) is equivalent to saying that a non-zero amount of work is done in moving the test particle around a small closed path in the \( x-y \) plane.

Alternatively one can think of the non-zero circulation of \( \mathbf{F} \) as causing a small test particle to rotate about its centre, with the axis of rotation in the direction of \( \nabla \times \mathbf{F} \).

More generally, \( n \cdot (\nabla \times \mathbf{a}) \) is a measure of the net circulation (per unit area) of the vector field \( \mathbf{a} \) about an infinitesimal area \( dS \) with normal \( n \).

\[
\begin{array}{cccc}
\text{dS} & \text{dS} & \text{dS} & \text{dS} \\
\text{curl } \mathbf{a} > 0 & \text{curl } \mathbf{a} < 0 & \text{curl } \mathbf{a} = 0 & \text{curl } \mathbf{a} > 0
\end{array}
\]

12.3 Stokes’ theorem

Let \( S \) be an open surface, bounded by a simple closed curve \( C \), and let \( \mathbf{a} \) be a vector field defined on \( S \), then

\[
\int_S (\nabla \times \mathbf{a}) \cdot d\mathbf{S} = \oint_C \mathbf{a} \cdot d\mathbf{r}
\]

where \( C \) is traversed in a right-hand sense about \( dS \). As usual, \( d\mathbf{S} = n \, dS \) where \( n \) is the unit normal to \( S \).
Proof: Divide the surface $S$ into $N$ adjacent small surfaces. Let $\delta S^{(i)} = \delta S^{(i)} \mathbf{n}^{(i)}$ be the vector element of area at $\mathbf{r}^{(i)}$, enclosed by the curve $\delta C^{(i)}$.

Start with the integral definition of curl

$$n \cdot (\nabla \times a) = \lim_{\delta S \to 0} \frac{1}{\delta S} \oint_{\delta C} a \cdot d\mathbf{r},$$

For a small but not infinitesimal open surface $\delta S^{(i)}$

$$\left( \nabla \times a \left( \mathbf{r}^{(i)} \right) \right) \cdot \mathbf{n}^{(i)} = \frac{1}{\delta S^{(i)}} \oint_{\delta C^{(i)}} a \cdot d\mathbf{r} + \epsilon^{(i)}$$

where $\epsilon^{(i)} \to 0$ as $\delta S^{(i)} \to 0$.

Multiply by $\delta S^{(i)}$ (before taking the limit), and sum over all $i$ to get

$$\sum_{i=1}^{N} \left( \nabla \times a \left( \mathbf{r}^{(i)} \right) \right) \cdot \mathbf{n}^{(i)} \delta S^{(i)} = \sum_{i=1}^{N} \oint_{\delta C^{(i)}} a \cdot d\mathbf{r} + \sum_{i=1}^{N} \epsilon^{(i)} \delta S^{(i)}$$

Since each small closed curve $\delta C^{(i)}$ is traversed in the same sense, then, from the diagram, all contributions to $\sum_{i=1}^{N} \oint_{\delta C^{(i)}} a \cdot d\mathbf{r}$ cancel, except on those curves where part of $\delta C^{(i)}$ lies on the curve $C$. For example, the line integrals along the common section of the two small closed curves $\delta C^{(1)}$ and $\delta C^{(2)}$ in the figure cancel exactly. Therefore

$$\sum_{i=1}^{N} \oint_{\delta C^{(i)}} a \cdot d\mathbf{r} = \oint_{C} a \cdot d\mathbf{r}$$

Hence, as $N \to \infty$,

$$\oint_{C} a \cdot d\mathbf{r} = \int_{S} (\nabla \times a) \cdot d\mathbf{S} = \int_{S} n \cdot (\nabla \times a) \ d\mathbf{S}$$

Mathematical note: For those worried about the ‘error term’, note that, for finite $N$, we can establish an upper bound

$$\left| \sum_{i=1}^{N} \epsilon^{(i)} \delta S^{(i)} \right| \leq S \max_{i} \{|\epsilon^{(i)}|\}$$
The RHS tends to zero in the limit $N \to \infty$, because $S$ is finite and $\epsilon^{(i)} \to 0$, $\forall i$. A similar analysis works in the proof of the divergence theorem.\textsuperscript{22}

### 12.4 Examples of the use of Stokes’ theorem

**Hemisphere:** Given the vector field $a = 4ye_x + xe_y + 2ze_z$, verify Stokes’ theorem for the (open) hemispherical surface $x^2 + y^2 + z^2 = R^2$ with $z > 0$.

In this case, we have $\nabla \times a = -3e_z$, and we have shown previously that $dS = R^2 \sin \theta \, d\theta \, d\phi \, e_r$ on the surface of a (hemi)sphere of radius $R$. Direct integration then gives

$$\int_S \nabla \times a \cdot dS = \int_S (-3e_z) \cdot R^2 \sin \theta \, d\theta \, d\phi \, e_r = -3R^2 \int_0^{\pi/2} \sin \theta \cos \theta \, d\theta \int_0^{2\pi} d\phi = -3\pi R^2$$

We can check our result using Stokes’ theorem. The closed curve $C$ bounding the hemisphere is a circle of radius $R$ in the $x-y$ plane. Parameterising this by $x = R \cos \phi$, $y = R \sin \phi$, $z = 0$, gives $dx = -R \sin \phi \, d\phi$, $dy = R \cos \phi \, d\phi$, and $a_x = 4y = 4R \sin \phi$, $a_y = x = R \cos \phi$, $a_z = 2z = 0$. Hence

$$\oint_C a \cdot dr = \oint_C (4y \, dx + x \, dy)$$

$$= \int_0^{2\pi} (-4R^2 \sin^2 \phi + R^2 \cos^2 \phi) \, d\phi = -3\pi R^2$$

**Planar areas:** Consider a planar surface $S$ parallel to the $x-y$ plane, bounded by a closed curve $C$, and let the vector field $a(x)$ be

$$a = \frac{1}{2} \left[ -y e_x + x e_y \right]$$

In this case $\nabla \times a = e_z$, and the vector element of area normal to the $x-y$ plane is $d\underline{S} = dS \, e_z$.

Hence

$$\int_S \nabla \times a \cdot d\underline{S} = \int_S e_z \cdot dS = \int_S dS = S$$

We can then use Stokes’ theorem to find the area of the surface

$$S = \oint_C a \cdot dr = \frac{1}{2} \oint_C (-y e_x + x e_y) \cdot (dx e_x + dy e_y)$$

which gives

$$S = \frac{1}{2} \oint_C (x \, dy - y \, dx)$$

**Example:** Find the area inside the curve

$$x^{2/3} + y^{2/3} = 1.$$
The curve can be parameterised by \( x = \cos^3 \phi, y = \sin^3 \phi \), for \( 0 \leq \phi \leq 2\pi \), so that
\[
\frac{dx}{d\phi} = -3 \cos^2 \phi \sin \phi, \quad \frac{dy}{d\phi} = 3 \sin^2 \phi \cos \phi
\]
which gives
\[
S = \frac{1}{2} \oint_C (x \, dy - y \, dx) = \frac{1}{2} \oint_C \left( \frac{dy}{d\phi} - \frac{dx}{d\phi} \right) \, d\phi
\]
\[
= \frac{1}{2} \int_0^{2\pi} \left( 3 \cos^4 \phi \sin^2 \phi + 3 \sin^4 \phi \cos^2 \phi \right) \, d\phi
\]
\[
= \frac{3}{2} \int_0^{2\pi} \sin^2 \phi \cos^2 \phi \, d\phi = \frac{3}{8} \int_0^{2\pi} \sin^2 2\phi \, d\phi = \frac{3\pi}{8}
\]

12.5 Corollaries of Stokes’ theorem

We may deduce several immediate consequences of Stokes’ theorem,
\[
\int_S \left( \nabla \times \mathbf{a} \right) \cdot dS = \oint_C \mathbf{a} \cdot d\mathbf{r}
\]
where \( C \) is the boundary (traversed in the anticlockwise direction) of the open surface \( S \).

(i) If \( \mathbf{a} = \mathbf{c} \), where \( \mathbf{c} \) is a constant vector, then \( \nabla \times \mathbf{a} = 0 \). Therefore \( \mathbf{c} \cdot \oint_C d\mathbf{r} = 0 \), and because \( \mathbf{c} \) is arbitrary, we have
\[
\oint_C d\mathbf{r} = 0
\]

(ii) Take \( \mathbf{a} = -Q(x, y) \mathbf{e}_x + P(x, y) \mathbf{e}_y \), and \( S \) to lie in the \( x-y \) plane with area \( A \). Then
\[
\nabla \times \mathbf{a} = \left( \partial_x P + \partial_y Q \right) \mathbf{e}_z \quad \text{and} \quad dS = dx \, dy \, \mathbf{e}_z
\]
\[
\mathbf{a} \cdot d\mathbf{r} = -Q \, dx + P \, dy
\]
so
\[
\int_A \left( \partial_x P + \partial_y Q \right) \, dx \, dy = \oint_C (-Q \, dx + P \, dy)
\]
which is again Green’s theorem in the plane, sometimes known as Stokes’ theorem in the plane or the two-dimensional divergence theorem.

Taking \( P = x/2 \) and \( Q = y/2 \) and gives the planar area result of the previous section. Indeed, Green’s theorem is a generalisation of this result.

(iii) Applying Stokes’ theorem to \( \mathbf{a} = \phi \mathbf{c} \) where \( \mathbf{c} \) is a constant vector, we have
\[
\nabla \times (\phi \mathbf{c}) = \left( \nabla \phi \right) \times \mathbf{c} + \phi \left( \nabla \times \mathbf{c} \right) = \left( \nabla \phi \right) \times \mathbf{c} + 0
\]
Hence
\[
\left( \nabla \times (\phi \mathbf{c}) \right) \cdot dS = \left( \left( \nabla \phi \right) \times \mathbf{c} \right) \cdot dS = \mathbf{c} \cdot (dS \times \nabla \phi)
\]
which gives
\[ \int_S (\nabla \times (\phi \mathbf{c})) \cdot d\mathbf{S} = \mathbf{c} \cdot \int_S d\mathbf{S} \times \nabla \phi = \mathbf{c} \cdot \oint_C \phi \mathbf{d}r \]

This holds for all constant vectors \( \mathbf{c} \), so
\[ \oint_C \phi \mathbf{d}r = \int_S d\mathbf{S} \times \nabla \phi = -\int_S \nabla \phi \times d\mathbf{S} \]

Such results are hard to remember, but as we have seen, they can be derived quite easily.
13 The scalar potential

A vector field $\mathbf{a}(\mathbf{r})$ is defined to be irrotational or conservative if its curl vanishes, i.e. if

$$\nabla \times \mathbf{a} = 0$$

13.1 Path independence of line integrals for conservative fields

Let $\nabla \times \mathbf{a} = 0$ everywhere in some region, and consider two (different) paths $C_1$ and $C_2$ from point $\mathbf{r}_0$ to point $\mathbf{r}$, say. Applying Stokes’ theorem to the open surface $S$ bounded by the closed path $C_1 - C_2$ gives

$$\int_S (\nabla \times \mathbf{a}) \cdot dS = 0 = \int_{C_1} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}' - \int_{C_2} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}'$$

where the $-$ve sign occurs in the second integral on the RHS because both paths are defined to go from $\mathbf{r}_0$ to $\mathbf{r}$. We use $\mathbf{r}'$ as integration variable to distinguish it from the integration limits $\mathbf{r}_0$ and $\mathbf{r}$.

Therefore, when $\nabla \times \mathbf{a} = 0$ everywhere in $S$, we have

$$\int_{C_1} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}' = \int_{C_2} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}'$$

This is true for any $S$, and therefore for any paths $C_1$ and $C_2$ from $\mathbf{r}_0$ to $\mathbf{r}$.

Clearly, the converse is also true: if the line integral between two points is path independent, then the line integral around any closed curve (connecting the two points) is zero, and Stokes’ theorem then gives $\nabla \times \mathbf{a} = 0$. We just reverse the steps of the argument above.

Therefore

$$\nabla \times \mathbf{a} = 0 \iff \int_{\mathbf{r}_0}^{\mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}' \text{ is path independent}$$

13.2 Scalar potential for conservative vector fields

Since the line integral of a conservative vector field between two fixed points $\mathbf{r}_0$ and $\mathbf{r}$ is path independent, it can be a function only of the end points of the path. Hence there must exist a function $\phi(\mathbf{r})$ such that

$$\phi(\mathbf{r}) - \phi(\mathbf{r}_0) = \int_{\mathbf{r}_0}^{\mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot d\mathbf{r}' . \quad (55)$$

The scalar field $\phi(\mathbf{r})$ is called the scalar potential of the vector field $\mathbf{a}(\mathbf{r})$.

It is useful to invert this equation (and to give a more conventional result) by considering
two neighbouring points \( \mathbf{r} \) and \( \mathbf{r} + \mathbf{d} \mathbf{r} \), for which

\[
\begin{align*}
\text{d} \phi &= \phi(\mathbf{r} + \mathbf{d} \mathbf{r}) - \phi(\mathbf{r}) \\
&= [\phi(\mathbf{r} + \mathbf{d} \mathbf{r}) - \phi(\mathbf{r} + \mathbf{r}_0)] - [\phi(\mathbf{r}) - \phi(\mathbf{r}_0)] \\
&= \int_{\mathbf{r}_0}^{\mathbf{r} + \mathbf{d} \mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot \text{d} \mathbf{r}' - \int_{\mathbf{r}_0}^{\mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot \text{d} \mathbf{r}' & \text{(using equation (55))} \\
&= \int_{\mathbf{r}_0}^{\mathbf{r} + \mathbf{d} \mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot \text{d} \mathbf{r}' & \text{along any path from } \mathbf{r} \text{ to } \mathbf{r} + \mathbf{d} \mathbf{r} \\
&= \mathbf{a}(\mathbf{r}) \cdot \int_{\mathbf{r}_0}^{\mathbf{r} + \mathbf{d} \mathbf{r}} \text{d} \mathbf{r}' + O(|\mathbf{d} \mathbf{r}|^2) & \text{(choosing the straight line from } \mathbf{r} \text{ to } \mathbf{r} + \mathbf{d} \mathbf{r}) \\
&= \mathbf{a}(\mathbf{r}) \cdot \left[ (\mathbf{r} + \mathbf{d} \mathbf{r}) - \mathbf{r}_0 \right] = \mathbf{a}(\mathbf{r}) \cdot \text{d} \mathbf{r}
\end{align*}
\]

because \( \mathbf{a}(\mathbf{r}) \) is approximately constant between \( \mathbf{r} \) and \( \mathbf{r} + \mathbf{d} \mathbf{r} \), and the correction term \( O(|\mathbf{d} \mathbf{r}|^2) \) can be ignored as \( \mathbf{d} \mathbf{r} \to 0 \).

But \( \text{d} \phi = \nabla \phi \cdot \text{d} \mathbf{r} \) (by definition), and so, since \( \text{d} \mathbf{r} \) is arbitrary, we must have

\[
\mathbf{a}(\mathbf{r}) = \nabla \phi(\mathbf{r})
\]

The converse is much easier to prove. If \( \mathbf{a} = \nabla \phi \), then \( \nabla \times \mathbf{a} = \nabla \times (\nabla \phi) \equiv 0 \).

Therefore

\[
\nabla \times \mathbf{a} = 0 \iff \mathbf{a} = \nabla \phi
\]

To determine whether a vector field is conservative, one simply checks whether \( \nabla \times \mathbf{a} = 0 \) (in the region of interest).

**NB:** The scalar potential \( \phi(\mathbf{r}) \) is only determined up to a constant. If \( \psi = \phi + \text{constant} \) then \( \nabla \psi = \nabla \phi \), so \( \psi \) is an equally good potential. The freedom in the constant corresponds to the freedom in choosing \( \mathbf{r}_0 \) when calculating the potential. So \( \phi(\mathbf{r}_0) \) in equation (55) is just an irrelevant constant. Equivalently, the absolute value of a scalar potential has no meaning, only potential differences are significant.

### 13.3 Finding scalar potentials

**Method (1): Integration along a straight line**

We have shown that the scalar potential \( \phi(\mathbf{r}) \) for a conservative vector field \( \mathbf{a}(\mathbf{r}) \) can be constructed from a line integral which is independent of the path of integration between the endpoints. A convenient way of evaluating such integrals is to integrate along a straight line from \( \mathbf{r}_0 \) to \( \mathbf{r} \). Depending on the convergence of the integral, there are two standard choices:

(i) \( \mathbf{r}_0 = 0 \): If \( \phi(\mathbf{r}) \) is finite or zero at \( \mathbf{r} = 0 \), we parameterise the straight line by \( \mathbf{r}' = \lambda \mathbf{r} \) and

\[
\int_{\mathbf{r}_0}^{\mathbf{r}} \mathbf{a}(\mathbf{r}') \cdot \text{d} \mathbf{r}' = \int_{\lambda = 0}^{\lambda = 1} \mathbf{a}(\lambda \mathbf{r}) \cdot \mathbf{r} \text{d} \lambda
\]
with $0 \leq \lambda \leq 1$. Thus $dr' = d\lambda r$, and hence

$$
\phi(r) = \int_0^\lambda a(r') \cdot dr' = \int_0^\lambda a(\lambda r) \cdot r \, d\lambda,
$$

(ii) $|r_0| = \infty$: If $\phi(r)$ is finite or zero as $|r| \to \infty$, we again parameterise the straight line by $r' = \lambda r$, but this time with $1 \leq \lambda < \infty$. Again, we have $dr' = d\lambda r$, and hence

$$
\phi(r) = \int_\infty^\lambda a(r') \cdot dr' = \int_\lambda^\infty a(\lambda r) \cdot r \, d\lambda.
$$

Example 1: Let $a(r) = (2xy + z^3)e_x + x^2 e_y + 3xz^2 e_z$.

First check that $\nabla \times a = 0$, so the field is conservative (exercise). Then

$$
\phi(r) = \int_0^\lambda a(r') \cdot dr' = \int_0^\lambda a(\lambda r) \cdot r \, d\lambda
$$

\begin{align*}
&= \int_0^\lambda \left[ (2\lambda^2 xy + \lambda^3 z^3) x + (\lambda^2 x^2) y + (\lambda^3 3xz^2) z \right] d\lambda \\
&= \frac{2}{3} x^2 y + \frac{1}{4} xz^3 + \frac{1}{3} x^2 y + \frac{3}{4} xz^3 \\
&= x^2 y + xz^3
\end{align*}

NB: Always check that your potential $\phi(r)$ satisfies $a(r) = \nabla \phi(r)$. (Exercise)

Example 2: Let $a(r) = 2(c \cdot r) \overline{r} + r^2 \overline{c}$ where $c$ is a constant vector.

$$
\nabla \times a = 2 \left[ (c \cdot r) \times \overline{r} + (c \cdot r) \nabla \times \overline{r} + (\nabla r^2) \times c \right] = 2 [c \times \overline{r} + 0] + 2 \overline{r} \times c = 0
$$

Then

$$
\phi(r) = \int_0^\lambda a(r') \cdot dr' = \int_0^\lambda a(\lambda r) \cdot (d\lambda r)
$$

\begin{align*}
&= \int_0^\lambda (2(c \cdot \lambda r) \lambda \overline{r} + \lambda^2 \overline{r}^2 \overline{c}) \cdot r \, d\lambda = (2(c \cdot \overline{r}) \overline{r} \cdot \overline{c} + r^2 (c \cdot \overline{r})) \int_0^\lambda \lambda^2 d\lambda \\
&= r^2 (c \cdot \overline{r})
\end{align*}

Integration along a straight line is a straightforward and fairly elegant method, and it’s generally applicable.

Method (2): Direct integration

Since $a = \nabla \phi$, we have

$$
\partial_x \phi = a_x(x, y, z) \quad \partial_y \phi = a_y(x, y, z) \quad \partial_z \phi = a_z(x, y, z)
$$
We can integrate these equations separately to give

\[
\phi(x, y, z) = \int_a^x a_x(x', y, z) \, dx' + f(y, z)
\]

\[
\phi(x, y, z) = \int_a^y a_y(x, y', z) \, dy' + g(x, z)
\]

\[
\phi(x, y, z) = \int_a^z a_z(x, y, z') \, dz' + h(x, y)
\]

and then determine the “constants” of integration \(f(y, z), g(x, z)\) and \(h(x, y)\) by consistency.

**Example 1 (revisited):** Let \(a = (2xy + z^3)e_x + x^2e_y + 3xz^2e_z\). Then

\[
\phi = x^2y + xz^3 + f(y, z)
\]

\[
\phi = x^2y + g(x, z)
\]

\[
\phi = xz^3 + h(x, y)
\]

These agree if we choose \(f(y, z) = 0, g(x, z) = xz^3\) and \(h(x, y) = x^2y\), hence

\[
\phi(r) = x^2y + xz^3
\]

as before. This method is straightforward but it’s rather clumsy for problems such as Example 2, which is typical of many Physics applications.

**Method (3): Direct integration “by inspection” (guessing)**

Sometimes the result can be spotted directly.

For example, if \(a(r) = (c \cdot r)c\) where \(c\) is a constant vector, then

\[
a(r) = (c \cdot r)c = (c \cdot r)\nabla(c \cdot r) = \nabla \left( \frac{1}{2} (c \cdot r)^2 + \text{constant} \right)
\]

**Example 2 (revisited)**

\[
a(r) = 2(c \cdot r)c + r^2c = (c \cdot r)\nabla r^2 + r^2 \nabla (c \cdot r) = \nabla \left( (c \cdot r) r^2 + \text{constant} \right)
\]

in agreement with what we had before if we choose the integration constant to be zero.

### 13.4 Conservative forces: conservation of energy

We now show how the name *conservative field* arises in Physics. Let the vector field \(\mathbf{F}(\mathbf{r})\) (assumed time-independent) be the total force acting on a particle of mass \(m\) at position \(\mathbf{r}\). We will show that for a conservative/irrotational force, where we can write

\[
\mathbf{F}(\mathbf{r}) = -\nabla V(\mathbf{r}),
\]

the total energy is *constant* in time. Note that the force is *minus* the gradient of the (scalar) potential. The minus sign is conventional.
Proof: Let \( r(t) \) be the position vector of a particle at time \( t \). Denote the first and second derivatives of \( r \) with respect to time by \( \dot{r} \) (velocity) and \( \ddot{r} \) (acceleration) respectively. The particle moves under the influence of Newton’s second law (N2):

\[
m\ddot{r} = F(r)
\]

In time \( dt \) the particle moves from \( r \) to \( r + dr \). From N2, we get

\[
m\dot{r} \cdot dr = F(r) \cdot dr = -\nabla V(r) \cdot dr
\]

Integrating this expression along the path of the particle starting from \( r_A \) at time \( t_A \), to \( r_B \) at time \( t_B \), gives

\[
m \int_{t_A}^{t_B} \dot{r} \cdot dr = - \int_{t_A}^{t_B} \nabla V(r) \cdot dr
\]

We can evaluate the left-hand side of equation (56)

\[
m \int_{t_A}^{t_B} \dot{r} \cdot dr = m \int_{t_A}^{t_B} \frac{d}{dt} (\dot{r} \cdot r) dt = \frac{1}{2}m \left[ |\dot{r}|^2 \right]_{t_A}^{t_B} = \frac{1}{2}m (v_B^2 - v_A^2)
\]

where \( v_A \) and \( v_B \) are the magnitudes of the particle’s velocity at points \( A \) and \( B \) respectively. The right-hand side of equation (56) gives

\[
- \int_{t_A}^{t_B} \nabla V(r) \cdot dr = - \int_{t_A}^{t_B} dV = V_A - V_B
\]

where \( V_A \) and \( V_B \) are the values of the potential \( V \) at \( r_A \) and \( r_B \), respectively. Therefore

\[
\frac{1}{2}m (v_B^2 - v_A^2) = V_A - V_B
\]

Rearranging, we get

\[
\frac{1}{2}mv_A^2 + V_A = \frac{1}{2}mv_B^2 + V_B
\]

Hence the total energy, defined as \( E \equiv \frac{1}{2}mv^2 + V \), is conserved – it’s constant in time.

(Choosing \( F = +\nabla V \) would lead to \( E \equiv \frac{1}{2}mv^2 - V \), a less desirable convention.)

Examples: Newtonian gravity and the electrostatic force are both conservative. Frictional forces are not conservative: energy is dissipated and work is done in traversing a closed path. In general, time-dependent forces are not conservative.

We now return to where we started in section (1.3).

### 13.5 Gravitation and Electrostatics (revisited)

The foundation of Newtonian Gravity is Newton’s Law of Gravitation. The force \( F(r) \) on a particle of mass \( m_1 \) at \( r \) due to a particle of mass \( m \) situated at the origin is given (in SI units) by

\[
F(r) = -Gmm_1 \frac{r}{r^3}
\]
where $G = 6.67259(85) \times 10^{-11} \text{Nm}^2\text{kg}^{-2}$ is Newton’s Gravitational Constant.

The gravitational field $\mathbf{G}(r)$ due to the mass at the origin is defined by

$$\mathbf{F}(r) \equiv m_1 \mathbf{G}(r) \quad \text{or} \quad \mathbf{G}(r) = -G m_1 \frac{r}{r^3}$$

where the test mass $m_1$ is so small that its gravitational field can be ignored. The gravitational field is conservative because

$$\nabla \times \left( \frac{r}{r^3} \right) = \nabla \left( \frac{1}{r^3} \right) \times r + \frac{1}{r^3} (\nabla \times r) = \left( -\frac{3r}{r^5} \right) \times r + 0 = 0$$

The gravitational potential defined by

$$\mathcal{G} = -\nabla \phi$$

can be obtained from equation (57) by spotting the direct integration, $\nabla \left( \frac{1}{r} \right) = -\frac{1}{r^3}$, giving

$$\phi = -\frac{G m}{r}$$

Alternatively, we may evaluate it explicitly by a line integral. Choosing $r_0$ at infinity gives

$$\phi(r) = -\int_{r_0}^{r} \mathcal{G}(r') \cdot dr' = -\int_{\infty}^{1} \mathcal{G}(\lambda r) \cdot d\lambda r = (-)^2 \int_{\infty}^{1} \frac{G m \cdot \lambda r \cdot r}{\lambda^2} \cdot d\lambda = -\frac{G m}{r}$$

**Note:** In this example, the vector field $\mathcal{G}$ is singular at the origin $r = 0$. This implies that we have to exclude the origin, so it’s not possible to obtain the scalar potential at $r$ by integration along a path from the origin. Instead we integrate from infinity, which in turn means that the gravitational potential at infinity is zero.

**Note:** Since $\mathbf{F} = m_1 \mathbf{G} = -\nabla (m_1 \phi)$, the potential energy of the mass $m_1$ is $V(r) = m_1 \phi(r)$. The distinction (a convention) between potential and potential energy is a common source of confusion.

**Electrostatics:** Coulomb’s Law states that the force $\mathbf{F}(r)$ on a particle of charge $q_1$ situated at $r$ in the electric field $\mathbf{E}(r)$ due to a particle of charge $q$ situated at the origin is given (in SI units) by

$$\mathbf{F} = q_1 \mathbf{E} = \frac{q_1 q}{4\pi\epsilon_0 r^3} \cdot ,$$

where $\epsilon_0 = 10^7/(4\pi c^2) = 8.854187 817 \cdots \times 10^{-12} \text{C}^2\text{N}^{-1}\text{m}^{-2}$ is called the permittivity of free space. Again the test charge $q_1$ is taken as small, so as not to disturb the electric field.

The electrostatic potential may be obtained by inspection, or by integrating $\mathbf{E} = -\nabla \phi$ from infinity to $r$.

$$\phi(r) = \frac{q}{4\pi\epsilon_0 r}$$

**Note** that electrostatics and gravitation are very similar mathematically, the only real difference being that the gravitational force between two masses is always attractive, whereas like charges repel.
13.6 The equations of Poisson and Laplace

In section [11.6], we derived Gauss’ Law and Maxwell’s first equation (ME1) for the electrostatic field

\[ \int_S E \cdot dS = \frac{Q}{\varepsilon_0} \quad \text{and} \quad \nabla \cdot E(r) = \frac{\rho(r)}{\varepsilon_0} \]

where \( \rho(r) \) is the charge density at \( r \), and \( Q = \int_V \rho(r) \, dV \) is the total charge in volume \( V \).

Writing \( E(r) = -\nabla \phi(r) \) and using Maxwell’s first equation gives Poisson’s equation

\[ \nabla^2 \phi = -\frac{\rho}{\varepsilon_0} \]

If \( \rho(r) = 0 \) everywhere in some region, we have

\[ \nabla^2 \phi = 0 \]

which is Laplace’s equation.

These partial differential equations are important in many branches of Physics and Mathematics. You will study (and solve) them next year.

14 The vector potential

We have shown that an irrotational vector field \( a(r) \), i.e. one that satisfies \( \nabla \times a = 0 \), can be written as the gradient of a scalar field, \( a = \nabla \phi \).

Under what conditions can we write a vector field \( B(r) \) as the curl of a vector field \( A(r) \)?

(i) If \( \nabla \cdot B = 0 \), it can be shown that a vector field \( A \) can be found such that \( B = \nabla \times A \).

(ii) The converse is easy to prove. If the field \( B \) can be written as \( B = \nabla \times A \), then

\[ \nabla \cdot B = \nabla \cdot (\nabla \times A) = 0 \]

because ‘div curl’ is always zero.

Hence

\[ \nabla \cdot B(r) = 0 \quad \Leftrightarrow \quad \text{There exists a field} \ A(r) \ \text{such that} \ B(r) = \nabla \times A(r) \]

The field \( A \) is called the vector potential for the solenoidal field \( B \).

For such a field \( B \), which is finite or zero at the origin, it can be shown that

\[ A(r) = -r \times \int_0^1 B(\lambda r) \, \lambda \, d\lambda \]  

(59)

is a vector potential for \( B(r) \).

**Example:** Find a vector potential \( A \) for the field \( B = c \times r \), where \( c \) is a constant vector.
It is easy to show that $\nabla \cdot B = 0$ (exercise). Equation (59) then gives

$$A(r) = -r \times \int_0^1 B(\lambda r) \, \lambda \, d\lambda = -r \times \int_0^1 (\mathbf{c} \times \lambda r) \lambda \, d\lambda$$

$$= -\left( r^2 \mathbf{c} - (r \cdot \mathbf{c}) r \right) \int_0^1 \lambda^2 \, d\lambda = \frac{1}{3} \left( (r \cdot \mathbf{c}) r - r^2 \mathbf{c} \right)$$

You should always check at the end that $A$ satisfies $B = \nabla \times A$ (exercise).

**Gauge invariance:** We can always add the gradient of an arbitrary scalar field $f(r)$ to the vector potential

$$A(r) \rightarrow A'(r) = A(r) + \nabla f(r)$$

without changing $B$.

We have $\nabla \times A' = \nabla \times A + 0$ because $\nabla \times \nabla f = 0$ for any scalar field (‘curl grad’ is always zero). This is called *gauge invariance* in electromagnetism, and is one of the most important symmetries in Physics. (*Electroweak gauge invariance* is (partly) broken by the Higgs field.)

### 14.1 Physical examples of vector potentials

**Magnetism:** Magnetic field lines do not have sources or sinks – the lines of a magnetic field are continuous. For example, the magnetic field lines around a straight current-carrying wire are *circles*. It can be shown that the magnetic field $B$ satisfies $\nabla \cdot B = 0$ everywhere, so we can write $B = \nabla \times A$, where $A$ is the *magnetic vector potential* (tutorial).

**Fluid mechanics:** For incompressible fluids (which have constant mass density, $\rho$) with no sources or sinks, we showed in section (11.5) that the velocity field $v(r)$ satisfies $\nabla \cdot v = 0$. In this case, there exists a velocity potential $\Psi(r)$ such that $v = \nabla \times \Psi$.

### 15 Orthogonal curvilinear coordinates

As we have seen, it is often convenient to work with coordinate systems other than Cartesian coordinates $\{x_i\}$, i.e. $(x_1, x_2, x_3)$ or $(x, y, z)$.

For example, spherical polar coordinates $(r, \theta, \phi)$ are defined by:

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

We shall set up a formalism to deal with rather general coordinate systems, of which spherical polars are a very important example.
Suppose we make a transformation from the Cartesian coordinates \((x_1, x_2, x_3)\) to the variables \((u_1, u_2, u_3)\), which are functions of the \(\{x_i\}\)

\[
    u_1 = u_1(x_1, x_2, x_3) \\
    u_2 = u_2(x_1, x_2, x_3) \\
    u_3 = u_3(x_1, x_2, x_3)
\]

If the variables \(\{u_i\}\) are single-valued functions of the variables \(\{x_i\}\), then we can make the inverse transformations,

\[
x_i = x_i(u_1, u_2, u_3) \quad \text{for } i = 1, 2, 3,
\]

except possibly at certain points.

A point may be specified by its Cartesian coordinates \(\{x_i\}\), or its curvilinear coordinates \(\{u_i\}\).

We may define the curvilinear coordinates by equations giving \(\{x_i\}\) as functions of \(\{u_i\}\), or vice versa.

- For Cartesian coordinates, the surfaces ‘\(x_i = \text{constant}\)’ \((i = 1, 2, 3)\) are planes, with (constant) normal vectors \(\boldsymbol{e}_i\) (the Cartesian basis vectors) intersecting at right angles.

- For curvilinear coordinates, the surfaces ‘\(u_i = \text{constant}\)’ do not, in general, have constant normal vectors, nor do they intersect at right angles. For example, in 2-D, we might have

From the definition of spherical polar coordinates \((r, \theta, \phi)\), we have

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

The surfaces of constant \(r, \theta, \) and \(\phi\) are

\[
    r = \text{constant} \Rightarrow \text{spheres centred at the origin} \\
    \theta = \text{constant} \Rightarrow \text{cones of semi-angle } \theta \text{ and axis along the } z\text{-axis} \\
    \phi = \text{constant} \Rightarrow \text{planes passing through the } z\text{-axis}
\]

Not all of these surfaces are planes, but they \textit{do} intersect at right angles.

\[\text{23Sometimes the curvilinear coordinates are called } (u, v, w), \text{ just as Cartesians are called } (x, y, z).\]
15.1 Orthogonal curvilinear coordinates

If the coordinate surfaces (surfaces of constant $u_i$), intersect at right angles, as in the above example of spherical polars, the curvilinear coordinates are said to be orthogonality.

Scale factors and basis vectors: Suppose the point $P$ has position vector $\mathbf{r} = \mathbf{r}(u_1, u_2, u_3)$. If we change curvilinear coordinate $u_1$ by $du_1$ (with $u_2$ and $u_3$ fixed), then $\mathbf{r} \rightarrow \mathbf{r} + d\mathbf{r}$, with

$$d\mathbf{r} = \partial_{u_1} \mathbf{r} \, du_1 \equiv h_1 \mathbf{e}_1 \, du_1$$

where we have defined the scale factor $h_1$ and the unit vector $\mathbf{e}_1$ by

$$h_1 = |\partial_{u_1} \mathbf{r}| \quad \text{and} \quad \mathbf{e}_1 = \frac{1}{h_1} \partial_{u_1} \mathbf{r}$$

- The scale factor $h_1$ gives the length $h_1 \, du_1$ of $d\mathbf{r}$ when we change $u_1 \rightarrow u_1 + du_1$.
- $\mathbf{e}_1$ is a unit vector in the direction of increasing $u_1$ (with fixed $u_2$ and $u_3$).

Similarly, we can define $h_i$ and $\mathbf{e}_i$ for $i = 2$ and 3.

In general, if we change a single $u_i$, keeping the other two fixed, we have

$$\partial_{u_i} \mathbf{r} = h_i \mathbf{e}_i \quad i = 1, 2, 3$$

- The unit vectors $\{\mathbf{e}_i\}$ are in general not constant vectors – their directions depend on the position vector $\mathbf{r}$, and hence on the curvilinear coordinates $\{u_i\}$. [They should perhaps be called $\{\mathbf{e}_{u_i}\}$ or $\{\mathbf{e}_x, \mathbf{e}_y, \mathbf{e}_z\}$ to avoid confusion with Cartesian basis vectors.]
- If the curvilinear unit vectors satisfy $\mathbf{e}_i \cdot \mathbf{e}_j = \delta_{ij}$, the $\{u_i\}$ are said to be orthogonal curvilinear coordinates, and the three unit vectors $\{\mathbf{e}_i\}$ form an orthonormal basis.

![Diagram of orthogonal curvilinear coordinates]

15.1.1 Examples of orthogonal curvilinear coordinates (OCCs)

Cartesian coordinates:

$$\mathbf{r} = x \mathbf{e}_x + y \mathbf{e}_y + z \mathbf{e}_z \quad \Rightarrow \quad h_x \mathbf{e}_x = \frac{\partial \mathbf{r}}{\partial x} = \mathbf{e}_x, \text{ etc.}$$

The scale factors are all unity, and the unit vectors point in the same direction everywhere.
Spherical polar coordinates: \( u_1 = r, u_2 = \theta, u_3 = \phi \) (in that order)

\[
\begin{align*}
    \vec{r} &= r \sin \theta \cos \phi \hat{e}_x + r \sin \theta \sin \phi \hat{e}_y + r \cos \theta \hat{e}_z \\
    \partial_r \vec{r} &= \sin \theta \cos \phi \hat{e}_x + \sin \theta \sin \phi \hat{e}_y + \cos \theta \hat{e}_z \\
    \partial_\theta \vec{r} &= r \cos \theta \cos \phi \hat{e}_x + \cos \theta \sin \phi \hat{e}_y - r \sin \theta \hat{e}_x \\
    \partial_\phi \vec{r} &= -r \sin \theta \sin \phi \hat{e}_x + r \sin \theta \cos \phi \hat{e}_y
\end{align*}
\]

Hence the unit vectors for spherical polars are

\[
\begin{align*}
    \hat{e}_r &= \frac{\sin \theta \cos \phi \hat{e}_x + \sin \theta \sin \phi \hat{e}_y + \cos \theta \hat{e}_z}{r} \\
    \hat{e}_\theta &= \cos \theta \cos \phi \hat{e}_x + \cos \theta \sin \phi \hat{e}_y - \sin \theta \hat{e}_z \\
    \hat{e}_\phi &= -\sin \phi \hat{e}_x + \cos \phi \hat{e}_y
\end{align*}
\]

These unit vectors are normal to the surfaces described above (spheres, cones and planes).

They are orthogonal:

\[
\hat{e}_r \cdot \hat{e}_\theta = \hat{e}_r \cdot \hat{e}_\phi = \hat{e}_\theta \cdot \hat{e}_\phi = 0
\]

And they form a right-handed orthonormal basis:

\[
\hat{e}_r \times \hat{e}_\theta = \hat{e}_\phi, \quad \hat{e}_\theta \times \hat{e}_\phi = \hat{e}_r, \quad \hat{e}_\phi \times \hat{e}_r = \hat{e}_\theta
\]

See also tutorial question (7.6).

Cylindrical coordinates: \( u_1 = \rho, u_2 = \phi, u_3 = z \) (in that order)

\[
\begin{align*}
    \vec{r} &= \rho \cos \phi \hat{e}_x + \rho \sin \phi \hat{e}_y + z \hat{e}_z \\
    \partial_\rho \vec{r} &= \cos \phi \hat{e}_x + \sin \phi \hat{e}_y \\
    \partial_\phi \vec{r} &= -\rho \sin \phi \hat{e}_x + \rho \cos \phi \hat{e}_y \\
    \partial_z \vec{r} &= \hat{e}_z
\end{align*}
\]

The scale factors are then (tutorial) \( h_\rho = 1, h_\phi = \rho, h_z = 1 \), and the basis vectors are

\[
\begin{align*}
    \hat{e}_\rho &= \cos \phi \hat{e}_x + \sin \phi \hat{e}_y \\
    \hat{e}_\phi &= -\sin \phi \hat{e}_x + \cos \phi \hat{e}_y \\
    \hat{e}_z &= \hat{e}_z
\end{align*}
\]

These unit vectors are normal to surfaces which are (respectively): cylinders centred on the \( z \)-axis (\( \rho = \text{constant} \)), planes through the \( z \)-axis (\( \phi = \text{constant} \)), planes perpendicular to the \( z \)-axis (\( z = \text{constant} \)), and they are clearly orthonormal.

15.2 Elements of length, area and volume in OCCs

Length: If we change \( u_1 \rightarrow u_1 + du_1 \), keeping \( u_2 \) and \( u_3 \) fixed, then \( \vec{r} \rightarrow \vec{r} + d\vec{r} \) where

\[
d\vec{r}_1 = h_1 \, du_1
\]

The infinitesimal element of length along \( \hat{e}_1 \) is \( h_1 \, du_1 \).

Clearly, the infinitesimal elements of length along the three curvilinear basis vectors \( \hat{e}_1, \hat{e}_2 \) and \( \hat{e}_3 \), respectively, are

\[
h_1 \, du_1 \quad h_2 \, du_2 \quad h_3 \, du_3
\]

If we change all three of the \( \{u_i\} \), then

\[
d\vec{r} = h_1 \, du_1 \hat{e}_1 + h_2 \, du_2 \hat{e}_2 + h_3 \, du_3 \hat{e}_3
\]
Arc length: If \( ds \) is the length of the infinitesimal vector \( dr \), then \((ds)^2 = dr \cdot dr\).

In Cartesian coordinates
\[
(ds)^2 = (dx)^2 + (dy)^2 + (dz)^2
\]

In curvilinear coordinates, if we change the \( i \)th coordinate \( u_i \) by \( du_i \), then
\[
(ds)^2 = dr \cdot dr = (h_i e_i du_i) \cdot (h_j e_j du_j) = h_i h_j (e_i \cdot e_j) \ du_i du_j,
\]
with Einstein summation convention implied as usual. For orthogonal curvilinear coordinates, we have \( e_i \cdot e_j = \delta_{ij} \), which tells us to set \( j = i \), and we can perform the sum over \( j \). Hence
\[
(ds)^2 = h_i^2 (du_i)^2 = h_1^2 \ du_1^2 + h_2^2 \ du_2^2 + h_3^2 \ du_3^2
\]

For spherical polars, \( h_r = 1 \), \( h_\theta = r \), \( h_\phi = r \sin \theta \), therefore
\[
(ds)^2 = (dr)^2 + r^2 (d\theta)^2 + r^2 \sin^2 \theta \ (d\phi)^2
\]

Vector Area:

If we let \( u_1 \rightarrow u_1 + du_1 \), then
\[
 r \rightarrow r + dr_1 \ 	ext{with} \ dr_1 = h_1 e_1 du_1
\]

If we let \( u_2 \rightarrow u_2 + du_2 \), then
\[
 r \rightarrow r + dr_2 \ 	ext{with} \ dr_2 = h_2 e_2 du_2
\]

The vector area of the infinitesimal parallelogram (actually a rectangle for OCCs) whose sides are the vectors \( dr_1 \) and \( dr_2 \) is
\[
dS = (dr_1) \times (dr_2) = (h_1 du_1 e_1) \times (h_2 du_2 e_2) = h_1 h_2 du_1 du_2 e_3,
\]
because \( e_1 \times e_2 = e_3 \) for orthogonal systems. Clearly, \( dS \) points in the direction of \( e_3 \) which is normal to the surfaces \( u_3 \) =constant.

The vector areas \( dS_1 \) and \( dS_2 \) are defined similarly.

Example: For spherical polars, if we vary \( \theta \) and \( \phi \), keeping \( r \) fixed, we obtain very easily the familiar result
\[
dS_r = (h_\theta d\theta e_\theta) \times (h_\phi d\phi e_\phi) = h_\theta h_\phi d\theta d\phi e_r = r^2 \sin \theta d\theta d\phi e_r
\]

Similarly, if we vary \( \phi \) and \( r \), keeping \( \theta \) fixed, we obtain the vector element of area on the cone of semi-angle \( \theta \), with its axis along the \( z \) axis
\[
dS_\theta = (h_\phi d\phi e_\phi) \times (h_r dr e_r) = h_\phi h_r d\phi dr e_\theta = r \sin \theta dr d\phi e_\theta
\]
and similarly for \( dS_\phi \). See also tutorial question (7.6).
**Volume:** The volume of the infinitesimal parallelepiped (actually a cuboid for OCCs) with edges $dr_1$, $dr_2$ and $dr_3$ is

$$dV = (dr_1 \times dr_2) \cdot dr_3 = (h_1 du_1 \xi_1) \times (h_2 du_2 \xi_2) \cdot (h_3 du_3 \xi_3)$$

which equals

$$dV = h_1 h_2 h_3 du_1 du_2 du_3$$  \hspace{1cm} (60)

because $(\xi_1 \times \xi_2) \cdot \xi_3 = 1$ for orthogonal curvilinear coordinates.

For spherical polars, we have

$$dV = h_r h_\theta h_\phi dr d\theta d\phi = r^2 \sin \theta dr d\theta d\phi$$

15.3 Components of a vector field in curvilinear coordinates

A vector field $\mathbf{a}(r)$ can be expressed in terms of *curvilinear components* $a_i$, defined by

$$a_i = \sum_{i=1}^{3} a_i(u_1, u_2, u_3) \xi_i$$

where $\xi_i$ is the $i^{th}$ curvilinear basis vector (which again should really be called $\xi_{u_i}$ to avoid confusion with the Cartesian basis vectors.)

For orthogonal curvilinear coordinates, the component $a_i$ can be obtained by taking the scalar product of $\mathbf{a}$ with the $i^{th}$ curvilinear basis vector $\xi_i$

$$a_i = a(r) \cdot \xi_i$$

NB $a_i$ must be expressed in terms of $u_1$, $u_2$, $u_3$ (not $x$, $y$, $z$) when working in the $\{u_i\}$ basis.

**Example:** If $\mathbf{a} = a \xi_x$ in Cartesians, then in spherical polars

$$a_r = \mathbf{a} \cdot \xi_r = (a \xi_x) \cdot (\sin \theta \cos \phi \xi_x + \sin \theta \sin \phi \xi_y + \cos \theta \xi_z) = a \sin \phi \cos \phi$$

Similarly, $a_\theta = a \cdot \xi_\theta$ and $a_\phi = a \cdot \xi_\phi$, and we obtain $\mathbf{a}$ in the spherical-polar basis (exercise)

$$a(r, \theta, \phi) = a (\sin \theta \cos \phi \xi_r + \cos \theta \cos \phi \xi_\theta - \sin \phi \xi_\phi)$$

- You can often spot the curvilinear components “by inspection”. See tutorial question (5.6) for an example of this.
- In general, one chooses the set of coordinates which matches most closely the *symmetry* of the problem.

15.4 Div, grad, curl and the Laplacian in orthogonal curvilinears

Below we obtain the expression of the expressions in the title in orthogonal curvilinear coordinates. In section 15.3 the expressions are collected in table and Eq. (69) for cylindrical and spherical coordinates respectively.
15.4.1 Gradient

In section (2) we defined the gradient in terms of the change in a scalar field \( f(r) \) when we let \( r \to r + dr \)
\[
\frac{df(r)}{dr} = \nabla f(r) \cdot dr \tag{61}
\]
Now consider writing \( f(r) \) in terms of orthogonal curvilinear coordinates, \( f(r) = f(u_1, u_2, u_3) \).
As usual, we denote the curvilinear basis vectors by \( \{e_1, e_2, e_3\} \).
Let \( u_1 \to u_1 + du_1, \ u_2 \to u_2 + du_2, \) and \( u_3 \to u_3 + du_3 \).
Using Taylor’s theorem, we have
\[
df = \partial_{u_1} f du_1 + \partial_{u_2} f du_2 + \partial_{u_3} f du_3 \tag{62}
\]
We can manipulate the RHS of this equation into the form of equation (61). Start with
\[
dr = h_1 du_1 e_1 + h_2 du_2 e_2 + h_3 du_3 e_3
\]
Now use orthogonality of the curvilinear basis vectors, \( e_i \cdot e_j = \delta_{ij} \), to rewrite equation (62) as
\[
df = \left( \frac{1}{h_1} e_1 \partial_{u_1} f + \frac{1}{h_2} e_2 \partial_{u_2} f + \frac{1}{h_3} e_3 \partial_{u_3} f \right) \cdot (h_1 du_1 + h_2 du_2 + h_3 du_3)
\]
Comparing this result with equation (61), which holds for all \( dr \), we obtain the del operator in curvilinear coordinates
\[
\nabla = \frac{1}{h_i} e_i \partial_{x_i} \tag{63}
\]
which reduces for Cartesian coordinates, with scale factors \( h_i = 1 \), to the expression (28).
For reference we quote the action of \( \nabla \) on a scalar function \( f \),
\[
\nabla f = \frac{1}{h_1} e_1 \partial_{u_1} f + \frac{1}{h_2} e_2 \partial_{u_2} f + \frac{1}{h_3} e_3 \partial_{u_3} f = \frac{1}{h_i} e_i \partial_{u_i} f \tag{64}
\]

15.4.2 Divergence

Let \( a(r) \) be a vector field which we write in orthogonal curvilinear coordinates as
\[
a(r) = \sum_{i=1}^{3} a_i(u_1, u_2, u_3) e_i
\]
\(^{24}\) We use \( f(r) \) rather than \( \phi(r) \) here in order to avoid confusion with the angle \( \phi \) in spherical polars.
\(^{25}\) \( a \) must be continuously differentiable.
where \( a_i \) are the components of \( \mathbf{a} \) in the curvilinear basis, and \( \mathbf{e}_i \) is the \( i^{th} \) curvilinear basis vector.

We obtain \( \nabla \cdot \mathbf{a} \) in orthogonal curvilinear systems using the integral definition of divergence

\[
\nabla \cdot \mathbf{a} = \lim_{\delta V \to 0} \frac{1}{\delta V} \int_\delta S \mathbf{a} \cdot d\mathbf{S},
\]

where \( \delta S \) is the closed surface bounding \( \delta V \).

Let the point \( P \) have curvilinear coordinates \((u_1, u_2, u_3)\).

Choose \( \delta V \) to be a small “cuboid” with its three edges \( \{\delta r_i\} \) along the basis vectors \( \{\mathbf{e}_i\} \) at \( P \):

\[
\begin{align*}
\delta r_1 &= h_1 \delta u_1 \mathbf{e}_1 \\
\delta r_2 &= h_2 \delta u_2 \mathbf{e}_2 \\
\delta r_3 &= h_3 \delta u_3 \mathbf{e}_3
\end{align*}
\]

The outward element of area on the face \( ABCD \) is \( d\mathbf{S} = +h_2 h_3 \, du_2 \, du_3 \, \mathbf{e}_1 \)

The outward element of area on the face \( PQRS \) is \( d\mathbf{S} = -h_2 h_3 \, du_2 \, du_3 \, \mathbf{e}_1 \)

The contributions to the surface integral from the faces \( ABCD \) and \( PQRS \) are then

\[
\begin{align*}
\int_{u_3}^{u_3 + \delta u_3} \int_{u_2}^{u_2 + \delta u_2} & \left\{ [a_1 h_2 h_3]_{ABCD} - [a_1 h_2 h_3]_{PQRS} \right\} \, du_2 \, du_3 \\
= & \int \int \left\{ [a_1 h_2 h_3]_{(u_1+\delta u_1, u_2, u_3)} - [a_1 h_2 h_3]_{(u_1, u_2, u_3)} \right\} \, du_2 \, du_3 \\
= & \int \int \left\{ \delta u_1 \left[ \partial_{u_1} (a_1 h_2 h_3) \right]_{(u_1, u_2, u_3)} \right\} \, du_2 \, du_3 \quad \text{(by Taylor’s theorem)} \\
= & \delta u_1 \, \delta u_2 \, \delta u_3 \, \left[ \partial_{u_1} (a_1 h_2 h_3) \right]_{(u_1, u_2, u_3)}
\end{align*}
\]

(65)

In the last step, we assumed that \( \delta V \) is small enough that the integrand is approximately constant over the range of integration. We then approximated the integrals over \( u_2 \) and \( u_3 \) by the integrand evaluated at the point \( P \),

\[
\delta u_1 \left[ \partial_{u_1} (a_1 h_2 h_3) \right]_{(u_1, u_2, u_3)}
\]

multiplied by the ranges of integration \( \delta u_2 \, \delta u_3 \).

The contributions of the other four faces to the integral over \( \delta S \) can be obtained similarly, or by cyclic permutations of the indices \( \{1, 2, 3\} \) in equation (65).

Dividing by the volume of the cuboid \( \delta V = h_1 h_2 h_3 \, \delta u_1 \, \delta u_2 \, \delta u_3 \), we obtain our final expression for \( \nabla \cdot \mathbf{a} \) in orthogonal curvilinear coordinates

\[
\nabla \cdot \mathbf{a} = \frac{1}{h_1 h_2 h_3} \left\{ \partial_{u_1} (a_1 h_2 h_3) + \partial_{u_2} (a_2 h_3 h_1) + \partial_{u_3} (a_3 h_1 h_2) \right\}
\]

(66)
For Cartesian coordinates, the scale factors are all unity, and we recover the usual expression for $\nabla \cdot a$ in Cartesians.

15.4.3 Curl

We obtain $\nabla \times a$ in orthogonal curvilinear coordinates using the line integral definition of curl.

The component of $\nabla \times a$ in the direction of the unit vector $\mathbf{n}$ is

$$n \cdot (\nabla \times a) = \lim_{\delta S \to 0} \frac{1}{\delta S} \oint_{\delta C} a \cdot dr$$

where $\delta S$ is a small planar surface, with unit normal $\mathbf{n}$, bounded by the closed curve $\delta C$.

Let $\delta S$ be a small rectangular surface parallel to the $\mathbf{e}_2 - \mathbf{e}_3$ plane with one corner at $r(u_1, u_2, u_3)$, and with edges

$$\delta r_2 = h_2 \Delta u_2 \mathbf{e}_2 \quad \text{and} \quad \delta r_3 = h_3 \Delta u_3 \mathbf{e}_3$$

which lie along the basis vectors, so that $\mathbf{n} = \mathbf{e}_1$.

The line integral around the curve $\delta C$ is the sum of the line integrals along the lines $1 \to 4$ respectively,

$$\oint_{\delta C} a \cdot dr = \int [a_2 h_2]_{(u_1, u_2, u_3)} du_2 + \int [a_3 h_3]_{(u_1, u_2 + \Delta u_2, u_3)} du_3
- \int [a_2 h_2]_{(u_1, u_2, u_3 + \Delta u_3)} du_2 - \int [a_3 h_3]_{(u_1, u_2, u_3)} du_3$$

Using Taylor’s theorem, we can write this as

$$\oint_{\delta C} a \cdot dr = \int \left\{ \delta u_2 \left[ \partial u_2 (a_3 h_3) \right]_{(u_1, u_2, u_3)} \right\} du_3 - \int \left\{ \delta u_3 \left[ \partial u_3 (a_2 h_2) \right]_{(u_1, u_2, u_3)} \right\} du_2$$

In each case, we approximate the integrals over $u_3$ and $u_2$ by the product of the integrand and the integration ranges $\Delta u_3$ and $\Delta u_2$, respectively.

$$\oint_{\delta C} a \cdot dr = \partial u_2 (a_3 h_3) \delta u_2 \delta u_3 - \partial u_3 (a_2 h_2) \delta u_3 \delta u_2$$

where all the $\{a_i\}$ and $\{h_i\}$ are evaluated at $r(u_1, u_2, u_3)$.

Finally, we divide by the area of the rectangle $\delta S = h_2 h_3 \delta u_2 \delta u_3$, to obtain

$$\mathbf{e}_1 \cdot (\nabla \times a) = (\nabla \times a)_1 = \frac{1}{h_2 h_3} \left\{ \partial u_2 (a_3 h_3) - \partial u_3 (a_2 h_2) \right\}$$

The components of $\nabla \times a$ in the directions of the curvilinear basis vectors $\mathbf{e}_2$ and $\mathbf{e}_3$ may be obtained similarly, or by cyclic permutations of the indices.

It is convenient to write the final result in the form
15.4.4 Laplacian of a scalar field

The action of the Laplacian operator on a scalar field $f(r)$ is defined by $\nabla^2 f = \nabla \cdot (\nabla f)$. Using the expression for $\nabla \cdot a$, with $a = \nabla f$, derived above, we find

$$\nabla^2 f = \frac{1}{h_1 h_2 h_3} \left\{ \partial_{u_1} \left( \frac{h_2 h_3}{h_1} \partial_{u_1} f \right) + \partial_{u_2} \left( \frac{h_3 h_1}{h_2} \partial_{u_2} f \right) + \partial_{u_3} \left( \frac{h_1 h_2}{h_3} \partial_{u_3} f \right) \right\}$$  \hspace{1cm} (67)

15.4.5 Laplacian of a vector field

We may define the Laplacian of a vector field $a(r)$ in curvilinear coordinates by rearranging identity (J) in section 5.5 $\nabla^2 a = \nabla \left( \nabla \cdot a \right) - \nabla \times \left( \nabla \times a \right)$. We shall refrain at this point from quoting explicit expressions.

The expression for the Laplacian of a scalar field in spherical polars is one of the most important results in the course, with applications in Quantum Mechanics, Electromagnetism, Optics, Meteorology, Fluid/Solid Mechanics, Cosmology, . . .
15.5 $\nabla$-expression in cylindrical and spherical coordinates

To this very end we summarise in table 1 the volume element $dV$ (60), the gradient on scalars and vector fields $\nabla f$ (64) and $\nabla \cdot a$ (66) and the Laplacian on a scalar $\nabla^2 f$ (67) in terms of cylindrical coordinates and spherical coordinates presented in section 15.1.1. The cross product on vector $\nabla a$ is given in (69) in cylindrical and spherical coordinates respectively. The vector $a$ is written in cylindrical and spherical coordinates as

$$
a = a_\rho e_\rho + a_\phi e_\phi + a_z e_z \quad \text{(cylindrical)}
$$

$$
a = a_r e_r + a_\theta e_\theta + a_\phi e_\phi \quad \text{(spherical)}
$$

where for cylindrical coordinates

$$
e_\rho = (\cos \phi, \sin \phi, 0), \quad e_\phi = (-\sin \phi, \cos \phi, 0), \quad e_z = (0, 0, 1)
$$

and for spherical coordinates

$$
e_r = (\sin \theta \cos \phi, \sin \theta \sin \phi, \cos \theta), \quad e_\theta = (\cos \theta \cos \phi, \cos \theta \sin \phi, -\sin \theta),
$$

$$
e_\phi = (-\sin \theta \sin \phi, \sin \theta \cos \phi, 0).
$$

<table>
<thead>
<tr>
<th>$h_i$</th>
<th>$h_\rho = 1$, $h_\phi = \rho$, $h_z = 1$</th>
<th>$h_r = 1$, $h_\theta = r$, $h_\phi = r \sin \theta$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$dV$</td>
<td>$\rho , dr , d\theta , d\phi$</td>
<td>$r^2 \sin \theta , dr , d\theta , d\phi$</td>
</tr>
<tr>
<td>$\nabla f$</td>
<td>$e_\rho \frac{\partial f}{\partial \rho} + \frac{1}{\rho} \frac{\partial f}{\partial \theta} + \frac{\partial f}{\partial z}$</td>
<td>$e_r \frac{\partial f}{\partial r} + \frac{\partial f}{\partial \theta} + \frac{1}{r \sin \theta} \frac{\partial f}{\partial \phi}$</td>
</tr>
<tr>
<td>$\nabla \cdot a$</td>
<td>$\frac{1}{\rho} \frac{\partial}{\partial \rho} (\rho a_\rho) + \frac{1}{\rho} \frac{\partial}{\partial \theta} a_\phi + \frac{\partial}{\partial z} a_z$</td>
<td>$\frac{1}{r} \frac{\partial}{\partial r} (r^2 a_r) + \frac{1}{r \sin \theta} \left{ \frac{\partial}{\partial \theta} (\sin \theta a_\theta) + \frac{\partial}{\partial \phi} (a_\phi) \right}$</td>
</tr>
<tr>
<td>$\nabla^2 f$</td>
<td>$\frac{1}{\rho} \frac{\partial}{\partial \rho} \left( \frac{\partial f}{\partial \rho} \right) + \frac{1}{\rho^2} \frac{\partial^2 f}{\partial \rho^2} + \frac{\partial^2 f}{\partial z^2}$</td>
<td>$\frac{1}{r^2} \frac{\partial}{\partial r} \left( \frac{r^2 \partial f}{\partial r} \right) + \frac{1}{r^2 \sin^2 \theta} \left{ \sin \theta \frac{\partial}{\partial \theta} \left( \sin \theta \frac{\partial f}{\partial \theta} \right) + \frac{\partial^2 f}{\partial \phi^2} \right}$</td>
</tr>
</tbody>
</table>

Table 1: Summary of del-expression for cylindrical and spherical coordinates which can be obtained from the generic formula for OCC described in the text above.

For cylindrical coordinates

$$
\nabla \times a|_{\text{cyl}} = \frac{1}{\rho} \begin{vmatrix}
    e_r & \rho e_\phi & e_z \\
    \partial_\rho & \partial_\phi & \partial_z \\
    a_\rho & \rho a_\phi & a_z
\end{vmatrix} = e_r \left( \frac{\partial a_z}{\partial \rho} - \frac{\partial a_\phi}{\partial z} \right) + e_\phi \left( \frac{\partial a_z}{\partial \phi} - \frac{\partial a_\rho}{\partial z} \right) + e_z \left( \frac{\partial (\rho a_\phi)}{\partial \rho} - \frac{\partial a_\phi}{\partial \phi} \right)
$$

For spherical coordinates

$$
\nabla \times a|_{\text{spherical}} = \frac{1}{r^2 \sin \theta} \begin{vmatrix}
    e_r & e_\theta & r \sin \theta e_\phi \\
    \partial_r & \partial_\theta & \partial_\phi \\
    a_r & r a_\theta & r \sin \theta a_\phi
\end{vmatrix} = \frac{e_r}{r \sin \theta} \left( \frac{\partial (\sin \theta a_\phi)}{\partial \theta} - \frac{\partial a_\theta}{\partial \phi} \right) + \epsilon_\theta \left( \frac{\partial a_\phi}{\partial \phi} - \frac{\partial a_\theta}{\partial r} \right) + \frac{e_\phi}{r} \left( \frac{\partial (ra_\theta)}{\partial r} - \frac{\partial a_\phi}{\partial \theta} \right)
$$
15.6 An alternative derivation of the divergence in OCC

For Cartesian coordinates $\nabla \cdot \vec{a}$ (34) was constructed by taking the scalar product of the del operator (28) with the vector $\vec{a}$ (6). This construction does apply to cylindrical and spherical coordinates as well but is not as straightforward as for Cartesian coordinates because the basis vector are coordinate dependent. This implies that the derivatives of the del operator do act on the basis vectors. For example the basis vectors $\vec{e}_\rho$ and $\vec{e}_\phi$ are functions of the angle $\phi$ and therefore the derivative with respect to $\phi$ will act on them in a non-trivial way,

$$
\partial_\phi \vec{e}_\rho = \partial_\phi (\cos \phi \vec{e}_1 + \sin \phi \vec{e}_2) = -\sin \phi \vec{e}_1 + \cos \phi \vec{e}_2 = \vec{e}_\phi
$$
$$
\partial_\phi \vec{e}_\phi = \partial_\phi (-\sin \phi \vec{e}_1 + \cos \phi \vec{e}_2) = -\cos \phi \vec{e}_1 - \sin \phi \vec{e}_2 = -\vec{e}_\rho.
$$

(70)

Applying the definition of the divergence (34), taking (70) into account, we get

$$
\nabla \cdot \vec{a} = \left(\vec{e}_\rho \partial_\rho + \frac{1}{\rho} \vec{e}_\phi \partial_\phi + \vec{e}_z \partial_\phi \right) \cdot \left( a_\rho \vec{e}_\rho + a_\phi \vec{e}_\phi + a_z \vec{e}_z \right)
$$

$$
= \frac{1}{\rho} a_\rho + \frac{1}{\rho} \partial_\phi a_\phi + \frac{\partial}{\partial z} a_z = \frac{1}{\rho} \frac{\partial}{\partial \rho} (\rho a_\rho) + \frac{1}{\rho} \frac{\partial}{\partial \phi} a_\phi + \frac{\partial}{\partial z} a_z
$$

(71)

the expression which is quoted in table 1. The extra term $\frac{1}{\rho} a_\rho$ is generated by using $\partial_\phi \vec{e}_\rho = \vec{e}_\phi$. The moral of this section is that we could have equally well derived all the expressions in table 1 by considering the nabla as a vector and contracting them with vectors.

Rederiving results in a different manner is reassuring and usually leads to additional insight and can be helpful in solving problems. With this remark it is time to end the course and I hope that you are now well equipped with vector differential calculus for the rest of your studies in physics and mathematics.

[The End]