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Quantum mechanics for many-particle systems

From standard methods to quantum computing and machine learning

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Preface

So, ultimately, in order to understand nature it may be necessary to have a deeper understanding of mathematical relationships. But the real reason is that the subject is enjoyable, and although we humans cut nature up in different ways, and we have different courses in different departments, such compartmentalization is really artificial, and we should take our intellectual pleasures where we find them. *Richard Feynman, The Laws of Thermodynamics.*

Why a preface you may ask? Isn't that just a mere exposition of a *raison d'être* of an author's choice of material, preferences, biases, teaching philosophy etc.? To a large extent I can answer in the affirmative to that. A preface ought to be personal. Indeed, what you will see in the various chapters of these notes represents how I perceive computational physics should be taught.

This set of lecture notes serves the scope of presenting to you and train you in an algorithmic approach to problems in the sciences, represented here by the unity of three disciplines, physics, mathematics and informatics. This trinity outlines the emerging field of computational physics.

Our insight in a physical system, combined with numerical mathematics gives us the rules for setting up an algorithm, viz. a set of rules for solving a particular problem. Our understanding of the physical system under study is obviously gauged by the natural laws at play, the initial conditions, boundary conditions and other external constraints which influence the given system. Having spelled out the physics, for example in the form of a set of coupled partial differential equations, we need efficient numerical methods in order to set up the final algorithm. This algorithm is in turn coded into a computer program and executed on available computing facilities. To develop such an algorithmic approach, you will be exposed to several physics cases, spanning from the classical pendulum to quantum mechanical systems. We will also present some of the most popular algorithms from numerical mathematics used to solve a plethora of problems in the sciences. Finally we will codify these algorithms using some of the most widely used programming languages, presently C, C++ and Fortran and its most recent standard Fortran 2008¹. However, a high-level and fully object-oriented language like Python is now emerging as a good alternative although C++ and Fortran still outperform Python when it comes to computational speed. In this text we offer an approach where one can write all programs in C/C++ or Fortran. We will also show you how to develop large programs in Python interfacing C++ and/or Fortran functions for those parts of the program which are CPU intensive. Such an approach allows you to structure the flow of data in a high-level language like Python while tasks of a mere repetitive and CPU intensive nature are left to low-level languages like C++ or Fortran. Python allows you also to smoothly interface your program with other software, such as plotting programs or operating system instructions.

¹ Throughout this text we refer to Fortran 2008 as Fortran, implying the latest standard.

A typical Python program you may end up writing contains everything from compiling and running your codes to preparing the body of a file for writing up your report.

Computer simulations are nowadays an integral part of contemporary basic and applied research in the sciences. Computation is becoming as important as theory and experiment. In physics, computational physics, theoretical physics and experimental physics are all equally important in our daily research and studies of physical systems. Physics is the unity of theory, experiment and computation². Moreover, the ability "to compute" forms part of the essential repertoire of research scientists. Several new fields within computational science have emerged and strengthened their positions in the last years, such as computational materials science, bioinformatics, computational mathematics and mechanics, computational chemistry and physics and so forth, just to mention a few. These fields underscore the importance of simulations as a means to gain novel insights into physical systems, especially for those cases where no analytical solutions can be found or an experiment is too complicated or expensive to carry out. To be able to simulate large quantal systems with many degrees of freedom such as strongly interacting electrons in a quantum dot will be of great importance for future directions in novel fields like nano-technology. This ability often combines knowledge from many different subjects, in our case essentially from the physical sciences, numerical mathematics, computing languages, topics from high-performance computing and some knowledge of computers.

In 1999, when I started this course at the department of physics in Oslo, computational physics and computational science in general were still perceived by the majority of physicists and scientists as topics dealing with just mere tools and number crunching, and not as subjects of their own. The computational background of most students enlisting for the course on computational physics could span from dedicated hackers and computer freaks to people who basically had never used a PC. The majority of undergraduate and graduate students had a very rudimentary knowledge of computational techniques and methods. Questions like 'do you know of better methods for numerical integration than the trapezoidal rule' were not uncommon. I do happen to know of colleagues who applied for time at a supercomputing centre because they needed to invert matrices of the size of $10^4 \times 10^4$ since they were using the trapezoidal rule to compute integrals. With Gaussian quadrature this dimensionality was easily reduced to matrix problems of the size of $10^2 \times 10^2$, with much better precision.

More than a decade later most students have now been exposed to a fairly uniform introduction to computers, basic programming skills and use of numerical exercises. Practically every undergraduate student in physics has now made a Matlab or Maple simulation of for example the pendulum, with or without chaotic motion. Nowadays most of you are familiar, through various undergraduate courses in physics and mathematics, with interpreted languages such as Maple, Matlab and/or Mathematica. In addition, the interest in scripting languages such as Python or Perl has increased considerably in recent years. The modern programmer would typically combine several tools, computing environments and programming languages. A typical example is the following. Suppose you are working on a project which demands extensive visualizations of the results. To obtain these results, that is to solve a physics problems like obtaining the density profile of a Bose-Einstein condensate, you need however a program which is fairly fast when computational speed matters. In this case you would most

² We mentioned previously the trinity of physics, mathematics and informatics. Viewing physics as the trinity of theory, experiment and simulations is yet another example. It is obviously tempting to go beyond the sciences. History shows that triunes, trinities and for example triple deities permeate the Indo-European cultures (and probably all human cultures), from the ancient Celts and Hindus to modern days. The ancient Celts revered many such trinues, their world was divided into earth, sea and air, nature was divided in animal, vegetable and mineral and the cardinal colours were red, yellow and blue, just to mention a few. As a curious digression, it was a Gaulish Celt, Hilary, philosopher and bishop of Poitiers (AD 315-367) in his work *De Trinitate* who formulated the Holy Trinity concept of Christianity, perhaps in order to accomodate millenia of human divination practice.

likely write a high-performance computing program using Monte Carlo methods in languages which are tailored for that. These are represented by programming languages like Fortran and C++. However, to visualize the results you would find interpreted languages like Matlab or scripting languages like Python extremely suitable for your tasks. You will therefore end up writing for example a script in Matlab which calls a Fortran or C++ program where the number crunching is done and then visualize the results of say a wave equation solver via Matlab's large library of visualization tools. Alternatively, you could organize everything into a Python or Perl script which does everything for you, calls the Fortran and/or C++ programs and performs the visualization in Matlab or Python. Used correctly, these tools, spanning from scripting languages to high-performance computing languages and visualization programs, speed up your capability to solve complicated problems. Being multilingual is thus an advantage which not only applies to our globalized modern society but to computing environments as well. This text shows you how to use C++ and Fortran as programming languages.

There is however more to the picture than meets the eye. Although interpreted languages like Matlab, Mathematica and Maple allow you nowadays to solve very complicated problems, and high-level languages like Python can be used to solve computational problems, computational speed and the capability to write an efficient code are topics which still do matter. To this end, the majority of scientists still use languages like C++ and Fortran to solve scientific problems. When you embark on a master or PhD thesis, you will most likely meet these high-performance computing languages. This course emphasizes thus the use of programming languages like Fortran, Python and C++ instead of interpreted ones like Matlab or Maple. You should however note that there are still large differences in computer time between for example numerical Python and a corresponding C++ program for many numerical applications in the physical sciences, with a code in C++ or Fortran being the fastest.

Computational speed is not the only reason for this choice of programming languages. Another important reason is that we feel that at a certain stage one needs to have some insights into the algorithm used, its stability conditions, possible pitfalls like loss of precision, ranges of applicability, the possibility to improve the algorithm and tailor it to special purposes etc etc. One of our major aims here is to present to you what we would dub 'the algorithmic approach', a set of rules for doing mathematics or a precise description of how to solve a problem. To devise an algorithm and thereafter write a code for solving physics problems is a marvelous way of gaining insight into complicated physical systems. The algorithm you end up writing reflects in essentially all cases your own understanding of the physics and the mathematics (the way you express yourself) of the problem. We do therefore devote quite some space to the algorithms behind various functions presented in the text. Especially, insight into how errors propagate and how to avoid them is a topic we would like you to pay special attention to. Only then can you avoid problems like underflow, overflow and loss of precision. Such a control is not always achievable with interpreted languages and canned functions where the underlying algorithm and/or code is not easily accessible. Although we will at various stages recommend the use of library routines for say linear algebra³, our belief is that one should understand what the given function does, at least to have a mere idea. With such a starting point, we strongly believe that it can be easier to develop more complicated programs on your own using Fortran, C++ or Python.

We have several other aims as well, namely:

- We would like to give you an opportunity to gain a deeper understanding of the physics you have learned in other courses. In most courses one is normally confronted with simple systems which provide exact solutions and mimic to a certain extent the realistic cases. Many are however the comments like 'why can't we do something else than the particle in

³ Such library functions are often tailored to a given machine's architecture and should accordingly run faster than user provided ones.

a box potential?'. In several of the projects we hope to present some more 'realistic' cases to solve by various numerical methods. This also means that we wish to give examples of how physics can be applied in a much broader context than it is discussed in the traditional physics undergraduate curriculum.

- To encourage you to "discover" physics in a way similar to how researchers learn in the context of research.
- Hopefully also to introduce numerical methods and new areas of physics that can be studied with the methods discussed.
- To teach structured programming in the context of doing science.
- The projects we propose are meant to mimic to a certain extent the situation encountered during a thesis or project work. You will typically have at your disposal 2-3 weeks to solve numerically a given project. In so doing you may need to do a literature study as well. Finally, we would like you to write a report for every project.

Our overall goal is to encourage you to learn about science through experience and by asking questions. Our objective is always understanding and the purpose of computing is further insight, not mere numbers! Simulations can often be considered as experiments. Rerunning a simulation need not be as costly as rerunning an experiment.

Needless to say, these lecture notes are upgraded continuously, from typos to new input. And we do always benefit from your comments, suggestions and ideas for making these notes better. It's through the scientific discourse and critics we advance. Moreover, I have benefitted immensely from many discussions with fellow colleagues and students. In particular I must mention Hans Petter Langtangen, Anders Malthe-Sørenssen, Knut Mørken and Øyvind Ryan, whose input during the last fifteen years has considerably improved these lecture notes. Furthermore, the time we have spent and keep spending together on the Computing in Science Education project at the University, is just marvelous. Thanks so much. Concerning the Computing in Science Education initiative, you can read more at <http://www.mn.uio.no/english/about/collaboration/cse/>.

Finally, I would like to add a petit note on referencing. These notes have evolved over many years and the idea is that they should end up in the format of a web-based learning environment for doing computational science. It will be fully free and hopefully represent a much more efficient way of conveying teaching material than traditional textbooks. I have not yet settled on a specific format, so any input is welcome. At present however, it is very easy for me to upgrade and improve the material on say a yearly basis, from simple typos to adding new material. When accessing the web page of the course, you will have noticed that you can obtain all source files for the programs discussed in the text. Many people have thus written to me about how they should properly reference this material and whether they can freely use it. My answer is rather simple. You are encouraged to use these codes, modify them, include them in publications, thesis work, your lectures etc. As long as your use is part of the dialectics of science you can use this material freely. However, since many weekends have elapsed in writing several of these programs, testing them, sweating over bugs, swearing in front of a `f*%@?%g` code which didn't compile properly ten minutes before monday morning's eight o'clock lecture etc etc, I would dearly appreciate in case you find these codes of any use, to reference them properly. That can be done in a simple way, refer to M. Hjorth-Jensen, *Computational Physics*, University of Oslo (2013). The weblink to the course should also be included. Hope it is not too much to ask for. Enjoy!

Acknowledgements

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Acronyms

Use the template *acronym.tex* together with the Springer document class SVMono (monograph-type books) or SVMult (edited books) to style your list(s) of abbreviations or symbols in the Springer layout.

Lists of abbreviations, symbols and the like are easily formatted with the help of the Springer-enhanced `description` environment.

ABC	Spelled-out abbreviation and definition
BABI	Spelled-out abbreviation and definition
CABR	Spelled-out abbreviation and definition

Part I

Linear algebra and second quantization

Chapter 1

Many-body Hamiltonians, basic linear algebra and Second Quantization

Definitions and notations

Before we proceed we need some definitions. We will assume that the interacting part of the Hamiltonian can be approximated by a two-body interaction. This means that our Hamiltonian is written as the sum of some onebody part and a twobody part

$$\hat{H} = \hat{H}_0 + \hat{H}_I = \sum_{i=1}^A \hat{h}_0(x_i) + \sum_{i<j}^A \hat{v}(r_{ij}), \quad (1.1)$$

with

$$H_0 = \sum_{i=1}^A \hat{h}_0(x_i). \quad (1.2)$$

The onebody part $u_{\text{ext}}(x_i)$ is normally approximated by a harmonic oscillator potential or the Coulomb interaction an electron feels from the nucleus. However, other potentials are fully possible, such as one derived from the self-consistent solution of the Hartree-Fock equations to be discussed here.

Our Hamiltonian is invariant under the permutation (interchange) of two particles. Since we deal with fermions however, the total wave function is antisymmetric. Let \hat{P} be an operator which interchanges two particles. Due to the symmetries we have ascribed to our Hamiltonian, this operator commutes with the total Hamiltonian,

$$[\hat{H}, \hat{P}] = 0,$$

meaning that $\Psi_\lambda(x_1, x_2, \dots, x_A)$ is an eigenfunction of \hat{P} as well, that is

$$\hat{P}_{ij} \Psi_\lambda(x_1, x_2, \dots, x_i, \dots, x_j, \dots, x_A) = \beta \Psi_\lambda(x_1, x_2, \dots, x_i, \dots, x_j, \dots, x_A),$$

where β is the eigenvalue of \hat{P} . We have introduced the suffix ij in order to indicate that we permute particles i and j . The Pauli principle tells us that the total wave function for a system of fermions has to be antisymmetric, resulting in the eigenvalue $\beta = -1$.

In our case we assume that we can approximate the exact eigenfunction with a Slater determinant

$$\Phi(x_1, x_2, \dots, x_A, \alpha, \beta, \dots, \sigma) = \frac{1}{\sqrt{A!}} \begin{vmatrix} \psi_\alpha(x_1) & \psi_\alpha(x_2) & \dots & \dots & \psi_\alpha(x_A) \\ \psi_\beta(x_1) & \psi_\beta(x_2) & \dots & \dots & \psi_\beta(x_A) \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \psi_\sigma(x_1) & \psi_\sigma(x_2) & \dots & \dots & \psi_\sigma(x_A) \end{vmatrix}, \quad (1.3)$$

where x_i stand for the coordinates and spin values of a particle i and $\alpha, \beta, \dots, \gamma$ are quantum numbers needed to describe remaining quantum numbers.

Brief reminder on some linear algebra properties.

Before we proceed with a more compact representation of a Slater determinant, we would like to repeat some linear algebra properties which will be useful for our derivations of the energy as function of a Slater determinant, Hartree-Fock theory and later the nuclear shell model.

The inverse of a matrix is defined by

$$\mathbf{A}^{-1} \cdot \mathbf{A} = I$$

A unitary matrix \mathbf{A} is one whose inverse is its adjoint

$$\mathbf{A}^{-1} = \mathbf{A}^\dagger$$

A real unitary matrix is called orthogonal and its inverse is equal to its transpose. A hermitian matrix is its own self-adjoint, that is

$$\mathbf{A} = \mathbf{A}^\dagger.$$

Relations	Name	matrix elements
$A = A^T$	symmetric	$a_{ij} = a_{ji}$
$A = (A^T)^{-1}$	real orthogonal	$\sum_k a_{ik} a_{jk} = \sum_k a_{ki} a_{kj} = \delta_{ij}$
$A = A^*$	real matrix	$a_{ij} = a_{ij}^*$
$A = A^\dagger$	hermitian	$a_{ij} = a_{ji}^*$
$A = (A^\dagger)^{-1}$	unitary	$\sum_k a_{ik} a_{jk}^* = \sum_k a_{ki}^* a_{kj} = \delta_{ij}$

Since we will deal with Fermions (identical and indistinguishable particles) we will form an ansatz for a given state in terms of so-called Slater determinants determined by a chosen basis of single-particle functions.

For a given $n \times n$ matrix \mathbf{A} we can write its determinant

$$\det(\mathbf{A}) = |\mathbf{A}| = \begin{vmatrix} a_{11} & a_{12} & \dots & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & \dots & a_{2n} \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ a_{n1} & a_{n2} & \dots & \dots & a_{nn} \end{vmatrix},$$

in a more compact form as

$$|\mathbf{A}| = \sum_{i=1}^{n!} (-1)^{p_i} \hat{P}_i a_{11} a_{22} \dots a_{nn},$$

where \hat{P}_i is a permutation operator which permutes the column indices $1, 2, 3, \dots, n$ and the sum runs over all $n!$ permutations. The quantity p_i represents the number of transpositions of column indices that are needed in order to bring a given permutation back to its initial ordering, in our case given by $a_{11} a_{22} \dots a_{nn}$ here.

A simple 2×2 determinant illustrates this. We have

$$\det(\mathbf{A}) = \begin{vmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{vmatrix} = (-1)^0 a_{11} a_{22} + (-1)^1 a_{12} a_{21},$$

where in the last term we have interchanged the column indices 1 and 2. The natural ordering we have chosen is $a_{11} a_{22}$.

Back to the derivation of the energy.

The single-particle function $\psi_\alpha(x_i)$ are eigenfunctions of the onebody Hamiltonian h_i , that is

$$\hat{h}_0(x_i) = \hat{t}(x_i) + \hat{u}_{\text{ext}}(x_i),$$

with eigenvalues

$$\hat{h}_0(x_i)\psi_\alpha(x_i) = (\hat{t}(x_i) + \hat{u}_{\text{ext}}(x_i))\psi_\alpha(x_i) = \varepsilon_\alpha\psi_\alpha(x_i).$$

The energies ε_α are the so-called non-interacting single-particle energies, or unperturbed energies. The total energy is in this case the sum over all single-particle energies, if no two-body or more complicated many-body interactions are present.

Let us denote the ground state energy by E_0 . According to the variational principle we have

$$E_0 \leq E[\Phi] = \int \Phi^* \hat{H} \Phi d\tau$$

where Φ is a trial function which we assume to be normalized

$$\int \Phi^* \Phi d\tau = 1,$$

where we have used the shorthand $d\tau = dx_1 dr_2 \dots dr_A$.

In the Hartree-Fock method the trial function is the Slater determinant of Eq. (1.3) which can be rewritten as

$$\Phi(x_1, x_2, \dots, x_A, \alpha, \beta, \dots, \nu) = \frac{1}{\sqrt{A!}} \sum_P (-)^P \hat{P} \psi_\alpha(x_1) \psi_\beta(x_2) \dots \psi_\nu(x_A) = \sqrt{A!} \hat{A} \Phi_H,$$

where we have introduced the antisymmetrization operator \hat{A} defined by the summation over all possible permutations of two particles.

It is defined as

$$\hat{A} = \frac{1}{A!} \sum_P (-)^P \hat{P}, \quad (1.4)$$

with p standing for the number of permutations. We have introduced for later use the so-called Hartree-function, defined by the simple product of all possible single-particle functions

$$\Phi_H(x_1, x_2, \dots, x_A, \alpha, \beta, \dots, \nu) = \psi_\alpha(x_1) \psi_\beta(x_2) \dots \psi_\nu(x_A).$$

Both \hat{H}_0 and \hat{H}_I are invariant under all possible permutations of any two particles and hence commute with \hat{A}

$$[H_0, \hat{A}] = [H_I, \hat{A}] = 0. \quad (1.5)$$

Furthermore, \hat{A} satisfies

$$\hat{A}^2 = \hat{A}, \quad (1.6)$$

since every permutation of the Slater determinant reproduces it.

The expectation value of \hat{H}_0

$$\int \Phi^* \hat{H}_0 \Phi d\tau = A! \int \Phi_H^* \hat{A} \hat{H}_0 \hat{A} \Phi_H d\tau$$

is readily reduced to

$$\int \Phi^* \hat{H}_0 \Phi d\tau = A! \int \Phi_H^* \hat{H}_0 \hat{A} \Phi_H d\tau,$$

where we have used Eqs. (1.5) and (1.6). The next step is to replace the antisymmetrization operator by its definition and to replace \hat{H}_0 with the sum of one-body operators

$$\int \Phi^* \hat{H}_0 \Phi d\tau = \sum_{i=1}^A \sum_p (-)^p \int \Phi_H^* \hat{h}_0 \hat{P} \Phi_H d\tau.$$

The integral vanishes if two or more particles are permuted in only one of the Hartree-functions Φ_H because the individual single-particle wave functions are orthogonal. We obtain then

$$\int \Phi^* \hat{H}_0 \Phi d\tau = \sum_{i=1}^A \int \Phi_H^* \hat{h}_0 \Phi_H d\tau.$$

Orthogonality of the single-particle functions allows us to further simplify the integral, and we arrive at the following expression for the expectation values of the sum of one-body Hamiltonians

$$\int \Phi^* \hat{H}_0 \Phi d\tau = \sum_{\mu=1}^A \int \psi_{\mu}^*(x) \hat{h}_0 \psi_{\mu}(x) dx d\mathbf{r}. \quad (1.7)$$

We introduce the following shorthand for the above integral

$$\langle \mu | \hat{h}_0 | \mu \rangle = \int \psi_{\mu}^*(x) \hat{h}_0 \psi_{\mu}(x) dx,$$

and rewrite Eq. (1.7) as

$$\int \Phi^* \hat{H}_0 \Phi d\tau = \sum_{\mu=1}^A \langle \mu | \hat{h}_0 | \mu \rangle. \quad (1.8)$$

The expectation value of the two-body part of the Hamiltonian is obtained in a similar manner. We have

$$\int \Phi^* \hat{H}_I \Phi d\tau = A! \int \Phi_H^* \hat{A} \hat{H}_I \hat{A} \Phi_H d\tau,$$

which reduces to

$$\int \Phi^* \hat{H}_I \Phi d\tau = \sum_{i \leq j=1}^A \sum_p (-)^p \int \Phi_H^* \hat{v}(r_{ij}) \hat{P} \Phi_H d\tau,$$

by following the same arguments as for the one-body Hamiltonian.

Because of the dependence on the inter-particle distance r_{ij} , permutations of any two particles no longer vanish, and we get

$$\int \Phi^* \hat{H}_I \Phi d\tau = \sum_{i < j=1}^A \int \Phi_H^* \hat{v}(r_{ij}) (1 - P_{ij}) \Phi_H d\tau.$$

where P_{ij} is the permutation operator that interchanges particle i and particle j . Again we use the assumption that the single-particle wave functions are orthogonal.

We obtain

$$\int \Phi^* \hat{H}_I \Phi d\tau = \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \left[\int \psi_{\mu}^*(x_i) \psi_{\nu}^*(x_j) \hat{v}(r_{ij}) \psi_{\mu}(x_i) \psi_{\nu}(x_j) dx_i dx_j \right. \quad (1.9)$$

$$\left. - \int \psi_{\mu}^*(x_i) \psi_{\nu}^*(x_j) \hat{v}(r_{ij}) \psi_{\nu}(x_i) \psi_{\mu}(x_j) dx_i dx_j \right]. \quad (1.10)$$

The first term is the so-called direct term. It is frequently also called the Hartree term, while the second is due to the Pauli principle and is called the exchange term or just the Fock term. The factor $1/2$ is introduced because we now run over all pairs twice.

The last equation allows us to introduce some further definitions. The single-particle wave functions $\psi_{\mu}(x)$, defined by the quantum numbers μ and x are defined as the overlap

$$\psi_{\alpha}(x) = \langle x | \alpha \rangle.$$

We introduce the following shorthands for the above two integrals

$$\langle \mu \nu | \hat{v} | \mu \nu \rangle = \int \psi_\mu^*(x_i) \psi_\nu^*(x_j) \hat{v}(r_{ij}) \psi_\mu(x_i) \psi_\nu(x_j) dx_i dx_j,$$

and

$$\langle \mu \nu | \hat{v} | \nu \mu \rangle = \int \psi_\mu^*(x_i) \psi_\nu^*(x_j) \hat{v}(r_{ij}) \psi_\nu(x_i) \psi_\mu(x_j) dx_i dx_j.$$

Preparing for later studies: varying the coefficients of a wave function expansion and orthogonal transformations

It is common to expand the single-particle functions in a known basis and vary the coefficients, that is, the new single-particle wave function is written as a linear expansion in terms of a fixed chosen orthogonal basis (for example the well-known harmonic oscillator functions or the hydrogen-like functions etc). We define our new single-particle basis (this is a normal approach for Hartree-Fock theory) by performing a unitary transformation on our previous basis (labelled with greek indices) as

$$\psi_p^{new} = \sum_{\lambda} C_{p\lambda} \phi_{\lambda}. \quad (1.11)$$

In this case we vary the coefficients $C_{p\lambda}$. If the basis has infinitely many solutions, we need to truncate the above sum. We assume that the basis ϕ_{λ} is orthogonal.

It is normal to choose a single-particle basis defined as the eigenfunctions of parts of the full Hamiltonian. The typical situation consists of the solutions of the one-body part of the Hamiltonian, that is we have

$$\hat{h}_0 \phi_{\lambda} = \epsilon_{\lambda} \phi_{\lambda}.$$

The single-particle wave functions $\phi_{\lambda}(\mathbf{r})$, defined by the quantum numbers λ and \mathbf{r} are defined as the overlap

$$\phi_{\lambda}(\mathbf{r}) = \langle \mathbf{r} | \lambda \rangle.$$

In deriving the Hartree-Fock equations, we will expand the single-particle functions in a known basis and vary the coefficients, that is, the new single-particle wave function is written as a linear expansion in terms of a fixed chosen orthogonal basis (for example the well-known harmonic oscillator functions or the hydrogen-like functions etc).

We stated that a unitary transformation keeps the orthogonality. To see this consider first a basis of vectors \mathbf{v}_i ,

$$\mathbf{v}_i = \begin{bmatrix} v_{i1} \\ \dots \\ \dots \\ v_{in} \end{bmatrix}$$

We assume that the basis is orthogonal, that is

$$\mathbf{v}_j^T \mathbf{v}_i = \delta_{ij}.$$

An orthogonal or unitary transformation

$$\mathbf{w}_i = \mathbf{U} \mathbf{v}_i,$$

preserves the dot product and orthogonality since

$$\mathbf{w}_j^T \mathbf{w}_i = (\mathbf{U} \mathbf{v}_j)^T \mathbf{U} \mathbf{v}_i = \mathbf{v}_j^T \mathbf{U}^T \mathbf{U} \mathbf{v}_i = \mathbf{v}_j^T \mathbf{v}_i = \delta_{ij}.$$

This means that if the coefficients $C_{p\lambda}$ belong to a unitary or orthogonal transformation (using the Dirac bra-ket notation)

$$|p\rangle = \sum_{\lambda} C_{p\lambda} |\lambda\rangle,$$

orthogonality is preserved, that is $\langle \alpha | \beta \rangle = \delta_{\alpha\beta}$ and $\langle p | q \rangle = \delta_{pq}$.

This property is extremely useful when we build up a basis of many-body Slater determinant based states.

Note also that although a basis $|\alpha\rangle$ contains an infinity of states, for practical calculations we have always to make some truncations.

Before we develop for example the Hartree-Fock equations, there is another very useful property of determinants that we will use both in connection with Hartree-Fock calculations and later shell-model calculations.

Consider the following determinant

$$\begin{vmatrix} \alpha_1 b_{11} + \alpha_2 s b_{12} & a_{12} \\ \alpha_1 b_{21} + \alpha_2 b_{22} & a_{22} \end{vmatrix} = \alpha_1 \begin{vmatrix} b_{11} & a_{12} \\ b_{21} & a_{22} \end{vmatrix} + \alpha_2 \begin{vmatrix} b_{12} & a_{12} \\ b_{22} & a_{22} \end{vmatrix}$$

We can generalize this to an $n \times n$ matrix and have

$$\begin{vmatrix} a_{11} & a_{12} & \dots & \sum_{k=1}^n c_k b_{1k} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & \sum_{k=1}^n c_k b_{2k} & \dots & a_{2n} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ a_{n1} & a_{n2} & \dots & \sum_{k=1}^n c_k b_{nk} & \dots & a_{nn} \end{vmatrix} = \sum_{k=1}^n c_k \begin{vmatrix} a_{11} & a_{12} & \dots & b_{1k} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & b_{2k} & \dots & a_{2n} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ a_{n1} & a_{n2} & \dots & b_{nk} & \dots & a_{nn} \end{vmatrix}.$$

This is a property we will use in our Hartree-Fock discussions.

We can generalize the previous results, now with all elements a_{ij} being given as functions of linear combinations of various coefficients c and elements b_{ij} ,

$$\begin{vmatrix} \sum_{k=1}^n b_{1k} c_{k1} & \sum_{k=1}^n b_{1k} c_{k2} & \dots & \sum_{k=1}^n b_{1k} c_{kj} & \dots & \sum_{k=1}^n b_{1k} c_{kn} \\ \sum_{k=1}^n b_{2k} c_{k1} & \sum_{k=1}^n b_{2k} c_{k2} & \dots & \sum_{k=1}^n b_{2k} c_{kj} & \dots & \sum_{k=1}^n b_{2k} c_{kn} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \sum_{k=1}^n b_{nk} c_{k1} & \sum_{k=1}^n b_{nk} c_{k2} & \dots & \sum_{k=1}^n b_{nk} c_{kj} & \dots & \sum_{k=1}^n b_{nk} c_{kn} \end{vmatrix} = \det(\mathbf{C}) \det(\mathbf{B}),$$

where $\det(\mathbf{C})$ and $\det(\mathbf{B})$ are the determinants of $n \times n$ matrices with elements c_{ij} and b_{ij} respectively. This is a property we will use in our Hartree-Fock discussions. Convince yourself about the correctness of the above expression by setting $n = 2$.

With our definition of the new basis in terms of an orthogonal basis we have

$$\psi_p(x) = \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x).$$

If the coefficients $C_{p\lambda}$ belong to an orthogonal or unitary matrix, the new basis is also orthogonal. Our Slater determinant in the new basis $\psi_p(x)$ is written as

$$\frac{1}{\sqrt{A!}} \begin{vmatrix} \psi_p(x_1) & \psi_p(x_2) & \dots & \dots & \psi_p(x_A) \\ \psi_q(x_1) & \psi_q(x_2) & \dots & \dots & \psi_q(x_A) \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \psi_t(x_1) & \psi_t(x_2) & \dots & \dots & \psi_t(x_A) \end{vmatrix} = \frac{1}{\sqrt{A!}} \begin{vmatrix} \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_A) \\ \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_A) \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_A) \end{vmatrix},$$

which is nothing but $\det(\mathbf{C})\det(\Phi)$, with $\det(\Phi)$ being the determinant given by the basis functions $\phi_\lambda(x)$.

In our discussions hereafter we will use our definitions of single-particle states above and below the Fermi (F) level given by the labels $ijkl \dots \leq F$ for so-called single-hole states and $abcd \dots > F$ for so-called particle states. For general single-particle states we employ the labels $pqrs \dots$.

The energy functional is

$$E[\Phi] = \sum_{\mu=1}^A \langle \mu | h | \mu \rangle + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \langle \mu \nu | \hat{v} | \mu \nu \rangle_{AS},$$

we found the expression for the energy functional in terms of the basis function $\phi_\lambda(\mathbf{r})$. We then varied the above energy functional with respect to the basis functions $|\mu\rangle$. Now we are interested in defining a new basis defined in terms of a chosen basis as defined in Eq. (1.11). We can then rewrite the energy functional as

$$E[\Phi^{New}] = \sum_{i=1}^A \langle i | h | i \rangle + \frac{1}{2} \sum_{ij=1}^A \langle ij | \hat{v} | ij \rangle_{AS}, \quad (1.12)$$

where Φ^{New} is the new Slater determinant defined by the new basis of Eq. (1.11).

Using Eq. (1.11) we can rewrite Eq. (1.12) as

$$E[\Psi] = \sum_{i=1}^A \sum_{\alpha\beta} C_{i\alpha}^* C_{i\beta} \langle \alpha | h | \beta \rangle + \frac{1}{2} \sum_{ij=1}^A \sum_{\alpha\beta\gamma\delta} C_{i\alpha}^* C_{j\beta}^* C_{i\gamma} C_{j\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle_{AS}. \quad (1.13)$$

Chapter 2

Second quantization

We introduce the time-independent operators a_α^\dagger and a_α which create and annihilate, respectively, a particle in the single-particle state φ_α . We define the fermion creation operator a_α^\dagger

$$a_\alpha^\dagger |0\rangle \equiv |\alpha\rangle, \quad (2.1)$$

and

$$a_\alpha^\dagger |\alpha_1 \dots \alpha_n\rangle_{\text{AS}} \equiv |\alpha \alpha_1 \dots \alpha_n\rangle_{\text{AS}} \quad (2.2)$$

In Eq. (2.1) the operator a_α^\dagger acts on the vacuum state $|0\rangle$, which does not contain any particles. Alternatively, we could define a closed-shell nucleus or atom as our new vacuum, but then we need to introduce the particle-hole formalism, see the discussion to come.

In Eq. (2.2) a_α^\dagger acts on an antisymmetric n -particle state and creates an antisymmetric $(n+1)$ -particle state, where the one-body state φ_α is occupied, under the condition that $\alpha \neq \alpha_1, \alpha_2, \dots, \alpha_n$. It follows that we can express an antisymmetric state as the product of the creation operators acting on the vacuum state.

$$|\alpha_1 \dots \alpha_n\rangle_{\text{AS}} = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle \quad (2.3)$$

It is easy to derive the commutation and anticommutation rules for the fermionic creation operators a_α^\dagger . Using the antisymmetry of the states (2.3)

$$|\alpha_1 \dots \alpha_i \dots \alpha_k \dots \alpha_n\rangle_{\text{AS}} = -|\alpha_1 \dots \alpha_k \dots \alpha_i \dots \alpha_n\rangle_{\text{AS}} \quad (2.4)$$

we obtain

$$a_{\alpha_i}^\dagger a_{\alpha_k}^\dagger = -a_{\alpha_k}^\dagger a_{\alpha_i}^\dagger \quad (2.5)$$

Using the Pauli principle

$$|\alpha_1 \dots \alpha_i \dots \alpha_i \dots \alpha_n\rangle_{\text{AS}} = 0 \quad (2.6)$$

it follows that

$$a_{\alpha_i}^\dagger a_{\alpha_i}^\dagger = 0. \quad (2.7)$$

If we combine Eqs. (2.5) and (2.7), we obtain the well-known anti-commutation rule

$$a_\alpha^\dagger a_\beta^\dagger + a_\beta^\dagger a_\alpha^\dagger \equiv \{a_\alpha^\dagger, a_\beta^\dagger\} = 0 \quad (2.8)$$

The hermitian conjugate of a_α^\dagger is

$$a_\alpha = (a_\alpha^\dagger)^\dagger \quad (2.9)$$

If we take the hermitian conjugate of Eq. (2.8), we arrive at

$$\{a_\alpha, a_\beta\} = 0 \quad (2.10)$$

What is the physical interpretation of the operator a_α and what is the effect of a_α on a given state $|\alpha_1 \alpha_2 \dots \alpha_n\rangle_{AS}$? Consider the following matrix element

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha | \alpha'_1 \alpha'_2 \dots \alpha'_m \rangle \quad (2.11)$$

where both sides are antisymmetric. We distinguish between two cases. The first (1) is when $\alpha \in \{\alpha_i\}$. Using the Pauli principle of Eq. (2.6) it follows

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha = 0 \quad (2.12)$$

The second (2) case is when $\alpha \notin \{\alpha_i\}$. It follows that an hermitian conjugation

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha = \langle \alpha \alpha_1 \alpha_2 \dots \alpha_n | \quad (2.13)$$

Eq. (2.13) holds for case (1) since the lefthand side is zero due to the Pauli principle. We write Eq. (2.11) as

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha | \alpha'_1 \alpha'_2 \dots \alpha'_m \rangle = \langle \alpha_1 \alpha_2 \dots \alpha_n | \alpha \alpha'_1 \alpha'_2 \dots \alpha'_m \rangle \quad (2.14)$$

Here we must have $m = n + 1$ if Eq. (2.14) has to be trivially different from zero.

For the last case, the minus and plus signs apply when the sequence $\alpha, \alpha_1, \alpha_2, \dots, \alpha_n$ and $\alpha'_1, \alpha'_2, \dots, \alpha'_{n+1}$ are related to each other via even and odd permutations. If we assume that $\alpha \notin \{\alpha_i\}$ we obtain

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha | \alpha'_1 \alpha'_2 \dots \alpha'_{n+1} \rangle = 0 \quad (2.15)$$

when $\alpha \in \{\alpha'_i\}$. If $\alpha \notin \{\alpha'_i\}$, we obtain

$$a_\alpha | \underbrace{\alpha'_1 \alpha'_2 \dots \alpha'_{n+1}}_{\neq \alpha} \rangle \neq 0 \quad (2.16)$$

and in particular

$$a_\alpha | 0 \rangle = 0 \quad (2.17)$$

If $\{\alpha \alpha_i\} = \{\alpha'_i\}$, performing the right permutations, the sequence $\alpha, \alpha_1, \alpha_2, \dots, \alpha_n$ is identical with the sequence $\alpha'_1, \alpha'_2, \dots, \alpha'_{n+1}$. This results in

$$\langle \alpha_1 \alpha_2 \dots \alpha_n | a_\alpha | \alpha \alpha_1 \alpha_2 \dots \alpha_n \rangle = 1 \quad (2.18)$$

and thus

$$a_\alpha | \alpha \alpha_1 \alpha_2 \dots \alpha_n \rangle = | \alpha_1 \alpha_2 \dots \alpha_n \rangle \quad (2.19)$$

The action of the operator a_α from the left on a state vector is to remove one particle in the state α . If the state vector does not contain the single-particle state α , the outcome of the operation is zero. The operator a_α is normally called for a destruction or annihilation operator.

The next step is to establish the commutator algebra of a_α^\dagger and a_β .

The action of the anti-commutator $\{a_\alpha^\dagger, a_\alpha\}$ on a given n -particle state is

$$\begin{aligned} a_\alpha^\dagger a_\alpha | \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha} \rangle &= 0 \\ a_\alpha a_\alpha^\dagger | \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha} \rangle &= a_\alpha | \underbrace{\alpha \alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha} \rangle = | \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha} \rangle \end{aligned} \quad (2.20)$$

if the single-particle state α is not contained in the state.

If it is present we arrive at

$$\begin{aligned}
a_\alpha^\dagger a_\alpha |\alpha_1 \alpha_2 \dots \alpha_k \alpha \alpha_{k+1} \dots \alpha_{n-1}\rangle &= a_\alpha^\dagger a_\alpha (-1)^k |\alpha \alpha_1 \alpha_2 \dots \alpha_{n-1}\rangle \\
&= (-1)^k |\alpha \alpha_1 \alpha_2 \dots \alpha_{n-1}\rangle = |\alpha_1 \alpha_2 \dots \alpha_k \alpha \alpha_{k+1} \dots \alpha_{n-1}\rangle \\
a_\alpha a_\alpha^\dagger |\alpha_1 \alpha_2 \dots \alpha_k \alpha \alpha_{k+1} \dots \alpha_{n-1}\rangle &= 0
\end{aligned} \tag{2.21}$$

From Eqs. (2.20) and (2.21) we arrive at

$$\{a_\alpha^\dagger, a_\alpha\} = a_\alpha^\dagger a_\alpha + a_\alpha a_\alpha^\dagger = 1 \tag{2.22}$$

The action of $\{a_\alpha^\dagger, a_\beta\}$, with $\alpha \neq \beta$ on a given state yields three possibilities. The first case is a state vector which contains both α and β , then either α or β and finally none of them.

The first case results in

$$\begin{aligned}
a_\alpha^\dagger a_\beta |\alpha \beta \alpha_1 \alpha_2 \dots \alpha_{n-2}\rangle &= 0 \\
a_\beta a_\alpha^\dagger |\alpha \beta \alpha_1 \alpha_2 \dots \alpha_{n-2}\rangle &= 0
\end{aligned} \tag{2.23}$$

while the second case gives

$$\begin{aligned}
a_\alpha^\dagger a_\beta |\beta \underbrace{\alpha_1 \alpha_2 \dots \alpha_{n-1}}_{\neq \alpha}\rangle &= |\alpha \underbrace{\alpha_1 \alpha_2 \dots \alpha_{n-1}}_{\neq \alpha}\rangle \\
a_\beta a_\alpha^\dagger |\beta \underbrace{\alpha_1 \alpha_2 \dots \alpha_{n-1}}_{\neq \alpha}\rangle &= a_\beta |\alpha \beta \underbrace{\alpha_1 \alpha_2 \dots \alpha_{n-1}}_{\neq \alpha}\rangle \\
&= -|\alpha \underbrace{\alpha_1 \alpha_2 \dots \alpha_{n-1}}_{\neq \alpha}\rangle
\end{aligned} \tag{2.24}$$

Finally if the state vector does not contain α and β

$$\begin{aligned}
a_\alpha^\dagger a_\beta |\underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha, \beta}\rangle &= 0 \\
a_\beta a_\alpha^\dagger |\underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha, \beta}\rangle &= a_\beta |\alpha \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha, \beta}\rangle = 0
\end{aligned} \tag{2.25}$$

For all three cases we have

$$\{a_\alpha^\dagger, a_\beta\} = a_\alpha^\dagger a_\beta + a_\beta a_\alpha^\dagger = 0, \quad \alpha \neq \beta \tag{2.26}$$

We can summarize our findings in Eqs. (2.22) and (2.26) as

$$\{a_\alpha^\dagger, a_\beta\} = \delta_{\alpha\beta} \tag{2.27}$$

with $\delta_{\alpha\beta}$ is the Kroenecker δ -symbol.

The properties of the creation and annihilation operators can be summarized as (for fermions)

$$a_\alpha^\dagger |0\rangle \equiv |\alpha\rangle,$$

and

$$a_\alpha^\dagger |\alpha_1 \dots \alpha_n\rangle_{\text{AS}} \equiv |\alpha \alpha_1 \dots \alpha_n\rangle_{\text{AS}}.$$

from which follows

$$|\alpha_1 \dots \alpha_n\rangle_{\text{AS}} = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle.$$

The hermitian conjugate has the following properties

$$a_\alpha = (a_\alpha^\dagger)^\dagger.$$

Finally we found

$$a_\alpha |\underbrace{\alpha'_1 \alpha'_2 \dots \alpha'_{n+1}}_{\neq \alpha}\rangle = 0, \quad \text{in particular } a_\alpha |0\rangle = 0,$$

and

$$a_\alpha |\alpha \alpha_1 \alpha_2 \dots \alpha_n\rangle = |\alpha_1 \alpha_2 \dots \alpha_n\rangle,$$

and the corresponding commutator algebra

$$\{a_\alpha^\dagger, a_\beta^\dagger\} = \{a_\alpha, a_\beta\} = 0 \quad \{a_\alpha^\dagger, a_\beta\} = \delta_{\alpha\beta}.$$

One-body operators in second quantization

A very useful operator is the so-called number-operator. Most physics cases we will study in this text conserve the total number of particles. The number operator is therefore a useful quantity which allows us to test that our many-body formalism conserves the number of particles. In for example (d, p) or (p, d) reactions it is important to be able to describe quantum mechanical states where particles get added or removed. A creation operator a_α^\dagger adds one particle to the single-particle state α of a give many-body state vector, while an annihilation operator a_α removes a particle from a single-particle state α .

Let us consider an operator proportional with $a_\alpha^\dagger a_\beta$ and $\alpha = \beta$. It acts on an n -particle state resulting in

$$a_\alpha^\dagger a_\alpha |\alpha_1 \alpha_2 \dots \alpha_n\rangle = \begin{cases} 0 & \alpha \notin \{\alpha_i\} \\ |\alpha_1 \alpha_2 \dots \alpha_n\rangle & \alpha \in \{\alpha_i\} \end{cases} \quad (2.28)$$

Summing over all possible one-particle states we arrive at

$$\left(\sum_\alpha a_\alpha^\dagger a_\alpha \right) |\alpha_1 \alpha_2 \dots \alpha_n\rangle = n |\alpha_1 \alpha_2 \dots \alpha_n\rangle \quad (2.29)$$

The operator

$$\hat{N} = \sum_\alpha a_\alpha^\dagger a_\alpha \quad (2.30)$$

is called the number operator since it counts the number of particles in a give state vector when it acts on the different single-particle states. It acts on one single-particle state at the time and falls therefore under category one-body operators. Next we look at another important one-body operator, namely \hat{H}_0 and study its operator form in the occupation number representation.

We want to obtain an expression for a one-body operator which conserves the number of particles. Here we study the one-body operator for the kinetic energy plus an eventual external one-body potential. The action of this operator on a particular n -body state with its pertinent expectation value has already been studied in coordinate space. In coordinate space the operator reads

$$\hat{H}_0 = \sum_i \hat{h}_0(x_i) \quad (2.31)$$

and the anti-symmetric n -particle Slater determinant is defined as

$$\Phi(x_1, x_2, \dots, x_n, \alpha_1, \alpha_2, \dots, \alpha_n) = \frac{1}{\sqrt{n!}} \sum_p (-1)^p \hat{P} \psi_{\alpha_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha_n}(x_n).$$

Defining

$$\hat{h}_0(x_i)\psi_{\alpha_i}(x_i) = \sum_{\alpha'_k} \psi_{\alpha'_k}(x_i) \langle \alpha'_k | \hat{h}_0 | \alpha_k \rangle \quad (2.32)$$

we can easily evaluate the action of \hat{H}_0 on each product of one-particle functions in Slater determinant. From Eq. (2.32) we obtain the following result without permuting any particle pair

$$\begin{aligned} & \left(\sum_i \hat{h}_0(x_i) \right) \psi_{\alpha_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha_n}(x_n) \\ = & \sum_{\alpha'_1} \langle \alpha'_1 | \hat{h}_0 | \alpha_1 \rangle \psi_{\alpha'_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha_n}(x_n) \\ + & \sum_{\alpha'_2} \langle \alpha'_2 | \hat{h}_0 | \alpha_2 \rangle \psi_{\alpha_1}(x_1) \psi_{\alpha'_2}(x_2) \dots \psi_{\alpha_n}(x_n) \\ + & \dots \\ + & \sum_{\alpha'_n} \langle \alpha'_n | \hat{h}_0 | \alpha_n \rangle \psi_{\alpha_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha'_n}(x_n) \end{aligned} \quad (2.33)$$

If we interchange particles 1 and 2 we obtain

$$\begin{aligned} & \left(\sum_i \hat{h}_0(x_i) \right) \psi_{\alpha_1}(x_2) \psi_{\alpha_2}(x_1) \dots \psi_{\alpha_n}(x_n) \\ = & \sum_{\alpha'_2} \langle \alpha'_2 | \hat{h}_0 | \alpha_2 \rangle \psi_{\alpha_1}(x_2) \psi_{\alpha'_2}(x_1) \dots \psi_{\alpha_n}(x_n) \\ + & \sum_{\alpha'_1} \langle \alpha'_1 | \hat{h}_0 | \alpha_1 \rangle \psi_{\alpha'_1}(x_2) \psi_{\alpha_2}(x_1) \dots \psi_{\alpha_n}(x_n) \\ + & \dots \\ + & \sum_{\alpha'_n} \langle \alpha'_n | \hat{h}_0 | \alpha_n \rangle \psi_{\alpha_1}(x_2) \psi_{\alpha_2}(x_1) \dots \psi_{\alpha'_n}(x_n) \end{aligned} \quad (2.34)$$

We can continue by computing all possible permutations. We rewrite also our Slater determinant in its second quantized form and skip the dependence on the quantum numbers x_i . Summing up all contributions and taking care of all phases $(-1)^p$ we arrive at

$$\begin{aligned} \hat{H}_0 | \alpha_1, \alpha_2, \dots, \alpha_n \rangle = & \sum_{\alpha'_1} \langle \alpha'_1 | \hat{h}_0 | \alpha_1 \rangle | \alpha'_1 \alpha_2 \dots \alpha_n \rangle \\ + & \sum_{\alpha'_2} \langle \alpha'_2 | \hat{h}_0 | \alpha_2 \rangle | \alpha_1 \alpha'_2 \dots \alpha_n \rangle \\ + & \dots \\ + & \sum_{\alpha'_n} \langle \alpha'_n | \hat{h}_0 | \alpha_n \rangle | \alpha_1 \alpha_2 \dots \alpha'_n \rangle \end{aligned} \quad (2.35)$$

In Eq. (2.35) we have expressed the action of the one-body operator of Eq. (2.31) on the n -body state in its second quantized form. This equation can be further manipulated if we use the properties of the creation and annihilation operator on each primed quantum number, that is

$$| \alpha_1 \alpha_2 \dots \alpha'_k \dots \alpha_n \rangle = a_{\alpha'_k}^\dagger a_{\alpha_k} | \alpha_1 \alpha_2 \dots \alpha_k \dots \alpha_n \rangle \quad (2.36)$$

Inserting this in the right-hand side of Eq. (2.35) results in

$$\begin{aligned}
\hat{H}_0 |\alpha_1 \alpha_2 \dots \alpha_n\rangle &= \sum_{\alpha'_1} \langle \alpha'_1 | \hat{h}_0 | \alpha_1 \rangle a_{\alpha'_1}^\dagger a_{\alpha_1} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&+ \sum_{\alpha'_2} \langle \alpha'_2 | \hat{h}_0 | \alpha_2 \rangle a_{\alpha'_2}^\dagger a_{\alpha_2} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&+ \dots \\
&+ \sum_{\alpha'_n} \langle \alpha'_n | \hat{h}_0 | \alpha_n \rangle a_{\alpha'_n}^\dagger a_{\alpha_n} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&= \sum_{\alpha, \beta} \langle \alpha | \hat{h}_0 | \beta \rangle a_\alpha^\dagger a_\beta |\alpha_1 \alpha_2 \dots \alpha_n\rangle
\end{aligned} \tag{2.37}$$

In the number occupation representation or second quantization we get the following expression for a one-body operator which conserves the number of particles

$$\hat{H}_0 = \sum_{\alpha \beta} \langle \alpha | \hat{h}_0 | \beta \rangle a_\alpha^\dagger a_\beta \tag{2.38}$$

Obviously, \hat{H}_0 can be replaced by any other one-body operator which preserved the number of particles. The structure of the operator is therefore not limited to say the kinetic or single-particle energy only.

The operator \hat{H}_0 takes a particle from the single-particle state β to the single-particle state α with a probability for the transition given by the expectation value $\langle \alpha | \hat{h}_0 | \beta \rangle$.

It is instructive to verify Eq. (2.38) by computing the expectation value of \hat{H}_0 between two single-particle states

$$\langle \alpha_1 | \hat{h}_0 | \alpha_2 \rangle = \sum_{\alpha \beta} \langle \alpha | \hat{h}_0 | \beta \rangle \langle 0 | a_{\alpha_1} a_\alpha^\dagger a_\beta a_{\alpha_2}^\dagger | 0 \rangle \tag{2.39}$$

Using the commutation relations for the creation and annihilation operators we have

$$a_{\alpha_1} a_\alpha^\dagger a_\beta a_{\alpha_2}^\dagger = (\delta_{\alpha \alpha_1} - a_\alpha^\dagger a_{\alpha_1}) (\delta_{\beta \alpha_2} - a_{\alpha_2}^\dagger a_\beta), \tag{2.40}$$

which results in

$$\langle 0 | a_{\alpha_1} a_\alpha^\dagger a_\beta a_{\alpha_2}^\dagger | 0 \rangle = \delta_{\alpha \alpha_1} \delta_{\beta \alpha_2} \tag{2.41}$$

and

$$\langle \alpha_1 | \hat{h}_0 | \alpha_2 \rangle = \sum_{\alpha \beta} \langle \alpha | \hat{h}_0 | \beta \rangle \delta_{\alpha \alpha_1} \delta_{\beta \alpha_2} = \langle \alpha_1 | \hat{h}_0 | \alpha_2 \rangle \tag{2.42}$$

Two-body operators in second quantization

Let us now derive the expression for our two-body interaction part, which also conserves the number of particles. We can proceed in exactly the same way as for the one-body operator. In the coordinate representation our two-body interaction part takes the following expression

$$\hat{H}_I = \sum_{i < j} V(x_i, x_j) \tag{2.43}$$

where the summation runs over distinct pairs. The term V can be an interaction model for the nucleon-nucleon interaction or the interaction between two electrons. It can also include additional two-body interaction terms.

The action of this operator on a product of two single-particle functions is defined as

$$V(x_i, x_j) \psi_{\alpha_k}(x_i) \psi_{\alpha_l}(x_j) = \sum_{\alpha'_k \alpha'_l} \psi'_{\alpha'_k}(x_i) \psi'_{\alpha'_l}(x_j) \langle \alpha'_k \alpha'_l | \hat{v} | \alpha_k \alpha_l \rangle \tag{2.44}$$

We can now let \hat{H}_I act on all terms in the linear combination for $|\alpha_1 \alpha_2 \dots \alpha_n\rangle$. Without any permutations we have

$$\begin{aligned}
& \left(\sum_{i < j} V(x_i, x_j) \right) \psi_{\alpha_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha_n}(x_n) \\
= & \sum_{\alpha'_1 \alpha'_2} \langle \alpha'_1 \alpha'_2 | \hat{v} | \alpha_1 \alpha_2 \rangle \psi'_{\alpha'_1}(x_1) \psi'_{\alpha'_2}(x_2) \dots \psi_{\alpha_n}(x_n) \\
& + \dots \\
& + \sum_{\alpha'_1 \alpha'_n} \langle \alpha'_1 \alpha'_n | \hat{v} | \alpha_1 \alpha_n \rangle \psi'_{\alpha'_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi'_{\alpha'_n}(x_n) \\
& + \dots \\
& + \sum_{\alpha'_2 \alpha'_n} \langle \alpha'_2 \alpha'_n | \hat{v} | \alpha_2 \alpha_n \rangle \psi_{\alpha_1}(x_1) \psi'_{\alpha'_2}(x_2) \dots \psi'_{\alpha'_n}(x_n) \\
& + \dots
\end{aligned} \tag{2.45}$$

where on the rhs we have a term for each distinct pairs.

For the other terms on the rhs we obtain similar expressions and summing over all terms we obtain

$$\begin{aligned}
H_I |\alpha_1 \alpha_2 \dots \alpha_n\rangle = & \sum_{\alpha'_1, \alpha'_2} \langle \alpha'_1 \alpha'_2 | \hat{v} | \alpha_1 \alpha_2 \rangle |\alpha'_1 \alpha'_2 \dots \alpha_n\rangle \\
& + \dots \\
& + \sum_{\alpha'_1, \alpha'_n} \langle \alpha'_1 \alpha'_n | \hat{v} | \alpha_1 \alpha_n \rangle |\alpha'_1 \alpha_2 \dots \alpha'_n\rangle \\
& + \dots \\
& + \sum_{\alpha'_2, \alpha'_n} \langle \alpha'_2 \alpha'_n | \hat{v} | \alpha_2 \alpha_n \rangle |\alpha_1 \alpha'_2 \dots \alpha'_n\rangle \\
& + \dots
\end{aligned} \tag{2.46}$$

We introduce second quantization via the relation

$$\begin{aligned}
& a_{\alpha'_k}^\dagger a_{\alpha'_l}^\dagger a_{\alpha_l} a_{\alpha_k} |\alpha_1 \alpha_2 \dots \alpha_k \dots \alpha_l \dots \alpha_n\rangle \\
= & (-1)^{k-1} (-1)^{l-2} a_{\alpha'_k}^\dagger a_{\alpha'_l}^\dagger a_{\alpha_l} a_{\alpha_k} |\alpha_k \alpha_l \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha_k, \alpha_l}\rangle \\
= & (-1)^{k-1} (-1)^{l-2} |\alpha'_k \alpha'_l \underbrace{\alpha_1 \alpha_2 \dots \alpha_n}_{\neq \alpha'_k, \alpha'_l}\rangle \\
= & |\alpha_1 \alpha_2 \dots \alpha'_k \dots \alpha'_l \dots \alpha_n\rangle
\end{aligned} \tag{2.47}$$

Inserting this in (2.46) gives

$$\begin{aligned}
H_I |\alpha_1 \alpha_2 \dots \alpha_n\rangle &= \sum_{\alpha'_1, \alpha'_2} \langle \alpha'_1 \alpha'_2 | \hat{v} | \alpha_1 \alpha_2 \rangle a_{\alpha'_1}^\dagger a_{\alpha'_2}^\dagger a_{\alpha_2} a_{\alpha_1} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&+ \dots \\
&= \sum_{\alpha'_1, \alpha'_n} \langle \alpha'_1 \alpha'_n | \hat{v} | \alpha_1 \alpha_n \rangle a_{\alpha'_1}^\dagger a_{\alpha'_n}^\dagger a_{\alpha_n} a_{\alpha_1} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&+ \dots \\
&= \sum_{\alpha'_2, \alpha'_n} \langle \alpha'_2 \alpha'_n | \hat{v} | \alpha_2 \alpha_n \rangle a_{\alpha'_2}^\dagger a_{\alpha'_n}^\dagger a_{\alpha_n} a_{\alpha_2} |\alpha_1 \alpha_2 \dots \alpha_n\rangle \\
&+ \dots \\
&= \sum'_{\alpha, \beta, \gamma, \delta} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma |\alpha_1 \alpha_2 \dots \alpha_n\rangle \quad (2.48)
\end{aligned}$$

Here we let \sum' indicate that the sums running over α and β run over all single-particle states, while the summations γ and δ run over all pairs of single-particle states. We wish to remove this restriction and since

$$\langle \alpha \beta | \hat{v} | \gamma \delta \rangle = \langle \beta \alpha | \hat{v} | \delta \gamma \rangle \quad (2.49)$$

we get

$$\sum_{\alpha \beta} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma = \sum_{\alpha \beta} \langle \beta \alpha | \hat{v} | \delta \gamma \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma \quad (2.50)$$

$$= \sum_{\alpha \beta} \langle \beta \alpha | \hat{v} | \delta \gamma \rangle a_\beta^\dagger a_\alpha^\dagger a_\gamma a_\delta \quad (2.51)$$

where we have used the anti-commutation rules.

Changing the summation indices α and β in (2.51) we obtain

$$\sum_{\alpha \beta} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma = \sum_{\alpha \beta} \langle \alpha \beta | \hat{v} | \delta \gamma \rangle a_\alpha^\dagger a_\beta^\dagger a_\gamma a_\delta \quad (2.52)$$

From this it follows that the restriction on the summation over γ and δ can be removed if we multiply with a factor $\frac{1}{2}$, resulting in

$$\hat{H}_I = \frac{1}{2} \sum_{\alpha \beta \gamma \delta} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma \quad (2.53)$$

where we sum freely over all single-particle states α, β, γ og δ .

With this expression we can now verify that the second quantization form of \hat{H}_I in Eq. (2.53) results in the same matrix between two anti-symmetrized two-particle states as its corresponding coordinate space representation. We have

$$\langle \alpha_1 \alpha_2 | \hat{H}_I | \beta_1 \beta_2 \rangle = \frac{1}{2} \sum_{\alpha \beta \gamma \delta} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle \langle 0 | a_{\alpha_2} a_{\alpha_1} a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma a_{\beta_1}^\dagger a_{\beta_2}^\dagger | 0 \rangle. \quad (2.54)$$

Using the commutation relations we get

$$\begin{aligned}
& a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma} a_{\beta_1}^{\dagger} a_{\beta_2}^{\dagger} \\
= & a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} (a_{\delta} \delta_{\gamma\beta_1} a_{\beta_2}^{\dagger} - a_{\delta} a_{\beta_1}^{\dagger} a_{\gamma} a_{\beta_2}^{\dagger}) \\
= & a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} (\delta_{\gamma\beta_1} \delta_{\delta\beta_2} - \delta_{\gamma\beta_1} a_{\beta_2}^{\dagger} a_{\delta} - a_{\delta} a_{\beta_1}^{\dagger} \delta_{\gamma\beta_2} + a_{\delta} a_{\beta_1}^{\dagger} a_{\beta_2}^{\dagger} a_{\gamma}) \\
= & a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} (\delta_{\gamma\beta_1} \delta_{\delta\beta_2} - \delta_{\gamma\beta_1} a_{\beta_2}^{\dagger} a_{\delta} \\
& - \delta_{\delta\beta_1} \delta_{\gamma\beta_2} + \delta_{\gamma\beta_2} a_{\beta_1}^{\dagger} a_{\delta} + a_{\delta} a_{\beta_1}^{\dagger} a_{\beta_2}^{\dagger} a_{\gamma}) \quad (2.55)
\end{aligned}$$

The vacuum expectation value of this product of operators becomes

$$\begin{aligned}
& \langle 0 | a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma} a_{\beta_1}^{\dagger} a_{\beta_2}^{\dagger} | 0 \rangle \\
= & (\delta_{\gamma\beta_1} \delta_{\delta\beta_2} - \delta_{\delta\beta_1} \delta_{\gamma\beta_2}) \langle 0 | a_{\alpha_2} a_{\alpha_1} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} | 0 \rangle \\
= & (\delta_{\gamma\beta_1} \delta_{\delta\beta_2} - \delta_{\delta\beta_1} \delta_{\gamma\beta_2}) (\delta_{\alpha\alpha_1} \delta_{\beta\alpha_2} - \delta_{\beta\alpha_1} \delta_{\alpha\alpha_2}) \quad (2.56)
\end{aligned}$$

Insertion of Eq. (2.56) in Eq. (2.54) results in

$$\begin{aligned}
\langle \alpha_1 \alpha_2 | \hat{H}_I | \beta_1 \beta_2 \rangle &= \frac{1}{2} [\langle \alpha_1 \alpha_2 | \hat{v} | \beta_1 \beta_2 \rangle - \langle \alpha_1 \alpha_2 | \hat{v} | \beta_2 \beta_1 \rangle \\
&\quad - \langle \alpha_2 \alpha_1 | \hat{v} | \beta_1 \beta_2 \rangle + \langle \alpha_2 \alpha_1 | \hat{v} | \beta_2 \beta_1 \rangle] \\
&= \langle \alpha_1 \alpha_2 | \hat{v} | \beta_1 \beta_2 \rangle - \langle \alpha_1 \alpha_2 | \hat{v} | \beta_2 \beta_1 \rangle \\
&= \langle \alpha_1 \alpha_2 | \hat{v} | \beta_1 \beta_2 \rangle_{\text{AS}}. \quad (2.57)
\end{aligned}$$

The two-body operator can also be expressed in terms of the anti-symmetrized matrix elements we discussed previously as

$$\begin{aligned}
\hat{H}_I &= \frac{1}{2} \sum_{\alpha\beta\gamma\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma} \\
&= \frac{1}{4} \sum_{\alpha\beta\gamma\delta} [\langle \alpha\beta | \hat{v} | \gamma\delta \rangle - \langle \alpha\beta | \hat{v} | \delta\gamma \rangle] a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma} \\
&= \frac{1}{4} \sum_{\alpha\beta\gamma\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle_{\text{AS}} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma} \quad (2.58)
\end{aligned}$$

The factors in front of the operator, either $\frac{1}{4}$ or $\frac{1}{2}$ tells whether we use antisymmetrized matrix elements or not.

We can now express the Hamiltonian operator for a many-fermion system in the occupation basis representation as

$$H = \sum_{\alpha,\beta} \langle \alpha | \hat{f} + \hat{u}_{\text{ext}} | \beta \rangle a_{\alpha}^{\dagger} a_{\beta} + \frac{1}{4} \sum_{\alpha\beta\gamma\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle_{\text{AS}} a_{\alpha}^{\dagger} a_{\beta}^{\dagger} a_{\delta} a_{\gamma}. \quad (2.59)$$

This is the form we will use in the rest of these lectures, assuming that we work with anti-symmetrized two-body matrix elements.

Particle-hole formalism

Second quantization is a useful and elegant formalism for constructing many-body states and quantum mechanical operators. One can express and translate many physical processes into simple pictures such as Feynman diagrams. Expectation values of many-body states are also easily calculated. However, although the equations are seemingly easy to set up, from a prac-

tical point of view, that is the solution of Schroedinger's equation, there is no particular gain. The many-body equation is equally hard to solve, irrespective of representation. The cliché that there is no free lunch brings us down to earth again. Note however that a transformation to a particular basis, for cases where the interaction obeys specific symmetries, can ease the solution of Schroedinger's equation.

But there is at least one important case where second quantization comes to our rescue. It is namely easy to introduce another reference state than the pure vacuum $|0\rangle$, where all single-particle states are active. With many particles present it is often useful to introduce another reference state than the vacuum state $|0\rangle$. We will label this state $|c\rangle$ (c for core) and as we will see it can reduce considerably the complexity and thereby the dimensionality of the many-body problem. It allows us to sum up to infinite order specific many-body correlations. The particle-hole representation is one of these handy representations.

In the original particle representation these states are products of the creation operators $a_{\alpha_i}^\dagger$ acting on the true vacuum $|0\rangle$. Following Eq. (2.3) we have

$$|\alpha_1 \alpha_2 \dots \alpha_{n-1} \alpha_n\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_{n-1}}^\dagger a_{\alpha_n}^\dagger |0\rangle \quad (2.60)$$

$$|\alpha_1 \alpha_2 \dots \alpha_{n-1} \alpha_n \alpha_{n+1}\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_{n-1}}^\dagger a_{\alpha_n}^\dagger a_{\alpha_{n+1}}^\dagger |0\rangle \quad (2.61)$$

$$|\alpha_1 \alpha_2 \dots \alpha_{n-1}\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_{n-1}}^\dagger |0\rangle \quad (2.62)$$

If we use Eq. (2.60) as our new reference state, we can simplify considerably the representation of this state

$$|c\rangle \equiv |\alpha_1 \alpha_2 \dots \alpha_{n-1} \alpha_n\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_{n-1}}^\dagger a_{\alpha_n}^\dagger |0\rangle \quad (2.63)$$

The new reference states for the $n+1$ and $n-1$ states can then be written as

$$|\alpha_1 \alpha_2 \dots \alpha_{n-1} \alpha_n \alpha_{n+1}\rangle = (-1)^n a_{\alpha_{n+1}}^\dagger |c\rangle \equiv (-1)^n |\alpha_{n+1}\rangle_c \quad (2.64)$$

$$|\alpha_1 \alpha_2 \dots \alpha_{n-1}\rangle = (-1)^{n-1} a_{\alpha_n} |c\rangle \equiv (-1)^{n-1} |\alpha_{n-1}\rangle_c \quad (2.65)$$

The first state has one additional particle with respect to the new vacuum state $|c\rangle$ and is normally referred to as a one-particle state or one particle added to the many-body reference state. The second state has one particle less than the reference vacuum state $|c\rangle$ and is referred to as a one-hole state. When dealing with a new reference state it is often convenient to introduce new creation and annihilation operators since we have from Eq. (2.65)

$$a_\alpha |c\rangle \neq 0 \quad (2.66)$$

since α is contained in $|c\rangle$, while for the true vacuum we have $a_\alpha |0\rangle = 0$ for all α .

The new reference state leads to the definition of new creation and annihilation operators which satisfy the following relations

$$b_\alpha |c\rangle = 0 \quad (2.67)$$

$$\{b_\alpha^\dagger, b_\beta^\dagger\} = \{b_\alpha, b_\beta\} = 0$$

$$\{b_\alpha^\dagger, b_\beta\} = \delta_{\alpha\beta} \quad (2.68)$$

We assume also that the new reference state is properly normalized

$$\langle c|c\rangle = 1 \quad (2.69)$$

The physical interpretation of these new operators is that of so-called quasiparticle states. This means that a state defined by the addition of one extra particle to a reference state $|c\rangle$ may not necessarily be interpreted as one particle coupled to a core. We define now new creation operators that act on a state α creating a new quasiparticle state

$$b_{\alpha}^{\dagger}|c\rangle = \begin{cases} a_{\alpha}^{\dagger}|c\rangle = |\alpha\rangle, & \alpha > F \\ a_{\alpha}|c\rangle = |\alpha^{-1}\rangle, & \alpha \leq F \end{cases} \quad (2.70)$$

where F is the Fermi level representing the last occupied single-particle orbit of the new reference state $|c\rangle$.

The annihilation is the hermitian conjugate of the creation operator

$$b_{\alpha} = (b_{\alpha}^{\dagger})^{\dagger},$$

resulting in

$$b_{\alpha}^{\dagger} = \begin{cases} a_{\alpha}^{\dagger} & \alpha > F \\ a_{\alpha} & \alpha \leq F \end{cases} \quad b_{\alpha} = \begin{cases} a_{\alpha} & \alpha > F \\ a_{\alpha}^{\dagger} & \alpha \leq F \end{cases} \quad (2.71)$$

With the new creation and annihilation operator we can now construct many-body quasi-particle states, with one-particle-one-hole states, two-particle-two-hole states etc in the same fashion as we previously constructed many-particle states. We can write a general particle-hole state as

$$|\beta_1\beta_2\ldots\beta_{n_p}\gamma_1^{-1}\gamma_2^{-1}\ldots\gamma_{n_h}^{-1}\rangle \equiv \underbrace{b_{\beta_1}^{\dagger}b_{\beta_2}^{\dagger}\ldots b_{\beta_{n_p}}^{\dagger}}_{>F} \underbrace{b_{\gamma_1}^{\dagger}b_{\gamma_2}^{\dagger}\ldots b_{\gamma_{n_h}}^{\dagger}}_{\leq F} |c\rangle \quad (2.72)$$

We can now rewrite our one-body and two-body operators in terms of the new creation and annihilation operators. The number operator becomes

$$\hat{N} = \sum_{\alpha} a_{\alpha}^{\dagger} a_{\alpha} = \sum_{\alpha > F} b_{\alpha}^{\dagger} b_{\alpha} + n_c - \sum_{\alpha \leq F} b_{\alpha}^{\dagger} b_{\alpha} \quad (2.73)$$

where n_c is the number of particle in the new vacuum state $|c\rangle$. The action of \hat{N} on a many-body state results in

$$\hat{N}|\beta_1\beta_2\ldots\beta_{n_p}\gamma_1^{-1}\gamma_2^{-1}\ldots\gamma_{n_h}^{-1}\rangle = (n_p + n_c - n_h)|\beta_1\beta_2\ldots\beta_{n_p}\gamma_1^{-1}\gamma_2^{-1}\ldots\gamma_{n_h}^{-1}\rangle \quad (2.74)$$

Here $n = n_p + n_c - n_h$ is the total number of particles in the quasi-particle state of Eq. (2.72). Note that \hat{N} counts the total number of particles present

$$N_{qp} = \sum_{\alpha} b_{\alpha}^{\dagger} b_{\alpha}, \quad (2.75)$$

gives us the number of quasi-particles as can be seen by computing

$$N_{qp}|\beta_1\beta_2\ldots\beta_{n_p}\gamma_1^{-1}\gamma_2^{-1}\ldots\gamma_{n_h}^{-1}\rangle = (n_p + n_h)|\beta_1\beta_2\ldots\beta_{n_p}\gamma_1^{-1}\gamma_2^{-1}\ldots\gamma_{n_h}^{-1}\rangle \quad (2.76)$$

where $n_{qp} = n_p + n_h$ is the total number of quasi-particles.

We express the one-body operator \hat{H}_0 in terms of the quasi-particle creation and annihilation operators, resulting in

$$\begin{aligned} \hat{H}_0 = & \sum_{\alpha\beta > F} \langle\alpha|\hat{h}_0|\beta\rangle b_{\alpha}^{\dagger} b_{\beta} + \sum_{\alpha > F, \beta \leq F} \left[\langle\alpha|\hat{h}_0|\beta\rangle b_{\alpha}^{\dagger} b_{\beta}^{\dagger} + \langle\beta|\hat{h}_0|\alpha\rangle b_{\beta} b_{\alpha} \right] \\ & + \sum_{\alpha \leq F} \langle\alpha|\hat{h}_0|\alpha\rangle - \sum_{\alpha\beta \leq F} \langle\beta|\hat{h}_0|\alpha\rangle b_{\alpha}^{\dagger} b_{\beta} \end{aligned} \quad (2.77)$$

The first term gives contribution only for particle states, while the last one contributes only for holestates. The second term can create or destroy a set of quasi-particles and the third term is the contribution from the vacuum state $|c\rangle$.

Before we continue with the expressions for the two-body operator, we introduce a nomenclature we will use for the rest of this text. It is inspired by the notation used in quantum chemistry. We reserve the labels i, j, k, \dots for hole states and a, b, c, \dots for states above F , viz. particle states. This means also that we will skip the constraint $\leq F$ or $> F$ in the summation symbols. Our operator \hat{H}_0 reads now

$$\begin{aligned} \hat{H}_0 = & \sum_{ab} \langle a|\hat{h}|b\rangle b_a^\dagger b_b + \sum_{ai} \left[\langle a|\hat{h}|i\rangle b_a^\dagger b_i^\dagger + \langle i|\hat{h}|a\rangle b_i b_a \right] \\ & + \sum_i \langle i|\hat{h}|i\rangle - \sum_{ij} \langle j|\hat{h}|i\rangle b_i^\dagger b_j \end{aligned} \quad (2.78)$$

The two-particle operator in the particle-hole formalism is more complicated since we have to translate four indices $\alpha\beta\gamma\delta$ to the possible combinations of particle and hole states. When performing the commutator algebra we can regroup the operator in five different terms

$$\hat{H}_I = \hat{H}_I^{(a)} + \hat{H}_I^{(b)} + \hat{H}_I^{(c)} + \hat{H}_I^{(d)} + \hat{H}_I^{(e)} \quad (2.79)$$

Using anti-symmetrized matrix elements, the term $\hat{H}_I^{(a)}$ is

$$\hat{H}_I^{(a)} = \frac{1}{4} \sum_{abcd} \langle ab|\hat{V}|cd\rangle b_a^\dagger b_b^\dagger b_d b_c \quad (2.80)$$

The next term $\hat{H}_I^{(b)}$ reads

$$\hat{H}_I^{(b)} = \frac{1}{4} \sum_{abci} \left(\langle ab|\hat{V}|ci\rangle b_a^\dagger b_b^\dagger b_i^\dagger b_c + \langle ai|\hat{V}|cb\rangle b_a^\dagger b_i b_b b_c \right) \quad (2.81)$$

This term conserves the number of quasiparticles but creates or removes a three-particle-one-hole state. For $\hat{H}_I^{(c)}$ we have

$$\begin{aligned} \hat{H}_I^{(c)} = & \frac{1}{4} \sum_{abij} \left(\langle ab|\hat{V}|ij\rangle b_a^\dagger b_b^\dagger b_j^\dagger b_i^\dagger + \langle ij|\hat{V}|ab\rangle b_a b_b b_j b_i \right) + \\ & \frac{1}{2} \sum_{abij} \langle ai|\hat{V}|bj\rangle b_a^\dagger b_j^\dagger b_b b_i + \frac{1}{2} \sum_{abi} \langle ai|\hat{V}|bi\rangle b_a^\dagger b_b. \end{aligned} \quad (2.82)$$

The first line stands for the creation of a two-particle-two-hole state, while the second line represents the creation to two one-particle-one-hole pairs while the last term represents a contribution to the particle single-particle energy from the hole states, that is an interaction between the particle states and the hole states within the new vacuum state. The fourth term reads

$$\begin{aligned} \hat{H}_I^{(d)} = & \frac{1}{4} \sum_{aijk} \left(\langle ai|\hat{V}|jk\rangle b_a^\dagger b_k^\dagger b_j^\dagger b_i + \langle ji|\hat{V}|ak\rangle b_k^\dagger b_j b_i b_a \right) + \\ & \frac{1}{4} \sum_{aij} \left(\langle ai|\hat{V}|ji\rangle b_a^\dagger b_j^\dagger + \langle ji|\hat{V}|ai\rangle - \langle ji|\hat{V}|ia\rangle b_j b_a \right). \end{aligned} \quad (2.83)$$

The terms in the first line stand for the creation of a particle-hole state interacting with hole states, we will label this as a two-hole-one-particle contribution. The remaining terms are a particle-hole state interacting with the holes in the vacuum state. Finally we have

$$\hat{H}_I^{(e)} = \frac{1}{4} \sum_{ijkl} \langle kl|\hat{V}|ij\rangle b_i^\dagger b_j^\dagger b_l b_k + \frac{1}{2} \sum_{ijk} \langle ij|\hat{V}|kj\rangle b_k^\dagger b_i + \frac{1}{2} \sum_{ij} \langle ij|\hat{V}|ij\rangle \quad (2.84)$$

The first terms represents the interaction between two holes while the second stands for the interaction between a hole and the remaining holes in the vacuum state. It represents a contribution to single-hole energy to first order. The last term collects all contributions to the energy of the ground state of a closed-shell system arising from hole-hole correlations.

Summarizing and defining a normal-ordered Hamiltonian

$$\Phi_{AS}(\alpha_1, \dots, \alpha_A; x_1, \dots, x_A) = \frac{1}{\sqrt{A}} \sum_{\hat{p}} (-1)^P \hat{p} \prod_{i=1}^A \psi_{\alpha_i}(x_i),$$

which is equivalent with $|\alpha_1 \dots \alpha_A\rangle = a_{\alpha_1}^\dagger \dots a_{\alpha_A}^\dagger |0\rangle$. We have also

$$a_p^\dagger |0\rangle = |p\rangle, \quad a_p |q\rangle = \delta_{pq} |0\rangle$$

$$\delta_{pq} = \{a_p, a_q^\dagger\},$$

and

$$0 = \{a_p^\dagger, a_q\} = \{a_p, a_q\} = \{a_p^\dagger, a_q^\dagger\}$$

$$|\Phi_0\rangle = |\alpha_1 \dots \alpha_A\rangle, \quad \alpha_1, \dots, \alpha_A \leq \alpha_F$$

$$\{a_p^\dagger, a_q\} = \delta_{pq}, p, q \leq \alpha_F$$

$$\{a_p, a_q^\dagger\} = \delta_{pq}, p, q > \alpha_F$$

with $i, j, \dots \leq \alpha_F$, $a, b, \dots > \alpha_F$, p, q, \dots – any

$$a_i |\Phi_0\rangle = |\Phi_i\rangle, \quad a_a^\dagger |\Phi_0\rangle = |\Phi^a\rangle$$

and

$$a_i^\dagger |\Phi_0\rangle = 0 \quad a_a |\Phi_0\rangle = 0$$

The one-body operator is defined as

$$\hat{F} = \sum_{pq} \langle p | \hat{f} | q \rangle a_p^\dagger a_q$$

while the two-body operator is defined as

$$\hat{V} = \frac{1}{4} \sum_{pqrs} \langle pq | \hat{v} | rs \rangle_{AS} a_p^\dagger a_q^\dagger a_s a_r$$

where we have defined the antisymmetric matrix elements

$$\langle pq | \hat{v} | rs \rangle_{AS} = \langle pq | \hat{v} | rs \rangle - \langle pq | \hat{v} | sr \rangle.$$

We can also define a three-body operator

$$\hat{V}_3 = \frac{1}{36} \sum_{pqrstu} \langle pqr | \hat{v}_3 | stu \rangle_{AS} a_p^\dagger a_q^\dagger a_r^\dagger a_u a_t a_s$$

with the antisymmetrized matrix element

$$\langle pqr|\hat{v}_3|stu\rangle_{AS} = \langle pqr|\hat{v}_3|stu\rangle + \langle pqr|\hat{v}_3|tus\rangle + \langle pqr|\hat{v}_3|ust\rangle - \langle pqr|\hat{v}_3|sut\rangle - \langle pqr|\hat{v}_3|tsu\rangle - \langle pqr|\hat{v}_3|uts\rangle. \quad (2.85)$$

Operators in second quantization

In the build-up of a shell-model or FCI code that is meant to tackle large dimensionalities is the action of the Hamiltonian \hat{H} on a Slater determinant represented in second quantization as

$$|\alpha_1 \dots \alpha_n\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle.$$

The time consuming part stems from the action of the Hamiltonian on the above determinant,

$$\left(\sum_{\alpha\beta} \langle \alpha|t+u|\beta\rangle a_{\alpha}^\dagger a_{\beta} + \frac{1}{4} \sum_{\alpha\beta\gamma\delta} \langle \alpha\beta|\hat{v}|\gamma\delta\rangle a_{\alpha}^\dagger a_{\beta}^\dagger a_{\delta} a_{\gamma} \right) a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle.$$

A practically useful way to implement this action is to encode a Slater determinant as a bit pattern.

Assume that we have at our disposal n different single-particle orbits $\alpha_0, \alpha_2, \dots, \alpha_{n-1}$ and that we can distribute among these orbits $N \leq n$ particles.

A Slater determinant can then be coded as an integer of n bits. As an example, if we have $n = 16$ single-particle states $\alpha_0, \alpha_1, \dots, \alpha_{15}$ and $N = 4$ fermions occupying the states $\alpha_3, \alpha_6, \alpha_{10}$ and α_{13} we could write this Slater determinant as

$$\Phi_A = a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle.$$

The unoccupied single-particle states have bit value 0 while the occupied ones are represented by bit state 1. In the binary notation we would write this 16 bits long integer as

$$\begin{array}{cccccccccccccccc} \alpha_0 & \alpha_1 & \alpha_2 & \alpha_3 & \alpha_4 & \alpha_5 & \alpha_6 & \alpha_7 & \alpha_8 & \alpha_9 & \alpha_{10} & \alpha_{11} & \alpha_{12} & \alpha_{13} & \alpha_{14} & \alpha_{15} \\ 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 & 0 \end{array}$$

which translates into the decimal number

$$2^3 + 2^6 + 2^{10} + 2^{13} = 9288.$$

We can thus encode a Slater determinant as a bit pattern.

With N particles that can be distributed over n single-particle states, the total number of Slater determinats (and defining thereby the dimensionality of the system) is

$$\dim(\mathcal{H}) = \binom{n}{N}.$$

The total number of bit patterns is 2^n .

We assume again that we have at our disposal n different single-particle orbits $\alpha_0, \alpha_2, \dots, \alpha_{n-1}$ and that we can distribute among these orbits $N \leq n$ particles. The ordering among these states is important as it defines the order of the creation operators. We will write the determinant

$$\Phi_A = a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

in a more compact way as

$$\Phi_{3,6,10,13} = |0001001000100100\rangle.$$

The action of a creation operator is thus

$$a_{\alpha_4}^\dagger \Phi_{3,6,10,13} = a_{\alpha_4}^\dagger |0001001000100100\rangle = a_{\alpha_4}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

which becomes

$$-a_{\alpha_3}^\dagger a_{\alpha_4}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle = -|0001101000100100\rangle.$$

Similarly

$$a_{\alpha_6}^\dagger \Phi_{3,6,10,13} = a_{\alpha_6}^\dagger |0001001000100100\rangle = a_{\alpha_6}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

which becomes

$$-a_{\alpha_4}^\dagger (a_{\alpha_6}^\dagger)^2 a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle = 0!$$

This gives a simple recipe:

- If one of the bits b_j is 1 and we act with a creation operator on this bit, we return a null vector
- If $b_j = 0$, we set it to 1 and return a sign factor $(-1)^l$, where l is the number of bits set before bit j .

Consider the action of $a_{\alpha_2}^\dagger$ on various slater determinants:

$$\begin{aligned} a_{\alpha_2}^\dagger \Phi_{00111} &= a_{\alpha_2}^\dagger |00111\rangle = 0 \times |00111\rangle \\ a_{\alpha_2}^\dagger \Phi_{01011} &= a_{\alpha_2}^\dagger |01011\rangle = (-1) \times |01111\rangle \\ a_{\alpha_2}^\dagger \Phi_{01101} &= a_{\alpha_2}^\dagger |01101\rangle = 0 \times |01101\rangle \\ a_{\alpha_2}^\dagger \Phi_{01110} &= a_{\alpha_2}^\dagger |01110\rangle = 0 \times |01110\rangle \\ a_{\alpha_2}^\dagger \Phi_{10011} &= a_{\alpha_2}^\dagger |10011\rangle = (-1) \times |10111\rangle \\ a_{\alpha_2}^\dagger \Phi_{10101} &= a_{\alpha_2}^\dagger |10101\rangle = 0 \times |10101\rangle \\ a_{\alpha_2}^\dagger \Phi_{10110} &= a_{\alpha_2}^\dagger |10110\rangle = 0 \times |10110\rangle \\ a_{\alpha_2}^\dagger \Phi_{11001} &= a_{\alpha_2}^\dagger |11001\rangle = (+1) \times |11101\rangle \\ a_{\alpha_2}^\dagger \Phi_{11010} &= a_{\alpha_2}^\dagger |11010\rangle = (+1) \times |11110\rangle \end{aligned}$$

What is the simplest way to obtain the phase when we act with one annihilation(creation) operator on the given Slater determinant representation?

We have an SD representation

$$\Phi_\Lambda = a_{\alpha_0}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

in a more compact way as

$$\Phi_{0,3,6,10,13} = |1001001000100100\rangle.$$

The action of

$$a_{\alpha_4}^\dagger a_{\alpha_0} \Phi_{0,3,6,10,13} = a_{\alpha_4}^\dagger |0001001000100100\rangle = a_{\alpha_4}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

which becomes

$$-a_{\alpha_3}^\dagger a_{\alpha_4}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle = -|0001101000100100\rangle.$$

The action

$$a_{\alpha_0} \Phi_{0,3,6,10,13} = |0001001000100100\rangle,$$

can be obtained by subtracting the logical sum (AND operation) of $\Phi_{0,3,6,10,13}$ and a word which represents only α_0 , that is

$$|1000000000000000\rangle,$$

from $\Phi_{0,3,6,10,13} = |1001001000100100\rangle$.

This operation gives $|0001001000100100\rangle$.

Similarly, we can form $a_{\alpha_4}^\dagger a_{\alpha_0} \Phi_{0,3,6,10,13}$, say, by adding $|0000100000000000\rangle$ to $a_{\alpha_0} \Phi_{0,3,6,10,13}$, first checking that their logical sum is zero in order to make sure that orbital α_4 is not already occupied.

It is trickier however to get the phase $(-1)^l$. One possibility is as follows

- Let S_1 be a word that represents the 1-bit to be removed and all others set to zero.

In the previous example $S_1 = |1000000000000000\rangle$

- Define S_2 as the similar word that represents the bit to be added, that is in our case

$S_2 = |0000100000000000\rangle$.

- Compute then $S = S_1 - S_2$, which here becomes

$$S = |0111000000000000\rangle$$

- Perform then the logical AND operation of S with the word containing

$$\Phi_{0,3,6,10,13} = |1001001000100100\rangle,$$

which results in $|0001000000000000\rangle$. Counting the number of 1-bits gives the phase. Here you need however an algorithm for bitcounting. Several efficient ones available.

Chapter 3

Full configuration interaction theory

Slater determinants as basis states, Repetition

The simplest possible choice for many-body wavefunctions are **product** wavefunctions. That is

$$\Psi(x_1, x_2, x_3, \dots, x_A) \approx \phi_1(x_1)\phi_2(x_2)\phi_3(x_3)\dots$$

because we are really only good at thinking about one particle at a time. Such product wavefunctions, without correlations, are easy to work with; for example, if the single-particle states $\phi_i(x)$ are orthonormal, then the product wavefunctions are easy to orthonormalize.

Similarly, computing matrix elements of operators are relatively easy, because the integrals factorize.

The price we pay is the lack of correlations, which we must build up by using many, many product wavefunctions. (Thus we have a trade-off: compact representation of correlations but difficult integrals versus easy integrals but many states required.)

Slater determinants as basis states, repetition

Because we have fermions, we are required to have antisymmetric wavefunctions, e.g.

$$\Psi(x_1, x_2, x_3, \dots, x_A) = -\Psi(x_2, x_1, x_3, \dots, x_A)$$

etc. This is accomplished formally by using the determinantal formalism

$$\Psi(x_1, x_2, \dots, x_A) = \frac{1}{\sqrt{A!}} \begin{vmatrix} \phi_1(x_1) & \phi_1(x_2) & \dots & \phi_1(x_A) \\ \phi_2(x_1) & \phi_2(x_2) & \dots & \phi_2(x_A) \\ \vdots & \vdots & \ddots & \vdots \\ \phi_A(x_1) & \phi_A(x_2) & \dots & \phi_A(x_A) \end{vmatrix}$$

Product wavefunction + antisymmetry = Slater determinant.

Slater determinants as basis states

$$\Psi(x_1, x_2, \dots, x_A) = \frac{1}{\sqrt{A!}} \begin{vmatrix} \phi_1(x_1) & \phi_1(x_2) & \dots & \phi_1(x_A) \\ \phi_2(x_1) & \phi_2(x_2) & \dots & \phi_2(x_A) \\ \vdots & \vdots & \ddots & \vdots \\ \phi_A(x_1) & \phi_A(x_2) & \dots & \phi_A(x_A) \end{vmatrix}$$

Properties of the determinant (interchange of any two rows or any two columns yields a change in sign; thus no two rows and no two columns can be the same) lead to the Pauli principle:

- No two particles can be at the same place (two columns the same); and
- No two particles can be in the same state (two rows the same).

Slater determinants as basis states

As a practical matter, however, Slater determinants beyond $N = 4$ quickly become unwieldy. Thus we turn to the **occupation representation** or **second quantization** to simplify calculations.

The occupation representation or number representation, using fermion **creation** and **annihilation** operators, is compact and efficient. It is also abstract and, at first encounter, not easy to internalize. It is inspired by other operator formalism, such as the ladder operators for the harmonic oscillator or for angular momentum, but unlike those cases, the operators **do not have coordinate space representations**.

Instead, one can think of fermion creation/annihilation operators as a game of symbols that compactly reproduces what one would do, albeit clumsily, with full coordinate-space Slater determinants.

Quick repetition of the occupation representation

We start with a set of orthonormal single-particle states $\{\phi_i(x)\}$. (Note: this requirement, and others, can be relaxed, but leads to a more involved formalism.) **Any** orthonormal set will do.

To each single-particle state $\phi_i(x)$ we associate a creation operator \hat{a}_i^\dagger and an annihilation operator \hat{a}_i .

When acting on the vacuum state $|0\rangle$, the creation operator \hat{a}_i^\dagger causes a particle to occupy the single-particle state $\phi_i(x)$:

$$\phi_i(x) \rightarrow \hat{a}_i^\dagger |0\rangle$$

Quick repetition of the occupation representation

But with multiple creation operators we can occupy multiple states:

$$\phi_i(x)\phi_j(x')\phi_k(x'') \rightarrow \hat{a}_i^\dagger \hat{a}_j^\dagger \hat{a}_k^\dagger |0\rangle.$$

Now we impose antisymmetry, by having the fermion operators satisfy **anticommutation relations**:

$$\hat{a}_i^\dagger \hat{a}_j^\dagger + \hat{a}_j^\dagger \hat{a}_i^\dagger = [\hat{a}_i^\dagger, \hat{a}_j^\dagger]_+ = \{\hat{a}_i^\dagger, \hat{a}_j^\dagger\} = 0$$

so that

$$\hat{a}_i^\dagger \hat{a}_j^\dagger = -\hat{a}_j^\dagger \hat{a}_i^\dagger$$

Quick repetition of the occupation representation

Because of this property, automatically $\hat{a}_i^\dagger \hat{a}_i^\dagger = 0$, enforcing the Pauli exclusion principle. Thus when writing a Slater determinant using creation operators,

$$\hat{a}_i^\dagger \hat{a}_j^\dagger \hat{a}_k^\dagger \dots |0\rangle$$

each index i, j, k, \dots must be unique.

For some relevant exercises with solutions see chapter 8 of Lecture Notes in Physics, volume 936.

Full Configuration Interaction Theory

We have defined the ansatz for the ground state as

$$|\Phi_0\rangle = \left(\prod_{i \leq F} \hat{a}_i^\dagger \right) |0\rangle,$$

where the index i defines different single-particle states up to the Fermi level. We have assumed that we have N fermions. A given one-particle-one-hole (1p1h) state can be written as

$$|\Phi_i^a\rangle = \hat{a}_a^\dagger \hat{a}_i |\Phi_0\rangle,$$

while a 2p2h state can be written as

$$|\Phi_{ij}^{ab}\rangle = \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_j \hat{a}_i |\Phi_0\rangle,$$

and a general $NpNh$ state as

$$|\Phi_{ijk\dots}^{abc\dots}\rangle = \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_c^\dagger \dots \hat{a}_k \hat{a}_j \hat{a}_i |\Phi_0\rangle.$$

Full Configuration Interaction Theory

We can then expand our exact state function for the ground state as

$$|\Psi_0\rangle = C_0|\Phi_0\rangle + \sum_{ai} C_i^a |\Phi_i^a\rangle + \sum_{abij} C_{ij}^{ab} |\Phi_{ij}^{ab}\rangle + \dots = (C_0 + \hat{C})|\Phi_0\rangle,$$

where we have introduced the so-called correlation operator

$$\hat{C} = \sum_{ai} C_i^a \hat{a}_a^\dagger \hat{a}_i + \sum_{abij} C_{ij}^{ab} \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_j \hat{a}_i + \dots$$

Since the normalization of Ψ_0 is at our disposal and since C_0 is by hypothesis non-zero, we may arbitrarily set $C_0 = 1$ with corresponding proportional changes in all other coefficients. Using this so-called intermediate normalization we have

$$\langle \Psi_0 | \Phi_0 \rangle = \langle \Phi_0 | \Phi_0 \rangle = 1,$$

resulting in

$$|\Psi_0\rangle = (1 + \hat{C})|\Phi_0\rangle.$$

Full Configuration Interaction Theory

We rewrite

$$|\Psi_0\rangle = C_0|\Phi_0\rangle + \sum_{ai} C_i^a |\Phi_i^a\rangle + \sum_{abij} C_{ij}^{ab} |\Phi_{ij}^{ab}\rangle + \dots,$$

in a more compact form as

$$|\Psi_0\rangle = \sum_{PH} C_H^P \Phi_H^P = \left(\sum_{PH} C_H^P \hat{A}_H^P \right) |\Phi_0\rangle,$$

where H stands for $0, 1, \dots, n$ hole states and P for $0, 1, \dots, n$ particle states. Our requirement of unit normalization gives

$$\langle \Psi_0 | \Phi_0 \rangle = \sum_{PH} |C_H^P|^2 = 1,$$

and the energy can be written as

$$E = \langle \Psi_0 | \hat{H} | \Phi_0 \rangle = \sum_{PP'HH'} C_H^{*P} \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle C_{H'}^{P'}.$$

Full Configuration Interaction Theory

Normally

$$E = \langle \Psi_0 | \hat{H} | \Phi_0 \rangle = \sum_{PP'HH'} C_H^{*P} \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle C_{H'}^{P'},$$

is solved by diagonalization setting up the Hamiltonian matrix defined by the basis of all possible Slater determinants. A diagonalization is equivalent to finding the variational minimum of

$$\langle \Psi_0 | \hat{H} | \Phi_0 \rangle - \lambda \langle \Psi_0 | \Phi_0 \rangle,$$

where λ is a variational multiplier to be identified with the energy of the system. The minimization process results in

$$\delta [\langle \Psi_0 | \hat{H} | \Phi_0 \rangle - \lambda \langle \Psi_0 | \Phi_0 \rangle] = \sum_{P'H'} \left\{ \delta[C_H^{*P}] \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle C_{H'}^{P'} + C_H^{*P} \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle \delta[C_{H'}^{P'}] - \lambda (\delta[C_H^{*P}] C_{H'}^{P'} + C_H^{*P} \delta[C_{H'}^{P'}]) \right\} = 0.$$

Since the coefficients $\delta[C_H^{*P}]$ and $\delta[C_{H'}^{P'}]$ are complex conjugates it is necessary and sufficient to require the quantities that multiply with $\delta[C_H^{*P}]$ to vanish.

Full Configuration Interaction Theory

This leads to

$$\sum_{P'H'} \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle C_{H'}^{P'} - \lambda C_H^P = 0,$$

for all sets of P and H .

If we then multiply by the corresponding C_H^{*P} and sum over PH we obtain

$$\sum_{PP'HH'} C_H^{*P} \langle \Phi_H^P | \hat{H} | \Phi_{H'}^{P'} \rangle C_{H'}^{P'} - \lambda \sum_{PH} |C_H^P|^2 = 0,$$

leading to the identification $\lambda = E$. This means that we have for all PH sets

$$\sum_{P'H'} \langle \Phi_H^P | \hat{H} - E | \Phi_{H'}^{P'} \rangle = 0. \quad (3.1)$$

Full Configuration Interaction Theory

An alternative way to derive the last equation is to start from

$$(\hat{H} - E) |\Psi_0\rangle = (\hat{H} - E) \sum_{P'H'} C_{H'}^{P'} |\Phi_{H'}^{P'}\rangle = 0,$$

and if this equation is successively projected against all Φ_H^P in the expansion of Ψ , then the last equation on the previous slide results. As stated previously, one solves this equation normally by diagonalization. If we are able to solve this equation exactly (that is numerically exactly) in a large Hilbert space (it will be truncated in terms of the number of single-particle states included in the definition of Slater determinants), it can then serve as a benchmark for other many-body methods which approximate the correlation operator \hat{C} .

Example of a Hamiltonian matrix

Suppose, as an example, that we have six fermions below the Fermi level. This means that we can make at most $6p - 6h$ excitations. If we have an infinity of single particle states above the Fermi level, we will obviously have an infinity of say $2p - 2h$ excitations. Each such way to configure the particles is called a **configuration**. We will always have to truncate in the basis of single-particle states. This gives us a finite number of possible Slater determinants. Our Hamiltonian matrix would then look like (where each block can have a large dimensionalities):

	$0p-0h$	$1p-1h$	$2p-2h$	$3p-3h$	$4p-4h$	$5p-5h$	$6p-6h$
$0p-0h$	x	x	x	0	0	0	0
$1p-1h$	x	x	x	x	0	0	0
$2p-2h$	x	x	x	x	x	0	0
$3p-3h$	0	x	x	x	x	x	0
$4p-4h$	0	0	x	x	x	x	x
$5p-5h$	0	0	0	x	x	x	x
$6p-6h$	0	0	0	0	x	x	x

with a two-body force. Why are there non-zero blocks of elements?

Example of a Hamiltonian matrix with a Hartree-Fock basis

If we use a Hartree-Fock basis, this corresponds to a particular unitary transformation where matrix elements of the type $\langle 0p-0h|\hat{H}|1p-1h\rangle = \langle \Phi_0|\hat{H}|\Phi_i^a\rangle = 0$ and our Hamiltonian matrix becomes

	$0p-0h$	$1p-1h$	$2p-2h$	$3p-3h$	$4p-4h$	$5p-5h$	$6p-6h$
$0p-0h$	\tilde{x}	0	\tilde{x}	0	0	0	0
$1p-1h$	0	\tilde{x}	\tilde{x}	\tilde{x}	0	0	0
$2p-2h$	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}	0	0
$3p-3h$	0	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}	0
$4p-4h$	0	0	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}
$5p-5h$	0	0	0	\tilde{x}	\tilde{x}	\tilde{x}	\tilde{x}
$6p-6h$	0	0	0	0	\tilde{x}	\tilde{x}	\tilde{x}

Shell-model jargon

If we do not make any truncations in the possible sets of Slater determinants (many-body states) we can make by distributing A nucleons among n single-particle states, we call such a calculation for **Full configuration interaction theory**

If we make truncations, we have different possibilities

- The standard nuclear shell-model. Here we define an effective Hilbert space with respect to a given core. The calculations are normally then performed for all many-body states that can be constructed from the effective Hilbert spaces. This approach requires a properly defined effective Hamiltonian

- We can truncate in the number of excitations. For example, we can limit the possible Slater determinants to only $1p - 1h$ and $2p - 2h$ excitations. This is called a configuration interaction calculation at the level of singles and doubles excitations, or just CISD.
- We can limit the number of excitations in terms of the excitation energies. If we do not define a core, this defines normally what is called the no-core shell-model approach.

What happens if we have a three-body interaction and a Hartree-Fock basis?

FCI and the exponential growth

Full configuration interaction theory calculations provide in principle, if we can diagonalize numerically, all states of interest. The dimensionality of the problem explodes however quickly.

The total number of Slater determinants which can be built with say N neutrons distributed among n single particle states is

$$\binom{n}{N} = \frac{n!}{(n-N)!N!}.$$

For a model space which comprises the first for major shells only $0s$, $0p$, $1s0d$ and $1p0f$ we have 40 single particle states for neutrons and protons. For the eight neutrons of oxygen-16 we would then have

$$\binom{40}{8} = \frac{40!}{(32)!8!} \sim 10^9,$$

and multiplying this with the number of proton Slater determinants we end up with approximately with a dimensionality d of $d \sim 10^{18}$.

Exponential wall

This number can be reduced if we look at specific symmetries only. However, the dimensionality explodes quickly!

- For Hamiltonian matrices of dimensionalities which are smaller than $d \sim 10^5$, we would use so-called direct methods for diagonalizing the Hamiltonian matrix
- For larger dimensionalities iterative eigenvalue solvers like Lanczos' method are used. The most efficient codes at present can handle matrices of $d \sim 10^{10}$.

A non-practical way of solving the eigenvalue problem

To see this, we look at the contributions arising from

$$\langle \Phi_H^P | = \langle \Phi_0 |$$

in Eq. (3.1), that is we multiply with $\langle \Phi_0 |$ from the left in

$$(\hat{H} - E) \sum_{P'H'} C_{H'}^{P'} |\Phi_{H'}^{P'}\rangle = 0.$$

If we assume that we have a two-body operator at most, Slater's rule gives then an equation for the correlation energy in terms of C_i^a and C_{ij}^{ab} only. We get then

$$\langle \Phi_0 | \hat{H} - E | \Phi_0 \rangle + \sum_{ai} \langle \Phi_0 | \hat{H} - E | \Phi_i^a \rangle C_i^a + \sum_{abij} \langle \Phi_0 | \hat{H} - E | \Phi_{ij}^{ab} \rangle C_{ij}^{ab} = 0,$$

or

$$E - E_0 = \Delta E = \sum_{ai} \langle \Phi_0 | \hat{H} | \Phi_i^a \rangle C_i^a + \sum_{abij} \langle \Phi_0 | \hat{H} | \Phi_{ij}^{ab} \rangle C_{ij}^{ab},$$

where the energy E_0 is the reference energy and ΔE defines the so-called correlation energy. The single-particle basis functions could be the results of a Hartree-Fock calculation or just the eigenstates of the non-interacting part of the Hamiltonian.

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$$\langle \Phi_0 | \hat{H} - E | \Phi_0 \rangle + \sum_{ai} \langle \Phi_0 | \hat{H} - E | \Phi_i^a \rangle C_i^a + \sum_{abij} \langle \Phi_0 | \hat{H} - E | \Phi_{ij}^{ab} \rangle C_{ij}^{ab} = 0,$$

or

$$E - E_0 = \Delta E = \sum_{ai} \langle \Phi_0 | \hat{H} | \Phi_i^a \rangle C_i^a + \sum_{abij} \langle \Phi_0 | \hat{H} | \Phi_{ij}^{ab} \rangle C_{ij}^{ab},$$

where the energy E_0 is the reference energy and ΔE defines the so-called correlation energy. The single-particle basis functions could be the results of a Hartree-Fock calculation or just the eigenstates of the non-interacting part of the Hamiltonian.

Rewriting the FCI equation

In our notes on Hartree-Fock calculations, we have already computed the matrix $\langle \Phi_0 | \hat{H} | \Phi_i^a \rangle$ and $\langle \Phi_0 | \hat{H} | \Phi_{ij}^{ab} \rangle$. If we are using a Hartree-Fock basis, then the matrix elements $\langle \Phi_0 | \hat{H} | \Phi_i^a \rangle = 0$ and we are left with a *correlation energy* given by

$$E - E_0 = \Delta E^{HF} = \sum_{abij} \langle \Phi_0 | \hat{H} | \Phi_{ij}^{ab} \rangle C_{ij}^{ab}.$$

Rewriting the FCI equation

Inserting the various matrix elements we can rewrite the previous equation as

$$\Delta E = \sum_{ai} \langle i | \hat{f} | a \rangle C_i^a + \sum_{abij} \langle ij | \hat{v} | ab \rangle C_{ij}^{ab}.$$

This equation determines the correlation energy but not the coefficients C .

Rewriting the FCI equation, does not stop here

We need more equations. Our next step is to set up

$$\langle \Phi_i^a | \hat{H} - E | \Phi_0 \rangle + \sum_{bj} \langle \Phi_i^a | \hat{H} - E | \Phi_j^b \rangle C_j^b + \sum_{bcjk} \langle \Phi_i^a | \hat{H} - E | \Phi_{jk}^{bc} \rangle C_{jk}^{bc} + \sum_{bcdjkl} \langle \Phi_i^a | \hat{H} - E | \Phi_{jkl}^{bcd} \rangle C_{jkl}^{bcd} = 0,$$

as this equation will allow us to find an expression for the coefficients C_i^a since we can rewrite this equation as

$$\langle i | \hat{f} | a \rangle + \langle \Phi_i^a | \hat{H} | \Phi_i^a \rangle C_i^a + \sum_{bj \neq ai} \langle \Phi_i^a | \hat{H} | \Phi_j^b \rangle C_j^b + \sum_{bcjk} \langle \Phi_i^a | \hat{H} | \Phi_{jk}^{bc} \rangle C_{jk}^{bc} + \sum_{bcdjkl} \langle \Phi_i^a | \hat{H} | \Phi_{jkl}^{bcd} \rangle C_{jkl}^{bcd} = E C_i^a.$$

Rewriting the FCI equation, please stop here

We see that on the right-hand side we have the energy E . This leads to a non-linear equation in the unknown coefficients. These equations are normally solved iteratively (that is we can start with a guess for the coefficients C_i^a). A common choice is to use perturbation theory for the first guess, setting thereby

$$C_i^a = \frac{\langle i | \hat{f} | a \rangle}{\epsilon_i - \epsilon_a}.$$

Rewriting the FCI equation, more to add

The observant reader will however see that we need an equation for C_{jk}^{bc} and C_{jkl}^{bcd} as well. To find equations for these coefficients we need then to continue our multiplications from the left with the various Φ_H^P terms.

For C_{jk}^{bc} we need then

$$\begin{aligned} & \langle \Phi_{ij}^{ab} | \hat{H} - E | \Phi_0 \rangle + \sum_{kc} \langle \Phi_{ij}^{ab} | \hat{H} - E | \Phi_k^c \rangle C_k^c + \\ & \sum_{cdkl} \langle \Phi_{ij}^{ab} | \hat{H} - E | \Phi_{kl}^{cd} \rangle C_{kl}^{cd} + \sum_{cdekml} \langle \Phi_{ij}^{ab} | \hat{H} - E | \Phi_{klm}^{cde} \rangle C_{klm}^{cde} + \sum_{cdefklmn} \langle \Phi_{ij}^{ab} | \hat{H} - E | \Phi_{klmn}^{cdef} \rangle C_{klmn}^{cdef} = 0, \end{aligned}$$

and we can isolate the coefficients C_{kl}^{cd} in a similar way as we did for the coefficients C_i^a .

Rewriting the FCI equation, more to add

A standard choice for the first iteration is to set

$$C_{ij}^{ab} = \frac{\langle ij | \hat{v} | ab \rangle}{\epsilon_i + \epsilon_j - \epsilon_a - \epsilon_b}.$$

At the end we can rewrite our solution of the Schroedinger equation in terms of n coupled equations for the coefficients C_H^P . This is a very cumbersome way of solving the equation. However, by using this iterative scheme we can illustrate how we can compute the various terms in the wave operator or correlation operator \hat{C} . We will later identify the calculation of the various terms C_H^P as parts of different many-body approximations to full CI. In particular, we can relate this non-linear scheme with Coupled Cluster theory and many-body perturbation theory.

Summarizing FCI and bringing in approximative methods

If we can diagonalize large matrices, FCI is the method of choice since:

- It gives all eigenvalues, ground state and excited states
- The eigenvectors are obtained directly from the coefficients C_H^P which result from the diagonalization
- We can compute easily expectation values of other operators, as well as transition probabilities
- Correlations are easy to understand in terms of contributions to a given operator beyond the Hartree-Fock contribution. This is the standard approach in many-body theory.

Definition of the correlation energy

The correlation energy is defined as, with a two-body Hamiltonian,

$$\Delta E = \sum_{ai} \langle i | \hat{f} | a \rangle C_i^a + \sum_{abij} \langle ij | \hat{v} | ab \rangle C_{ij}^{ab}.$$

The coefficients C result from the solution of the eigenvalue problem. The energy of say the ground state is then

$$E = E_{ref} + \Delta E,$$

where the so-called reference energy is the energy we obtain from a Hartree-Fock calculation, that is

$$E_{ref} = \langle \Phi_0 | \hat{H} | \Phi_0 \rangle.$$

FCI equation and the coefficients

However, as we have seen, even for a small case like the four first major shells and a nucleus like oxygen-16, the dimensionality becomes quickly intractable. If we wish to include single-particle states that reflect weakly bound systems, we need a much larger single-particle basis. We need thus approximative methods that sum specific correlations to infinite order.

Popular methods are

- Many-body perturbation theory (in essence a Taylor expansion)
- Coupled cluster theory (coupled non-linear equations)
- Green's function approaches (matrix inversion)
- Similarity group transformation methods (coupled ordinary differential equations)

All these methods start normally with a Hartree-Fock basis as the calculational basis.

Important ingredients to have in codes

- Be able to validate and verify the algorithms.
- Include concepts like unit testing. Gives the possibility to test and validate several or all parts of the code.
- Validation and verification are then included *naturally* and one can develop a better attitude to what is meant with an ethically sound scientific approach.

A structured approach to solving problems

In the steps that lead to the development of clean code you should think of

1. How to structure a code in terms of functions (use IDEs or advanced text editors like sublime or atom)
2. How to make a module
3. How to read input data flexibly from the command line or files
4. How to create graphical/web user interfaces
5. How to write unit tests
6. How to refactor code in terms of classes (instead of functions only)
7. How to conduct and automate large-scale numerical experiments
8. How to write scientific reports in various formats (L^AT_EX, HTML, doconce)

Additional benefits

Many of the above aspects will save you a lot of time when you incrementally extend software over time from simpler to more complicated problems. In particular, you will benefit from many good habits:

1. New code is added in a modular fashion to a library (modules)
2. Programs are run through convenient user interfaces
3. It takes one quick command to let all your code undergo heavy testing
4. Tedious manual work with running programs is automated,
5. Your scientific investigations are reproducible, scientific reports with top quality typesetting are produced both for paper and electronic devices. Use version control software like git and repositories like github

Unit Testing

Unit Testing is the practice of testing the smallest testable parts, called units, of an application individually and independently to determine if they behave exactly as expected.

Unit tests (short code fragments) are usually written such that they can be preformed at any time during the development to continually verify the behavior of the code.

In this way, possible bugs will be identified early in the development cycle, making the debugging at later stages much easier.

Unit Testing, benefits

There are many benefits associated with Unit Testing, such as

- It increases confidence in changing and maintaining code. Big changes can be made to the code quickly, since the tests will ensure that everything still is working properly.
- Since the code needs to be modular to make Unit Testing possible, the code will be easier to reuse. This improves the code design.
- Debugging is easier, since when a test fails, only the latest changes need to be debugged.
 - Different parts of a project can be tested without the need to wait for the other parts to be available.
- A unit test can serve as a documentation on the functionality of a unit of the code.

Simple example of unit test

Look up the guide on how to install unit tests for c++ at course webpage. This is the version with classes.

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
include <unittest++/UnitTest++.h>
class MyMultiplyClass public: double multiply(double x, double y) return x * y; ;
TEST(MyMath) MyMultiplyClass my; CHECK_EQUAL(56,my.multiply(7,8));
int main() return UnitTest::RunAllTests();
```

Simple example of unit test

And without classes

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
include <unittest++/UnitTest++.h>
double multiply(double x, double y) return x * y;
TEST(MyMath) CHECK_EQUAL(56,multiply(7,8));
int main() return UnitTest::RunAllTests();
```

For Fortran users, the link at <http://sourceforge.net/projects/fortranxunit/> contains a similar software for unit testing. For Python go to <https://docs.python.org/2/library/unittest.html>.

Unit tests

There are many types of **unit test** libraries. One which is very popular with C++ programmers is Catch

Catch is header only. All you need to do is drop the file(s) somewhere reachable from your project - either in some central location you can set your header search path to find, or directly into your project tree itself!

This is a particularly good option for other Open-Source projects that want to use Catch for their test suite.

Examples

Computing factorials

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
inline unsigned int Factorial( unsigned int number ) return number > 1 ? Factorial(number-
1)*number : 1;
```

Factorial Example

Simple test where we put everything in a single file

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
define CATCH_CONFIG_MAIN // This tells Catch to provide a main() and include "catch.hpp" inline unsigned int Factorial( unsigned int number ) return
TEST_CASE("Factorials are computed", "[factorial]") REQUIRE( Factorial(0) == 1 ); REQUIRE( Factorial(1) == 1 ); REQUIRE( Factorial(2) == 2 ); REQUIRE( Factorial(3) == 6 ); REQUIRE( Factorial(4) == 24 );
```

This will compile to a complete executable which responds to command line arguments. If you just run it with no arguments it will execute all test cases (in this case there is just one), report any failures, report a summary of how many tests passed and failed and return the number of failed tests.

What did we do (1)?

All we did was

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
define
```

one identifier and

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
include
```

one header and we got everything - even an implementation of main() that will respond to command line arguments. Once you have more than one file with unit tests in you'll just need to

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
include "catch.hpp"
```

and go. Usually it's a good idea to have a dedicated implementation file that just has

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
define CATCH_CONFIG_MAINinclude"catch.hpp".
```

You can also provide your own implementation of main and drive Catch yourself.

What did we do (2)?

We introduce test cases with the

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
TEST_CASE
macro.
```

The test name must be unique. You can run sets of tests by specifying a wildcarded test name or a tag expression. All we did was **define** one identifier and **include** one header and we got everything.

We write our individual test assertions using the

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]c++
REQUIRE
macro.
```

Unit test summary and testing approach

Three levels of tests

1. Microscopic level: testing small parts of code, use often unit test libraries
2. Mesoscopic level: testing the integration of various parts of your code
3. Macroscopic level: testing that the final result is ok

Coding Recommendations

Writing clean and clear code is an art and reflects your understanding of

1. derivation, verification, and implementation of algorithms
2. what can go wrong with algorithms
3. overview of important, known algorithms
4. how algorithms are used to solve mathematical problems
5. reproducible science and ethics
6. algorithmic thinking for gaining deeper insights about scientific problems

Computing is understanding and your understanding is reflected in your abilities to write clear and clean code.

Summary and recommendations

Some simple hints and tips in order to write clean and clear code

1. Spell out the algorithm and have a top-down approach to the flow of data
2. Start with coding as close as possible to eventual mathematical expressions
3. Use meaningful names for variables
4. Split tasks in simple functions and modules/classes
5. Functions should return as few as possible variables
6. Use unit tests and make sure your codes are producing the correct results
7. Where possible use symbolic coding to autogenerate code and check results
8. Make a proper timing of your algorithms
9. Use version control and make your science reproducible
10. Use IDEs or smart editors with debugging and analysis tools.
11. Automatize your computations interfacing high-level and compiled languages like C++ and Fortran.
12.

Building a many-body basis

Here we will discuss how we can set up a single-particle basis which we can use in the various parts of our projects, from the simple pairing model to infinite nuclear matter. We will use here the simple pairing model to illustrate in particular how to set up a single-particle basis. We will also use this to discuss standard FCI approaches like:

1. Standard shell-model basis in one or two major shells
2. Full CI in a given basis and no truncations
3. CISD and CISDT approximations
4. No-core shell model and truncation in excitation energy

Building a many-body basis

An important step in an FCI code is to construct the many-body basis.

While the formalism is independent of the choice of basis, the **effectiveness** of a calculation will certainly be basis dependent.

Furthermore there are common conventions useful to know.

First, the single-particle basis has angular momentum as a good quantum number. You can imagine the single-particle wavefunctions being generated by a one-body Hamiltonian, for example a harmonic oscillator. Modifications include harmonic oscillator plus spin-orbit splitting, or self-consistent mean-field potentials, or the Woods-Saxon potential which mocks up the self-consistent mean-field. For nuclei, the harmonic oscillator, modified by spin-orbit splitting, provides a useful language for describing single-particle states.

Building a many-body basis

Each single-particle state is labeled by the following quantum numbers:

- Orbital angular momentum l
- Intrinsic spin $s = 1/2$ for protons and neutrons
- Angular momentum $j = l \pm 1/2$
- z -component j_z (or m)
- Some labeling of the radial wavefunction, typically n the number of nodes in the radial wavefunction, but in the case of harmonic oscillator one can also use the principal quantum number N , where the harmonic oscillator energy is $(N + 3/2)\omega$.

In this format one labels states by $n(l)_j$, with (l) replaced by a letter: s for $l = 0$, p for $l = 1$, d for $l = 2$, f for $l = 3$, and thenceforth alphabetical.

Building a many-body basis

In practice the single-particle space has to be severely truncated. This truncation is typically based upon the single-particle energies, which is the effective energy from a mean-field potential.

Sometimes we freeze the core and only consider a valence space. For example, one may assume a frozen ^4He core, with two protons and two neutrons in the $0s_{1/2}$ shell, and then only allow active particles in the $0p_{1/2}$ and $0p_{3/2}$ orbits.

Another example is a frozen ^{16}O core, with eight protons and eight neutrons filling the $0s_{1/2}$, $0p_{1/2}$ and $0p_{3/2}$ orbits, with valence particles in the $0d_{5/2}$, $1s_{1/2}$ and $0d_{3/2}$ orbits.

Sometimes we refer to nuclei by the valence space where their last nucleons go. So, for example, we call ^{12}C a p -shell nucleus, while ^{26}Al is an sd -shell nucleus and ^{56}Fe is a pf -shell nucleus.

Building a many-body basis

There are different kinds of truncations.

- For example, one can start with ‘filled’ orbits (almost always the lowest), and then allow one, two, three... particles excited out of those filled orbits. These are called 1p-1h, 2p-2h, 3p-3h excitations.
- Alternately, one can state a maximal orbit and allow all possible configurations with particles occupying states up to that maximum. This is called *full configuration*.
- Finally, for particular use in nuclear physics, there is the *energy* truncation, also called the $N\Omega$ or N_{\max} truncation.

Building a many-body basis

Here one works in a harmonic oscillator basis, with each major oscillator shell assigned a principal quantum number $N = 0, 1, 2, 3, \dots$. The $N\Omega$ or N_{max} truncation: Any configuration is given an noninteracting energy, which is the sum of the single-particle harmonic oscillator energies. (Thus this ignores spin-orbit splitting.)

Excited state are labeled relative to the lowest configuration by the number of harmonic oscillator quanta.

This truncation is useful because if one includes *all* configuration up to some N_{max} , and has a translationally invariant interaction, then the intrinsic motion and the center-of-mass motion factor. In other words, we can know exactly the center-of-mass wavefunction.

In almost all cases, the many-body Hamiltonian is rotationally invariant. This means it commutes with the operators \hat{J}^2, \hat{J}_z and so eigenstates will have good J, M . Furthermore, the eigenenergies do not depend upon the orientation M .

Therefore we can choose to construct a many-body basis which has fixed M ; this is called an M -scheme basis.

Alternately, one can construct a many-body basis which has fixed J , or a J -scheme basis.

Building a many-body basis

The Hamiltonian matrix will have smaller dimensions (a factor of 10 or more) in the J -scheme than in the M -scheme. On the other hand, as we'll show in the next slide, the M -scheme is very easy to construct with Slater determinants, while the J -scheme basis states, and thus the matrix elements, are more complicated, almost always being linear combinations of M -scheme states. J -scheme bases are important and useful, but we'll focus on the simpler M -scheme.

The quantum number m is additive (because the underlying group is Abelian): if a Slater determinant $\hat{a}_i^\dagger \hat{a}_j^\dagger \hat{a}_k^\dagger \dots |0\rangle$ is built from single-particle states all with good m , then the total

$$M = m_i + m_j + m_k + \dots$$

This is *not* true of J , because the angular momentum group $SU(2)$ is not Abelian.

Building a many-body basis

The upshot is that

- It is easy to construct a Slater determinant with good total M ;
- It is trivial to calculate M for each Slater determinant;
- So it is easy to construct an M -scheme basis with fixed total M .

Note that the individual M -scheme basis states will *not*, in general, have good total J . Because the Hamiltonian is rotationally invariant, however, the eigenstates will have good J . (The situation is muddled when one has states of different J that are nonetheless degenerate.)

Building a many-body basis

Example: two $j = 1/2$ orbits

Index	n	l	j	m_j
1	0	0	1/2	-1/2
2	0	0	1/2	1/2
3	1	0	1/2	-1/2
4	1	0	1/2	1/2

Note that the order is arbitrary.

Building a many-body basis

There are $\binom{4}{2} = 6$ two-particle states, which we list with the total M :

Occupied	M
1,2	0
1,3	-1
1,4	0
2,3	0
2,4	1
3,4	0

There are 4 states with $M = 0$, and 1 each with $M = \pm 1$.

Building a many-body basis

As another example, consider using only single particle states from the $0d_{5/2}$ space. They have the following quantum numbers

Index	n	l	j	m_j
1	0	2	5/2	-5/2
2	0	2	5/2	-3/2
3	0	2	5/2	-1/2
4	0	2	5/2	1/2
5	0	2	5/2	3/2
6	0	2	5/2	5/2

Building a many-body basis

There are $\binom{6}{2} = 15$ two-particle states, which we list with the total M :

Occupied M	Occupied M	Occupied M
1,2 -4	2,3 -2	3,5 1
1,3 -3	2,4 -1	3,6 2
1,4 -2	2,5 0	4,5 2
1,5 -1	2,6 1	4,6 3
1,6 0	3,4 0	5,6 4

There are 3 states with $M = 0$, 2 with $M = 1$, and so on.

Shell-model project

The first step is to construct the M -scheme basis of Slater determinants. Here M -scheme means the total J_z of the many-body states is fixed.

The steps could be:

- Read in a user-supplied file of single-particle states (examples can be given) or just code these internally;
- Ask for the total M of the system and the number of particles N ;
- Construct all the N -particle states with given M . You will validate the code by comparing both the number of states and specific states.

Shell-model project

The format of a possible input file could be

Index	n	l	$2j$	$2m_j$
1	1	0	1	-1
2	1	0	1	1
3	0	2	3	-3
4	0	2	3	-1
5	0	2	3	1
6	0	2	3	3
7	0	2	5	-5
8	0	2	5	-3
9	0	2	5	-1
10	0	2	5	1
11	0	2	5	3
12	0	2	5	5

This represents the $1s_{1/2}0d_{3/2}0d_{5/2}$ valence space, or just the sd -space. There are twelve single-particle states, labeled by an overall index, and which have associated quantum numbers the

number of radial nodes, the orbital angular momentum l , and the angular momentum j and third component j_z . To keep everything as integers, we could store $2 \times j$ and $2 \times j_z$.

Shell-model project

To read in the single-particle states you need to:

- Open the file
 - Read the number of single-particle states (in the above example, 12); allocate memory; all you need is a single array storing $2 \times j_z$ for each state, labeled by the index.
- Read in the quantum numbers and store $2 \times j_z$ (and anything else you happen to want).

Shell-model project

The next step is to read in the number of particles N and the fixed total M (or, actually, $2 \times M$). For this project we assume only a single species of particles, say neutrons, although this can be relaxed. **Note:** Although it is often a good idea to try to write a more general code, given the short time allotted we would suggest you keep your ambition in check, at least in the initial phases of the project.

You should probably write an error trap to make sure N and M are congruent; if N is even, then $2 \times M$ should be even, and if N is odd then $2 \times M$ should be odd.

Shell-model project

The final step is to generate the set of N -particle Slater determinants with fixed M . The Slater determinants will be stored in occupation representation. Although in many codes this representation is done compactly in bit notation with ones and zeros, but for greater transparency and simplicity we will list the occupied single particle states.

Hence we can store the Slater determinant basis states as $sd(i, j)$, that is an array of dimension N_{SD} , the number of Slater determinants, by N , the number of occupied state. So if for the 7th Slater determinant the 2nd, 3rd, and 9th single-particle states are occupied, then $sd(7, 1) = 2$, $sd(7, 2) = 3$, and $sd(7, 3) = 9$.

Shell-model project

We can construct an occupation representation of Slater determinants by the *odometer* method. Consider $N_{sp} = 12$ and $N = 4$. Start with the first 4 states occupied, that is:

- $sd(1,:) = 1, 2, 3, 4$ (also written as $|1, 2, 3, 4\rangle$)

Now increase the last occupancy recursively:

- $sd(2,:) = 1, 2, 3, 5$
- $sd(3,:) = 1, 2, 3, 6$
- $sd(4,:) = 1, 2, 3, 7$
- ...
- $sd(9,:) = 1, 2, 3, 12$

Then start over with

- $sd(10,:) = 1, 2, 4, 5$

and again increase the rightmost digit

- $sd(11,:) = 1, 2, 4, 6$
- $sd(12,:) = 1, 2, 4, 7$
- ...
- $sd(17,:) = 1, 2, 4, 12$

Shell-model project

When we restrict ourselves to an M -scheme basis, we could choose two paths. The first is simplest (and simplest is often best, at least in the first draft of a code): generate all possible Slater determinants, and then extract from this initial list a list of those Slater determinants with a given M . (You will need to write a short function or routine that computes M for any given occupation.)

Alternately, and not too difficult, is to run the odometer routine twice: each time, as a Slater determinant is calculated, compute M , but do not store the Slater determinants except the current one. You can then count up the number of Slater determinants with a chosen M . Then allocated storage for the Slater determinants, and run the odometer algorithm again, this time storing Slater determinants with the desired M (this can be done with a simple logical flag).

Shell-model project

Some example solutions: Let's begin with a simple case, the $0d_{5/2}$ space containing six single-particle states

Index	n	l	j	m_j
1	0	2	5/2	-5/2
2	0	2	5/2	-3/2
3	0	2	5/2	-1/2
4	0	2	5/2	1/2
5	0	2	5/2	3/2
6	0	2	5/2	5/2

For two particles, there are a total of 15 states, which we list here with the total M :

- $|1,2\rangle, M = -4, |1,3\rangle, M = -3$
- $|1,4\rangle, M = -2, |1,5\rangle, M = -1$
- $|1,5\rangle, M = 0, |2,3\rangle, M = -2$
- $|2,4\rangle, M = -1, |2,5\rangle, M = 0$
- $|2,6\rangle, M = 1, |3,4\rangle, M = 0$
- $|3,5\rangle, M = 1, |3,6\rangle, M = 2$
- $|4,5\rangle, M = 2, |4,6\rangle, M = 3$
- $|5,6\rangle, M = 4$

Of these, there are only 3 states with $M = 0$.

Shell-model project

You should try by hand to show that in this same single-particle space, that for $N = 3$ there are 3 states with $M = 1/2$ and for $N = 4$ there are also only 3 states with $M = 0$.

To test your code, confirm the above.

Also, for the sd -space given above, for $N = 2$ there are 14 states with $M = 0$, for $N = 3$ there are 37 states with $M = 1/2$, for $N = 4$ there are 81 states with $M = 0$.

Shell-model project

For our project, we will only consider the pairing model. A simple space is the $(1/2)^2$ space with four single-particle states

Index	n	l	s	m_s
1	0	0	1/2	-1/2
2	0	0	1/2	1/2
3	1	0	1/2	-1/2
4	1	0	1/2	1/2

For $N = 2$ there are 4 states with $M = 0$; show this by hand and confirm your code reproduces it.

Shell-model project

Another, slightly more challenging space is the $(1/2)^4$ space, that is, with eight single-particle states we have

Index	n	l	s	m_s
1	0	0	1/2	-1/2
2	0	0	1/2	1/2
3	1	0	1/2	-1/2
4	1	0	1/2	1/2
5	2	0	1/2	-1/2
6	2	0	1/2	1/2
7	3	0	1/2	-1/2
8	3	0	1/2	1/2

For $N = 2$ there are 16 states with $M = 0$; for $N = 3$ there are 24 states with $M = 1/2$, and for $N = 4$ there are 36 states with $M = 0$.

Shell-model project

In the shell-model context we can interpret this as 4 $s_{1/2}$ levels, with $m = \pm 1/2$, we can also think of these as simple four pairs, $\pm k, k = 1, 2, 3, 4$. Later on we will assign single-particle energies, depending on the radial quantum number n , that is, $\epsilon_k = |k|\delta$ so that they are equally spaced.

Shell-model project

For application in the pairing model we can go further and consider only states with no “broken pairs,” that is, if $+k$ is filled (or $m = +1/2$, so is $-k$ ($m = -1/2$). If you want, you can write your code to accept only these, and obtain the following six states:

- $|1, 2, 3, 4\rangle$,
- $|1, 2, 5, 6\rangle$,
- $|1, 2, 7, 8\rangle$,
- $|3, 4, 5, 6\rangle$,
- $|3, 4, 7, 8\rangle$,
- $|5, 6, 7, 8\rangle$

Shell-model project

Hints for coding.

- Write small modules (routines/functions) ; avoid big functions that do everything. (But not too small.)
- Use Unit tests! Write lots of error traps, even for things that are ‘obvious.’
- Document as you go along. The Unit tests serve as documentation. For each function write a header that includes:
 1. Main purpose of function and/or unit test
 2. names and brief explanation of input variables, if any
 3. names and brief explanation of output variables, if any
 4. functions called by this function
 5. called by which functions

Shell-model project

Hints for coding

- Unit tests will save time. Use also IDEs for debugging. If you insist on brute force debugging, print out intermediate values. It’s almost impossible to debug a code by looking at it—the code will almost always win a ‘staring contest.’
- Validate code with SIMPLE CASES. Validate early and often. Unit tests!!

The number one mistake is using a too complex a system to test. For example , if you are computing particles in a potential in a box, try removing the potential—you should get particles in a box. And start with one particle, then two, then three... Don’t start with eight particles.

Shell-model project

Our recommended occupation representation, e.g. $|1,2,4,8\rangle$, is easy to code, but numerically inefficient when one has hundreds of millions of Slater determinants.

In state-of-the-art shell-model codes, one generally uses bit representation, i.e. $|1101000100\dots\rangle$ where one stores the Slater determinant as a single (or a small number of) integer.

This is much more compact, but more intricate to code with considerable more overhead. There exist bit-manipulation functions. We will discuss these in more detail at the beginning of the third week.

Example case: pairing Hamiltonian

We consider a space with 2Ω single-particle states, with each state labeled by $k = 1, 2, 3, \dots, \Omega$ and $m = \pm 1/2$. The convention is that the state with $k > 0$ has $m = +1/2$ while $-k$ has $m = -1/2$.

The Hamiltonian we consider is

$$\hat{H} = -G\hat{P}_+\hat{P}_-,$$

where

$$\hat{P}_+ = \sum_{k>0} \hat{a}_k^\dagger \hat{a}_{-k}^\dagger.$$

and $\hat{P}_- = (\hat{P}_+)^\dagger$.

This problem can be solved using what is called the quasi-spin formalism to obtain the exact results. Thereafter we will try again using the explicit Slater determinant formalism.

Example case: pairing Hamiltonian

One can show (and this is part of the project) that

$$[\hat{P}_+, \hat{P}_-] = \sum_{k>0} \left(\hat{a}_k^\dagger \hat{a}_k + \hat{a}_{-k}^\dagger \hat{a}_{-k} - 1 \right) = \hat{N} - \Omega.$$

Now define

$$\hat{P}_z = \frac{1}{2}(\hat{N} - \Omega).$$

Finally you can show

$$[\hat{P}_z, \hat{P}_\pm] = \pm \hat{P}_\pm.$$

This means the operators \hat{P}_\pm, \hat{P}_z form a so-called $SU(2)$ algebra, and we can use all our insights about angular momentum, even though there is no actual angular momentum involved.

So we rewrite the Hamiltonian to make this explicit:

$$\hat{H} = -G\hat{P}_+\hat{P}_- = -G(\hat{P}^2 - \hat{P}_z^2 + \hat{P}_z)$$

Example case: pairing Hamiltonian

Because of the $SU(2)$ algebra, we know that the eigenvalues of \hat{P}^2 must be of the form $p(p+1)$, with p either integer or half-integer, and the eigenvalues of \hat{P}_z are m_p with $p \geq |m_p|$, with m_p also integer or half-integer.

But because $\hat{P}_z = (1/2)(\hat{N} - \Omega)$, we know that for N particles the value $m_p = (N - \Omega)/2$. Furthermore, the values of m_p range from $-\Omega/2$ (for $N = 0$) to $+\Omega/2$ (for $N = 2\Omega$, with all states filled).

We deduce the maximal $p = \Omega/2$ and for a given n the values range of p range from $|N - \Omega|/2$ to $\Omega/2$ in steps of 1 (for an even number of particles)

Following Racah we introduce the notation $p = (\Omega - \nu)/2$ where $\nu = 0, 2, 4, \dots, \Omega - |N - \Omega|$. With this it is easy to deduce that the eigenvalues of the pairing Hamiltonian are

$$-G(N - \nu)(2\Omega + 2 - N - \nu)/4$$

This also works for N odd, with $\nu = 1, 3, 5, \dots$.

Example case: pairing Hamiltonian

Let's take a specific example: $\Omega = 3$ so there are 6 single-particle states, and $N = 3$, with $\nu = 1, 3$. Therefore there are two distinct eigenvalues,

$$E = -2G, 0$$

Now let's work this out explicitly. The single particle degrees of freedom are defined as

Index	k	m
1	1	-1/2
2	-1	1/2
3	2	-1/2
4	-2	1/2
5	3	-1/2
6	-3	1/2

There are $\binom{6}{3} = 20$ three-particle states, but there are 9 states with $M = +1/2$, namely $|1, 2, 3\rangle, |1, 2, 5\rangle, |1, 4, 6\rangle, |2, 3, 4\rangle, |2, 3, 6\rangle, |2, 4, 5\rangle, |2, 5, 6\rangle, |3, 4, 6\rangle, |4, 5, 6\rangle$.

Example case: pairing Hamiltonian

In this basis, the operator

$$\hat{P}_+ = \hat{a}_1^\dagger \hat{a}_2^\dagger + \hat{a}_3^\dagger \hat{a}_4^\dagger + \hat{a}_5^\dagger \hat{a}_6^\dagger$$

From this we can determine that

$$\hat{P}_- |1, 4, 6\rangle = \hat{P}_- |2, 3, 6\rangle = \hat{P}_- |2, 4, 5\rangle = 0$$

so those states all have eigenvalue 0.

Example case: pairing Hamiltonian

Now for further example,

$$\hat{P}_- |1, 2, 3\rangle = |3\rangle$$

so

$$\hat{P}_+ \hat{P}_- |1, 2, 3\rangle = |1, 2, 3\rangle + |3, 4, 3\rangle + |5, 6, 3\rangle$$

The second term vanishes because state 3 is occupied twice, and reordering the last term we get

$$\hat{P}_+ \hat{P}_- |1, 2, 3\rangle = |1, 2, 3\rangle + |3, 5, 6\rangle$$

without picking up a phase.

Example case: pairing Hamiltonian

Continuing in this fashion, with the previous ordering of the many-body states ($|1, 2, 3\rangle, |1, 2, 5\rangle, |1, 4, 6\rangle, |2, 3, 4\rangle, |2, 3, 6\rangle$), the Hamiltonian matrix of this system is

$$H = -G \begin{pmatrix} 1 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \\ 0 & 1 & 0 & 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 0 & 1 & 0 \\ 1 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \end{pmatrix}$$

This is useful for our project. One can by hand confirm that there are 3 eigenvalues $-2G$ and 6 with value zero.

Example case: pairing Hamiltonian

Another example Using the $(1/2)^4$ single-particle space, resulting in eight single-particle states

Index	n	l	s	m_s
1	0	0	1/2	-1/2
2	0	0	1/2	1/2
3	1	0	1/2	-1/2
4	1	0	1/2	1/2
5	2	0	1/2	-1/2
6	2	0	1/2	1/2
7	3	0	1/2	-1/2
8	3	0	1/2	1/2

and then taking only 4-particle, $M = 0$ states that have no 'broken pairs', there are six basis Slater determinants:

- $|1, 2, 3, 4\rangle$,
- $|1, 2, 5, 6\rangle$,

- $|1, 2, 7, 8\rangle$,
- $|3, 4, 5, 6\rangle$,
- $|3, 4, 7, 8\rangle$,
- $|5, 6, 7, 8\rangle$

Example case: pairing Hamiltonian

Now we take the following Hamiltonian

$$\hat{H} = \sum_n n\delta\hat{N}_n - G\hat{P}^\dagger\hat{P}$$

where

$$\hat{N}_n = \hat{a}_{n,m=+1/2}^\dagger \hat{a}_{n,m=+1/2} + \hat{a}_{n,m=-1/2}^\dagger \hat{a}_{n,m=-1/2}$$

and

$$\hat{P}^\dagger = \sum_n \hat{a}_{n,m=+1/2}^\dagger \hat{a}_{n,m=-1/2}^\dagger$$

We can write down the 6×6 Hamiltonian in the basis from the prior slide:

$$H = \begin{pmatrix} 2\delta - 2G & -G & -G & -G & -G & 0 \\ -G & 4\delta - 2G & -G & -G & -0 & -G \\ -G & -G & 6\delta - 2G & 0 & -G & -G \\ -G & -G & 0 & 6\delta - 2G & -G & -G \\ -G & 0 & -G & -G & 8\delta - 2G & -G \\ 0 & -G & -G & -G & -G & 10\delta - 2G \end{pmatrix}$$

(You should check by hand that this is correct.)

For $\delta = 0$ we have the closed form solution of the g.s. energy given by $-6G$.

Building a Hamiltonian matrix

The goal is to compute the matrix elements of the Hamiltonian, specifically matrix elements between many-body states (Slater determinants) of two-body operators

$$\sum_{p < q, r < s} V_{pqrs} \hat{a}_p^\dagger \hat{a}_q^\dagger \hat{a}_s \hat{a}_r$$

In particular we will need to compute

$$\langle \beta | \hat{a}_p^\dagger \hat{a}_q^\dagger \hat{a}_s \hat{a}_r | \alpha \rangle$$

where α, β are indices labeling Slater determinants and p, q, r, s label single-particle states.

Building a Hamiltonian matrix

Note: there are other, more efficient ways to do this than the method we describe, but you will be able to produce a working code quickly.

As we coded in the first step, a Slater determinant $|\alpha\rangle$ with index α is a list of N occupied single-particle states $i_1 < i_2 < i_3 \dots i_N$.

Furthermore, for the two-body matrix elements V_{pqrs} we normally assume $p < q$ and $r < s$. For our specific project, the interaction is much simpler and you can use this to simplify considerably the setup of a shell-model code for project 2.

What follows here is a more general, but still brute force, approach.

Building a Hamiltonian matrix

Write a function that:

1. Has as input the single-particle indices p, q, r, s for the two-body operator and the index α for the ket Slater determinant;
2. Returns the index β of the unique (if any) Slater determinant such that

$$|\beta\rangle = \pm \hat{a}_p^\dagger \hat{a}_q^\dagger \hat{a}_s \hat{a}_r |\alpha\rangle$$

as well as the phase

This is equivalent to computing

$$\langle \beta | \hat{a}_p^\dagger \hat{a}_q^\dagger \hat{a}_s \hat{a}_r | \alpha \rangle$$

Building a Hamiltonian matrix, first step

The first step can take as input an initial Slater determinant (whose position in the list of basis Slater determinants is α) written as an ordered list of occupied single-particle states, e.g. 1, 2, 5, 8, and the indices p, q, r, s from the two-body operator.

It will return another final Slater determinant if the single-particle states r and s are occupied, else it will return an empty Slater determinant (all zeroes).

If r and s are in the list of occupied single particle states, then replace the initial single-particle states ij as $i \rightarrow r$ and $j \rightarrow s$.

Building a Hamiltonian matrix, second step

The second step will take the final Slater determinant from the first step (if not empty), and then order by pairwise permutations (i.e., if the Slater determinant is i_1, i_2, i_3, \dots , then if $i_n > i_{n+1}$, interchange $i_n \leftrightarrow i_{n+1}$).

Building a Hamiltonian matrix

It will also output a phase. If any two single-particle occupancies are repeated, the phase is 0. Otherwise it is +1 for an even permutation and -1 for an odd permutation to bring the final Slater determinant into ascending order, $j_1 < j_2 < j_3 \dots$.

Building a Hamiltonian matrix

Example: Suppose in the sd single-particle space that the initial Slater determinant is 1, 3, 9, 12. If $p, q, r, s = 2, 8, 1, 12$, then after the first step the final Slater determinant is 2, 3, 9, 8. The second step will return 2, 3, 8, 9 and a phase of -1, because an odd number of interchanges is required.

Building a Hamiltonian matrix

Example: Suppose in the sd single-particle space that the initial Slater determinant is 1, 3, 9, 12. If $p, q, r, s = 3, 8, 1, 12$, then after the first step the final Slater determinant is 3, 3, 9, 8, but after the second step the phase is 0 because the single-particle state 3 is occupied twice.

Lastly, the final step takes the ordered final Slater determinant and we search through the basis list to determine its index in the many-body basis, that is, β .

Building a Hamiltonian matrix

The Hamiltonian is then stored as an $N_{SD} \times N_{SD}$ array of real numbers, which can be allocated once you have created the many-body basis and know N_{SD} .

Building a Hamiltonian matrix

1. Initialize $H(\alpha, \beta) = 0.0$
2. Set up an outer loop over β
3. Loop over $\alpha = 1, NSD$
4. For each α , loop over $a = 1, nbme$ and fetch $V(a)$ and the single-particle indices p, q, r, s
5. If $V(a) = 0$ skip. Otherwise, apply $\hat{a}_p^\dagger \hat{a}_q^\dagger \hat{a}_s \hat{a}_r$ to the Slater determinant labeled by α .
6. Find, if any, the label β of the resulting Slater determinant and the phase (which is 0, +1, -1).
7. If phase $\neq 0$, then update $H(\alpha, \beta)$ as $H(\alpha, \beta) + phase * V(a)$. The sum is important because multiple operators might contribute to the same matrix element.
8. Continue loop over a
9. Continue loop over α .
10. End the outer loop over β .

You should force the resulting matrix H to be symmetric. To do this, when updating $H(\alpha, \beta)$, if $\alpha \neq \beta$, also update $H(\beta, \alpha)$.

Building a Hamiltonian matrix

You will also need to include the single-particle energies. This is easy: they only contribute to diagonal matrix elements, that is, $H(\alpha, \alpha)$. Simply find the occupied single-particle states i and add the corresponding $\varepsilon(i)$.

Hamiltonian matrix without the bit representation

Consider the many-body state Ψ_λ expressed as linear combinations of Slater determinants (SD) of orthonormal single-particle states $\phi(\mathbf{r})$:

$$\Psi_\lambda = \sum_i C_{\lambda i} SD_i \quad (3.2)$$

Using the Slater-Condon rules the matrix elements of any one-body (\mathcal{O}_∞) or two-body (\mathcal{O}_∞) operator expressed in the determinant space have simple expressions involving one- and two-fermion integrals in our given single-particle basis. The diagonal elements are given by:

$$\begin{aligned} \langle SD | \mathcal{O}_\infty | \mathcal{SD} \rangle &= \sum_{i \in SD} \langle \phi_i | \mathcal{O}_\infty | \phi_i \rangle \\ \langle SD | \mathcal{O}_\infty | \mathcal{SD} \rangle &= \frac{1}{2} \sum_{(i,j) \in SD} \langle \phi_i \phi_j | \mathcal{O}_\infty | \phi_i \phi_j \rangle - \\ &\quad \langle \phi_i \phi_j | \mathcal{O}_\infty | \phi_j \phi_i \rangle \end{aligned} \quad (3.3)$$

Hamiltonian matrix without the bit representation, one and two-body operators

For two determinants which differ only by the substitution of single-particle states i with a single-particle state j :

$$\begin{aligned}\langle SD|\mathcal{O}_\infty|\mathcal{SD}\rangle &= \langle \phi_i|\mathcal{O}_\infty|\phi_j\rangle \\ \langle SD|\mathcal{O}_\infty|\mathcal{SD}\rangle &= \sum_{k \in SD} \langle \phi_i\phi_k|\mathcal{O}_\infty|\phi_j\phi_k\rangle - \langle \phi_j\phi_k|\mathcal{O}_\infty|\phi_i\phi_k\rangle\end{aligned}\tag{3.4}$$

For two determinants which differ by two single-particle states

$$\begin{aligned}\langle SD|\mathcal{O}_\infty|\mathcal{SD}\rangle &= 0 \\ \langle SD|\mathcal{O}_\infty|\mathcal{SD}\rangle &= \langle \phi_i\phi_k|\mathcal{O}_\infty|\phi_j\phi_l\rangle - \langle \phi_j\phi_l|\mathcal{O}_\infty|\phi_i\phi_k\rangle\end{aligned}\tag{3.5}$$

All other matrix elements involving determinants with more than two substitutions are zero.

Strategies for setting up an algorithm

An efficient implementation of these rules requires

- to find the number of single-particle state substitutions between two determinants
- to find which single-particle states are involved in the substitution
- to compute the phase factor if a reordering of the single-particle states has occurred

We can solve this problem using our odometric approach or alternatively using a bit representation as discussed below and in more detail in

- Scemama and Gimer's article (Fortran codes)
- Simen Kvaal's article on how to build an FCI code (C++ code)

We recommend in particular the article by Simen Kvaal. It contains nice general classes for creation and annihilation operators as well as the calculation of the phase (see below).

Computing expectation values and transitions in the shell-model

When we diagonalize the Hamiltonian matrix, the eigenvectors are the coefficients $C_{\lambda i}$ used to express the many-body state Ψ_λ in terms of a linear combinations of Slater determinants (SD) of orthonormal single-particle states $\phi(\mathbf{r})$.

With these eigenvectors we can compute say the transition likelihood of a one-body operator as

$$\langle \Psi_\lambda|\mathcal{O}_\infty|\Psi_\sigma\rangle = \sum_{\lambda'} C_{\lambda'}^* \mathcal{C}_\sigma \langle \mathcal{SD}|\mathcal{O}_\infty|\mathcal{SD}\rangle.$$

Writing the one-body operator in second quantization as

$$\mathcal{O}_\infty = \sum_{\sqrt{\lambda}} \langle \sqrt{\lambda} | \lambda_\infty | \blacksquare \rangle a_{\sqrt{\lambda}}^\dagger a_{\sqrt{\lambda}},$$

we have

$$\langle \Psi_\lambda | \mathcal{O}_\infty | \Theta_\sigma \rangle = \sum_{\sqrt{\lambda}} \langle \sqrt{\lambda} | \lambda_\infty | \blacksquare \rangle \sum_{\lambda} \mathcal{C}_\lambda^* \mathcal{C}_\sigma \langle \mathcal{SD} | a_{\sqrt{\lambda}}^\dagger a_{\sqrt{\lambda}} | \mathcal{SD} \rangle.$$

Computing expectation values and transitions in the shell-model and spectroscopic factors

The terms we need to evaluate then are just the elements

$$\langle SD_i | a_p^\dagger a_q | SD_j \rangle,$$

which can be rewritten in terms of spectroscopic factors by inserting a complete set of Slater determinants as

$$\langle SD_i | a_p^\dagger a_q | SD_j \rangle = \sum_l \langle SD_i | a_p^\dagger | SD_l \rangle \langle SD_l | a_q | SD_j \rangle,$$

where $\langle SD_l | a_q | SD_j \rangle$ are the spectroscopic factors. These can be easily evaluated in m -scheme. Using the Wigner-Eckart theorem we can transform these to a J -coupled scheme through so-called reduced matrix elements.

Operators in second quantization

In the build-up of a shell-model or FCI code that is meant to tackle large dimensionalities we need to deal with the action of the Hamiltonian \hat{H} on a Slater determinant represented in second quantization as

$$|\alpha_1 \dots \alpha_n\rangle = a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle.$$

The time consuming part stems from the action of the Hamiltonian on the above determinant,

$$\left(\sum_{\alpha\beta} \langle \alpha | t + u | \beta \rangle a_\alpha^\dagger a_\beta + \frac{1}{4} \sum_{\alpha\beta\gamma\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle a_\alpha^\dagger a_\beta^\dagger a_\delta a_\gamma \right) a_{\alpha_1}^\dagger a_{\alpha_2}^\dagger \dots a_{\alpha_n}^\dagger |0\rangle.$$

A practically useful way to implement this action is to encode a Slater determinant as a bit pattern.

Operators in second quantization

Assume that we have at our disposal n different single-particle states $\alpha_0, \alpha_2, \dots, \alpha_{n-1}$ and that we can distribute among these states $N \leq n$ particles.

A Slater determinant can then be coded as an integer of n bits. As an example, if we have $n = 16$ single-particle states $\alpha_0, \alpha_1, \dots, \alpha_{15}$ and $N = 4$ fermions occupying the states $\alpha_3, \alpha_6, \alpha_{10}$ and α_{13} we could write this Slater determinant as

$$\Phi_\Lambda = a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle.$$

The unoccupied single-particle states have bit value 0 while the occupied ones are represented by bit state 1. In the binary notation we would write this 16 bits long integer as

$$\begin{array}{cccccccccccccccc} \alpha_0 & \alpha_1 & \alpha_2 & \alpha_3 & \alpha_4 & \alpha_5 & \alpha_6 & \alpha_7 & \alpha_8 & \alpha_9 & \alpha_{10} & \alpha_{11} & \alpha_{12} & \alpha_{13} & \alpha_{14} & \alpha_{15} \\ 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 1 & 0 & 0 \end{array}$$

which translates into the decimal number

$$2^3 + 2^6 + 2^{10} + 2^{13} = 9288.$$

We can thus encode a Slater determinant as a bit pattern.

Operators in second quantization

With N particles that can be distributed over n single-particle states, the total number of Slater determinats (and defining thereby the dimensionality of the system) is

$$\dim(\mathcal{H}) = \binom{n}{N}.$$

The total number of bit patterns is 2^n .

Operators in second quantization

We assume again that we have at our disposal n different single-particle orbits $\alpha_0, \alpha_2, \dots, \alpha_{n-1}$ and that we can distribute among these orbits $N \leq n$ particles. The ordering among these states is important as it defines the order of the creation operators. We will write the determinant

$$\Phi_\Lambda = a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

in a more compact way as

$$\Phi_{3,6,10,13} = |0001001000100100\rangle.$$

The action of a creation operator is thus

$$a_{\alpha_4}^\dagger \Phi_{3,6,10,13} = a_{\alpha_4}^\dagger |0001001000100100\rangle = a_{\alpha_4}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

which becomes

$$-a_{\alpha_3}^\dagger a_{\alpha_4}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle = -|0001101000100100\rangle.$$

Operators in second quantization

Similarly

$$a_{\alpha_6}^\dagger \Phi_{3,6,10,13} = a_{\alpha_6}^\dagger |0001001000100100\rangle = a_{\alpha_6}^\dagger a_{\alpha_3}^\dagger a_{\alpha_6}^\dagger a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle,$$

which becomes

$$-a_{\alpha_4}^\dagger (a_{\alpha_6}^\dagger)^2 a_{\alpha_{10}}^\dagger a_{\alpha_{13}}^\dagger |0\rangle = 0!$$

This gives a simple recipe:

- If one of the bits b_j is 1 and we act with a creation operator on this bit, we return a null vector
- If $b_j = 0$, we set it to 1 and return a sign factor $(-1)^l$, where l is the number of bits set before bit j .

Operators in second quantization

Consider the action of $a_{\alpha_2}^\dagger$ on various slater determinants:

$$\begin{aligned} a_{\alpha_2}^\dagger \Phi_{00111} &= a_{\alpha_2}^\dagger |00111\rangle = 0 \times |00111\rangle \\ a_{\alpha_2}^\dagger \Phi_{01011} &= a_{\alpha_2}^\dagger |01011\rangle = (-1) \times |01111\rangle \\ a_{\alpha_2}^\dagger \Phi_{01101} &= a_{\alpha_2}^\dagger |01101\rangle = 0 \times |01101\rangle \\ a_{\alpha_2}^\dagger \Phi_{01110} &= a_{\alpha_2}^\dagger |01110\rangle = 0 \times |01110\rangle \\ a_{\alpha_2}^\dagger \Phi_{10011} &= a_{\alpha_2}^\dagger |10011\rangle = (-1) \times |10111\rangle \\ a_{\alpha_2}^\dagger \Phi_{10101} &= a_{\alpha_2}^\dagger |10101\rangle = 0 \times |10101\rangle \\ a_{\alpha_2}^\dagger \Phi_{10110} &= a_{\alpha_2}^\dagger |10110\rangle = 0 \times |10110\rangle \\ a_{\alpha_2}^\dagger \Phi_{11001} &= a_{\alpha_2}^\dagger |11001\rangle = (+1) \times |11101\rangle \\ a_{\alpha_2}^\dagger \Phi_{11010} &= a_{\alpha_2}^\dagger |11010\rangle = (+1) \times |11110\rangle \end{aligned}$$

What is the simplest way to obtain the phase when we act with one annihilation(creation) operator on the given Slater determinant representation?

Operators in second quantization

We have an SD representation

$$\Phi_{\Lambda} = a_{\alpha_0}^{\dagger} a_{\alpha_3}^{\dagger} a_{\alpha_6}^{\dagger} a_{\alpha_{10}}^{\dagger} a_{\alpha_{13}}^{\dagger} |0\rangle,$$

in a more compact way as

$$\Phi_{0,3,6,10,13} = |1001001000100100\rangle.$$

The action of

$$a_{\alpha_4}^{\dagger} a_{\alpha_0} \Phi_{0,3,6,10,13} = a_{\alpha_4}^{\dagger} |0001001000100100\rangle = a_{\alpha_4}^{\dagger} a_{\alpha_3}^{\dagger} a_{\alpha_6}^{\dagger} a_{\alpha_{10}}^{\dagger} a_{\alpha_{13}}^{\dagger} |0\rangle,$$

which becomes

$$-a_{\alpha_3}^{\dagger} a_{\alpha_4}^{\dagger} a_{\alpha_6}^{\dagger} a_{\alpha_{10}}^{\dagger} a_{\alpha_{13}}^{\dagger} |0\rangle = -|0001101000100100\rangle.$$

Operators in second quantization

The action

$$a_{\alpha_0} \Phi_{0,3,6,10,13} = |0001001000100100\rangle,$$

can be obtained by subtracting the logical sum (AND operation) of $\Phi_{0,3,6,10,13}$ and a word which represents only α_0 , that is

$$|1000000000000000\rangle,$$

from $\Phi_{0,3,6,10,13} = |1001001000100100\rangle$.

This operation gives $|0001001000100100\rangle$.

Similarly, we can form $a_{\alpha_4}^{\dagger} a_{\alpha_0} \Phi_{0,3,6,10,13}$, say, by adding $|0000100000000000\rangle$ to $a_{\alpha_0} \Phi_{0,3,6,10,13}$, first checking that their logical sum is zero in order to make sure that the state α_4 is not already occupied.

Operators in second quantization

It is trickier however to get the phase $(-1)^l$. One possibility is as follows

- Let S_1 be a word that represents the 1-bit to be removed and all others set to zero.

In the previous example $S_1 = |1000000000000000\rangle$

- Define S_2 as the similar word that represents the bit to be added, that is in our case

$$S_2 = |0000100000000000\rangle.$$

- Compute then $S = S_1 - S_2$, which here becomes

$$S = |0111000000000000\rangle$$

- Perform then the logical AND operation of S with the word containing

$$\Phi_{0,3,6,10,13} = |1001001000100100\rangle,$$

which results in $|0001000000000000\rangle$. Counting the number of 1-bits gives the phase. Here you need however an algorithm for bitcounting.

Bit counting

We include here a python program which may aid in this direction. It uses bit manipulation functions from <http://wiki.python.org/moin/BitManipulation>.

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
import math

""" A simple Python class for Slater determinant manipulation Bit-manipulation stolen
from:
http://wiki.python.org/moin/BitManipulation """
bitCount() counts the number of bits set (not an optimal function)
def bitCount(int_type) : """Countbitssetininteger""" count = 0 while(int_type) : int_type = int_type - 1 count += 1
return(count)
testBit() returns a nonzero result, 2**offset, if the bit at 'offset' is one.
def testBit(int_type,offset) : mask = 1 << offset return(int_type & mask) >> offset
setBit() returns an integer with the bit at 'offset' set to 1.
def setBit(int_type,offset) : mask = 1 << offset return(int_type | mask)
clearBit() returns an integer with the bit at 'offset' cleared.
def clearBit(int_type,offset) : mask = (1 << offset) return(int_type & ~mask)
toggleBit() returns an integer with the bit at 'offset' inverted, 0 -> 1 and 1 -> 0.
def toggleBit(int_type,offset) : mask = 1 << offset return(int_type ^ mask)
binary string made from number
def bin0(s): return str(s) if s<=1 else bin0(s>>1) + str(s1)
def bin(s, L = 0): ss = bin0(s) if L > 0: return '0'*(L-len(ss)) + ss else: return ss
class Slater: """ Class for Slater determinants """ def __init__(self): self.word = int(0)
def create(self, j): print "c:" + str(j) + "|" + bin(self.word) + " >=", Assume bit j is set, then we return zero. s =
0 Check if bit j is set. isset = testBit(self.word, j) if isset == 0: bits = bitCount(self.word((1 << j) - 1)) s = pow(-1, bits) self.word =
setBit(self.word, j)
print str(s) + " x |" + bin(self.word) + ">" return s
def annihilate(self, j): print "c:" + str(j) + "|" + bin(self.word) + " >=", Assume bit j is not set, then we return zero. s =
0 Check if bit j is set. isset = testBit(self.word, j) if isset == 1: bits = bitCount(self.word((1 << j) - 1)) s = pow(-1, bits) self.word =
clearBit(self.word, j)
print str(s) + " x |" + bin(self.word) + ">" return s
Do some testing:
phi = Slater() phi.create(0) phi.create(1) phi.create(2) phi.create(3)
print
s = phi.annihilate(2) s = phi.create(7) s = phi.annihilate(0) s = phi.create(200)
```

Eigenvalue problems, basic definitions

Let us consider the matrix \mathbf{A} of dimension n . The eigenvalues of \mathbf{A} are defined through the matrix equation

$$\mathbf{A}\mathbf{x}^{(v)} = \lambda^{(v)}\mathbf{x}^{(v)},$$

where $\lambda^{(v)}$ are the eigenvalues and $\mathbf{x}^{(v)}$ the corresponding eigenvectors. Unless otherwise stated, when we use the wording eigenvector we mean the right eigenvector. The left eigen-

value problem is defined as

$$\mathbf{x}_L^{(v)} \mathbf{A} = \lambda^{(v)} \mathbf{x}_L^{(v)}$$

The above right eigenvector problem is equivalent to a set of n equations with n unknowns x_i .

Eigenvalue problems, basic definitions

The eigenvalue problem can be rewritten as

$$(\mathbf{A} - \lambda^{(v)} \mathbf{I}) \mathbf{x}^{(v)} = 0,$$

with \mathbf{I} being the unity matrix. This equation provides a solution to the problem if and only if the determinant is zero, namely

$$|\mathbf{A} - \lambda^{(v)} \mathbf{I}| = 0,$$

which in turn means that the determinant is a polynomial of degree n in λ and in general we will have n distinct zeros.

Eigenvalue problems, basic definitions

The eigenvalues of a matrix $\mathbf{A} \in \mathbb{C}^{n \times n}$ are thus the n roots of its characteristic polynomial

$$P(\lambda) = \det(\lambda \mathbf{I} - \mathbf{A}),$$

or

$$P(\lambda) = \prod_{i=1}^n (\lambda_i - \lambda).$$

The set of these roots is called the spectrum and is denoted as $\lambda(\mathbf{A})$. If $\lambda(\mathbf{A}) = \{\lambda_1, \lambda_2, \dots, \lambda_n\}$ then we have

$$\det(\mathbf{A}) = \lambda_1 \lambda_2 \dots \lambda_n,$$

and if we define the trace of \mathbf{A} as

$$\text{Tr}(\mathbf{A}) = \sum_{i=1}^n a_{ii}$$

then

$$\text{Tr}(\mathbf{A}) = \lambda_1 + \lambda_2 + \dots + \lambda_n.$$

Abel-Ruffini Impossibility Theorem

The *Abel-Ruffini* theorem (also known as Abel's impossibility theorem) states that there is no general solution in radicals to polynomial equations of degree five or higher.

The content of this theorem is frequently misunderstood. It does not assert that higher-degree polynomial equations are unsolvable. In fact, if the polynomial has real or complex coefficients, and we allow complex solutions, then every polynomial equation has solutions; this is the fundamental theorem of algebra. Although these solutions cannot always be computed exactly with radicals, they can be computed to any desired degree of accuracy using numerical methods such as the Newton-Raphson method or Laguerre method, and in this way they are no different from solutions to polynomial equations of the second, third, or fourth degrees.

The theorem only concerns the form that such a solution must take. The content of the theorem is that the solution of a higher-degree equation cannot in all cases be expressed in terms of the polynomial coefficients with a finite number of operations of addition, subtraction, multiplication, division and root extraction. Some polynomials of arbitrary degree, of which the simplest nontrivial example is the monomial equation $ax^n = b$, are always solvable with a radical.

Abel-Ruffini Impossibility Theorem

The *Abel-Ruffini* theorem says that there are some fifth-degree equations whose solution cannot be so expressed. The equation $x^5 - x + 1 = 0$ is an example. Some other fifth degree equations can be solved by radicals, for example $x^5 - x^4 - x + 1 = 0$. The precise criterion that distinguishes between those equations that can be solved by radicals and those that cannot was given by Galois and is now part of Galois theory: a polynomial equation can be solved by radicals if and only if its Galois group is a solvable group.

Today, in the modern algebraic context, we say that second, third and fourth degree polynomial equations can always be solved by radicals because the symmetric groups S_2, S_3 and S_4 are solvable groups, whereas S_n is not solvable for $n \geq 5$.

Eigenvalue problems, basic definitions

In the present discussion we assume that our matrix is real and symmetric, that is $\mathbf{A} \in \mathbb{R}^{n \times n}$. The matrix \mathbf{A} has n eigenvalues $\lambda_1 \dots \lambda_n$ (distinct or not). Let \mathbf{D} be the diagonal matrix with the eigenvalues on the diagonal

$$\mathbf{D} = \begin{pmatrix} \lambda_1 & 0 & 0 & 0 & \dots & 0 & 0 \\ 0 & \lambda_2 & 0 & 0 & \dots & 0 & 0 \\ 0 & 0 & \lambda_3 & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots & \dots & \dots & \dots \\ 0 & \dots & \dots & \dots & \dots & \lambda_{n-1} & 0 \\ 0 & \dots & \dots & \dots & \dots & 0 & \lambda_n \end{pmatrix}.$$

If \mathbf{A} is real and symmetric then there exists a real orthogonal matrix \mathbf{S} such that

$$\mathbf{S}^T \mathbf{A} \mathbf{S} = \text{diag}(\lambda_1, \lambda_2, \dots, \lambda_n),$$

and for $j = 1 : n$ we have $\mathbf{A}\mathbf{S}(:, j) = \lambda_j \mathbf{S}(:, j)$.

Eigenvalue problems, basic definitions

To obtain the eigenvalues of $\mathbf{A} \in \mathbb{R}^{n \times n}$, the strategy is to perform a series of similarity transformations on the original matrix \mathbf{A} , in order to reduce it either into a diagonal form as above or into a tridiagonal form.

We say that a matrix \mathbf{B} is a similarity transform of \mathbf{A} if

$$\mathbf{B} = \mathbf{S}^T \mathbf{A} \mathbf{S}, \quad \text{where} \quad \mathbf{S}^T \mathbf{S} = \mathbf{S}^{-1} \mathbf{S} = \mathbf{I}.$$

The importance of a similarity transformation lies in the fact that the resulting matrix has the same eigenvalues, but the eigenvectors are in general different.

Eigenvalue problems, basic definitions

To prove this we start with the eigenvalue problem and a similarity transformed matrix \mathbf{B} .

$$\mathbf{A}\mathbf{x} = \lambda \mathbf{x} \quad \text{and} \quad \mathbf{B} = \mathbf{S}^T \mathbf{A} \mathbf{S}.$$

We multiply the first equation on the left by \mathbf{S}^T and insert $\mathbf{S}^T \mathbf{S} = \mathbf{I}$ between \mathbf{A} and \mathbf{x} . Then we get

$$(\mathbf{S}^T \mathbf{A} \mathbf{S})(\mathbf{S}^T \mathbf{x}) = \lambda \mathbf{S}^T \mathbf{x}, \quad (3.6)$$

which is the same as

$$\mathbf{B}(\mathbf{S}^T \mathbf{x}) = \lambda (\mathbf{S}^T \mathbf{x}).$$

The variable λ is an eigenvalue of \mathbf{B} as well, but with eigenvector $\mathbf{S}^T \mathbf{x}$.

Eigenvalue problems, basic definitions

The basic philosophy is to

- Either apply subsequent similarity transformations (direct method) so that

$$\mathbf{S}_N^T \dots \mathbf{S}_1^T \mathbf{A} \mathbf{S}_1 \dots \mathbf{S}_N = \mathbf{D}, \quad (3.7)$$

- Or apply subsequent similarity transformations so that \mathbf{A} becomes tridiagonal (Householder) or upper/lower triangular (the QR method to be discussed later).
- Thereafter, techniques for obtaining eigenvalues from tridiagonal matrices can be used.
- Or use so-called power methods
- Or use iterative methods (Krylov, Lanczos, Arnoldi). These methods are popular for huge matrix problems.

Discussion of methods for eigenvalues

The general overview.

One speaks normally of two main approaches to solving the eigenvalue problem.

- The first is the formal method, involving determinants and the characteristic polynomial. This proves how many eigenvalues there are, and is the way most of you learned about how to solve the eigenvalue problem, but for matrices of dimensions greater than 2 or 3, it is rather impractical.
- The other general approach is to use similarity or unitary transformations to reduce a matrix to diagonal form. This is normally done in two steps: first reduce to for example a *tridiagonal* form, and then to diagonal form. The main algorithms we will discuss in detail, Jacobi's and Householder's (so-called direct method) and Lanczos algorithms (an iterative method), follow this methodology.

Eigenvalues methods

Direct or non-iterative methods require for matrices of dimensionality $n \times n$ typically $O(n^3)$ operations. These methods are normally called standard methods and are used for dimensionalities $n \sim 10^5$ or smaller. A brief historical overview

Year	n	
1950	$n = 20$	(Wilkinson)
1965	$n = 200$	(Forsythe et al.)
1980	$n = 2000$	Linpac
1995	$n = 20000$	Lapack
This decade	$n \sim 10^5$	Lapack

shows that in the course of 60 years the dimension that direct diagonalization methods can handle has increased by almost a factor of 10^4 (note this is for serial versions). However, it pales beside the progress achieved by computer hardware, from flops to petaflops, a factor of almost 10^{15} . We see clearly played out in history the $O(n^3)$ bottleneck of direct matrix algorithms.

Sloppily speaking, when $n \sim 10^4$ is cubed we have $O(10^{12})$ operations, which is smaller than the 10^{15} increase in flops.

Discussion of methods for eigenvalues

If the matrix to diagonalize is large and sparse, direct methods simply become impractical, also because many of the direct methods tend to destroy sparsity. As a result large dense matrices may arise during the diagonalization procedure. The idea behind iterative methods is to project the n -dimensional problem in smaller spaces, so-called Krylov subspaces. Given a matrix \mathbf{A} and a vector \mathbf{v} , the associated Krylov sequences of vectors (and thereby subspaces) $\mathbf{v}, \mathbf{A}\mathbf{v}, \mathbf{A}^2\mathbf{v}, \mathbf{A}^3\mathbf{v}, \dots$, represent successively larger Krylov subspaces.

Matrix	$\mathbf{Ax} = \mathbf{b}$	$\mathbf{Ax} = \lambda \mathbf{x}$
$\mathbf{A} = \mathbf{A}^*$	Conjugate gradient	Lanczos
$\mathbf{A} \neq \mathbf{A}^*$	GMRES etc	Arnoldi

Eigenvalues and Lanczos' method

Basic features with a real symmetric matrix (and normally huge $n > 10^6$ and sparse) \hat{A} of dimension $n \times n$:

- Lanczos' algorithm generates a sequence of real tridiagonal matrices T_k of dimension $k \times k$ with $k \leq n$, with the property that the extremal eigenvalues of T_k are progressively better estimates of \hat{A} ' extremal eigenvalues.* The method converges to the extremal eigenvalues.
- The similarity transformation is

$$\hat{T} = \hat{Q}^T \hat{A} \hat{Q},$$

with the first vector $\hat{Q}\hat{e}_1 = \hat{q}_1$.

We are going to solve iteratively

$$\hat{T} = \hat{Q}^T \hat{A} \hat{Q},$$

with the first vector $\hat{Q}\hat{e}_1 = \hat{q}_1$. We can write out the matrix \hat{Q} in terms of its column vectors

$$\hat{Q} = [\hat{q}_1 \hat{q}_2 \dots \hat{q}_n].$$

Eigenvalues and Lanczos' method, tridiagonal matrix

The matrix

$$\hat{T} = \hat{Q}^T \hat{A} \hat{Q},$$

can be written as

$$\hat{T} = \begin{pmatrix} \alpha_1 & \beta_1 & 0 & \dots & \dots & 0 \\ \beta_1 & \alpha_2 & \beta_2 & 0 & \dots & 0 \\ 0 & \beta_2 & \alpha_3 & \beta_3 & \dots & 0 \\ \dots & \dots & \dots & \dots & \dots & 0 \\ \dots & & \beta_{n-2} & \alpha_{n-1} & \beta_{n-1} & \\ 0 & \dots & \dots & 0 & \beta_{n-1} & \alpha_n \end{pmatrix}$$

Eigenvalues and Lanczos' method, tridiagonal and orthogonal matrices

Using the fact that

$$\hat{Q}\hat{Q}^T = \hat{I},$$

we can rewrite

$$\hat{T} = \hat{Q}^T \hat{A} \hat{Q},$$

as

$$\hat{Q} \hat{T} = \hat{A} \hat{Q}.$$

Eigenvalues and Lanczos' method

If we equate columns

$$\hat{T} = \begin{pmatrix} \alpha_1 & \beta_1 & 0 & \dots & \dots & 0 \\ \beta_1 & \alpha_2 & \beta_2 & 0 & \dots & 0 \\ 0 & \beta_2 & \alpha_3 & \beta_3 & \dots & 0 \\ \dots & \dots & \dots & \dots & \dots & 0 \\ \dots & & & \beta_{n-2} & \alpha_{n-1} & \beta_{n-1} \\ 0 & \dots & \dots & 0 & \beta_{n-1} & \alpha_n \end{pmatrix}$$

we obtain

$$\hat{A} \hat{q}_k = \beta_{k-1} \hat{q}_{k-1} + \alpha_k \hat{q}_k + \beta_k \hat{q}_{k+1}.$$

Eigenvalues and Lanczos' method, defining the Lanczos' vectors

We have thus

$$\hat{A} \hat{q}_k = \beta_{k-1} \hat{q}_{k-1} + \alpha_k \hat{q}_k + \beta_k \hat{q}_{k+1},$$

with $\beta_0 \hat{q}_0 = 0$ for $k = 1 : n - 1$. Remember that the vectors \hat{q}_k are orthonormal and this implies

$$\alpha_k = \hat{q}_k^T \hat{A} \hat{q}_k,$$

and these vectors are called Lanczos vectors.

Eigenvalues and Lanczos' method, basic steps

We have thus

$$\hat{A} \hat{q}_k = \beta_{k-1} \hat{q}_{k-1} + \alpha_k \hat{q}_k + \beta_k \hat{q}_{k+1},$$

with $\beta_0 \hat{q}_0 = 0$ for $k = 1 : n - 1$ and

$$\alpha_k = \hat{q}_k^T \hat{A} \hat{q}_k.$$

If

$$\hat{r}_k = (\hat{A} - \alpha_k \hat{I}) \hat{q}_k - \beta_{k-1} \hat{q}_{k-1},$$

is non-zero, then

$$\hat{q}_{k+1} = \hat{r}_k / \beta_k,$$

with $\beta_k = \pm ||\hat{r}_k||_2$.

Part II
Mean-field theories and post Hartree-Fock
methods

Chapter 4

Hartree-Fock methods

Why Hartree-Fock? Derivation of Hartree-Fock equations in coordinate space

Hartree-Fock (HF) theory is an algorithm for finding an approximative expression for the ground state of a given Hamiltonian. The basic ingredients are

- Define a single-particle basis $\{\psi_\alpha\}$ so that

$$\hat{h}^{\text{HF}} \psi_\alpha = \varepsilon_\alpha \psi_\alpha$$

with the Hartree-Fock Hamiltonian defined as

$$\hat{h}^{\text{HF}} = \hat{t} + \hat{u}_{\text{ext}} + \hat{u}^{\text{HF}}$$

- The term \hat{u}^{HF} is a single-particle potential to be determined by the HF algorithm.
- The HF algorithm means to choose \hat{u}^{HF} in order to have

$$\langle \hat{H} \rangle = E^{\text{HF}} = \langle \Phi_0 | \hat{H} | \Phi_0 \rangle$$

that is to find a local minimum with a Slater determinant Φ_0 being the ansatz for the ground state.

- The variational principle ensures that $E^{\text{HF}} \geq E_0$, with E_0 the exact ground state energy.

We will show that the Hartree-Fock Hamiltonian \hat{h}^{HF} equals our definition of the operator \hat{f} discussed in connection with the new definition of the normal-ordered Hamiltonian (see later lectures), that is we have, for a specific matrix element

$$\langle p | \hat{h}^{\text{HF}} | q \rangle = \langle p | \hat{f} | q \rangle = \langle p | \hat{t} + \hat{u}_{\text{ext}} | q \rangle + \sum_{i \leq F} \langle pi | \hat{V} | qi \rangle_{AS},$$

meaning that

$$\langle p | \hat{u}^{\text{HF}} | q \rangle = \sum_{i \leq F} \langle pi | \hat{V} | qi \rangle_{AS}.$$

The so-called Hartree-Fock potential \hat{u}^{HF} brings an explicit medium dependence due to the summation over all single-particle states below the Fermi level F . It brings also in an explicit dependence on the two-body interaction (in nuclear physics we can also have complicated three- or higher-body forces). The two-body interaction, with its contribution from the other bystanding fermions, creates an effective mean field in which a given fermion moves, in addition to the external potential \hat{u}_{ext} which confines the motion of the fermion. For systems like nuclei, there is no external confining potential. Nuclei are examples of self-bound systems,

where the binding arises due to the intrinsic nature of the strong force. For nuclear systems thus, there would be no external one-body potential in the Hartree-Fock Hamiltonian.

Variational Calculus and Lagrangian Multipliers

The calculus of variations involves problems where the quantity to be minimized or maximized is an integral.

In the general case we have an integral of the type

$$E[\Phi] = \int_a^b f(\Phi(x), \frac{\partial \Phi}{\partial x}, x) dx,$$

where E is the quantity which is sought minimized or maximized. The problem is that although f is a function of the variables Φ , $\partial \Phi / \partial x$ and x , the exact dependence of Φ on x is not known. This means again that even though the integral has fixed limits a and b , the path of integration is not known. In our case the unknown quantities are the single-particle wave functions and we wish to choose an integration path which makes the functional $E[\Phi]$ stationary. This means that we want to find minima, or maxima or saddle points. In physics we search normally for minima. Our task is therefore to find the minimum of $E[\Phi]$ so that its variation δE is zero subject to specific constraints. In our case the constraints appear as the integral which expresses the orthogonality of the single-particle wave functions. The constraints can be treated via the technique of Lagrangian multipliers

Let us specialize to the expectation value of the energy for one particle in three-dimensions. This expectation value reads

$$E = \int dx dy dz \psi^*(x, y, z) \hat{H} \psi(x, y, z),$$

with the constraint

$$\int dx dy dz \psi^*(x, y, z) \psi(x, y, z) = 1,$$

and a Hamiltonian

$$\hat{H} = -\frac{1}{2} \nabla^2 + V(x, y, z).$$

We will, for the sake of notational convenience, skip the variables x, y, z below, and write for example $V(x, y, z) = V$.

The integral involving the kinetic energy can be written as, with the function ψ vanishing strongly for large values of x, y, z (given here by the limits a and b),

$$\int_a^b dx dy dz \psi^* \left(-\frac{1}{2} \nabla^2 \right) \psi dx dy dz = \psi^* \nabla \psi \Big|_a^b + \int_a^b dx dy dz \frac{1}{2} \nabla \psi^* \nabla \psi.$$

We will drop the limits a and b in the remaining discussion. Inserting this expression into the expectation value for the energy and taking the variational minimum we obtain

$$\delta E = \delta \left\{ \int dx dy dz \left(\frac{1}{2} \nabla \psi^* \nabla \psi + V \psi^* \psi \right) \right\} = 0.$$

The constraint appears in integral form as

$$\int dx dy dz \psi^* \psi = \text{constant},$$

and multiplying with a Lagrangian multiplier λ and taking the variational minimum we obtain the final variational equation

$$\delta \left\{ \int dx dy dz \left(\frac{1}{2} \nabla \psi^* \nabla \psi + V \psi^* \psi - \lambda \psi^* \psi \right) \right\} = 0.$$

We introduce the function f

$$f = \frac{1}{2} \nabla \psi^* \nabla \psi + V \psi^* \psi - \lambda \psi^* \psi = \frac{1}{2} (\psi_x^* \psi_x + \psi_y^* \psi_y + \psi_z^* \psi_z) + V \psi^* \psi - \lambda \psi^* \psi,$$

where we have skipped the dependence on x, y, z and introduced the shorthand ψ_x , ψ_y and ψ_z for the various derivatives.

For ψ^* the Euler-Lagrange equations yield

$$\frac{\partial f}{\partial \psi^*} - \frac{\partial}{\partial x} \frac{\partial f}{\partial \psi_x^*} - \frac{\partial}{\partial y} \frac{\partial f}{\partial \psi_y^*} - \frac{\partial}{\partial z} \frac{\partial f}{\partial \psi_z^*} = 0,$$

which results in

$$-\frac{1}{2} (\psi_{xx} + \psi_{yy} + \psi_{zz}) + V \psi = \lambda \psi.$$

We can then identify the Lagrangian multiplier as the energy of the system. The last equation is nothing but the standard Schroedinger equation and the variational approach discussed here provides a powerful method for obtaining approximate solutions of the wave function.

Derivation of Hartree-Fock equations in coordinate space

Let us denote the ground state energy by E_0 . According to the variational principle we have

$$E_0 \leq E[\Phi] = \int \Phi^* \hat{H} \Phi d\tau$$

where Φ is a trial function which we assume to be normalized

$$\int \Phi^* \Phi d\tau = 1,$$

where we have used the shorthand $d\tau = dx_1 dx_2 \dots dx_A$.

In the Hartree-Fock method the trial function is a Slater determinant which can be rewritten as

$$\Psi(x_1, x_2, \dots, x_A, \alpha, \beta, \dots, \nu) = \frac{1}{\sqrt{A!}} \sum_P (-)^P P \psi_\alpha(x_1) \psi_\beta(x_2) \dots \psi_\nu(x_A) = \sqrt{A!} \hat{A} \Phi_H,$$

where we have introduced the anti-symmetrization operator \hat{A} defined by the summation over all possible permutations p of two fermions. It is defined as

$$\hat{A} = \frac{1}{A!} \sum_p (-)^p \hat{P},$$

with the the Hartree-function given by the simple product of all possible single-particle function

$$\Phi_H(x_1, x_2, \dots, x_A, \alpha, \beta, \dots, \nu) = \psi_\alpha(x_1) \psi_\beta(x_2) \dots \psi_\nu(x_A).$$

Our functional is written as

$$E[\Phi] = \sum_{\mu=1}^A \int \psi_{\mu}^*(x_i) \hat{h}_0(x_i) \psi_{\mu}(x_i) dx_i + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \left[\int \psi_{\mu}^*(x_i) \psi_{\nu}^*(x_j) \hat{v}(r_{ij}) \psi_{\mu}(x_i) \psi_{\nu}(x_j) dx_i dx_j - \int \psi_{\mu}^*(x_i) \psi_{\nu}^*(x_j) \hat{v}(r_{ij}) \psi_{\nu}(x_i) \psi_{\mu}(x_j) dx_i dx_j \right]$$

The more compact version reads

$$E[\Phi] = \sum_{\mu} \langle \mu | \hat{h}_0 | \mu \rangle + \frac{1}{2} \sum_{\mu\nu} [\langle \mu\nu | \hat{v} | \mu\nu \rangle - \langle \nu\mu | \hat{v} | \mu\nu \rangle].$$

Since the interaction is invariant under the interchange of two particles it means for example that we have

$$\langle \mu\nu | \hat{v} | \mu\nu \rangle = \langle \nu\mu | \hat{v} | \nu\mu \rangle,$$

or in the more general case

$$\langle \mu\nu | \hat{v} | \sigma\tau \rangle = \langle \nu\mu | \hat{v} | \tau\sigma \rangle.$$

The direct and exchange matrix elements can be brought together if we define the anti-symmetrized matrix element

$$\langle \mu\nu | \hat{v} | \mu\nu \rangle_{AS} = \langle \mu\nu | \hat{v} | \mu\nu \rangle - \langle \mu\nu | \hat{v} | \nu\mu \rangle,$$

or for a general matrix element

$$\langle \mu\nu | \hat{v} | \sigma\tau \rangle_{AS} = \langle \mu\nu | \hat{v} | \sigma\tau \rangle - \langle \mu\nu | \hat{v} | \tau\sigma \rangle.$$

It has the symmetry property

$$\langle \mu\nu | \hat{v} | \sigma\tau \rangle_{AS} = -\langle \mu\nu | \hat{v} | \tau\sigma \rangle_{AS} = -\langle \nu\mu | \hat{v} | \sigma\tau \rangle_{AS}.$$

The antisymmetric matrix element is also hermitian, implying

$$\langle \mu\nu | \hat{v} | \sigma\tau \rangle_{AS} = \langle \sigma\tau | \hat{v} | \mu\nu \rangle_{AS}.$$

With these notations we rewrite the Hartree-Fock functional as

$$\int \Phi^* \hat{H}_I \Phi d\tau = \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \langle \mu\nu | \hat{v} | \mu\nu \rangle_{AS}. \quad (4.1)$$

Adding the contribution from the one-body operator \hat{H}_0 to (??) we obtain the energy functional

$$E[\Phi] = \sum_{\mu=1}^A \langle \mu | h | \mu \rangle + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \langle \mu\nu | \hat{v} | \mu\nu \rangle_{AS}. \quad (4.2)$$

In our coordinate space derivations below we will spell out the Hartree-Fock equations in terms of their integrals.

If we generalize the Euler-Lagrange equations to more variables and introduce N^2 Lagrange multipliers which we denote by $\varepsilon_{\mu\nu}$, we can write the variational equation for the functional of E

$$\delta E - \sum_{\mu\nu} \varepsilon_{\mu\nu} \delta \int \psi_{\mu}^* \psi_{\nu} = 0.$$

For the orthogonal wave functions ψ_i this reduces to

$$\delta E - \sum_{\mu=1}^A \varepsilon_{\mu} \delta \int \psi_{\mu}^* \psi_{\mu} = 0.$$

Variation with respect to the single-particle wave functions ψ_μ yields then

$$\begin{aligned} & \sum_{\mu=1}^A \int \delta \psi_\mu^* \hat{h}_0(x_i) \psi_\mu dx_i + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \left[\int \delta \psi_\mu^* \psi_\nu^* \hat{v}(r_{ij}) \psi_\mu \psi_\nu dx_i dx_j - \int \delta \psi_\mu^* \psi_\nu^* \hat{v}(r_{ij}) \psi_\nu \psi_\mu dx_i dx_j \right] + \\ & \sum_{\mu=1}^A \int \psi_\mu^* \hat{h}_0(x_i) \delta \psi_\mu dx_i + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \left[\int \psi_\mu^* \psi_\nu^* \hat{v}(r_{ij}) \delta \psi_\mu \psi_\nu dx_i dx_j - \int \psi_\mu^* \psi_\nu^* \hat{v}(r_{ij}) \psi_\nu \delta \psi_\mu dx_i dx_j \right] - \sum_{\mu=1}^A E_\mu \int \delta \psi_\mu^* \psi_\mu dx_i - \sum_{\mu=1}^A E_\mu \int \psi_\mu^* \delta \psi_\mu dx_i \end{aligned}$$

Although the variations $\delta\psi$ and $\delta\psi^*$ are not independent, they may in fact be treated as such, so that the terms dependent on either $\delta\psi$ and $\delta\psi^*$ individually may be set equal to zero. To see this, simply replace the arbitrary variation $\delta\psi$ by $i\delta\psi$, so that $\delta\psi^*$ is replaced by $-i\delta\psi^*$, and combine the two equations. We thus arrive at the Hartree-Fock equations

$$\left[-\frac{1}{2} \nabla_i^2 + \sum_{\nu=1}^A \int \psi_\nu^*(x_j) \hat{v}(r_{ij}) \psi_\nu(x_j) dx_j \right] \psi_\mu(x_i) - \left[\sum_{\nu=1}^A \int \psi_\nu^*(x_j) \hat{v}(r_{ij}) \psi_\mu(x_j) dx_j \right] \psi_\nu(x_i) = \varepsilon_\mu \psi_\mu(x_i). \quad (4.3)$$

Notice that the integration $\int dx_j$ implies an integration over the spatial coordinates \mathbf{r}_j and a summation over the spin-coordinate of fermion j . We note that the factor of $1/2$ in front of the sum involving the two-body interaction, has been removed. This is due to the fact that we need to vary both $\delta\psi_\mu^*$ and $\delta\psi_\nu^*$. Using the symmetry properties of the two-body interaction and interchanging μ and ν as summation indices, we obtain two identical terms.

The two first terms in the last equation are the one-body kinetic energy and the electron-nucleus potential. The third or *direct* term is the averaged electronic repulsion of the other electrons. As written, the term includes the *self-interaction* of electrons when $\mu = \nu$. The self-interaction is cancelled in the fourth term, or the *exchange* term. The exchange term results from our inclusion of the Pauli principle and the assumed determinantal form of the wave-function. Equation (??), in addition to the kinetic energy and the attraction from the atomic nucleus that confines the motion of a single electron, represents now the motion of a single-particle modified by the two-body interaction. The additional contribution to the Schroedinger equation due to the two-body interaction, represents a mean field set up by all the other bystanding electrons, the latter given by the sum over all single-particle states occupied by N electrons.

The Hartree-Fock equation is an example of an integro-differential equation. These equations involve repeated calculations of integrals, in addition to the solution of a set of coupled differential equations. The Hartree-Fock equations can also be rewritten in terms of an eigenvalue problem. The solution of an eigenvalue problem represents often a more practical algorithm and the solution of coupled integro-differential equations. This alternative derivation of the Hartree-Fock equations is given below.

Analysis of Hartree-Fock equations in coordinate space

A theoretically convenient form of the Hartree-Fock equation is to regard the direct and exchange operator defined through

$$V_\mu^d(x_i) = \int \psi_\mu^*(x_j) \hat{v}(r_{ij}) \psi_\mu(x_j) dx_j$$

and

$$V_\mu^{ex}(x_i) g(x_i) = \left(\int \psi_\mu^*(x_j) \hat{v}(r_{ij}) g(x_j) dx_j \right) \psi_\mu(x_i),$$

respectively.

The function $g(x_i)$ is an arbitrary function, and by the substitution $g(x_i) = \psi_v(x_i)$ we get

$$V_\mu^{ex}(x_i)\psi_v(x_i) = \left(\int \psi_\mu^*(x_j) \hat{v}(r_{ij}) \psi_v(x_j) dx_j \right) \psi_\mu(x_i).$$

We may then rewrite the Hartree-Fock equations as

$$\hat{h}^{HF}(x_i)\psi_v(x_i) = \varepsilon_v \psi_v(x_i),$$

with

$$\hat{h}^{HF}(x_i) = \hat{h}_0(x_i) + \sum_{\mu=1}^A V_\mu^d(x_i) - \sum_{\mu=1}^A V_\mu^{ex}(x_i),$$

and where $\hat{h}_0(i)$ is the one-body part. The latter is normally chosen as a part which yields solutions in closed form. The harmonic oscillator is a classical problem thereof. We normally rewrite the last equation as

$$\hat{h}^{HF}(x_i) = \hat{h}_0(x_i) + \hat{u}^{HF}(x_i).$$

Hartree-Fock by varying the coefficients of a wave function expansion

Another possibility is to expand the single-particle functions in a known basis and vary the coefficients, that is, the new single-particle wave function is written as a linear expansion in terms of a fixed chosen orthogonal basis (for example the well-known harmonic oscillator functions or the hydrogen-like functions etc). We define our new Hartree-Fock single-particle basis by performing a unitary transformation on our previous basis (labelled with greek indices) as

$$\psi_p^{HF} = \sum_{\lambda} C_{p\lambda} \phi_{\lambda}. \quad (4.4)$$

In this case we vary the coefficients $C_{p\lambda}$. If the basis has infinitely many solutions, we need to truncate the above sum. We assume that the basis ϕ_{λ} is orthogonal.

It is normal to choose a single-particle basis defined as the eigenfunctions of parts of the full Hamiltonian. The typical situation consists of the solutions of the one-body part of the Hamiltonian, that is we have

$$\hat{h}_0 \phi_{\lambda} = \varepsilon_{\lambda} \phi_{\lambda}.$$

The single-particle wave functions $\phi_{\lambda}(\mathbf{r})$, defined by the quantum numbers λ and \mathbf{r} are defined as the overlap

$$\phi_{\lambda}(\mathbf{r}) = \langle \mathbf{r} | \lambda \rangle.$$

In deriving the Hartree-Fock equations, we will expand the single-particle functions in a known basis and vary the coefficients, that is, the new single-particle wave function is written as a linear expansion in terms of a fixed chosen orthogonal basis (for example the well-known harmonic oscillator functions or the hydrogen-like functions etc).

We stated that a unitary transformation keeps the orthogonality. To see this consider first a basis of vectors \mathbf{v}_i ,

$$\mathbf{v}_i = \begin{bmatrix} v_{i1} \\ \vdots \\ \vdots \\ v_{in} \end{bmatrix}$$

We assume that the basis is orthogonal, that is

$$\mathbf{v}_j^T \mathbf{v}_i = \delta_{ij}.$$

An orthogonal or unitary transformation

$$\mathbf{w}_i = \mathbf{U} \mathbf{v}_i,$$

preserves the dot product and orthogonality since

$$\mathbf{w}_j^T \mathbf{w}_i = (\mathbf{U} \mathbf{v}_j)^T \mathbf{U} \mathbf{v}_i = \mathbf{v}_j^T \mathbf{U}^T \mathbf{U} \mathbf{v}_i = \mathbf{v}_j^T \mathbf{v}_i = \delta_{ij}.$$

This means that if the coefficients $C_{p\lambda}$ belong to a unitary or orthogonal transformation (using the Dirac bra-ket notation)

$$|p\rangle = \sum_{\lambda} C_{p\lambda} |\lambda\rangle,$$

orthogonality is preserved, that is $\langle \alpha | \beta \rangle = \delta_{\alpha\beta}$ and $\langle p | q \rangle = \delta_{pq}$.

This property is extremely useful when we build up a basis of many-body Slater determinant based states.

Note also that although a basis $|\alpha\rangle$ contains an infinity of states, for practical calculations we have always to make some truncations.

Before we develop the Hartree-Fock equations, there is another very useful property of determinants that we will use both in connection with Hartree-Fock calculations and later shell-model calculations.

Consider the following determinant

$$\begin{vmatrix} \alpha_1 b_{11} + \alpha_2 b_{12} & a_{12} \\ \alpha_1 b_{21} + \alpha_2 b_{22} & a_{22} \end{vmatrix} = \alpha_1 \begin{vmatrix} b_{11} & a_{12} \\ b_{21} & a_{22} \end{vmatrix} + \alpha_2 \begin{vmatrix} b_{12} & a_{12} \\ b_{22} & a_{22} \end{vmatrix}$$

We can generalize this to an $n \times n$ matrix and have

$$\begin{vmatrix} a_{11} & a_{12} & \dots & \sum_{k=1}^n c_k b_{1k} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & \sum_{k=1}^n c_k b_{2k} & \dots & a_{2n} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ a_{n1} & a_{n2} & \dots & \sum_{k=1}^n c_k b_{nk} & \dots & a_{nn} \end{vmatrix} = \sum_{k=1}^n c_k \begin{vmatrix} a_{11} & a_{12} & \dots & b_{1k} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & b_{2k} & \dots & a_{2n} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ a_{n1} & a_{n2} & \dots & b_{nk} & \dots & a_{nn} \end{vmatrix}.$$

This is a property we will use in our Hartree-Fock discussions.

We can generalize the previous results, now with all elements a_{ij} being given as functions of linear combinations of various coefficients c and elements b_{ij} ,

$$\begin{vmatrix} \sum_{k=1}^n b_{1k} c_{k1} & \sum_{k=1}^n b_{1k} c_{k2} & \dots & \sum_{k=1}^n b_{1k} c_{kj} & \dots & \sum_{k=1}^n b_{1k} c_{kn} \\ \sum_{k=1}^n b_{2k} c_{k1} & \sum_{k=1}^n b_{2k} c_{k2} & \dots & \sum_{k=1}^n b_{2k} c_{kj} & \dots & \sum_{k=1}^n b_{2k} c_{kn} \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ \sum_{k=1}^n b_{nk} c_{k1} & \sum_{k=1}^n b_{nk} c_{k2} & \dots & \sum_{k=1}^n b_{nk} c_{kj} & \dots & \sum_{k=1}^n b_{nk} c_{kn} \end{vmatrix} = \det(\mathbf{C}) \det(\mathbf{B}),$$

where $\det(\mathbf{C})$ and $\det(\mathbf{B})$ are the determinants of $n \times n$ matrices with elements c_{ij} and b_{ij} respectively. This is a property we will use in our Hartree-Fock discussions. Convince yourself about the correctness of the above expression by setting $n = 2$.

With our definition of the new basis in terms of an orthogonal basis we have

$$\psi_p(x) = \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x).$$

If the coefficients $C_{p\lambda}$ belong to an orthogonal or unitary matrix, the new basis is also orthogonal. Our Slater determinant in the new basis $\psi_p(x)$ is written as

$$\frac{1}{\sqrt{A!}} \begin{vmatrix} \psi_p(x_1) & \psi_p(x_2) & \dots & \dots & \psi_p(x_A) \\ \psi_q(x_1) & \psi_q(x_2) & \dots & \dots & \psi_q(x_A) \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \psi_t(x_1) & \psi_t(x_2) & \dots & \dots & \psi_t(x_A) \end{vmatrix} = \frac{1}{\sqrt{A!}} \begin{vmatrix} \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{p\lambda} \phi_{\lambda}(x_A) \\ \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{q\lambda} \phi_{\lambda}(x_A) \\ \dots & \dots & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_1) & \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_2) & \dots & \dots & \sum_{\lambda} C_{t\lambda} \phi_{\lambda}(x_A) \end{vmatrix},$$

which is nothing but $\det(\mathbf{C})\det(\Phi)$, with $\det(\Phi)$ being the determinant given by the basis functions $\phi_{\lambda}(x)$.

In our discussions hereafter we will use our definitions of single-particle states above and below the Fermi (F) level given by the labels $ijkl \dots \leq F$ for so-called single-hole states and $abcd \dots > F$ for so-called particle states. For general single-particle states we employ the labels $pqrs \dots$.

In Eq. (??), restated here

$$E[\Phi] = \sum_{\mu=1}^A \langle \mu | h | \mu \rangle + \frac{1}{2} \sum_{\mu=1}^A \sum_{\nu=1}^A \langle \mu \nu | \hat{v} | \mu \nu \rangle_{AS},$$

we found the expression for the energy functional in terms of the basis function $\phi_{\lambda}(\mathbf{r})$. We then varied the above energy functional with respect to the basis functions $|\mu\rangle$. Now we are interested in defining a new basis defined in terms of a chosen basis as defined in Eq. (1.11). We can then rewrite the energy functional as

$$E[\Phi^{HF}] = \sum_{i=1}^A \langle i | h | i \rangle + \frac{1}{2} \sum_{ij=1}^A \langle ij | \hat{v} | ij \rangle_{AS}, \quad (4.5)$$

where Φ^{HF} is the new Slater determinant defined by the new basis of Eq. (1.11).

Using Eq. (1.11) we can rewrite Eq. (1.12) as

$$E[\Psi] = \sum_{i=1}^A \sum_{\alpha\beta} C_{i\alpha}^* C_{i\beta} \langle \alpha | h | \beta \rangle + \frac{1}{2} \sum_{ij=1}^A \sum_{\alpha\beta\gamma\delta} C_{i\alpha}^* C_{j\beta}^* C_{i\gamma} C_{j\delta} \langle \alpha\beta | \hat{v} | \gamma\delta \rangle_{AS}. \quad (4.6)$$

We wish now to minimize the above functional. We introduce again a set of Lagrange multipliers, noting that since $\langle i | j \rangle = \delta_{i,j}$ and $\langle \alpha | \beta \rangle = \delta_{\alpha,\beta}$, the coefficients $C_{i\gamma}$ obey the relation

$$\langle i | j \rangle = \delta_{i,j} = \sum_{\alpha\beta} C_{i\alpha}^* C_{j\beta} \langle \alpha | \beta \rangle = \sum_{\alpha} C_{i\alpha}^* C_{j\alpha},$$

which allows us to define a functional to be minimized that reads

$$F[\Phi^{HF}] = E[\Phi^{HF}] - \sum_{i=1}^A \epsilon_i \sum_{\alpha} C_{i\alpha}^* C_{i\alpha}. \quad (4.7)$$

Minimizing with respect to $C_{i\alpha}^*$, remembering that the equations for $C_{i\alpha}^*$ and $C_{i\alpha}$ can be written as two independent equations, we obtain

$$\frac{d}{dC_{i\alpha}^*} \left[E[\Phi^{HF}] - \sum_j \epsilon_j \sum_{\alpha} C_{j\alpha}^* C_{j\alpha} \right] = 0,$$

which yields for every single-particle state i and index α (recalling that the coefficients $C_{i\alpha}$ are matrix elements of a unitary (or orthogonal for a real symmetric matrix) matrix) the following

Hartree-Fock equations

$$\sum_{\beta} C_{i\beta} \langle \alpha | h | \beta \rangle + \sum_{j=1}^A \sum_{\gamma\delta} C_{j\beta}^* C_{j\delta} C_{i\gamma} \langle \alpha \beta | \hat{v} | \gamma \delta \rangle_{AS} = \epsilon_i^{HF} C_{i\alpha}.$$

We can rewrite this equation as (changing dummy variables)

$$\sum_{\beta} \left\{ \langle \alpha | h | \beta \rangle + \sum_j \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | \hat{v} | \beta \delta \rangle_{AS} \right\} C_{i\beta} = \epsilon_i^{HF} C_{i\alpha}.$$

Note that the sums over greek indices run over the number of basis set functions (in principle an infinite number).

Defining

$$h_{\alpha\beta}^{HF} = \langle \alpha | h | \beta \rangle + \sum_{j=1}^A \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | \hat{v} | \beta \delta \rangle_{AS},$$

we can rewrite the new equations as

$$\sum_{\beta} h_{\alpha\beta}^{HF} C_{i\beta} = \epsilon_i^{HF} C_{i\alpha}. \quad (4.8)$$

The latter is nothing but a standard eigenvalue problem. Compared with Eq. (??), we see that we do not need to compute any integrals in an iterative procedure for solving the equations. It suffices to tabulate the matrix elements $\langle \alpha | h | \beta \rangle$ and $\langle \alpha \gamma | \hat{v} | \beta \delta \rangle_{AS}$ once and for all. Successive iterations require thus only a look-up in tables over one-body and two-body matrix elements. These details will be discussed below when we solve the Hartree-Fock equations numerical.

Hartree-Fock algorithm

Our Hartree-Fock matrix is thus

$$\hat{h}_{\alpha\beta}^{HF} = \langle \alpha | \hat{h}_0 | \beta \rangle + \sum_{j=1}^A \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | \hat{v} | \beta \delta \rangle_{AS}.$$

The Hartree-Fock equations are solved in an iterative waym starting with a guess for the coefficients $C_{j\gamma} = \delta_{j,\gamma}$ and solving the equations by diagonalization till the new single-particle energies ϵ_i^{HF} do not change anymore by a prefixed quantity.

Normally we assume that the single-particle basis $|\beta\rangle$ forms an eigenbasis for the operator \hat{h}_0 , meaning that the Hartree-Fock matrix becomes

$$\hat{h}_{\alpha\beta}^{HF} = \epsilon_{\alpha} \delta_{\alpha,\beta} + \sum_{j=1}^A \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | \hat{v} | \beta \delta \rangle_{AS}.$$

The Hartree-Fock eigenvalue problem

$$\sum_{\beta} \hat{h}_{\alpha\beta}^{HF} C_{i\beta} = \epsilon_i^{HF} C_{i\alpha},$$

can be written out in a more compact form as

$$\hat{h}^{HF} \hat{C} = \epsilon^{HF} \hat{C}.$$

The Hartree-Fock equations are, in their simplest form, solved in an iterative way, starting with a guess for the coefficients $C_{i\alpha}$. We label the coefficients as $C_{i\alpha}^{(n)}$, where the subscript n stands for iteration n . To set up the algorithm we can proceed as follows:

- We start with a guess $C_{i\alpha}^{(0)} = \delta_{i,\alpha}$. Alternatively, we could have used random starting values as long as the vectors are normalized. Another possibility is to give states below the Fermi level a larger weight.
- The Hartree-Fock matrix simplifies then to (assuming that the coefficients $C_{i\alpha}$ are real)

$$\hat{h}_{\alpha\beta}^{HF} = \epsilon_{\alpha}\delta_{\alpha,\beta} + \sum_{j=1}^A \sum_{\gamma\delta} C_{j\gamma}^{(0)} C_{j\delta}^{(0)} \langle \alpha\gamma | \hat{v} | \beta\delta \rangle_{AS}.$$

Solving the Hartree-Fock eigenvalue problem yields then new eigenvectors $C_{i\alpha}^{(1)}$ and eigenvalues $\epsilon_i^{HF(1)}$.

- With the new eigenvalues we can set up a new Hartree-Fock potential

$$\sum_{j=1}^A \sum_{\gamma\delta} C_{j\gamma}^{(1)} C_{j\delta}^{(1)} \langle \alpha\gamma | \hat{v} | \beta\delta \rangle_{AS}.$$

The diagonalization with the new Hartree-Fock potential yields new eigenvectors and eigenvalues. This process is continued till for example

$$\frac{\sum_p |\epsilon_i^{(n)} - \epsilon_i^{(n-1)}|}{m} \leq \lambda,$$

where λ is a user prefixed quantity ($\lambda \sim 10^{-8}$ or smaller) and p runs over all calculated single-particle energies and m is the number of single-particle states.

Analysis of Hartree-Fock equations and Koopman's theorem

We can rewrite the ground state energy by adding and subtracting $\hat{u}^{HF}(x_i)$

$$E_0^{HF} = \langle \Phi_0 | \hat{H} | \Phi_0 \rangle = \sum_{i \leq F} \langle i | \hat{h}_0 + \hat{u}^{HF} | j \rangle + \frac{1}{2} \sum_{i \leq F} \sum_{j \leq F} [\langle ij | \hat{v} | ij \rangle - \langle ij | \hat{v} | ji \rangle] - \sum_{i \leq F} \langle i | \hat{u}^{HF} | i \rangle,$$

which results in

$$E_0^{HF} = \sum_{i \leq F} \epsilon_i^{HF} + \frac{1}{2} \sum_{i \leq F} \sum_{j \leq F} [\langle ij | \hat{v} | ij \rangle - \langle ij | \hat{v} | ji \rangle] - \sum_{i \leq F} \langle i | \hat{u}^{HF} | i \rangle.$$

Our single-particle states $ijk \dots$ are now single-particle states obtained from the solution of the Hartree-Fock equations.

Using our definition of the Hartree-Fock single-particle energies we obtain then the following expression for the total ground-state energy

$$E_0^{HF} = \sum_{i \leq F} \epsilon_i - \frac{1}{2} \sum_{i \leq F} \sum_{j \leq F} [\langle ij | \hat{v} | ij \rangle - \langle ij | \hat{v} | ji \rangle].$$

This form will be used in our discussion of Koopman's theorem.

In the atomic physics case we have

$$E[\Phi^{\text{HF}}(N)] = \sum_{i=1}^H \langle i | \hat{h}_0 | i \rangle + \frac{1}{2} \sum_{ij=1}^N \langle ij | \hat{v} | ij \rangle_{AS},$$

where $\Phi^{\text{HF}}(N)$ is the new Slater determinant defined by the new basis of Eq. (1.11) for N electrons (same Z). If we assume that the single-particle wave functions in the new basis do not change when we remove one electron or add one electron, we can then define the corresponding energy for the $N - 1$ systems as

$$E[\Phi^{\text{HF}}(N-1)] = \sum_{i=1; i \neq k}^N \langle i | \hat{h}_0 | i \rangle + \frac{1}{2} \sum_{ij=1; i, j \neq k}^N \langle ij | \hat{v} | ij \rangle_{AS},$$

where we have removed a single-particle state $k \leq F$, that is a state below the Fermi level.

Calculating the difference

$$E[\Phi^{\text{HF}}(N)] - E[\Phi^{\text{HF}}(N-1)] = \langle k | \hat{h}_0 | k \rangle + \frac{1}{2} \sum_{i=1; i \neq k}^N \langle ik | \hat{v} | ik \rangle_{AS} + \frac{1}{2} \sum_{j=1; j \neq k}^N \langle kj | \hat{v} | kj \rangle_{AS},$$

we obtain

$$E[\Phi^{\text{HF}}(N)] - E[\Phi^{\text{HF}}(N-1)] = \langle k | \hat{h}_0 | k \rangle + \sum_{j=1}^N \langle kj | \hat{v} | kj \rangle_{AS}$$

which is just our definition of the Hartree-Fock single-particle energy

$$E[\Phi^{\text{HF}}(N)] - E[\Phi^{\text{HF}}(N-1)] = \epsilon_k^{\text{HF}}$$

Similarly, we can now compute the difference (we label the single-particle states above the Fermi level as $abcd > F$)

$$E[\Phi^{\text{HF}}(N+1)] - E[\Phi^{\text{HF}}(N)] = \epsilon_a^{\text{HF}}.$$

These two equations can thus be used to the electron affinity or ionization energies, respectively. Koopman's theorem states that for example the ionization energy of a closed-shell system is given by the energy of the highest occupied single-particle state. If we assume that changing the number of electrons from N to $N + 1$ does not change the Hartree-Fock single-particle energies and eigenfunctions, then Koopman's theorem simply states that the ionization energy of an atom is given by the single-particle energy of the last bound state. In a similar way, we can also define the electron affinities.

As an example, consider a simple model for atomic sodium, Na. Neutral sodium has eleven electrons, with the weakest bound one being confined the $3s$ single-particle quantum numbers. The energy needed to remove an electron from neutral sodium is rather small, 5.1391 eV, a feature which pertains to all alkali metals. Having performed a Hartree-Fock calculation for neutral sodium would then allows us to compute the ionization energy by using the single-particle energy for the $3s$ states, namely $\epsilon_{3s}^{\text{HF}}$.

From these considerations, we see that Hartree-Fock theory allows us to make a connection between experimental observables (here ionization and affinity energies) and the underlying interactions between particles. In this sense, we are now linking the dynamics and structure of a many-body system with the laws of motion which govern the system. Our approach is a reductionistic one, meaning that we want to understand the laws of motion in terms of the particles or degrees of freedom which we believe are the fundamental ones. Our Slater determinant, being constructed as the product of various single-particle functions, follows this philosophy.

With similar arguments as in atomic physics, we can now use Hartree-Fock theory to make a link between nuclear forces and separation energies. Changing to nuclear system, we define

$$E[\Phi^{\text{HF}}(A)] = \sum_{i=1}^A \langle i | \hat{h}_0 | i \rangle + \frac{1}{2} \sum_{ij=1}^A \langle ij | \hat{v} | ij \rangle_{AS},$$

where $\Phi^{\text{HF}}(A)$ is the new Slater determinant defined by the new basis of Eq. (1.11) for A nucleons, where $A = N + Z$, with N now being the number of neutrons and Z the number of protons. If we assume again that the single-particle wave functions in the new basis do not change from a nucleus with A nucleons to a nucleus with $A - 1$ nucleons, we can then define the corresponding energy for the $A - 1$ systems as

$$E[\Phi^{\text{HF}}(A - 1)] = \sum_{i=1; i \neq k}^A \langle i | \hat{h}_0 | i \rangle + \frac{1}{2} \sum_{ij=1; i, j \neq k}^A \langle ij | \hat{v} | ij \rangle_{AS},$$

where we have removed a single-particle state $k \leq F$, that is a state below the Fermi level.

Calculating the difference

$$E[\Phi^{\text{HF}}(A)] - E[\Phi^{\text{HF}}(A - 1)] = \langle k | \hat{h}_0 | k \rangle + \frac{1}{2} \sum_{i=1; i \neq k}^A \langle ik | \hat{v} | ik \rangle_{AS} + \frac{1}{2} \sum_{j=1; j \neq k}^A \langle kj | \hat{v} | kj \rangle_{AS},$$

which becomes

$$E[\Phi^{\text{HF}}(A)] - E[\Phi^{\text{HF}}(A - 1)] = \langle k | \hat{h}_0 | k \rangle + \sum_{j=1}^A \langle kj | \hat{v} | kj \rangle_{AS}$$

which is just our definition of the Hartree-Fock single-particle energy

$$E[\Phi^{\text{HF}}(A)] - E[\Phi^{\text{HF}}(A - 1)] = \epsilon_k^{\text{HF}}$$

Similarly, we can now compute the difference (recall that the single-particle states $abcd > F$)

$$E[\Phi^{\text{HF}}(A + 1)] - E[\Phi^{\text{HF}}(A)] = \epsilon_a^{\text{HF}}.$$

If we then recall that the binding energy differences

$$BE(A) - BE(A - 1) \quad \text{and} \quad BE(A + 1) - BE(A),$$

define the separation energies, we see that the Hartree-Fock single-particle energies can be used to define separation energies. We have thus our first link between nuclear forces (included in the potential energy term) and an observable quantity defined by differences in binding energies.

We have thus the following interpretations (if the single-particle fields do not change)

$$BE(A) - BE(A - 1) \approx E[\Phi^{\text{HF}}(A)] - E[\Phi^{\text{HF}}(A - 1)] = \epsilon_k^{\text{HF}},$$

and

$$BE(A + 1) - BE(A) \approx E[\Phi^{\text{HF}}(A + 1)] - E[\Phi^{\text{HF}}(A)] = \epsilon_a^{\text{HF}}.$$

If we use ^{16}O as our closed-shell nucleus, we could then interpret the separation energy

$$BE(^{16}\text{O}) - BE(^{15}\text{O}) \approx \epsilon_{0p_{1/2}^v}^{\text{HF}},$$

and

$$BE(^{16}\text{O}) - BE(^{15}\text{N}) \approx \epsilon_{0p_{1/2}^\pi}^{\text{HF}}.$$

Similarlry, we could interpret

$$BE(^{17}\text{O}) - BE(^{16}\text{O}) \approx \epsilon_{0d_{5/2}^v}^{\text{HF}},$$

and

$$BE(^{17}\text{F}) - BE(^{16}\text{O}) \approx \epsilon_{0d_{5/2}^\pi}^{\text{HF}}.$$

We can continue like this for all $A \pm 1$ nuclei where A is a good closed-shell (or subshell closure) nucleus. Examples are ^{22}O , ^{24}O , ^{40}Ca , ^{48}Ca , ^{52}Ca , ^{54}Ca , ^{56}Ni , ^{68}Ni , ^{78}Ni , ^{90}Zr , ^{88}Sr , ^{100}Sn , ^{132}Sn and ^{208}Pb , to mention some possible cases.

We can thus make our first interpretation of the separation energies in terms of the simplest possible many-body theory. If we also recall that the so-called energy gap for neutrons (or protons) is defined as

$$\Delta S_n = 2BE(N, Z) - BE(N-1, Z) - BE(N+1, Z),$$

for neutrons and the corresponding gap for protons

$$\Delta S_p = 2BE(N, Z) - BE(N, Z-1) - BE(N, Z+1),$$

we can define the neutron and proton energy gaps for ^{16}O as

$$\Delta S_v = \epsilon_{0d_{5/2}^v}^{\text{HF}} - \epsilon_{0p_{1/2}^v}^{\text{HF}},$$

and

$$\Delta S_\pi = \epsilon_{0d_{5/2}^\pi}^{\text{HF}} - \epsilon_{0p_{1/2}^\pi}^{\text{HF}}.$$

*Exercise : Derivation of Hartree-Fock equations

Consider a Slater determinant built up of single-particle orbitals ψ_λ , with $\lambda = 1, 2, \dots, N$.

The unitary transformation

$$\psi_a = \sum_\lambda C_{a\lambda} \phi_\lambda,$$

brings us into the new basis. The new basis has quantum numbers $a = 1, 2, \dots, N$.

a) Show that the new basis is orthonormal.

b) Show that the new Slater determinant constructed from the new single-particle wave functions can be written as the determinant based on the previous basis and the determinant of the matrix C .

c) Show that the old and the new Slater determinants are equal up to a complex constant with absolute value unity.

Hint.

Use the fact that C is a unitary matrix.

*Exercise : Derivation of Hartree-Fock equations

Consider the Slater determinant

$$\Phi_0 = \frac{1}{\sqrt{n!}} \sum_p (-)^p P \prod_{i=1}^n \psi_{\alpha_i}(x_i).$$

A small variation in this function is given by

$$\delta \Phi_0 = \frac{1}{\sqrt{n!}} \sum_p (-)^p P \psi_{\alpha_1}(x_1) \psi_{\alpha_2}(x_2) \dots \psi_{\alpha_{i-1}}(x_{i-1}) (\delta \psi_{\alpha_i}(x_i)) \psi_{\alpha_{i+1}}(x_{i+1}) \dots \psi_{\alpha_n}(x_n).$$

a) Show that

$$\langle \delta \Phi_0 | \sum_{i=1}^n \{t(x_i) + u(x_i)\} + \frac{1}{2} \sum_{i \neq j=1}^n v(x_i, x_j) | \Phi_0 \rangle = \sum_{i=1}^n \langle \delta \psi_{\alpha_i} | \hat{t} + \hat{u} | \phi_{\alpha_i} \rangle + \sum_{i \neq j=1}^n \{ \langle \delta \psi_{\alpha_i} \psi_{\alpha_j} | \hat{v} | \psi_{\alpha_i} \psi_{\alpha_j} \rangle - \langle \delta \psi_{\alpha_i} \psi_{\alpha_j} | \hat{v} | \psi_{\alpha_j} \psi_{\alpha_i} \rangle \}$$

*Exercise : Developing a Hartree-Fock program

Neutron drops are a powerful theoretical laboratory for testing, validating and improving nuclear structure models. Indeed, all approaches to nuclear structure, from ab initio theory to shell model to density functional theory are applicable in such systems. We will, therefore, use neutron drops as a test system for setting up a Hartree-Fock code. This program can later be extended to studies of the binding energy of nuclei like ^{16}O or ^{40}Ca . The single-particle energies obtained by solving the Hartree-Fock equations can then be directly related to experimental separation energies. Since Hartree-Fock theory is the starting point for several many-body techniques (density functional theory, random-phase approximation, shell-model etc), the aim here is to develop a computer program to solve the Hartree-Fock equations in a given single-particle basis, here the harmonic oscillator.

The Hamiltonian for a system of N neutron drops confined in a harmonic potential reads

$$\hat{H} = \sum_{i=1}^N \frac{\hat{p}_i^2}{2m} + \sum_{i=1}^N \frac{1}{2} m \omega r_i^2 + \sum_{i < j} \hat{V}_{ij},$$

with $\hbar^2/2m = 20.73 \text{ fm}^2$, $mc^2 = 938.90590 \text{ MeV}$, and \hat{V}_{ij} is the two-body interaction potential whose matrix elements are precalculated and to be read in by you.

The Hartree-Fock algorithm can be broken down as follows. We recall that our Hartree-Fock matrix is

$$\hat{h}_{\alpha\beta}^{HF} = \langle \alpha | \hat{h}_0 | \beta \rangle + \sum_{j=1}^N \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | V | \beta \delta \rangle_{AS}.$$

Normally we assume that the single-particle basis $|\beta\rangle$ forms an eigenbasis for the operator \hat{h}_0 (this is our case), meaning that the Hartree-Fock matrix becomes

$$\hat{h}_{\alpha\beta}^{HF} = \varepsilon_\alpha \delta_{\alpha,\beta} + \sum_{j=1}^N \sum_{\gamma\delta} C_{j\gamma}^* C_{j\delta} \langle \alpha \gamma | V | \beta \delta \rangle_{AS}.$$

The Hartree-Fock eigenvalue problem

$$\sum_{\beta} \hat{h}_{\alpha\beta}^{HF} C_{i\beta} = \varepsilon_i^{HF} C_{i\alpha},$$

can be written out in a more compact form as

$$\hat{h}^{HF} \hat{C} = \varepsilon^{HF} \hat{C}.$$

The equations are often rewritten in terms of a so-called density matrix, which is defined as

$$\rho_{\gamma\delta} = \sum_{i=1}^N \langle \gamma | i \rangle \langle i | \delta \rangle = \sum_{i=1}^N C_{i\gamma} C_{i\delta}^*. \quad (4.9)$$

It means that we can rewrite the Hartree-Fock Hamiltonian as

$$\hat{h}_{\alpha\beta}^{HF} = \varepsilon_\alpha \delta_{\alpha,\beta} + \sum_{\gamma\delta} \rho_{\gamma\delta} \langle \alpha \gamma | V | \beta \delta \rangle_{AS}.$$

It is convenient to use the density matrix since we can precalculate in every iteration the product of two eigenvector components C .

Note that $\langle \alpha | \hat{h}_0 | \beta \rangle$ denotes the matrix elements of the one-body part of the starting hamiltonian. For self-bound nuclei $\langle \alpha | \hat{h}_0 | \beta \rangle$ is the kinetic energy, whereas for neutron drops, $\langle \alpha | \hat{h}_0 | \beta \rangle$ represents the harmonic oscillator hamiltonian since the system is confined in a harmonic trap. If we are working in a harmonic oscillator basis with the same ω as the trapping potential, then $\langle \alpha | \hat{h}_0 | \beta \rangle$ is diagonal.

The python program shows how one can, in a brute force way read in matrix elements in m -scheme and compute the Hartree-Fock single-particle energies for four major shells. The interaction which has been used is the so-called N3LO interaction of Machleidt and Entem using the Similarity Renormalization Group approach method to renormalize the interaction, using an oscillator energy $\omega = 10$ MeV.

The nucleon-nucleon two-body matrix elements are in m -scheme and are fully anti-symmetrized. The Hartree-Fock programs uses the density matrix discussed above in order to compute the Hartree-Fock matrix. Here we display the Hartree-Fock part only, assuming that single-particle data and two-body matrix elements have already been read in.

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
import numpy as np from decimal import Decimal expectation value for the one body
part, Harmonic oscillator in three dimensions def onebody(i, n, l): homega = 10.0 return
homega*(2*n[i] + l[i] + 1.5)
```

```
if name == '__main__':
    Nparticles = 16 """ Read quantum numbers from file """ index = [] n = [] l = [] j = [] mj =
    [] tz = [] spOrbitals = 0 with open("nucleispnumbers.dat", "r") as qnumfile: for line in qnum-
    file: nums = line.split() if len(nums) != 0: index.append(int(nums[0])) n.append(int(nums[1]))
    l.append(int(nums[2])) j.append(int(nums[3])) mj.append(int(nums[4])) tz.append(int(nums[5]))
    spOrbitals += 1
    """ Read two-nucleon interaction elements (integrals) from file, brute force 4-dim array """
    nninteraction = np.zeros([spOrbitals, spOrbitals, spOrbitals, spOrbitals]) with open("nucleitwobody.dat",
    "r") as infile: for line in infile: number = line.split() a = int(number[0]) - 1 b = int(number[1]) -
    1 c = int(number[2]) - 1 d = int(number[3]) - 1 nninteraction[a][b][c][d] = Decimal(number[4])
    """ Set up single-particle integral """ singleparticleH = np.zeros(spOrbitals) for i in range(spOrbitals):
    singleparticleH[i] = Decimal(onebody(i, n, l))
    """ Star HF-iterations, preparing variables and density matrix """
    """ Coefficients for setting up density matrix, assuming only one along the diagonals """
    C = np.eye(spOrbitals) HF coefficients DensityMatrix = np.zeros([spOrbitals, spOrbitals]) for
    gamma in range(spOrbitals): for delta in range(spOrbitals): sum = 0.0 for i in range(Nparticles):
    sum += C[gamma][i]*C[delta][i] DensityMatrix[gamma][delta] = Decimal(sum) maxHFIter
    = 100 epsilon = 1.0e-5 difference = 1.0 hfcount = 0 oldenergies = np.zeros(spOrbitals) newenergies =
    np.zeros(spOrbitals) while hfcount < maxHFIter and difference > epsilon: print("Iteration HF matrix = np.zeros([spOrbitals, spOrbitals]) for
    for beta in range(spOrbitals): """ If tests for three-dimensional systems, including isospin conservation """ if l[alpha] !=
    l[beta] and j[alpha] != j[beta] and mj[alpha] != mj[beta] and tz[alpha] != tz[beta]: continue """ Setting up the Fock matrix using the density matrix and
    scheme """ sumFockTerm = 0.0 for gamma in range(spOrbitals): for delta in range(spOrbitals): if (mj[alpha] + mj[gamma]) !=
    (mj[beta] + mj[delta]) and (tz[alpha] + tz[gamma]) != (tz[beta] + tz[delta]): continue sumFockTerm += DensityMatrix[gamma][delta] *
    nninteraction[alpha][gamma][beta][delta] HF matrix[alpha][beta] = Decimal(sumFockTerm) """ Adding the one-body term, here plain harmonic os-
    alpha: HF matrix[alpha][alpha] += singleparticleH[alpha] spenergies, C = np.linalg.eigh(HF matrix) """ Setting up new density matrix in m-
    scheme """ DensityMatrix = np.zeros([spOrbitals, spOrbitals]) for gamma in range(spOrbitals): for delta in range(spOrbitals):
    sum = 0.0 for i in range(Nparticles): sum += C[gamma][i]*C[delta][i] DensityMatrix[gamma][delta] = Decimal(sum) newenergies =
    spenergies """ Brute force computation of difference between previous and new sp HF energies """ sum = 0.0 for i in range(spOrbitals):
    sum += (abs(newenergies[i] - oldenergies[i]))/spOrbitals difference = sum oldenergies = newenergies print("Single -
    particle energies, ordering may have changed") for i in range(spOrbitals): print('0: 4d1: .4f'.format(i, Decimal(oldenergies[i]))) hfcount +=
    1
```

Running the program, one finds that the lowest-lying states for a nucleus like ^{16}O , we see that the nucleon-nucleon force brings a natural spin-orbit splitting for the $0p$ states (or

other states except the s -states). Since we are using the m -scheme for our calculations, we observe that there are several states with the same eigenvalues. The number of eigenvalues corresponds to the degeneracy $2j + 1$ and is well respected in our calculations, as see from the table here.

The values of the lowest-lying states are (π for protons and ν for neutrons)

Quantum numbers	Energy [MeV]
$0s_{1/2}^{\pi}$	-40.4602
$0s_{1/2}^{\pi}$	-40.4602
$0s_{1/2}^{\nu}$	-40.6426
$0s_{1/2}^{\nu}$	-40.6426
$0p_{1/2}^{\pi}$	-6.7133
$0p_{1/2}^{\pi}$	-6.7133
$0p_{1/2}^{\nu}$	-6.8403
$0p_{1/2}^{\nu}$	-6.8403
$0p_{3/2}^{\pi}$	-11.5886
$0p_{3/2}^{\pi}$	-11.5886
$0p_{3/2}^{\pi}$	-11.5886
$0p_{3/2}^{\pi}$	-11.5886
$0p_{3/2}^{\nu}$	-11.7201
$0p_{3/2}^{\nu}$	-11.7201
$0p_{3/2}^{\nu}$	-11.7201
$0p_{3/2}^{\nu}$	-11.7201
$0d_{5/2}^{\pi}$	18.7589
$0d_{5/2}^{\nu}$	18.8082

We can use these results to attempt our first link with experimental data, namely to compute the shell gap or the separation energies. The shell gap for neutrons is given by

$$\Delta S_n = 2BE(N, Z) - BE(N - 1, Z) - BE(N + 1, Z).$$

For ^{16}O we have an experimental value for the shell gap of 11.51 MeV for neutrons, while our Hartree-Fock calculations result in 25.65 MeV. This means that correlations beyond a simple Hartree-Fock calculation with a two-body force play an important role in nuclear physics. The splitting between the $0p_{3/2}^{\nu}$ and the $0p_{1/2}^{\nu}$ state is 4.88 MeV, while the experimental value for the gap between the ground state $1/2^-$ and the first excited $3/2^-$ states is 6.08 MeV. The two-nucleon spin-orbit force plays a central role here. In our discussion of nuclear forces we will see how the spin-orbit force comes into play here.

Hartree-Fock in second quantization and stability of HF solution

We wish now to derive the Hartree-Fock equations using our second-quantized formalism and study the stability of the equations. Our ansatz for the ground state of the system is approximated as (this is our representation of a Slater determinant in second quantization)

$$|\Phi_0\rangle = |c\rangle = a_i^\dagger a_j^\dagger \dots a_l^\dagger |0\rangle.$$

We wish to determine \hat{u}^{HF} so that $E_0^{HF} = \langle c | \hat{H} | c \rangle$ becomes a local minimum.

In our analysis here we will need Thouless' theorem, which states that an arbitrary Slater determinant $|c'\rangle$ which is not orthogonal to a determinant $|c\rangle = \prod_{i=1}^n a_{\alpha_i}^\dagger |0\rangle$, can be written as

$$|c'\rangle = \exp \left\{ \sum_{a>F} \sum_{i\leq F} C_{ai} a_a^\dagger a_i \right\} |c\rangle$$

Let us give a simple proof of Thouless' theorem. The theorem states that we can make a linear combination of particle-hole excitations with respect to a given reference state $|c\rangle$. With this linear combination, we can make a new Slater determinant $|c'\rangle$ which is not orthogonal to $|c\rangle$, that is

$$\langle c|c'\rangle \neq 0.$$

To show this we need some intermediate steps. The exponential product of two operators $\exp \hat{A} \times \exp \hat{B}$ is equal to $\exp(\hat{A} + \hat{B})$ only if the two operators commute, that is

$$[\hat{A}, \hat{B}] = 0.$$

Thouless' theorem

If the operators do not commute, we need to resort to the Baker-Campbell-Hauersdorf. This relation states that

$$\exp \hat{C} = \exp \hat{A} \exp \hat{B},$$

with

$$\hat{C} = \hat{A} + \hat{B} + \frac{1}{2}[\hat{A}, \hat{B}] + \frac{1}{12}[[\hat{A}, \hat{B}], \hat{B}] - \frac{1}{12}[[\hat{A}, \hat{B}], \hat{A}] + \dots$$

From these relations, we note that in our expression for $|c'\rangle$ we have commutators of the type

$$[a_a^\dagger a_i, a_b^\dagger a_j],$$

and it is easy to convince oneself that these commutators, or higher powers thereof, are all zero. This means that we can write out our new representation of a Slater determinant as

$$|c'\rangle = \exp \left\{ \sum_{a>F} \sum_{i\leq F} C_{ai} a_a^\dagger a_i \right\} |c\rangle = \prod_i \left\{ 1 + \sum_{a>F} C_{ai} a_a^\dagger a_i + \left(\sum_{a>F} C_{ai} a_a^\dagger a_i \right)^2 + \dots \right\} |c\rangle$$

We note that

$$\prod_i \sum_{a>F} C_{ai} a_a^\dagger a_i \sum_{b>F} C_{bi} a_b^\dagger a_i |c\rangle = 0,$$

and all higher-order powers of these combinations of creation and annihilation operators disappear due to the fact that $(a_i)^n |c\rangle = 0$ when $n > 1$. This allows us to rewrite the expression for $|c'\rangle$ as

$$|c'\rangle = \prod_i \left\{ 1 + \sum_{a>F} C_{ai} a_a^\dagger a_i \right\} |c\rangle,$$

which we can rewrite as

$$|c'\rangle = \prod_i \left\{ 1 + \sum_{a>F} C_{ai} a_a^\dagger a_i \right\} |a_{i_1}^\dagger a_{i_2}^\dagger \dots a_{i_n}^\dagger |0\rangle.$$

The last equation can be written as

$$|c'\rangle = \prod_i \left\{ 1 + \sum_{a>F} C_{ai} a_a^\dagger a_i \right\} |a_1^\dagger a_2^\dagger \dots a_n^\dagger |0\rangle = \left(1 + \sum_{a>F} C_{ai_1} a_a^\dagger a_{i_1} \right) a_{i_1}^\dagger \quad (4.10)$$

$$\times \left(1 + \sum_{a>F} C_{ai_2} a_a^\dagger a_{i_2} \right) a_{i_2}^\dagger \dots |0\rangle = \prod_i \left(a_i^\dagger + \sum_{a>F} C_{ai} a_a^\dagger \right) |0\rangle. \quad (4.11)$$

New operators

If we define a new creation operator

$$b_i^\dagger = a_i^\dagger + \sum_{a>F} C_{ai} a_a^\dagger, \quad (4.12)$$

we have

$$|c'\rangle = \prod_i b_i^\dagger |0\rangle = \prod_i \left(a_i^\dagger + \sum_{a>F} C_{ai} a_a^\dagger \right) |0\rangle,$$

meaning that the new representation of the Slater determinant in second quantization, $|c'\rangle$, looks like our previous ones. However, this representation is not general enough since we have a restriction on the sum over single-particle states in Eq. (??). The single-particle states have all to be above the Fermi level. The question then is whether we can construct a general representation of a Slater determinant with a creation operator

$$\tilde{b}_i^\dagger = \sum_p f_{ip} a_p^\dagger,$$

where f_{ip} is a matrix element of a unitary matrix which transforms our creation and annihilation operators a^\dagger and a to \tilde{b}^\dagger and \tilde{b} . These new operators define a new representation of a Slater determinant as

$$|\tilde{c}\rangle = \prod_i \tilde{b}_i^\dagger |0\rangle.$$

Showing that $|\tilde{c}\rangle = |c'\rangle$

We need to show that $|\tilde{c}\rangle = |c'\rangle$. We need also to assume that the new state is not orthogonal to $|c\rangle$, that is $\langle c|\tilde{c}\rangle \neq 0$. From this it follows that

$$\langle c|\tilde{c}\rangle = \langle 0|a_{i_n} \dots a_{i_1} \left(\sum_{p=i_1}^{i_n} f_{ip} a_p^\dagger \right) \left(\sum_{q=i_1}^{i_n} f_{iq} a_q^\dagger \right) \dots \left(\sum_{t=i_1}^{i_n} f_{it} a_t^\dagger \right) |0\rangle,$$

which is nothing but the determinant $\det(f_{ip})$ which we can, using the intermediate normalization condition, normalize to one, that is

$$\det(f_{ip}) = 1,$$

meaning that f has an inverse defined as (since we are dealing with orthogonal, and in our case unitary as well, transformations)

$$\sum_k f_{ik} f_{kj}^{-1} = \delta_{ij},$$

and

$$\sum_j f_{ij}^{-1} f_{jk} = \delta_{ik}.$$

Using these relations we can then define the linear combination of creation (and annihilation as well) operators as

$$\sum_i f_{ki}^{-1} \tilde{b}_i^\dagger = \sum_i f_{ki}^{-1} \sum_{p=i_1}^{\infty} f_{ip} a_p^\dagger = a_k^\dagger + \sum_i \sum_{p=i_{n+1}}^{\infty} f_{ki}^{-1} f_{ip} a_p^\dagger.$$

Defining

$$c_{kp} = \sum_{i \leq F} f_{ki}^{-1} f_{ip},$$

we can redefine

$$a_k^\dagger + \sum_i \sum_{p=i_{n+1}}^{\infty} f_{ki}^{-1} f_{ip} a_p^\dagger = a_k^\dagger + \sum_{p=i_{n+1}}^{\infty} c_{kp} a_p^\dagger = b_k^\dagger,$$

our starting point. We have shown that our general representation of a Slater determinant

$$|\tilde{c}\rangle = \prod_i \tilde{b}_i^\dagger |0\rangle = |c'\rangle = \prod_i b_i^\dagger |0\rangle,$$

with

$$b_k^\dagger = a_k^\dagger + \sum_{p=i_{n+1}}^{\infty} c_{kp} a_p^\dagger.$$

This means that we can actually write an ansatz for the ground state of the system as a linear combination of terms which contain the ansatz itself $|c\rangle$ with an admixture from an infinity of one-particle-one-hole states. The latter has important consequences when we wish to interpret the Hartree-Fock equations and their stability. We can rewrite the new representation as

$$|c'\rangle = |c\rangle + |\delta c\rangle,$$

where $|\delta c\rangle$ can now be interpreted as a small variation. If we approximate this term with contributions from one-particle-one-hole ($1p-1h$) states only, we arrive at

$$|c'\rangle = \left(1 + \sum_{ai} \delta C_{ai} a_a^\dagger a_i \right) |c\rangle.$$

In our derivation of the Hartree-Fock equations we have shown that

$$\langle \delta c | \hat{H} | c \rangle = 0,$$

which means that we have to satisfy

$$\langle c | \sum_{ai} \delta C_{ai} \{ a_a^\dagger a_i \} \hat{H} | c \rangle = 0.$$

With this as a background, we are now ready to study the stability of the Hartree-Fock equations.

Hartree-Fock in second quantization and stability of HF solution

The variational condition for deriving the Hartree-Fock equations guarantees only that the expectation value $\langle c | \hat{H} | c \rangle$ has an extreme value, not necessarily a minimum. To figure out

whether the extreme value we have found is a minimum, we can use second quantization to analyze our results and find a criterion for the above expectation value to a local minimum. We will use Thouless' theorem and show that

$$\frac{\langle c' | \hat{H} | c' \rangle}{\langle c' | c' \rangle} \geq \langle c | \hat{H} | c \rangle = E_0,$$

with

$$|c'\rangle = |c\rangle + |\delta c\rangle.$$

Using Thouless' theorem we can write out $|c'\rangle$ as

$$|c'\rangle = \exp \left\{ \sum_{a>F} \sum_{i\leq F} \delta C_{ai} a_a^\dagger a_i \right\} |c\rangle \quad (4.13)$$

$$= \left\{ 1 + \sum_{a>F} \sum_{i\leq F} \delta C_{ai} a_a^\dagger a_i + \frac{1}{2!} \sum_{ab>F} \sum_{ij\leq F} \delta C_{ai} \delta C_{bj} a_a^\dagger a_i a_b^\dagger a_j + \dots \right\} \quad (4.14)$$

where the amplitudes δC are small.

The norm of $|c'\rangle$ is given by (using the intermediate normalization condition $\langle c' | c \rangle = 1$)

$$\langle c' | c' \rangle = 1 + \sum_{a>F} \sum_{i\leq F} |\delta C_{ai}|^2 + O(\delta C_{ai}^3).$$

The expectation value for the energy is now given by (using the Hartree-Fock condition)

$$\langle c' | \hat{H} | c' \rangle = \langle c | \hat{H} | c \rangle + \sum_{ab>F} \sum_{ij\leq F} \delta C_{ai}^* \delta C_{bj} \langle c | a_i^\dagger a_a \hat{H} a_b^\dagger a_j | c \rangle +$$

$$\frac{1}{2!} \sum_{ab>F} \sum_{ij\leq F} \delta C_{ai} \delta C_{bj} \langle c | \hat{H} a_a^\dagger a_i a_b^\dagger a_j | c \rangle + \frac{1}{2!} \sum_{ab>F} \sum_{ij\leq F} \delta C_{ai}^* \delta C_{bj}^* \langle c | a_j^\dagger a_b a_i^\dagger a_a \hat{H} | c \rangle + \dots$$

We have already calculated the second term on the right-hand side of the previous equation

$$\langle c | \left(\{a_i^\dagger a_a\} \hat{H} \{a_b^\dagger a_j\} \right) | c \rangle = \sum_{pq} \sum_{ijab} \delta C_{ai}^* \delta C_{bj} \langle p | \hat{h}_0 | q \rangle \langle c | \left(\{a_i^\dagger a_a\} \{a_p^\dagger a_q\} \{a_b^\dagger a_j\} \right) | c \rangle \quad (4.15)$$

$$+ \frac{1}{4} \sum_{pqrs} \sum_{ijab} \delta C_{ai}^* \delta C_{bj} \langle pq | \hat{v} | rs \rangle \langle c | \left(\{a_i^\dagger a_a\} \{a_p^\dagger a_q^\dagger a_s a_r\} \{a_b^\dagger a_j\} \right) | c \rangle, \quad (4.16)$$

resulting in

$$E_0 \sum_{ai} |\delta C_{ai}|^2 + \sum_{ai} |\delta C_{ai}|^2 (\epsilon_a - \epsilon_i) - \sum_{ijab} \langle aj | \hat{v} | bi \rangle \delta C_{ai}^* \delta C_{bj}.$$

$$\frac{1}{2!} \langle c | \left(\{a_j^\dagger a_b\} \{a_i^\dagger a_a\} \hat{V}_N \right) | c \rangle = \frac{1}{2!} \langle c | \left(\hat{V}_N \{a_a^\dagger a_i\} \{a_b^\dagger a_j\} \right)^\dagger | c \rangle$$

which is nothing but

$$\frac{1}{2!} \langle c | \left(\hat{V}_N \{a_a^\dagger a_i\} \{a_b^\dagger a_j\} \right) | c \rangle^* = \frac{1}{2} \sum_{ijab} (\langle ij | \hat{v} | ab \rangle)^* \delta C_{ai}^* \delta C_{bj}$$

or

$$\frac{1}{2} \sum_{ijab} (\langle ab | \hat{v} | ij \rangle) \delta C_{ai}^* \delta C_{bj}$$

where we have used the relation

$$\langle a | \hat{A} | b \rangle = (\langle b | \hat{A}^\dagger | a \rangle)^*$$

due to the hermiticity of \hat{H} and \hat{V} .

We define two matrix elements

$$A_{ai,bj} = -\langle aj|\hat{v}bi\rangle$$

and

$$B_{ai,bj} = \langle ab|\hat{v}|ij\rangle$$

both being anti-symmetrized.

With these definitions we write out the energy as

$$\langle c'|H|c'\rangle = \left(1 + \sum_{ai} |\delta C_{ai}|^2\right) \langle c|H|c\rangle + \sum_{ai} |\delta C_{ai}|^2 (\epsilon_a^{HF} - \epsilon_i^{HF}) + \sum_{ijab} A_{ai,bj} \delta C_{ai}^* \delta C_{bj} + \quad (4.17)$$

$$\frac{1}{2} \sum_{ijab} B_{ai,bj}^* \delta C_{ai} \delta C_{bj} + \frac{1}{2} \sum_{ijab} B_{ai,bj} \delta C_{ai}^* \delta C_{bj}^* + O(\delta C_{ai}^3), \quad (4.18)$$

which can be rewritten as

$$\langle c'|H|c'\rangle = \left(1 + \sum_{ai} |\delta C_{ai}|^2\right) \langle c|H|c\rangle + \Delta E + O(\delta C_{ai}^3),$$

and skipping higher-order terms we arrived

$$\frac{\langle c'|\hat{H}|c'\rangle}{\langle c'|c'\rangle} = E_0 + \frac{\Delta E}{(1 + \sum_{ai} |\delta C_{ai}|^2)}.$$

We have defined

$$\Delta E = \frac{1}{2} \langle \chi | \hat{M} | \chi \rangle$$

with the vectors

$$\chi = [\delta C \quad \delta C^*]^T$$

and the matrix

$$\hat{M} = \begin{pmatrix} \Delta + A & B \\ B^* & \Delta + A^* \end{pmatrix},$$

with $\Delta_{ai,bj} = (\epsilon_a - \epsilon_i) \delta_{ab} \delta_{ij}$.

The condition

$$\Delta E = \frac{1}{2} \langle \chi | \hat{M} | \chi \rangle \geq 0$$

for an arbitrary vector

$$\chi = [\delta C \quad \delta C^*]^T$$

means that all eigenvalues of the matrix have to be larger than or equal zero. A necessary (but no sufficient) condition is that the matrix elements (for all ai)

$$(\epsilon_a - \epsilon_i) \delta_{ab} \delta_{ij} + A_{ai,bj} \geq 0.$$

This equation can be used as a first test of the stability of the Hartree-Fock equation.

Chapter 5

Many-body perturbation theory

We assume here that we are only interested in the ground state of the system and expand the exact wave function in term of a series of Slater determinants

$$|\Psi_0\rangle = |\Phi_0\rangle + \sum_{m=1}^{\infty} C_m |\Phi_m\rangle,$$

where we have assumed that the true ground state is dominated by the solution of the unperturbed problem, that is

$$\hat{H}_0 |\Phi_0\rangle = W_0 |\Phi_0\rangle.$$

The state $|\Psi_0\rangle$ is not normalized, rather we have used an intermediate normalization $\langle \Phi_0 | \Psi_0 \rangle = 1$ since we have $\langle \Phi_0 | \Phi_0 \rangle = 1$.

The Schroedinger equation is

$$\hat{H} |\Psi_0\rangle = E |\Psi_0\rangle,$$

and multiplying the latter from the left with $\langle \Phi_0 |$ gives

$$\langle \Phi_0 | \hat{H} | \Psi_0 \rangle = E \langle \Phi_0 | \Psi_0 \rangle = E,$$

and subtracting from this equation

$$\langle \Psi_0 | \hat{H}_0 | \Phi_0 \rangle = W_0 \langle \Psi_0 | \Phi_0 \rangle = W_0,$$

and using the fact that the both operators \hat{H} and \hat{H}_0 are hermitian results in

$$\Delta E = E - W_0 = \langle \Phi_0 | \hat{H}_I | \Psi_0 \rangle,$$

which is an exact result. We call this quantity the correlation energy.

This equation forms the starting point for all perturbative derivations. However, as it stands it represents nothing but a mere formal rewriting of Schroedinger's equation and is not of much practical use. The exact wave function $|\Psi_0\rangle$ is unknown. In order to obtain a perturbative expansion, we need to expand the exact wave function in terms of the interaction \hat{H}_I .

Here we have assumed that our model space defined by the operator \hat{P} is one-dimensional, meaning that

$$\hat{P} = |\Phi_0\rangle \langle \Phi_0|,$$

and

$$\hat{Q} = \sum_{m=1}^{\infty} |\Phi_m\rangle \langle \Phi_m|.$$

We can thus rewrite the exact wave function as

$$|\Psi_0\rangle = (\hat{P} + \hat{Q})|\Psi_0\rangle = |\Phi_0\rangle + \hat{Q}|\Psi_0\rangle.$$

Going back to the Schrödinger equation, we can rewrite it as, adding and subtracting a term $\omega|\Psi_0\rangle$ as

$$(\omega - \hat{H}_0)|\Psi_0\rangle = (\omega - E + \hat{H}_I)|\Psi_0\rangle,$$

where ω is an energy variable to be specified later.

We assume also that the resolvent of $(\omega - \hat{H}_0)$ exists, that is it has an inverse which defined the unperturbed Green's function as

$$(\omega - \hat{H}_0)^{-1} = \frac{1}{(\omega - \hat{H}_0)}.$$

We can rewrite Schroedinger's equation as

$$|\Psi_0\rangle = \frac{1}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) |\Psi_0\rangle,$$

and multiplying from the left with \hat{Q} results in

$$\hat{Q}|\Psi_0\rangle = \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) |\Psi_0\rangle,$$

which is possible since we have defined the operator \hat{Q} in terms of the eigenfunctions of \hat{H} .

These operators commute meaning that

$$\hat{Q} \frac{1}{(\omega - \hat{H}_0)} \hat{Q} = \hat{Q} \frac{1}{(\omega - \hat{H}_0)} = \frac{\hat{Q}}{(\omega - \hat{H}_0)}.$$

With these definitions we can in turn define the wave function as

$$|\Psi_0\rangle = |\Phi_0\rangle + \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) |\Psi_0\rangle.$$

This equation is again nothing but a formal rewrite of Schrödinger's equation and does not represent a practical calculational scheme. It is a non-linear equation in two unknown quantities, the energy E and the exact wave function $|\Psi_0\rangle$. We can however start with a guess for $|\Psi_0\rangle$ on the right hand side of the last equation.

The most common choice is to start with the function which is expected to exhibit the largest overlap with the wave function we are searching after, namely $|\Phi_0\rangle$. This can again be inserted in the solution for $|\Psi_0\rangle$ in an iterative fashion and if we continue along these lines we end up with

$$|\Psi_0\rangle = \sum_{i=0}^{\infty} \left\{ \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) \right\}^i |\Phi_0\rangle,$$

for the wave function and

$$\Delta E = \sum_{i=0}^{\infty} \langle \Phi_0 | \hat{H}_I \left\{ \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) \right\}^i | \Phi_0 \rangle,$$

which is now a perturbative expansion of the exact energy in terms of the interaction \hat{H}_I and the unperturbed wave function $|\Phi_0\rangle$.

In our equations for $|\Psi_0\rangle$ and ΔE in terms of the unperturbed solutions $|\Phi_i\rangle$ we have still an undetermined parameter ω and a dependency on the exact energy E . Not much has been gained thus from a practical computational point of view.

In Brilluoin-Wigner perturbation theory it is customary to set $\omega = E$. This results in the following perturbative expansion for the energy ΔE

$$\begin{aligned}\Delta E &= \sum_{i=0}^{\infty} \langle \Phi_0 | \hat{H}_I \left\{ \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) \right\}^i | \Phi_0 \rangle = \\ &\langle \Phi_0 | \left(\hat{H}_I + \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I + \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I + \dots \right) | \Phi_0 \rangle. \\ \Delta E &= \sum_{i=0}^{\infty} \langle \Phi_0 | \hat{H}_I \left\{ \frac{\hat{Q}}{\omega - \hat{H}_0} (\omega - E + \hat{H}_I) \right\}^i | \Phi_0 \rangle = \\ &\langle \Phi_0 | \left(\hat{H}_I + \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I + \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I \frac{\hat{Q}}{E - \hat{H}_0} \hat{H}_I + \dots \right) | \Phi_0 \rangle.\end{aligned}$$

This expression depends however on the exact energy E and is again not very convenient from a practical point of view. It can obviously be solved iteratively, by starting with a guess for E and then solve till some kind of self-consistency criterion has been reached.

Actually, the above expression is nothing but a rewrite again of the full Schrödinger equation.

Defining $e = E - \hat{H}_0$ and recalling that \hat{H}_0 commutes with \hat{Q} by construction and that \hat{Q} is an idempotent operator $\hat{Q}^2 = \hat{Q}$. Using this equation in the above expansion for ΔE we can write the denominator

$$\begin{aligned}\frac{\hat{Q}}{e - \hat{Q}\hat{H}_I\hat{Q}} &= \\ \hat{Q} \left[\frac{1}{e} + \frac{1}{e} \hat{Q}\hat{H}_I\hat{Q} \frac{1}{e} + \frac{1}{e} \hat{Q}\hat{H}_I\hat{Q} \frac{1}{e} \hat{Q}\hat{H}_I\hat{Q} \frac{1}{e} + \dots \right] \hat{Q}.\end{aligned}$$

Inserted in the expression for ΔE leads to

$$\Delta E = \langle \Phi_0 | \hat{H}_I + \hat{H}_I \hat{Q} \frac{1}{E - \hat{H}_0 - \hat{Q}\hat{H}_I\hat{Q}} \hat{Q}\hat{H}_I | \Phi_0 \rangle.$$

In RS perturbation theory we set $\omega = W_0$ and obtain the following expression for the energy difference

$$\begin{aligned}\Delta E &= \sum_{i=0}^{\infty} \langle \Phi_0 | \hat{H}_I \left\{ \frac{\hat{Q}}{W_0 - \hat{H}_0} (\hat{H}_I - \Delta E) \right\}^i | \Phi_0 \rangle = \\ &\langle \Phi_0 | \left(\hat{H}_I + \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} (\hat{H}_I - \Delta E) + \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} (\hat{H}_I - \Delta E) \frac{\hat{Q}}{W_0 - \hat{H}_0} (\hat{H}_I - \Delta E) + \dots \right) | \Phi_0 \rangle.\end{aligned}$$

Recalling that \hat{Q} commutes with \hat{H}_0 and since ΔE is a constant we obtain that

$$\hat{Q}\Delta E|\Phi_0\rangle = \hat{Q}\Delta E|\hat{Q}\Phi_0\rangle = 0.$$

Inserting this results in the expression for the energy results in

$$\Delta E = \langle \Phi_0 | \left(\hat{H}_I + \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I + \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} (\hat{H}_I - \Delta E) \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I + \dots \right) | \Phi_0 \rangle.$$

We can now this expression in terms of a perturbative expression in terms of \hat{H}_I where we iterate the last expression in terms of ΔE

$$\Delta E = \sum_{i=1}^{\infty} \Delta E^{(i)}.$$

We get the following expression for $\Delta E^{(i)}$

$$\Delta E^{(1)} = \langle \Phi_0 | \hat{H}_I | \Phi_0 \rangle,$$

which is just the contribution to first order in perturbation theory,

$$\Delta E^{(2)} = \langle \Phi_0 | \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I | \Phi_0 \rangle,$$

which is the contribution to second order.

$$\Delta E^{(3)} = \langle \Phi_0 | \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I | \Phi_0 \rangle - \langle \Phi_0 | \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \langle \Phi_0 | \hat{H}_I | \Phi_0 \rangle \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I | \Phi_0 \rangle,$$

being the third-order contribution.

Interpreting the correlation energy and the wave operator

In the shell-model lectures we showed that we could rewrite the exact state function for say the ground state, as a linear expansion in terms of all possible Slater determinants. That is, we define the ansatz for the ground state as

$$|\Phi_0\rangle = \left(\prod_{i \leq F} \hat{a}_i^\dagger \right) |0\rangle,$$

where the index i defines different single-particle states up to the Fermi level. We have assumed that we have N fermions. A given one-particle-one-hole ($1p1h$) state can be written as

$$|\Phi_i^a\rangle = \hat{a}_a^\dagger \hat{a}_i |\Phi_0\rangle,$$

while a $2p2h$ state can be written as

$$|\Phi_{ij}^{ab}\rangle = \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_j \hat{a}_i |\Phi_0\rangle,$$

and a general $ApAh$ state as

$$|\Phi_{ijk\dots}^{abc\dots}\rangle = \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_c^\dagger \dots \hat{a}_k \hat{a}_j \hat{a}_i |\Phi_0\rangle.$$

We use letters $ijkl\dots$ for states below the Fermi level and $abcd\dots$ for states above the Fermi level. A general single-particle state is given by letters $pqrs\dots$.

We can then expand our exact state function for the ground state as

$$|\Psi_0\rangle = C_0 |\Phi_0\rangle + \sum_{ai} C_i^a |\Phi_i^a\rangle + \sum_{abij} C_{ij}^{ab} |\Phi_{ij}^{ab}\rangle + \dots = (C_0 + \hat{C}) |\Phi_0\rangle,$$

where we have introduced the so-called correlation operator

$$\hat{C} = \sum_{ai} C_i^a \hat{a}_a^\dagger \hat{a}_i + \sum_{abij} C_{ij}^{ab} \hat{a}_a^\dagger \hat{a}_b^\dagger \hat{a}_j \hat{a}_i + \dots$$

Since the normalization of Ψ_0 is at our disposal and since C_0 is by hypothesis non-zero, we may arbitrarily set $C_0 = 1$ with corresponding proportional changes in all other coefficients. Using this so-called intermediate normalization we have

$$\langle \Psi_0 | \Phi_0 \rangle = \langle \Phi_0 | \Phi_0 \rangle = 1,$$

resulting in

$$|\Psi_0\rangle = (1 + \hat{C})|\Phi_0\rangle.$$

In a shell-model calculation, the unknown coefficients in \hat{C} are the eigenvectors which result from the diagonalization of the Hamiltonian matrix.

How can we use perturbation theory to determine the same coefficients? Let us study the contributions to second order in the interaction, namely

$$\Delta E^{(2)} = \langle \Phi_0 | \hat{H}_I \frac{\hat{Q}}{W_0 - \hat{H}_0} \hat{H}_I | \Phi_0 \rangle.$$

The intermediate states given by \hat{Q} can at most be of a $2p - 2h$ nature if we have a two-body Hamiltonian. This means that second order in the perturbation theory can have $1p - 1h$ and $2p - 2h$ at most as intermediate states. When we diagonalize, these contributions are included to infinite order. This means that higher-orders in perturbation theory bring in more complicated correlations.

If we limit the attention to a Hartree-Fock basis, then we have that $\langle \Phi_0 | \hat{H}_I | 2p - 2h \rangle$ is the only contribution and the contribution to the energy reduces to

$$\Delta E^{(2)} = \frac{1}{4} \sum_{abij} \langle ij | \hat{v} | ab \rangle \frac{\langle ab | \hat{v} | ij \rangle}{\epsilon_i + \epsilon_j - \epsilon_a - \epsilon_b}.$$

If we compare this to the correlation energy obtained from full configuration interaction theory with a Hartree-Fock basis, we found that

$$E - E_0 = \Delta E = \sum_{abij} \langle ij | \hat{v} | ab \rangle C_{ij}^{ab},$$

where the energy E_0 is the reference energy and ΔE defines the so-called correlation energy.

We see that if we set

$$C_{ij}^{ab} = \frac{1}{4} \frac{\langle ab | \hat{v} | ij \rangle}{\epsilon_i + \epsilon_j - \epsilon_a - \epsilon_b},$$

we have a perfect agreement between FCI and MBPT. However, FCI includes such $2p - 2h$ correlations to infinite order. In order to make a meaningful comparison we would at least need to sum such correlations to infinite order in perturbation theory.

Summing up, we can see that

- MBPT introduces order-by-order specific correlations and we make comparisons with exact calculations like FCI
- At every order, we can calculate all contributions since they are well-known and either tabulated or calculated on the fly.
- MBPT is a non-variational theory and there is no guarantee that higher orders will improve the convergence.
- However, since FCI calculations are limited by the size of the Hamiltonian matrices to diagonalize (today's most efficient codes can attach dimensionalities of ten billion basis states, MBPT can function as an approximative method which gives a straightforward (but tedious) calculation recipe.
- MBPT has been widely used to compute effective interactions for the nuclear shell-model.
- But there are better methods which sum to infinite order important correlations. Coupled cluster theory is one of these methods.

Part III
Monte Carlo Methods

Chapter 6

Variational Monte Carlo methods

Quantum Monte Carlo Motivation

We start with the variational principle. Given a hamiltonian H and a trial wave function Ψ_T , the variational principle states that the expectation value of $\langle H \rangle$, defined through

$$E[H] = \langle H \rangle = \frac{\int dR \Psi_T^*(R) H(R) \Psi_T(R)}{\int dR \Psi_T^*(R) \Psi_T(R)},$$

is an upper bound to the ground state energy E_0 of the hamiltonian H , that is

$$E_0 \leq \langle H \rangle.$$

In general, the integrals involved in the calculation of various expectation values are multi-dimensional ones. Traditional integration methods such as the Gauss-Legendre will not be adequate for say the computation of the energy of a many-body system.

The trial wave function can be expanded in the eigenstates of the hamiltonian since they form a complete set, viz.,

$$\Psi_T(R) = \sum_i a_i \Psi_i(R),$$

and assuming the set of eigenfunctions to be normalized one obtains

$$\frac{\sum_{nm} a_m^* a_n \int dR \Psi_m^*(R) H(R) \Psi_n(R)}{\sum_{nm} a_m^* a_n \int dR \Psi_m^*(R) \Psi_n(R)} = \frac{\sum_n a_n^2 E_n}{\sum_n a_n^2} \geq E_0,$$

where we used that $H(R) \Psi_n(R) = E_n \Psi_n(R)$. In general, the integrals involved in the calculation of various expectation values are multi-dimensional ones. The variational principle yields the lowest state of a given symmetry.

In most cases, a wave function has only small values in large parts of configuration space, and a straightforward procedure which uses homogenously distributed random points in configuration space will most likely lead to poor results. This may suggest that some kind of importance sampling combined with e.g., the Metropolis algorithm may be a more efficient way of obtaining the ground state energy. The hope is then that those regions of configurations space where the wave function assumes appreciable values are sampled more efficiently.

The tedious part in a VMC calculation is the search for the variational minimum. A good knowledge of the system is required in order to carry out reasonable VMC calculations. This is not always the case, and often VMC calculations serve rather as the starting point for so-called diffusion Monte Carlo calculations (DMC). DMC is a way of solving exactly the many-body Schroedinger equation by means of a stochastic procedure. A good guess on the binding

energy and its wave function is however necessary. A carefully performed VMC calculation can aid in this context.

The basic recipe in a VMC calculation consists of the following elements:

- Construct first a trial wave function $\psi_T(R, \alpha)$, for a many-body system consisting of N particles located at positions $R = (R_1, \dots, R_N)$. The trial wave function depends on α variational parameters $\alpha = (\alpha_1, \dots, \alpha_M)$.
- Then we evaluate the expectation value of the hamiltonian H

$$E[H] = \langle H \rangle = \frac{\int dR \Psi_T^*(R, \alpha) H(R) \Psi_T(R, \alpha)}{\int dR \Psi_T^*(R, \alpha) \Psi_T(R, \alpha)}.$$

- Thereafter we vary α according to some minimization algorithm and return to the first step.

With a trial wave function $\psi_T(R)$ we can in turn construct the quantum mechanical probability distribution

$$P(R) = \frac{|\psi_T(R)|^2}{\int |\psi_T(R)|^2 dR}.$$

This is our new probability distribution function (PDF). The approximation to the expectation value of the Hamiltonian is now

$$E[H(\alpha)] = \frac{\int dR \Psi_T^*(R, \alpha) H(R) \Psi_T(R, \alpha)}{\int dR \Psi_T^*(R, \alpha) \Psi_T(R, \alpha)}.$$

Define a new quantity

$$E_L(R, \alpha) = \frac{1}{\psi_T(R, \alpha)} H \psi_T(R, \alpha),$$

called the local energy, which, together with our trial PDF yields

$$E[H(\alpha)] = \int P(R) E_L(R) dR \approx \frac{1}{N} \sum_{i=1}^N P(R_i, \alpha) E_L(R_i, \alpha)$$

with N being the number of Monte Carlo samples.

The Algorithm for performing a variational Monte Carlo calculations runs thus as this

- Initialisation: Fix the number of Monte Carlo steps. Choose an initial R and variational parameters α and calculate $|\psi_T^\alpha(R)|^2$.
- Initialise the energy and the variance and start the Monte Carlo calculation.
 - Calculate a trial position $R_p = R + r * \text{step}$ where r is a random variable $r \in [0, 1]$.
 - Metropolis algorithm to accept or reject this move $w = P(R_p)/P(R)$.
 - If the step is accepted, then we set $R = R_p$.
 - Update averages
- Finish and compute final averages.

Observe that the jumping in space is governed by the variable *step*. This is Called brute-force sampling. Need importance sampling to get more relevant sampling, see lectures below.

Quantum Monte Carlo: hydrogen atom.

The radial Schroedinger equation for the hydrogen atom can be written as

$$-\frac{\hbar^2}{2m} \frac{\partial^2 u(r)}{\partial r^2} - \left(\frac{ke^2}{r} - \frac{l(l+1)\hbar^2}{2mr^2} \right) u(r) = E u(r),$$

or with dimensionless variables

$$-\frac{1}{2} \frac{\partial^2 u(\rho)}{\partial \rho^2} - \frac{u(\rho)}{\rho} + \frac{l(l+1)}{2\rho^2} u(\rho) - \lambda u(\rho) = 0,$$

with the hamiltonian

$$H = -\frac{1}{2} \frac{\partial^2}{\partial \rho^2} - \frac{1}{\rho} + \frac{l(l+1)}{2\rho^2}.$$

Use variational parameter α in the trial wave function

$$u_T^\alpha(\rho) = \alpha \rho e^{-\alpha \rho}.$$

Inserting this wave function into the expression for the local energy E_L gives

$$E_L(\rho) = -\frac{1}{\rho} - \frac{\alpha}{2} \left(\alpha - \frac{2}{\rho} \right).$$

A simple variational Monte Carlo calculation results in

α	$\langle H \rangle$	σ^2	σ/\sqrt{N}
7.00000E-01	-4.57759E-01	4.51201E-02	6.71715E-04
8.00000E-01	-4.81461E-01	3.05736E-02	5.52934E-04
9.00000E-01	-4.95899E-01	8.20497E-03	2.86443E-04
1.00000E-00	-5.00000E-01	0.00000E+00	0.00000E+00
1.10000E+00	-4.93738E-01	1.16989E-02	3.42036E-04
1.20000E+00	-4.75563E-01	8.85899E-02	9.41222E-04
1.30000E+00	-4.54341E-01	1.45171E-01	1.20487E-03

We note that at $\alpha = 1$ we obtain the exact result, and the variance is zero, as it should. The reason is that we then have the exact wave function, and the action of the hamiltonian on the wave function

$$H\psi = \text{constant} \times \psi,$$

yields just a constant. The integral which defines various expectation values involving moments of the hamiltonian becomes then

$$\langle H^n \rangle = \frac{\int dR \Psi_T^*(R) H^n(R) \Psi_T(R)}{\int dR \Psi_T^*(R) \Psi_T(R)} = \text{constant} \times \frac{\int dR \Psi_T^*(R) \Psi_T(R)}{\int dR \Psi_T^*(R) \Psi_T(R)} = \text{constant}.$$

This gives an important information: the exact wave function leads to zero variance!

Variation is then performed by minimizing both the energy and the variance.

For bosons in a harmonic oscillator-like trap we will use is a spherical (S) or an elliptical (E) harmonic trap in one, two and finally three dimensions, with the latter given by

$$V_{ext}(\mathbf{r}) = \begin{cases} \frac{1}{2} m \omega_{ho}^2 r^2 & (S) \\ \frac{1}{2} m [\omega_{ho}^2 (x^2 + y^2) + \omega_z^2 z^2] & (E) \end{cases} \quad (6.1)$$

where (S) stands for symmetric and

$$\hat{H} = \sum_i^N \left(\frac{-\hbar^2}{2m} \nabla_i^2 + V_{ext}(\mathbf{r}_i) \right) + \sum_{i < j}^N V_{int}(\mathbf{r}_i, \mathbf{r}_j), \quad (6.2)$$

as the two-body Hamiltonian of the system.

We will represent the inter-boson interaction by a pairwise, repulsive potential

$$V_{int}(|\mathbf{r}_i - \mathbf{r}_j|) = \begin{cases} \infty & |\mathbf{r}_i - \mathbf{r}_j| \leq a \\ 0 & |\mathbf{r}_i - \mathbf{r}_j| > a \end{cases} \quad (6.3)$$

where a is the so-called hard-core diameter of the bosons. Clearly, $V_{int}(|\mathbf{r}_i - \mathbf{r}_j|)$ is zero if the bosons are separated by a distance $|\mathbf{r}_i - \mathbf{r}_j|$ greater than a but infinite if they attempt to come within a distance $|\mathbf{r}_i - \mathbf{r}_j| \leq a$.

Our trial wave function for the ground state with N atoms is given by

$$\Psi_T(\mathbf{R}) = \Psi_T(\mathbf{r}_1, \mathbf{r}_2, \dots, \mathbf{r}_N, \alpha, \beta) = \prod_i g(\alpha, \beta, \mathbf{r}_i) \prod_{i < j} f(a, |\mathbf{r}_i - \mathbf{r}_j|), \quad (6.4)$$

where α and β are variational parameters. The single-particle wave function is proportional to the harmonic oscillator function for the ground state

$$g(\alpha, \beta, \mathbf{r}_i) = \exp[-\alpha(x_i^2 + y_i^2 + \beta z_i^2)]. \quad (6.5)$$

For spherical traps we have $\beta = 1$ and for non-interacting bosons ($a = 0$) we have $\alpha = 1/2a_{ho}^2$. The correlation wave function is

$$f(a, |\mathbf{r}_i - \mathbf{r}_j|) = \begin{cases} 0 & |\mathbf{r}_i - \mathbf{r}_j| \leq a \\ (1 - \frac{a}{|\mathbf{r}_i - \mathbf{r}_j|}) & |\mathbf{r}_i - \mathbf{r}_j| > a. \end{cases} \quad (6.6)$$

A simple Python code that solves the two-boson or two-fermion case in two-dimensions.

[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
Importing various packages from math import exp, sqrt from random import random, seed im-
port numpy as np import matplotlib.pyplot as plt from mpl_toolkits.mplot3d import Axes3D from matplotlib import cm from matplotlib

Trial wave function for quantum dots in two dims def WaveFunction(r,alpha,beta): r1 = r[0,0]**2 + r[0,1]**2 r2 = r[1,0]**2 + r[1,1]**2 r12 = sqrt((r[0,0]-r[1,0])**2 + (r[0,1]-r[1,1])**2) deno = r12/(1+beta*r12) return exp(-0.5*alpha*(r1+r2)+deno)

Local energy for quantum dots in two dims, using analytical local energy def LocalEnergy(r,alpha,beta):

r1 = (r[0,0]**2 + r[0,1]**2) r2 = (r[1,0]**2 + r[1,1]**2) r12 = sqrt((r[0,0]-r[1,0])**2 + (r[0,1]-r[1,1])**2) deno = 1.0/(1+beta*r12) deno2 = deno*deno return 0.5*(1-alpha*alpha)*(r1 + r2) + 2.0*alpha + 1.0/r12+deno2*(alpha*r12-deno2+2*beta*deno-1.0/r12)

The Monte Carlo sampling with the Metropolis algo def MonteCarloSampling():

NumberMCCycles= 100000 StepSize = 1.0 positions PositionOld = np.zeros((NumberParticles,Dimension), np.double) PositionNew = np.zeros((NumberParticles,Dimension), np.double) seed for rng generator seed() start variational parameter alpha = 0.9 for ia in range(MaxVariations): alpha += .025 AlphaValues[ia] = alpha beta = 0.2 for jb in range(MaxVariations): beta += .01 BetaValues[jb] = beta energy = energy2 = 0.0 DeltaE = 0.0 Initial position for i in range(NumberParticles): for j in range(Dimension): PositionOld[i,j] = StepSize * (random() - .5) wfold = WaveFunction(PositionOld,alpha,beta)

Loop over MC MCCycles for MCCycle in range(NumberMCCycles): Trial position for i in range(NumberParticles): for j in range(Dimension): PositionNew[i,j] = PositionOld[i,j] + StepSize * (random() - .5) wfnew = WaveFunction(PositionNew,alpha,beta)

Metropolis test to see whether we accept the move if random() < wfnew**2 / wfold**2: PositionOld = PositionNew.copy() wfold = wfnew DeltaE = LocalEnergy(PositionOld,alpha,beta) energy += DeltaE energy2 += DeltaE**2

We calculate mean, variance and error ... energy /= NumberMCCycles energy2 /= NumberMCCycles variance = energy2 - energy**2 error = sqrt(variance/NumberMCCycles) Energies[ia,jb] = energy return Energies, AlphaValues, BetaValues

Here starts the main program with variable declarations NumberParticles = 2 Dimension = 2 MaxVariations = 10 Energies = np.zeros((MaxVariations,MaxVariations)) AlphaValues =

```
np.zeros(MaxVariations) BetaValues = np.zeros(MaxVariations) (Energies, AlphaValues, BetaValues) = MonteCarloSampling()
```

```
Prepare for plots fig = plt.figure() ax = fig.gca(projection='3d') Plot the surface. X, Y = np.meshgrid(AlphaValues, BetaValues) surf = ax.plot_surface(X, Y, Energies, cmap = cm.coolwarm, linewidth = 0, antialiased = False) Customizethezaxis.zmin = np.matrix(Energies).min() zmax = np.matrix(Energies).max() ax.set_zlim(zmin, zmax) ax.set_xlabel(r'\beta') ax.set_ylabel(r'\alpha') ax.set_zlabel(r'\langle E \rangle') ax.xaxis.set_major_locator(LinearLocator(10)) ax.xaxis.set_major_formatter(FormatStrFormatter('%A 0.5, aspect = 5)) plt.show()
```

Quantum Monte Carlo: the helium atom

The helium atom consists of two electrons and a nucleus with charge $Z = 2$. The contribution to the potential energy due to the attraction from the nucleus is

$$-\frac{2ke^2}{r_1} - \frac{2ke^2}{r_2},$$

and if we add the repulsion arising from the two interacting electrons, we obtain the potential energy

$$V(r_1, r_2) = -\frac{2ke^2}{r_1} - \frac{2ke^2}{r_2} + \frac{ke^2}{r_{12}},$$

with the electrons separated at a distance $r_{12} = |r_1 - r_2|$.

The hamiltonian becomes then

$$\hat{H} = -\frac{\nabla_1^2}{2m} - \frac{\nabla_2^2}{2m} - \frac{2ke^2}{r_1} - \frac{2ke^2}{r_2} + \frac{ke^2}{r_{12}},$$

and Schroedingers equation reads

$$\hat{H}\psi = E\psi.$$

All observables are evaluated with respect to the probability distribution

$$P(R) = \frac{|\psi_T(R)|^2}{\int |\psi_T(R)|^2 dR}.$$

generated by the trial wave function. The trial wave function must approximate an exact eigenstate in order that accurate results are to be obtained.

Choice of trial wave function for Helium: Assume $r_1 \rightarrow 0$.

$$E_L(R) = \frac{1}{\psi_T(R)} H \psi_T(R) = \frac{1}{\psi_T(R)} \left(-\frac{1}{2} \nabla_1^2 - \frac{Z}{r_1} \right) \psi_T(R) + \text{finite terms}.$$

$$E_L(R) = \frac{1}{\mathbf{R}_T(r_1)} \left(-\frac{1}{2} \frac{d^2}{dr_1^2} - \frac{1}{r_1} \frac{d}{dr_1} - \frac{Z}{r_1} \right) \mathbf{R}_T(r_1) + \text{finite terms}$$

For small values of r_1 , the terms which dominate are

$$\lim_{r_1 \rightarrow 0} E_L(R) = \frac{1}{\mathbf{R}_T(r_1)} \left(-\frac{1}{r_1} \frac{d}{dr_1} - \frac{Z}{r_1} \right) \mathbf{R}_T(r_1),$$

since the second derivative does not diverge due to the finiteness of Ψ at the origin.

This results in

$$\frac{1}{\mathbf{R}_T(r_1)} \frac{d\mathbf{R}_T(r_1)}{dr_1} = -Z,$$

and

$$\mathbf{R}_T(r_1) \propto e^{-Zr_1}.$$

A similar condition applies to electron 2 as well. For orbital momenta $l > 0$ we have

$$\frac{1}{\mathbf{R}_T(r)} \frac{d\mathbf{R}_T(r)}{dr} = -\frac{Z}{l+1}.$$

Similarly, studying the case $r_{12} \rightarrow 0$ we can write a possible trial wave function as

$$\psi_T(R) = e^{-\alpha(r_1+r_2)} e^{\beta r_{12}}.$$

The last equation can be generalized to

$$\psi_T(R) = \phi(r_1)\phi(r_2)\dots\phi(r_N) \prod_{i<j} f(r_{ij}),$$

for a system with N electrons or particles.

During the development of our code we need to make several checks. It is also very instructive to compute a closed form expression for the local energy. Since our wave function is rather simple it is straightforward to find an analytic expressions. Consider first the case of the simple helium function

$$\Psi_T(r_1, r_2) = e^{-\alpha(r_1+r_2)}$$

The local energy is for this case

$$E_{L1} = (\alpha - Z) \left(\frac{1}{r_1} + \frac{1}{r_2} \right) + \frac{1}{r_{12}} - \alpha^2$$

which gives an expectation value for the local energy given by

$$\langle E_{L1} \rangle = \alpha^2 - 2\alpha \left(Z - \frac{5}{16} \right)$$

With closed form formulae we can speed up the computation of the correlation. In our case we write it as

$$\Psi_C = \exp \left\{ \sum_{i<j} \frac{ar_{ij}}{1 + \beta r_{ij}} \right\},$$

which means that the gradient needed for the so-called quantum force and local energy can be calculated analytically. This will speed up your code since the computation of the correlation part and the Slater determinant are the most time consuming parts in your code.

We will refer to this correlation function as Ψ_C or the *linear Pade-Jastrow*.

We can test this by computing the local energy for our helium wave function

$$\psi_T(r_1, r_2) = \exp(-\alpha(r_1 + r_2)) \exp\left(\frac{r_{12}}{2(1 + \beta r_{12})}\right),$$

with α and β as variational parameters.

The local energy is for this case

$$E_{L2} = E_{L1} + \frac{1}{2(1 + \beta r_{12})^2} \left\{ \frac{\alpha(r_1 + r_2)}{r_{12}} \left(1 - \frac{r_1 r_2}{r_{12}} \right) - \frac{1}{2(1 + \beta r_{12})^2} - \frac{2}{r_{12}} + \frac{2\beta}{1 + \beta r_{12}} \right\}$$

It is very useful to test your code against these expressions. It means also that you don't need to compute a derivative numerically as discussed in the code example below.

For the computation of various derivatives with different types of wave functions, you will find it useful to use python with symbolic python, that is sympy, see online manual. Using sympy allows you autogenerate both Latex code as well c++, python or Fortran codes. Here you will find some simple examples. We choose the $2s$ hydrogen-orbital (not normalized) as an example

$$\phi_{2s}(r) = (Zr - 2) \exp\left(-\frac{1}{2}Zr\right),$$

with $r^2 = x^2 + y^2 + z^2$.

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
from sympy import symbols, diff, exp, sqrt
x, y, z, Z = symbols('x y z Z')
r = sqrt(x*x + y*y + z*z)
phi = (Z*r - 2)*exp(-Z*r/2)
phi.diff(phi, x)
```

This doesn't look very nice, but sympy provides several functions that allow for improving and simplifying the output.

We can improve our output by factorizing and substituting expressions

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
from sympy import symbols, diff, exp, sqrt, factor, Symbol, printing
x, y, z, Z = symbols('x y z Z')
r = sqrt(x*x + y*y + z*z)
phi = (Z*r - 2)*exp(-Z*r/2)
R = Symbol('r')
Creates a symbolic equivalent of r
print latex and c++ code
print printing.latex(diff(phi, x).factor().subs(r, R))
print printing.ccode(diff(phi, x).factor().subs(r, R))
```

We can in turn look at second derivatives

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
from sympy import symbols, diff, exp, sqrt, factor, Symbol, printing
x, y, z, Z = symbols('x y z Z')
r = sqrt(x*x + y*y + z*z)
phi = (Z*r - 2)*exp(-Z*r/2)
R = Symbol('r')
Creates a symbolic equivalent of r
(diff(diff(phi, x), x) + diff(diff(phi, y), y) + diff(diff(phi, z), z)).factor().subs(r, R)
Collect the Z values
(diff(diff(phi, x), x) + diff(diff(phi, y), y) + diff(diff(phi, z), z)).factor().collect(Z).subs(r, R)
Factorize also the r**2 terms
(diff(diff(phi, x), x) + diff(diff(phi, y), y) + diff(diff(phi, z), z)).factor().collect(Z).subs(r, R).subs(r**2, R**2).factor()
print printing.ccode((diff(diff(phi, x), x) + diff(diff(phi, y), y) + diff(diff(phi, z), z)).factor().collect(Z).subs(r, R).subs(r**2, R**2).factor())
```

With some practice this allows one to be able to check one's own calculation and translate automatically into code lines.

The Metropolis algorithm

The Metropolis algorithm, see the original article was invented by Metropolis et. al and is often simply called the Metropolis algorithm. It is a method to sample a normalized probability distribution by a stochastic process. We define $\mathbf{P}_i^{(n)}$ to be the probability for finding the system in the state i at step n . The algorithm is then

- Sample a possible new state j with some probability $T_{i \rightarrow j}$.
- Accept the new state j with probability $A_{i \rightarrow j}$ and use it as the next sample. With probability $1 - A_{i \rightarrow j}$ the move is rejected and the original state i is used again as a sample.

We wish to derive the required properties of T and A such that $\mathbf{P}_i^{(n \rightarrow \infty)} \rightarrow p_i$ so that starting from any distribution, the method converges to the correct distribution. Note that the description here is for a discrete probability distribution. Replacing probabilities p_i with expressions like $p(x_i)dx_i$ will take all of these over to the corresponding continuum expressions.

The dynamical equation for $\mathbf{P}_i^{(n)}$ can be written directly from the description above. The probability of being in the state i at step n is given by the probability of being in any state j

at the previous step, and making an accepted transition to i added to the probability of being in the state i , making a transition to any state j and rejecting the move:

$$\mathbf{P}_i^{(n)} = \sum_j \left[\mathbf{P}_j^{(n-1)} T_{j \rightarrow i} A_{j \rightarrow i} + \mathbf{P}_i^{(n-1)} T_{i \rightarrow j} (1 - A_{i \rightarrow j}) \right].$$

Since the probability of making some transition must be 1, $\sum_j T_{i \rightarrow j} = 1$, and the above equation becomes

$$\mathbf{P}_i^{(n)} = \mathbf{P}_i^{(n-1)} + \sum_j \left[\mathbf{P}_j^{(n-1)} T_{j \rightarrow i} A_{j \rightarrow i} - \mathbf{P}_i^{(n-1)} T_{i \rightarrow j} A_{i \rightarrow j} \right].$$

For large n we require that $\mathbf{P}_i^{(n \rightarrow \infty)} = p_i$, the desired probability distribution. Taking this limit, gives the balance requirement

$$\sum_j [p_j T_{j \rightarrow i} A_{j \rightarrow i} - p_i T_{i \rightarrow j} A_{i \rightarrow j}] = 0.$$

The balance requirement is very weak. Typically the much stronger detailed balance requirement is enforced, that is rather than the sum being set to zero, we set each term separately to zero and use this to determine the acceptance probabilities. Rearranging, the result is

$$\frac{A_{j \rightarrow i}}{A_{i \rightarrow j}} = \frac{p_i T_{i \rightarrow j}}{p_j T_{j \rightarrow i}}.$$

The Metropolis choice is to maximize the A values, that is

$$A_{j \rightarrow i} = \min \left(1, \frac{p_i T_{i \rightarrow j}}{p_j T_{j \rightarrow i}} \right).$$

Other choices are possible, but they all correspond to multiplying $A_{i \rightarrow j}$ and $A_{j \rightarrow i}$ by the same constant smaller than unity.¹

Having chosen the acceptance probabilities, we have guaranteed that if the $\mathbf{P}_i^{(n)}$ has equilibrated, that is if it is equal to p_i , it will remain equilibrated. Next we need to find the circumstances for convergence to equilibrium.

The dynamical equation can be written as

$$\mathbf{P}_i^{(n)} = \sum_j M_{ij} \mathbf{P}_j^{(n-1)}$$

with the matrix M given by

$$M_{ij} = \delta_{ij} \left[1 - \sum_k T_{i \rightarrow k} A_{i \rightarrow k} \right] + T_{j \rightarrow i} A_{j \rightarrow i}.$$

Summing over i shows that $\sum_i M_{ij} = 1$, and since $\sum_k T_{i \rightarrow k} = 1$, and $A_{i \rightarrow k} \leq 1$, the elements of the matrix satisfy $M_{ij} \geq 0$. The matrix M is therefore a stochastic matrix.

The Metropolis method is simply the power method for computing the right eigenvector of M with the largest magnitude eigenvalue. By construction, the correct probability distribution is a right eigenvector with eigenvalue 1. Therefore, for the Metropolis method to converge to this result, we must show that M has only one eigenvalue with this magnitude, and all other eigenvalues are smaller.

¹ The penalty function method uses just such a factor to compensate for p_i that are evaluated stochastically and are therefore noisy.

Importance sampling

We need to replace the brute force Metropolis algorithm with a walk in coordinate space biased by the trial wave function. This approach is based on the Fokker-Planck equation and the Langevin equation for generating a trajectory in coordinate space. The link between the Fokker-Planck equation and the Langevin equations are explained, only partly, in the slides below. An excellent reference on topics like Brownian motion, Markov chains, the Fokker-Planck equation and the Langevin equation is the text by Van Kampen. Here we will focus first on the implementation part first.

For a diffusion process characterized by a time-dependent probability density $P(x, t)$ in one dimension the Fokker-Planck equation reads (for one particle /walker)

$$\frac{\partial P}{\partial t} = D \frac{\partial}{\partial x} \left(\frac{\partial}{\partial x} - F \right) P(x, t),$$

where F is a drift term and D is the diffusion coefficient.

The new positions in coordinate space are given as the solutions of the Langevin equation using Euler's method, namely, we go from the Langevin equation

$$\frac{\partial x(t)}{\partial t} = DF(x(t)) + \eta,$$

with η a random variable, yielding a new position

$$y = x + DF(x)\Delta t + \xi \sqrt{\Delta t},$$

where ξ is gaussian random variable and Δt is a chosen time step. The quantity D is, in atomic units, equal to $1/2$ and comes from the factor $1/2$ in the kinetic energy operator. Note that Δt is to be viewed as a parameter. Values of $\Delta t \in [0.001, 0.01]$ yield in general rather stable values of the ground state energy.

The process of isotropic diffusion characterized by a time-dependent probability density $P(\mathbf{x}, t)$ obeys (as an approximation) the so-called Fokker-Planck equation

$$\frac{\partial P}{\partial t} = \sum_i D \frac{\partial}{\partial \mathbf{x}_i} \left(\frac{\partial}{\partial \mathbf{x}_i} - \mathbf{F}_i \right) P(\mathbf{x}, t),$$

where \mathbf{F}_i is the i^{th} component of the drift term (drift velocity) caused by an external potential, and D is the diffusion coefficient. The convergence to a stationary probability density can be obtained by setting the left hand side to zero. The resulting equation will be satisfied if and only if all the terms of the sum are equal zero,

$$\frac{\partial^2 P}{\partial \mathbf{x}_i^2} = P \frac{\partial}{\partial \mathbf{x}_i} \mathbf{F}_i + \mathbf{F}_i \frac{\partial}{\partial \mathbf{x}_i} P.$$

The drift vector should be of the form $\mathbf{F} = g(\mathbf{x}) \frac{\partial P}{\partial \mathbf{x}}$. Then,

$$\frac{\partial^2 P}{\partial \mathbf{x}_i^2} = P \frac{\partial g}{\partial P} \left(\frac{\partial P}{\partial \mathbf{x}_i} \right)^2 + P g \frac{\partial^2 P}{\partial \mathbf{x}_i^2} + g \left(\frac{\partial P}{\partial \mathbf{x}_i} \right)^2.$$

The condition of stationary density means that the left hand side equals zero. In other words, the terms containing first and second derivatives have to cancel each other. It is possible only if $g = \frac{1}{P}$, which yields

$$\mathbf{F} = 2 \frac{1}{\Psi_T} \nabla \Psi_T,$$

which is known as the so-called *quantum force*. This term is responsible for pushing the walker towards regions of configuration space where the trial wave function is large, increasing the efficiency of the simulation in contrast to the Metropolis algorithm where the walker has the same probability of moving in every direction.

The Fokker-Planck equation yields a (the solution to the equation) transition probability given by the Green's function

$$G(y, x, \Delta t) = \frac{1}{(4\pi D \Delta t)^{3N/2}} \exp\left(-(y - x - D \Delta t F(x))^2 / 4D \Delta t\right)$$

which in turn means that our brute force Metropolis algorithm

$$A(y, x) = \min(1, q(y, x)),$$

with $q(y, x) = |\Psi_T(y)|^2 / |\Psi_T(x)|^2$ is now replaced by the Metropolis-Hastings algorithm as well as Hasting's article,

$$q(y, x) = \frac{G(x, y, \Delta t) |\Psi_T(y)|^2}{G(y, x, \Delta t) |\Psi_T(x)|^2}$$

Importance sampling, program elements

The general derivative formula of the Jastrow factor is (the subscript C stands for Correlation)

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{\partial g_{ik}}{\partial x_k} + \sum_{i=k+1}^N \frac{\partial g_{ki}}{\partial x_k}$$

However, with our written in way which can be reused later as

$$\Psi_C = \prod_{i < j} g(r_{ij}) = \exp \left\{ \sum_{i < j} f(r_{ij}) \right\},$$

the gradient needed for the quantum force and local energy is easy to compute. The function $f(r_{ij})$ will depends on the system under study. In the equations below we will keep this general form.

In the Metropolis/Hasting algorithm, the *acceptance ratio* determines the probability for a particle to be accepted at a new position. The ratio of the trial wave functions evaluated at the new and current positions is given by (OB for the onebody part)

$$R \equiv \frac{\Psi_T^{new}}{\Psi_T^{old}} = \frac{\Psi_{OB}^{new} \Psi_C^{new}}{\Psi_{OB}^{old} \Psi_C^{old}}$$

Here Ψ_{OB} is our onebody part (Slater determinant or product of boson single-particle states) while Ψ_C is our correlation function, or Jastrow factor. We need to optimize the $\nabla \Psi_T / \Psi_T$ ratio and the second derivative as well, that is the $\nabla^2 \Psi_T / \Psi_T$ ratio. The first is needed when we compute the so-called quantum force in importance sampling. The second is needed when we compute the kinetic energy term of the local energy.

$$\frac{\nabla \Psi}{\Psi} = \frac{\nabla(\Psi_{OB} \Psi_C)}{\Psi_{OB} \Psi_C} = \frac{\Psi_C \nabla \Psi_{OB} + \Psi_{OB} \nabla \Psi_C}{\Psi_{OB} \Psi_C} = \frac{\nabla \Psi_{OB}}{\Psi_{OB}} + \frac{\nabla \Psi_C}{\Psi_C}$$

The expectation value of the kinetic energy expressed in atomic units for electron i is

$$\langle \hat{K}_i \rangle = -\frac{1}{2} \frac{\langle \Psi | \nabla_i^2 | \Psi \rangle}{\langle \Psi | \Psi \rangle},$$

$$\hat{K}_i = -\frac{1}{2} \frac{\nabla_i^2 \Psi}{\Psi}.$$

The second derivative which enters the definition of the local energy is

$$\frac{\nabla^2 \Psi}{\Psi} = \frac{\nabla^2 \Psi_{OB}}{\Psi_{OB}} + \frac{\nabla^2 \Psi_C}{\Psi_C} + 2 \frac{\nabla \Psi_{OB}}{\Psi_{OB}} \cdot \frac{\nabla \Psi_C}{\Psi_C}$$

We discuss here how to calculate these quantities in an optimal way,

We have defined the correlated function as

$$\Psi_C = \prod_{i < j} g(r_{ij}) = \prod_{i < j}^N g(r_{ij}) = \prod_{i=1}^N \prod_{j=i+1}^N g(r_{ij}),$$

with $r_{ij} = |\mathbf{r}_i - \mathbf{r}_j| = \sqrt{(x_i - x_j)^2 + (y_i - y_j)^2 + (z_i - z_j)^2}$ in three dimensions or $r_{ij} = |\mathbf{r}_i - \mathbf{r}_j| = \sqrt{(x_i - x_j)^2 + (y_i - y_j)^2}$ if we work with two-dimensional systems.

In our particular case we have

$$\Psi_C = \prod_{i < j} g(r_{ij}) = \exp \left\{ \sum_{i < j} f(r_{ij}) \right\}.$$

The total number of different relative distances r_{ij} is $N(N-1)/2$. In a matrix storage format, the relative distances form a strictly upper triangular matrix

$$\mathbf{r} \equiv \begin{pmatrix} 0 & r_{1,2} & r_{1,3} & \cdots & r_{1,N} \\ \vdots & 0 & r_{2,3} & \cdots & r_{2,N} \\ \vdots & \vdots & 0 & \ddots & \vdots \\ \vdots & \vdots & \vdots & \ddots & r_{N-1,N} \\ 0 & 0 & 0 & \cdots & 0 \end{pmatrix}.$$

This applies to $\mathbf{g} = \mathbf{g}(r_{ij})$ as well.

In our algorithm we will move one particle at the time, say the k th-particle. This sampling will be seen to be particularly efficient when we are going to compute a Slater determinant.

We have that the ratio between Jastrow factors R_C is given by

$$R_C = \frac{\Psi_C^{\text{new}}}{\Psi_C^{\text{cur}}} = \prod_{i=1}^{k-1} \frac{g_{ik}^{\text{new}}}{g_{ik}^{\text{cur}}} \prod_{i=k+1}^N \frac{g_{ki}^{\text{new}}}{g_{ki}^{\text{cur}}}.$$

For the Pade-Jastrow form

$$R_C = \frac{\Psi_C^{\text{new}}}{\Psi_C^{\text{cur}}} = \frac{\exp U_{\text{new}}}{\exp U_{\text{cur}}} = \exp \Delta U,$$

where

$$\Delta U = \sum_{i=1}^{k-1} (f_{ik}^{\text{new}} - f_{ik}^{\text{cur}}) + \sum_{i=k+1}^N (f_{ki}^{\text{new}} - f_{ki}^{\text{cur}})$$

One needs to develop a special algorithm that runs only through the elements of the upper triangular matrix \mathbf{g} and have k as an index.

The expression to be derived in the following is of interest when computing the quantum force and the kinetic energy. It has the form

$$\frac{\nabla_i \Psi_C}{\Psi_C} = \frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_i},$$

for all dimensions and with i running over all particles.

For the first derivative only $N - 1$ terms survive the ratio because the g -terms that are not differentiated cancel with their corresponding ones in the denominator. Then,

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{1}{g_{ik}} \frac{\partial g_{ik}}{\partial x_k} + \sum_{i=k+1}^N \frac{1}{g_{ki}} \frac{\partial g_{ki}}{\partial x_k}.$$

An equivalent equation is obtained for the exponential form after replacing g_{ij} by $\exp(f_{ij})$, yielding:

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{\partial g_{ik}}{\partial x_k} + \sum_{i=k+1}^N \frac{\partial g_{ki}}{\partial x_k},$$

with both expressions scaling as $\mathcal{O}(N)$.

Using the identity

$$\frac{\partial}{\partial x_i} g_{ij} = -\frac{\partial}{\partial x_j} g_{ij},$$

we get expressions where all the derivatives acting on the particle are represented by the *second* index of g :

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{1}{g_{ik}} \frac{\partial g_{ik}}{\partial x_k} - \sum_{i=k+1}^N \frac{1}{g_{ki}} \frac{\partial g_{ki}}{\partial x_i},$$

and for the exponential case:

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{\partial g_{ik}}{\partial x_k} - \sum_{i=k+1}^N \frac{\partial g_{ki}}{\partial x_i}.$$

For correlation forms depending only on the scalar distances r_{ij} we can use the chain rule. Noting that

$$\frac{\partial g_{ij}}{\partial x_j} = \frac{\partial g_{ij}}{\partial r_{ij}} \frac{\partial r_{ij}}{\partial x_j} = \frac{x_j - x_i}{r_{ij}} \frac{\partial g_{ij}}{\partial r_{ij}},$$

we arrive at

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{1}{g_{ik}} \frac{\mathbf{r}_{ik}}{r_{ik}} \frac{\partial g_{ik}}{\partial r_{ik}} - \sum_{i=k+1}^N \frac{1}{g_{ki}} \frac{\mathbf{r}_{ki}}{r_{ki}} \frac{\partial g_{ki}}{\partial r_{ki}}.$$

Note that for the Pade-Jastrow form we can set $g_{ij} \equiv g(r_{ij}) = e^{f(r_{ij})} = e^{f_{ij}}$ and

$$\frac{\partial g_{ij}}{\partial r_{ij}} = g_{ij} \frac{\partial f_{ij}}{\partial r_{ij}}.$$

Therefore,

$$\frac{1}{\Psi_C} \frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1} \frac{\mathbf{r}_{ik}}{r_{ik}} \frac{\partial f_{ik}}{\partial r_{ik}} - \sum_{i=k+1}^N \frac{\mathbf{r}_{ki}}{r_{ki}} \frac{\partial f_{ki}}{\partial r_{ki}},$$

where

$$\mathbf{r}_{ij} = |\mathbf{r}_j - \mathbf{r}_i| = (x_j - x_i)\mathbf{e}_1 + (y_j - y_i)\mathbf{e}_2 + (z_j - z_i)\mathbf{e}_3$$

is the relative distance.

The second derivative of the Jastrow factor divided by the Jastrow factor (the way it enters the kinetic energy) is

$$\left[\frac{\nabla^2 \Psi_C}{\Psi_C} \right]_x = 2 \sum_{k=1}^N \sum_{i=1}^{k-1} \frac{\partial^2 g_{ik}}{\partial x_k^2} + \sum_{k=1}^N \left(\sum_{i=1}^{k-1} \frac{\partial g_{ik}}{\partial x_k} - \sum_{i=k+1}^N \frac{\partial g_{ki}}{\partial x_i} \right)^2$$

But we have a simple form for the function, namely

$$\Psi_C = \prod_{i < j} \exp f(r_{ij}),$$

and it is easy to see that for particle k we have

$$\frac{\nabla_k^2 \Psi_C}{\Psi_C} = \sum_{ij \neq k} \frac{(\mathbf{r}_k - \mathbf{r}_i)(\mathbf{r}_k - \mathbf{r}_j)}{r_{ki} r_{kj}} f'(r_{ki}) f'(r_{kj}) + \sum_{j \neq k} \left(f''(r_{kj}) + \frac{2}{r_{kj}} f'(r_{kj}) \right)$$

Importance sampling, Fokker-Planck and Langevin equations

A stochastic process is simply a function of two variables, one is the time, the other is a stochastic variable X , defined by specifying

- the set $\{x\}$ of possible values for X ;
- the probability distribution, $w_X(x)$, over this set, or briefly $w(x)$

The set of values $\{x\}$ for X may be discrete, or continuous. If the set of values is continuous, then $w_X(x)$ is a probability density so that $w_X(x)dx$ is the probability that one finds the stochastic variable X to have values in the range $[x, x+dx]$.

An arbitrary number of other stochastic variables may be derived from X . For example, any Y given by a mapping of X , is also a stochastic variable. The mapping may also be time-dependent, that is, the mapping depends on an additional variable t

$$Y_X(t) = f(X, t).$$

The quantity $Y_X(t)$ is called a random function, or, since t often is time, a stochastic process. A stochastic process is a function of two variables, one is the time, the other is a stochastic variable X . Let x be one of the possible values of X then

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

is a function of t , called a sample function or realization of the process. In physics one considers the stochastic process to be an ensemble of such sample functions.

For many physical systems initial distributions of a stochastic variable y tend to equilibrium distributions: $w(y, t) \rightarrow w_0(y)$ as $t \rightarrow \infty$. In equilibrium detailed balance constrains the transition rates

$$W(y \rightarrow y')w(y) = W(y' \rightarrow y)w_0(y),$$

where $W(y' \rightarrow y)$ is the probability, per unit time, that the system changes from a state $|y\rangle$, characterized by the value y for the stochastic variable Y , to a state $|y'\rangle$.

Note that for a system in equilibrium the transition rate $W(y' \rightarrow y)$ and the reverse $W(y \rightarrow y')$ may be very different.

Consider, for instance, a simple system that has only two energy levels $\varepsilon_0 = 0$ and $\varepsilon_1 = \Delta E$.

For a system governed by the Boltzmann distribution we find (the partition function has been taken out)

$$W(0 \rightarrow 1) \exp(-\varepsilon_0/kT) = W(1 \rightarrow 0) \exp(-\varepsilon_1/kT)$$

We get then

$$\frac{W(1 \rightarrow 0)}{W(0 \rightarrow 1)} = \exp(-(\Delta E/kT)),$$

which goes to zero when T tends to zero.

If we assume a discrete set of events, our initial probability distribution function can be given by

$$w_i(0) = \delta_{i,0},$$

and its time-development after a given time step $\Delta t = \varepsilon$ is

$$w_i(t) = \sum_j W(j \rightarrow i) w_j(t=0).$$

The continuous analog to $w_i(0)$ is

$$w(\mathbf{x}) \rightarrow \delta(\mathbf{x}),$$

where we now have generalized the one-dimensional position x to a generic-dimensional vector \mathbf{x} . The Kroenecker δ function is replaced by the δ distribution function $\delta(\mathbf{x})$ at $t = 0$.

The transition from a state j to a state i is now replaced by a transition to a state with position \mathbf{y} from a state with position \mathbf{x} . The discrete sum of transition probabilities can then be replaced by an integral and we obtain the new distribution at a time $t + \Delta t$ as

$$w(\mathbf{y}, t + \Delta t) = \int W(\mathbf{y}, t + \Delta t | \mathbf{x}, t) w(\mathbf{x}, t) d\mathbf{x},$$

and after m time steps we have

$$w(\mathbf{y}, t + m\Delta t) = \int W(\mathbf{y}, t + m\Delta t | \mathbf{x}, t) w(\mathbf{x}, t) d\mathbf{x}.$$

When equilibrium is reached we have

$$w(\mathbf{y}) = \int W(\mathbf{y} | \mathbf{x}, t) w(\mathbf{x}) d\mathbf{x},$$

that is no time-dependence. Note our change of notation for W

We can solve the equation for $w(\mathbf{y}, t)$ by making a Fourier transform to momentum space. The PDF $w(\mathbf{x}, t)$ is related to its Fourier transform $\tilde{w}(\mathbf{k}, t)$ through

$$w(\mathbf{x}, t) = \int_{-\infty}^{\infty} d\mathbf{k} \exp(i\mathbf{k}\mathbf{x}) \tilde{w}(\mathbf{k}, t),$$

and using the definition of the δ -function

$$\delta(\mathbf{x}) = \frac{1}{2\pi} \int_{-\infty}^{\infty} d\mathbf{k} \exp(i\mathbf{k}\mathbf{x}),$$

we see that

$$\tilde{w}(\mathbf{k}, 0) = 1/2\pi.$$

We can then use the Fourier-transformed diffusion equation

$$\frac{\partial \tilde{w}(\mathbf{k}, t)}{\partial t} = -D\mathbf{k}^2 \tilde{w}(\mathbf{k}, t),$$

with the obvious solution

$$\tilde{w}(\mathbf{k}, t) = \tilde{w}(\mathbf{k}, 0) \exp[-(D\mathbf{k}^2 t)] = \frac{1}{2\pi} \exp[-(D\mathbf{k}^2 t)].$$

With the Fourier transform we obtain

$$w(\mathbf{x}, t) = \int_{-\infty}^{\infty} d\mathbf{k} \exp[i\mathbf{k}\mathbf{x}] \frac{1}{2\pi} \exp[-(D\mathbf{k}^2 t)] = \frac{1}{\sqrt{4\pi Dt}} \exp[-(\mathbf{x}^2/4Dt)],$$

with the normalization condition

$$\int_{-\infty}^{\infty} w(\mathbf{x}, t) d\mathbf{x} = 1.$$

The solution represents the probability of finding our random walker at position \mathbf{x} at time t if the initial distribution was placed at $\mathbf{x} = 0$ at $t = 0$.

There is another interesting feature worth observing. The discrete transition probability W itself is given by a binomial distribution. The results from the central limit theorem state that transition probability in the limit $n \rightarrow \infty$ converges to the normal distribution. It is then possible to show that

$$W(il - jl, n\epsilon) \rightarrow W(\mathbf{y}, t + \Delta t | \mathbf{x}, t) = \frac{1}{\sqrt{4\pi D \Delta t}} \exp[-((\mathbf{y} - \mathbf{x})^2/4D\Delta t)],$$

and that it satisfies the normalization condition and is itself a solution to the diffusion equation.

Let us now assume that we have three PDFs for times $t_0 < t' < t$, that is $w(\mathbf{x}_0, t_0)$, $w(\mathbf{x}', t')$ and $w(\mathbf{x}, t)$. We have then

$$w(\mathbf{x}, t) = \int_{-\infty}^{\infty} W(\mathbf{x}, t | \mathbf{x}', t') w(\mathbf{x}', t') d\mathbf{x}',$$

and

$$w(\mathbf{x}, t) = \int_{-\infty}^{\infty} W(\mathbf{x}, t | \mathbf{x}_0, t_0) w(\mathbf{x}_0, t_0) d\mathbf{x}_0,$$

and

$$w(\mathbf{x}', t') = \int_{-\infty}^{\infty} W(\mathbf{x}', t' | \mathbf{x}_0, t_0) w(\mathbf{x}_0, t_0) d\mathbf{x}_0.$$

We can combine these equations and arrive at the famous Einstein-Smoluchenski-Kolmogorov-Chapman (ESKC) relation

$$W(\mathbf{x}, t | \mathbf{x}_0, t_0) = \int_{-\infty}^{\infty} W(\mathbf{x}, t | \mathbf{x}', t') W(\mathbf{x}', t' | \mathbf{x}_0, t_0) d\mathbf{x}'.$$

We can replace the spatial dependence with a dependence upon say the velocity (or momentum), that is we have

$$W(\mathbf{v}, t | \mathbf{v}_0, t_0) = \int_{-\infty}^{\infty} W(\mathbf{v}, t | \mathbf{v}', t') W(\mathbf{v}', t' | \mathbf{v}_0, t_0) d\mathbf{x}'.$$

We will now derive the Fokker-Planck equation. We start from the ESKC equation

$$W(\mathbf{x}, t | \mathbf{x}_0, t_0) = \int_{-\infty}^{\infty} W(\mathbf{x}, t | \mathbf{x}', t') W(\mathbf{x}', t' | \mathbf{x}_0, t_0) d\mathbf{x}'.$$

Define $s = t' - t_0$, $\tau = t - t'$ and $t - t_0 = s + \tau$. We have then

$$W(\mathbf{x}, s + \tau | \mathbf{x}_0) = \int_{-\infty}^{\infty} W(\mathbf{x}, \tau | \mathbf{x}') W(\mathbf{x}', s | \mathbf{x}_0) d\mathbf{x}'.$$

Assume now that τ is very small so that we can make an expansion in terms of a small step xi , with $\mathbf{x}' = \mathbf{x} - \xi$, that is

$$W(\mathbf{x}, s | \mathbf{x}_0) + \frac{\partial W}{\partial s} \tau + O(\tau^2) = \int_{-\infty}^{\infty} W(\mathbf{x}, \tau | \mathbf{x} - \xi) W(\mathbf{x} - \xi, s | \mathbf{x}_0) d\mathbf{x}'.$$

We assume that $W(\mathbf{x}, \tau | \mathbf{x} - \xi)$ takes non-negligible values only when ξ is small. This is just another way of stating the Master equation!!

We say thus that \mathbf{x} changes only by a small amount in the time interval τ . This means that we can make a Taylor expansion in terms of ξ , that is we expand

$$W(\mathbf{x}, \tau | \mathbf{x} - \xi) W(\mathbf{x} - \xi, s | \mathbf{x}_0) = \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} [W(\mathbf{x} + \xi, \tau | \mathbf{x}) W(\mathbf{x}, s | \mathbf{x}_0)].$$

We can then rewrite the ESKC equation as

$$\frac{\partial W}{\partial s} \tau = -W(\mathbf{x}, s | \mathbf{x}_0) + \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} \left[W(\mathbf{x}, s | \mathbf{x}_0) \int_{-\infty}^{\infty} \xi^n W(\mathbf{x} + \xi, \tau | \mathbf{x}) d\xi \right].$$

We have neglected higher powers of τ and have used that for $n = 0$ we get simply $W(\mathbf{x}, s | \mathbf{x}_0)$ due to normalization.

We say thus that \mathbf{x} changes only by a small amount in the time interval τ . This means that we can make a Taylor expansion in terms of ξ , that is we expand

$$W(\mathbf{x}, \tau | \mathbf{x} - \xi) W(\mathbf{x} - \xi, s | \mathbf{x}_0) = \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} [W(\mathbf{x} + \xi, \tau | \mathbf{x}) W(\mathbf{x}, s | \mathbf{x}_0)].$$

We can then rewrite the ESKC equation as

$$\frac{\partial W(\mathbf{x}, s | \mathbf{x}_0)}{\partial s} \tau = -W(\mathbf{x}, s | \mathbf{x}_0) + \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} \left[W(\mathbf{x}, s | \mathbf{x}_0) \int_{-\infty}^{\infty} \xi^n W(\mathbf{x} + \xi, \tau | \mathbf{x}) d\xi \right].$$

We have neglected higher powers of τ and have used that for $n = 0$ we get simply $W(\mathbf{x}, s | \mathbf{x}_0)$ due to normalization.

We simplify the above by introducing the moments

$$M_n = \frac{1}{\tau} \int_{-\infty}^{\infty} \xi^n W(\mathbf{x} + \xi, \tau | \mathbf{x}) d\xi = \frac{\langle [\Delta x(\tau)]^n \rangle}{\tau},$$

resulting in

$$\frac{\partial W(\mathbf{x}, s | \mathbf{x}_0)}{\partial s} = \sum_{n=1}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} [W(\mathbf{x}, s | \mathbf{x}_0) M_n].$$

When $\tau \rightarrow 0$ we assume that $\langle [\Delta x(\tau)]^n \rangle \rightarrow 0$ more rapidly than τ itself if $n > 2$. When τ is much larger than the standard correlation time of system then M_n for $n > 2$ can normally be neglected. This means that fluctuations become negligible at large time scales.

If we neglect such terms we can rewrite the ESKC equation as

$$\frac{\partial W(\mathbf{x}, s | \mathbf{x}_0)}{\partial s} = -\frac{\partial M_1 W(\mathbf{x}, s | \mathbf{x}_0)}{\partial x} + \frac{1}{2} \frac{\partial^2 M_2 W(\mathbf{x}, s | \mathbf{x}_0)}{\partial x^2}.$$

In a more compact form we have

$$\frac{\partial W}{\partial s} = -\frac{\partial M_1 W}{\partial x} + \frac{1}{2} \frac{\partial^2 M_2 W}{\partial x^2},$$

which is the Fokker-Planck equation! It is trivial to replace position with velocity (momentum).

Consider a particle suspended in a liquid. On its path through the liquid it will continuously collide with the liquid molecules. Because on average the particle will collide more often on the front side than on the back side, it will experience a systematic force proportional with its velocity, and directed opposite to its velocity. Besides this systematic force the particle will experience a stochastic force $\mathbf{F}(t)$. The equations of motion are

- $\frac{d\mathbf{r}}{dt} = \mathbf{v}$ and
- $\frac{d\mathbf{v}}{dt} = -\xi\mathbf{v} + \mathbf{F}$.

From hydrodynamics we know that the friction constant ξ is given by

$$\xi = 6\pi\eta a/m$$

where η is the viscosity of the solvent and a is the radius of the particle .

Solving the second equation in the previous slide we get

$$\mathbf{v}(t) = \mathbf{v}_0 e^{-\xi t} + \int_0^t d\tau e^{-\xi(t-\tau)} \mathbf{F}(\tau).$$

If we want to get some useful information out of this, we have to average over all possible realizations of $\mathbf{F}(t)$, with the initial velocity as a condition. A useful quantity for example is

$$\begin{aligned} \langle \mathbf{v}(t) \cdot \mathbf{v}(t) \rangle_{\mathbf{v}_0} &= v_0^2 e^{-2\xi t} + 2 \int_0^t d\tau e^{-\xi(2t-\tau)} \mathbf{v}_0 \cdot \langle \mathbf{F}(\tau) \rangle_{\mathbf{v}_0} \\ &+ \int_0^t d\tau' \int_0^t d\tau e^{-\xi(2t-\tau-\tau')} \langle \mathbf{F}(\tau) \cdot \mathbf{F}(\tau') \rangle_{\mathbf{v}_0}. \end{aligned}$$

In order to continue we have to make some assumptions about the conditional averages of the stochastic forces. In view of the chaotic character of the stochastic forces the following assumptions seem to be appropriate

$$\langle \mathbf{F}(t) \rangle = 0,$$

and

$$\langle \mathbf{F}(t) \cdot \mathbf{F}(t') \rangle_{\mathbf{v}_0} = C_{\mathbf{v}_0} \delta(t - t').$$

We omit the subscript \mathbf{v}_0 , when the quantity of interest turns out to be independent of \mathbf{v}_0 . Using the last three equations we get

$$\langle \mathbf{v}(t) \cdot \mathbf{v}(t) \rangle_{\mathbf{v}_0} = v_0^2 e^{-2\xi t} + \frac{C_{\mathbf{v}_0}}{2\xi} (1 - e^{-2\xi t}).$$

For large t this should be equal to $3kT/m$, from which it follows that

$$\langle \mathbf{F}(t) \cdot \mathbf{F}(t') \rangle = 6 \frac{kT}{m} \xi \delta(t - t').$$

This result is called the fluctuation-dissipation theorem .

Integrating

$$\mathbf{v}(t) = \mathbf{v}_0 e^{-\xi t} + \int_0^t d\tau e^{-\xi(t-\tau)} \mathbf{F}(\tau),$$

we get

$$\mathbf{r}(t) = \mathbf{r}_0 + \mathbf{v}_0 \frac{1}{\xi} (1 - e^{-\xi t}) + \int_0^t d\tau \int_0^\tau \tau' e^{-\xi(\tau-\tau')} \mathbf{F}(\tau'),$$

from which we calculate the mean square displacement

$$\langle (\mathbf{r}(t) - \mathbf{r}_0)^2 \rangle_{\mathbf{v}_0} = \frac{v_0^2}{\xi^2} (1 - e^{-\xi t})^2 + \frac{3kT}{m\xi^2} (2\xi t - 3 + 4e^{-\xi t} - e^{-2\xi t}).$$

For very large t this becomes

$$\langle (\mathbf{r}(t) - \mathbf{r}_0)^2 \rangle = \frac{6kT}{m\xi} t$$

from which we get the Einstein relation

$$D = \frac{kT}{m\xi}$$

where we have used $\langle (\mathbf{r}(t) - \mathbf{r}_0)^2 \rangle = 6Dt$.

Code example for two electrons in a quantum dots

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
2-electron VMC code for 2dim quantum dot with importance sampling Using gaussian rng
for new positions and Metropolis- Hastings No energy minimization from math import exp,
sqrt from random import random, seed, normalvariate import numpy as np import mat-
plotlib.pyplot as plt from mpl_toolkits.mplot3d import Axes3D from matplotlib import cm from matplotlib.ticker import LinearLocator, Form
```

```
Read name of output file from command line if len(sys.argv) == 2: outfilename = sys.argv[1]
else: print('Name of output file must be given as command line argument.') outfile =
open(outfilename, 'w')
```

```
Trial wave function for the 2-electron quantum dot in two dims def WaveFunction(r,alpha,beta):
r1 = r[0,0]**2 + r[0,1]**2 r2 = r[1,0]**2 + r[1,1]**2 r12 = sqrt((r[0,0]-r[1,0])**2 + (r[0,1]-
r[1,1])**2) deno = r12/(1+beta*r12) return exp(-0.5*alpha*(r1+r2)+deno)
```

```
Local energy for the 2-electron quantum dot in two dims, using analytical local energy def
LocalEnergy(r,alpha,beta):
```

```
r1 = (r[0,0]**2 + r[0,1]**2) r2 = (r[1,0]**2 + r[1,1]**2) r12 = sqrt((r[0,0]-r[1,0])**2 +
(r[0,1]-r[1,1])**2) deno = 1.0/(1+beta*r12) deno2 = deno*deno return 0.5*(1-alpha*alpha)*(r1
+ r2) + 2.0*alpha + 1.0/r12+deno2*(alpha*r12-deno2+2*beta*deno-1.0/r12)
```

```
Setting up the quantum force for the two-electron quantum dot, recall that it is a vector
def QuantumForce(r,alpha,beta):
```

```
qforce = np.zeros((NumberParticles,Dimension), np.double) r12 = sqrt((r[0,0]-r[1,0])**2 +
(r[0,1]-r[1,1])**2) deno = 1.0/(1+beta*r12) qforce[0,:] = -2*r[0,:]*alpha*(r[0,:]-r[1,:])*deno*deno/r12
qforce[1,:] = -2*r[1,:]*alpha*(r[1,:]-r[0,:])*deno*deno/r12 return qforce
```

```
The Monte Carlo sampling with the Metropolis algo jit decorator tells Numba to compile
this function. The argument types will be inferred by Numba when function is called. @jit()
def MonteCarloSampling():
```

```
NumberMCCycles= 100000 Parameters in the Fokker-Planck simulation of the quantum
force D = 0.5 TimeStep = 0.05 positions PositionOld = np.zeros((NumberParticles,Dimension),
np.double) PositionNew = np.zeros((NumberParticles,Dimension), np.double) Quantum force
QuantumForceOld = np.zeros((NumberParticles,Dimension), np.double) QuantumForceNew
= np.zeros((NumberParticles,Dimension), np.double)
```

```
seed for rng generator seed() start variational parameter loops, two parameters here alpha
= 0.9 for ia in range(MaxVariations): alpha += .025 AlphaValues[ia] = alpha beta = 0.2 for jb
in range(MaxVariations): beta += .01 BetaValues[jb] = beta energy = energy2 = 0.0 DeltaE =
0.0 Initial position for i in range(NumberParticles): for j in range(Dimension): PositionOld[i,j]
= normalvariate(0.0,1.0)*sqrt(TimeStep) wfold = WaveFunction(PositionOld,alpha,beta) Quan-
tumForceOld = QuantumForce(PositionOld,alpha, beta)
```

```
Loop over MC MCCycles for MCCycle in range(NumberMCCycles): Trial position mov-
ing one particle at the time for i in range(NumberParticles): for j in range(Dimension):
PositionNew[i,j] = PositionOld[i,j]+normalvariate(0.0,1.0)*sqrt(TimeStep)+ QuantumForce-
Old[i,j]*TimeStep*D wfnew = WaveFunction(PositionNew,alpha,beta) QuantumForceNew =
QuantumForce(PositionNew,alpha, beta) GreensFunction = 0.0 for j in range(Dimension):
GreensFunction += 0.5*(QuantumForceOld[i,j]+QuantumForceNew[i,j])*(D*TimeStep*0.5*(QuantumForceOld[i,j]-
QuantumForceNew[i,j])- PositionNew[i,j]+PositionOld[i,j])
```


GreensFunction = exp(GreensFunction) ProbabilityRatio = GreensFunction*wfnew**2/wfold**2
 Metropolis-Hastings test to see whether we accept the move if random() <= ProbabilityRatio:
 for j in range(Dimension): PositionOld[i,j] = PositionNew[i,j] QuantumForceOld[i,j] = QuantumForceNew[i,j] wfold = wfnew DeltaE = LocalEnergy(PositionOld,alpha,beta) energy += DeltaE energy2 += DeltaE**2 We calculate mean, variance and error (no blocking applied)
 energy /= NumberMCCycles energy2 /= NumberMCCycles variance = energy2 - energy**2
 error = sqrt(variance/NumberMCCycles) Energies[ia,jb] = energy outfile.write('return Energies, AlphaValues, BetaValues

Here starts the main program with variable declarations NumberParticles = 2 Dimension = 2 MaxVariations = 10 Energies = np.zeros((MaxVariations,MaxVariations)) AlphaValues = np.zeros(MaxVariations) BetaValues = np.zeros(MaxVariations) (Energies, AlphaValues, BetaValues) = MonteCarloSampling() outfile.close() Prepare for plots fig = plt.figure() ax = fig.gca(projection='3d') Plot the surface. X, Y = np.meshgrid(AlphaValues, BetaValues) surf = ax.plot_surface(X, Y, Energies, cmap = cm.coolwarm, linewidth=0, antialiased=False) Customizethezaxis.zmin = np.matrix(Energies).min() zmax = np.matrix(Energies).max() ax.set_zlim(zmin, zmax) ax.set_xlabel(r'\alpha') ax.set_ylabel(r'\beta') ax.set_zlabel(r'\langle E \rangle') ax.xaxis.set_major_locator(LinearLocator(10)) ax.xaxis.set_major_formatter(FormatStrFormatter('Addcolorbarwhichm 0.5, aspect = 5)) plt.show()

Bringing the gradient optimization.

The simple one-particle case in a harmonic oscillator trap

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
Gradient descent stepping with analytical derivative import numpy as np from scipy.optimize
import minimize def DerivativeE(x): return x-1.0/(4*x*x*x);
def Energy(x): return x*x*0.5+1.0/(8*x*x); x0 = 1.0 eta = 0.1 Niterations = 100
for iter in range(Niterations): gradients = DerivativeE(x0) x0 -= eta*gradients
print(x0)
```

```
[fontsize=,linenos=false,mathescape,baselinestretch=1.0,fontfamily=tt,xleftmargin=7mm]python
2-electron VMC code for 2dim quantum dot with importance sampling Using gaussian rng
for new positions and Metropolis- Hastings from math import exp, sqrt from random im-
port random, seed, normalvariate import numpy as np import matplotlib.pyplot as plt from
mpl_toolkits.mplot3d import Axes3D from matplotlib import cm from matplotlib.ticker import LinearLocator, FormatStrFormatter import sys from
```

Trial wave function for the 2-electron quantum dot in two dims def WaveFunction(r,alpha):
 r1 = r[0,0]**2 + r[0,1]**2 r2 = r[1,0]**2 + r[1,1]**2 return exp(-0.5*alpha*(r1+r2))

Local energy for the 2-electron quantum dot in two dims, using analytical local energy def
 LocalEnergy(r,alpha):

r1 = (r[0,0]**2 + r[0,1]**2) r2 = (r[1,0]**2 + r[1,1]**2) return 0.5*(1-alpha*alpha)*(r1 +
 r2) + 2.0*alpha

Derivate of wave function ansatz as function of variational parameters def DerivativeW-
 Fansatz(r,alpha):

r1 = (r[0,0]**2 + r[0,1]**2) r2 = (r[1,0]**2 + r[1,1]**2) WfDer = -(r1+r2) return WfDer

Setting up the quantum force for the two-electron quantum dot, recall that it is a vector
 def QuantumForce(r,alpha):

qforce = np.zeros((NumberParticles,Dimension), np.double) qforce[0,:] = -2*r[0,:]*alpha
 qforce[1,:] = -2*r[1,:]*alpha return qforce

Computing the derivative of the energy and the energy jit decorator tells Numba to com-
 pile this function. The argument types will be inferred by Numba when function is called.
 @jit def EnergyMinimization(alpha):

NumberMCCycles= 1000 Parameters in the Fokker-Planck simulation of the quantum
 force D = 0.5 TimeStep = 0.05 positions PositionOld = np.zeros((NumberParticles,Dimension),

```

np.double) PositionNew = np.zeros((NumberParticles,Dimension), np.double) Quantum force
QuantumForceOld = np.zeros((NumberParticles,Dimension), np.double) QuantumForceNew
= np.zeros((NumberParticles,Dimension), np.double)
seed for rng generator seed() energy = 0.0 DeltaE = 0.0 EnergyDer = 0.0 DeltaPsi = 0.0
DerivativePsiE = 0.0 Initial position for i in range(NumberParticles): for j in range(Dimension):
PositionOld[i,j] = normalvariate(0.0,1.0)*sqrt(TimeStep) wfold = WaveFunction(PositionOld,alpha)
QuantumForceOld = QuantumForce(PositionOld,alpha)
Loop over MC MCcycles for MCcycle in range(NumberMCcycles): Trial position mov-
ing one particle at the time for i in range(NumberParticles): for j in range(Dimension):
PositionNew[i,j] = PositionOld[i,j]+normalvariate(0.0,1.0)*sqrt(TimeStep)+ QuantumForce-
Old[i,j]*TimeStep*D wfnew = WaveFunction(PositionNew,alpha) QuantumForceNew = Quan-
tumForce(PositionNew,alpha) GreensFunction = 0.0 for j in range(Dimension): GreensFunc-
tion += 0.5*(QuantumForceOld[i,j]+QuantumForceNew[i,j])* (D*TimeStep*0.5*(QuantumForceOld[i,j]-
QuantumForceNew[i,j])- PositionNew[i,j]+PositionOld[i,j])
GreensFunction = exp(GreensFunction) ProbabilityRatio = GreensFunction*wfnew**2/wfold**2
Metropolis-Hastings test to see whether we accept the move if random() <= ProbabilityRatio:
for j in range(Dimension): PositionOld[i,j] = PositionNew[i,j] QuantumForceOld[i,j] = Quan-
tumForceNew[i,j] wfold = wfnew DeltaE = LocalEnergy(PositionOld,alpha) DeltaPsi = Deriva-
tiveWFansatz(PositionOld,alpha) energy += DeltaE DerivativePsiE += DeltaPsi*DeltaE
We calculate mean, variance and error (no blocking applied) energy /= NumberMCcycles
DerivativePsiE /= NumberMCcycles DeltaPsi /= NumberMCcycles EnergyDer = 2*(DerivativePsiE-
DeltaPsi*energy) return energy, EnergyDer
Here starts the main program with variable declarations NumberParticles = 2 Dimension =
2 guess for variational parameters x0 = 1.5 Set up iteration using stochastic gradient method
Energy = 0 ; EnergyDer = 0 Energy, EnergyDer = EnergyMinimization(x0) print(Energy, En-
ergyDer)
eta = 0.01 Niterations = 100
for iter in range(Niterations): gradients = EnergyDer x0 -= eta*gradients Energy, Energy-
Der = EnergyMinimization(x0)
print(x0)

```

VMC for fermions: Efficient calculation of Slater determinants

The potentially most time-consuming part is the evaluation of the gradient and the Laplacian of an N -particle Slater determinant.

We have to differentiate the determinant with respect to all spatial coordinates of all particles. A brute force differentiation would involve $N \cdot d$ evaluations of the entire determinant which would even worsen the already undesirable time scaling, making it $Nd \cdot O(N^3) \sim O(d \cdot N^4)$.

This poses serious hindrances to the overall efficiency of our code.

The efficiency can be improved however if we move only one electron at the time. The Slater determinant matrix \hat{D} is defined by the matrix elements

$$d_{ij} = \phi_j(x_i)$$

where $\phi_j(\mathbf{r}_i)$ is a single particle wave function. The columns correspond to the position of a given particle while the rows stand for the various quantum numbers.

What we need to realize is that when differentiating a Slater determinant with respect to some given coordinate, only one row of the corresponding Slater matrix is changed.

Therefore, by recalculating the whole determinant we risk producing redundant information. The solution turns out to be an algorithm that requires to keep track of the *inverse* of the Slater matrix.

Let the current position in phase space be represented by the $(N \cdot d)$ -element vector \mathbf{r}^{old} and the new suggested position by the vector \mathbf{r}^{new} .

The inverse of \hat{D} can be expressed in terms of its cofactors C_{ij} and its determinant (this our notation for a determinant) $|\hat{D}|$:

$$d_{ij}^{-1} = \frac{C_{ji}}{|\hat{D}|} \quad (6.7)$$

Notice that the interchanged indices indicate that the matrix of cofactors is to be transposed.

If \hat{D} is invertible, then we must obviously have $\hat{D}^{-1}\hat{D} = \mathbf{1}$, or explicitly in terms of the individual elements of \hat{D} and \hat{D}^{-1} :

$$\sum_{k=1}^N d_{ik} d_{kj}^{-1} = \delta_{ij} \quad (6.8)$$

Consider the ratio, which we shall call R , between $|\hat{D}(\mathbf{r}^{\text{new}})|$ and $|\hat{D}(\mathbf{r}^{\text{old}})|$. By definition, each of these determinants can individually be expressed in terms of the i -th row of its cofactor matrix

$$R = \frac{|\hat{D}(\mathbf{r}^{\text{new}})|}{|\hat{D}(\mathbf{r}^{\text{old}})|} = \frac{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{new}}) C_{ij}(\mathbf{r}^{\text{new}})}{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{old}}) C_{ij}(\mathbf{r}^{\text{old}})} \quad (6.9)$$

Suppose now that we move only one particle at a time, meaning that \mathbf{r}^{new} differs from \mathbf{r}^{old} by the position of only one, say the i -th, particle. This means that $\hat{D}(\mathbf{r}^{\text{new}})$ and $\hat{D}(\mathbf{r}^{\text{old}})$ differ only by the entries of the i -th row. Recall also that the i -th row of a cofactor matrix \hat{C} is independent of the entries of the i -th row of its corresponding matrix \hat{D} . In this particular case we therefore get that the i -th row of $\hat{C}(\mathbf{r}^{\text{new}})$ and $\hat{C}(\mathbf{r}^{\text{old}})$ must be equal. Explicitly, we have:

$$C_{ij}(\mathbf{r}^{\text{new}}) = C_{ij}(\mathbf{r}^{\text{old}}) \quad \forall j \in \{1, \dots, N\} \quad (6.10)$$

Inserting this into the numerator of eq. (??) and using eq. (??) to substitute the cofactors with the elements of the inverse matrix, we get:

$$R = \frac{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{new}}) C_{ij}(\mathbf{r}^{\text{old}})}{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{old}}) C_{ij}(\mathbf{r}^{\text{old}})} = \frac{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{new}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}})}{\sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{old}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}})} \quad (6.11)$$

Now by eq. (??) the denominator of the rightmost expression must be unity, so that we finally arrive at:

$$R = \sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{new}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}}) = \sum_{j=1}^N \phi_j(\mathbf{r}_i^{\text{new}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}}) \quad (6.12)$$

What this means is that in order to get the ratio when only the i -th particle has been moved, we only need to calculate the dot product of the vector $(\phi_1(\mathbf{r}_i^{\text{new}}), \dots, \phi_N(\mathbf{r}_i^{\text{new}}))$ of single particle wave functions evaluated at this new position with the i -th column of the inverse matrix \hat{D}^{-1} evaluated at the original position. Such an operation has a time scaling of $O(N)$. The only extra thing we need to do is to maintain the inverse matrix $\hat{D}^{-1}(\mathbf{x}^{\text{old}})$.

If the new position \mathbf{r}^{new} is accepted, then the inverse matrix can be suitably updated by an algorithm having a time scaling of $O(N^2)$. This algorithm goes as follows. First we update all but the i -th column of \hat{D}^{-1} . For each column $j \neq i$, we first calculate the quantity:

$$S_j = (\hat{D}(\mathbf{r}^{\text{new}}) \times \hat{D}^{-1}(\mathbf{r}^{\text{old}}))_{ij} = \sum_{l=1}^N d_{il}(\mathbf{r}^{\text{new}}) d_{lj}^{-1}(\mathbf{r}^{\text{old}}) \quad (6.13)$$

The new elements of the j -th column of \hat{D}^{-1} are then given by:

$$d_{kj}^{-1}(\mathbf{r}^{\text{new}}) = d_{kj}^{-1}(\mathbf{r}^{\text{old}}) - \frac{S_j}{R} d_{ki}^{-1}(\mathbf{r}^{\text{old}}) \quad \forall \begin{matrix} k \in \{1, \dots, N\} \\ j \neq i \end{matrix} \quad (6.14)$$

Finally the elements of the i -th column of \hat{D}^{-1} are updated simply as follows:

$$d_{ki}^{-1}(\mathbf{r}^{\text{new}}) = \frac{1}{R} d_{ki}^{-1}(\mathbf{r}^{\text{old}}) \quad \forall \quad k \in \{1, \dots, N\} \quad (6.15)$$

We see from these formulas that the time scaling of an update of \hat{D}^{-1} after changing one row of \hat{D} is $O(N^2)$.

The scheme is also applicable for the calculation of the ratios involving derivatives. It turns out that differentiating the Slater determinant with respect to the coordinates of a single particle \mathbf{r}_i changes only the i -th row of the corresponding Slater matrix.

The gradient and the Laplacian.

The gradient and the Laplacian can therefore be calculated as follows:

$$\frac{\nabla_i |\hat{D}(\mathbf{r})|}{|\hat{D}(\mathbf{r})|} = \sum_{j=1}^N \nabla_i d_{ij}(\mathbf{r}) d_{ji}^{-1}(\mathbf{r}) = \sum_{j=1}^N \nabla_i \phi_j(\mathbf{r}_i) d_{ji}^{-1}(\mathbf{r})$$

and

$$\frac{\nabla_i^2 |\hat{D}(\mathbf{r})|}{|\hat{D}(\mathbf{r})|} = \sum_{j=1}^N \nabla_i^2 d_{ij}(\mathbf{r}) d_{ji}^{-1}(\mathbf{r}) = \sum_{j=1}^N \nabla_i^2 \phi_j(\mathbf{r}_i) d_{ji}^{-1}(\mathbf{r})$$

Thus, to calculate all the derivatives of the Slater determinant, we only need the derivatives of the single particle wave functions ($\nabla_i \phi_j(\mathbf{r}_i)$ and $\nabla_i^2 \phi_j(\mathbf{r}_i)$) and the elements of the corresponding inverse Slater matrix ($\hat{D}^{-1}(\mathbf{r}_i)$). A calculation of a single derivative is by the above result an $O(N)$ operation. Since there are $d \cdot N$ derivatives, the time scaling of the total evaluation becomes $O(d \cdot N^2)$. With an $O(N^2)$ updating algorithm for the inverse matrix, the total scaling is no worse, which is far better than the brute force approach yielding $O(d \cdot N^4)$.

Important note: In most cases you end with closed form expressions for the single-particle wave functions. It is then useful to calculate the various derivatives and make separate functions for them.

The Slater determinant takes the form

$$\Phi(\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_3, \mathbf{r}_4, \alpha, \beta, \gamma, \delta) = \frac{1}{\sqrt{4!}} \begin{vmatrix} \psi_{100\uparrow}(\mathbf{r}_1) & \psi_{100\uparrow}(\mathbf{r}_2) & \psi_{100\uparrow}(\mathbf{r}_3) & \psi_{100\uparrow}(\mathbf{r}_4) \\ \psi_{100\downarrow}(\mathbf{r}_1) & \psi_{100\downarrow}(\mathbf{r}_2) & \psi_{100\downarrow}(\mathbf{r}_3) & \psi_{100\downarrow}(\mathbf{r}_4) \\ \psi_{200\uparrow}(\mathbf{r}_1) & \psi_{200\uparrow}(\mathbf{r}_2) & \psi_{200\uparrow}(\mathbf{r}_3) & \psi_{200\uparrow}(\mathbf{r}_4) \\ \psi_{200\downarrow}(\mathbf{r}_1) & \psi_{200\downarrow}(\mathbf{r}_2) & \psi_{200\downarrow}(\mathbf{r}_3) & \psi_{200\downarrow}(\mathbf{r}_4) \end{vmatrix}.$$

The Slater determinant as written is zero since the spatial wave functions for the spin up and spin down states are equal. But we can rewrite it as the product of two Slater determinants, one for spin up and one for spin down.

We can rewrite it as

$$\begin{aligned} \Phi(\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_3, \mathbf{r}_4, \alpha, \beta, \gamma, \delta) = & |\uparrow(1,2)| \downarrow(3,4) - |\uparrow(1,3)| \downarrow(2,4) \\ & - |\uparrow(1,4)| \downarrow(3,2) + |\uparrow(2,3)| \downarrow(1,4) - |\uparrow(2,4)| \downarrow(1,3) \\ & + |\uparrow(3,4)| \downarrow(1,2), \end{aligned}$$

where we have defined

$$|\uparrow|(1,2) = \frac{1}{\sqrt{2}} \begin{vmatrix} \psi_{100\uparrow}(\mathbf{r}_1) & \psi_{100\uparrow}(\mathbf{r}_2) \\ \psi_{200\uparrow}(\mathbf{r}_1) & \psi_{200\uparrow}(\mathbf{r}_2) \end{vmatrix},$$

and

$$|\downarrow|(3,4) = \frac{1}{\sqrt{2}} \begin{vmatrix} \psi_{100\downarrow}(\mathbf{r}_3) & \psi_{100\downarrow}(\mathbf{r}_4) \\ \psi_{200\downarrow}(\mathbf{r}_3) & \psi_{200\downarrow}(\mathbf{r}_4) \end{vmatrix}.$$

We want to avoid to sum over spin variables, in particular when the interaction does not depend on spin.

It can be shown, see for example Moskowitz and Kalos, *Int. J. Quantum Chem.* **20** 1107 (1981), that for the variational energy we can approximate the Slater determinant as

$$\Phi(\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_3, \mathbf{r}_4, \alpha, \beta, \gamma, \delta) \propto |\uparrow|(1,2) |\downarrow|(3,4),$$

or more generally as

$$\Phi(\mathbf{r}_1, \mathbf{r}_2, \dots, \mathbf{r}_N) \propto |\uparrow| |\downarrow|,$$

where we have the Slater determinant as the product of a spin up part involving the number of electrons with spin up only (2 for beryllium and 5 for neon) and a spin down part involving the electrons with spin down.

This ansatz is not antisymmetric under the exchange of electrons with opposite spins but it can be shown (show this) that it gives the same expectation value for the energy as the full Slater determinant.

As long as the Hamiltonian is spin independent, the above is correct. It is rather straightforward to see this if you go back to the equations for the energy discussed earlier this semester.

We will thus factorize the full determinant $|\hat{D}|$ into two smaller ones, where each can be identified with \uparrow and \downarrow respectively:

$$|\hat{D}| = |\hat{D}|_{\uparrow} \cdot |\hat{D}|_{\downarrow}$$

The combined dimensionality of the two smaller determinants equals the dimensionality of the full determinant. Such a factorization is advantageous in that it makes it possible to perform the calculation of the ratio R and the updating of the inverse matrix separately for $|\hat{D}|_{\uparrow}$ and $|\hat{D}|_{\downarrow}$:

$$\frac{|\hat{D}|_{\uparrow}^{\text{new}}}{|\hat{D}|_{\uparrow}^{\text{old}}} = \frac{|\hat{D}|_{\uparrow}^{\text{new}}}{|\hat{D}|_{\uparrow}^{\text{old}}} \cdot \frac{|\hat{D}|_{\downarrow}^{\text{new}}}{|\hat{D}|_{\downarrow}^{\text{old}}}$$

This reduces the calculation time by a constant factor. The maximal time reduction happens in a system of equal numbers of \uparrow and \downarrow particles, so that the two factorized determinants are half the size of the original one.

Consider the case of moving only one particle at a time which originally had the following time scaling for one transition:

$$O_R(N) + O_{\text{inverse}}(N^2)$$

For the factorized determinants one of the two determinants is obviously unaffected by the change so that it cancels from the ratio R .

Therefore, only one determinant of size $N/2$ is involved in each calculation of R and update of the inverse matrix. The scaling of each transition then becomes:

$$O_R(N/2) + O_{\text{inverse}}(N^2/4)$$

and the time scaling when the transitions for all N particles are put together:

$$O_R(N^2/2) + O_{\text{inverse}}(N^3/4)$$

which gives the same reduction as in the case of moving all particles at once.

Computing the ratios discussed above requires that we maintain the inverse of the Slater matrix evaluated at the current position. Each time a trial position is accepted, the row number i of the Slater matrix changes and updating its inverse has to be carried out. Getting the inverse of an $N \times N$ matrix by Gaussian elimination has a complexity of order of $\mathcal{O}(N^3)$ operations, a luxury that we cannot afford for each time a particle move is accepted. We will use the expression

$$d_{kj}^{-1}(\mathbf{x}^{\text{new}}) = \begin{cases} d_{kj}^{-1}(\mathbf{x}^{\text{old}}) - \frac{d_{ki}^{-1}(\mathbf{x}^{\text{old}})}{R} \sum_{l=1}^N d_{il}(\mathbf{x}^{\text{new}}) d_{lj}^{-1}(\mathbf{x}^{\text{old}}) & \text{if } j \neq i \\ \frac{d_{ki}^{-1}(\mathbf{x}^{\text{old}})}{R} \sum_{l=1}^N d_{il}(\mathbf{x}^{\text{old}}) d_{lj}^{-1}(\mathbf{x}^{\text{old}}) & \text{if } j = i \end{cases}$$

This equation scales as $\mathcal{O}(N^2)$. The evaluation of the determinant of an $N \times N$ matrix by standard Gaussian elimination requires $\mathcal{O}(N^3)$ calculations. As there are Nd independent coordinates we need to evaluate Nd Slater determinants for the gradient (quantum force) and Nd for the Laplacian (kinetic energy). With the updating algorithm we need only to invert the Slater determinant matrix once. This can be done by standard LU decomposition methods.

Expectation value of the kinetic energy.

The expectation value of the kinetic energy expressed in atomic units for electron i is

$$\begin{aligned} \langle \hat{K}_i \rangle &= -\frac{1}{2} \frac{\langle \Psi | \nabla_i^2 | \Psi \rangle}{\langle \Psi | \Psi \rangle}, \\ K_i &= -\frac{1}{2} \frac{\nabla_i^2 \Psi}{\Psi}. \end{aligned} \quad (6.16)$$

$$\begin{aligned} \frac{\nabla^2 \Psi}{\Psi} &= \frac{\nabla^2(\Psi_D \Psi_C)}{\Psi_D \Psi_C} = \frac{\nabla \cdot [\nabla(\Psi_D \Psi_C)]}{\Psi_D \Psi_C} = \frac{\nabla \cdot [\Psi_C \nabla \Psi_D + \Psi_D \nabla \Psi_C]}{\Psi_D \Psi_C} \\ &= \frac{\nabla \Psi_C \cdot \nabla \Psi_D + \Psi_C \nabla^2 \Psi_D + \nabla \Psi_D \cdot \nabla \Psi_C + \Psi_D \nabla^2 \Psi_C}{\Psi_D \Psi_C} \end{aligned} \quad (6.17)$$

$$\frac{\nabla^2 \Psi}{\Psi} = \frac{\nabla^2 \Psi_D}{\Psi_D} + \frac{\nabla^2 \Psi_C}{\Psi_C} + 2 \frac{\nabla \Psi_D}{\Psi_D} \cdot \frac{\nabla \Psi_C}{\Psi_C} \quad (6.18)$$

The second derivative of the Jastrow factor divided by the Jastrow factor (the way it enters the kinetic energy) is

$$\left[\frac{\nabla^2 \Psi_C}{\Psi_C} \right]_x = 2 \sum_{k=1}^N \sum_{i=1}^{k-1} \frac{\partial^2 g_{ik}}{\partial x_k^2} + \sum_{k=1}^N \left(\sum_{i=1}^{k-1} \frac{\partial g_{ik}}{\partial x_k} - \sum_{i=k+1}^N \frac{\partial g_{ki}}{\partial x_i} \right)^2$$

But we have a simple form for the function, namely

$$\Psi_C = \prod_{i < j} \exp f(r_{ij}) = \exp \left\{ \sum_{i < j} \frac{a r_{ij}}{1 + \beta r_{ij}} \right\},$$

and it is easy to see that for particle k we have

$$\frac{\nabla_k^2 \Psi_C}{\Psi_C} = \sum_{ij \neq k} \frac{(\mathbf{r}_k - \mathbf{r}_i)(\mathbf{r}_k - \mathbf{r}_j)}{r_{ki} r_{kj}} f'(r_{ki}) f'(r_{kj}) + \sum_{j \neq k} \left(f''(r_{kj}) + \frac{2}{r_{kj}} f'(r_{kj}) \right)$$

Using

$$f(r_{ij}) = \frac{ar_{ij}}{1 + \beta r_{ij}},$$

and $g'(r_{kj}) = dg(r_{kj})/dr_{kj}$ and $g''(r_{kj}) = d^2g(r_{kj})/dr_{kj}^2$ we find that for particle k we have

$$\frac{\nabla_k^2 \Psi_C}{\Psi_C} = \sum_{ij \neq k} \frac{(\mathbf{r}_k - \mathbf{r}_i)(\mathbf{r}_k - \mathbf{r}_j)}{r_{ki} r_{kj}} \frac{a}{(1 + \beta r_{ki})^2} \frac{a}{(1 + \beta r_{kj})^2} + \sum_{j \neq k} \left(\frac{2a}{r_{kj}(1 + \beta r_{kj})^2} - \frac{2a\beta}{(1 + \beta r_{kj})^3} \right)$$

The gradient and Laplacian can be calculated as follows:

$$\frac{\nabla_i |\hat{D}(\mathbf{r})|}{|\hat{D}(\mathbf{r})|} = \sum_{j=1}^N \nabla_i d_{ij}(\mathbf{r}) d_{ji}^{-1}(\mathbf{r}) = \sum_{j=1}^N \nabla_i \phi_j(\mathbf{r}_i) d_{ji}^{-1}(\mathbf{r})$$

and

$$\frac{\nabla_i^2 |\hat{D}(\mathbf{r})|}{|\hat{D}(\mathbf{r})|} = \sum_{j=1}^N \nabla_i^2 d_{ij}(\mathbf{r}) d_{ji}^{-1}(\mathbf{r}) = \sum_{j=1}^N \nabla_i^2 \phi_j(\mathbf{r}_i) d_{ji}^{-1}(\mathbf{r})$$

The gradient for the determinant is

$$\frac{\nabla_i |\hat{D}(\mathbf{r})|}{|\hat{D}(\mathbf{r})|} = \sum_{j=1}^N \nabla_i d_{ij}(\mathbf{r}) d_{ji}^{-1}(\mathbf{r}) = \sum_{j=1}^N \nabla_i \phi_j(\mathbf{r}_i) d_{ji}^{-1}(\mathbf{r}).$$

We have

$$\Psi_C = \prod_{i < j} g(r_{ij}) = \exp \left\{ \sum_{i < j} \frac{ar_{ij}}{1 + \beta r_{ij}} \right\},$$

the gradient needed for the quantum force and local energy is easy to compute. We get for particle k

$$\frac{\nabla_k \Psi_C}{\Psi_C} = \sum_{j \neq k} \frac{\mathbf{r}_{kj}}{r_{kj}} \frac{a}{(1 + \beta r_{kj})^2},$$

which is rather easy to code. Remember to sum over all particles when you compute the local energy.

We need to compute the ratio between wave functions, in particular for the Slater determinants.

$$R = \sum_{j=1}^N d_{ij}(\mathbf{r}^{\text{new}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}}) = \sum_{j=1}^N \phi_j(\mathbf{r}_i^{\text{new}}) d_{ji}^{-1}(\mathbf{r}^{\text{old}})$$

What this means is that in order to get the ratio when only the i -th particle has been moved, we only need to calculate the dot product of the vector $(\phi_1(\mathbf{r}_i^{\text{new}}), \dots, \phi_N(\mathbf{r}_i^{\text{new}}))$ of single particle wave functions evaluated at this new position with the i -th column of the inverse matrix \hat{D}^{-1} evaluated at the original position. Such an operation has a time scaling of $O(N)$. The only extra thing we need to do is to maintain the inverse matrix $\hat{D}^{-1}(\mathbf{x}^{\text{old}})$.

Part IV
Machine Learning

Part V
Quantum Computing

