Maxwell’s equations for magnets

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Abstract
Magnetostatic fields in accelerators are conventionally described in terms of multipoles. We show that in two dimensions, multipole fields do provide solutions of Maxwell’s equations, and we consider the distributions of electric currents and geometries of ferromagnetic materials required (in idealized situations) to generate specified multipole fields. Then, we consider how to determine the multipole components in a given field. Finally, we show how the two-dimensional multipole description may be extended to three dimensions; this allows fringe fields, or the main fields in such devices as undulators and wigglers, to be expressed in terms of a set of modes, where each mode provides a solution to Maxwell’s equations.

1 Maxwell’s equations
Maxwell’s equations may be written in differential form as follows:

\[ \text{div } \vec{D} = \rho, \quad (1) \]
\[ \text{div } \vec{B} = 0, \quad (2) \]
\[ \text{curl } \vec{H} = \vec{J} + \frac{\partial \vec{D}}{\partial t}, \quad (3) \]
\[ \text{curl } \vec{E} = -\frac{\partial \vec{B}}{\partial t}. \quad (4) \]

The fields \( \vec{B} \) (magnetic flux density) and \( \vec{E} \) (electric field strength) determine the force on a particle of charge \( q \) travelling with velocity \( \vec{v} \) (the Lorentz force equation):

\[ \vec{F} = q \left( \vec{E} + \vec{v} \times \vec{B} \right). \]

The electric displacement \( \vec{D} \) and magnetic intensity \( \vec{H} \) are related to the electric field and magnetic flux density by:

\[ \vec{D} = \varepsilon \vec{E}, \]
\[ \vec{B} = \mu \vec{H}. \]

The electric permittivity \( \varepsilon \) and magnetic permeability \( \mu \) depend on the medium within which the fields exist. The values of these quantities in vacuum are fundamental physical constants. In SI units:

\[ \mu_0 = 4\pi \times 10^{-7} \text{ Hm}^{-1}, \]
\[ \varepsilon_0 = \frac{1}{\mu_0 c^2}, \]

where \( c \) is the speed of light in vacuum. The permittivity and permeability of a material characterize the response of that material to electric and magnetic fields. In simplified models, they are often regarded as constants for a given material; however, in reality the permittivity and permeability can have a complicated dependence on the fields that are present. Note that the relative permittivity \( \varepsilon_r \) and the relative permeability \( \mu_r \) are frequently used. These are dimensionless quantities, defined by:

\[ \varepsilon_r = \frac{\varepsilon}{\varepsilon_0}, \quad \mu_r = \frac{\mu}{\mu_0}. \quad (5) \]
That is, the relative permittivity is the permittivity of a material relative to the permittivity of free space, and similarly for the relative permeability.

The quantities \( \rho \) and \( \vec{J} \) are respectively the electric charge density (charge per unit volume) and electric current density (\( \vec{J} \cdot \vec{n} \) is the charge crossing unit area perpendicular to unit vector \( \vec{n} \) per unit time). Equations (2) and (4) are independent of \( \rho \) and \( \vec{J} \), and are generally referred to as the “homogeneous” equations; the other two equations, (1) and (3), are dependent on \( \rho \) and \( \vec{J} \), and are generally referred to as the “inhomogeneous” equations. The charge density and current density may be regarded as sources of electromagnetic fields. When the charge density and current density are specified (as functions of space, and, generally, time), one can integrate Maxwell’s equations (1)–(3) to find possible electric and magnetic fields in the system. Usually, however, the solution one finds by integration is not unique: for example, the field within an accelerator dipole magnet may be modified by propagating an electromagnetic wave through the magnet. However, by imposing certain constraints (for example, that the fields within a magnet are independent of time) it is possible to obtain a unique solution for the fields in a given system of electric charges and currents.

Most realistic situations are sufficiently complicated that solutions to Maxwell’s equations cannot be obtained analytically. A variety of computer codes exist to provide solutions numerically, once the charges, currents, and properties of the materials present are all specified, see for example References [1–3]. Solving for the fields in realistic (three-dimensional) systems often requires a reasonable amount of computing power; some sophisticated techniques have been developed for solving Maxwell’s equations numerically with good efficiency [4]. We do not consider such techniques here, but focus instead on the analytical solutions that may be obtained in idealized situations. Although the solutions in such cases may not be sufficiently accurate to complete the design of a real accelerator magnet, the analytical solutions do provide a useful basis for describing the fields in real magnets, and provide also some important connections with the beam dynamics in an accelerator.

An important feature of Maxwell’s equations is that, for systems containing materials with constant permittivity and permeability (i.e. permittivity and permeability that are independent of the fields present), the equations are linear in the fields and sources. That is, each term in the equations involves a field or a source to (at most) the first power, and products of fields or sources do not appear. As a consequence, the principle of superposition applies: if \( \vec{B}_1 \) and \( \vec{B}_2 \) are solutions of Maxwell’s equations with the current densities \( \vec{J}_1 \) and \( \vec{J}_2 \), then the field \( \vec{B}_T = \vec{B}_1 + \vec{B}_2 \) will be a solution of Maxwell’s equations, with the source given by the total current density \( \vec{J}_T = \vec{J}_1 + \vec{J}_2 \). This means that it is possible to represent complicated fields as superpositions of simpler fields. An important and widely used analysis technique for accelerator magnets is to decompose the field (determined from either a magnetic model, or from measurements of the field in an actual magnet) into a set of multipoles. While it is often the ideal to produce a field consisting of a single multipole component, this is never perfectly achieved in practice: the multipole decomposition indicates the extent to which components other than the “desired” multipole are present. Multipole decompositions also produce useful information for modelling the beam dynamics. Although the principle of superposition strictly only applies in systems where the permittivity and permeability are independent of the fields, it is always possible to perform a multipole decomposition of the fields in free space (e.g. in the interior of a vacuum chamber), since in that region the permittivity and permeability are constants. However, it should be remembered that for nonlinear materials (where the permeability, for example, depends on the magnetic field strength), the field inside the material comprising the magnet will not necessarily be that expected if one were simply to add together the fields corresponding to the multipole components.

Solutions to Maxwell’s equations lead to a rich diversity of phenomena, including the fields around charges and currents in certain simple configurations, and the generation, transmission and absorption of electromagnetic radiation. Many existing texts cover these phenomena in detail; see, for example, the authoritative text by Jackson [5]. Therefore, we consider only briefly the electric field around a point charge and the magnetic field around a long straight wire carrying a uniform current: our main purpose
here is to remind the reader of two important integral theorems (Gauss’ theorem, and Stokes’ theorem), of which we shall make use later. In the following sections, we will discuss analytical solutions to Maxwell’s equations for situations relevant to some of the types of magnets commonly used in accelerators. These include multipoles (dipoles, quadrupoles, sextupoles, and so on), solenoids, and insertion devices (undulators and wigglers). We consider only static fields. We begin with two-dimensional fields, that is fields that are independent of one coordinate (generally, the coordinate representing the direction of motion of the beam). We will show that multipole fields are indeed solutions of Maxwell’s equations, and we will derive the current distributions needed to generate “pure” multipole fields. We then discuss multipole decompositions, and compare techniques for determining the multipole components present in a given field from numerical field data (from a model, or from measurements). Finally, we consider how the two-dimensional multipole decomposition may be extended to three-dimensional fields, to include (for example) insertion devices, and fringe fields in multipole magnets.

2 Integral theorems and the physical interpretation of Maxwell’s equations

2.1 Gauss’ theorem and Coulomb’s law

Gauss’ theorem states that for any smooth vector field \( \vec{a} \):

\[
\int_V \text{div} \, \vec{a} \, dV = \int_{\partial V} \vec{a} \cdot d\vec{S},
\]

where \( V \) is a volume bounded by the closed surface \( \partial V \). Note that the area element \( d\vec{S} \) is oriented to point out of \( V \).

Gauss’ theorem is helpful for obtaining physical interpretations of two of Maxwell’s equations, (1) and (2). First, applying Gauss’ theorem to (1) gives:

\[
\int_V \text{div} \, \vec{D} \, dV = \int_{\partial V} \vec{E} \cdot d\vec{S} = q,
\]

where \( q = \int_V \rho \, dV \) is the total charge enclosed by \( \partial V \).

Suppose that we have a single isolated point charge in an homogeneous, isotropic medium with constant permittivity \( \varepsilon \). In this case, it is interesting to take \( \partial V \) to be a sphere of radius \( r \). By symmetry, the magnitude of the electric field must be the same at all points on \( \partial V \), and must be normal to the surface at each point. Then, we can perform the surface integral in (6):

\[
\int_{\partial V} \vec{D} \cdot d\vec{S} = 4\pi r^2 D.
\]

This is illustrated in Fig. 1: the outer circle represents a cross-section of a sphere (\( \partial V \)) enclosing volume \( V \), with the charge \( q \) at its centre. The black arrows in Fig. 1 represent the electric field lines, which are everywhere perpendicular to the surface \( \partial V \). Since \( \vec{D} = \varepsilon \vec{E} \), we find Coulomb’s law for the magnitude of the electric field around a point charge:

\[
E = \frac{q}{4\pi \varepsilon r^2}.
\]

Applied to Maxwell’s equation (2), Gauss’ theorem leads to:

\[
\int_V \text{div} \, \vec{B} \, dV = \int_{\partial V} \vec{B} \cdot d\vec{S} = 0.
\]

In other words, the magnetic flux integrated over any closed surface must equal zero – at least, until we discover magnetic monopoles. Lines of magnetic flux occur in closed loops; whereas lines of electric field can start (and end) on electric charges.
2.2 Stokes’ theorem and Ampère’s law

Stokes’ theorem states that for any smooth vector field $\vec{a}$:

$$\int_S \text{curl} \vec{a} \cdot d\vec{S} = \int_{\partial S} \vec{a} \cdot d\vec{l},$$

where the loop $\partial S$ bounds the surface $S$. Applied to Maxwell’s equation (3), Stokes’ theorem leads to:

$$\int_{\partial S} \vec{H} \cdot d\vec{l} = \int_S \vec{J} \cdot d\vec{S},$$

which is Ampère’s law. From Ampère’s law, we can derive an expression for the strength of the magnetic field around a long, straight wire carrying current $I$. The magnetic field must have rotational symmetry around the wire. There are two possibilities: a radial field, or a field consisting of closed concentric loops centred on the wire (or some superposition of these fields). A radial field would violate Maxwell’s equation (2). Therefore, the field must consist of closed concentric loops; and by considering a circular loop of radius $r$, we can perform the integral in Eq. (8):

$$2\pi r H = I,$$

where $I$ is the total current carried by the wire. In this case, the line integral is performed around a loop $\partial S$ centered on the wire, and in a plane perpendicular to the wire: essentially, this corresponds to one of the magnetic field lines, see Fig. 2. The total current passing through the surface $S$ bounded by the loop $\partial S$ is simply the total current $I$.

In an homogeneous, isotropic medium with constant permeability $\mu$, $\vec{B} = \mu_0 \vec{H}$, and we obtain the expression for the magnetic flux density at distance $r$ from the wire:

$$B = \frac{I}{2\pi \mu r}.$$ (9)

This result will be useful when we come to consider how to generate specified multipole fields from current distributions.

Finally, applying Stokes’ theorem to the homogeneous Maxwell’s equation (4), we find:

$$\int_{\partial S} \vec{E} \cdot d\vec{l} = -\frac{\partial}{\partial t} \int_S \vec{B} \cdot d\vec{S}.$$ (10)
Defining the electromotive force $\mathcal{E}$ as the integral of the electric field around a closed loop, and the magnetic flux $\Phi$ as the integral of the magnetic flux density over the surface bounded by the loop, Eq. (10) gives:

$$\mathcal{E} = -\frac{\partial \Phi}{\partial t},$$

which is Faraday’s law of electromagnetic induction. Faraday’s law is significant for magnets with time-dependent fields, such as pulsed magnets (used for injection and extraction), and magnets that are “ramped” (for example, when changing the beam energy in a storage ring). The change in magnetic field will induce a voltage across the coil of the magnet, that must be taken into account when designing the power supply. Also, the induced voltages can induce eddy currents in the core of the magnet, or in the coils themselves, leading to heating. This is an issue for superconducting magnets, which must be ramped slowly to avoid quenching [6].

2.3 Boundary conditions

Gauss’ theorem and Stokes’ theorem can be applied to Maxwell’s equations to derive constraints on the behaviour of electromagnetic fields at boundaries between different materials. Here, we shall focus on the boundary conditions on the magnetic field: these conditions will be useful when we consider multipole fields in iron-dominated magnets.

Consider first a short cylinder or “pill box” that crosses the boundary between two media, with the flat ends of the cylinder parallel to the boundary, see Fig. 3 (a). Applying Gauss’ theorem to Maxwell’s equation (2) gives:

$$\int_V \text{div} \mathbf{B} \, dV = \int_{\partial V} \mathbf{B} \cdot d\mathbf{S} = 0,$$

where the boundary $\partial V$ encloses the volume $V$ within the cylinder. If we take the limit where the length of the cylinder ($2h$) approaches zero, then the only contributions to the surface integral come from the flat ends; if these have infinitesimal area $dS$, then since the orientations of these surfaces are in opposite directions on opposite sides of the boundary, and parallel to the normal component of the magnetic field, we find:

$$-B_{1\perp} \, dS + B_{2\perp} \, dS = 0,$$

where $B_{1\perp}$ and $B_{2\perp}$ are the normal components of the magnetic flux density on either side of the boundary. Hence:

$$B_{1\perp} = B_{2\perp}.$$
In other words, the normal component of the magnetic flux density is continuous across a boundary.

A second boundary condition, this time on the component of the magnetic field parallel to a boundary, can be obtained by applying Stokes’ theorem to Maxwell’s equation (3). In particular, we consider a surface $S$ bounded by a loop $\partial S$ that crosses the boundary of the material, see Fig. 3 (b). If we integrate both sides of Eq. (3) over that surface, and apply Stokes’ theorem (7), we find:

$$\int_S \nabla \times \vec{H} \cdot d\vec{S} = \int_{\partial S} \vec{H} \cdot d\vec{l} = I + \frac{\partial}{\partial t} \int_S \vec{D} \cdot d\vec{S},$$

where $I$ is the total current flowing through the surface $S$. Now, let the surface $S$ take the form of a thin strip, with the short ends perpendicular to the boundary, and the long ends parallel to the boundary. In the limit that the length of the short ends goes to zero, the area of $S$ goes to zero: both the current flowing through the surface $S$, and the electric displacement integrated over $S$ become zero. However, there are still contributions to the integral of $\vec{H}$ around $\partial S$ from the long sides of the strip. Thus, we find that:

$$H_{1\parallel} = H_{2\parallel}, \quad (13)$$

where $H_{1\parallel}$ is the component of the magnetic intensity parallel to the boundary at a point on one side of the boundary, and $H_{2\parallel}$ is the component of the magnetic intensity parallel to the boundary at a nearby point on the other side of the boundary. In other words, the tangential component of the magnetic intensity $\vec{H}$ is continuous across a boundary.

We can derive a stronger constraint on the magnetic field at a boundary in the case that the material on one side of the boundary has infinite permeability (which can provide a reasonable model for some ferromagnetic materials). Since $\vec{B} = \mu \vec{H}$, it follows from (13) that:

$$\frac{B_{1\parallel}}{\mu_1} = \frac{B_{2\parallel}}{\mu_2},$$

and in the limit $\mu_2 \to \infty$, while $\mu_1$ remains finite, we must have:

$$B_{1\parallel} = 0. \quad (14)$$
In other words, the magnetic flux density at the surface of a material of infinite permeability must be perpendicular to that surface. Of course, the permeability of a material characterizes its response to an applied external magnetic field: in the case that the permeability is infinite, a material placed in an external magnetic field acquires a magnetization that exactly cancels any component of the external field at the surface of the material.

3 Two-dimensional multipole fields

Consider a region of space free of charges and currents; for example, the interior of an accelerator vacuum chamber (at least, in an ideal case, and when the beam is not present). If we further exclude propagating electromagnetic waves, then any magnetic field generated by steady currents outside the vacuum chamber must satisfy:

\[
\text{div} \, \vec{B} = 0, \tag{15}
\]
\[
\text{curl} \, \vec{B} = 0. \tag{16}
\]

Eq. (15) is just Maxwell’s equation (2), and Eq. (16) follows from Maxwell’s equation (3) given that \( \vec{J} = 0, \vec{B} = \mu_0 \vec{H} \), and derivatives with respect to time vanish.

We shall show that a magnetic field \( \vec{B} = (B_x, B_y, B_z) \) with \( B_z \) constant, and \( B_x, B_y \) given by:

\[
B_y + i B_x = C_n (x + iy)^{n-1} \tag{17}
\]

where \( i = \sqrt{-1} \) and \( C_n \) is a (complex) constant, satisfies Eqs. (15) and (16). Note that the field components \( B_x \) and \( B_y \) are real, and are obtained from the imaginary and real parts of the right hand side of Eq. (17). To show that the above field satisfies Eqs. (15) and (16), we apply the differential operator

\[
\frac{\partial}{\partial x} + i \frac{\partial}{\partial y} \tag{18}
\]

to each side of Eq. (17). Applied to the left hand side, we find:

\[
\left( \frac{\partial}{\partial x} + i \frac{\partial}{\partial y} \right) (B_y + i B_x) = \left( \frac{\partial B_y}{\partial x} - \frac{\partial B_x}{\partial y} \right) + i \left( \frac{\partial B_x}{\partial x} + \frac{\partial B_y}{\partial y} \right). \tag{19}
\]

Applied to the right hand side of Eq. (17), the differential operator (18) gives:

\[
\left( \frac{\partial}{\partial x} + i \frac{\partial}{\partial y} \right) C_n (x + iy)^{n-1} = C_n (n-1)(x + iy)^{n-2} + i^2 C_n (n-1)(x + iy)^{n-2} = 0. \tag{20}
\]

Combining Eqs. (17), (19) and (20), we find:

\[
\frac{\partial B_x}{\partial x} + \frac{\partial B_y}{\partial y} = 0,
\]
\[
\frac{\partial B_y}{\partial x} - \frac{\partial B_x}{\partial y} = 0.
\]

Finally, we note that \( B_z \) is constant, so any derivatives of \( B_z \) vanish; furthermore, \( B_x \) and \( B_y \) are independent of \( z \), so any derivatives of these coordinates with respect to \( z \) vanish. Thus, we conclude that for the field (17):

\[
\text{div} \, \vec{B} = 0, \tag{21}
\]
\[
\text{curl} \, \vec{B} = 0. \tag{22}
\]

and that this field is therefore a solution to Maxwell’s equations within the vacuum chamber. Of course, this analysis tells us only that the field is a possible physical field: it does not tell us how to generate such a field. The problem of generating a field of the form Eq. (17) we shall consider in Section 4.
Fields of the form \((17)\) are known as multipole fields. The index \(n\) (an integer) indicates the order of the multipole: \(n = 1\) is a dipole field, \(n = 2\) is a quadrupole field, \(n = 3\) is a sextupole field, and so on. A solenoid field has \(C_n = 0\) for all \(n\), and \(B_z\) non-zero; usually, a solenoid field is not considered a multipole field, and we assume (unless stated otherwise) that \(B_z = 0\) in a multipole magnet. Note that we can apply the principle of superposition to deduce that a more general magnetic field can be constructed by adding together a set of multipole fields:

\[
B_y + iB_x = \sum_{n=1}^{\infty} C_n(x + iy)^{n-1}.
\]

A “pure” multipole field of order \(n\) has \(C_n \neq 0\) for only that one value of \(n\).

The coefficients \(C_n\) in Eq. 23 characterise the strength and orientation of each multipole component in a two-dimensional magnetic field. It is sometimes more convenient to express the field using polar coordinates, rather than Cartesian coordinates. Writing \(x = r \cos \theta\) and \(y = r \sin \theta\), we see that Eq. (23) becomes:

\[
B_y + iB_x = \sum_{n=1}^{\infty} C_n r^{n-1} e^{i(n-1)\theta}.
\]

By writing the multipole expansion in this form, we see immediately that the strength of the field in a pure multipole of order \(n\) varies as \(r^{n-1}\) with distance from the magnetic axis. We can go a stage further, and express the field in terms of polar components:

\[
B_y + iB_x = B_r \sin \theta + B_\theta \cos \theta + iB_r \cos \theta - iB_\theta \sin \theta = (B_\theta + iB_r) e^{-i\theta},
\]

thus:

\[
B_\theta + iB_r = \sum_{n=1}^{\infty} C_n r^{n-1} e^{i(n-1)\theta}.
\]

By writing the field in this form, we see that for a pure multipole of order \(n\), rotation of the magnet through \(\pi/n\) around the \(z\) axis simply changes the sign of the field. We also see that if we write:

\[
C_n = |C_n| e^{i\phi_n}
\]

then the value of \(\phi_n\) (the phase of \(C_n\)) determines the orientation of the field. Conventionally, a pure multipole with \(\phi_n = 0\) is known as a “normal” multipole, while a pure multipole with \(\phi_n = \pi/2\) is known as a “skew” multipole.

The units of \(C_n\) depend on the order of the multipole. In SI units, for a dipole, the units of \(C_1\) are tesla (T); for a quadrupole, the units of \(C_2\) are Tm\(^{-1}\); for a sextupole, the units of \(C_3\) are Tm\(^{-2}\), and so on. It is sometimes preferred to specify multipole components in dimensionless units. In that case, we introduce a reference field, \(B_{\text{ref}}\), and a reference radius, \(R_{\text{ref}}\). The multipole expansion is then written:

\[
B_y + iB_x = B_{\text{ref}} \sum_{n=1}^{\infty} (a_n + ib_n) \left( \frac{x + iy}{R_{\text{ref}}} \right)^{n-1}.
\]

This is a standard notation for multipole fields, see for example [7]. In polar coordinates:

\[
B_y + iB_x = B_{\text{ref}} \sum_{n=1}^{\infty} (a_n + ib_n) \left( \frac{r}{R_{\text{ref}}} \right)^{n-1} e^{i(n-1)\theta}.
\]

The reference field and reference radius can be chosen arbitrarily, but must be specified if the coefficients \(a_n\) and \(b_n\) are to be interpreted fully.
Fig. 4: “Pure” multipole fields. Top: dipole. Middle: quadrupole. Bottom: sextupole. Fields on the left are normal ($a_n$ positive); those on the right are skew ($b_n$ positive). The positive $y$ axis is vertically up; the positive $x$ axis is horizontal and to the right.
Note that for a pure multipole field of order \(n\), the coefficients \(a_n\) and \(b_n\) are related to the derivates of the field components with respect to the \(x\) and \(y\) coordinates. Thus, for a normal multipole:

\[
\frac{\partial^{n-1}B_y}{\partial x^{n-1}} = (n - 1)! \frac{B_{\text{ref}}}{R_{\text{ref}}^{n-1}} a_n,
\]

and for a skew multipole:

\[
\frac{\partial^{n-1}B_x}{\partial x^{n-1}} = (n - 1)! \frac{B_{\text{ref}}}{R_{\text{ref}}^{n-1}} b_n.
\]

A normal dipole has a uniform vertical field; a normal quadrupole has a vertical field for \(y = 0\), that increases linearly with \(x\); a normal sextupole has a vertical field for \(y = 0\) that increases as the square of \(x\); and so on.

4 Generating multipole fields

Given a system of electric charges and currents, we can integrate Maxwell’s equations to find the electric and magnetic fields generated by those charges and currents. In general, the integration must be done numerically; but for simple systems it is possible to find analytical solutions. We considered two such cases in Section 2: the electric field around an isolated point charge, and the magnetic field around a long straight wire carrying a constant current.

It turns out that we can combine the magnetic fields from long, straight, parallel wires to generate pure multipole fields. It is also possible to generate pure multipole fields using high-permeability materials with the appropriate geometry. We consider both methods in this section. For the moment, we deal with “idealised” geometries without practical constraints. We discuss the impact of some of the practical limitations in later sections.

4.1 Current distribution for a multipole field

Our goal is to determine a current distribution that will generate a pure multipole field of specified order. As a first step, we derive the multipole components in the field around a long straight wire carrying a uniform current. We already know, from Ampère’s law (9) that the field at distance \(r\) from a long straight wire carrying current \(I\) in free space has magnitude given by:

\[
B = \frac{I}{2\pi \mu_0 r},
\]

and that the direction of the field describes a circle centred on the wire. To derive the multipole components in the field, we first derive an expression for the field components at an arbitrary point \((x, y)\) from a wire carrying current \(I\), passing through a point \((x_0, y_0)\) and parallel to the \(z\) axis.

Since we are working in two dimensions, we can represent the components of a vector by the real and imaginary parts of a complex number. Thus, the vector from \((x_0, y_0)\) to a point \((x, y)\) is given by \(re^{i\theta} - r_0e^{i\theta_0}\), and the magnitude of the field at \((x, y)\) is:

\[
B = \frac{I}{2\pi \mu_0} \frac{1}{|re^{i\theta} - r_0e^{i\theta_0}|}.
\]

The geometry is shown in Fig. 5. The direction of the field is perpendicular to the line from \((x_0, y_0)\) to \((x, y)\). Since a rotation through 90° can be represented by a multiplication by \(i\), we can write:

\[
B_x + iB_y = \frac{I}{2\pi \mu_0} \frac{i (re^{i\theta} - r_0e^{i\theta_0})}{|re^{i\theta} - r_0e^{i\theta_0}|^2},
\]
Fig. 5: Geometry for calculation of multipole components in the field around a long, straight wire carrying a uniform current. The wire passes through $r_0 e^{i\theta_0}$, and is parallel to the z axis (the direction of the current is pointing out of the page).

and hence:

$$B_y + iB_x = \frac{I}{2\pi\mu_0} \frac{(r e^{-i\theta} - r_0 e^{-i\theta_0})}{|r e^{i\theta} - r_0 e^{i\theta_0}|^2} = \frac{I}{2\pi\mu_0} \frac{1}{r e^{i\theta} - r_0 e^{i\theta_0}}.$$ 

Now, we write the magnetic field as:

$$B_y + iB_x = -\frac{I}{2\pi\mu_0 r_0} \frac{e^{-i\theta_0}}{1 - \frac{r}{r_0} e^{i(\theta - \theta_0)}}.$$ 

and use the Taylor series expansion for $(1 - \zeta)^{-1}$, where $\zeta$ is a complex number with $|\zeta| < 1$:

$$\frac{1}{1 - \zeta} = \sum_{n=0}^{\infty} \zeta^n,$$

to write:

$$B_y + iB_x = -\frac{I}{2\pi\mu_0 r_0} e^{-i\theta_0} \sum_{n=1}^{\infty} \left( \frac{r}{r_0} \right)^{n-1} e^{i(n-1)(\theta - \theta_0)}.$$  \hspace{1cm} (27)

Eq. (27) is valid for $r < r_0$. Comparing with the standard multipole expansion, Eq. (26), we see that if we choose for the reference field $B_{ref}$ and the reference radius $R_{ref}$:

$$B_{ref} = \frac{I}{2\pi\mu_0 r_0},$$

$$R_{ref} = r_0,$$

then the coefficients for the multipole components in the field are given by:

$$b_n + i a_n = -e^{-i\theta_0}.$$ 

The field around a long straight wire can be represented as an infinite sum over all multipoles.
Now we consider a current flowing on the surface of a cylinder of radius \( r_0 \). Suppose that the current flowing in a section of the cylinder at angle \( \theta_0 \) and subtending angle \( d\theta_0 \) at the origin is \( I(\theta_0) \, d\theta_0 \). By the principle of superposition, we can obtain the total field by summing the contributions from the currents at all values of \( \theta_0 \):

\[
B_y + iB_x = -\frac{1}{2\pi\mu_0 r_0} \sum_{n=1}^{\infty} \left( \frac{r}{r_0} \right)^{n-1} e^{i(n-1)\theta} \int_0^{2\pi} e^{-i\theta_0} I(\theta_0) \, d\theta_0.
\]  

We see that the multipole components are related to the Fourier components in the current distribution over the cylinder of radius \( r_0 \). In particular, if we consider a current distribution with just a single Fourier component:

\[
I(\theta_0) = I_0 \cos(n_0 \theta_0 - \phi),
\]  

the integral in the right hand side of Eq. (28) vanishes except for \( n = n_0 \), and we find:

\[
B_y + iB_x = -\frac{I_0}{2\pi\mu_0 r_0} \left( \frac{r}{r_0} \right)^{n_0-1} e^{i(n_0-1)\theta} \pi e^{-i\phi}.
\]

The current distribution (29) generates a pure multipole field of order \( n_0 \). If we choose, as before:

\[
B_{\text{ref}} = \frac{I_0}{2\pi\mu_0 r_0}, \quad R_{\text{ref}} = r_0,
\]

then the multipole coefficients are:

\[
b_n + ia_n = -\pi e^{-i\phi}.
\]

The parameter \( \phi \) gives the “angle” of the current distribution. For \( \phi = 0 \) or \( \phi = \pi \), the current generates a normal multipole; for \( \phi = \pm\pi/2 \), the current generates a skew multipole.

The fact that a sinusoidal current distribution on a cylinder can generate a pure multipole field is not simply of academic interest. By winding wires in an appropriate pattern on a cylinder, it is possible to approximate a sinusoidal current distribution closely enough to produce a multipole field of acceptable quality for many applications. Usually, several layers of windings are used with a different pattern of wires in each layer, to improve the approximation to a sinusoidal current distribution. Superconducting wires can be used to achieve strong fields: an example of superconducting quadrupoles in the LHC is shown in Fig. 7.

4.2 Geometry of an iron-dominated multipole magnet

Normal-conducting magnets usually use iron cores to increase the flux density achieved by a given current. In such a magnet, the shape of the magnetic field depends mainly on the geometry of the iron. In this section, we shall derive the geometry required to generate a pure multipole of given order. To simplify the problem, we will make some approximations: in particular, we shall assume that the iron core has uniform cross-section and infinite extent along \( z \); that there are no limits to the iron in \( x \) or \( y \); and that the iron has infinite permeability. The field in a more realistic magnet will generally need to be calculated numerically; however, the characteristics derived from our idealized model are often a good starting point for the design of an iron-dominated multipole magnet.

We base our analysis on the magnetic scalar potential, \( \varphi \), which is related to the magnetic field \( \vec{B} \) by:

\[
\vec{B} = -\nabla \varphi.
\]

Note that the curl of the field in this case is zero, for any function \( \varphi \): this is a consequence of the mathematical properties of the grad and curl operators. Therefore, it follows from Maxwell’s equation (3)
Fig. 6: Current distributions for generating pure multipole fields. Top: dipole. Middle: quadrupole. Bottom: sextupole. Fields on the left are normal ($a_n$, positive); those on the right are skew ($b_n$, positive). The positive $y$ axis is vertically up; the positive $x$ axis is horizontal and to the right. The deviation of the red line from the circular boundary shows the local current density. Current is flowing in the positive $z$ direction (out of the page) for increased radius, and in the negative $z$ direction for reduced radius.
that a magnetic field can only be derived from a scalar potential if: (i) there is no current density at the location where the field is to be calculated; (ii) there is no time-dependent electric displacement at the location where the field is to be calculated. Where there exists an electric current or a time-dependent electric field, it is more appropriate to use a vector potential (in which case, the magnetic flux density is found from the curl of the vector potential). However, for multipole fields, we have already shown that both the divergence and the curl of the field vanish, Eqs. (21) and (22). Since the curl of the grad of any function is identically zero, Eq. (22) is automatically satisfied for any field \( \vec{B} \) derived using (30). From Eq. (21), we find:

\[
\nabla^2 \phi = 0, \quad (31)
\]

where \( \nabla^2 \) is the laplacian operator. Eq. (31) is Poisson’s equation: the scalar potential in a particular case is found by solving this equation with given boundary conditions.

To determine the geometry of iron required to generate a pure multipole field, we shall start by writing down the scalar potential for a pure multipole field. Since the magnetic flux density \( \vec{B} \) is obtained from the gradient of the scalar potential, the flux density at any point must be perpendicular to a surface of constant scalar potential. However, we already know, from Eq. (14), that the magnetic flux density at the surface of a material with infinite permeability must be perpendicular to that surface. Hence, to generate a pure multipole field in a magnet containing material of infinite permeability, we just need to shape the material so that its surface follows a surface of constant magnetic scalar potential for the required field.

We therefore look for a potential \( \phi \) that satisfies:

\[
- \left( \frac{\partial}{\partial y} + i \frac{\partial}{\partial x} \right) \phi = B_y + i B_x = C_n (x + iy)^{n-1}.
\]

As we shall now show, an appropriate solution is:

\[
\phi = - |C_n| \frac{r^n}{n} \sin(n\theta - \phi_n) \quad (32)
\]

where:

\[
x + iy = re^{i\theta},
\]
and so:
\[
\begin{align*}
  x &= r \cos \theta, \\
  y &= r \sin \theta.
\end{align*}
\]

That Eq. (32) is indeed the potential for a pure multipole of order \( n \) can be shown as follows. In polar coordinates, the gradient can be written:
\[
\nabla \varphi = \hat{r} \frac{\partial \varphi}{\partial r} + \hat{\theta} \frac{\partial \varphi}{\partial \theta},
\]
(33)

where \( \hat{r} \) and \( \hat{\theta} \) are unit vectors in the directions of increasing \( r \) and \( \theta \), respectively. Using:
\[
\begin{align*}
  \hat{r} &= \hat{x} \cos \theta + \hat{y} \sin \theta, \\
  \hat{\theta} &= -\hat{x} \sin \theta + \hat{y} \cos \theta,
\end{align*}
\]

it follows from Eq. (33) that:
\[
-\nabla \varphi = (\hat{x} \cos \theta + \hat{y} \sin \theta) |C_n| r^{n-1} \sin(n\theta - \phi_n) - (\hat{x} \sin \theta - \hat{y} \cos \theta) |C_n| r^{n-1} \cos(n\theta - \phi_n),
\]
\[
= \hat{x} \sin ((n-1)\theta - \phi_n) |C_n| r^{n-1} + \hat{y} \cos ((n-1)\theta - \phi_n) |C_n| r^{n-1}.
\]

Thus, the field derived from the potential (32) can be written:
\[
B_y + iB_x = |C_n| e^{-i\phi_n} r^{n-1} e^{i(n-1)\theta}.
\]

Therefore, if:
\[
C_n = |C_n| e^{-i\phi_n},
\]

then:
\[
B_y + iB_x = C_n r^{n-1} e^{i(n-1)\theta} = C_n (x + iy)^{n-1},
\]

and we see that the potential (32) does indeed generate a pure multipole field of order \( n \).

From the above argument, we can immediately conclude that to generate a pure multipole field, we can shape a high permeability material such that the surface of the material follows the curve given (in parametric form, with parameter \( \theta \)) by:
\[
r^n \sin(n\theta - \phi_n) = r_0^n,
\]
(34)

where \( r_0 \) is a constant giving the minimum distance between the surface of the material and the origin. The cross-sections of iron-dominated multipole magnets of orders 1, 2 and 3 are shown in Fig. 8. Note that \( r \to \infty \) for \( n\theta - \phi_n \to \text{integer} \times \pi \). Treating each region between infinite values of \( r \) as a separate pole, we see that a pure multipole of order \( n \) has \( 2n \) poles. We also see that the potential changes sign when moving from one pole to either adjacent pole: that is, poles alternate between “north” and “south”. The field must be generated by currents flowing along wires between the poles, parallel to the \( z \) axis: to avoid direct contribution from the field around the wires, these wires should be located a large (in fact, infinite) distance from the origin.

Note that it is possible to determine the shape of the pole face for a magnet containing any specified set of multipoles, by summing the potentials for the different multipole components, and then solving for \( r \) as a function of \( \theta \), for a fixed value of the scalar potential. Magnets designed to have more than one multipole component are often known as “combined function” magnets. Perhaps the most common type of combined function magnet is a dipole with a quadrupole component: such magnets can be used to steer and focus a beam simultaneously. The shape of the pole faces and the field lines in a dipole with (strong) quadrupole component is shown in Fig. 9.
Fig. 8: Pole shapes for generating pure multipole fields. Top: dipole. Middle: quadrupole. Bottom: sextupole. Fields on the left are normal ($a_n$ positive); those on the right are skew ($b_n$ positive). The positive $y$ axis is vertically up; the positive $x$ axis is horizontal and to the right. The poles, shown as black (north) or grey (south), are constructed from material with infinite permeability.
In practice, some variation from the “ideal” geometry is needed, to account for the fact that the material used in the magnet has finite permeability, and finite extent transversely and longitudinally. The wires carrying the current that generates the magnetic flux are arranged in coils around each pole; as we shall see, the strength of the field is determined by the number of ampere-turns in each coil. An iron-dominated electromagnetic quadrupole is shown in Fig. 10.

To complete our discussion of methods to generate multipole fields, we derive an expression for the field strength in an iron dominated magnet with a given number of ampere-turns in the coil around each pole. To do this, we consider a line integral as shown in Fig. 11. In the figure, we show a quadrupole; however the generalisation to other orders of multipole is straightforward. Note that, in principle, the coils carrying the electric current, and the line segment $C_3$, are an infinite distance from the origin (the centre of the magnet).

Using Maxwell’s equation (3), with constant (zero) electric displacement, and integrating over the surface $S$ bounded by the curve $C_1 + C_2 + C_3$ gives:

$$\int_S \text{curl} \, \vec{H} \cdot d\vec{S} = \int_S \vec{J} \cdot d\vec{S} = -NI.$$

Note that the surface is oriented so that the normal is parallel to the positive $z$ axis; and the coil around
each pole consists of $N$ turns of wire carrying current $I$. Applying Stokes’ theorem (7) gives:

$$\int_{C_1} \vec{H} \cdot d\vec{l} + \int_{C_2} \vec{H} \cdot d\vec{l} + \int_{C_3} \vec{H} \cdot d\vec{l} = -NI.$$  

We know, from Eq. (12), that the normal component of the magnetic flux density $\vec{B}$ is continuous across a boundary. Then, since $\vec{B} = \mu \vec{H}$, it follows that for a finite field between the poles, and for $\mu \to \infty$, the magnetic intensity $\vec{H}$ vanishes within the poles. Also, the field is perpendicular to the line segment $C_2$. Thus, the only part of the integral that makes a non-zero contribution, is the integral along $C_1$ from the face of the pole to the origin. Hence:

$$\int_0^{r_0} \frac{B_r(r)}{\rho_0} dr = NI. \quad (35)$$

The contour $C_1$ is chosen so that along this contour, the field has only a radial component, parallel to the contour. From Eq. (24), we see that for a multipole of order $n$, along this contour we have:

$$B_r = |C_n| r^{n-1}.$$  

Thus, we find by performing the integral in Eq. (35):

$$|C_n| = \mu_0 NI \frac{n}{r_0^n}.$$  

For a normal multipole, the field is given by:

$$B_y + iB_x = \frac{\mu_0 n NI}{r_0} \left( \frac{x + iy}{r_0} \right)^{n-1}.\quad (36)$$

For example, in a normal quadrupole ($n = 2$), the field gradient is given by:

$$\frac{\partial B_y}{\partial x} = \frac{2\mu_0 NI}{r_0^2}. \quad (36)$$
5 Multipole decomposition

In the previous section, we derived the current density distributions and material geometries needed to generate a pure multipole field of a given order. However, the distributions and geometries required are not perfectly achievable in practice: the currents and materials have infinite longitudinal extent; and we require either a current that exists purely on the surface of a cylinder, or infinite permeability materials with infinite transverse extent.

Real multipole magnets, therefore, will not consist of a single multipole component, but a superposition of (in general) an infinite number of multipole fields. The exact shape of the field can have a significant impact on the beam dynamics in an accelerator. In many simulation codes for accelerator beam dynamics, the magnets are specified by the multipole coefficients: this is because simple techniques exist for approximating the effect, for example, of sextupole, octupole, and other higher-order components in the field of a quadrupole magnet. The question then arises how to determine the multipole components in a given magnetic field.

At this point, we can make a distinction between the design field of a magnet, and the field that exists within a fabricated magnet. The design field is one that is still in some sense “ideal”; though the design field for a quadrupole magnet (for example) will contain other multipole components, because the design has to respect practical constraints, i.e. the magnet will have finite longitudinal and transverse extent, any currents will flow in wires of non-zero dimension, and any materials present will have finite (and often non-linear) permeability. Usually, one attempts to optimize the design to minimize the strengths of the multipole components apart from the one required: the residual strengths are generally known as systematic multipole errors. These errors will be present in any fabricated magnet, although, because of construction tolerances, the errors will vary between any two magnets of the same nominal design. The differences between the multipole components in the design field and the components in a particular magnet are known as random multipole errors.

The effects of systematic and random multipole errors on an accelerator, and hence the specification of upper limits on these quantities, can usually only be properly understood by running beam dynamics simulations. Therefore, accelerator magnet (and lattice) design often proceeds iteratively. Some initial estimate of the limits on the errors is often needed to guide the magnet design; but then any design that is developed must be studied by further beam dynamics simulations to determine whether improvements are needed.

It is therefore important to be able to determine the multipole components in a magnetic field from numerical field data: these data may come from either a magnetic model (i.e. from the design of a magnet), or from measurements on a real device. There are different procedures that can be used to achieve the “decomposition” of a field into its multipole components. In this section, we shall consider methods based on Cartesian and polar representations of two-dimensional fields (i.e. fields that are independent of the longitudinal coordinate). In Section 6 we shall consider decompositions of three-dimensional fields (i.e. fields that have explicit dependence on longitudinal as well as transverse coordinates). However, we first consider an important concept in the discussion of multipole field errors, namely how the symmetry of a multipole magnet leads to “allowed” and “forbidden” higher-order multipoles.

5.1 Multipole symmetry, “allowed” and “forbidden” higher-order multipoles

A pure multipole field of order \( n \) can be written:

\[
B_y + iB_x = |C_n| e^{-i\phi_n} r^{n-1} e^{i(n-1)\theta}.
\]  

The parameter \( \phi_n \) characterises the angular orientation of the magnet around the \( z \) axis. In particular, from Eq. (34), we see that a change in \( \phi_n \) by \( n\alpha \) is equivalent to a rotation of the coordinates (a change in \( \theta \)) by \(-\alpha\). Thus, a rotation of a magnet around the \( z \) axis by angle \( \alpha \) may be represented by a change in \( \phi_n \) by \( n\alpha \). In particular, if the magnet is rotated by \( \pi/n \), then from Eq. (37), we see that the field at
any point simply changes sign:

\[
\text{if } \phi_n \mapsto \phi_n + \pi, \text{ then } \vec{B} \mapsto -\vec{B}.
\]  

(38)

This property of the magnetic field is imposed by the symmetry of the magnet. In a real magnet, it will not be satisfied exactly, because random variations in the geometry will break the symmetry. However, it is possible to maintain the symmetry exactly in the design of the magnet; this means that although higher order multipoles will in general be present, only those multipoles satisfying the symmetry constraint (38) can be present. These are the “allowed” multipoles. Other multipoles, which must be completely absent, are the “forbidden” multipoles.

We can derive a simple expression for the allowed multipoles in a magnet designed with symmetry for a multipole of order \( n \). Consider an additional multipole (a “systematic error”) in this field, of order \( m \). By the principle of superposition, the total field can be written as:

\[
B_y + iB_x = |C_n|e^{-i\phi_n}r^{n-1}e^{i(n-1)\theta} + |C_m|e^{-i\phi_m}r^{m-1}e^{i(m-1)\theta}.
\]

The geometry is such that under a rotation about the \( z \) axis through \( \pi/n \), the magnet looks the same, except that all currents have reversed direction: therefore the field simply changes sign. Under this rotation \( \phi_n \mapsto \phi_n + \pi \); however, \( \phi_m \mapsto \phi_m + m\pi/n \). This means that we must have:

\[
e^{-i\frac{m\pi}{n}} = -1.
\]

Therefore, \( m/n \) must be an odd integer. Assuming that \( m \neq n \) (i.e. the multipole error is of a different order than the “main” multipole field), then:

\[
\frac{m}{n} = 3, 5, 7, \ldots
\]  

(39)

Thus, for a dipole, the allowed higher order multipoles are sextupole, decapole, etc.; for a quadrupole, the allowed higher order multipoles are dodecapole, 20-pole, etc. The fact that the allowed higher order multipoles have an order given by an odd integer multiplied by the order of the main multipole is a consequence of the fact that magnetic poles always occur in north-south pairs. This is illustrated for a quadrupole in Fig. 12; here, we see that to maintain the correct rotational symmetry (with the field changing sign under a rotation through \( \pi/2 \)) the first higher-order multipole must be constructed by “splitting” each main pole into three, then into five, and so on.

The field in a real magnet will contain all higher order multipoles, not just the ones allowed by symmetry. However, it is often the case that the allowed multipoles dominate over the forbidden multipoles.

### 5.2 Fitting multipoles: Cartesian basis

Suppose we have obtained a set of numerical field data, either from a magnetic model, or from measurements on a real magnet. To determine the effect of the field on the beam dynamics in an accelerator, it is helpful to know the multipole components in the field. One way to compute the multipole components is to fit a polynomial to the field data. For example, if we consider a normal multipole (coefficients \( C_n \) are all real), the vertical field along the \( x \) axis (i.e. for \( y = 0 \)) is given by:

\[
B_y = \sum_{n=1}^{\infty} C_n x^{n-1}.
\]  

(40)

The number of data points determines the highest order multipole that can be fitted. Fitting may be achieved using, for example, a routine that minimises the squares of the residuals between the data and the fitted function. However, although this procedure can, in principle, produce good results, it is not
very robust. In particular, the presence of multipoles of higher order than those included in the fit can affect the values determined for those multipoles that are included in the fit. We can illustrate this as follows.

Let us construct a quadrupole field \((n = 2)\), and add to it higher order multipoles of order 3, 4, 5 and 6. The values of the coefficients \(a_n\) (actual values, and fitted values in two different cases) are given in Table 1. The field \(B_y/B_{ref}\) is plotted as a function of \(x/R_{ref}\) in Fig. 13: the field data are shown as points, while the fit, including multipoles up to order 6, is shown as a line. Also shown is the deviation \(\Delta B_y/B_{ref}\) from an ideal quadrupole field, i.e. \(\Delta B_y\) is the contribution of the higher order multipoles. We see that if we base the fit on all the multipoles that are present (i.e. up to order 6), then we obtain accurate values for all multipole coefficients.

However, in general, multipoles of all orders are present, while our fit is based on a finite number of multipoles. If we try to fit the data in our illustrative case using multipoles up to order 5 only (i.e. omitting the order 6 multipole that is present), then we see that there is an impact on the accuracy with which we determine the lower-order multipoles. This can be seen in the final column of Table 1: there is even an error in the value that we determine for the quadrupole strength. When we plot the fit against the field data, we see that there is some small residual deviation between the data and the fit: this is to be expected, since the function we are using to obtain the fit does not match exactly the function used to generate the data. Although not visible in the total field, plotted in Fig. 14, the difference between the fit and the data is apparent in the plot of the deviation from the quadrupole field.

Our concern is that the presence of higher-order multipoles has affected the accuracy with which

\[ \begin{array}{|c|c|c|c|} 
\hline 
\text{n} & \text{actual coefficient } a_n & \text{fitted coefficient } a_n \\
\hline 
\text{(n ≤ 6)} & \text{(n ≤ 5)} \\
\hline 
2 & 1.000 & 1.000 & 0.9972 \\
3 & 0.010 & 0.010 & 0.0100 \\
4 & 0.001 & 0.001 & 0.0131 \\
5 & 0.010 & 0.010 & 0.0100 \\
6 & 0.010 & 0.010 & — \\
\hline 
\end{array} \]

Fig. 12: Normal quadrupole field (left) and dodecapole field (right). The dodecapole is the first higher-order multipole with the same rotational symmetry as the quadrupole (under a rotation by \(\pi/2\), north and south poles interchange).
we determine the lower-order multipoles, even down to the quadrupole field strength. This can have significant implications for beam dynamics: the effect of a linear focusing error in a beam line (from some variation in the quadrupole strength) can be very different from the effects of higher-order multipole errors. For example, if one is measuring the betatron tunes or the beta functions in a storage ring, these can be very sensitive to linear focusing errors, and relatively insensitive to higher order multipoles. Determining the multipole coefficients using a polynomial fit can lead to inaccurate predictions of the linear behaviour of the beam line, depending on the higher-order multipoles present in the magnets.

The problem is that we have based our fit on monomials, i.e. powers of $x$. Our fit is a sum of these monomials, with coefficients determined from the data. However, it is possible to obtain a fit to data generated using one monomial, with a different monomial. For example, if one constructs data which is purely linear in $x$, then one can obtain a fit using a monomial $x^3$ (even though the fit will not be as good as one obtained using a monomial $x$). Mathematically, the basis functions we are using (monomials in this case) are not orthogonal: the coefficients we determine depend on the which set of basis functions we choose to use. A more robust technique would use basis functions that are orthogonal, i.e. the coefficients we determine will be the same, no matter which set of functions we choose. Fortunately, there exists an appropriate set of functions that provides an orthogonal basis for multipole fields. We discuss this basis in the following section, 5.3. The advantage of orthogonal basis functions is that the coefficients we determine for different terms in the fit are independent of which terms we include in the fit; for example, the quadrupole strength that we find in a particular magnet will be the same, irrespective of which higher-order terms we include in the fit, and which higher-order terms are actually present.
5.3 Fitting multipoles: polar basis

From Eq. (24) we know that the field in a multipole magnet can be written in polar coordinates as:

\[ B_\theta + i B_r = \sum_{n=1}^{\infty} C_n r^{n-1} e^{in\theta}. \]

We see that if we make a set of measurements of \( B_r \) and \( B_\theta \) at different values of \( \theta \) and fixed radial distance \( r \), then we can obtain the coefficients \( C_n \) by a discrete Fourier transform.

Suppose we make \( M \) measurements of the field, at \( \theta = \theta_m \), where:

\[ \theta_m = \frac{2\pi m}{M}, \quad m = 0, 1, 2 \ldots M - 1. \]  \( \text{(41)} \)

We write the measurement at \( \theta = \theta_m \) as \( B_m \); note that \( B_m \) is a complex number, whose real and imaginary parts are given by the azimuthal and radial components of the field at \( \theta = \theta_m \).

Now we construct, for a chosen integer \( n' \):

\[ \sum_{m=0}^{M-1} B_m e^{-2\pi in' \frac{m}{M}} = \sum_{m=0}^{M-1} \sum_{n=1}^{\infty} C_n r_0^{n-1} e^{2\pi i(n-n') \frac{m}{M}}, \]

where \( r_0 \) is the radial distance at which the field measurements are made. The summation over \( m \) on the right hand side vanishes, unless \( n = n' \). Thus, we can write:

\[ \sum_{m=0}^{M-1} B_m e^{-2\pi in' \frac{m}{M}} = MC_n r_0^{n'-1}. \]

If we relabel \( n' \) as \( n \), then we see that the multipole coefficients \( C_n \) are given by:

\[ C_n = \frac{1}{Mr_0^{n-1}} \sum_{m=0}^{M-1} B_m e^{-2\pi in \frac{m}{M}}. \]  \( \text{(42)} \)

The advantage of this technique over that in section 5.2 is that the basis functions used to construct the fit are of the form \( e^{in\theta} \), for integer \( n \). These functions are orthogonal: mathematically, this means that:

\[ \int_0^{2\pi} e^{in\theta} e^{-in'\theta} \, d\theta = 2\pi \delta_{nn'}, \]

where the Kronecker delta function \( \delta_{nn'} = 1 \) if \( n = n' \), and \( \delta_{nn'} = 0 \) if \( n \neq n' \). The important consequence for us, is that the value we determine for any given multipole using Eq. (42) is independent of the presence of any other multipoles, of higher or lower order.

A further advantage of using the polar basis instead of the Cartesian basis, comes from the dependence of the field on the radial distance. Suppose that the field data are measured (or obtained from a model) with accuracy \( \Delta B_m \). Then the accuracy in the multipole coefficients will be:

\[ \Delta C_n \approx \frac{\Delta B_m}{r_0^{n-1}}. \]

We obtain better accuracy in the multipole coefficients if we choose the radius \( r_0 \) on which the measurements are made, to be as large as possible. Furthermore, the accuracy in the fitted field will be:

\[ \Delta B \approx \Delta C_n \left( \frac{r}{r_0} \right)^{n-1}. \]

We obtain improved accuracy in the field for \( r < r_0 \); but the accuracy reduces quickly (particularly for higher-order multipoles) for \( r > r_0 \). It is important to choose the radial distance \( r_0 \) large enough to enclose all particles likely to pass through the magnet, otherwise results from tracking may not be accurate.
5.4 Multipole decomposition: some comments

In this section, we have considered two techniques for deriving the multipole components of twodimensional magnetostatic fields. We have seen that while the multipole components can be obtained, in principle, from a simple least-squares fit of a polynomial to the field components along one or other of the coordinate axes, there are advantages to basing the fit on field data obtained on a circle enclosing the origin, with as large a radius as possible. In the next section, we shall see how the idea of a multipole expansion can be generalised to three dimensions, and how a multipole decomposition can be performed in that case. However, it is worth pausing to consider in a little more detail some of the reasons for wishing to represent a field as a set of (multipole) modes.

It is of course possible to represent a magnetic field using a set of numerical field data, giving the three field components on points forming a “mesh” covering the region of interest. In some ways, this is a very convenient representation, since it is the one usually provided directly by a magnetic modelling code: further processing is usually required to arrive at other representations. However, while a numerical field map in two dimensions is often a practical representation, in three dimensions the amount of data in even a relatively simple magnet can become extremely large, especially if a high resolution is required for the mesh. A multipole representation, on the other hand, provides the description of a magnetic field as a relatively small set of coefficients, from which the field components at any point can be reconstructed, using the basis functions. In other words, a multipole representation is more “portable” than a numerical field map.

Secondly, a representation based on a multipole expansion lends itself to further manipulation in ways that a numerical field map does not. For example, any noise in the data (from measurement or computational errors) can be “smoothed” by suppression of higher-order modes. Conversely, random errors can be introduced into data based on a model with perfect symmetry by introducing multipole coefficients corresponding to “forbidden” harmonics. There will of course be issues surrounding the suppression or enhancement of errors by adjusting the multipole coefficients; however, one benefit of this approach is that for any set of multipole coefficients, the field is at least a physical field, in the sense of satisfying Maxwell’s equations. The same will, for example, not usually be true if a general smoothing algorithm is applied to a numerical field map.

Finally, one of the main motivations for performing a multipole decomposition of a field is to provide data in a format appropriate for many beam dynamics codes. Accuracy is one criterion often important for beam dynamics codes: efficiency is another. Characterisation of a storage ring frequently requires tracking of thousands of particles over hundreds or thousands of turns, through a beam line that can easily consist of hundreds of magnetic elements. Numerical integration of the equations of motion for a particle in a numerical field map is generally too slow to be a practical method. There are many techniques that can be used to improve the efficiency of particle tracking in accelerators: one of the most common is the “thin lens” method. The dynamical effects of dipoles and quadrupoles usually need to be represented with high accuracy. Fortunately, for these magnets, it is possible to write down accurate solutions to the equations of motion in closed form, allowing tracking through a magnet of given length to be performed in a single step. The same is not true for sextupoles, or higher-order multipoles; however, it is usually sufficiently accurate to represent such magnets by a model in which the length of the magnet approaches zero, but where the integrated strength (the multipole coefficient multiplied by the length) remains constant. For such a “thin lens” it is possible to write down exact solutions to the equations of motion, allowing tracking again to be performed in a single step. A quadrupole with higher-order multipole errors can be represented as a “long” perfect quadrupole field, with a set of “thin” multipoles at one end, or at the centre. However, construction of such a model for a tracking code requires a multipole decomposition of the field obtained from a magnet modelling code.

We should emphasise that in our discussion of multipole decomposition, here and in sections 5.2 and 5.3, we have made no clear distinction between field data obtained from a computational model, or from measurement of a real magnet. Of course, it is much easier to obtain the data required from a
computational model: it is then quite straightforward to perform the required decomposition to determine the values of the various multipole coefficients. Unfortunately, the data do not include manufacturing errors, which can be very important. Measurements provide more realistic data: however, many other issues need to be addressed, including accuracy of field measurements, alignment of the measurement instruments with respect to the magnet, etc. Such issues are beyond the scope of our discussion.

6 Three-dimensional fields

In the previous sections, we have restricted ourselves to the case where the magnetic field is independent of the longitudinal coordinate. The multipole modes that we can use for such fields actually provide a good description for many accelerator multipole magnets, even though such magnets of course have finite length. The ends or “fringe fields” of dipoles, quadrupoles and so on, where the field strengths often vary rapidly with longitudinal position, cannot be accurately represented by two-dimensional fields; however, in many accelerators, only the fringe fields of dipoles have a significant impact on the dynamics.

However, there are cases where a full three-dimensional description of a magnetic field is desirable, or even necessary. For example, the fields of insertion devices (wigglers and undulators) are often represented as a sequence of short dipoles of alternating polarity; however, where the period becomes small compared with the aperture, the three-dimensional nature of the field can start to have effects that cannot be ignored. There can even be cases where “conventional” multipoles designed for special situations (for example, where very wide aperture is required, and where the length of the magnet needs to be short, because of space constraints) can have fringe fields that affect the dynamics to a significant extent.

It is therefore of somewhat more than purely academic interest to consider how two-dimensional multipole representations may be generalised to three dimensions. As usual, there are many different ways to approach the problem: the method that is used will often depend on the problem to be solved. In the following sections, we describe two rather general methods that may be of use in many situations arising in accelerators. First, we consider a field expansion based on Cartesian modes. While this provides some nice illustrations, the Cartesian expansion does have some disadvantages. To address these disadvantages, we describe how a field expansion based on polar coordinates can be performed.

6.1 Cartesian modes

Consider the field given by:

\[
\begin{align*}
B_x &= -B_0 \frac{k_x}{k_y} \sin k_x x \sinh k_y y \sin k_z z, \\
B_y &= B_0 \cos k_x x \cosh k_y y \sin k_z z, \\
B_z &= B_0 \frac{k_z}{k_y} \cos k_x x \sinh k_y y \cos k_z z.
\end{align*}
\]

As may easily be verified, this field satisfies:

\[
\text{curl } \vec{B} = 0.
\]

Furthermore, the equation:

\[
\text{div } \vec{B} = 0
\]

is satisfied if:

\[
k_y^2 = k_x^2 + k_z^2. \quad (46)
\]

We conclude that, as long as the constraint (46) is satisfied, that the fields (43)–(45) provide solutions to Maxwell’s equations in regions with constant permeability, and static (or zero) electric fields. Of course, it is possible to find similar sets of equations but with different “phase” along each of the coordinate axes, and with the hyperbolic trigonometric function appearing for the dependence on \(x\) or \(z\), rather than \(y\).
By superposing fields, with appropriate variations on the form given by Eqs. (43)–(45), it is possible to construct quite general three-dimensional magnetic fields. For example, a slightly more general field than that given by Eqs. (43)–(45) can be obtained simply by superposing fields with different mode numbers and amplitudes:

\[
B_x = - \int \int \tilde{B}(k_x, k_z) \frac{k_x}{k_y} \sin k_x x \sin k_y y \sin k_z z \, dk_x \, dk_z, \tag{47}
\]

\[
B_y = \int \int \tilde{B}(k_x, k_z) \cos k_x x \cosh k_y y \sin k_z z \, dk_x \, dk_z, \tag{48}
\]

\[
B_z = \int \int \tilde{B}(k_x, k_z) \frac{k_z}{k_y} \cos k_x x \sinh k_y y \cos k_z z \, dk_x \, dk_z. \tag{49}
\]

In this form, we see already how to perform a mode decomposition, i.e. how we can determine the coefficients \( \tilde{B}(k_x, k_z) \) as functions of the “mode numbers” \( k_x \) and \( k_z \). If we consider in particular the vertical field component on the plane \( y = y_0 \), then we have from (48):

\[
\frac{B_y}{\cosh k_y y_0} = \int \int \tilde{B}(k_x, k_z) \cos k_x x \sin k_z z \, dk_x \, dk_z.
\]

Hence, \( \tilde{B}(k_x, k_z) \) may be obtained from an inverse Fourier transform of \( B_y(x, z) / \cosh k_y y_0 \). Given field data on a grid over \( x \) and \( z \), then we can perform numerically an inverse discrete Fourier transform, to obtain a set of coefficients \( \tilde{B}(k_x, k_z) \). Note that once we have obtained these coefficients, then we can reconstruct all field components at all points in space. This is an important consequence of the strong constraints on the fields provided by Maxwell’s equations: in general, for a static field, if we know how one field component varies over a two-dimensional plane, then we can deduce how all the field components vary over all space (on and off the plane).

Let us consider an example. To keep things simple, we shall again work with the case where the field is independent of one coordinate: now, however, we shall assume that the fields are independent of the horizontal transverse, rather than the longitudinal coordinate. This may be a suitable model for a planar wiggler or undulator with very wide poles. The model may of course be extended to include dependence of the fields on the horizontal transverse coordinate: although our immediate example strictly deals with a two-dimensional field, the extension to three dimensions is quite straightforward.

Suppose that the mode amplitude function \( \tilde{B}(k_x, k_z) \) has the form:

\[
\tilde{B}(k_x, k_z) = \delta(k_x) \tilde{B}(k_z),
\tag{50}
\]

where \( \delta(k_x) \) is the Dirac delta function. The delta function has the property that, for any function \( f(k_x) \):

\[
\int_{-\infty}^{\infty} \delta(k_x) f(k_x) \, dk_x = f(0).
\]

Using (50) in Eqs. (47) to (49) gives:

\[
B_x = 0,
\]

\[
B_y = \int \tilde{B}(k_z) \cosh k_z y \sin k_z z \, dk_z,
\]

\[
B_z = \int \tilde{B}(k_z) \sinh k_z y \cos k_z z \, dk_z.
\]

There is no horizontal transverse field component, and the vertical and longitudinal field components have no dependence on \( x \): we have a two-dimensional field. In a plane defined by a particular value for the vertical coordinate, \( y = y_0 \), the vertical field component is given by:

\[
B_y(z) = \cosh k_z y_0 \int \tilde{B}(k_z) \sin k_z z \, dk_z.
\]
The mode amplitude function can be obtained from a Fourier transform of the vertical component of the magnetic field on the plane \( y = y_0 \). Usually, we will have a finite set of field data, obtained from a magnet modelling code, or from measurements on a real device.

Suppose that we have a data set of \( 2M + 1 \) vertical field measurements, taken at locations:

\[
y = y_0, \quad z = \frac{m}{M} \hat{z},
\]

where \( m \) is an integer in the range \(-M \leq m \leq M\). The field at any point is given by:

\[
B_y(y, z) = \sum_{m=-M}^{M} \tilde{B}_m \cosh mkz y \sin mkz z,
\]

where:

\[
k_z = \frac{2\pi}{2\hat{z}}.
\]

Note that in this case, the field is antisymmetric about \( z = 0 \), i.e.

\[
B_y(y, -z) = -B_y(y, z).
\]

The mode amplitudes \( \tilde{B}_m \) are obtained by:

\[
\tilde{B}_m = \frac{1}{\cosh mkz y_0} \frac{1}{2M} \sum_{m'=-M}^{M} B_y(y_0, z) \sin m' k z z.
\]

Note that, because of the antisymmetry of the field:

\[
\tilde{B}_{-m} = -\tilde{B}_m.
\]

As a specific numerical example, let us construct an “artificial” data set along a line \( y = y_0 = 0.25 \), and with \( \hat{z} = 3 \). The data are constructed using a function that gives a sinusoidal variation in the field along \( z \) up to \( |z| < 1.25 \); then a continuous and smooth (continuous first derivative) fall-off to zero field for \( |z| > 1.5 \). For this numerical example, we do not worry unduly about units: the reader may assume lengths in cm, fields in kG, or any other preferred units. Initially, we take \( M = 40 \), i.e. we assume we have 81 measurements of the field (or, we have computed the field from a model at 81 equally-spaced points along \( z \); strictly speaking, because we are dealing with the case that the field is antisymmetric in \( z \), we need only half this number of field measurements or computations). The field “data”, the fitted field (reconstructed using the mode amplitudes) in \( y = 0.25 \), and the mode amplitudes, are shown in Fig. 15.

It is interesting to compare with the situation where we have only 31 field measurements or computations, i.e. \( M = 15 \). Using the same function that we used to construct the data set with 81 data points, we produce the fit and the mode amplitudes shown in Fig. 16. Comparing Figs. 15 and 16, we see that in both cases, the fitted field does pass exactly through all the data points. This is a necessary consequence of the fit, which is based on a discrete Fourier transform of the data points. However, using only 31 data points, there is a significant oscillation of the fitted field between the data points in the region \( 1.5 < |z| < 3 \), where the field is actually zero (by construction). This is a consequence of the fact that we have “truncated” some modes with non-negligible amplitude. The mode amplitudes in both cases are the same for mode numbers \(-15 \leq m \leq 15\); but with 81 data points we can determine amplitudes for a larger number of modes, which gives us a more accurate interpolation between the data points.

Having obtained fits to the field in the plane \( y = 0.25 \), we can reconstruct the field at any point, on or off the plane. It is often of interest to look at the mid-plane; usually, this is defined by \( y = 0 \). In this plane, we do not have any field data. However, we can compare the field produced by the fits with 81 data points and with 31 data points: these fields are shown in Fig. 17. We see that the fit based on
Fig. 15: Left: Field data (points) and fit (line) in a magnet with dependence of the field on longitudinal coordinate \( z \). Right: Mode amplitudes. The field data consist of 81 measurements (or computations) at equally-spaced points from \( z = -3 \) to \( z = +3 \), and \( y = 0.25 \).

Fig. 16: Left: Field data (points) and fit (line) in a magnet with dependence of the field on longitudinal coordinate \( z \). Right: Mode amplitudes. The field data consist of 31 measurements (or computations) at equally-spaced points from \( z = -3 \) to \( z = +3 \), and \( y = 0.25 \).

Fig. 17: Field on the plane \( y = 0 \) determined from fits to the data shown in Figs. 15 and 16. Left: fit determined from data set with 81 data points. Right: fit determined from data set with 31 data points.

31 data points produces an essentially identical field on \( y = 0 \) as the fit based on 81 data points. (The data points in each case are taken on the plane \( y = 0.25 \)). This is a consequence of the “suppression” of higher-order modes, that arises from the hyperbolic dependence of the field on the \( y \) coordinate.

To emphasise the significance of the hyperbolic dependence of the field on the vertical coordinate, we can look at the variation of \( B_y \) with \( y \), for a given value of \( z \). We choose \( z = 0.25 \), which corresponds to a peak in the vertical field component as a function of \( z \). The variation of \( B_y \) with \( y \) for the two cases (fit based on 81 data points, and fit based on 31 data points) is shown in Fig. 18.
Fig. 18: Vertical field component as a function of $y$, for $z = 0.25$. The field is determined from fits to the data shown in Figs. 15 and 16. Left: fit determined from data set with 81 data points. Right: fit determined from data set with 31 data points.

Fig. 19: Vertical field component as a function of $y$, for $z = 0.25$. The field is determined from fits to the data shown in Figs. 15 and 16. Left: fit determined from data set with 81 data points. Right: fit determined from data set with 31 data points.

Up to $y = 0.25$, the two fits give essentially the same field. However, if we try to extrapolate beyond this plane (the plane on which the fit was performed), we see dramatically different behaviour. Fig. 19 compares the vertical field component obtained from the two fits (81 data points and 31 data points), again at $z = 0.25$, but now with a range of $y$ from 0 to 0.5. In one case (81 data points), the field increases to a maximum before dropping rapidly. In the other case (31 data points), the field increases monotonically over the range. The reason for the different behaviour is the additional modes in the fit to the set of 81 data points. These higher order modes make only a small additional contribution to the field for $|y| < 0.25$; but for values of the vertical coordinate beyond this value, because of the hyperbolic dependence of $y$, the contribution of these modes becomes increasingly significant, and eventually, dominant.

The behaviour of the field fits for $|y| > 0.25$ is a clear illustration of why it is dangerous to extrapolate the fit beyond the region enclosed by the plane of the fit. In this case, because of the symmetry in the vertical direction, the region enclosed is between the planes $y = -0.25$ and $y = +0.25$. The “safe” region is also bounded in $z$, by $z = -0.3$ and $z = +0.3$; because we use discrete mode numbers in $z$, the fitted fields will in fact be periodic in $z$, and will repeat with period $z = 0.6$. In general, there will be similar periodicity in $x$; however, in this particular example, we analysed a field that was independent of $x$, so the “safe” region of the fit is unbounded in $x$. 

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6.2 Cylindrical modes

The Cartesian modes discussed in Section 6.1 are often useful for describing fields in insertion devices, particularly those that have weak variation of the field with $x$, and periodic behaviour in $z$ (over some range): because the modes “reflect” the geometry, it is often possible to achieve good fits to a given field using a small number of modes. To maximise the region over which the fit is reliable, one needs to choose a plane with a value of $y$ as large as possible, with $x$ and $z$ extending out as far as possible on this plane. For a planar undulator or wiggler, it is often possible to choose a plane close to the pole tips in which $x$ in particular extends over the entire vacuum chamber.

However, for other geometries, the Cartesian modes may not provide a convenient basis. For example, if the magnet has a circular aperture, then the plane that provides the largest range in $x$ is the mid-plane, $y = 0$, and as $y$ increases, the available range in $x$ decreases. To base the fit on the Cartesian basis requires some compromise between the range of reliability in the horizontal transverse and vertical directions.

Fortunately, it is possible to choose an alternative basis for magnets with circular aperture, in which the field fit can be based on the surface of a cylinder inscribed through the magnet. In that case, the radius of the cylinder can be close to the aperture limit, maximising the range of reliability of the fit. The appropriate modes in this case are most easily expressed in cylindrical polar coordinates.

A field with zero divergence and curl (and hence satisfying Maxwell’s equations for static fields in regions with uniform permeability) is given by:

\[ B_r = \int dk_z \sum_n \tilde{B}_n(k_z) I_n(k_z r) \sin n \theta \cos k_z z, \]  
\[ B_\theta = \int dk_z \sum_n \tilde{B}_n(k_z) \frac{n}{k_z r} I_n(k_z r) \cos n \theta \cos k_z z, \]  
\[ B_z = -\int dk_z \sum_n \tilde{B}_n(k_z) I_n(k_z r) \sin n \theta \sin k_z z. \]  

Here, $I_n(k_z r)$ is the modified Bessel function of the first kind, of order $n$. Modified Bessel functions of the first kind for order $n = 0$ to $n = 3$ are plotted in Fig. 20. For small values of the argument $\xi$, the modified Bessel function of order $n$ has the series expansion:

\[ I_n(\xi) = \frac{\xi^n}{2^n n! (1 + n)} + O(n + 1). \]  

For larger values of the argument, the modified Bessel functions $I_n(\xi)$ increase exponentially. This is significant: it means that if we fit a field to data on the surface of a cylinder of given radius, then residuals of the fit will decrease exponentially within the cylinder towards $r = 0$, and increase exponentially outside the cylinder with increasing $r$. The “safe” region of the fit will be within the cylinder.

Note that Eqs. (52)–(54) may be generalised to include different “phases” in the azimuthal angle $\theta$ and the longitudinal coordinate $z$.

An attractive feature of the polar basis is that it is possible to draw a direct connection between the three-dimensional modes in this basis and the multipole components in a two-dimensional field. Consider a mode amplitude $\tilde{B}_n(k_z)$ given (for some particular value of $n$) by:

\[ \tilde{B}_n(k_z) = 2^n \Gamma(1 + n) C_n \frac{\delta(k_z)}{nk_z^{n-1}}, \]  

where $\delta()$ is the Dirac delta function, and $C_n$ is a constant. Substituting these mode amplitudes into Eqs. (52)–(54), using the expansion (55), and performing the integral over $k_z$ gives:

\[ B_r = \sum_n C_n r^{n-1} \sin n \theta. \]
Comparing with Eq. (24), we see that this is a multipole field of order \( n \). Thus, a two-dimensional multipole field is a special case of a three-dimensional field (52)–(54), with mode coefficient given by Eq. (56).

In general, the mode coefficients \( \tilde{B}_n(k_z) \) may be obtained by a Fourier transform of the field on the surface of a cylinder of given radius. For example, it follows from Eq. (52) that:

\[
\frac{B_r}{I_n(k_z r)} = \int dk_z \sum_n \tilde{B}_n(k_z) \sin n\theta \cos k_z z.
\]

An elegant feature of the polar basis, as compared to the Cartesian basis discussed in Section 6.1, is that the modes reflect the real periodicity of the field in the angle coordinate \( \theta \). In the Cartesian basis, the modes were periodic in \( x \), although the field, in general, would not have any periodicity in \( x \).

Since the mode coefficients \( \tilde{B}_n(k_z) \) are related to the multipole coefficients in a two-dimensional field, we can use these coefficients to extend the idea of a multipole to a three-dimensional field. Strictly speaking, the mode coefficients \( \tilde{B}_n(k_z) \) are related to the field by a two-dimensional Fourier transform; however, we can perform a one-dimensional inverse Fourier transform (in the \( z \) variable) to obtain a set of functions which represent, in some sense, the “multipole components” of a three-dimensional field as a function of \( z \). Here, we use the term “multipole components” rather loosely, since a multipole field is strictly defined only in the two-dimensional case (i.e. for a field that is independent of the longitudinal coordinate). A quantity that is perhaps easier to interpret is the contribution to the field at any point made by the mode coefficients \( \tilde{B}_n(k_z) \) with a given \( n \). For \( n = 1 \), the field components \( B_r \) and \( B_\theta \) at any point in \( z \) will behave as for a dipole field; for \( n = 2 \), \( B_r \) and \( B_\theta \) will behave as for a quadrupole field, and so on.

As an illustrative example, we consider the field in a specific device: the wiggler in a damping ring for TESLA (a proposed linear collider) [10]. This wiggler has a peak field of 1.6 T and period 400 mm; the total length of wiggler in each of the TESLA damping rings would be over 400 m. The field in the wiggler has been studied extensively, because of concerns that dynamical effects associated with the nonlinear components in the field would limit the acceptance of the damping ring [11]. A model was constructed for one quarter period of the magnet, which allowed the field at any point within the body of the magnet to be computed. Effects associated with the ends of the wiggler were neglected, but could in principle be included in the study. By performing a mode decomposition using the techniques described
above, it was possible to construct an accurate dynamical model allowing fast tracking to characterise the acceptance of the damping ring. The methods used for the dynamical analysis are beyond the scope of the present discussion; however, we present the results of the analysis relating directly to the field, to illustrate the methods described in this section.

A model of the wiggler was used to compute the magnetic field on a mesh of points bounded by a cylinder of radius 9 mm, within one quarter period of the wiggler. Although all field components were computed on the mesh, which covered the interior of the cylinder as well as the surface, only the radial field component on the surface of the cylinder was used to calculate the mode amplitudes. The fit can be validated by comparing the field “predicted” by the fit with the field data (from the computational model) not used directly in the fitting procedure. A fit achieved using 7 azimuthal and 100 longitudinal modes is shown in Fig. 21. Each plot shows the variation of the vertical field as a function of one Cartesian coordinate, with the other two coordinates fixed at zero. In the vertical direction, the range shown is from the mid-plane of the wiggler to close to the pole tip. Note that the variation in the field in the transverse ($x$ and $y$) directions is very small, less than 0.1% of the maximum field. It appears from Fig. 21, that there is very good agreement between the fit (line) and the field data (circles) within the cylinder on the surface of which the fit was performed.

The quality of the fit can be further illustrated by plotting the residuals, i.e. the difference between the fitted field and the field data. The residuals for the vertical field component on two horizontal planes, $y = 0$ mm and $y = 6$ mm are shown in Fig. 22. Note that to produced “smooth” surface plots, we interpolate between the mesh points used in the computational model. On the mid-plane of the wiggler,
the residuals are less than 1 gauss (recall that the peak field is 1.6 T); the region shown in the left hand plot in Fig. 22 lies entirely within the surface of the cylinder used in the fit. On the plane \( y = 6 \text{ mm} \), the residuals are somewhat larger, and show an exponential increase for large values of the horizontal transverse coordinate: but note that for values of \( x \) larger than about 6.7 mm, the points in the plot are outside the surface of the cylinder used for the fit. In the longitudinal direction, the residuals appear to be dominated by very high frequency modes: this suggests that it may be possible to reduce the residuals still further by increasing the number of longitudinal modes used in the fit. However, this fit was considered to be of sufficient quality to allow an accurate determination of the effect of the wiggler on the beam dynamics to be made.

The tools used for study of the beam dynamics were based on the mode coefficients determined by the fitting procedure. Once a fit has been obtained and shown to be of good quality, then, strictly speaking, further analysis of the field is not required. However, it is interesting to compute, from the mode amplitudes, the contribution to the field in the wiggler from different “multipole” components, as a function of longitudinal position. As described above, the contribution of a multipole of order \( n \) is obtained by a one-dimensional (in the longitudinal dimension) inverse Fourier transform of the mode amplitudes \( \tilde{B}_n(k_z) \). To obtain non-zero values for the contributions from multipoles higher than order \( n = 1 \) (dipole), we need to choose non-zero values for either the \( x \) or \( y \) coordinates at which we compute the field. We choose (arbitrarily) \( x = 8 \text{ mm} \), and \( y = 0 \text{ mm} \). The contributions to the vertical field component from multipoles of order 1 through 7 are shown in Fig. 23. Note that multipoles of even order are forbidden by the symmetry of the wiggler (see Section 5.1). We see from Fig. 23 that the dominant contribution by far is, as expected, the dipole component. The sextupole component is not insignificant; the contributions of higher order multipoles are extremely small, and the high-frequency “oscillation” as a function of longitudinal position is probably unphysical, and the result of noise in the fitting.

It is worth making a few final remarks about mode decompositions for three-dimensional fields. First, as already mentioned, in many cases a full three-dimensional mode decomposition will not be necessary. While this does provide a detailed description of the field in a form suitable for beam dynamics studies, three-dimensional decompositions do rely on a large number of accurate and detailed field measurements. While such “measurements” may be conveniently obtained from a model, it may be difficult or impractical to make such measurements on a real magnet. Fortunately, in many cases, a two-dimensional field description in terms of multipoles is sufficient. Generally, a three-dimensional analysis only need be undertaken where there are grounds to believe that the three-dimensional nature of
the field is likely to have a significant impact on the beam dynamics.

Second, we have already emphasised that to obtain an accurate description of the field within some region in terms of a mode decomposition, the mode amplitudes should be determined by a fit on a surface enclosing the region of interest. Outside the region bounded by the surface of the fit, the fitted field can be expected to diverge exponentially from the real field. However, in choosing the surface for the fit, the geometry of the magnet will impose some constraints. A magnet with a wide rectangular aperture may lend itself to a description using a Cartesian basis (fitting on the surface of a rectangular box); a circular aperture, however, is more likely to require use of a polar basis (fitting on the surface of a cylinder with circular cross-section). Both cases have been described above. It may be appropriate in other cases to perform a fit on the surface of a cylinder with elliptical cross-section. The basis functions in this case involve Mathieu functions. For further details, the reader is referred to work by Dragt [12] and by Dragt and Mitchell [13].

Appendices

A The vector potential

Our analysis of iron-dominated multipole magnets in Section 4.2 was based on the magnetic scalar potential, \( \varphi \). The magnetic flux density can be derived from a scalar potential:

\[
\vec{B} = -\nabla \varphi
\]

in the case that the flux density has vanishing divergence and curl:

\[
\text{div} \, \vec{B} = \text{curl} \, \vec{B} = 0.
\]

More generally (in particular, where the flux density has non-vanishing curl) one derives the magnetic flux density from a vector potential \( \vec{A} \), using:

\[
\vec{B} = \text{curl} \, \vec{A}.
\]
Although we have not required the vector potential in our discussion of Maxwell’s equations for accelerator magnets, it is sometimes used in analysis of beam dynamics. In particular, descriptions of the dynamics based on Hamiltonian mechanics generally use the vector potential rather than the magnetic flux density or the magnetic scalar potential. We therefore include here a brief discussion of the vector potential, paying attention to aspects relevant to the descriptions we have developed for two-dimensional and three-dimensional magnet fields.

First, we note that the divergence of any curl is identically zero:
\[
\text{div} \, \text{curl} \, \vec{V} \equiv 0,
\]
for any differentiable vector field \( \vec{V} \). Thus, if we write \( \vec{B} = \text{curl} \, \vec{A} \), then Maxwell’s equation (2):
\[
\text{div} \, \vec{B} = 0,
\]
is automatically satisfied. Maxwell’s equation (3) in uniform media (constant permeability), with zero current and static electric fields gives:
\[
\text{curl} \, \vec{B} = \mu \vec{J},
\]
where \( \vec{J} \) is the current density. This leads to the equation for the vector potential:
\[
\text{curl} \, \text{curl} \, \vec{A} \equiv \text{grad} \, (\text{div} \, \vec{A}) - \nabla^2 \vec{A} = \mu \vec{J}.
\]
Eq. (A.3) is a second-order differential equation for the vector potential in a medium with permeability \( \mu \), where the current density is \( \vec{J} \). This appears harder to solve than the first-order differential equation for the magnetic flux density, Eq. (A.2). However, Eq. (A.3) may be simplified significantly, if we apply an appropriate \textit{gauge condition}. To understand what this means, recall that the magnetic flux density is given by the curl of the vector potential, and that the curl of the gradient of any scalar field is identically zero. Thus, we can add the gradient of a scalar field to a vector potential, and obtain a new vector potential that gives the same flux density as the old one. That is, if:
\[
\vec{A} = \text{curl} \, \vec{A},
\]
and:
\[
\vec{A}' = \vec{A} + \text{grad} \, \psi,
\]
for an arbitrary differentiable scalar field \( \psi \), then:
\[
\text{curl} \, \vec{A}' = \text{curl} \, \vec{A} + \text{curl} \, \text{grad} \, \psi = \text{curl} \, \vec{A} = \vec{B}.
\]
In other words, the vector potential \( \vec{A}' \) leads to exactly the same flux density as the vector potential \( \vec{A} \). Since the dynamics of a given system are determined by the fields rather than the potentials, either \( \vec{A}' \) or \( \vec{A} \) is a valid choice for the description of the system. Eq. (A.4) is known as a \textit{gauge transformation}. The consequence of having the freedom to make a gauge transformation means that the vector potential for any given system is not uniquely defined: given some particular vector potential, it is always possible to make a gauge tranformation without any change in the physical observables of a system. The analogue in the case of electric fields, of course, is that the “zero” of the electric scalar potential can be chosen arbitrarily: only \textit{changes} in potential (i.e. energy) are observable, so given some particular scalar potential field, it is possible to add a constant (that is, a quantity independent of position) and obtain a new scalar potential that gives the same physical observables as the original scalar potential.

For magnetostatic fields, we can use a gauge transformation to simplify Eq. (A.3). Suppose we have obtained a vector potential \( \vec{A} \) for some particular physical system. Define a scalar field \( \psi \), which satisfies:
\[
\nabla^2 \psi = -\text{div} \, \vec{A}.
\]

Then define:
\[ \vec{A}' = \vec{A} + \text{grad} \psi. \]
Since \( \vec{A}' \) and \( \vec{A} \) are related by a gauge transformation, they lead to the same magnetic flux density, and the same physical observables for the system. However, the divergence of \( \vec{A}' \) vanishes:
\[ \text{div} \vec{A}' = \text{div} \vec{A} + \text{div} \text{grad} \psi = -\nabla^2 \psi + \nabla^2 \psi = 0, \]
where we have used Eq. (A.5). Thus, given any vector potential, we can make a gauge transformation to find a new vector potential that gives the same magnetic flux density, but has vanishing divergence. The \textit{gauge condition}:
\[ \text{div} \vec{A} = 0, \quad (A.6) \]
is known as the \textit{Coulomb gauge}. It is possible to work with other gauge conditions (for example, for time-dependent electromagnetic fields the Lorenz gauge condition is often more appropriate); however, for our present purposes, the Coulomb gauge leads to a simplification of Eq. (A.3), which now becomes:
\[ \nabla^2 \vec{A} = -\mu \vec{J}. \quad (A.7) \]
Eq. (A.7) is Poisson’s equation for a vector field. Note that despite being a second-order differential equation, it is in a sense simpler than Maxwell’s equation (3), since we have “decoupled” the components of the vectors; that is, we have a set of three uncoupled second-order differential equations, where each equation relates a component of the vector potential to the corresponding component of the current density. Eq. (A.7) has the solution:
\[ \vec{A}(\vec{r}) = -\frac{\mu}{4\pi} \int \frac{\vec{J}(\vec{r}')}{|\vec{r} - \vec{r}'|} d^3 r'. \]
In this form, we see that the potential at a point in space is inversely proportional to the distance from the source.

Now, consider the potential given by:
\[ A_x = 0, \quad A_y = 0, \quad A_z = -\text{Re} \frac{C_n(x + iy)^n}{n}. \quad (A.8) \]
Taking derivatives, we find that:
\[ \frac{\partial A_z}{\partial x} = -\text{Re} C_n(x + iy)^{n-1}, \]
\[ \frac{\partial A_z}{\partial y} = \text{Im} C_n(x + iy)^{n-1}. \]
Then, since \( A_x \) and \( A_y \) are zero, we have:
\[ \vec{B} = \text{curl} \vec{A} = \left( \frac{\partial A_z}{\partial y}, -\frac{\partial A_z}{\partial x}, 0 \right). \]
Hence:
\[ B_y + iB_x = C_n(x + iy)^{n-1}, \quad (A.9) \]
which is just the multipole field. Thus, Eq. (A.8) is a potential that gives a multipole field. Note also that, since \( A_z \) is independent of \( z \), this potential satisfies the Coulomb gauge condition (A.6):
\[ \text{div} \vec{A} = \frac{\partial A_x}{\partial x} + \frac{\partial A_y}{\partial y} + \frac{\partial A_z}{\partial z} = 0. \]
An advantage of working with the vector potential in the Coulomb gauge is that, for multipole fields, the transverse components of the vector potential are both zero. This simplifies, to some extent, the Hamiltonian equations of motion for a particle moving through a multipole field. However, note that the longitudinal component $B_z$ of the magnetic flux density is zero in this case. To generate a solenoidal field, with $B_z$ equal to a non-zero constant, we need to introduce non-zero components for $A_x$, or $A_y$, or both. For example, a solenoid field with flux density $B_{\text{sol}}$ may be derived from the vector potential:

$$A_x = -\frac{1}{2} B_{\text{sol}} y, \quad A_y = \frac{1}{2} B_{\text{sol}} x.$$ 

Let us return for a moment to the case of multipole fields. If we work in a gauge in which the transverse components of the vector potential are both zero, then the field components are given by:

$$B_y = -\frac{\partial A_z}{\partial x}, \quad B_x = \frac{\partial A_z}{\partial y}.$$ 

From these expressions, we see that if we take any two points with the same $y$ coordinate, then the difference in the vector potential between these two points is given by the "flux" passing through a line between these points:

$$\Delta A_z = -\int B_y \, dx.$$ 

Similarly for any two points with the same $x$ coordinate:

$$\Delta A_z = \int B_x \, dy.$$ 

In general, for a field that is independent of $z$, and working in a gauge where $A_x = A_y = 0$, we can write:

$$\Delta A_z = \frac{\Delta \Phi}{\Delta z}, \quad (A.10)$$

where $\Delta A_z$ is the change in the vector potential between two points $P_1$ and $P_2$ in a given plane $z = z_0$; and $\Delta \Phi$ is the magnetic flux through a rectangular “loop” with vertices $P_1$, $P_2$, $P_3$ and $P_4$: see Fig. A.1. $P_3$ and $P_4$ are points obtained by transporting $P_1$ and $P_2$ a distance $\Delta z$ parallel to the $z$ axis. Eq. (A.10) can also be obtained by applying Stokes’ theorem to the loop $P_1 P_2 P_3 P_4$, with the relationship (A.1) between $\vec{B}$ and $\vec{A}$:

$$\int A \cdot d\vec{l} = \int \text{curl} \, \vec{A} \cdot d\vec{S} = \int \vec{B} \cdot d\vec{S},$$

hence:

$$A_z(P_2) \Delta z - A_z(P_1) \Delta z = \Delta \Phi.$$

Finally, we give the vector potentials corresponding to three-dimensional fields. In the Cartesian basis, with the field given by Eqs. (43)–(45), a possible vector potential (in the Coulomb gauge) is:

$$A_x = 0,$$

$$A_y = B_0 \frac{k_z}{k_x k_y} \sin k_x x \sinh k_y y \cos k_z z,$$

$$A_z = -B_0 \frac{1}{k_x} \sin k_x x \cosh k_y y \sin k_z z.$$ 

In the polar basis, with the field given by Eqs. (52)–(54), a possible vector potential is:

$$A_r = -\frac{r}{m} I_m(k_z r) \cos n\theta \sin k_z z,$$

$$A_\theta = 0,$$

$$A_z = -\frac{r}{2m} I'_m(k_z r) \cos n\theta \sin k_z z.$$ 

However, note that this potential does not satisfy the Coulomb gauge condition.
Fig. A.1: Interpretation of the vector potential in a two-dimensional magnetic field (i.e. a field that is independent of \(z\)). The change in the vector potential between \(P_1\) and \(P_2\) is equal to the flux of the magnetic field through the loop \(P_1P_2P_3P_4\), divided by \(\Delta z\).

References