
Stefaan Cottenier

August 7, 2013
Contents

Foreword(s) v
Foreword to the second edition (2013) v
Foreword to the first edition (2002) v

Citing this text vii

Useful references for further reading ix

1 Density Functional Theory as a way to solve the quantum many body problem 1
1.1 Level 1: The Born-Oppenheimer approximation 1
1.2 Level 2: Density Functional Theory 2
1.2.1 The theorems of Hohenberg and Kohn 2
1.2.2 The Kohn-Sham equations 4
1.2.3 The exchange-correlation functional 7
1.3 Level 3: Solving the equations 8

2 The pseudopotential method (in brief) 11

3 The APW method 15

4 The LAPW method 21
4.1 The regular LAPW method 21
4.2 LAPW with Local Orbitals (LAPW+LO) 22

5 The APW+lo method 25
5.1 The ‘pure’ APW+lo basis set 25
5.2 Mixed LAPW/APW+lo basis sets 26
5.3 APW+lo with Local Orbitals (APW+lo+LO) 26

6 The PAW method (in brief) 27

7 Examples for WIEN2k 31
7.1 Linearization energies 31
7.1.1 Initialization 32
7.1.2 Interpreting input 32
7.1.3 Interpreting output 35
7.1.4 The -in1orig and -in1new options 40
Foreword(s)

Foreword to the second edition (2013)

A lot has happened in the 10 years since the first version of this text was published. There is no longer a lack of introductory books that are explicitly meant for users of DFT-codes – some of these works are listed now on page ix. The PAW method has emerged as a mainstream method, next to the former binity of pseudopotential methods and linearized methods – this is reflected in the new Chap. 6 that treats PAW in the same brief way as Chap. 2 treats the pseudopotential method. And finally, the WIEN2k code has evolved to such an extent that much of the content of the chapter with WIEN2k examples/illustrations had become irrelevant (and even confusing for new users). That’s why this chapter (numbered Chap. 7 in this version) has been completely rewritten, using the WIEN2k 13.1 version.

All other chapters and appendices have been left basically unchanged. I would do formulate some parts now a little different compared to 10 years ago. However, it would be much harder now for me to imagine the questions and needs of a beginning user then it was 10 years ago. Leaving those chapters untouched guarantees that this text keeps it spirit as a digestible piece of background info for new users, while the updates described above prevent it from becoming outdated.

Foreword to the first edition (2002)

As you can see in the reference section (page ix), there are more then a few excellent books, reviews and papers on all topics covered in the present text. Why would there be any need for yet another text? It is my feeling that much of the existing material has been written for an audience that was already rather familiar with theoretical and/or computational condensed matter physics. With the increasing availability of more and better general-purpose condensed matter codes – such as WIEN2k – increasingly more people with an explicit experimental background start to use these codes. Especially those new experimental users that have no theoretician nearby, might find the existing reading material either too short and fast or too voluminous and detailed to start with. This text tries to fill that gap: it is meant to give a rather detailed overview of DFT and the family of (L)APW-methods, but starts from the
knowledge that every physicist, chemist and engineer has acquired at an undergraduate level. A large appendix section even brushes up the most important of these undergraduate items, with emphasis on those aspects that will play a role in the main text. From the basic material of the appendices, the reader is taken at the hand, and is step by step introduced into DFT and the (L)APW-methods. How the different elements explained in the text appear in the in- and output of the WIEN2k-code, is demonstrated in the example section.

Chapter 1 on Density Functional Theory can be read independently from the chapters 2 to 5 on (L)APW-based band structure methods. Readers that are particularly interested in one of both main subjects, can without much harm read only the relevant chapters.

This text is by no means a course to teach you how to use the WIEN2k-code from scratch. Such a course can be found in the User’s Guide accompanying the code, and this text should be used parallel to it. The emphasis of the User’s Guide is more on how to work with WIEN2k. This text explains the formalism behind the code, and should make clear why this or that quantity appears in the in- or output of WIEN2k, and what its influence on a calculation can be. This text is meant to help a beginner to acquire a deeper understanding about what she or he is doing with the code.

Much emphasis has been put on mathematical rigour. Therefore you will find many formulae that contain tons of indices. At first sight, this might scare some people. Nevertheless, this crowded notation is meant as a learning tool: try hard to understand the meaning and necessity of every piece of notation. This will prevent you from deviating into misconceptions. As a compensation for this loaded notation, steps between subsequent formulae are kept very small and should be really easy to understand. This text is by no means complete or error-free. If you find an error or inconsistency, please e-mail Stefaan.Cottenier@fys.kuleuven.ac.be or stefaan.cottenier@ugent.be. If you – as a beginner – feel that important topics are not covered, you might suggest items to add. Or if you – as an expert – feel the same, you might even write them ...

Stefaan Cottenier
Instituut voor Kern- en Stralingsfysica (IKS)
Katholieke Universiteit Leuven
Celestijnenlaan 200 D
B-3001 Leuven
Belgium

Stefaan.Cottenier@fys.kuleuven.ac.be
http://www.fys.kuleuven.ac.be/iks/nvsf/nvsf.html

Center for Molecular Modeling (CMM) &
Department of Materials Science and Engineering (DMS)

Ghent University
Technologiepark 903
BE-9052 Zwijnaarde
Belgium

stefaan.cottenier@ugent.be
http://molmod.ugent.be
Citing this text

If you found this text useful, it might occur that you want to cite it in your own work. In that case, please use exactly the following citation:


The fact that people do cite this work, is the only formal argument I can use to justify the time spent to writing it.
Useful references for further reading

For each of the main sections of this text, several good books, reviews or original papers exist in which you can read in-depth about the subject. Several of these sources are used in this text, a few times verbatim.

- Books written for DFT code users:

  *Modeling materials using Density Functional Theory*
  John. R. Kitchin

  A recipe oriented book, “with the intention of giving you enough information and knowledge to start your research”. Relies on VASP and the ASE environment.

  *Density Functional Theory: a practical introduction*
  D. Sholl and J. Steckel
  Wiley (2009)

  “The application of density functional theory (DFT) calculations is rapidly becoming a standard tool for diverse materials modeling problems in physics, chemistry, materials science, and multiple branches of engineering. Although a number of highly detailed books and articles on the theoretical foundations of DFT are available, it remains difficult for a newcomer to these methods to rapidly learn the tools that allow him or her to actually perform calculations that are now routine in the fields listed above. This book aims to fill this gap by guiding the reader through the applications of DFT that might be considered the core of continually growing scientific literature based on these methods.”

  *Atomistic Computer Simulations*
  V. Brázdová and D. R. Bowler
  Wiley (2012)

  “With this book, we are not aiming to provide a detailed guide to the theory underlying the many different atomistic simulation techniques available. Nor are we trying to give a recipe book of how to write or implement particular algorithms. We have written a textbook that aims to provide the reader with the knowledge required to design and execute any project involving atomistic simulations.”
Electronic Structure – basic theory and practical methods
Richard M. Martin

One of the first books that tried to bridge the gap between the formal theory and practical calculations (with most emphasis on the formal aspects, though).

Quantum-Mechanical ab initio calculation of the properties of crystalline materials

An early book that covers DFT and many of the technical features that pop up in actual computational work (k-space, integration methods, computational aspects, examples of computational studies on structural properties, phonons, thermodynamical quantities,...).

- On Density Functional Theory
  A good general review of DFT and its applications:
  Atomic-scale computational materials science

  An extensive text that focusses on technical as well as on conceptual aspects:
  A Bird’s eye view on Density Functional Theory
  Free download at http://arxiv.org/abs/condmat/0211443

  Two books that go really in-depth:
  Density Functional Theory
  Density Functional Theory of Atoms and Molecules

- On the pseudopotential method
  Pseudopotential methods in condensed matter applications

- On the APW method
  Three old but nice and detailed sources:
  Augmented Plane Wave method: a guide to performing electronic structure calculations
  T.L. Loucks, Benjamin, New York (1967)
  L.F. Mattheis, J.H. Wood, A.C. Switendick,
  Methods in Computational Physics 8 (1968) 64
  The calculation of electronic energy bands by the augmented plane wave method
  J.O. Dimmock in Solid State Physics, H. Ehrenreich, F. Seitz, D. Turnbull (editors),
• On the LAPW method
  The basic LAPW-book that covers in detail many of the technical aspects:

  *Planewaves, pseudopotentials and the LAPW-method (2nd edition)*,  

  The original LAPW papers:

  *Linear methods in band theory*  

  *Use of energy derivative of the radial solution in an augmented plane wave method: application to copper*  

• The APW+lo method
  The original papers, that cover also (L)APW:

  *An alternative way of linearizing the augmented plane-wave method*  

  *Efficient linearization of the augmented plane-wave method*  
  G.K.H. Madsen, P. Blaha, K. Schwarz, E. Sjöstedt, L. Nordström,  
  *Physical Review B* 64 (2001) 195134

• **wien2k-specific reviews**

  *Electronic structure calculations of solids using the wienn2k package for material science*  
  K. Schwarz, P. Blaha, G.K.H. Madsen,  
  *Computer Physics Communications* 147 (2002) 71

  *DFT calculations of solids with LAPW and wienn2k*  

  *Electronic structure of solids with wienn2k*  
Chapter 1

Density Functional Theory as a way to solve the quantum many body problem

A solid is a collection of heavy, positively charged particles (nuclei) and lighter, negatively charged particles (electrons). If we have N nuclei, we are dealing with a problem of N+ZN electromagnetically interacting particles. This is a many-body problem, and because the particles are so light, quantum mechanics is needed: a quantum many body problem. The exact many-particle hamiltonian for this system is:

\[
\hat{H} = -\frac{\hbar^2}{2} \sum_i \nabla^2 \frac{\hat{R}_i}{M_i} - \frac{\hbar^2}{2} \sum_i \nabla^2 \frac{\vec{r}_i}{m_e} - \frac{1}{4\pi\epsilon_0} \sum_{i,j} \frac{e^2 Z_i}{|\hat{R}_i - \vec{r}_j|} + \frac{1}{8\pi\epsilon_0} \sum_{i\neq j} \frac{e^2}{|\vec{r}_i - \vec{r}_j|} + \frac{1}{8\pi\epsilon_0} \sum_{i\neq j} \frac{e^2 Z_i Z_j}{|\hat{R}_i - \hat{R}_j|}
\] (1.1)

The mass of the nucleus at \(\hat{R}_i\) is \(M_i\), the electrons have mass \(m_e\) and are at \(\vec{r}_i\). The first term is the kinetic energy operator for the nuclei, the second for the electrons. The last three terms describe the Coulomb interaction between electrons and nuclei, between electrons and other electrons, and between nuclei and other nuclei. It is out of question to solve this problem exactly. In order to find acceptable approximate eigenstates, we will need to make approximations at 3 different levels.

1.1 Level 1: The Born-Oppenheimer approximation

The nuclei are much heavier and therefore much slower than the electrons. We can hence ‘freeze’ them at fixed positions and assume the electrons to be in instantaneous equilibrium with them. In other words: only the electrons are kept as players in our many body problem. The nuclei are deprived from this status, and reduced to a given source of positive charge, they become ‘external’ to the electron cloud. After having applied this approximation, we are left with a collection of NZ interacting negative particles, moving in the (now external or given) potential of the nuclei.

What are the consequences of the Born-Oppenheimer approximation on the hamiltonian 1.1? The nuclei do not move any more, their kinetic energy is zero and the first term disappears. The last term reduces to a constant. We are left with the kinetic energy of the electron gas, the
potential energy due to electron-electron interactions and the potential energy of the electrons in the (now external) potential of the nuclei. We write this formally as:

$$\hat{H} = \hat{T} + \hat{V} + \hat{V}_{ext}$$  \hfill (1.2)

It is interesting to note here that the kinetic and electron-electron terms of 1.2 depend only on the fact that we are dealing with a many-electron system (and not with a many-proton system for instance, where the strong nuclear force would play a role). They are independent of the particular kind of many-electron system (Br$_2$ or H$_2$O molecule, Cu or Fe, bcc-Fe or fcc-Fe, ...). This part is universal. System-specific information (which nuclei, and on which positions) is given entirely by $\hat{V}_{ext}$.

### 1.2 Level 2: Density Functional Theory

The quantum many body problem obtained after the first level approximation (Born-Oppenheimer) is much simpler than the original one, but still far too difficult to solve. Several methods exist to reduce equation 1.2 to an approximate but tractable form. A historically very important one is the Hartree-Fock method (HF), described in many condensed matter textbooks. It performs very well for atoms and molecules, and is therefore used a lot in quantum chemistry. For solids it is less accurate, however. We will not treat HF, but explain a more modern and probably also more powerful method: Density Functional Theory (DFT)

Although its history goes back to the early thirties of the 20th century, DFT has been formally established in 1964 by two theorems due to Hohenberg and Kohn$^2$.

#### 1.2.1 The theorems of Hohenberg and Kohn

The traditional formulation of the two theorems of Hohenberg and Kohn is as follows:

First theorem: There is a one-to-one correspondence between the ground-state density $\rho(\vec{r})$ of a many-electron system (atom, molecule, solid) and the external potential $V_{ext}$. An immediate consequence is that the ground-state expectation value of any observable $\hat{O}$ is a unique functional of the exact ground-state electron density:

$$\langle \Psi | \hat{O} | \Psi \rangle = O[\rho]$$  \hfill (1.3)

Second theorem: For $\hat{O}$ being the hamiltonian $\hat{H}$, the ground-state total energy functional $H[\rho] \equiv E_{V_{ext}}[\rho]$ is of the form

$$E_{V_{ext}}[\rho] = \left< \Psi | \hat{T} + \hat{V} | \Psi \right> + \left< \Psi | \hat{V}_{ext} | \Psi \right>$$  \hfill (1.4)

$$= F_{HK}[\rho] + \int \rho(\vec{r}) V_{ext}(\vec{r}) d\vec{r}$$  \hfill (1.5)

---

$^1$As Hartree-Fock, DFT is a general method to solve the quantum many body problem. It can be applied not only to the electron gas as we use here, but also to the proton-neutron gas in order to construct nuclear models, or to the nucleus-electron gas without the Born-Oppenheimer approximation to describe solids with light elements.

where the Hohenberg-Kohn density functional $F_{HK}[\rho]$ is universal for any many-electron system. $E_{V_{ext}}[\rho]$ reaches its minimal value (equal to the ground-state total energy) for the ground-state density corresponding to $V_{ext}$.

We do not prove these theorems here, but ponder a few implications of the three keywords invertibility (one-to-one correspondence $\rho \leftrightarrow V_{ext}$), universality and variational access (minimal value).

First, the one-to-one correspondence between ground-state density and external potential is intriguing. It is obvious that a given many-electron system has a unique external potential, which by the hamiltonian 1.2 and the Schrödinger equation yields a unique ground-state many particle wave function. From this wave function, the corresponding electron density is easily found. An external potential hence leads in a well-defined way to a unique ground-state density corresponding to it. But intuitively it looks like the density contains less information than the wave function. If this would be true, it would not be possible to find a unique external potential if only a ground-state density is given. The first theorem of Hohenberg and Kohn tells exactly that this is possible! The density contains as much information as the wave function does (i.e. everything you could possibly know about an atom, molecule or solid). All observable quantities can be retrieved therefore in a unique way from the density only, i.e. they can be written as functionals of the density.

Second, the universality of $F_{HK}[\rho]$. Equation 1.5 is easily written down by using the density operator $3$, and supposing the ground-state density is known, the contribution to the total energy from the external potential can be exactly calculated. An explicit expression for the Hohenberg-Kohn functional $F_{HK}$ is not known. But anyway, because $F_{HK}$ does not contain information on the nuclei and their position, it is a universal functional for any many-electron system. This means that in principle an expression for $F_{HK}[\rho]$ exists which can be used for every atom, molecule or solid which can be imagined.

Third, the second theorem makes it possible to use the variational principle of Rayleigh-Ritz in order to find the ground-state density. Out of the infinite number of possible densities, the one which minimizes $E_{V_{ext}}[\rho]$ is the ground-state density corresponding to the external potential $V_{ext}(\vec{r})$. Of course, this can be done only if (an approximation to) $F_{HK}[\rho]$ is known. But having found $\rho$, all knowledge about the system is within reach...

It is useful to stress the meaning of the energy functional $E_{V_{ext}}[\rho]$ once more. When it is

---

3The density operator $\hat{\rho}(\vec{r})$ for an $N$-particle system is defined as

$$\hat{\rho}(\vec{r}) = \sum_{i=1}^{N} \delta(\vec{r}_i - \vec{r})$$ (1.6)

such that its evaluation for a many body wave function $\Psi$ yields the density:

$$\rho(\vec{r}) = \langle \Psi(\vec{r}_1, \vec{r}_2, ..., \vec{r}_N) | \hat{\rho}(\vec{r}) | \Psi(\vec{r}_1, \vec{r}_2, ..., \vec{r}_N) \rangle$$ (1.7)

$$= \langle \Psi(\vec{r}_1, \vec{r}_2, ..., \vec{r}_N) | \sum_{i=1}^{N} \delta(\vec{r}_i - \vec{r}) | \Psi(\vec{r}_1, \vec{r}_2, ..., \vec{r}_N) \rangle$$ (1.8)

$$= \sum_{i=1}^{N} \int \Psi^*(\vec{r}_1, \vec{r}_2, ..., \vec{r}_i \equiv \vec{r}, ..., \vec{r}_N) \Psi(\vec{r}_1, \vec{r}_2, ..., \vec{r}_i \equiv \vec{r}, ..., \vec{r}_N) \, d\vec{r}_1 \, d\vec{r}_2 \, ... \, d\vec{r}_i \, ... \, d\vec{r}_N$$ (1.9)
evaluated for the density \( \rho \) corresponding to the particular \( V_{\text{ext}} \) for this solid, it gives the ground state energy. When it is evaluated for any other density however, the resulting number has no physical meaning\(^4\).

### 1.2.2 The Kohn-Sham equations

The equations of Kohn and Sham, published in 1965, turn DFT into a practical tool\(^5\). They are a practical procedure to obtain the ground state density. Let us first rewrite the Hohenberg-Kohn functional. The *correlation energy* is defined as this part of the total energy which is present in the exact solution, but absent in the Hartree-Fock solution. The total energy functionals \( E_e[\rho] \) and \( E_{HF}[\rho] \) corresponding to the exact and Hartree-Fock Hamiltonians respectively, are\(^6\):

\[
E_e = T + V \\
E_{HF} = T_0 + (V_H + V_x)
\]

Here \( T \) and \( V \) are the exact kinetic and electron-electron potential energy functionals, \( T_0 \) is the functional for the kinetic energy of a non-interacting electron gas, \( V_H \) stands for the Hartree contribution and \( V_x \) for the exchange contribution. By subtracting 1.11 from 1.10, the functional for the correlation contribution appears to be:

\[
V_c = T - T_0
\]

The exchange contribution to the total energy is defined as the part which is present in the Hartree-Fock solution, but absent in the Hartree solution. Obviously, with the Hartree functional given by

\[
E_H = T_0 + V_H
\]

it can be defined as

\[
V_x = V - V_H
\]

With this knowledge, we can rewrite the Hohenberg-Kohn functional in the following way:

\[
F_{HK} = T + V + T_0 - T_0 \\
= T_0 + V + (T - T_0) \\
= T_0 + V + V_e + V_H - V_H \\
= T_0 + V_H + V_c + (V - V_H) \\
= T_0 + V_H + (V_x + V_c)
\]

\(^4\)This is too crude, but the details are subtle: any density \( \rho_i(\vec{r}) \), other than the ground-state density, which extremalizes \( E_{V_{\text{ext}}}[\rho] \) is an excited-state density, with \( E_i = E_{V_{\text{ext}}}[\rho = \rho_i] \) the corresponding energy. The inverse statement is not true: not all excited-state densities extremalize \( E_{V_{\text{ext}}}[\rho] \). If \( \rho_j(\vec{r}) \) is such a density, \( E_j = E_{V_{\text{ext}}}[\rho = \rho_j] \) is a lower bound for the energy of the excited state.


\(^6\)For easier notation, we drop for a while the square brackets.
Here $V_{xc}$ is the exchange-correlation energy functional. We don’t know it formally, as it contains the difficult exchange and correlation contributions only. If we assume for a while that we do know $V_{xc}$, we can write explicitly the energy functional:

$$E_{\text{ext}}[\rho] = T_0[\rho] + V_H[\rho] + V_{xc}[\rho] + V_{\text{ext}}[\rho]$$

One could use now the second Hohenberg-Kohn theorem to find the ground state density, but then we would have won nothing by our transformation. Instead, you can interpret the above expression also as the energy functional of a non-interacting classical electron gas, subject to two external potentials: one due to the nuclei, and one due to exchange and correlation effects. The corresponding Hamiltonian – called the Kohn-Sham Hamiltonian – is

$$\hat{H}_{KS} = \hat{T}_0 + \hat{V}_H + \hat{V}_{xc} + \hat{V}_{\text{ext}}$$

where the exchange-correlation potential is given by the functional derivative

$$\hat{V}_{xc} = \frac{\delta V_{xc}[\rho]}{\delta \rho}$$

The theorem of Kohn and Sham can now be formulated as follows:

The exact ground-state density $\rho(\vec{r})$ of an N-electron system is

$$\rho(\vec{r}) = \sum_{i=1}^{N} \phi_i(\vec{r})^* \phi_i(\vec{r})$$

where the single-particle wave functions $\phi_i(\vec{r})$ are the N lowest-energy solutions of the Kohn-Sham equation

$$\hat{H}_{KS} \phi_i = \epsilon_i \phi_i$$

And now we did won a lot. To find the ground-state density, we don’t need to use the second Hohenberg-Kohn theorem any more, but we can rely on solving familiar Schrödinger-like non-interacting single-particle equations. The alternative of using the regular Schrödinger equation, would have led to a far more difficult system of coupled differential equations, because of the electron-electron interaction.

Be aware that the single-particle wave functions $\phi_i$ are not the wave functions of electrons! They describe mathematical quasi-particles, without a direct physical meaning. Only the overall density of these quasi-particles is guaranteed to be equal to the true electron density. Also the single-particle energies $\epsilon_i$ are not single-electron energies.

Both the Hartree operator $V_H$ and the exchange-correlation operator $V_{xc}$ depend on the density $\rho(\vec{r})$, which in turn depends on the $\phi_i$, which are being searched. This means we are dealing with a self-consistency problem: the solutions ($\phi_i$) determine the original equation ($V_H$ and $V_{xc}$ in $H_{KS}$), and the equation cannot be written down and solved before its solution is
Figure 1.1: Flow chart for the $n^{th}$ iteration in the selfconsistent procedure to solve Hartree-Fock or Kohn-Sham equations.
known. An iterative procedure is needed to escape from this paradox (see fig. 1.1). Some starting density \( \rho_0 \) is guessed, and a hamiltonian \( \hat{H}_{KS1} \) is constructed with it. The eigenvalue problem is solved, and results in a set of \( \phi_1 \) from which a density \( \rho_1 \) can be derived. Most probably \( \rho_0 \) will differ from \( \rho_1 \). Now \( \rho_1 \) is used to construct \( \hat{H}_{KS2} \), which will yield a \( \rho_2 \), etc. The procedure can be set up in such a way that this series will converge to a density \( \rho_f \) which generates a \( \hat{H}_{KSf} \) which yields as solution again \( \rho_f \): this final density is then consistent with the hamiltonian.

### 1.2.3 The exchange-correlation functional

The Kohn-Sham scheme described above was exact: apart from the preceding Born-Oppenheimer approximation, no other approximations were made. But we neglected so far the fact that we do not know the exchange-correlation functional. It is here that approximations enter our theory.

A widely used approximation – called the Local Density Approximation (LDA) – is to postulate that the exchange-correlation functional has the following form:

\[
E_{xc}^{LDA} = \int \rho(\vec{r}) \epsilon_{xc}(\rho(\vec{r})) \, d\vec{r} \tag{1.21}
\]

The function (not: functional) \( \epsilon_{xc}(\rho) \) for the homogeneous electron gas has been defined in Appendix E, and is numerically known. This postulate is somehow reasonable: it means that the exchange-correlation energy due to a particular density \( \rho(\vec{r}) \) could be found by dividing the material in infinitesimally small volumes with a constant density. Each such volume contributes to the total exchange correlation energy by an amount equal to the exchange correlation energy of an identical volume filled with a homogeneous electron gas, that has the same overall density as the original material has in this volume (see Fig. 1.2). No law of nature guarantees that the true \( E_{xc} \) is of this form, it is only a reasonable guess. By construction, LDA is expected to perform well for systems with a slowly varying density. But rather surprisingly, it appears to be very accurate in many other (realistic) cases too.

A next logical step to improve on LDA, is to make the exchange-correlation contribution of every infinitesimal volume not only dependent on the local density in that volume, but also on the density in the neighbouring volumes. In other words, the gradient of the density will play a role. This approximation is therefore called the Generalized\(^7\) Gradient Approximation (GGA). Although GGA performs in general slightly better than LDA, there are a few drawbacks. There is only one LDA exchange-correlation functional, because there is a unique definition for \( \epsilon_{xc} \). But there is some freedom to incorporate the density gradient, and therefore several versions of GGA exist (first drawback). Moreover, in practice one often fits a candidate GGA-functional with (hopefully only a few) free parameters to a large set of experimental data on atoms and molecules. The best values for these parameters are fixed then, and the functional is ready to be used routinely in solids. Therefore such a GGA-calculation is strictly spoken not an \textit{ab initio} calculation, as some experimental information is used (second drawback). Nevertheless, there exist GGA’s that are parameter free.

\(^7\)“Generalized”, because a straightforward gradient would result in a functional which violates some relationships which can be proven to be exact for the true functional (and for LDA)
Figure 1.2: Illustration of the idea behind the LDA postulate. Every infinitesimally small volume of the material contributes to the exchange-correlation energy with an amount equal to the contribution of a homogeneous electron gas that occupies that same infinitesimally small volume, and that has the same (overall) charge density as the charge density of the original material in that volume. The horizontal axis is proportional to the density of the homogeneous electron gas. The vertical axis displays the exchange-correlation energy of the homogeneous electron gas. (picture taken from Martin Fuchs, lecture notes, ICTP Workshop August 1999)

1.3 Level 3: Solving the equations

Irrespective whether one has used HF or DFT as level 2 approximation, one ends up with an infinite set of one-electron equations of the following type (m is an integer number that counts the members of the set):

\[
\left( -\frac{\hbar^2}{2m_e} \nabla_m^2 + \frac{e^2}{4\pi \epsilon_0} \int \frac{\rho(\vec{r}')}{|\vec{r} - \vec{r}'|} \, d\vec{r}' + V_\alpha + V_{ext} \right) \phi_m(\vec{r}) = \epsilon_m \phi_m(\vec{r}) \quad (1.22)
\]

We call \( \hat{H}_{sp} \) the single-particle Hamiltonian. For HF, \( V_\alpha \) is the exchange operator. The \( \phi_m \) are true one-electron (or single-particle) orbitals for HF. Exchange is treated exactly, but correlation effects are not included at all. They can be added only in elaborations on the HF-method.

For DFT, \( V_\alpha \) is the exchange-correlation operator, in the L(S)DA, GGA or another approximation. Exchange and correlation are both treated, but both approximately. The \( \phi_m \) are
mathematical single-particle orbitals.

The similarity between the Hartree-Fock and Kohn-Sham equations means that the same mathematical techniques can be used to solve them. ‘Solving’ in most methods means that we want to find the coefficients $c^m_p$ needed to express $\phi_m$ in a given basis set $\phi^b_p$:

$$\phi_m = \sum_{p=1}^{P} c^m_p \phi^b_p \quad (1.23)$$

The wave functions $\phi_m$ belong to a function space which has an infinite dimension, $P$ is therefore in principle infinite. In practice one works with a limited set of basis functions. Such a limited basis will never be able to describe $\phi_m$ exactly, but one could try to find a basis that can generate a function that is ‘close’ to $\phi_m$.

Having chosen a basis (and hence a finite value for $P$) we realize that we can tackle the equations 1.22 as an eigenvalue problem, in the same way as described in Appendix C. For a given $m$, substitute equation 1.23 in 1.22, and left-multiply with $\langle \phi^b_i | (i = 1, \ldots, P)$. As in equation C.8, this leads to

$$\begin{bmatrix}
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & \ldots & \ldots & \ldots & \ldots \\
\vdots & \vdots & \vdots & \vdots & \vdots
\end{bmatrix}
\begin{bmatrix}
c^m_1 \\
c^m_2 \\
\vdots \\
c^m_P
\end{bmatrix}
= 
\begin{bmatrix}
0 \\
\vdots \\
0
\end{bmatrix} \quad (1.24)$$

We recognize here the matrix elements of the single-particle hamiltonian in the basis states, and the overlap matrix elements $S_{ij}$. Remember that the overlap matrix is a unit matrix if the basis set is orthonormal. Diagonalization of the hamiltonian matrix will lead to $P$ eigenvalues and $P$ sets of coefficients that express each of the $P$ eigenfunctions in the given basis (if more eigenfunctions are needed, $P$ must be increased). The larger $P$, the better the approximation of the eigenfunction, but the more time-consuming the diagonalization of the matrix in equation 1.24.

In the following chapters, we will elaborate on this quick sketch. We will see that the label $m$ stands for the quantum numbers $(n, \vec{k})$. It will be possible and advantageous to use a different basis set for each $\vec{k}$, which will enormously reduce the size of $P$, but will lead to as many different eigenvalue problems to solve as there are $\vec{k}$-vectors needed to be known.

What is a good basis set? If the functions of the basis set are very similar to $\phi_m$, one needs only a few of them to accurately describe the wave function, and hence $P$ and the matrix size are small. Such a basis set is called efficient. However, this assumes that you know the solution of your problem almost before you start solving it. Such a basis set can therefore never be very general: for some specific problems it will very quickly yield the solution, but for the majority of cases it will poorly describe the eigenfunctions. In the latter cases the required $P$ is much higher than what is affordable, and limiting $P$ would lead to approximate eigenfunctions that are not acceptable. These approximations carry too much properties from the basis function, and such a basis set is therefore called biased. The art of theoretical condensed matter physics is to find a basis set that is a simultaneously efficient and unbiased. In the following chapters, two families of basis sets will be described – plane waves and augmented plane waves – that each in their own way try to realize this compromise.
Chapter 2

The pseudopotential method (in brief)

At the end of section 1.3, we formulated two principal requirements for a basis set in which we want to expand the eigenstates of the solid state hamiltonian: the basis set should be *unbiased* (=it should not in a hidden way force the solution into a built-in direction) and *efficient* (= $P$ in equation 1.23 should be as low as possible). Furthermore, it would be nice if the basis functions are mathematically simple. This makes both theory development and programming work easier. A basis set that is certainly unbiased and simple is the plane wave basis set that was introduced in section A.4. We saw in equation A.22 that any eigenfunction $\psi_n(\vec{r})$ of a periodic hamiltonian can be expressed exactly in this basis set by means of an infinite set of coefficients $c_{n, \vec{k}}$:

$$\psi_n(\vec{r}) = \sum_{\vec{K}} c_{n, \vec{k} + \vec{K}} e^{i(\vec{k} + \vec{K}) \cdot \vec{r}}$$  \hspace{1cm} (2.1)$$

This has to be compared with the general formulation in equation 1.23, where $m$ stands for $(n, \vec{k})$ and $p$ for $\vec{k} + \vec{K}$. One basis function for $\psi_n(\vec{r})$ (or $\vec{K}$) is therefore:

$$\phi_{\vec{K}}(\vec{r}) = |\vec{K}\rangle = e^{i(\vec{k} + \vec{K}) \cdot \vec{r}}$$  \hspace{1cm} (2.2)$$

Note that this basis-set is $\vec{k}$-dependent: all eigenstates $\psi_n(\vec{r})$ that have the same $\vec{k}$ but a different $n$ will be expressed in the basis set with this particular value of $\vec{k}$. For eigenstates with another $\vec{k}$, a new basis set using that other $\vec{k}$ has to be used.

In practice we cannot work with an infinite basis set, and will have to limit it somehow (the $P$ in equation 1.23). For plane waves, this can be easily done by limiting the set to all $\vec{K}$ with $K \leq K_{max}$. This corresponds to a sphere with radius $K_{max}$ centered at the origin of reciprocal space. All reciprocal lattice vectors that are inside this sphere are taken into the basis set. Instead of $K_{max}$, often the free electron energy corresponding to $K_{max}$ is specified, and is called the *cut-off energy*:

$$E_{cut} = \frac{\hbar^2 K_{max}^2}{2m_e}$$  \hspace{1cm} (2.3)$$

Plane waves are orthogonal:

$$\langle \vec{K}_1 \mid \vec{K}_2 \rangle = \int e^{i(\vec{K}_2 - \vec{K}_1) \cdot \vec{r}} d^3 \vec{r}$$  \hspace{1cm} (2.4)$$

$$= \delta(\vec{K}_2 - \vec{K}_1)$$  \hspace{1cm} (2.5)$$
and therefore the overlap matrix $S$ in equations C.18 or 1.24 is a unit matrix. The matrix elements of the effective single-particle hamiltonian (Kohn-Sham hamiltonian) are calculated in the plane wave basis, and the resulting matrix is diagonalized (see section C.3 for a description of the full procedure). To each eigenvalue $\epsilon^{n,\vec{k}}$ an eigenvector $[c_{\vec{K}}^n]_{P\times 1}$ of $P$ values for $c_{\vec{K}}^n$ corresponds. Actually, an infinite number of eigenvectors, as multiplying with any real number gives again an eigenvector (see Appendix C). This degree of freedom is used to take that $[c_{\vec{K}}^n]_{P\times 1}$ for which – by equation 2.1 – $\phi_{\vec{k}}^n$ is a normalized function.

So far, so good. But: we were looking for the eigenvalue $\epsilon^{n,\vec{k}}$ of $\psi_{\vec{k}}^n$, and instead of one eigenvalue we found $P$ of them, each with their own set of coefficients and hence each leading to another eigenfunction $\psi_{\vec{k}}^n$. What has happened? Each eigenfunction is uniquely labeled by $n$ and $\vec{k}$. Eq. 2.1 applies for all our solutions with the same $\vec{k}$. We are therefore forced to conclude that for all these eigenfunctions, $n$ – which does not explicitly appears in the equations – must be different. In one go, we have found $P$ different eigenfunctions with the same $\vec{k}$ but a different band index!

This process has to be repeated for as many $\vec{k}$-points that are needed for a densely enough sampling of the first Brillouin zone.

We did not yet discuss the efficiency of a plane wave basis set. Will we need many or only a few plane waves? The number of plane waves is determined by the smallest length scales that are to be described in real space. Consider the radial part of a 3s wave function in Ca (fig. 2-a). Near the nucleus, the wave function shows steep behaviour. In order to describe the sharp part between 0 and the minimum at 0.1 Å plane waves with a period as small as roughly an order of magnitude less than this distance are needed (0.01 Å or $10^{-12}$ m). This sets the order of magnitude for $K_{max}$ to $K_{max} = 2\pi/(10^{-12} \text{ m}) = 6.3 \cdot 10^{12} \text{ m}^{-1}$. Calculate the volume of a sphere with this radius ($10^{39} \text{ m}^{-3}$), divide it by the volume of the first Brillouin zone ($2^3 \cdot (\vec{a}^* \times \vec{b}^*)$, about $9.2 \cdot 10^{30} \text{ m}^{-3}$ for a cubic lattice with a lattice constant of 3 Å), and because the first Brillouin zone contains $1 \vec{K}$ you find that about $10^8$ plane waves are needed. That would require diagonalization of (many) $(10^8 \times 10^8)$-matrices, which is way beyond the capability of even supercomputers.
Should we conclude now that a plane wave basis set cannot be used? The most oscillating part of the wave functions are the tails that reach out into the region close to the nucleus. But this region of the solid is quite shielded from the more outer regions of the atoms where chemistry happens, and the electrons will not behave very differently from free atom electrons here. One can therefore replace the potential in these inner regions by a pseudopotential, that is designed to yield very smooth tails of wavefunctions inside the atom (hence only a few plane waves are needed). More to the outer regions of the atoms, the pseudopotential continuously evolves into the true potential, such that this region of the crystal behaves as if nothing happened. In this way, it is possible to use an (ultrasoft) plane wave basis set for realistic cases with a cut-off energy $E_{cut} = 20 \text{ Ry} = 272 \text{ eV}$. This corresponds to $K_{max} \approx 4.5 \text{ au}^{-1}$, which for our example of a cubic lattice with lattice constant 3 Å calls for about 270 plane waves\footnote{Only accidentally similar to $E_{cut}$ in eV.}. A basis set of this size is manageable.

There is no unique recipe to construct a pseudopotential for a particular element. An infinite amount of choices are possible. Two criteria to judge whether a particular pseudopotential is good, are softness and transferability. A pseudopotential is called soft when few plane waves are needed. A specific class of pseudopotentials is even called ultrasoft, because of the very small amount of plane waves it calls for. Generally, making a pseudopotential soft means that it gets tailored for an element in a specific environment. What one really wants is a potential that can be used in whatever environment (molecule, cluster, solid, surface, insulator, metal,...) the corresponding element can be. Such a pseudopotential is called transferable. The art of creating good pseudopotentials is to find potentials that are both (ultra)soft and transferable.
Chapter 3

The APW method

Although the pseudopotential method is extremely useful, there are reasons why alternatives could be attractive. Is the introduction of the pseudopotential completely innocent? What do you do if you are interested in information that is inherently contained in the region near the nucleus (hyperfine fields for instance, or core level excitations)? Can the basis set be made more efficient? Therefore, we will search for a basis set that uses other functions than plane waves, and that does not require the introduction of a pseudopotential. Such a basis set will have to be more efficient, but of course we do not want it to be biased. Our first example of this will be the Augmented Plane Wave (APW) basis set. Right from the beginning it has to be said that the APW-method itself is of no practical use any more today. But for didactical reasons it is advantageous to discuss APW first, before going to its successors, LAPW and APW+lo.

The ideas that lead to the APW basis set are very similar to what made us to introduce the pseudopotential. In the region far away from the nuclei, the electrons are more or less ‘free’. Free electrons are described by plane waves\(^1\). Close to the nuclei, the electrons behave quite as they were in a free atom, and they could be described more efficiently by atomic like functions. Space is therefore divided now in two regions: around each atom\(^2\) a sphere with radius \(R_\alpha\) is drawn (call it \(S_\alpha\)). Such a sphere is often called a muffin tin sphere, the part of space occupied by the spheres is the muffin tin region. The remaining space outside the spheres is called the interstitial region (call it \(I\)). One augmented plane wave (APW) used in the expansion of \(\psi^\alpha_n\) is defined as:

\[
\phi^\alpha_{\vec{k}\vec{K}}(\vec{r}, E) = \begin{cases} 
\frac{1}{\sqrt{V}} e^{i(\vec{k} + \vec{K}) \cdot \vec{r}} & \vec{r} \in I \\
\sum_{\ell, m} A^{n, \vec{k} + \vec{K}}_{\ell m} u^\alpha_\ell (r') Y^\ell_m (\hat{r}') & \vec{r} \in S_\alpha
\end{cases}
\]  

(3.1)

The symbols \(\vec{k}, \vec{K}\) and \(\vec{r}\) keep their usual meaning, \(V\) is the volume of the unit cell. Note that the APW basis set is \(\vec{k}\)-dependent, as was the plane wave basis set. The position inside the spheres is given with respect to the center of each sphere by \(\vec{r}' = \vec{r} - \vec{r}_\alpha\) (see fig. 3). The length of \(\vec{r}'\) is \(r'\), and the angles \(\theta'\) and \(\phi'\) specifying the direction of \(\vec{r}'\) in spherical coordinates, are indicated as \(\hat{r}'\). The \(Y^\ell_m\) are spherical harmonics. The \(A^{n, \vec{k} + \vec{K}}_{\ell m}\) are yet undetermined parameters, as is \(E\). The latter has the dimension of energy. The \(u^\alpha_\ell\) are solutions to the radial part of

\(^1\)Plane waves are eigenfunctions of a hamiltonian with zero potential.

\(^2\)Mark the different atoms in the unit cell by a label \(\alpha\). This label is different for all atoms in the unit cell, not just for all inequivalent atoms.
the Schrödinger equation for a free atom $\alpha$, and this at the energy $E$. For a true free atom, the boundary condition that $u^{\alpha}_{\ell}(r, E)$ should vanish for $r \rightarrow \infty$, limits the number of energies $E$ for which a solution $u^{\alpha}_{\ell}$ can be found. But as this boundary condition does not apply here, we can find a numerical solution for any $E$. Hence, the $u^{\alpha}_{\ell}$ themselves do not correspond to something physical, but that doesn’t harm: they are only part of a basis function, not of the searched eigenfunction itself. And because they are close to how the actual eigenfunction will look like in that region of the crystal, they will do their job as basis function very efficiently.

If an eigenfunction would be discontinuous, its kinetic energy would not be well-defined. Such a situation can therefore never happen, and we have to require that the plane wave outside the sphere matches the function inside the sphere over the complete surface of the sphere (in value, not in slope). That seems a weird thing to do: a plane wave is oscillating and has a unique direction built in, how can it match another function based on spherical harmonics over the entire surface of a sphere? To see how this is possible, we expand the plane wave in spherical harmonics about the origin of the sphere of atom $\alpha$:

$$\frac{1}{\sqrt{V}}e^{i(\vec{k}+\vec{K}) \cdot \vec{r}} = \frac{4\pi}{\sqrt{V}}e^{i(\vec{k}+\vec{K}) \cdot \vec{r}_\alpha} \sum_{\ell, m} j_\ell \left( |\vec{k} + \vec{K}| |\vec{r}'| \right) Y^*_{m}(\hat{\vec{k}} + \hat{\vec{K}}) Y^\ell_{m}(\hat{\vec{r}}')$$  \hspace{1cm} (3.2)

$j_\ell(x)$ is the Bessel function of order $\ell$. Requiring this at the sphere boundary (where $\vec{r}' = \vec{R}_\alpha$, which defines $\vec{R}_\alpha$) to be equal to the $\ell m$-part of Eq. 3.1 easily yields:

$$A_{\ell m}^{\alpha, \vec{k}+\vec{K}} = \frac{4\pi i \ell e^{i(\vec{k}+\vec{K}) \cdot r_\alpha}}{\sqrt{V} u^\alpha_\ell(\vec{R}_\alpha, E)} j_\ell \left( |\vec{k} + \vec{K}| R_\alpha \right) Y^*_{m}(\vec{k} + \vec{K})$$  \hspace{1cm} (3.3)

This uniquely defines the $A_{\ell m}^{\alpha, \vec{k}+\vec{K}}$, apart from the still undetermined $E$. In principle there are an infinite number of terms in Eq. 3.2, which would force us to use an infinite number of $A_{\ell m}^{\alpha, \vec{k}+\vec{K}}$
Figure 3.2: In a) and b), a great circle for some arbitrary value of $\phi$ is drawn. The points on the circle for which a) $Y_{m=2}^{\ell=2}$ or b) $Y_{m=1}^{\ell=2}$ is zero, are indicated by white points. In a) there are no other points on which $Y_{m=2}^{\ell=2}$ is zero. In b) all other such points for $Y_{m=1}^{\ell=2}$ fall on a horizontal circle with $\theta = \pi/2$ (point-dashed circle). In both cases, the number of zeros along a great circle is at most $2 \cdot \ell$. Remember: $Y_{m=2}^{\ell=2} = \frac{1}{4} \sqrt{\frac{15}{2\pi}} \sin^2 \theta e^{2i\phi}$ and $Y_{m=1}^{\ell=2} = -\sqrt{\frac{15}{8\pi}} \sin \theta \cos \theta e^{i\phi}$.

in order to create the matching. In practice we will have to truncate at some value $\ell_{\text{max}}$. What would be a reasonable choice? For a given $\ell_{\text{max}}$, $Y_m^{\ell_{\text{max}}}(\theta, \phi)$ can have at most $2\ell_{\text{max}}$ nodes along a great circle (i.e. $\theta = 0 \rightarrow 2\pi$ for any fixed $\phi$) of the $\alpha$-sphere (see fig. 3). Converted into nodes per unit of length, this is $2\ell_{\text{max}}/(2\pi R_\alpha) = \ell_{\text{max}}/(\pi R_\alpha)$. If a plane wave should match with this, there should be plane waves with at least a similar number of nodes per unit of length available. The plane wave with the shortest period $2\pi/K_{\text{max}}$ has $2/(2\pi/K_{\text{max}}) = K_{\text{max}}/\pi$ nodes per unit of length. The cut-off for the plane waves ($K_{\text{max}}$) and for the angular functions ($\ell_{\text{max}}$) are of comparable quality if the number of nodes per unit of length is identical. This yields the condition $R_\alpha K_{\text{max}} = \ell_{\text{max}}$. This allows to determine a good $\ell_{\text{max}}$ for a given $K_{\text{max}}$. A finite value for $\ell_{\text{max}}$ means that for each APW the matching at the sphere boundaries will not be exact, but good enough to work with. It is not useful to make $\ell_{\text{max}}$ larger than the condition $R_\alpha K_{\text{max}}$ requires, as it would lead to unstable behaviour at the sphere boundary (you can compare this with fitting a polynomial of high order through a limited number of points: the fit will be ‘perfect’, but not very meaningful). Therefore, it is also clear now that the muffin tin radii for the different atoms should not be too different: if they were, a value for $\ell_{\text{max}}$ that is suitable for each atom would not exist.

Now you should be able to visualize the meaning of a single APW $\phi_{\vec{k}_K}^\alpha(\vec{r}, E)$ of equation 3.1: it is an oscillating function that runs through the unit cell. Whenever it encounters an atom on its path, the simple oscillating behaviour is changed into something more complex inside the muffin tin sphere of that atom. Nevertheless, the function values inside and outside the sphere smoothly match, which is taken care of by a set of $\sum_{\ell=1}^{\ell_{\text{max}}} 2\ell_{\text{max}} + 1$ coefficients $A_{\ell m}^{\alpha, \vec{k} + \vec{K}}$ that is different for each atom (the atom determines $\alpha$, the APW under consideration determines $\vec{k}$ and $\vec{K}$, all $\ell$ up to $\ell_{\text{max}}$ are present, with the corresponding values of $m$).

At first sight, it looks like we can now use the APW’s as a basis set, and proceed in the same
way as for the plane wave basis set in order to determine the coefficients $c_{n,k}$ in the expansion of the searched eigenfunction. However, this does not work. We did not settle the parameter $E$ yet. It turns out that in order to describe an eigenstate $\psi_{n,k}(\vec{r})$ accurately with APW’s, one has to set $E$ equal to the eigenvalue (or band energy) $\epsilon_{n,k}$ of that state. But this is exactly what we are trying to determine! We are hence forced to start with a guessed value for $\epsilon_{k}^{n}$ and take this as $E$. Now we can determine the APW’s, and construct the hamiltonian matrix elements and overlap matrix (the APW’s are not orthogonal). The secular equation is determined, and our guessed $\epsilon_{k}^{n}$ should be a root of it. Usually it is not, hence we have to try a second guess. Due to this new $E$, the APW’s have to be determined again, and similarly for all matrix elements. With the help of root determination algorithms, this guessing continues until a root – say $\epsilon_{k}^{(n=1)}$ – is found. And then the whole procedure starts over for $\epsilon_{k}^{(n=2)}$, etc. (see fig. 3 for a suggestive visualization of the roots of a secular equation, and fig. 3 for a flow chart of the APW method).

In practice, $K_{\text{max}} \approx 3.5 \text{ au}^{-1}$ is needed for sufficient accuracy. This is less than the typical value of 5.5 for plane waves and pseudopotentials. As on page 12, the basis set size can be estimated to be about $P = 131$ for APW, compared to roughly $P = 270$ for plane waves. The calculation time (mainly determined by matrix diagonalization) scales with the third power of the basis set size, which would suggest APW to be 10 times faster than pseudopotentials. However, with a plane wave basis set, $P$ eigenvalues are found by a single diagonalization, while with APW one diagonalization is needed for every eigenvalue. This makes the APW method inherently slow, much slower than the pseudopotential method.
Determine accuracy by choosing $K_{\text{max}}$ and $\ell_{\text{max}}$

For $l (n, k)$

Choose a trial $E$

$\forall \vec{K} : |\vec{K}| \leq |\vec{K}_{\text{max}}| :$ determine the corresponding APW

Calculate the matrix elements of $H$ and $S$

Determine the secular equation

E is not a root

E is a root

Determine the coefficients $c^{n,k}_k$

Figure 3.4: Flowchart of the APW method.
Chapter 4

The LAPW method

4.1 The regular LAPW method

The problem with the APW method was that the \( u_\ell(r', E) \) have to be constructed at the – yet unknown – eigenenergy \( E = \epsilon_k^n \) of the searched eigenstate. It would be helpful if we were able to recover \( u_\ell(r', \epsilon_k^n) \) on the fly from known quantities. That is exactly what the Linearized Augmented Plane Wave method enables us to do. If we have calculated \( u_\ell \) at some energy \( E_0 \), we could make a Taylor expansion to find it at energies not far away from it:

\[
\begin{aligned}
    u_\ell(r', \epsilon_k^n) &= u_\ell(r', E_0) + (E_0 - \epsilon_k^n) \frac{\partial u_\ell(r', E)}{\partial E} \bigg|_{E=E_0} + O((E_0 - \epsilon_k^n)^2)
\end{aligned}
\]  

Substituting the first two terms of the expansion in the APW for a fixed \( E_0 \) gives the definition of an LAPW. This has a price: the energy difference \( (E_0 - \epsilon_k^n) \) is unknown, and hence a yet undetermined \( B_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}} \) has to be introduced:

\[
\phi_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}}(\mathbf{r}) = \begin{cases} 
\frac{1}{\sqrt{V}} e^{i (\mathbf{k}+\mathbf{K}) \cdot \mathbf{r}} & \mathbf{r} \in I \\
\sum_{\ell, m} \left(A_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}} u_\ell^{\alpha}(r', E_0) + B_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}} \dot{u}_\ell^{\alpha}(r', E_0)\right) Y_m^\ell(\hat{r}') & \mathbf{r} \in S_{\alpha}
\end{cases}
\]

Equation 4.2 is not the final definition of an LAPW yet. Imagine we want to describe an eigenstate \( \psi_{\ell m}^{\alpha, n} \) that has predominantly p-character (\( \ell = 1 \)) for atom \( \alpha \). This means that in its expansion in LAPW’s, the \( A_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}} \) and \( B_{\ell m}^{\alpha, \mathbf{k}+\mathbf{K}} \) are large. It is therefore advantageous to choose \( E_0 \) near the centre of the p-band (how this can be done will be made clear in the examples of Chapter 7). In this way, the \( O((E_0 - \epsilon_k^n)^2) \)-term in equation 4.1 will remain small, and cutting after the linear term is certainly allowed. We can repeat this argument for every physically important \( \ell \) (s-, p-, d- and f-states, i.e. up to \( \ell = 3 \)) and for every atom. As a result, we should not choose one universal \( E_0 \), but a set of well-chosen \( E_{\ell m}^{\alpha, 1} \) up to \( \ell = 3 \) (the meaning of the index ‘1’ will become
clear in section 4.2. For higher $\ell$, a fixed value can be kept. The final definition of an LAPW is then:

$$
\phi_\vec{k}^\ell_\vec{K}(\vec{r}) = \begin{cases} 
\frac{1}{\sqrt{V_e}} e^{i(\vec{k}+\vec{K}) \cdot \vec{r}} & \vec{r} \in I \\
\sum_{\ell, m} \left( A_{\ell m}^\alpha \hat{u}_\ell^\alpha (\vec{r}', E_{1,\ell}^\alpha) + B_{\ell m}^\alpha \hat{\dot{u}}_\ell^\alpha (\vec{r}', E_{1,\ell}^\alpha) \right) Y_{\ell m}(\hat{\vec{r}}') & \vec{r} \in S_\alpha
\end{cases}
$$

(4.3)

With the $E_{1,\ell}^\alpha$ being fixed, the basis functions can be calculated once and for all. The same procedure as used for the plane wave basis set can now be applied. One diagonalization will yield $P$ different band energies for this $\vec{k}$.

The accuracy of a plane wave basis set was determined by $K_{\text{max}}$. For the APW or LAPW basis set, it is not incorrect to use the same criterion. However, a better quantity to judge the accuracy here is the product $R_{\alpha}^{\text{min}} K_{\text{max}}$ between the smallest muffin tin radius and $K_{\text{max}}$. This can be understood as follows. If the smallest muffin tin radius is increased, the closest point a plane wave can come to a nucleus moves farther away from the nucleus. The part of the wave function that need not to be described with plane waves any more, in general will have displayed the steepest behaviour, steeper than anywhere else in the interstitial region (it was closest to the nucleus). Less plane waves are needed to describe the remaining, smoother parts of the wave function. $K_{\text{max}}$ can be reduced, and a good rule of thumb is that the product $R_{\alpha}^{\text{min}} K_{\text{max}}$ should remain constant in order to have comparable accuracy. Reducing $K_{\text{max}}$ means reducing the size of the matrices, and because matrix diagonalization is very expensive, a larger $R_{\alpha}^{\text{min}}$ can significantly reduce the computation time. $R_{\alpha}^{\text{min}}$ cannot be too large on the other hand, as the spherical harmonics are not suited to describe the wave functions in the region far away from the nuclei.

Compared to a plane wave basis set, the LAPW basis set can be much smaller. The required $K_{\text{max}}$ turns out to be $K_{\text{max}} = \frac{7.5 + 9.0}{R_{\alpha}^{\text{min}}} \approx 4 \text{ au}^{-1}$, depending on the desired accuracy. This yields $P \approx 195$ as basis set size, compared to $P \approx 270$ for plane waves. The calculation time (mainly determined by matrix diagonalization) scales with the third power of the basis set size, which makes LAPW in this respect about 2 to 3 times faster than plane waves. There are other aspects however that slow down LAPW\(^1\), such that in the end it is comparable in speed with plane waves.

### 4.2 LAPW with Local Orbitals (LAPW+LO)

It was not explicitly stated so far which electron states are calculated with the LAPW method. Does it make sense to calculate the 1s orbital of Fe in bcc-Fe? No, because this electron is extremely well bound to the nucleus (-514 Ry), and will behave almost exactly as if it were in a free Fe atom. Such a state is called a core state. The criterion for a core state is that it does not participate directly in chemical bonding with other atoms. Therefore, it must be contained entirely in the muffin tin sphere. States that leak out of the muffin tin sphere, are called valence states. Valence states participate in chemical bonds, and these states are treated by LAPW. Core states are treated as in free atoms, but subject to the potential due to the valence states.

When applying this definition, it frequently happens that states with the same $\ell$ but a different principal quantum number $n$ are both valence states. For instance, due to hybridisation,

---

\(^1\)The non-orthogonal basis set for instance.
Fe in bcc-Fe will have a non-negligible amount of 4p-character in its valence states that are about 0.2 Ry below the Fermi level. But the 3p-states that are 4.3 Ry below the Fermi level are not entirely confined in the core too. Such low-lying valence states are called semi-core states. It is not clear how $E_{1,\ell}^{Fe}$ should be chosen: close to 3p, close to 4p, at an intermediate value, ...? None of the choices is optimal. This dilemma is solved by adding another type of basis function to the LAPW basis set, called a local orbital (LO). A local orbital is defined as:

$$
\phi_{\alpha,LO}^{\ell m}(\vec{r}) = \begin{cases} 
0 & \vec{r} \notin S_\alpha \\
A_{\ell m}^{\alpha,LO}u_\ell^\alpha(r', E_{1,\ell}^\alpha) + B_{\ell m}^{\alpha,LO}\hat{u}_\ell^\alpha(r', E_{1,\ell}^\alpha) + C_{\ell m}^{\alpha,LO}u_\ell^\alpha(r', E_{2,\ell}^\alpha) \right) Y_{\ell m}(\hat{r}') & \vec{r} \in S_\alpha 
\end{cases}
$$

(4.4)

A local orbital is defined for a particular $\ell$ and $m$, and for a particular atom $\alpha$. A local orbital is zero in the interstitial region and in the muffin tin spheres of other atoms, hence its name local orbital. In the muffin tin sphere of atom $\alpha$, the same $u_\ell^\alpha(r', E_{1,\ell}^\alpha)$ and $\hat{u}_\ell^\alpha(r', E_{1,\ell}^\alpha)$ as in the LAPW basis set are used, with as linearization energy $E_{1,\ell}^\alpha$ a value suitable for the highest of the two valence states (4p in our example)$^2$. The lower valence state – that is much more free-atom-like – is sharply peaked at an energy $E_{2,\ell}^\alpha$. A single radial function $u_\ell^\alpha(r', E_{2,\ell}^\alpha)$ at that same energy will be sufficient to describe it. Local orbitals are not connected to plane waves in the interstitial region, they have hence no $\vec{k}$- or $\vec{K}$-dependence. The three coefficients $A_{\ell m}^{\alpha,LO}$, $B_{\ell m}^{\alpha,LO}$ and $C_{\ell m}^{\alpha,LO}$ are determined by requiring that the LO is normalized, and has zero value and zero slope at the muffin tin boundary (= it does not leak out of the muffin tin sphere).

Adding local orbitals increases the LAPW basis set size. If for each atom local orbitals for p- and d-states are added, the basis set increases with 3+5=8 functions per atom$^3$ in the unit cell. This number is rather small compared to typical LAPW basis set sizes of a few hundred functions. The slightly increased computational time is a small price to be paid for the much better accuracy that local orbitals offer, and therefore they are always used.

---

$^2$Linearization energies for two atoms that are equivalent, are taken to be equal.

$^3$Indeed, the more atoms in the unit cell, the more LO’s have to be added. In contrast to this, the number of LAPW’s does not depend on the number of atoms in the unit cell, but – for a fixed $R_{\alpha}^\text{min}K_{\text{max}}$ and cell symmetry – on the volume of the unit cell, independent on how many atoms appear in it (more atoms mean just more sets of coefficients $A_{\ell m}^{\alpha,\vec{k}+\vec{R}}$ and $B_{\ell m}^{\alpha,\vec{k}+\vec{R}}$).
Chapter 5

The APW+lo method

5.1 The ‘pure’ APW+lo basis set

The problem with the APW method was the energy dependence of the basis set. This energy dependence could be removed in the LAPW+LO method, at the cost of a somewhat larger basis set size\(^1\). In the APW+lo method that will be described now, the basis set will be energy independent and still have the same size as in the APW method. In this sense, APW+lo combines the good features of APW and LAPW+LO.

The APW+lo basis set contains two kinds of functions. The first kind are APW’s, with a set of fixed energies \(E_{1,\ell}^\alpha\):

\[
\phi_{R}^\ell(\vec{r}) = \begin{cases} 
\frac{1}{\sqrt{V}}e^{i(\vec{k}+\vec{R})\cdot\vec{r}} & \vec{r} \in I \\
\sum_\ell, m A_{\ell m}^{\alpha, R} u_{\ell}^\alpha(\vec{r}', E_{1,\ell}^\alpha) Y_m^\ell(\hat{r}') & \vec{r} \in S_\alpha
\end{cases}
\]  

(5.1)

We know from the discussion in Chapter 3 that with fixed energies this basis set does not give a good description of the eigenfunctions. The basis set is therefore augmented with a second type of functions. These are local orbitals, but another type as the one used in connection with the LAPW method. We abbreviate them therefore as ‘lo’ instead of as ‘LO’. Their definition is:

\[
\phi_{\alpha,lo}^\ell(\vec{r}) = \begin{cases} 
0 & \vec{r} \notin S_\alpha \\
\left( A_{\ell m}^{\alpha,lo} u_{\ell}^\alpha(\vec{r}', E_{1,\ell}^\alpha) + B_{\ell m}^{\alpha,lo} \hat{u}_{\ell}^\alpha(\vec{r}', E_{1,\ell}^\alpha) \right) Y_m^\ell(\hat{r}') & \vec{r} \in S_\alpha
\end{cases}
\]  

(5.2)

The same set of energies \(E_{1,\ell}^\alpha\) is used as for the corresponding APW’s (although this is not strictly needed). The two coefficients \(A_{\ell m}^{\alpha,lo}\) and \(B_{\ell m}^{\alpha,lo}\) are determined by normalization, and by requiring that the local orbital has zero value at the muffin tin boundary (not zero slope). Hence, both the APW and the local orbital are continuous at the sphere boundary, but for both their first derivative is discontinuous.

For accurate results, the APW+lo basis set appears to require a size that is comparable to the APW method \((K_{max} \approx 3.5 \text{ au}^{-1}, P \approx 130)\). This is less than in the LAPW+LO method

\(^1\)Due to LAPW itself, and due to adding local orbitals.
(\(K_{\text{max}} \approx 4 \text{ au}^{-1}, P \approx 200\)). Nevertheless, \(P\) eigenvalues are obtained by a single diagonalization, as in LAPW+LO.

## 5.2 Mixed LAPW/APW+lo basis sets

The reason why LAPW needs a higher \(K_{\text{max}}\) than APW+lo can be traced back to a few states that are hard for LAPW. They are:

- Valence d- and f-states.
- States in atoms that have a muffin tin sphere that is much smaller than others spheres in the unit cell.

It is advantageous to treat those states with APW+lo, and keep using LAPW for all other states. Why? Using APW+lo for a state means that per atom \(2\ell + 1\) local orbitals are added to the basis set. This makes an APW+lo basis set for the same \(R_{\alpha K_{\text{max}}}^{\text{min}}\) considerably larger than the LAPW basis set. This is compensated by the fact that a lower \(R_{\alpha K_{\text{max}}}^{\text{min}}\) is needed for accurate results, but nevertheless, it is better to use these extra basis functions only there where they are useful. Such an approach leads to a mixed LAPW/APW+lo basis set: for all atoms \(\alpha\) and values of \(\ell\), equation 4.2 is used. But for one or more atoms \(\alpha_0\) \((r \in S_{\alpha_0})\) and one or more \(\ell_0\), equation 5.1 is used. Corresponding \(\phi_{\ell m}^{\alpha_0}\) according to equation 5.2 are then added to the basis set. Such a mixed basis set is the recommended choice in **WIEN2k**.

## 5.3 APW+lo with Local Orbitals (APW+lo+LO)

With APW+lo basis functions, the same problem with semi-core states appears as described in section 4.2. It will be remediated in the same way: by adding Local Orbitals (LO). The definition of a local orbital for APW+lo is:

\[
\phi_{\alpha,LO}^{\ell m}(\vec{r}) = \begin{cases} 
0 & \vec{r} \notin S_\alpha \\
(A_{\ell m}^{\alpha,LO} u_\ell^\alpha (r', E_1^\alpha) + C_{\ell m}^{\alpha,LO} u_\ell^\alpha (r', E_2^\alpha)) Y_{m}^{\ell} (\hat{r}') & \vec{r} \in S_\alpha
\end{cases}
\]

(5.3)

In contrast to the LO for LAPW, there is no derivative of \(u_\ell^\alpha\) here. The two coefficients \(A_{\ell m}^{\alpha,LO}\) and \(C_{\ell m}^{\alpha,LO}\) are determined by the requirement that the LO is normalized, and has zero value (not zero slope) at the sphere boundary.
Chapter 6

The PAW method (in brief)

Pseudopotentials with plane waves (Chap. 2) and the APW-based methods (Chaps. 3-5) are two strategies to obtain basis sets that are simultaneously unbiased and efficient. A third major way to achieve the same goal is the Projector Augmented Wave method (PAW). There are several excellent descriptions of PAW available 1, and it would not be very useful to write in this chapter yet another one. Instead, only a short overview will be given, with as main purpose to point to the similarities and differences with respect to pseudopotential and APW-based methods. In contrast to the standard PAW literature, the story will be told in a ‘backward’ way – this helps to put a different emphasis.

Just as any of the other methods mentioned before, the goal of PAW is to express the single particle all-electron Kohn-Sham wave functions $\psi_n(\vec{r})$. PAW does this by writing the all-electron wave function as a sum of a few other functions, each of which can be expressed in a natural way in a basis. The recipe which PAW uses for this, is the following one:

$$ \psi_n(\vec{r}) = \tilde{\psi}_n(\vec{r}) + \sum_a \sum_i \left( \phi_i^a(\vec{r}) - \tilde{\phi}_i^a(\vec{r}) \right) \left( \phi_i \right)_1^3 \left( \tilde{\psi}_n \right)_2^3 \langle \tilde{p}_i^a | \tilde{\psi}_n \rangle (6.1) $$

We’ll understand this recipe step by step. The recipe depends on three kinds of objects – indicated by numbers in Eq. 6.1 – for which a choice has to be made. Once the choice is made, it will be possible to write down the right-hand side of Eq. 6.1 explicitly.

The first term at the right-hand side of Eq. 6.1 is the pseudo wave function $\tilde{\psi}_n(\vec{r})$. This function is identical to the all-electron single-particle Kohn-Sham wave function $\psi_n(\vec{r})$, except inside the augmentation spheres. The remaining terms at the right-hand side of Eq. 6.1 express that difference. The atom-centered augmentation spheres play a role that is similar to the muffin tin spheres in the APW-based methods. The index ‘$a$’ in Eq. 6.1 runs over all augmentation spheres (atoms) in the unit cell. Inside the augmentation spheres, the pseudo wave function will be much smoother than the all-electron wave function, which is similar to what happens in the pseudopotential method. The relation between the all-electron and pseudo wave function can be formally expressed by transformation operators $T^a$, which operate on their corresponding

---

augmentation sphere:
\[
\psi_n(\vec{r}) = \left(1 + \sum_a T_a^{a} \right) \tilde{\psi}_n(\vec{r}) \tag{6.2}
\]

These transformation operators are such that they transform the pseudo wave function inside sphere \(a\) into the all-electron wave function. How these operators look like, is fully determined by the choice of the three numbered objects in Eq. 6.1. We’ll examine these now:

1. \(\phi_i^a(\vec{r}) : all-electron partial waves\). These are basis functions (labeled by an index \(i\)) in which the all-electron wave function can be expanded. As this expansion will only be needed within the augmentation spheres, a different choice of basis functions can be made for each atom, whence the superscript \(a\). One possible choice for these basis functions are radial solutions of the Schrödinger equation for the free atom (see App. D). This is quite similar to what is done in APW-based methods within the muffin tin spheres.

The same operator \(T\) that is used to transform all-electron and pseudo wave functions into each other, can be used to create pseudo versions of the all-electron partial waves, which are called \(pseudo partial waves\ \tilde{\phi}_i^a(\vec{r})\):
\[
\phi_i^a(\vec{r}) = (1 + T^a) \tilde{\phi}_i^a(\vec{r}) \tag{6.3}
\]

(As we look at one specific atom now, there is no sum over atoms). Remember, however, that the form of \(T^a\) was not yet specified. We will do this now, by choosing a form for the pseudo partial waves. Indeed, by rewriting the previous equation in the following way:
\[
T^a \tilde{\phi}_i^a = \phi_i^a - \tilde{\phi}_i^a \tag{6.4}
\]

we see that the action of \(T^a\) is fully defined as soon as the all-electron partial waves and pseudo partial waves are chosen. The all-electron partial waves have been chosen above (see 1). The choice of the pseudo partial waves is split into two steps: a different choice is made inside and outside the augmentation spheres:

2. \(\tilde{\phi}_i^a(\vec{r}) : pseudo partial waves \rightarrow \text{outside augmentation spheres}\). When being outside the augmentation spheres, the pseudo partial waves are chosen to be identical to the all-electron partial waves.

3. \(\tilde{p}_i^a : projector functions \rightarrow \text{pseudo partial waves inside augmentation spheres}\). The pseudo partial waves can be chosen in such a way that they form a basis for the pseudo wave function within the augmentation sphere:
\[
\tilde{\psi}_n(\vec{r}) = \sum_i C_{ni}^a \tilde{\phi}_i^a(\vec{r}) \quad \vec{r} \in \text{sphere }'a' \tag{6.5}
\]

(it can easily be shown that \(\psi_n\) can be expanded in the \(\phi_i^a\) by exactly the same coefficients \(C_{ni}^a\)). If we choose a set of projector functions that are orthonormal to the pseudo partial waves (\(\langle \tilde{p}_i^a | \tilde{\phi}_j^a \rangle = \delta_{ij}\)), then the coefficients \(C_{ni}^a\) turn out to be identical to the dot product between the projector function and the pseudo wave function:
\[
\langle \tilde{p}_i^a | \tilde{\psi}_n \rangle = C_{ni}^a \tag{6.6}
\]
Selecting appropriate projector functions for each element of the periodic table is a kind of art, and can have considerable effect on the speed and accuracy of the DFT calculation. PAW-based DFT codes typically provide one or more choices of projector functions for every element of the periodic table (called ‘POTPAW files’ in VASP, ‘PAW setups’ in GPAW, etc.).

Having these three ingredients specified, Eq. 6.1 can be explicitly written down. Apart from the smooth pseudo wave function $\tilde{\psi}_n$, at first sight. However, the choices made for $\phi$ and $\tilde{\phi}$ (and $\tilde{p}$, which co-determines $\tilde{\phi}$) do unambiguously determine $\tilde{\psi}_n$ too. Indeed, our choice of $\phi$ and $\tilde{\phi}$ has determined the operator $T$ in Eq. 6.2. We can use the latter equation to substitute $\psi_n$ in the Kohn-Sham equations by $T \psi_n$, after which the Kohn-Sham equations can be rearranged in a new equation, completely known, that can be solved for $\tilde{\psi}_n$. Therefore, we do indeed explicitly know every single ingredient of the right-hand side of Eq. 6.1.

One can get more insight into the physical meaning of Eq. 6.1 by rearranging it somewhat. Consider:

$$
\sum_i \langle \tilde{p}_i^a | \tilde{\psi}_n \rangle \phi_i^a(r) = \psi_n^a(r)
$$

(6.7)

As this is a sum over all-electron partial waves, it can be interpreted as the all-electron Kohn-Sham wave function within the augmentation sphere $a$ – including any steep parts or fast oscillations. The newly introduced symbol $\psi_n^a(r)$ means just that. Equivalently, consider:

$$
\sum_i \langle \tilde{p}_i^a | \tilde{\psi}_n \rangle \tilde{\phi}_i^a(r) = \tilde{\psi}_n^a(r)
$$

(6.8)

This is a sum over pseudo partial waves. The newly introduced symbol $\tilde{\psi}_n^a(r)$ represents therefore the smooth part of the all-electron wave function, within augmentation sphere $a$.

With these two new symbols, Eq. 6.1 can be written as the following sum of three types of terms – each of these corresponds to one of the cartoons in Fig. 6.1:

$$
\psi_n(r) = \tilde{\psi}_n(r) + \sum_a \psi_n^a(r) - \sum_a \tilde{\psi}_n^a(r)
$$

(6.9)

The interpretation is as follows. The all-electron single-particle Kohn-Sham wave function (steep near the nuclei and smooth elsewhere) can be written as the sum of (1) a pseudo wave function that is smooth everywhere, to which is added (2) a steep function that is defined only within each augmentation sphere, and from which is subtracted (3) the smooth part, only within the spheres. (1) can be conveniently represented in a plane wave basis, (2) is already represented in the basis of partial waves, and (3) is already represented in the basis of pseudo partial waves.
Figure 6.1: A slide taken from a presentation by P. Blöchl, the inventor of the PAW method. This slide is a graphical summary of Eq. 6.9.
Chapter 7

Examples for WIEN2k

The purpose of this chapter is to illustrate the meaning of (mainly) linearization energies, basis set size and k-point sampling by scrutinizing some input and output files of WIEN2k. It is not a substitute for the WIEN2k User’s Guide, yet is complimentary to it. Only command line input is used here, but reproducing this by w2web is straightforward. Version 13.1 of WIEN2k is used throughout.

7.1 Linearization energies

In this example, the input file case.in1 for WIEN2k and some output that is related to linearization energies and to the basis set will be discussed. The file case.in1 contains crucial information for the LAPW and APW+lo methods. As our example, we take fcc-Ag (F-type lattice with atom at (0, 0, 0), or alternatively space group Fm\(\overline{3}m\) (nr. 225)) at the experimental lattice constant \(a = 7.7218\) au, and with a muffin tin radius \(R_{mt} = 2.50\) au. You can produce the corresponding case.struct by w2web, or copy it from here:

```plaintext
fcc-Ag
F   LATTICE,NONEQUIV.ATOMS 1
MODE OF CALC=RELA
 7.721800 7.721800 7.721800 90.000000 90.000000 90.000000
ATOM -1: X=0.00000000 Y=0.00000000 Z=0.00000000
   MULT= 1   ISPLIT=15
Ag   NPT= 781  RO=.000010000  RMT= 2.500000  Z: 47.000000
LOCAL ROT MATRIX:
   1.0000000 0.0000000 0.0000000
   0.0000000 1.0000000 0.0000000
   0.0000000 0.0000000 1.0000000
0   NUMBER OF SYMMETRY OPERATIONS
```
7.1.1 Initialization

Initialize this Ag case by:

```
init_lapw -b -vxc 13 -ecut -6.0 -rkmax 7.0 -numk 5000
```

In this way, you have selected the PBE exchange-correlation functional (-vxc 13, which is the default and could therefore have been omitted), a separation energy of -6.0 Ry to distinguish between core and valence electrons (default, but see hereafter), a basis set size corresponding to $R_{\text{mt}}^{\text{min}} K_{\text{max}} = 7.0$ (-rkmax 7.0, again the default and could have been omitted) and 5000 k-points in the full first Brillouin zone (-numk 5000).

7.1.2 Interpreting input

The `init_lapw` command has created quite a lot of files. Have a look into `case.outputst`, and search for the following section:

<table>
<thead>
<tr>
<th>E-up(Ry)</th>
<th>E-dn(Ry)</th>
<th>Occupancy</th>
<th>q/sphere</th>
<th>core-state</th>
</tr>
</thead>
<tbody>
<tr>
<td>1S -1860.657050</td>
<td>-1860.655641</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2S -273.879379</td>
<td>-273.879742</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2P* -254.505707</td>
<td>-254.505702</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2P -241.619285</td>
<td>-241.619366</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3S -50.209917</td>
<td>-50.209130</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3P* -42.504700</td>
<td>-42.503996</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3P -40.230546</td>
<td>-40.229780</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3D* -26.594440</td>
<td>-26.594150</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3D -26.136504</td>
<td>-26.136194</td>
<td>3.00</td>
<td>3.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>4S -6.978668</td>
<td>-6.979164</td>
<td>1.00</td>
<td>1.00</td>
<td>0.9996 T</td>
</tr>
<tr>
<td>4P* -4.557328</td>
<td>-4.556808</td>
<td>1.00</td>
<td>1.00</td>
<td>0.9984 F</td>
</tr>
<tr>
<td>4P -4.183849</td>
<td>-4.182902</td>
<td>2.00</td>
<td>2.00</td>
<td>0.9977 F</td>
</tr>
<tr>
<td>4D* -0.576134</td>
<td>-0.561530</td>
<td>2.00</td>
<td>2.00</td>
<td>0.9295 F</td>
</tr>
<tr>
<td>4D -0.535644</td>
<td>-0.520181</td>
<td>3.00</td>
<td>3.00</td>
<td>0.9228 F</td>
</tr>
<tr>
<td>5S -0.348003</td>
<td>-0.290980</td>
<td>1.00</td>
<td>0.00</td>
<td>0.3372 F</td>
</tr>
</tbody>
</table>

TOTAL CORE-CHARGE: 30.000000
TOTAL CORE-CHARGE INSIDE SPHERE: 29.999275
TOTAL CORE-CHARGE OUTSIDE SPHERE: 0.000725

You see here the eigenvalues$^2$ as calculated for the free Ag-atom, for spin-up (2\textsuperscript{nd} column) and spin-down (3\textsuperscript{rd} column), together with the number of electrons in these orbitals (4\textsuperscript{th} and 5\textsuperscript{th} columns). The 6\textsuperscript{th} column gives the fraction of the total charge (up+down) within the muffin tin sphere: 3d electrons lie entirely within the muffin tin sphere, and can therefore be taken as core electrons, while nearly 8% of the 3d electrons are outside the muffin tin sphere – they must be taken as valence electrons. The separation between orbitals fully confined within the

$^1$It is assumed you have a basic understanding about performing WIEN2k calculations, i.e. you have read at least the Quick Start section in the User’s Guide.

$^2$The ‘*’ in the first column refers to the two different values for the relativistic quantum number $\kappa$. for instance, 2P* and 2P correspond to $p_{1/2}$ ($\kappa = 1$) and $p_{3/2}$ ($\kappa = -2$) orbitals.
muffin tin spheres and orbitals that aren’t, lies between 3d (at -26 Ry) and 4s (at -6.98 Ry). We had taken the borderline to distinguish between core and valence electrons to be -6.0 Ry (by taking \texttt{-ecut -6.0}). This is reflected in the 7th column where all orbitals with an energy more negative than -6 Ry are labeled by \texttt{T} to indicated that they are core states. These core states (1s+2s+2p+3s+3p+3d+4s) should contain 30 electrons, a number that is plot as \texttt{TOTAL CORE-CHARGE}. The lines thereafter give the actual electron count with the muffin tin sphere (which turns out to be 29.999275), as well as the difference between the actual count and the expected number of electrons (0.000725).

APW-based methods are not very tolerant against core electrons that leak out of the muffin tin spheres. A case with less than 0.001 electrons outside the spheres is OK, but it will always be more accurate to reduce this number even further. This can be done by lowering the separation energy to -7.5 Ry, in order to tell \texttt{WIEN2k} to consider the 4s electrons as valence state:

\texttt{init\_lapw -b -vxc 13 -ecut -7.5 -rkmax 7.0 -numk 8000}

In that case, this section of \texttt{case\_output} would look like this:

<table>
<thead>
<tr>
<th>E-up (Ry)</th>
<th>E-dn (Ry)</th>
<th>Occupancy</th>
<th>q/sphere</th>
<th>core-state</th>
</tr>
</thead>
<tbody>
<tr>
<td>1S -1860.657050</td>
<td>-1860.655641</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2S -273.879379</td>
<td>-273.879742</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2P* -254.505707</td>
<td>-254.505702</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>2P -241.619285</td>
<td>-241.619366</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3S -50.209917</td>
<td>-50.209130</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3P* -42.504700</td>
<td>-42.503996</td>
<td>1.00</td>
<td>1.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3P -40.230546</td>
<td>-40.229780</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3D* -26.594440</td>
<td>-26.594150</td>
<td>2.00</td>
<td>2.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>3D -26.136504</td>
<td>-26.136194</td>
<td>3.00</td>
<td>3.00</td>
<td>1.0000 T</td>
</tr>
<tr>
<td>4S -6.978668</td>
<td>-6.979164</td>
<td>1.00</td>
<td>1.00</td>
<td>0.9996 F</td>
</tr>
<tr>
<td>4P* -4.557328</td>
<td>-4.556808</td>
<td>1.00</td>
<td>1.00</td>
<td>0.9984 F</td>
</tr>
<tr>
<td>4P -4.183849</td>
<td>-4.182902</td>
<td>2.00</td>
<td>2.00</td>
<td>0.9977 F</td>
</tr>
<tr>
<td>4D* -0.576134</td>
<td>-0.561530</td>
<td>2.00</td>
<td>2.00</td>
<td>0.9295 F</td>
</tr>
<tr>
<td>4D -0.535644</td>
<td>-0.520181</td>
<td>3.00</td>
<td>3.00</td>
<td>0.9228 F</td>
</tr>
<tr>
<td>5S -0.348003</td>
<td>-0.290980</td>
<td>1.00</td>
<td>0.00</td>
<td>0.3372 F</td>
</tr>
</tbody>
</table>

\texttt{TOTAL CORE-CHARGE: 28.000000}
\texttt{TOTAL CORE-CHARGE INSIDE SPHERE: 28.000000}
\texttt{TOTAL CORE-CHARGE OUTSIDE SPHERE: 0.000000}

The \texttt{F} in the final column indicates that the 4s electrons are now considered as valence states. Only 28 core electrons are expected now, and indeed, all of these are completely inside the muffin tin sphere. With a separation energy of -4.0 Ry (\texttt{-ecut -4.0}), about 0.013 electrons (=a fraction of 0.998 of six 4p-electrons) would be outside the muffin tin sphere. This is unacceptably large, and would trigger \texttt{WIEN2k} to produce a warning in the screen output, identifying the problematic orbitals:
WARNING: 0.013 Ag CORE electrons leak out of MT-sphere !!!!
WARNING: touch .lcore and run scf-cycle with core density superposition
WARNING: Or: rerun lstart with lower E-core separation energy
WARNING: ORBITAL: 4S -6.979 -6.979
WARNING: ORBITAL: 4P* -4.557 -4.557
WARNING: ORBITAL: 4P -4.184 -4.183

We continue with the case where the initialization was done with -ecut -7.0, and inspect the file case.in1:

```
WFFIL EF= 0.50000 (WFFIL, WFPRI, ENFIL, SUPWF)

7.0 10 4 (R-MT*K-MAX; MAX L IN WF, V-NMT
0.30 5 0 (GLOBAL E-PARAMETER WITH n OTHER CHOICES, global APW/LAPW)
0 0.30 0.000 CONT 1
0 -6.68 0.001 STOP 1
1 0.30 0.000 CONT 1
1 -3.88 0.001 STOP 1
2 0.30 0.005 CONT 1

K-VECTORS FROM UNIT:4 -10.0 1.5 21 emin / de (emax=Ef+de) / nband
```

The value 7.0 in the first line – which is due to choosing -rkmax 7.0 – is the product of the minimal muffin-tin radius $R_{\text{mt}}^{\text{min}}$ (here 2.5 au) and the $K_{\text{max}}$ discussed in Chapter 4. By 7.0, we have asked for $K_{\text{max}} = 7.0/2.5 = 2.8$ au$^{-1}$.

The number ‘10’ to the right of $R_{\text{mt}}^{\text{min}}K_{\text{max}}$ is the $l_{\text{max}}$ from page 17. The ‘4’ to the right of it, is not discussed in the preceding sections (it is highest $l$ in the expansion of the potential in lattice harmonics, a kind of spherical harmonics that is adapted to the symmetry of the lattice).

In the next line, the ‘0.30’ is the linearization energy $E_{1,\ell}^{\alpha}$ ($\ell = 0, \ldots, l_{\text{max}}, \alpha = \text{Ag}$) in Eqs. 4.3 or 5.1+5.2, expressed in Ry $^3$. The ‘5’ to the right means that for 5 values of $\ell$ an exception will be made. These 5 exceptions will be listed below. The ‘0’ to the right means that for all $\ell$ a LAPW is used, except if explicitly stated otherwise (with ‘1’, an APW would be used everywhere).

The following 5 lines list exceptions to this general choice of $E_{1,\ell}^{\alpha}$. The first number in each line is $\ell$. If the same $\ell$ is given twice, then the first one specifies a (L)APW(+lo) basis function (Eqs. 4.3 or 5.1+5.2) while the second one specifies a Local Orbital (LO) (Eqs. 4.4 or 5.3). The second number of each line is the linearization energy: $E_{1,\ell}^{\alpha}$ in case of a (L)APW(+lo) (Eqs. 4.3 or 5.1+5.2), and $E_{2,\ell}^{\alpha}$ in case of a Local Orbital (LO) (Eqs. 4.4 or 5.3). The third column gives a small number 0.00x. In case it is non-zero, an optimal linearization energy will be searched with incremental energy steps of 0.00x (see footnote 7 on page 38). This procedure is used only for semi-core states (=at strongly negative energies) or for partially filled valence states (like Ag-4d). How optimal linearization energies are found for states that are empty in the free atom (as Ag-5p and all $\ell \geq 3$), will be described on page 38. The CONT means that even if no acceptable linearization energy would be found, the program will nevertheless continue. If this would happen for a semi-core state – where it should be no problem at all to find a good linearization energy – it indicates a serious problem. Therefore, those lines carry a STOP keyword, which tells WIEN2k to stop if no good linearization energy would be found. The final

---

$^3$See footnote 6 on page 36 to know how the zero-point on this energy scale is defined.
column indicates whether for this particular state a LAPW (‘0’) or an APW+lo (‘1’) will be used. This overrules the general choice specified before for all \( \ell \). In this particular example, LAPW basis functions are used for \( 3 \leq \ell \leq 10 \), while APW+lo basis functions are used for \( 0 \leq \ell \leq 2 \).

### 7.1.3 Interpreting output

So far, we looked at input files only. Run now a scf-cycle by the command:

```
run_lapw -cc 0.000005
```

This will take 13 iterations (and about 2 minutes on a normal pc) to provide you with the selfconsistent solution. By:

```
grepline :RKM '*.scf' 1
```

you will find some information on the size of the basis set that has been used:

```
case.scf::RKM : MATRIX SIZE 64LOs: 13 RKM= 6.75 WEIGHT= 1.00 PGR:
```

The value RKM= 6.75 corresponds to our -rkmax 7.0 (due to discrete jumps of \( K_{\text{max}} \), 6.75 is the value that comes closest to 7.0). The MATRIX SIZE is the total number of basis functions: 4 LO’s (1+3, for the 3s and 3p semi-core states, see case.in1 above) and 9 lo’s (1+3+5 for \( \ell \) being either 0, 1 or 2, i.e. all states that are described by a APW+lo), summing to 4 + 9 = 13 – the latter is indeed the number given by L0s: 13. A total of 64 − 13 = 51 basis functions are left for the LAPW’s and the APW part of the APW+lo’s. Following the reasoning on page 12, the estimated basis set size is 43, which is indeed consistent.

In order to help our understanding of the linearization energies, produce now Density of States plots for this case. Prepare the following case.int:

---

\(^4\)The volume of the first Brillouin zone for fcc-Ag can be calculated to be \( V_{BZ} = 8\pi^3/V = 2.155\text{au}^{-3} \). Here \( V \) is the volume of the unit cell in direct space, which is for the fcc structure one quarter of the volume of the conventional cube (\( V = (7.7218\text{au})^3/4 \)). This volume can always be found at the :VOL label in case.scf. Alternatively, \( V_{BZ} \) can be found by equations A.14 to A.16, considering that \( \vec{a} = (0, \frac{a}{2}, 0) \), \( \vec{b} = (\frac{a}{2}, 0, \frac{a}{2}) \) and \( \vec{c} = (\frac{a}{2}, \frac{a}{2}, 0) \) with \( a \) the lattice constant.

\(^5\)Our value of 43 is an estimate, not an exact value. The unit cell of Ag is quite small, hence the distances between \( \vec{K} \)-points are large. \( K_{\text{max}} = 2.8 \text{au}^{-1} \) is not very large. Our estimate approximates the number of \( \vec{K} \)-points in a sphere with radius \( K_{\text{max}} \), and for a small radius the surface region is rather important relative to the volume. Details of the reciprocal lattice will determine whether certain \( \vec{K} \) are just inside or just outside this sphere. For larger unit cells (try with the Ag lattice constant increased to 12 au) or larger \( K_{\text{max}} \) (try with \( R_{\text{min}} K_{\text{max}} = 9 \)), the estimate reproduces almost the exact basis set size. A question to test yourself: explain why in both of these cases the basis set size increases (182 and 102, respectively).
fcc Ag
-0.50 0.002 1.500 0.003 # EMIN, DE, EMAX, Gauss-broadening(>de)
7 N 0.000 # NUMBER OF DOS-CASES below, G/L/B broadening (Ry)
 0 1 total # atom, case=column in qtl-header, label
   1 Ag-tot
   1 2 Ag-s
   1 3 Ag-p
   1 4 Ag-d
   1 7 Ag-f
2 1 Inter
SUM: 0 2 # NUMBER OF SUMMATIONS, max-nr-of summands
2 5 # this sums dos-cases 2+5 from the input above

After having run:

x lapw2 -qtl
x tetra

this results into the total DOS of the crystal (4th line of case.int), the total DOS inside the Ag muffin tin sphere (5th line), the s-DOS of Ag inside the sphere (6th line), similarly for the p-, d- and f-DOS of Ag (lines 7-9), and the total DOS in the interstitial region (10th line). The data for these densities of states can be found in case.dos1 and are shown in Fig. 7.1.3. The Fermi energy (:FER) – indicated by a dashed vertical line in fig. 7.1.3 – is about 0.465 Ry:

```
grep :FER case.scf
```

<table>
<thead>
<tr>
<th>:FER</th>
<th>ENERGY(TETRAH.M.)</th>
</tr>
</thead>
<tbody>
<tr>
<td>:FER</td>
<td>0.6030637284</td>
</tr>
<tr>
<td>:FER</td>
<td>0.595461790</td>
</tr>
<tr>
<td>:FER</td>
<td>0.5780000291</td>
</tr>
<tr>
<td>:FER</td>
<td>0.5172412001</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4592686207</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4641690407</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4649683335</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4652388160</td>
</tr>
<tr>
<td>:FER</td>
<td>0.46534998651</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4653598853</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4653487985</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4653488903</td>
</tr>
<tr>
<td>:FER</td>
<td>0.4653487700</td>
</tr>
</tbody>
</table>

---

6 In a free atom, the zero of the energy scale is chosen such that the potential is zero at infinite distance from the atom. In an infinite lattice, this is not possible. In WIEN2k, the zero corresponds to the average potential in the interstitial region. Whenever an energy is expressed in Ry, it is relative to this zero-level. The DOS can be chosen to be plotted against a Ry-scale (and then this zero is used) or against a scale in eV. In the latter case, the Fermi energy is taken zero by default.
Figure 7.1: DOS for fcc-Ag. The dashed line marks the Fermi energy. Mind the two different vertical scales: up to 4 states per Ry per atom and up to 80 states per Ry per atom. This picture shows 6 out of the 7 cases listed in the case.int file on page 35.
After having done the selfconsistent calculation, *case.in1* will look like this:

```plaintext
WFFIL  EF=.4653487700  (WFFIL, WFPRI, ENFIL, SUPWF)
7.0  10  4 (R-MT*K-MAX; MAX L IN WF, V-NMT
  0.30  0  0  (GLOBAL E-PARAMETER WITH n OTHER CHOICES, global APW/LAPW)
0  0.30  0.000 CONT 1
0 -6.68  0.001 STOP 1
1  0.30  0.000 CONT 1
1 -3.88  0.001 STOP 1
2  0.30  0.005 CONT 1
K-VECTORS FROM UNIT:4  -10.0  1.5  21  emin / de (emax=Ef+de) / nband
```

It is identical to the file as we saw it before running the selfconsistent calculation, except for the fact that the Fermi energy is written in the top line. As *case.in1* is input for the *lapw1* routine, the latter can now use information about the Fermi energy: all linearization energies with ‘0.000’ (i.e. those for which no search for an optimal value was done) will be put 0.2 Ry below the Fermi energy. The actual values for every linearization energy can be read from *case.scf* (we look at the final iteration here):

```
ATOMIC SPHERE DEPENDENT PARAMETERS FOR ATOM  Ag
:e__0001:  OVERALL ENERGY PARAMETER IS  0.2653
OVERALL BASIS SET ON ATOM IS LAPW
:E0_0001: E( 0)=  0.2653
  APW+lo
:E0_0001: E( 0)=   -6.3015  E(BOTTOM)=  -6.312  E(TOP)=  -6.291 3 4 137
  LOCAL ORBITAL
:E1_0001: E( 1)=  0.2653
  APW+lo
:E1_0001: E( 1)=   -3.6220  E(BOTTOM)=  -3.661  E(TOP)=  -3.583 2 3 126
  LOCAL ORBITAL
:E2_0001: E( 2)=  0.1282  E(BOTTOM)=  -0.161  E(TOP)=   0.418 1 2 134
  APW+lo
```

The **OVERALL ENERGY PARAMETER** of 0.2653 Ry is the $E_{\ell}^{\alpha}$ for all $\ell \leq \ell_{\text{max}}$ for which no exception was defined. The same value is used for the two APW+lo basis functions for $\ell = 0$ (5s) and $\ell = 1$ (5p). The value 0.2653 is indeed the Fermi energy minus 0.2 Ry. For the 3 orbitals where a search with a non-zero increment was performed it, the actual linearization energies turn out to be -6.3015 (4s), -3.6220 (4p) and 0.1282 (3d).

The DOS-pictures in Fig. 7.1.3 confirm that those linearization energies make sense: the occupied part of the valence band covers a range from -0.05 Ry to 0.46 Ry (=EF), and the global linearization energy of 0.26 Ry is close to the center of this energy range. The latter was a requirement for a meaningful linearization energy (see Chap. 4). The d-states cover a smaller range, from about -0.05 Ry to 0.25 Ry. They had a non-zero increment of 0.005 Ry in *case.in1*, which allowed to search for a value of the linearization energy that lies closer to the

---

7 See the Usersguide, section about lapw1, to see how this is determined: the $E_{\ell}$ for which $u_{\ell=2}(r, E_{\ell})$ has zero slope at $r = R_{mt}$ is taken as the bottom of the band (bonding state). The $E_{\ell}$ for which $u_{\ell=2}(r, E_{\ell})$ has zero value at $r = R_{mt}$ is taken as the top of the band (antibonding state).
center of this particular band: \(0.1282\) Ry was found, which is indeed a value at the center of the 4d band. The semi-core states have eigenvalues that form a very narrow band, with a width of \(0.1\) Ry or less (see the \(E(\text{BOTTOM})\) and \(E(\text{TOP})\) values, which are close to each other). They would appear much more to the left on Fig. 7.1.3, and would show up as a very sharp peak. This is equivalent to saying that there is hardly any dispersion (\(=k\)-point dependence) for the eigenvalues of those orbitals. Therefore, one can identify those orbitals among the eigenvalues that are printed for one k-point (the one near the origin) in \textit{case.scf}:

\[
\begin{array}{cccc}
K &=& 0.00000 & 0.00000 & 0.00000 & 1 \\
\text{:RKM : MATRIX SIZE} &=& 64\text{LOs:} & 13 & \text{RKM=} & 6.75 & \text{WEIGHT=} & 1.00 & \text{PGR:} \\
\text{EIGENVALUES ARE:} & & & & & \\
\text{:EIG00001:} &=& -6.3009210 & -3.6087514 & -3.6087514 & -3.6087514 & -0.1121691 \\
\text{:EIG00006:} &=& 0.1138505 & 0.1138505 & 0.1138505 & 0.1859000 & 0.1859000 \\
\text{:EIG00011:} &=& 1.6414487 & 1.8971107 & 1.8971107 & 1.8971107 & \\
\end{array}
\]

Indeed, the single value at \(-6.30\) Ry is the 4s-orbital (2 electrons, degenerate for up and down), and the three values at \(-3.61\) Ry are the 4p orbitals (6 electrons, degenerate for up and down). The same information can be read from the following section of \textit{case.scf}:

\[
\begin{array}{c}
\text{Bandranges (emin - emax) and occupancy:} \\
\text{:BAN00001:} &=& 1 & -6.300921 & -6.299224 & 2.00000000 \\
\text{:BAN00002:} &=& 2 & -3.620654 & -3.608751 & 2.00000000 \\
\text{:BAN00003:} &=& 3 & -3.616564 & -3.608751 & 2.00000000 \\
\text{:BAN00004:} &=& 4 & -3.613515 & -3.608751 & 2.00000000 \\
\text{:BAN00005:} &=& 5 & -0.112169 & 0.072029 & 2.00000000 \\
\text{:BAN00006:} &=& 6 & 0.000638 & 0.124086 & 2.00000000 \\
\text{:BAN00007:} &=& 7 & 0.089604 & 0.243048 & 2.00000000 \\
\text{:BAN00008:} &=& 8 & 0.112973 & 0.260414 & 2.00000000 \\
\text{:BAN00009:} &=& 9 & 0.160076 & 0.261012 & 2.00000000 \\
\text{:BAN00010:} &=& 10 & 0.185053 & 0.384880 & 0.99999970 \\
\text{:BAN00011:} &=& 11 & 0.732345 & 1.687525 & 0.00000000 \\
\text{:BAN00012:} &=& 12 & 0.991261 & 1.897111 & 0.00000000 \\
\text{:BAN00013:} &=& 13 & 1.092676 & 1.938713 & 0.00000000 \\
\text{:BAN00014:} &=& 14 & 1.673633 & 1.964410 & 0.00000000 \\
\text{:BAN00015:} &=& 15 & 1.761252 & 1.954693 & 0.00000000 \\
\text{Energy to separate low and high energystates:} &=& -0.16217 \\
\end{array}
\]

It lists the minimal and maximal eigenvalue of every band (i.e. not just for a single k-point, but by considering all k-points), as well as the occupancy of the band (full occupancy is 2 for non-spinpolarized calculations, and 1 for spinpolarized calculations). In agreement with what we saw for the example k-point \(0.00000 \ 0.00000 \ 0.00000 \), there are two (4s) electrons in a very narrow energy range (-6.301, -6.299), and six (4p) electrons in the range (-3.621, -3.609). Still eleven electrons are left, which correspond to ten 4d- and one 5s-electron. They appear in bands that show a non-negligible dispersion, which means that such a band cannot assigned uniquely any more to a free atom orbital any more. These dispersed bands are the ones that

\footnote{In \textit{case.scf} only a limited number of bands is shown. If you want to see all of them, you’ll find this in \textit{case.output2} or in \textit{case.outputt}.}
appear in the DOS plots in fig. 7.1.3, and the linearization energies are chosen such that they fall “near the center of” the occupied part of those bands.

### 7.1.4 The -in1orig and -in1new options

It might happen one wants to keep some or all linearization energies fixed. This can be done by using the -in1orig switch with run(sp) lapw. In that case, case.in1 will be neglected and case.in1Orig will be used instead (the latter file has to be created if not present, for instance as a copy of case.in1). Linearization energies with a non-zero step size in case.in1Orig will be searched in the usual way. Linearization energies where the step size was put to 0.000 in case.in1Orig will be kept fixed to the value given in case.in1. The latter holds true for the global linearization energy as well.

There is an alternative procedure ⁹ to select linearization energies, which is used if the -in1new option is specified. For this procedure, WIEN2k needs to know which states are true valence states, and which states are semi-core states (that can be described by LO’s). The semi-core states are separated by a gap from the continuous region of valence states. By the following line in case.in2:

```
-13.0 19.0 0.50 0.05 EMIN, NE, ESEPERMIN, ESEPER0
```

the program is told that it should start to look at \( E_{\text{seper}}^0 = 0.05 \) Ry below the Fermi energy, for a gap between states with a width of at least \( E_{\text{seper}}^{\text{min}} = 0.50 \) Ry. In our Ag-example where the Fermi energy is about 0.46 Ry, it will start to search at 0.46-0.05=0.41 Ry. Then it goes downwards, and at about -0.05 Ry it finds a gap that extends down to -3.6 Ry: above this gap are the true valence states, below are the semi-core states that are treated by LO’s.

In case.scf, the line labeled with :EPLxxx (‘xxx’ is the number of the atom) keeps track of the average energy for every \( \ell \) for atom ‘xxx’: the arithmetic mean between the bottom and top of the band. The line labeled with :EPHxxx does the same for the valence states, but unlike for the semi-core states the average energy is now the arithmetic mean between the bottom of the band and the Fermi energy (because higher states are not occupied).

All this is done always, but has no effect unless the switch -in1new N is given with run_lapw, with N a integer number. With this switch, after N iterations a new case.in1 is made based on the information in the :EPLxxx and :EPHxxx from the previous iteration. The old case.in1 is copied to case.in1Orig. Using the -in1new switch for our case of fcc-Ag:

```
save_lapw -a Ag-default
run_lapw -cc 0.000005 -in1new 2
```

leads after 20 additional iterations to to the following linearization energies (from case.scf):

```
-13.0 19.0 0.50 0.05 EMIN, NE, ESEPERMIN, ESEPER0
```

---

⁹The scheme described on the preceding pages is the default behaviour from WIEN2k 10.1 onwards. In version 9.2, this procedure was used only if the switch -in1ef was used. Since version 10.1, the switch -in1ef has no meaning any longer. Prior to version 9.1, the default behaviour was to keep the global linearization energies as well as the ones with a zero step size (0.000) in case.in1 fixed at 0.30 Ry – i.e. implicitly assuming a typical value of 0.50 Ry for the Fermi energy.
ATOMIC SPHERE DEPENDENT PARAMETERS FOR ATOM Ag

\[ e_{0001} \]: OVERALL ENERGY PARAMETER IS 0.2653

OVERALL BASIS SET ON ATOM IS LAPW

\[ E_{0001} \]: \( E(0) = 0.1230 \)

\[ E_{0001} \]: \( E(0) = -6.3014 \) \( E(\text{BOTTOM}) = -6.312 \) \( E(\text{TOP}) = -6.291 \)

\[ E_{1001} \]: \( E(1) = 0.2110 \)

\[ E_{1001} \]: \( E(1) = -3.6219 \) \( E(\text{BOTTOM}) = -3.661 \) \( E(\text{TOP}) = -3.583 \)

\[ E_{2001} \]: \( E(2) = 0.1470 \)

The global linearization energy (OVERALL ENERGY PARAMETER) is still taken to be 0.2 Ry below the Fermi energy, and the 4s and 4p semi-core states are treated in the same way as before – no differences here. For the 5s-, 5p- and 4d-states, the \texttt{-in1new} procedure was used. This gives minor changes for 5p and 4d, while for 5s the linearization energy is a little smaller and indeed closer to the center of the occupied part of the s-band (see Fig. 7.1.3). The freedom to give each of these states a linearization energy that falls at the center of each individual band, makes the \texttt{-in1new} procedure slightly more accurate. This is reflected in the total energy of \(-10634.58512460\) Ry, which is slightly lower (=better) than the previous value of \(-10634.58510644\) Ry. In nearly all cases, however, this almost uniform shift in total energy will not affect any observable property (=equilibrium volumes will be identical, etc.). As the \texttt{-in1new} procedure can be misled by particular band structures, it is not recommended to use it as a default. Use it only if you know what you’re doing and if you have specific reasons to do so.

7.1.5 high-lying local orbitals

In some cases, you may want to describe the unoccupied states better (for instance, for spectroscopy studies). This can be done by adding Local Orbitals with a linearization energy that lies above the Fermi energy. You can fix this linearization energy at a value you have determined after inspecting the DOS, or you can let \texttt{wienn2k} determine it as 0.2 Ry added to the Fermi energy. The latter is done by adding a local orbital with the default 0.30 Ry linearization energy and no searching (0.000). In our fcc-Ag example:

\[
\text{WFFIL EF=4.699838177 (WFFIL, WFPRI, ENFIL, SUPWF)}
\]

\[
7.0 10 4 (R-MT*K-MAX; MAX L IN WF, V-NMT)
\]

\[
0.30 6 0 (GLOBAL E-PARAMETER WITH \text{n} OTHER CHOICES, global APW/LAPW)
\]

\[
0 0.30 0.000 \text{ CONT 1}
\]

\[
0 -6.68 0.001 \text{ STOP 1}
\]

\[
1 0.30 0.000 \text{ CONT 1}
\]

\[
1 -3.88 0.001 \text{ STOP 1}
\]

\[
2 0.30 0.05 \text{ CONT 1}
\]

\[
2 0.30 0.000 \text{ CONT 1}
\]

\[
K-VECTORS \text{ FROM UNIT:} 4 -10.0 1.5 21 \text{ emin / de (emax=Ef+de) / nband}
\]
(mind the ‘6’ in the third line, as there are now 6 ‘exceptions’). This will lead to the following choice of linearization energies (from case.scf):

ATOMIC SPHERE DEPENDENT PARAMETERS FOR ATOM  Ag
:e__0001: OVERALL ENERGY PARAMETER IS  0.2700
OVERALL BASIS SET ON ATOM IS LAPW
:E0_0001: E( 0)=  0.2700
APW+lo
:E0_0001: E( 0)= -6.2997 E(BOTTOM)= -6.310 E(TOP)= -6.289  3  4  135
LOCAL ORBITAL
:E1_0001: E( 1)=  0.2700
APW+lo
:E1_0001: E( 1)= -3.6204 E(BOTTOM)= -3.660 E(TOP)= -3.581  2  3  125
LOCAL ORBITAL
:E2_0001: E( 2)=  0.1294 E(BOTTOM)= -0.160 E(TOP)=  0.419  1  2  137
APW+lo
:E2_0001: E( 2)=  0.6700
LOCAL ORBITAL

5s and 5p have a linearization energy of 0.27 Ry (0.2 below E_F), 4d has an adapted linearization energy at 0.13 Ry (as before), while the extra d-LO is at 0.67 Ry (0.2 above E_F).

## 7.2 Finding the best $R_{mt}K_{max}$ and k-mesh

### 7.2.1 Why do we need this?

We want to use a basis set that is large enough to be accurate, but only just large enough in order not to waste computer time. The same holds for the density of k-mesh. Therefore, you need to do some testing before you go to real calculations for your problem. Often it is possible to find a small analog to your problem, where you can relatively quickly do the tests. For instance, if you want to calculate properties of an impurity in a solid, you need a large supercell. But the pure solid is described by a much smaller cell. You can determine a good basis set and k-mesh size for this small cell\(^{10}\), and extrapolate it to the supercell. It is very important to determine a good basis set size and k-mesh before doing lengthy calculations and/or interpreting physical results. If you don’t, you might be wasting lots of computer time (if you use higher precision than you need) or you are producing nonsense instead of physics (if you use less precision than you need).

### 7.2.2 General procedure

A procedure will be described to make a fair choice of basis set size and k-point density. We will take non-magnetic hcp-Fe as an example. This is the crystal structure that an iron crystal adopts when it is subject to an external pressure of more than 12 GPa. Its case.struct is:

\(^{10}\)And perhaps double-check your conclusion for a cell with 50% impurities, which will probably be a small cell as well.
hcp-Fe (exp, 12.7 GPa)
H LATTICE,NONEQUIV.ATOMS: 1 194 P63/mmc
MODE OF CALC=RELA unit=bohr
  4.673293  4.673293  7.487095  90.000000  90.000000120.000000
ATOM  1: X=0.33333333 Y=0.66666667 Z=0.75000000
  MULT= 2 ISPLIT=15
  1: X=0.66666667 Y=0.33333333 Z=0.25000000
Fe1 NPT= 781 R0=.000005000 RMT= 2.20000 Z: 26.
LOCAL ROT MATRIX: 1.0000000 0.0000000 0.0000000
  0.0000000 1.0000000 0.0000000
  0.0000000 0.0000000 1.0000000
  0 NUMBER OF SYMMETRY OPERATIONS

Make a default initialization (see before), only asking for a small number of k-points\(^{11}\) (in this case: 6 k-points in the irreducible part of the first Brillouin zone – tail case.klist):

```
init_lapw -b -vxc 13 -ecut -7.0 -rkmax 7.0 -numk 30
```

Run this to convergence by

```
run_lapw -cc 0.000001
```

Record the total energy and the main component of the electric-field gradient tensor\(^{12}\):

```
grep :ENE case.scf
:ENE : ********** TOTAL ENERGY IN Ry = -5091.14367665
```

```
grep :EFG001 case.scf
:EFG001: EFG = 3.54133 *10**21 V / m**2
```

Save this calculation (save_lapw -a hcp-Fe_RKM70,k0006), run x kgen and double the total number of k-points from 30 to 60. Rerun, save again. Repeat for 120, 240, 480,\ldots k-points. You will see how every time the calculation takes more time. Finally, collect information about the number of k-points in the irreducible part of the Brillouin zone, the total energy, and the electric-field gradient:

```
grepline :KPT '*.scf' 1
grepline :ENE '*.scf' 1
grepline :EFG001 '*.scf' 1
```

\(^{11}\)Play with the \texttt{-numk} value – i.e. the number of k-points in the full Brillouin zone – until you find the value you need to have at least 4 k-points in the irreducible part of the Brillouin zone. You can find this in case.klist. Once you have run \texttt{init_lapw} once, you can do the k-mesh generation separately by \texttt{x kgen}.

\(^{12}\)The electric-field gradient is monitored here because it is a rather sensitive quantity. If your basis set size and k-mesh are sufficient to have the field gradient converged, they will be good enough for many other properties as well. Generally speaking, you should monitor the quantity you are most interested in, for instance a force (in particular if you will do optimization of internal coordinates), a magnetic moment, \ldots
which yields:

```
hcp-Fe_RKM70_k0006.scf::KPT  :  NUMBER OF K-POINTS:  6
hcp-Fe_RKM70_k0008.scf::KPT  :  NUMBER OF K-POINTS:  8
hcp-Fe_RKM70_k0014.scf::KPT  :  NUMBER OF K-POINTS: 14
hcp-Fe_RKM70_k0024.scf::KPT  :  NUMBER OF K-POINTS: 24
hcp-Fe_RKM70_k0036.scf::KPT  :  NUMBER OF K-POINTS: 36
hcp-Fe_RKM70_k0076.scf::KPT  :  NUMBER OF K-POINTS: 76
hcp-Fe_RKM70_k0135.scf::KPT  :  NUMBER OF K-POINTS: 135
hcp-Fe_RKM70_k0240.scf::KPT  :  NUMBER OF K-POINTS: 240
hcp-Fe_RKM70_k0427.scf::KPT  :  NUMBER OF K-POINTS: 427
hcp-Fe_RKM70_k0819.scf::KPT  :  NUMBER OF K-POINTS: 819
hcp-Fe_RKM70_k1540.scf::KPT  :  NUMBER OF K-POINTS: 1540
hcp-Fe_RKM70_k3038.scf::KPT  :  NUMBER OF K-POINTS: 3038
```

```
hcp-Fe_RKM70_k0006.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.14367665
hcp-Fe_RKM70_k0008.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17971354
hcp-Fe_RKM70_k0014.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17312332
hcp-Fe_RKM70_k0024.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17381364
hcp-Fe_RKM70_k0036.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17725858
hcp-Fe_RKM70_k0076.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17587588
hcp-Fe_RKM70_k0135.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17633794
hcp-Fe_RKM70_k0240.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17609663
hcp-Fe_RKM70_k0427.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17618021
hcp-Fe_RKM70_k0819.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17617023
hcp-Fe_RKM70_k1540.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17615363
hcp-Fe_RKM70_k3038.scf::ENE   : ********** TOTAL ENERGY IN Ry = -5091.17616443
```

```
hcp-Fe_RKM70_k0006.scf::EFG001: EFG  =  3.54133  *10**21 V / m**2
hcp-Fe_RKM70_k0008.scf::EFG001: EFG  =  0.27851  *10**21 V / m**2
hcp-Fe_RKM70_k0014.scf::EFG001: EFG  = -0.46525  *10**21 V / m**2
hcp-Fe_RKM70_k0024.scf::EFG001: EFG  = -0.04741  *10**21 V / m**2
hcp-Fe_RKM70_k0036.scf::EFG001: EFG  = -0.46080  *10**21 V / m**2
hcp-Fe_RKM70_k0076.scf::EFG001: EFG  = -0.41362  *10**21 V / m**2
hcp-Fe_RKM70_k0135.scf::EFG001: EFG  = -0.31932  *10**21 V / m**2
hcp-Fe_RKM70_k0240.scf::EFG001: EFG  = -0.28007  *10**21 V / m**2
hcp-Fe_RKM70_k0427.scf::EFG001: EFG  = -0.26982  *10**21 V / m**2
hcp-Fe_RKM70_k0819.scf::EFG001: EFG  = -0.25721  *10**21 V / m**2
hcp-Fe_RKM70_k1540.scf::EFG001: EFG  = -0.25695  *10**21 V / m**2
hcp-Fe_RKM70_k3038.scf::EFG001: EFG  = -0.25181  *10**21 V / m**2
```

Fig. 7.2 shows this information graphically: the total energy (left axis) and electric-field gradient (right axis) as a function of the number of k-points in the irreducible part of the first Brillouin zone. The same information is given a second time, but now as the logarithm of the difference between a quantity at a given number of k-points and at the largest number of k-points. The latter way of representing shows immediately the number of converged digits. Some conclusions that can be drawn from this test are:

1. Using only 6 or 8 k-points does produce nonsense for hcp-Fe.
2. From 14 k-points onwards, at least the order of magnitude is right.
3. If you are satisfied with 0.001 Ry accuracy in total energy, then at least 36 k-points are required.
4. If you are satisfied with 0.0001 Ry accuracy in total energy, then at least 240 k-points are required.

5. If you need 0.00001 Ry accuracy in total energy, then (probably) 1540 or more k-points are required.

6. If you need to know only the first digit of the electric-field gradient, then 135 k-points are sufficient.

7. If you need to know the second digit as well, then 1540 k-points will be a safer choice.

Let us assume that for the problem you are investigating, you need the kind of accuracy that corresponds to 240 k-points (this corresponds to 3840 k-points in the full Brillouin zone). We’ll fix this value from here on, and will vary now the basis set size from 5.5 (=very small basis set) to 10.0 (=very large basis set), in steps of 0.5 (change the value in case.in1, run, save, change again, etc.) Then examine basis set size, total energy and electric-field gradient by grepline for :RKM, :ENE and :EFG001:

```
hcp-Fe_RKM55_k0240.scf::RKM : MATRIX SIZE 61LOs: 26 RKM= 5.03 WEIGHT= 1.00 PGR:
hcp-Fe_RKM60_k0240.scf::RKM : MATRIX SIZE 69LOs: 26 RKM= 5.92 WEIGHT= 1.00 PGR:
hcp-Fe_RKM65_k0240.scf::RKM : MATRIX SIZE 81LOs: 26 RKM= 6.20 WEIGHT= 1.00 PGR:
hcp-Fe_RKM70_k0240.scf::RKM : MATRIX SIZE 111LOs: 26 RKM= 6.97 WEIGHT= 1.00 PGR:
hcp-Fe_RKM75_k0240.scf::RKM : MATRIX SIZE 125LOs: 26 RKM= 7.38 WEIGHT= 1.00 PGR:
hcp-Fe_RKM80_k0240.scf::RKM : MATRIX SIZE 137LOs: 26 RKM= 7.77 WEIGHT= 1.00 PGR:
hcp-Fe_RKM85_k0240.scf::RKM : MATRIX SIZE 161LOs: 26 RKM= 8.14 WEIGHT= 1.00 PGR:
hcp-Fe_RKM90_k0240.scf::RKM : MATRIX SIZE 173LOs: 26 RKM= 8.79 WEIGHT= 1.00 PGR:
hcp-Fe_RKM95_k0240.scf::RKM : MATRIX SIZE 223LOs: 26 RKM= 9.46 WEIGHT= 1.00 PGR:
hcp-Fe_RKM10_k0240.scf::RKM : MATRIX SIZE 259LOs: 26 RKM= 9.84 WEIGHT= 1.00 PGR:
```

The largest basis set is 4 times larger than the smallest basis set, and as the calculation time scales with the third power of the basis set size, the last case takes $4^3 = 64$ times more time than the first one.
Figure 7.3: Top left: Total energy (left axis) and EFG (right axis) for hcp-Fe as a function of basis set size. Top right: the same EFG data, expressed as the log of the difference with the value at the largest basis set size. Bottom left: the difference between the total energies of two hcp-Fe unit cells that differ 5% in volume. Bottom right: logarithm of the difference between the latter quantity and its value at the largest basis set size (whence, ”difference difference” in the vertical axis is not a typo).

We knew that with 240 k-points the electric-field gradient is known with an accuracy of about 0.05. That means it does not make sense to ask for a basis set that offers a better accuracy. Even a small basis set as RKM=6.5 offers this accuracy. Therefore, initialization with \(-kmax 6.5 -numk 3840\) is sufficient to obtain an electric-field gradient that is within 0.05 of the exact

---

13One could go even as low as RKM=6.0, but as this value is on the threshold (RKM=5.5 goes nuts), it is safer to take a slightly larger value.
PBE-value \(^{14}\). If you rather want to have the field gradient converged up to 0.01, then an initialization with `-rmax 8.5 -numk 30720` is more appropriate.

The total energy as a function of basis set size, is given hereafter and also in Fig. 7.3:

\[
\begin{align*}
\text{hcp-Fe\_RKM55\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5090.96714470} \\
\text{hcp-Fe\_RKM60\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.07322057} \\
\text{hcp-Fe\_RKM65\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.13862964} \\
\text{hcp-Fe\_RKM70\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.17609664} \\
\text{hcp-Fe\_RKM75\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.1957998} \\
\text{hcp-Fe\_RKM80\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20377312} \\
\text{hcp-Fe\_RKM85\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20759715} \\
\text{hcp-Fe\_RKM90\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.21033657} \\
\text{hcp-Fe\_RKM95\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.21072273} \\
\text{hcp-Fe\_RKM10\_k0240.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.21072273} \\
\end{align*}
\]

In contrast to all previous series, the latter table does not show a damped scatter, but a clearly monotonously evolving function. If we would apply the same argument as for the electric-field gradient, then the conclusion would be that even an extremely large basis set equivalent to RKM=10.0 would not warrant the same 0.0001 Ry accuracy as a medium-sized k-mesh of 240 points. The reason for this lies in the variational principle behind the Kohn-Sham equations: with a finite basis set, the numerically found total energy is an upper bound for the exact total energy. The larger the basis set, the better you approximate the exact value, i.e. the lower the numerically found total energy will be. Very often, however, we do not need absolute accuracy of the total energy. Many properties depend on differences \(^{15}\) between total energies, or on derivatives \(^{16}\) of total energies. Such properties are insensitive to a rigid shift of all total energies. What one really needs, is a basis set that is sufficiently large to know energy differences up to a desired accuracy. To that end, reduce the volume of the hcp-Fe unit cell by, say, 5% \((a=b=4.594069 \text{ a.u.}, c=7.360171 \text{ a.u.})\), and determine once again the total energy as a function of basis set size:

\[
\begin{align*}
\text{hcp-Fe\_RKM55\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5090.97640700} \\
\text{hcp-Fe\_RKM60\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.07687257} \\
\text{hcp-Fe\_RKM65\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.14030136} \\
\text{hcp-Fe\_RKM70\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.17634711} \\
\text{hcp-Fe\_RKM75\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.19470946} \\
\text{hcp-Fe\_RKM80\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20332407} \\
\text{hcp-Fe\_RKM85\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20707643} \\
\text{hcp-Fe\_RKM90\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20887104} \\
\text{hcp-Fe\_RKM95\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.20972091} \\
\text{hcp-Fe\_RKM10\_k0240\_volm05.scf::ENE} & : ****** TOTAL ENERGY IN Ry = \text{-5091.21010215} \\
\end{align*}
\]

This yields another monotonously behaving curve, that asymptotically approaches the exact total energy. Finally, plot the difference between both total energies as function of basis set size (bottom halve of Fig. 7.3): from \(RK_{\max}=8.5\) onwards, the change in energy differences is smaller than 0.0001 Ry. Initializing by `-rmax 8.5 -numk 3840` therefore guarantees that energy differences are accurate up to 0.1 mRy for one unit cell of non-magnetic hcp-Fe.

\(^{14}\)Indeed, the basis set size and the k-mesh density determine the numerical accuracy only. With a very large basis set and very dense k-mesh, you will obtain the exact value for the particular exchange-correlation functional that was taken. It does not guarantee that the numerically converged result will be close to the exact value of a property as it appears in Nature. See http://dx.doi.org/10.1080/10408436.2013.772503, http://arxiv.org/abs/1204.2733 or http://molmod.ugent.be/deltacodesDFT for an elaboration of this issue.

\(^{15}\)For instance, the equilibrium volume of a unit cell determined from the minimum of a E(V)-curve will not change if the entire curve is vertically shifted by a constant value.

\(^{16}\)Forces, for instance.
7.2.3 \( R_{mt}K_{max}-\text{values per element: overview} \)

The wien2k FAQ at http://www.wien2k.at/reg_user/faq/rkmax.html describes a procedure that offers a good hint for the \( R_{mt}K_{max} \)-value that is needed for a fair calculation, and this without having to do test calculations. This procedure is repeated here, after which we will compare its prediction with what we found before for hcp-Fe by explicit testing. The procedure goes as follows:

1. Take your case.struct, and have setrmt_lapw determine the muffin tin radii (setrmt_lapw case, then cp case.struct setrmt case.struct).
2. Identify the element in your newly obtained case.struct that has the smallest muffin tin sphere.
3. Look up this element in the table at http://www.wien2k.at/reg_user/faq/rkmax.html, and use the corresponding value of \( R_{mt}K_{max} \) for your calculation.

For our example of hcp-Fe, setrmt_lapw would suggest muffin tin radii of 2.29 a.u. for all Fe atoms, and obviously Fe is the element with the lowest muffin tin radius in this unit cell. The table suggest \( R_{mt}K_{max}=8.0 \) for Fe. Our explicit tests in the previous paragraph suggest a slightly larger value (8.5), which is not surprising as we asked for a rather high accuracy of 0.1 mRy/unit cell for energy differences.

This quick procedure is not meant as a substitute for explicit testing, but it is useful as it shows right away in which region you should test.

7.2.4 Translation of tests to other cells

Explicit tests as in Sec. 7.2.2 are usually performed only for small, representative cases. How can we translate the conclusions to the larger cells that are likely to be the actual topic of your computational project?

7.2.4.1 k-mesh

We had concluded that 240 k-points in the irreducible part of the Brillouin zone (or 3840 k-points in the full Brillouin zone) was the size of k-mesh we need for non-magnetic hcp-Fe. Imagine, however, that our goal is not to study hcp-Fe as such, but Fe-Ni alloys with a few at.% of Ni. A first model for this could be a \( 3 \times 3 \times 2 \) supercell (36 atoms), where one of the Fe atoms is replaced by Ni. How do you specify a k-mesh for this supercell that has the same quality (density) as for the small cell?

The k-mesh samples the first Brillouin zone, and the latter is inversely proportional to the size of the unit cell. If the unit cell is large in one direction, this will correspond to a narrow edge of the Brillouin zone, and vice-versa. As a consequence, a narrow edge of the Brillouin zone will require less k-points for the same sampling density compared to a longer edge. As hcp-Fe has a unit cell with two ‘short’ edges (\( a=b=4.673293 \) a.u.) and one ‘long’ edge (\( c=7.487095 \) a.u.), its Brillouin zone has two long edges and one short one. The latter is reflected in the k-mesh for hcp-Fe, which we can read at the first line of case.klist:

\[
3840 \text{ k, div: ( 19 19 10)}
\]
A $19 \times 19 \times 10$-mesh: 19 k-points along each long edge, only 10 along the short edge. This points the way to the choice of k-mesh for a supercell: making the ‘a’ lattice parameter three times as long will make the corresponding Brillouin zone edge three times shorter – therefore, three times less k-points along that direction will be needed. For a $3 \times 3 \times 2$-supercell, a $19/3 \times 19/3 \times 10/2 = 6 \times 6 \times 5$ mesh would be sufficient.

Generate first the supercell by `x supercell` (no shift, no vacuum). Replace one of 36 Fe-atoms by a Ni atom, and do

```
x sgrou
```

```
cp case.struct_sgroup case.struct
```

Your `case.struct` now has the following lattice parameters and cell angles:

```
14.019879 14.019879 14.974190 90.000000 90.000000 120.000000
```

which is indeed $3 \times 3 \times 2$. Run now `x kgen` and try various numbers of k-points until you have about a $6 \times 6 \times 5$ mesh. You’ll get it for 200 points in the full Brillouin zone (or 90 in the irreducible part). This will provide you with the same quality as the $19 \times 19 \times 10$-mesh in the original hcp-Fe cell.

### 7.2.4.2 $R_{mt}K_{\text{max}}$

In contrast to the k-mesh, $R_{mt}K_{\text{max}}$ should often be left unchanged when going to supercells or when modifying the structure. The general rule is: if you want to compare different calculations with each other, then the muffin tin radii for identical elements should be identical in both calculations. Only if in one of the calculations the minimal muffin tin radius is different (because it lacks a particular element), then $R_{mt}K_{\text{max}}$ has to be adapted. An instructive example for MgO and Mg(OH)$_2$ is given at [http://www.wien2k.at/reg_user/faq/rkmax.html](http://www.wien2k.at/reg_user/faq/rkmax.html).

$R_{mt}K_{\text{max}}$ is a fairly robust indicator of the quality of a basis set for a particular element. That is the reason why a table as at [http://www.wien2k.at/reg_user/faq/rkmax.html](http://www.wien2k.at/reg_user/faq/rkmax.html) is possible at all: one value per element, regardless of the crystal structure and the specific $R_{mt}$ value. This is different from plane wave + pseudopotential codes, where the basis set quality is determined by $K_{\text{max}}$, not by $R_{mt}K_{\text{max}}$. The physical explanation for this difference was formulated by John E. Pask in the wien2k mailing list as follows \(^{17}\): \textit{(F)or a given atom, the wavefunctions become more rapidly oscillatory, in both angular and radial directions, as one approaches the nucleus. Hence, a fixed $RK_{\text{max}}$ naturally increases/decreases) planewave resolution ($K_{\text{max}}$) at $R_{mt}$ as $R_{mt}$ is decreased/increased, as needed to accommodate the more/less rapid oscillations in the wavefunctions at $R_{mt}$. Are the oscillations exactly twice as rapid at 1/2 the $R_{mt}$? Certainly not. So the scaling is not perfect. But it provides a far better indicator of convergence, as a single number, for a range of sphere sizes, than $K_{\text{max}}$ alone.}

\[^{17}\text{http://www.mail-archive.com/wien@zeus.theochem.tuwien.ac.at/msg08818.html}\]
Appendix A

Fourier transforms, plane waves, the reciprocal lattice and Bloch’s theorem

A.1 Fourier transforms

Most often, one’s first encounter with Fourier transforms is in the context of time dependent functions. The Fourier transform $\mathcal{F}$ of a function $f(t)$ is a function $F(\omega)$ in frequency domain, and is defined as:

$$F(\omega) = \mathcal{F}\{f\} = \int_{-\infty}^{\infty} f(t)e^{-i\omega t}dt$$  \hspace{1cm} (A.1)

$$f(t) = \mathcal{F}^{-1}\{F\} = \frac{1}{2\pi} \int_{-\infty}^{\infty} F(\omega)e^{i\omega t}d\omega$$ \hspace{1cm} (A.2)

If one would like to construct $f(t)$ by adding different functions $e^{i\omega t}$, then $F(\omega)$ tells what the weight of each function is. For instance, for $f(t) = \cos(\omega_0 t)$ it can be shown that:

$$F(\omega) = \frac{1}{2}\delta(\omega_0 + \omega) + \frac{1}{2}\delta(\omega_0 - \omega)$$ \hspace{1cm} (A.3)

(see fig. A.1) Therefore, $\cos(\omega_0 t)$ must be equal to a sum of two $e^{i\omega t}$-functions, each with a weight given above:

$$\cos(\omega_0 t) = \frac{1}{2}e^{i(-\omega_0)t} + \frac{1}{2}e^{i\omega_0 t}$$ \hspace{1cm} (A.4)

which indeed is true, considering the definition:

$$e^{i\omega t} = \cos(\omega t) + i\sin(\omega t)$$ \hspace{1cm} (A.5)

In this example, there is a discrete set of $\omega$-values where $F(\omega)$ is non-zero. In general, this is true for any periodic function: its Fourier transform is nonzero at a discrete set of frequencies (although there might be an infinite number of frequencies in the set). A periodic function can therefore be written as a sum of functions $e^{i\omega t}$. If $f(t)$ is aperiodic, $F(\omega)$ is non-zero over a continuous and often infinite range, and $f(t)$ can therefore be written only as an integral over $e^{i\omega t}$.
A.2 Plane waves

The concept of Fourier transforms can be translated for use with functions $f(\vec{r})$ in real space. The role of $\omega$ will be taken over by so-called reciprocal vectors $\vec{g}$ that have as dimension $1/\text{length}$. In the same way as the infinite amount of position vectors continuously fills real space, so the infinite amount of reciprocal vectors continuously fills reciprocal space. Fourier and inverse Fourier transforms between real and reciprocal space are now defined as:

$$F(\vec{g}) = \mathcal{F}\{f\} = \int f(\vec{r}) e^{-i\vec{g} \cdot \vec{r}} d^3\vec{r}$$

(A.6)

$$f(\vec{r}) = \mathcal{F}^{-1}\{F\} = \frac{-1}{2\pi} \int F(\vec{g}) e^{i\vec{g} \cdot \vec{r}} d^3\vec{g}$$

(A.7)

A particular kind of function in real space is a plane wave, defined as:

$$f(\vec{r}) = e^{i\vec{g}_0 \cdot \vec{r}}$$

(A.8)

with $\vec{g}_0$ any vector in reciprocal space. Convince yourself that a plane wave is periodic in space with wave vector $\vec{g}_0$ (= it takes the same value if you go along a distance $2\pi/|\vec{g}_0|$ parallel with $\vec{g}_0$), and that for any plane perpendicular to $\vec{g}_0$ it has the same value everywhere in that plane. The Fourier transform of a plane wave is non-zero at a single point $\vec{g}_0$ in reciprocal space only:

$$F(\vec{g}) = \int e^{i(\vec{g}_0 - \vec{g}) \cdot \vec{r}} d^3\vec{g}$$

(A.9)

$$= \delta(\vec{g}_0 - \vec{g})$$

(A.10)

which is very logical, as we need only one function $e^{i\vec{g}_0 \cdot \vec{r}}$ at $\vec{g} = \vec{g}_0$ with weight 1 to ‘build’ $e^{i\vec{g}_0 \cdot \vec{r}}$. The shorter the period of the plane wave in real space, the further away the point indicated by $\vec{g}_0$ lies from the origin of reciprocal space.

Functions that are periodic in real space, will have a Fourier transform that is non-zero only at discrete points in reciprocal space (possibly an infinite number of points). For aperiodic functions in real space, the Fourier transform will be non-zero over a continuous volume of reciprocal space.
A.3 The reciprocal lattice

So far, we did not make any distinction between different points in real or reciprocal space\textsuperscript{1}, none of them were ‘special’. Now, we introduce an infinite and regular (Bravais) lattice of special points in real space, which we call the real lattice or direct lattice. As an example, take the following basis vectors:

\[
\vec{a} = (5, 0, 0) \quad |\vec{a}| = 5 \\
\vec{b} = (5\sqrt{2}, 5\sqrt{2}, 0) \quad |\vec{b}| = 10 \\
\vec{c} = (0, 0, 1) \quad |\vec{c}| = 1
\]  

(A.11)

\[
\vec{a}^* = \frac{2\pi}{\vec{a} \cdot (\vec{b} \times \vec{c})} \quad a^* = \frac{2\pi}{5\sqrt{2}} (\frac{\sqrt{2}}{2}, -\frac{\sqrt{2}}{2}, 0) \quad \frac{2\pi}{|a^*|} = \frac{5}{\sqrt{2}} \approx 3.536
\]  

(A.14)

\[
\vec{b}^* = \frac{2\pi}{\vec{b} \cdot (\vec{c} \times \vec{a})} \quad b^* = \frac{2\pi}{10\sqrt{2}} (0, 1, 0) \quad \frac{2\pi}{|b^*|} = \frac{10}{\sqrt{2}} \approx 7.071
\]  

(A.15)

\[
\vec{c}^* = \frac{2\pi}{\vec{c} \cdot (\vec{a} \times \vec{b})} \quad c^* = 2\pi (0, 0, 1) \quad \frac{2\pi}{|c^*|} = 1
\]  

(A.16)

A cut along the XY-plane of this lattice is shown in fig. A.2-a. Define now three basis vectors to construct a lattice of special points in reciprocal space in the following way:

\[a^* = 2\pi \frac{\vec{b} \times \vec{c}}{\vec{a} \cdot (\vec{b} \times \vec{c})} \quad a^* = \frac{2\pi}{5\sqrt{2}} \left(\frac{\sqrt{2}}{2}, -\frac{\sqrt{2}}{2}, 0\right) \quad \frac{2\pi}{|a^*|} = \frac{5}{\sqrt{2}} \approx 3.536\]

\[b^* = 2\pi \frac{\vec{c} \times \vec{a}}{\vec{b} \cdot (\vec{c} \times \vec{a})} \quad b^* = \frac{2\pi}{10\sqrt{2}} (0, 1, 0) \quad \frac{2\pi}{|b^*|} = \frac{10}{\sqrt{2}} \approx 7.071\]

\[c^* = 2\pi \frac{\vec{a} \times \vec{b}}{\vec{c} \cdot (\vec{a} \times \vec{b})} \quad c^* = 2\pi (0, 0, 1) \quad \frac{2\pi}{|c^*|} = 1\]

(A.16)

This definition depends on the details of the real lattice. The lattice of special points in reciprocal space defined by these basis vectors is called the reciprocal lattice corresponding to this particular real lattice. We will indicate these special points by $\vec{K}$, reciprocal lattice vectors. In Fig. A.2-a a cut through the XY-plane of our real and reciprocal lattices are shown superimposed on each other\textsuperscript{2}, Fig. A.2-b shows a detail of the reciprocal lattice close to the origin. As for any point in reciprocal space, also the points of the reciprocal lattice correspond to plane waves in real space. Fig. A.2-c shows wave fronts of the three waves corresponding to $\vec{a}^*$, $3\vec{a}^*$ and $\vec{a}^* + \vec{b}^*$. Apparently, the plane waves corresponding to the reciprocal lattice have a special property: they are commensurate with the real lattice. This means: if you are on a particular wave front and you make a displacement along any $n\vec{a} + m\vec{b} + p\vec{c}$, then you end up on a wavefront with exactly the same value. An alternative definition of the reciprocal lattice could therefore be: the collection of all points in reciprocal space that correspond to plane waves that are commensurate with the real lattice.

\textsuperscript{1}Be careful in the remainder of this section not to confuse the words space and lattice: they have clearly distinct meanings!

\textsuperscript{2}Note that such a plot depends on the length unit chosen: if $a = 5$ mm in real space, $a^* = (2\pi)/(5/\sqrt{2}) \approx 1.78$ mm\textsuperscript{-1} in reciprocal space. In Fig. A.2-a, this would be represented as 1.78 mm. But if $a$ would be measured in meter, $a = 0.05$ m would correspond to 177.7 m. This at first sight strange behaviour is due to the fact that we want to display the reciprocal lattice with inverse length as dimension into our real world where we have to draw lines with real lengths. Fig. A.2 is made with an arbitrary length unit in which $\vec{a}$ has length 5 and $\vec{b}$ has length 10. Anyway, this affects only the appearance of the plot and not the physical results: whatever units are used, the wavelengths (which have dimension of length) are found to be the same.
Figure A.2: 

a) Lattice and reciprocal lattice. 

b) The first Brillouin zone. 

b*) Wave fronts of plane waves that are commensurate with the real lattice.
As for any lattice, several equivalent primitive unit cells can be defined for the reciprocal lattice. One such choice is a cell that contains all points in reciprocal space that are closer to the origin than to any other point of the reciprocal lattice (see Fig. A.2-b). This particular choice of primitive unit cell is called the first Brillouin zone. Analogously, the second Brillouin zone is defined as containing all points that have the origin of the reciprocal lattice as their second nearest neighbour. Any point in reciprocal space can uniquely be assigned to some $n^{th}$ Brillouin zone. By definition of the concept ‘unit cell’, any point (in the $n^{th}$ Brillouin zone) can be connected to a point in the first Brillouin zone by a vector of the reciprocal lattice.

A.4 Bloch’s theorem

As is explained in B, the (discrete) set of eigenfunctions of a quantum-mechanical Hamiltonian can always be labeled by quantum numbers. Bloch’s theorem tells how this can be done for a Hamiltonian that has a lattice periodicity. The theorem says: Any eigenfunction $\psi(\vec{r})$ can be written as a product of a function $u_{\vec{g}}(\vec{r})$ that has the periodicity of the lattice, and a plane wave $e^{i\vec{g} \cdot \vec{r}}$ with $\vec{g}$ any $3$-vector in reciprocal space.

$$\psi(\vec{r}) = u_{\vec{g}}(\vec{r}) e^{i\vec{g} \cdot \vec{r}}$$  \hspace{1cm} (A.17)

As there are an infinite number of vectors in reciprocal space, there are an infinite number of eigenstates of such a Hamiltonian. The reciprocal vectors $\vec{g}$ serve as labels (quantum numbers) of the eigenstates, and we could therefore rename $\psi(\vec{r})$ into $\psi_{\vec{g}}(\vec{r})$.

Usually this is not done exactly in this way however. Every $\vec{g}$ can be written as the sum of a vector in the first Brillouin zone (we call this $\vec{k}$) and a reciprocal lattice vector $\vec{K}$:

$$\vec{g} = \vec{k} + \vec{K} \hspace{1cm} (A.18)$$

Bloch’s theorem can now be rewritten as follows:

$$\psi_{\vec{g}}(\vec{r}) = \left\{ u_{\vec{g}}(\vec{r}) e^{i\vec{K} \cdot \vec{r}} \right\} e^{i\vec{k} \cdot \vec{r}} \hspace{1cm} (A.19)$$

Because $e^{i\vec{K} \cdot \vec{r}}$ is commensurate with the lattice, the function between brackets still has the periodicity of the lattice. We could rename it into $u_{\vec{k}}^n(\vec{r})$, where $n$ indicates the number of the Brillouin zone where $\vec{g}$ was in. Indeed, $n$ and $\vec{k}$ contain the same information as $\vec{g}$, and can therefore be used as an alternative way of labeling. For $n = 1$, $\vec{k}$ and $\vec{g}$ are identical. For $\vec{g}$ in the second Brillouin zone, we reuse the same set of $\vec{k}$, but $n$ is increased to 2, etc. For each $\vec{k}$, an infinite number of $n$ is possible. The parameter $n$ is called the band index.

---

3Strictly spoken only for an infinitely large lattice every $\vec{g}$ is allowed. If a finite crystal and periodic boundary conditions are used, the allowed $\vec{g}$ form a dense mesh in reciprocal space: in the first Brillouin zone, there are as many allowed $\vec{g}$ as there are unit cells contained in the finite (real) lattice. For a crystal of macroscopic size, this is a huge number. For any $\vec{g}$ there will be an allowed $\vec{g}$ very nearby, and therefore we can for practical purposes assume that all $\vec{g}$ are allowed. The exact property that distinguishes the allowed $\vec{g}$ from the not-allowed ones, is that they correspond to plane waves that are commensurate with the boundaries of the finite crystal.

4Pay attention to the consistent names that are use: $\vec{g}$ for any reciprocal vector, $\vec{K}$ for a reciprocal lattice vector and $\vec{k}$ for a vector in the first Brillouin zone. Other texts may use other notations.
Bloch’s theorem can now be restated in its most frequently used form: \( \psi_{n\vec{k}}(\vec{r}) \) can be written as a product of a function \( u_{n\vec{k}}(\vec{r}) \) that has the periodicity of the lattice, and a plane wave \( e^{i\vec{k} \cdot \vec{r}} \) with \( \vec{k} \) any vector in the first Brillouin zone:

\[
\psi_{n\vec{k}}(\vec{r}) = u_{n\vec{k}}(\vec{r}) e^{i\vec{k} \cdot \vec{r}} \tag{A.20}
\]

Actually we have now split off a known part \( e^{i\vec{k} \cdot \vec{r}} \) from the eigenstate, such that only the unknown remainder \( u_{n\vec{k}}(\vec{r}) \) has to be determined. An important advantage is that this part is now known to have the periodicity of the lattice (which is not true for \( \psi_{n\vec{k}}(\vec{r}) \) itself). If a plane wave basis set is used, it can therefore be written as a sum (no integral needed!) over plane waves that have this same periodicity, and these precisely are the plane waves corresponding to reciprocal lattice vectors:

\[
u_{n\vec{k}}(\vec{r}) = \sum_{\vec{K}} c_{n\vec{k},\vec{K}} e^{i\vec{K} \cdot \vec{r}} \tag{A.21}
\]

The expansion of \( \psi_{n\vec{k}}(\vec{r}) \) in the same basis is then

\[
\psi_{n\vec{k}}(\vec{r}) = \sum_{\vec{K}} c_{n\vec{k},\vec{K}} e^{i(\vec{k} + \vec{K}) \cdot \vec{r}} \tag{A.22}
\]

and what have to be searched are the coefficients \( c_{n\vec{k},\vec{K}} \).
Appendix B

Quantum numbers and the Density Of States

B.1 Familiar examples

In quantum mechanics, every physical situation that one wants to study is completely defined by its Hamiltonian $\hat{H}$. Every stationary solution of such a problem is described by a state $\psi_{\kappa}$ that is an eigenstate of that Hamiltonian. The energy of that solution is the corresponding eigenvalue $E_{\kappa}$:

$$\hat{H}\psi_{\kappa} = E_{\kappa}\psi_{\kappa}$$  \hspace{1cm} (B.1)

This is called the (time-independent) Schrödinger equation. In a physical situation, there are always boundary conditions. This limits the possible eigenvalues, and only a discrete but infinite set $E_{\kappa}$ remains. The symbols $\kappa$ stands for one or more quantum numbers that are used to label the different eigenfunctions and eigenvalues that satisfy the boundary conditions.

Some well-known examples:

- **The one-dimensional harmonic quantum oscillator**
  A (light) particle with mass $M$ moves in a one-dimensional harmonic potential $V(x) = Cx^2/2$. The Schrödinger equation is:

$$\left( -\frac{\hbar^2}{2M} \frac{d^2}{dx^2} + \frac{Cx^2}{2} \right) \psi_n(x) = E_n\psi_n(x)$$  \hspace{1cm} (B.2)

The boundary condition is that the particle is bound: its probability to appear at $x \to \infty$ is zero. The discrete set of eigenvalues $E_n$ and eigenfunctions $\psi_n(x)$ can be labeled by a single quantum number $n$ ($n = 0, 1, 2, \ldots$) and are

$$E_n = \left(n + \frac{1}{2}\right)\hbar\nu$$  \hspace{1cm} (B.3)

$$\psi_n(x) = \sqrt{\frac{1}{\sqrt{\pi}2^n n!}} e^{-\xi^2/2} H_n(\xi)$$  \hspace{1cm} (B.4)
The vibration frequency $\nu$ is fully determined by $C$ and $M$. The variable $\xi = \sqrt{\alpha} x$, with $\alpha$ fully determined by $C$ and $M$. $H_n(\xi)$ is the Hermite polynomial of order $n$.

Fig. B.1 shows the eigenvalues as a function of the quantum number $n$, which is called the eigenvalue spectrum. In the same picture, the following function of the energy is shown:

$$g(E) = \sum_n \delta(E - E_n)$$

(B.5)

Here $\delta$ is the Dirac delta function. The function $g(E)$ is called the Density Of States (DOS). Whenever $E$ equals an eigenvalue, the DOS is non-zero. If two eigenvalues would be degenerate (which does not happen in this example), the DOS would be twice as high for that energy. If in some energy region many eigenvalues would exist, the DOS would be non-zero at many places in that region. Obviously, the eigenvalue spectrum contains the full physical information about the problem. The DOS contains less information than the eigenvalue spectrum, but carries still some of the characteristics of the problem. The DOS will serve as an easy ‘fingerprint’ of the problem in cases where the eigenvalue spectrum is difficult to visualize (see soon).

- A single electron orbiting a given nucleus with infinite mass

This case is treated in Appendix D. The Hamiltonian is given in equations D.1 and D.2. Boundary conditions for the 3 separate differential equations are that the solutions are single-valued (for $\Phi$), everywhere finite (for $\Theta$) and vanishing for $r \to \infty$ (for $u$). Eigenfunctions and eigenvalues are labeled by the quantum numbers $n$, $\ell$ and $m$.

Fig. B.2 shows the eigenvalue spectrum, that is clearly different from the case of the harmonic oscillator. For a given $n$, there are $n$ values of $\ell$ ranging from 0 to $n - 1$ that are degenerate. Each of these $\ell$ itself is $2\ell + 1$-fold degenerate due to $m$. That makes for each $n$ a $\sum_{\ell=0}^{n-1}(2\ell + 1)$-fold degeneracy, which is reflected in the increasing value of the DOS. The energies where the DOS is non-zero, are also not evenly spaced any more.

\[1\] Try to construct an imaginary eigenvalue spectrum that would yield exactly the same DOS as in fig. B.1.
A free particle in a one-dimensional box with length $L$

The Schrödinger equation is:

$$\left( -\frac{\hbar^2}{2M} \frac{d^2}{dx^2} \right) \psi_p(x) = E_p \psi_p(x)$$  \hspace{1cm} (B.6)

with $M$ the mass of the particle. The boundary condition is that the particle cannot leave the box. Eigenvalues and eigenstates are labeled by a quantum number $p$ that can take the values

$$p = \frac{2\pi n}{L}$$  \hspace{1cm} (B.7)

with $n$ any integer (negative, zero or positive). Eigenvalues and eigenstates are:

$$E_p = \frac{\hbar^2 p^2}{2M}$$  \hspace{1cm} (B.8)

$$\psi_p(x) = \frac{1}{\sqrt{L}} e^{ipx}$$  \hspace{1cm} (B.9)

The eigenvalue spectrum and the DOS are shown in fig. B.3. Apart from $\psi_{p=0}$, all states are doubly degenerate. When going to higher energies, the intervals between eigenvalues grow larger.

B.2 Crystalline solids

A general property of crystalline solids is that the potential due to the nuclei is periodic:

$$V(\vec{r} + \vec{R}) = V(\vec{r})$$  \hspace{1cm} (B.10)
Figure B.3: Left: Eigenvalues as a function of the label (quantum number) \( p \) for a particle in a one-dimensional box. Right: DOS as a function of energy for a particle in one-dimensional box (the picture is rotated over 90° in order to make comparison with the energy axis of the left picture easy).

with \( \vec{R} \) any vector of the Bravais lattice\(^2\). It can be shown that for the kinetic part of the Hamiltonian, this condition is always fulfilled. Therefore, if the potential is periodic, the total Hamiltonian is. We look now at the case where periodic boundary conditions are used for the macroscopic piece of crystal: an infinite solid is built by filling space periodically with copies of the crystal (fig. B.4). Under these conditions, eigenvalues and eigenfunctions of a periodic Hamiltonian can be labeled with the quantum ‘numbers’ \( n \) and \( \vec{k} \). \( n = 1, 2, 3, \ldots \) and \( \vec{k} \) is any vector in the first Brillouin zone that corresponds to a plane wave that is commensurate with the macroscopic crystal (fig. B.4). For every valid \( \vec{k} \), all values of \( n \) occur. The number of valid \( \vec{k} \)-vectors is equal to the number of unit cells in the macroscopic crystal. For real solids, this number is huge (order of magnitude \( 10^{23} \)). Therefore, the \( \vec{k} \)-vectors are very close to each other, and we can assume that they continuously fill the first Brillouin zone. Note the difference between the \( \vec{K} \)-vectors from section A.3 that build the reciprocal lattice: they correspond to plane waves that are commensurate with the unit cell of the crystal. The special \( \vec{k} \)-vectors that are selected inside the first Brillouin zone here, are commensurate with the entire macroscopic piece of crystal. Eigenfunctions are written as \( \psi_{n \vec{k}} \), eigenvalues as \( \epsilon_{\vec{k}}^n \) (or \( E_{\vec{k}}^n \)).

When we want to visualize the spectrum of eigenvalues, a problem is met: the eigenvalues are labeled by 4 independent numbers \( (n, k_x, k_y, k_z) \), such that 5 dimensions would be needed to make a plot as in the preceding examples. Obviously, this cannot be done. An alternative procedure is to select a path through the first Brillouin zone, and plot for every \( n \) the energies \( \epsilon_{\vec{k}}^n \) of the \( \vec{k} \) that happen to fall onto that path (see fig. B.5 for an example). The way the path is selected, is largely conventional. Straight intervals of the path are given conventional capital letters (often Greek) as names, and special points in the first Brillouin zone are given conventional capital letters (often Latin). Figures as fig. B.5 (in WIEN2K they are called ‘spaghetti plots’) represent most but not all essential information of the eigenvalue spectrum.

As in the previous cases, it is possible to plot the DOS here. The DOS is a function of the energy only, and hence not suffers from too high dimensions. Due to the presence of the

\(^2\)The Bravais lattice contains all essential information on the underlying symmetry of the crystal lattice.
Figure B.4: Periodic boundary conditions and wave fronts of plane waves that are commensurate with the entire crystal.

(almost) continuous extra quantum number $\vec{k}$, the definition of the DOS has to be extended:

$$g(E) = \frac{2}{V_{BZ}} \sum_n \int \delta(\epsilon - \epsilon_n^{\vec{k}}) d\vec{k}$$  \hspace{1cm} (B.11)

Here $V_{BZ}$ is the volume of the first Brillouin zone, and the integral is over the first Brillouin zone. The factor 2 is included to account explicitly for (degenerate) spin. An example of a DOS – which is a continuous function now due to $\vec{k}$ – is shown in fig. B.5. When there are flat regions in the spaghetti plot, many $\vec{k}$ have the same energy. This will show up as a peak in the DOS.
Figure B.5: Bottom: The first Brillouin zone of a bcc lattice, with 4 high symmetry points indicated: $\Gamma$, H, N and P. These points are connected by lines, labeled $\Delta$, G, $\Sigma$ and $\Lambda$. These 4 lines define a path through the Brillouin zone. Top left: the eigenvalues of all the $k$-vectors along the path are plotted, for some bcc compound (bcc-Fe, spin up). This so-called band structure plot contains less information than the full eigenvalues-versus-$\vec{k}$ plot, but the latter would require 5 dimensions. The fact that for every $\vec{k}$ multiple eigenvalues are possible, is due to different $n$. Top right: the DOS for the same material, plotted in such a way that the energy scale matches with the energy scale of the band structure plot. The DOS too contains less information than the full eigenvalues-versus-$\vec{k}$ plot, but compared to the band structure other information is omitted or included: the DOS is an integral over the full Brillouin zone (and not limited to a onedimensional path only), but on the other hand does not give information for a specific $\vec{k}$. Peaks in the DOS mean that many $\vec{k}$-vectors have the same eigenvalue. If the path goes through the region of the Brillouin zone that is responsible for this, we see relatively flat (horizontal) lines in the left picture, at the same energies where a peak is in the DOS. Four peaks are labeled 1-4 in the DOS, and are related to features of the band structure plot.
Appendix C

The eigenvalue problem

C.1 Eigenvalues and eigenvectors

Consider a vector space \((\mathbb{R}^n, n = 2)\) with a basis set \((\vec{e}_1, \vec{e}_2)\) that need not to be orthogonal. With respect to this basis, every vector \(\vec{x}\) can be uniquely characterized by two numbers \((x_1, x_2)\). With \(\hat{H}\), we denote an operator that can operate on \(\vec{x}\), resulting in a new vector \(\vec{y} = (y_1, y_2)\):

\[
\hat{H}\vec{x} = \vec{y}
\]  

(C.1)

Given \(\hat{H}\), we want to find all vectors \(\vec{x}\) that under \(\hat{H}\) transform into a vector that is parallel with itself:

\[
\hat{H}\vec{x} = \lambda \vec{x}
\]  

(C.2)

We rewrite this condition using the basis vectors, and will search for the \((x_1, x_2)\) that satisfy it:

\[
\hat{H}(x_1\vec{e}_1 + x_2\vec{e}_2) = \lambda (x_1\vec{e}_1 + x_2\vec{e}_2)
\]  

(C.3)

Left-multiply this equation with \(\vec{e}_1\):

\[
x_1 \vec{e}_1 \cdot (\hat{H}\vec{e}_1) + x_2 \vec{e}_1 \cdot (\hat{H}\vec{e}_2) = \lambda (x_1 \vec{e}_1 \cdot \vec{e}_1 + x_2 \vec{e}_1 \cdot \vec{e}_2)
\]  

(C.4)

Do the same with \(\vec{e}_2\):

\[
x_1 \vec{e}_2 \cdot (\hat{H}\vec{e}_1) + x_2 \vec{e}_2 \cdot (\hat{H}\vec{e}_2) = \lambda (x_1 \vec{e}_2 \cdot \vec{e}_1 + x_2 \vec{e}_2 \cdot \vec{e}_2)
\]  

(C.5)

Equations C.4 and C.5 can be summarized in matrix notation as

\[
\begin{bmatrix}
\vec{e}_1 \cdot (\hat{H}\vec{e}_1) & \vec{e}_1 \cdot (\hat{H}\vec{e}_2) \\
\vec{e}_2 \cdot (\hat{H}\vec{e}_1) & \vec{e}_2 \cdot (\hat{H}\vec{e}_2)
\end{bmatrix}
\begin{bmatrix}
x_1 \\
x_2
\end{bmatrix}
- \lambda
\begin{bmatrix}
\vec{e}_1 \cdot \vec{e}_1 & \vec{e}_1 \cdot \vec{e}_2 \\
\vec{e}_2 \cdot \vec{e}_1 & \vec{e}_2 \cdot \vec{e}_2
\end{bmatrix}
\begin{bmatrix}
x_1 \\
x_2
\end{bmatrix}
= \begin{bmatrix}
0 \\
0
\end{bmatrix}
\]  

(C.6)

The \(ij^{th}\) element (call it \(H_{ij}\)) of the first matrix is a number, as it is the result of a dot product. This matrix is completely determined if the action of \(\hat{H}\) on the basis vectors is known\(^1\). The

\(^1\)It is straightforward to show that with this matrix \(\hat{H}\vec{x} = \vec{y}\) can be expressed as

\[
\begin{bmatrix}
H_{11} & H_{12} \\
H_{21} & H_{22}
\end{bmatrix}
\begin{bmatrix}
x_1 \\
x_2
\end{bmatrix}
= \begin{bmatrix}
y_1 \\
y_2
\end{bmatrix}
\]  

(C.7)

If the \(H_{ij}\) are known, the effect of \(\hat{H}\) on any vector \(\vec{x}\) can be found by a simple matrix multiplication.
elements of the second matrix (call them \( S_{ij} \)) are determined by the basis only. This matrix is called the overlap matrix. With this notation, our 2 equations to find suitable \((x_1, x_2)\) become:

\[
\begin{pmatrix}
H_{11} - \lambda S_{11} & H_{12} - \lambda S_{12} \\
H_{21} - \lambda S_{21} & H_{22} - \lambda S_{22}
\end{pmatrix}
\begin{pmatrix}
x_1 \\
x_2
\end{pmatrix}
= 
\begin{pmatrix}
0 \\
0
\end{pmatrix}
\tag{C.8}
\]

The left matrix is fully known, apart from \( \lambda \) that is a parameter. For every value of \( \lambda \), equation C.8 can be solved for \( x_1 \) and \( x_2 \). For most values of \( \lambda \), the determinant of the matrix will be different from zero. Equation C.8 will have then one unique solution, which is obviously \((0, 0)\). This vector is parallel with the original \( \vec{x} \), but in a trivial way. Other, more interesting solutions can occur only for those values of \( \lambda \) that give a vanishing determinant:

\[
\begin{vmatrix}
H_{11} - \lambda S_{11} & H_{12} - \lambda S_{12} \\
H_{21} - \lambda S_{21} & H_{22} - \lambda S_{22}
\end{vmatrix}
= 0
\tag{C.9}
\]

The above equation is called the secular equation of \( \hat{H} \). It is a polynomial equation in \( \lambda \), with as highest power the dimension of the space \( n \) (2 in this case). The roots\(^2\) of the secular equation are called the eigenvalues of \( \hat{H} \). If \( \lambda = \lambda_1 \) is an eigenvalue, \([\hat{H} - \lambda_1 S][x] = 0\] has an infinite number of solutions: indeed, if \((x_1 = a, x_2 = b)\) is a solution, \((\beta a, \beta b)\) is one too, for \( \beta \) any real number. This conforms with our initial requirement: if \( \hat{H} \) transforms \( \vec{x} \) into a vector that is parallel with \( \vec{x} \), it will do the same with any vector \( \beta \vec{x} \). These vectors are called the eigenvectors of \( \hat{H} \), belonging to the eigenvalue \( \lambda_1 \). Usually a unit vector among them is chosen to represent this set of vectors. Eigenvectors belonging to different eigenvalues, can be shown to be perpendicular\(^3\). Therefore, unit vectors belonging to \( n \) different eigenvalues can be taken as an orthonormal basis \((\vec{e}_{\lambda_1}^1, \vec{e}_{\lambda_2}^2)\) for the vector space.

What will be the matrix representation of the operator \( \hat{H} \), written in this new basis? According to equation C.6, the matrix elements are:

\[
H_{ij} = \vec{e}_{\lambda_i}^i \cdot (\hat{H} \vec{e}_{\lambda_j}^j)
= \vec{e}_{\lambda_i}^i \cdot (\lambda_j \vec{e}_{\lambda_j}^j)
= \lambda_j \delta_{ij}
\tag{C.10}
\]

\[
[H] = \begin{bmatrix}
\lambda_1 & 0 \\
0 & \lambda_2
\end{bmatrix}
\tag{C.13}
\]

We could now try to find the eigenvalues and eigenvectors of this operator in the new basis, using the same procedure as described above. Obviously we will find the same \( \lambda_1 \) and \( \lambda_2 \) as before, with as eigenvectors \((1, 0)\) and \((0, 1)\) (i.e. the old eigenvectors which are now the basis vectors). This illustrates that eigenvalues and eigenvectors are intrinsic properties of the operator, and do not depend on the choice of the basis.

If the basis we originally started with was already orthonormal, then the overlap matrix will be a unit matrix. This simplifies equation C.6.

\(^2\)There will be at most \( n \) roots. In all cases that will be of interest for us, there will be exactly \( n \) roots.

\(^3\)At least if the matrix representing the operator is hermitian, which is the case for ‘well-behaved’ operators.
It could happen that 2 of the roots of the secular equation coincide. In that case, that eigenvalue leads to a plane instead of a line of eigenvectors. There is freedom then to choose two mutually perpendicular basis vectors in that plane, but still they are perpendicular to eigenvectors of the other – different – eigenvalues.

This idea can be extended to vector spaces of infinite dimension (e.g. function spaces, in which single particle wave functions are described).

C.2 Basis transformations

Consider an operator $\hat{A}$ with matrix representation $A^4$, in a vector space $\mathbb{R}^n$ with a normalized but not necessarily orthogonal basis $(\vec{e}_1^\alpha, \ldots, \vec{e}_n^\alpha)$. This operator transforms every vector into a new vector. Also the basis vectors are transformed, and in this way you can understand that the $j^{th}$ column of $A$ contains the coefficients that express the transformed $j^{th}$ basis vector in the original basis:

$$
\begin{bmatrix}
\cdots & a_{1j} & \cdots \\
\cdots & \vdots & \cdots \\
\cdots & a_{jj} & \cdots \\
\cdots & \vdots & \cdots \\
\cdots & a_{nj} & \cdots \\
\end{bmatrix}_{n\times n}
\begin{bmatrix}
e_{ij}^\alpha = 0 \\
e_{ij}^\alpha = 1 \\
e_{nj}^\alpha = 0 \\
\end{bmatrix}_{n\times 1}
= 
\begin{bmatrix}
a_{1j} \\
a_{jj} \\
a_{nj} \\
\end{bmatrix}_{n\times 1}
$$

(C.14)

This provides a practical recipe to find the matrix representation of an operator about which we want that it transforms the basis vectors $\vec{e}_j^\alpha$ in a new set $\vec{e}_j^\beta$: express each $\vec{e}_j^\beta$ in the old basis $\vec{e}_j^\alpha$, and these expansion coefficients form the columns of $A$. Note that in this way a non-orthogonal basis can be transformed into an orthogonal basis. If both bases are orthonormal, than the matrix $A$ has special properties (it is a unitary matrix).

Conversely, the $j^{th}$ column of the inverse matrix $A^{-1}$ contains the coefficients that express the vector $\vec{e}_j^\alpha$ of the old basis in the new basis of the $\vec{e}_j^\beta$. This is illustrated in fig. C.1 for an operator

$$
A = \begin{bmatrix}
\frac{\sqrt{2}}{2} & -\frac{\sqrt{2}}{2} \\
\frac{\sqrt{2}}{2} & \frac{\sqrt{2}}{2}
\end{bmatrix}
A^{-1} = \begin{bmatrix}
\frac{\sqrt{2}}{2} & \frac{\sqrt{2}}{2} \\
-\frac{\sqrt{2}}{2} & \frac{\sqrt{2}}{2}
\end{bmatrix}
$$

(C.15)

that rotates every vector of $\mathbb{R}^2$ counterclockwise over 45°. The old basis ($\vec{e}_1^\alpha$, $\vec{e}_2^\alpha$) is transformed in the new basis ($\vec{e}_1^\beta$, $\vec{e}_2^\beta$). The coordinates of these 4 vectors in both bases are:

<table>
<thead>
<tr>
<th></th>
<th>old</th>
<th>new</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\vec{e}_1^\alpha$</td>
<td>(1, 0)</td>
<td>$(\frac{\sqrt{2}}{2}, -\frac{\sqrt{2}}{2})$</td>
</tr>
<tr>
<td>$\vec{e}_2^\alpha$</td>
<td>(0, 1)</td>
<td>$(\frac{\sqrt{2}}{2}, \frac{\sqrt{2}}{2})$</td>
</tr>
<tr>
<td>$\vec{e}_1^\beta$</td>
<td>$(\frac{\sqrt{2}}{2}, \frac{\sqrt{2}}{2})$</td>
<td>(1, 0)</td>
</tr>
<tr>
<td>$\vec{e}_2^\beta$</td>
<td>$(-\frac{\sqrt{2}}{2}, \frac{\sqrt{2}}{2})$</td>
<td>(0, 1)</td>
</tr>
</tbody>
</table>

$^4\hat{A}$ is a general operator, not necessarily identical to the $\hat{H}$ we were discussing before. But everything that is said about $\hat{A}$, applies for $\hat{H}$ too.
Check with this table and with fig. C.1 that the following expression

\[ AX = Y \] (C.16)

\( X \) and \( Y \) are \((n \times 1)\)-matrices) can be interpreted in these two different ways:

- First interpretation: \( X \) contains the coefficients of any vector expressed in the old basis. Then \( Y \) will contain the coefficients of the transformation of \( X \) under \( \hat{A} \), expressed in the old basis.

- Second interpretation: \( X \) contains the coefficients of any vector expressed in the new basis. Then \( Y \) will contain the coefficients of this same vector, expressed in the old basis.

Similarly, this expression

\[ A^{-1}X = Y \] (C.17)

can be interpreted in two ways too:

- First interpretation: \( X \) contains the coefficients of the transformation under \( \hat{A} \) of any vector, expressed in the old basis. Then \( Y \) will contain the coefficients of this original vector (=before transformation) expressed in the old basis.

- Second interpretation: \( X \) contains the coefficients of any vector expressed in the old basis. Then \( Y \) will contain the coefficients of this same vector, expressed in the new basis.

### C.3 A practical procedure to find eigenvalues and eigenvectors

Given a vector space \( \mathbb{R}^n \) with a normalized but not necessarily orthogonal basis, we want to find the \( n \) eigenvalues and the \( n \) eigenvectors of a given operator \( \hat{H} \). The matrix elements \( H_{ij} \)
and $S_{ij}$ of the matrix representation $H$ of the operator and of the overlap matrix $S$ of the basis set, can be obtained as described before. Both are $n \times n$ matrices. What we want to determine are the coefficients $c_{ij}$ that express in the given basis the $j^{th}$ normalized eigenvector. In total, $n \times n$ coefficients $c_{ij}$ are needed to specify all eigenvectors, and they can be arranged in an $n \times n$ matrix $C$. Equation C.8 can then be extended to:

$$HC = SCE$$  \hspace{1cm} (C.18)

where $E$ is an $n \times n$ diagonal matrix that contains as its $j^{th}$ diagonal element the eigenvalue corresponding to the $j^{th}$ eigenvector. In equation C.18, the $n^2$ coefficients $c_{ij}$ are unknown, as are the $n$ eigenvalues $E_{jj}$. Equation C.18 represents $n^2$ independent equations, and $n^2 + n$ unknowns have to be determined. The remaining $n$ equations come from the requirement that in each class of eigenvectors, we search a normalized one:

$$[c_{1j} \ c_{2j} \ ... \ c_{nj}]_{1 \times n} \begin{bmatrix} c_{1j} \\ c_{2j} \\ \vdots \\ c_{nj} \end{bmatrix}_{n \times 1} = 1$$  \hspace{1cm} (C.19)

Equations C.18 and C.19 define the problem to be solved.

Now, we left-multiply equation C.18 by an $n \times n$ matrix $A^{-1}$. At two other places, we insert a multiplication by the unit matrix $II = AA^{-1}$:

$$A^{-1}HA_{\ H_0} A^{-1}C_{\ C_0} = A^{-1}SA_{\ S_0} A^{-1}C_{\ C_0} E$$  \hspace{1cm} (C.20)

It can be proven that for ‘well-behaved’ operators $\hat{H}$, always a matrix $A$ can be found such that $H_0 = A^{-1}HA$ is a diagonal matrix. We know that in such a case the elements on the diagonal must be the eigenvalues of $\hat{H}$. Consequently, $H_0 = E$. But if $H_0$ is diagonal, we know that it must be expressed in a basis that consists of its eigenvectors. Therefore, the operator $A$ must transform the original basis into the new basis of eigenvectors. Because $C$ contained the (unknown) coefficients of the eigenvectors in the old basis, $C_0 = A^{-1}C$ contains the coefficients of the eigenvectors in the basis of eigenvectors: $C_0$ must be diagonal. (We cannot say that $C_0$ should be the unit matrix: according to equation C.18 it should contain coefficients of eigenvectors, not of normalized eigenvectors.) Equation C.20 is identical to equation C.18, but written in the new basis of eigenvectors:

$$H_0C_0 = S_0C_0E$$  \hspace{1cm} (C.21)

with $H_0 = E$ and $C_0$ diagonal. Because the eigenvectors are mutually orthogonal, this new basis is orthonormal and hence $S_0$ is a unit matrix. The final expression

$$EC_0 = C_0E$$  \hspace{1cm} (C.22)

---

5Note that the coefficients that describe the $j^{th}$ normalized eigenvector form the $j^{th}$ column of $C$.

6Check for $n = 2$ that equation C.18 is indeed equivalent to twice equation C.8, once for both eigenvalues/eigenvectors.

7In an orthonormal basis ($S$ is a unit matrix), this reduces to the more familiar formula for the square of the length of a vector.
is true, because multiplication among diagonal matrices is commutative.

We conclude that if we can find the matrix $A$ that diagonalizes $H$, the eigenvalues are read from its diagonal, and the eigenvectors can be calculated from

$$C = AC_0$$  \hspace{1cm} (C.23)

where the $n$ undetermined values on the diagonal of $C_0$ are fixed by the still unused normalization requirements C.19. Diagonalizing the Hamiltonian matrix is the key to solve the eigenvalue problem. A lot of standard numerical procedures are available to perform a diagonalization. An important common feature is that the time it takes to execute them on a computer scales with $n^3$. 
Appendix D

Solutions of the radial part of the Schrödinger equation

Many text books on quantum mechanics explicitly solve the Schrödinger equation for a single electron with charge $-e$ and mass $m_e$ in the Coulomb potential of a nucleus with charge $Ze$ and infinite mass\footnote{An infinite mass avoids the rotation of nucleus and electron around their common center of mass.}. In spherical coordinates $\vec{r} = (r, \theta, \phi)$, the Schrödinger equation is:

$$-\frac{\hbar^2}{2m_e} \nabla^2 \psi(r, \theta, \phi) + V(r) = E \psi(r, \theta, \phi) \quad (D.1)$$

with

$$V(r) = -\frac{Ze^2}{4\pi\epsilon_0 r} \quad (D.2)$$

$$\nabla^2 = \frac{1}{r^2} \frac{\partial}{\partial r} \left( r^2 \frac{\partial}{\partial r} \right) + \frac{1}{r^2 \sin^2 \theta} \frac{\partial^2}{\partial \phi^2} + \frac{1}{r^2 \sin \theta} \frac{\partial}{\partial \theta} \left( \sin \theta \frac{\partial}{\partial \theta} \right) \quad (D.3)$$

This differential equation can be shown to be separable in 3 independent differential equations, each of them involving only $r$, $\theta$ or $\phi$ respectively. General solutions of equation D.1 will therefore be of the form

$$\psi(r, \theta, \phi) = u(r) \Theta(\theta) \Phi(\phi) \quad (D.4)$$

The three differential equations of which $\Phi(\phi)$, $\Theta(\theta)$ and $u(r)$ are solutions, are:

$$\frac{d^2\Phi}{d\phi^2} = -m^2\Phi \quad (D.5)$$

$$-\frac{1}{\sin \theta} \frac{d}{d\theta} \left( \sin \theta \frac{d\Theta}{d\theta} \right) + \frac{m^2\Theta}{\sin^2 \theta} = \ell(\ell + 1) \Theta \quad (D.6)$$

$$\frac{1}{r^2} \frac{d}{dr} \left( r^2 \frac{du}{dr} \right) + \frac{2m_e}{\hbar^2} (E - V(r)) u = \ell(\ell + 1) \frac{u}{r^2} \quad (D.7)$$

Here $\ell$ and $m$ are dimensionless parameters. In order to have well-behaved solutions, it can be shown that $\ell = 0, 1, 2, 3, \ldots$ and $m = -\ell, -\ell + 1, \ldots, \ell - 1, \ell$. The product of $\Theta(\theta)$ and $\Phi(\phi)$ is called a spherical harmonic $Y_{\ell m}(\theta, \phi)$:

$$Y_{\ell m}(\theta, \phi) = \Theta(\theta) \Phi(\phi) \quad (D.8)$$
Spherical harmonics can be found tabulated for many $(\ell, m)$ at various places. The full solution of equation D.1 can now be written as

$$\psi(r, \theta, \phi) = u_\ell(r) Y_\ell^m(\theta, \phi)$$  \hspace{1cm} (D.9)

with $u_\ell(r)$ a solution of equation D.7 with the proper value of $\ell$. Equation D.7 is the radial Schrödinger equation.

What will be the solutions of equation D.7? Using these transformations of variables:

$$\rho = 2\beta r$$  \hspace{1cm} (D.10)

$$\beta^2 = -\frac{2m_eE}{\hbar^2}$$  \hspace{1cm} (D.11)

$$\gamma = \frac{m_eZe^2}{\hbar^2\beta}$$  \hspace{1cm} (D.12)

it can be shown that the solutions have the form

$$u_\ell(r) = e^{-\rho/2} \rho^\ell G_\gamma(\rho)$$  \hspace{1cm} (D.13)

where $G_\gamma(\rho)$ is defined by the above relation. For the case of a free atom which we are considering here, the boundary condition for the radial equation is that $u_\ell(r) \to 0$ as $r \to \infty$ (or equivalently, as $\rho \to \infty$). This can be shown to be possible only if $\gamma = n$, with $n$ any of the integers

$$n = \ell + 1, \ell + 2, \ell + 3, \ldots$$  \hspace{1cm} (D.14)

For these specific values of $\gamma$, the $G_\gamma(\rho)$ – now written $G_n(\rho)$ – are called the associated Laguerre functions. The radial solutions $u(r)$ carry a double index $n\ell$ now ($u_{n\ell}(r)$), and the solutions of the entire Schrödinger equation D.1 are labeled by $n$, $\ell$ and $m$:

$$\psi_{n\ell m}(r, \theta, \phi) = u_{n\ell}(r) Y_\ell^m(\theta, \phi)$$  \hspace{1cm} (D.15)

The corresponding discrete set of energies (eigenvalues) is given by

$$E_{n\ell m} = E_n = -\frac{m_eZ^2e^4}{2\hbar^2n^2} \quad n = 1, 2, 3, 4, \ldots$$  \hspace{1cm} (D.16)

which is degenerate in $\ell$ and $m$.

In the APW, LAPW+LO and APW+lo methods, we will use radial functions $u_\ell(r)$ that are solutions of equation D.7. But we will need these functions only inside the muffin tin spheres. What their behaviour for $r \to \infty$ is, is not important for us: we will never use these functions in that region. The boundary condition that $u_\ell(r)$ should vanish at infinity therefore disappears. Hence, the requirement that $\gamma$ must be integer (leading to a discrete set of energies) disappears too. If no boundary condition is used, solutions $u_\ell(r)$ can be found for every energy $E$. Note that the subscript $n$ is dropped now, as it has no meaning without a boundary condition.
Appendix E

The homogeneous electron gas

The homogeneous electron gas, uniform electron gas or jellium model is an imaginary solid where all nuclear charge is homogeneously smeared out over space. This material is completely isotropic, and identical on every length scale. Therefore the electron density is constant:

\[ \rho(\vec{r}) \equiv \rho \equiv \frac{N}{V} \]  

(E.1)

with \( N \) the number of electrons in the material, and \( V \) its volume. The parameter \( \rho \) is the only thing we need to specify a particular homogeneous electron gas completely.

If the electrons do not interact, we are in the case of the free electron gas, which can be solved analytically in a straightforward way. The problem is much more difficult for an interacting electron gas. Here numerical calculations for the total energy are possible by quantum Monte Carlo. Subtracting the non-interacting kinetic energy and the Hartree energy gives a numerical result for the exchange-correlation energy \( \epsilon_{xc} \). If this is done for several densities \( \rho \), the function \( \epsilon_{xc}(\rho) \) is obtained. Note that \( \epsilon_{xc} \) is a function of \( \rho \), not a functional.
Appendix F

Functionals

F.1 Definition and Examples

To say that \( F \) is a functional of the function \( f(\tau) \) means that \( F \) is a number whose value depends on the form of the function \( f(\tau) \). In other words: for every function \( f \) (defined on the same domain, given by \( \tau \)) the functional \( F[f] \) returns a unique number:

\[
F : \mathcal{F} \mapsto \mathbb{R} : f \mapsto F[f]
\]  

(F.1)

(\( \mathcal{F} \) is a space of functions) A functional is therefore a ‘function of functions’. Some examples of functionals are:

- \( F[f] = \int_{-\infty}^{+\infty} f(\tau) e^{-\tau^2} d\tau \)

Clearly, if \( f \) is changed the value of the integral and therefore \( F[f] \) is changed. Note that \( F \) does not depend on the coordinates \( \tau \) of \( f \), only the form of \( f \) matters.

- \( F[f] = f(0) \)

This very simple functional depends only on the value of \( f \) at \( \tau = 0 \). Another \( f \) gives another value for \( F[f] \).

- The integral of a function is a functional. Another function gives another integral, and hence another value for the functional.

- The total energy of a quantum system is a functional of the density, as stated by the first theorem of Hohenberg and Kohn. The density is a function defined on \( \mathbb{R}^3 \).

F.2 Functional derivatives

The generalization of the derivative \( \frac{df(x)}{dx} \) of a function \( f(x) \), is the functional derivative \( \frac{\delta F[\rho]}{\delta \rho} \) of the functional \( F[\rho] \). As an example, we calculate here the functional derivative of the nucleus-electron part of the total-energy functional used in Chapter 1. The same method can be applied to functionals of other form too.

Consider a functional

\[
E_{V_{\text{ext}}}^{n-e}[\rho] = \int_{a}^{b} \rho(x) V_{\text{ext}}(x) \, dx
\]  

(F.2)
where $V_{\text{ext}}(x)$ is a given function. The integral runs over the 1-dimensional volume $[a, b]$ (it is straightforward to extend the following reasoning to functions $\rho(\vec{r})$ and $V_{\text{ext}}(\vec{r})$ and integrals over 3-dimensional volumes, as is needed in Chapter 1).

The variation of $E_{V_{\text{ext}}}^{n-e}[\rho]$ with respect to its argument $\rho(x)$ is given by

$$
\delta E_{V_{\text{ext}}}^{n-e}[\rho] = E_{V_{\text{ext}}}^{n-e}[\rho + \delta \rho] - E_{V_{\text{ext}}}^{n-e}[\rho] = \left( \int_a^b \rho(x) V_{\text{ext}}(x) \, dx + \int_a^b \delta \rho(x) V_{\text{ext}}(x) \, dx \right) - \int_a^b \rho(x) V_{\text{ext}}(x) \, dx \tag{F.4}
$$

$$
= \int_a^b \delta \rho(x) V_{\text{ext}}(x) \, dx \tag{F.5}
$$

In standard calculus, the differential $df$ of a function $f(\vec{x})$ defined on an $n$-dimensional domain, is given by the chain rule:

$$
df = \sum_{i=1}^{n} \frac{\partial f}{\partial x_i} \, dx_i \tag{F.6}
$$

In the calculus of functionals, the dimension of the domain ($\mathcal{F}$) is infinite. If one divides the interval $[a, b]$ into small enough sub-intervals defined by the set $\{x_1 = a, x_2, \ldots, x_{n-1}, x_n = b\}$ so that the variation $\delta \rho(x)$ is approximately constant in each interval, then equation F.5 is given approximately by

$$
\delta E_{V_{\text{ext}}}^{n-e}[\rho] \approx \sum_{i=1}^{p} V_{\text{ext}}(x_i) \left[ \delta \rho(x_i) \Delta x \right] \tag{F.7}
$$

In the limit of $p \to \infty$, equation F.5 is recovered. Comparing equations F.5 and F.6 one can identify

$$
dx_i \rightarrow \delta \rho(x) \, dx \tag{F.8}
$$

$$
\frac{\partial f}{\partial x_i} \rightarrow V_{\text{ext}}(x) = \frac{\delta E_{V_{\text{ext}}}^{n-e}[\rho]}{\delta \rho} \tag{F.9}
$$

Generalization to functionals

$$
E_{V_{\text{ext}}}^{n-e}[\rho] = \int_{\text{vol}} \rho(\vec{r}) V_{\text{ext}}(\vec{r}) \, d\vec{r} \tag{F.10}
$$

defined with functions $\rho(\vec{r}) : \mathbb{R}^3 \to \mathbb{R}$ leads to the following formula for the functional derivative:

$$
\frac{\delta E_{V_{\text{ext}}}^{n-e}[\rho]}{\delta \rho} = V_{\text{ext}}(\vec{r}) \tag{F.11}
$$