

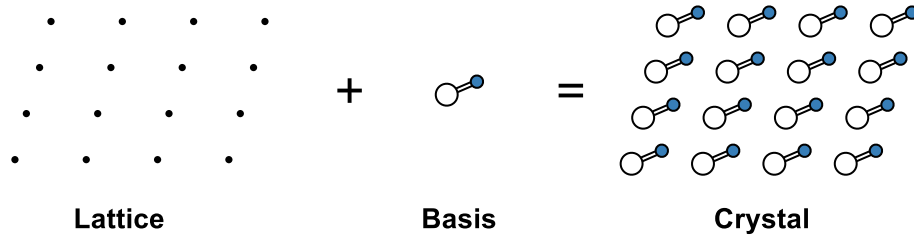
Solid State Physics

May 31, 2018

Chapter 1

Crystal Structure

1.1 Direct lattice



The mathematical description of the crystal consists of two parts: the *lattice* which is a periodic grid of points extending over space and the *basis* which in this context is the set of ions repeating at every lattice point.

The lattice can be described by the set of so called *lattice vectors* $\{\mathbf{R}\}$ given by

$$\mathbf{R} \equiv \mathbf{R}_{n_1 n_2 n_3} = n_1 \mathbf{a}_1 + n_2 \mathbf{a}_2 + n_3 \mathbf{a}_3, \quad (1.1)$$

where $n_i \in \mathbb{Z}$ and \mathbf{a}_i are linearly independent vectors in the real space called the *primitive vectors*. The primitive vectors have the dimensions of length owing to which the lattice spanned by them is referred to as the *direct lattice*. It should be noted the choice of the primitive vectors is not unique; any non-collinear set of the lattice vectors can be used.

The lattice can be infinite or finite depending on the allowed values of n_i . Quite often it is convenient to deal with infinite lattice which extends over all spatial dimensions. In such case the lattice is known as the *Bravais lattice*.

An important concept of *unit cell* builds upon the lattice. A unit cell is a volume of space that can fill the entire space without overlaps or leaving gaps behind when translated by a suitable set of Bravais lattice vectors. A *primitive unit cell* is a special set of unit cells which enclose precisely one lattice point. Therefore, if the number density of lattice points is n , the primitive unit cell volume v is

$$v = \frac{1}{n} \quad (1.2)$$

regardless the shape of the cell. Sometimes it is, however, easier to work with a *conventional unit cell* instead. For example, a primitive unit cells of a BCC (body-centered cubic) lattice can be difficult to work with since their angles are not orthogonal. The usual choice is to

use a cubical cell containing two lattice points instead. Note that whereas primitive unit cells can be translated by all the lattice vectors without any overlaps, the same is not true for the conventional ones.

There is, however, a common used way to choose a primitive unit cell that has the full symmetry of the lattice. Consider a single lattice point and take all the points in its vicinity that are nearer to it than any other point in the lattice. The volume covered by those points leads to a unique primitive unit cell which is known as the *Wigner-Seitz cell*. The concept is worth remembering as it has an important role in the subsequent discussion.

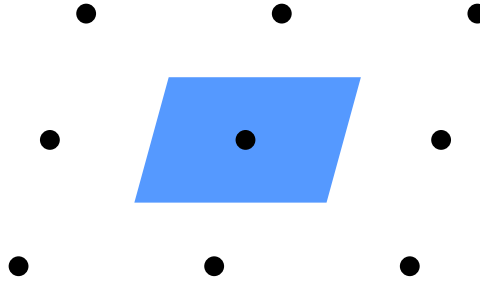


Figure 1.1: Wigner-Seitz cell of an oblique 2D-lattice. The blue area covers the points closest to the lattice point in the center.

1.2 Reciprocal Lattice

In the scope of this for a general plane wave in terms of position \mathbf{r} and time t is written as

$$f(\mathbf{r}, t) = Ae^{i\mathbf{k}\cdot\mathbf{r} - i\omega t}. \quad (1.3)$$

The wavevector \mathbf{k} points at the direction of propagation of the wave with the length of $k \equiv |\mathbf{k}| = 2\pi/\lambda$, where λ is the wavelength. The *angular frequency* (often just frequency) ω is related to k by $\omega = vk$, where v is the *phase velocity* of the wave.

Dropping out the time-dependent part of the phase $e^{-i\omega t}$, consider a plane wave $e^{i\mathbf{k}\cdot\mathbf{r}}$ with respect to the direct lattice vectors \mathbf{R} defined by Eq. (1.1). In order for the wave to have exactly the same periodicity as the lattice, it needs to have the same value at every single lattice point. This condition can be written as

$$e^{i\mathbf{k}\cdot\mathbf{R}} = 1 \quad \text{or} \quad \mathbf{k} \cdot \mathbf{R} = 2\pi n, \quad n \in \mathbb{Z} \quad (1.4)$$

for all \mathbf{R} . The infinite set of wavevectors \mathbf{k} which fulfil the condition, defines the so-called *reciprocal lattice* which is an extremely useful concept in solid state physics which we shall see in the subsequent chapters. The reciprocal lattice vectors are usually denoted by \mathbf{G} .¹

Given the primitive vectors \mathbf{a}_i of the direct lattice, any reciprocal lattice can be written as

$$\mathbf{G} \equiv \mathbf{G}_{hkl} = h\mathbf{b}_1 + k\mathbf{b}_2 + l\mathbf{b}_3, \quad (1.5)$$

¹Some sources use \mathbf{h} instead.

where $h, k, l \in \mathbb{Z}$ and the reciprocal primitive vectors \mathbf{b}_i are

$$\mathbf{b}_1 = 2\pi \frac{\mathbf{a}_2 \times \mathbf{a}_3}{\mathbf{a}_1 \cdot \mathbf{a}_2 \times \mathbf{a}_3} \quad \mathbf{b}_2 = 2\pi \frac{\mathbf{a}_3 \times \mathbf{a}_1}{\mathbf{a}_1 \cdot \mathbf{a}_2 \times \mathbf{a}_3} \quad \mathbf{b}_3 = 2\pi \frac{\mathbf{a}_1 \times \mathbf{a}_2}{\mathbf{a}_1 \cdot \mathbf{a}_2 \times \mathbf{a}_3}. \quad (1.6)$$

Thus we see that the set of all \mathbf{G} itself forms a Bravais lattice in the *reciprocal space* (with linear dimensions measured in units of inverse length).

The reciprocal lattice can be also understood from an alternative viewpoint. We may describe the direct lattice in terms of the Dirac delta function as follows:

$$\rho_R(\mathbf{r}) = \sum_{\mathbf{R}} \delta(\mathbf{r} - \mathbf{R}), \quad (1.7)$$

where the sum goes over all lattice vectors. Therefore $\rho_R(\mathbf{r})$ is function whose value is an (integrable) infinity at the lattice points and zero elsewhere. Now taking the Fourier transform² of the lattice we obtain

$$\mathcal{F}[\rho_R](\mathbf{k}) = \sum_{\mathbf{R}} \int d\mathbf{r} \delta(\mathbf{r} - \mathbf{R}) e^{-i\mathbf{k} \cdot \mathbf{r}} = \sum_{\mathbf{R}} e^{-i\mathbf{k} \cdot \mathbf{R}}. \quad (1.8)$$

The result can be interpreted as the Fourier expansion of a function periodic in terms of \mathbf{k} . Since the coefficient is constant for each term, we see that the series is that of the Dirac delta, which is repeated at every \mathbf{k} which fulfils $e^{i\mathbf{k} \cdot \mathbf{R}} = 1$. Thus in terms of the reciprocal lattice vectors \mathbf{G} , we may write

$$\mathcal{F}[\rho_R](\mathbf{k}) = (2\pi)^3 \sum_{\mathbf{G}} \delta(\mathbf{k} - \mathbf{G}), \quad (1.9)$$

where the factor of $(2\pi)^3$ comes from the normalisation. Therefore we have shown that the reciprocal lattice can be obtained from the direct lattice *via* the Fourier transform and *vice versa*. This exactly the same procedure how one makes the conversion between the momentum and position in quantum mechanics. We'll see later how the reciprocal space and momentum are related, but it is already now useful to keep in mind, that the reciprocal space is basically just the momentum space in suitable units.

²In this document, the convention is that

$$\mathcal{F}[f](\mathbf{k}) = \int d\mathbf{r} f(\mathbf{r}) e^{-i\mathbf{k} \cdot \mathbf{r}} \quad \mathcal{F}^{-1}[F](\mathbf{r}) = \frac{1}{(2\pi)^3} \int d\mathbf{k} F(\mathbf{k}) e^{i\mathbf{k} \cdot \mathbf{r}}$$

Chapter 2

Diffraction of X-rays

2.1 Elastic Scattering

Consider a charge density subject to a linearly polarized (classical) electromagnetic magnetic wave whose electric part is given by $\mathbf{E}_0 \cos(\mathbf{k} \cdot \mathbf{r} - \omega t)$. However, since the Maxwell equations and classical mechanics do not mix the real and imaginary parts, we may for mathematical convenience work with the complex exponential plane wave

$$\mathbf{E}(\mathbf{r}, t) = \mathbf{E}_0 e^{i\mathbf{k} \cdot \mathbf{r} - i\omega t}, \quad (2.1)$$

from which the physical field is recovered by taking its real part. Neglecting the effect of the magnetic field, the Lorentz force affecting an miniscule charge dq in an infinitesimal volume element is

$$d\mathbf{F} = dq \mathbf{E}. \quad (2.2)$$

Substituting the force from Newton's II law $\mathbf{F} = m\ddot{\mathbf{r}} \Rightarrow d\mathbf{F} = dm\ddot{\mathbf{r}}$, we obtain

$$\ddot{\mathbf{r}} = \frac{dq}{dm} \mathbf{E}. \quad (2.3)$$

Since the nuclei of atoms are orders of magnitude more massive than electrons, we may ignore their movement and thus contribution to the scattering. Thus we can concentrate on the electronic contribution for which $dq/dm = -e/m$, where e is the elementary charge and m is the electron mass. Therefore by substituting the wavefield (2.1) to Eq. (2.3), we get

$$\ddot{\mathbf{r}} = -\frac{e}{m} \mathbf{E}_0 e^{i\mathbf{k} \cdot \mathbf{r} - i\omega t} \quad (2.4)$$

Integrating once with respect to time, we find that the velocity of an infinitesimal charge dq initially at rest is

$$\dot{\mathbf{r}} = -\frac{ie}{m\omega} \mathbf{E}_0 e^{i\mathbf{k} \cdot \mathbf{r} - i\omega t}. \quad (2.5)$$

Therefore the current density \mathbf{J} owing to oscillating electrons becomes

$$\mathbf{J}(\mathbf{r}, t) = -en(\mathbf{r})\dot{\mathbf{r}} = \frac{ie^2}{m\omega} n(\mathbf{r}) \mathbf{E}_0 e^{i\mathbf{k} \cdot \mathbf{r} - i\omega t} \quad (2.6)$$

where $n(\mathbf{r})$ is the number density of electrons without the influence of external wavefield.

What kind of wavefield does the oscillating charge density emit? For that we need to solve the wave equation for the vector potential \mathbf{A} :

$$\nabla^2 \mathbf{A} - \frac{1}{c^2} \frac{\partial^2 \mathbf{A}}{\partial t^2} = -\mu_0 \mathbf{J} \quad (2.7)$$

Since the $\nabla^2 - c^{-2} \partial_t^2$ is a linear operator, we seek the solution using the *Green's function*. The Green's function $G(\mathbf{r}, t; \mathbf{r}', t')$ is the solution to the problem

$$\nabla^2 G(\mathbf{r}, t; \mathbf{r}', t') - \frac{1}{c^2} \frac{\partial^2 G(\mathbf{r}, t; \mathbf{r}', t')}{\partial t^2} = -\delta(\mathbf{r} - \mathbf{r}') \delta(t - t'). \quad (2.8)$$

The solution to the inhomogeneous problem (2.7) is then

$$\mathbf{A}(\mathbf{r}, t) = \mu_0 \int d\mathbf{r}' dt' \mathbf{J}(\mathbf{r}', t') G(\mathbf{r}, t; \mathbf{r}', t'). \quad (2.9)$$

Since the retarded Green's function describing the solutions propagating away from the source is

$$G(\mathbf{r}, t; \mathbf{r}', t') = \frac{1}{4\pi} \frac{\delta(t' - t_r)}{|\mathbf{r} - \mathbf{r}'|}, \quad (2.10)$$

where the retarded time $t_r = t - |\mathbf{r} - \mathbf{r}'|/c$ takes into account the finite speed of light, the formal solution to the wave equation is

$$\mathbf{A} = \frac{\mu_0}{4\pi} \int d\mathbf{r}' dt' \frac{\mathbf{J}(\mathbf{r}', t')}{|\mathbf{r} - \mathbf{r}'|} \delta(t' - t_r). \quad (2.11)$$

Substituting \mathbf{J} from Eq. (2.6) to Eq. (2.11), we get

$$\begin{aligned} \mathbf{A} &= \frac{\mu_0}{4\pi} \frac{ie^2}{m\omega} \mathbf{E}_0 \int d\mathbf{r}' dt' \frac{\delta(t' - t_r)}{|\mathbf{r} - \mathbf{r}'|} n(\mathbf{r}') e^{i\mathbf{k} \cdot \mathbf{r}' - i\omega t'} \\ &= \frac{\mu_0}{4\pi} \frac{ie^2}{m\omega} \mathbf{E}_0 \int d\mathbf{r}' \frac{n(\mathbf{r}')}{|\mathbf{r} - \mathbf{r}'|} e^{i\mathbf{k} \cdot \mathbf{r}' - i\omega t_r} \end{aligned} \quad (2.12)$$

Considering distances considerably larger than atomic dimensions, we may approximate $|\mathbf{r} - \mathbf{r}'|^{-1} \approx r^{-1}$. In addition, by substituting $t_r = t - |\mathbf{r} - \mathbf{r}'|/c$, we find

$$\begin{aligned} \mathbf{A} &= \frac{\mu_0}{4\pi} \frac{ie^2}{m\omega} \frac{1}{r} \mathbf{E}_0 \int d\mathbf{r}' n(\mathbf{r}') e^{i\mathbf{k} \cdot \mathbf{r}' + i\omega |\mathbf{r} - \mathbf{r}'|/c - i\omega t} \\ &= \frac{\mu_0}{4\pi} \frac{ie^2}{m\omega} \frac{1}{r} \mathbf{E}_0 e^{-i\omega t} \int d\mathbf{r}' n(\mathbf{r}') e^{i\mathbf{k} \cdot \mathbf{r}' + ik|\mathbf{r} - \mathbf{r}'|}, \end{aligned} \quad (2.13)$$

where in the last step $k = \omega/c$ has been used. By denoting the outgoing wavevector by \mathbf{k}' , we recognize that

$$k|\mathbf{r} - \mathbf{r}'| = \mathbf{k}' \cdot (\mathbf{r} - \mathbf{r}') \quad (2.14)$$

In principle, the direction of \mathbf{k}' depends on \mathbf{r}' but far away from the atom we can regard it being constant. Thus

$$\mathbf{A} = \frac{\mu_0}{4\pi} \frac{ie^2}{m\omega} \frac{1}{r} \mathbf{E}_0 e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t} \int d\mathbf{r}' n(\mathbf{r}') e^{i(\mathbf{k} - \mathbf{k}') \cdot \mathbf{r}'}. \quad (2.15)$$

Since we haven't calculated the scalar potential φ , we can't directly obtain the electric field $\mathbf{E} = -\nabla\varphi - \partial\mathbf{A}/\partial t$ from the potential. Therefore we first compute the magnetic field

$$\mathbf{B}' = \nabla \times \mathbf{A} = -\frac{\mu_0}{4\pi} \frac{e^2}{m\omega} \frac{1}{r} \mathbf{k}' \times \mathbf{E}_0 e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t} \int d\mathbf{r}' n(\mathbf{r}') e^{i(\mathbf{k} - \mathbf{k}') \cdot \mathbf{r}'}, \quad (2.16)$$

where we have benefited from the fact that $e^{i\mathbf{k}' \cdot \mathbf{r}}$ varies considerably faster than r^{-1} , thus dominating the derivative. The electric field is now found using the Maxwell's equations. Assuming the plane wave form $\mathbf{E}' e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t}$, the time derivative of the field is

$$\frac{\partial}{\partial t} \mathbf{E}' e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t} = -i\omega \mathbf{E}' e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t}. \quad (2.17)$$

In the absence of free currents, $\partial \mathbf{E}' / \partial t = c^2 \nabla \times \mathbf{B}'$. Therefore for the plane waves in question $\mathbf{E}' = -c^2 / \omega \mathbf{k}' \times \mathbf{B}'$. Thus

$$\mathbf{E}' = \frac{1}{4\pi\epsilon_0} \frac{e^2}{mc^2} \frac{1}{r} \frac{\mathbf{k}' \times (\mathbf{k}' \times \mathbf{E}_0)}{k^2} e^{i\mathbf{k}' \cdot \mathbf{r} - i\omega t} \int d\mathbf{r}' n(\mathbf{r}') e^{i(\mathbf{k} - \mathbf{k}') \cdot \mathbf{r}'}, \quad (2.18)$$

where $c^2 = 1/\epsilon_0\mu_0$ has been used. The quantity we measure in the scattering experiments is not the electric field, however, but the time-averaged intensity of the field. Since the intensity is proportional to $|\mathbf{E}|^2$, the ratio of the scattered intensity I' to the incident intensity I is

$$\frac{I'}{I} = \frac{|\mathbf{E}'|^2}{|\mathbf{E}|^2} \quad (2.19)$$

Since

$$\begin{aligned} \mathbf{k}' \times (\mathbf{k}' \times \mathbf{E}_0) &= (\mathbf{k}' \cdot \mathbf{E}_0) \mathbf{k}' - k^2 \mathbf{E}_0 = k |\mathbf{E}_0| \cos \chi \mathbf{k}' - k^2 \mathbf{E}_0 \\ \Rightarrow |\mathbf{k}' \times (\mathbf{k}' \times \mathbf{E}_0)|^2 &= k^2 |\mathbf{E}_0|^2 \cos^2 \chi \mathbf{k}' \cdot \mathbf{k}' - 2k^3 |\mathbf{E}_0| \cos \chi \mathbf{k}' \cdot \mathbf{E}_0 + k^4 |\mathbf{E}_0|^2 \\ &= (1 - \cos^2 \chi) k^4 |\mathbf{E}_0|^2 = \sin^2 \chi k^4 |\mathbf{E}_0|^2, \end{aligned} \quad (2.20)$$

where χ is the angle between the polarization of the incident wave and the scattering vector, we obtain

$$\frac{I'}{I} = \left(\frac{e^2}{4\pi\epsilon_0 mc^2} \right)^2 \frac{1}{r^2} \sin^2 \chi \left| \int d\mathbf{r}' n(\mathbf{r}') e^{i(\mathbf{k} - \mathbf{k}') \cdot \mathbf{r}'} \right|^2. \quad (2.21)$$

The expression of intensity ratio still contains the factor of r^{-2} which is basically a parameter depending on the experiment setup, not the system itself. In scattering studies, it is customary to use the so-called *differential scattering cross section* $d\sigma/d\Omega$ as the quantity of interest instead. By definition $d\sigma/d\Omega$ is the fraction of intensity (or particles) scattered into a given direction. Since the infinitesimal area $dA = r^2 d\Omega$, where $d\Omega$ is the solid angle corresponding to the area, we find that the differential cross-section is

$$\frac{d\sigma}{d\Omega} = \left(\frac{e^2}{4\pi\epsilon_0 mc^2} \right)^2 \sin^2 \chi \left| \int d\mathbf{r}' n(\mathbf{r}') e^{-i\mathbf{q} \cdot \mathbf{r}'} \right|^2. \quad (2.22)$$

where the *scattering vector* $\mathbf{q} = \mathbf{k}' - \mathbf{k}$.

Chapter 3

Magnetism

3.1 Pauli equation

A suitable starting point to investigate magnetism in atomic systems is the Pauli equation, which describes the dynamics of a non-relativistic spin-1/2 particle in an electromagnetic field. A single electron with the mass m and the charge $-e$ the standard form of the time-independent equation is

$$\left[\frac{(\mathbf{p} + e\mathbf{A})^2}{2m} + \frac{e\hbar}{m} \mathbf{s} \cdot \mathbf{B} - e\phi \right] |\psi\rangle = E|\psi\rangle, \quad (3.1)$$

where ϕ is the electric scalar potential, \mathbf{A} is the (unquantized) vector potential, the magnetic field $\mathbf{B} = \nabla \times \mathbf{A}$ and $\mathbf{s} = (s_x, s_y, s_z)$ are the components of the (unitless) spin-1/2 operator. Quite often s_i are given in terms of Pauli matrices:

$$s_x = \frac{1}{2} \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \quad s_y = \frac{1}{2} \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix} \quad s_z = \frac{1}{2} \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}. \quad (3.2)$$

In such a representation the state $|\psi\rangle$ is represented by a two-component spinor

$$|\psi\rangle = \begin{pmatrix} \psi_+ \\ \psi_- \end{pmatrix} = \psi_+ |\uparrow\rangle + \psi_- |\downarrow\rangle. \quad (3.3)$$

Let us consider an ion with N electrons. The Pauli equation now reads

$$\sum_{i=1}^N \left[\frac{(\mathbf{p}_i + e\mathbf{A}_i)^2}{2m} + \frac{e\hbar}{m} \mathbf{s}_i \cdot \mathbf{B}_i - e\phi_i \right] |\psi\rangle = E|\psi\rangle, \quad (3.4)$$

where

$$\phi_i = \frac{1}{4\pi\epsilon_0} \left(\frac{Z_n e}{|\mathbf{r}_i|} - \frac{1}{2} \sum_{j \neq i} \frac{e}{|\mathbf{r}_i - \mathbf{r}_j|} \right) \quad (3.5)$$

is the electrostatic potential due to the nuclear charge and other electrons¹. \mathbf{r}_i is the position vector of i th electron. Assuming that \mathbf{A}_i are in Coulomb gauge *i.e.* $\nabla \cdot \mathbf{A}_i = 0$, we may expand²

$$(\mathbf{p}_i + e\mathbf{A}_i)^2 = \mathbf{p}_i^2 + \underbrace{e\mathbf{p}_i \cdot \mathbf{A}_i + 2e\mathbf{A}_i \cdot \mathbf{p}_i}_{=0} + e^2 \mathbf{A}_i^2 = \mathbf{p}_i^2 + 2e\mathbf{A}_i \cdot \mathbf{p}_i + e^2 \mathbf{A}_i^2. \quad (3.6)$$

¹Factor of $\frac{1}{2}$ is added to avoid double counting

²Note that \mathbf{p}_i and \mathbf{A}_i are operators that do not commute. The obtained result was derived using a test function.

Therefore the total Hamiltonian can be written as $\hat{H} = \hat{H}_0 + \hat{H}_M$, where

$$\hat{H}_0 = \sum_i \frac{\mathbf{p}_i^2}{2m} - e\phi_i, \quad (3.7)$$

is the electronic part without magnetic interactions and

$$\hat{H}_M = \sum_i \frac{e}{m} \mathbf{A}_i \cdot \mathbf{p}_i + \frac{e^2 \mathbf{A}_i^2}{2m} + \frac{e\hbar}{m} \mathbf{s}_i \cdot \mathbf{B}_i, \quad (3.8)$$

describes the contribution of the magnetic field.

3.2 Magnetic energy

Omitting the spin-orbit interaction, consider an ion which is put into a homogeneous magnetic field $\mathbf{B}_i = \mathbf{B}_0$. In this case the vector potential can be chosen so that

$$\mathbf{A}_i = -\frac{1}{2} \mathbf{r}_i \times \mathbf{B}_0 \quad (3.9)$$

Therefore³

$$\mathbf{A}_i \cdot \mathbf{p}_i = -\frac{1}{2} (\mathbf{r}_i \times \mathbf{B}_0) \cdot \mathbf{p}_i = \frac{1}{2} \mathbf{B}_0 \cdot (\mathbf{r}_i \times \mathbf{p}_i) = \frac{\hbar}{2} \mathbf{B}_0 \cdot \mathbf{l}_i, \quad (3.10)$$

where in the last step we have defined a unitless orbital angular momentum operator $\hbar \mathbf{l}_i = \mathbf{r}_i \times \mathbf{p}_i$. Without the loss of generality, we may choose z -axis to be parallel with \mathbf{B}_0 . Thus

$$\mathbf{A}_i \cdot \mathbf{p}_i = \frac{\hbar B_0}{2} l_{z,i}, \quad (3.11)$$

where $l_{z,i}$ is the z -component of \mathbf{l}_i . Similarly

$$\mathbf{A}_i^2 = \frac{1}{4} (\mathbf{r}_i \times \mathbf{B}_0)^2 = \frac{1}{4} (y_i B_0 \hat{x} - x_i B_0 \hat{y})^2 = \frac{B_0^2}{4} (x_i^2 + y_i^2) \quad (3.12)$$

Substituting the derived expressions for $\mathbf{A}_i \cdot \mathbf{p}_i$ and \mathbf{A}_i^2 , and $\mathbf{s}_i \cdot \mathbf{B}_i = B_0 s_{z,i}$ into Eq. (3.8), we obtain

$$\hat{H}_M = \sum_i \frac{e\hbar B_0}{2m} l_{z,i} + \frac{e^2 B_0^2}{8m} (x_i^2 + y_i^2) + \frac{e\hbar B_0}{m} s_{z,i}. \quad (3.13)$$

We may simplify the result by defining the Bohr magneton

$$\mu_B = \frac{e\hbar}{2m} = 5.788 \cdot 10^{-5} \frac{\text{eV}}{\text{T}} = 9.274 \cdot 10^{-24} \frac{\text{J}}{\text{T}}. \quad (3.14)$$

Thus

$$\hat{H}_M = \mu_B B_0 (L_z + 2S_z) + \frac{e^2 B_0^2}{8m} \sum_i (x_i^2 + y_i^2), \quad (3.15)$$

where $L_z = \sum_i l_{z,i}$ and $S_z = \sum_i s_{z,i}$ are the z -components of the total orbital angular momentum and spin angular momentum operators \mathbf{L} and \mathbf{S} , respectively.

Since magnetic energy is rather small compared to the energy of an unperturbed atom, we may obtain it by using the perturbation theory. Up to the first order the magnetic energy is

$$\Delta E_M = \mu_B B_0 \langle \psi | L_z + 2S_z | \psi \rangle + \frac{e^2 B_0^2}{8m} \sum_i \langle \psi | x_i^2 + y_i^2 | \psi \rangle \quad (3.16)$$

³Again, the derivation using a test function.

If we deal with ions with a one electron shell being half filled, or molecules that do not have an axis of symmetry parallel to the magnetic field, we have to include also the second order contribution

$$\sum_{\psi' \neq \psi} \frac{|\langle \psi | L_z + 2S_z | \psi' \rangle|^2}{E_\psi - E_{\psi'}}, \quad (3.17)$$

where the sum goes over all the excited states. This term is the source of so called *Van Vleck paramagnetism* which will not be considered further in these notes.

3.3 Magnetization and thermodynamics

The magnetization \mathbf{M} connects the magnetic flux density \mathbf{B} and the magnetic field strength \mathbf{H} as follows

$$\mathbf{B} = \mu_0(\mathbf{H} + \mathbf{M}), \quad (3.18)$$

where $\mu_0 = 4\pi \cdot 10^{-7}$ H/m is the permeability of the vacuum. The quantity which is the measure of "magnetisability" of the material as a response to an external field is known as the magnetic susceptibility χ , which can be defined as⁴

$$\chi = \frac{\partial M}{\partial H} \quad (3.19)$$

We now take a brief tour into thermodynamics because it is a necessary tool for calculating the magnetic properties. The partition function Z is defined as

$$Z = \sum_n e^{-E_n/k_B T}. \quad (3.20)$$

Since the magnetic energy is given by $E = -\mathbf{m} \cdot \mathbf{B} = -m_z B_0$, we find that

$$\frac{\partial Z}{\partial B_0} = \sum_n -\frac{1}{k_B T} e^{-E_n/k_B T} \frac{\partial E_n}{\partial B_0} = \frac{1}{k_B T} \sum_n m_{z,n} e^{-E_n/k_B T}. \quad (3.21)$$

Thus the magnetization of the system is

$$M = \frac{\langle m_z \rangle}{V} = \frac{1}{V} \frac{k_B T}{Z} \frac{\partial Z}{\partial B_0} \quad (3.22)$$

Since the Helmholtz energy $A = -k_B T \ln Z$, the former can be simplified to

$$M = \frac{\langle m_z \rangle}{V} = -\frac{1}{V} \frac{\partial A}{\partial B_0}. \quad (3.23)$$

Thus the magnetic susceptibility can be obtained from A as follows

$$\chi = \frac{\partial M}{\partial H} = -\frac{1}{V} \frac{\partial^2 A}{\partial H \partial B_0}. \quad (3.24)$$

If χ is small, then $\mathbf{B} \approx \mu_0 \mathbf{H}$ and we get

$$\chi \approx \mu_0 \frac{\partial M}{\partial B} = -\frac{\mu_0}{V} \frac{\partial^2 A}{\partial B_0^2} \quad (3.25)$$

⁴Here χ is a scalar quantity but in general it can be a tensor $\chi_{ij} = \partial M_i / \partial H_j$

3.4 Diamagnetism

Consider an ion which contains only filled shells. For such an ion $L^2|\psi\rangle = S^2|\psi\rangle = 0$ and thus the magnetic energy (3.16) becomes

$$\Delta E_M = \frac{e^2 B_0^2}{8m} \sum_i \langle \psi | x_i^2 + y_i^2 | \psi \rangle. \quad (3.26)$$

In low temperatures, only the ground state of the system is populated: the Helmholtz energy of N ions is simply $A = N\Delta E_M$. Therefore the magnetization of such ions in the volume V is

$$M = -\frac{N}{V} \frac{\partial \Delta E_M}{\partial B_0} = -\frac{N}{V} \frac{e^2 B_0}{4m} \sum_i \langle \psi | x_i^2 + y_i^2 | \psi \rangle, \quad (3.27)$$

and the susceptibility

$$\chi = -\frac{\mu_0 N}{V} \frac{\partial^2 \Delta E_M}{\partial B_0^2} = -\frac{N}{V} \frac{\mu_0 e^2}{4m} \sum_i \langle \psi | x_i^2 + y_i^2 | \psi \rangle. \quad (3.28)$$

Since the shells are filled, we may assume that the ion is spherically symmetric, thus $\langle \psi | x_i^2 | \psi \rangle = \langle \psi | y_i^2 | \psi \rangle = \frac{1}{3} \langle \psi | r_i^2 | \psi \rangle$. Furthermore, we may write $\sum_i \langle \psi | r_i^2 | \psi \rangle = Z_e r^2$, where Z_e is the number of electrons and r is the effective ionic radius. Thus the magnetization and susceptibility become

$$M_{dia} = -\frac{N}{V} \frac{Z_e r^2 e^2}{12m} B_0, \quad \chi_{dia} = -\frac{N}{V} \frac{\mu_0 e^2}{12m} Z_e r^2. \quad (3.29)$$

The magnetization and susceptibility are negative, so it goes against the applied magnetic field. Thus it is fitting to call the macroscopic magnetic response *diamagnetic*⁵. More specifically, the derived term is known as the *Larmor diamagnetism*.

3.5 Paramagnetism part I: Electron spin

Before we consider the general case, we first consider paramagnetism of ions which have a single unpaired electron for which $L = 0$. Therefore the permanent magnetic moment of the ion is completely due to the electron spin. Thus the magnetic energy of the ion given by Eq. (3.16) becomes

$$\Delta E_M = 2\mu_B B_0 \langle \psi | S_z | \psi \rangle, \quad (3.30)$$

where the diamagnetic term is dropped out since it is significantly smaller than the paramagnetic contribution for the magnetic field strengths encountered in all practical situations. Hiding all the unnecessary information about the ion under the carpet, we may describe its quantum states in a simple manner with the eigenstates of S_z : $|\uparrow\rangle$ and $|\downarrow\rangle$. Therefore we are left with a two-level system with the energy levels $E_{\pm} = \pm\mu_B B_0$.⁶ The partition function of the system is thus

$$Z = e^{-E_-/k_B T} + e^{-E_+/k_B T} = e^{\mu_B B_0/k_B T} + e^{-\mu_B B_0/k_B T} = 2 \cosh\left(\frac{\mu_B B_0}{k_B T}\right) \quad (3.31)$$

⁵ *dia* = against, across

⁶ It is worth noting, that the higher energy level corresponds to $|\uparrow\rangle$ state for which the spin is parallel with the magnetic field. This is because the magnetic moment associated with the spin $\boldsymbol{\mu}_S = -2\mu_B \mathbf{S}$ points to the opposite direction.

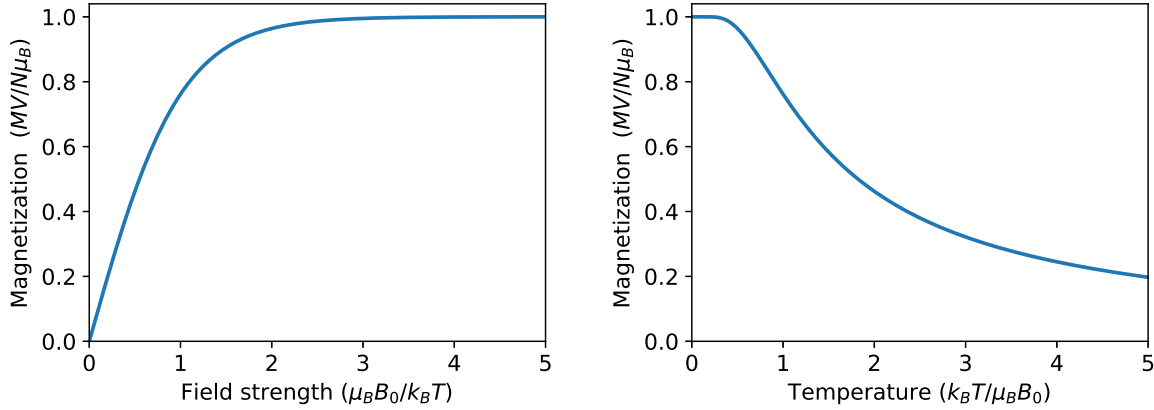


Figure 3.1: Left: Magnetization of non-interacting spin-1/2 ions as a function of field strength; Right: Magnetization as a function of temperature.

Therefore the expectation value of the Helmholtz energy is

$$\langle A \rangle = -k_B T \ln \left[\cosh \left(\frac{\mu_B B_0}{k_B T} \right) \right] - k_B T \ln 2 \quad (3.32)$$

Now consider a system consisting of N such non-interacting spin-1/2 ions in volume V . The Helmholtz energy, being an additive quantity, is now simply N times the expectation value of a single ion⁷, *i.e.*

$$A = N \langle A \rangle = -N k_B T \ln \left[\cosh \left(\frac{\mu_B B_0}{k_B T} \right) \right] - N k_B T \ln 2 \quad (3.33)$$

Thus the magnetization is

$$M = \frac{N k_B T}{V} \frac{\partial}{\partial B_0} \ln \left[\cosh \left(\frac{\mu_B B_0}{k_B T} \right) \right] = \frac{N \mu_B}{V} \tanh \left(\frac{\mu_B B_0}{k_B T} \right) \quad (3.34)$$

The magnetization as functions of field strength and temperature are presented in Figure 3.1. As can be seen, in the absence of external field, there is no net magnetization as the thermal motion randomizes the orientations of the spins. Increasing the field strength, we see that the magnetization behaves linearly in terms of B_0 when the magnetic energy per spin $\mu_B B_0$ is considerably smaller than the thermal energy $k_B T$. Eventually, when the magnetic energy overcomes the thermal energy, we see the saturation of magnetization towards the value $N \mu_B / V$ which corresponds to the situation where all magnetic moments are aligned with the magnetic field.

In the room temperature $T = 300$ K, the magnetic field required to overcome the thermal energy is the order of 500 T, which is an extremely strong field! Therefore we may usually employ the linear approximation for magnetization:

$$M \approx \frac{N \mu_B^2}{V} \frac{B_0}{k_B T} \quad (3.35)$$

For weak magnetic fields, the susceptibility then becomes

$$\chi \approx \mu_0 \frac{\partial M}{\partial B_0} = \frac{N \mu_0 \mu_B^2}{k_B V} \frac{1}{T} \quad (3.36)$$

⁷Strictly speaking, this is also the expectation value of A , but when N is large, the distribution of any thermodynamic quantity is heavily peaked around its expectation value.

The susceptibility of the paramagnetic system consisting of spin-1/2 particles is thus found to be inversely proportional to the temperature under the weak field assumption. The behaviour is known as *Curie's law* and applies for other paramagnetic systems as well. However, the constant of proportionality differs from system to system since most ions are not spin-1/2 particles. We will generalize the formalism in the subsequent sections.

3.6 Total angular momentum J

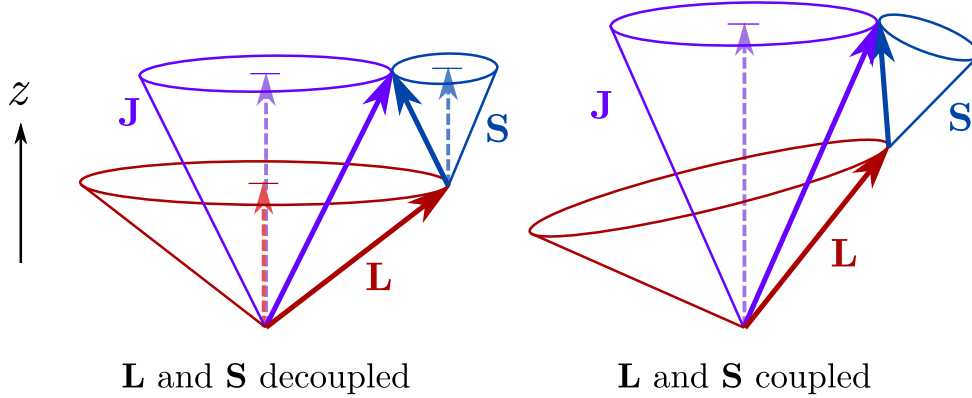


Figure 3.2: The effect of the LS-coupling to the summation of the angular momenta.

Consider a system where orbital and spin angular momenta are decoupled from each other. The states of such a system are consequently the eigenstates of operators \mathbf{L}^2 , \mathbf{S}^2 , L_z , and S_z . Thus the states can be characterised by the quantum numbers L , S , m_L , and m_S or $|LSm_Lm_S\rangle$. However, the situation changes if we add in the spin-orbit interaction. If the interaction strength is weak compared to the Coulomb interaction, we may consider its effect as a small perturbation. Therefore we may still use the L and S as the good quantum numbers. However, the same does not apply to m_L and m_S .

In absence of the spin-orbit coupling, the directions of \mathbf{L} and \mathbf{S} are independent allowing the determination of their z -components simultaneously. On the other hand, introducing the spin-orbit interaction Hamiltonian of the form

$$\hat{H}_{LS} = \Lambda \mathbf{L} \cdot \mathbf{S} \quad (3.37)$$

fixes the angle between \mathbf{L} and \mathbf{S} owing to which the energy eigenstates are no longer eigenstates of L_z and S_z . Stated in more mathematical terms, we may write $\mathbf{L} \cdot \mathbf{S} = L_x S_x + L_y S_y + L_z S_z$ since \mathbf{L} and \mathbf{S} commute. Whereas $|LSm_Lm_S\rangle$ is an eigenstate for L_z and S_z operators, it not so for x - and y -components. Thus the magnetic quantum numbers m_L and m_S are not good quantum numbers for describing the eigenstates of LS-coupled system.

To tackle the problem, we define the total angular momentum operator

$$\mathbf{J} = \mathbf{L} + \mathbf{S}. \quad (3.38)$$

Since \mathbf{L} and \mathbf{S} commute, also \mathbf{J} commutes with them. This means that an eigenstate of \mathbf{L}^2 and \mathbf{S}^2 is also an eigenstate of \mathbf{J}^2 . It can also be shown that such an eigenstate is also an eigenstate

of J_z labelled by m_J . Since $J_z = L_z + S_z$ it is easy to see that the state $|LSm_Lm_S\rangle$ is also an eigenstate of J_z . However, *the converse is not true*. This is illustrated in Figure 3.2.

Now since $\mathbf{J}^2 = \mathbf{L}^2 + 2\mathbf{L} \cdot \mathbf{S} + \mathbf{S}^2$, we find that the interaction Hamiltonian \hat{H}_{LS}

$$\Lambda \mathbf{L} \cdot \mathbf{S} = \frac{\Lambda}{2}(\mathbf{J}^2 - \mathbf{L}^2 - \mathbf{S}^2) \quad (3.39)$$

which is conformant with the earlier argument that for the weak spin-orbit coupling the eigenstates are still the eigenstates of \mathbf{L}^2 and \mathbf{S}^2 . Therefore we arrive at the conclusion that the good quantum numbers describing the LS-coupled system are J , L , S , and m_J .

3.7 Magnetic moment μ_J and Landé g-factor

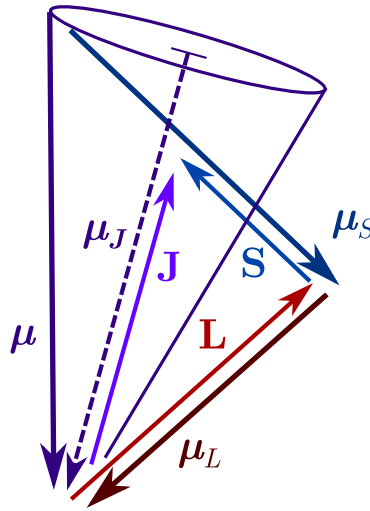


Figure 3.3: Summation of the orbital and spin magnetic moments

What is the magnetic moment associated with the total angular momentum then? As depicted in Figure 3.3, the total magnetic moment is not parallel with \mathbf{J} since the g-factors of the orbital and spin magnetic moments differ. However, since \mathbf{L} and \mathbf{S} precess around \mathbf{J} so does the total magnetic moment. Therefore it has definitive component μ_J (anti)parallel to \mathbf{J} and thus we find that the expectation value of the magnetic momentum μ is

$$\langle JLSm_J | \mu | JLSm_J \rangle = \langle JLSm_J | \mu_J | JLSm_J \rangle \quad (3.40)$$

Therefore it also holds that

$$\langle JLSm_J | \mu \cdot \mathbf{J} | JLSm_J \rangle = \langle JLSm_J | \mu_J \cdot \mathbf{J} | JLSm_J \rangle \quad (3.41)$$

Since $\mu_J \parallel \mathbf{J}$, we may write

$$\mu_J = -g_J \mu_B \mathbf{J}, \quad (3.42)$$

where g_J is the *Landé g-factor* which can be interpreted as the effective number of Bohr magnetons pointing at the direction of \mathbf{J} . Now the right hand side of Eq. (3.41) becomes

$$\langle JLSm_J | \mu_J \cdot \mathbf{J} | JLSm_J \rangle = -g_J \mu_B \langle JLSm_J | \mathbf{J}^2 | JLSm_J \rangle = -g_J \mu_B J(J+1) \quad (3.43)$$

Similarly, since $\boldsymbol{\mu} = -\mu_B(\mathbf{L} + 2\mathbf{S}) = -\mu_B(\mathbf{J} + \mathbf{S})$, the left hand side of Eq. (3.41) becomes

$$\begin{aligned}\langle JLSm_J | \boldsymbol{\mu}_J \cdot \mathbf{J} | JLSm_J \rangle &= -\mu_B \langle JLSm_J | \mathbf{J}^2 + \mathbf{S}^2 + \mathbf{S} \cdot \mathbf{L} | JLSm_J \rangle \\ &= -\mu_B \frac{1}{2} \langle JLSm_J | 3\mathbf{J}^2 + \mathbf{S}^2 - \mathbf{L}^2 | JLSm_J \rangle \\ &= -\mu_B \frac{3J(J+1) + S(S+1) - L(L+1)}{2},\end{aligned}\tag{3.44}$$

where $\mathbf{J}^2 = \mathbf{L}^2 + 2\mathbf{L} \cdot \mathbf{S} + \mathbf{S}^2$ has been used. Thus by combining the Equations (3.43) and (3.43) we find that the Landé g-factor is

$$g_L = \frac{3}{2} + \frac{S(S+1) - L(L+1)}{2J(J+1)}\tag{3.45}$$

3.8 Paramagnetism part II: Full ion

We are now equipped to calculate the magnetization and the susceptibility of considering a full ion.

Bibliography

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