Variational Monte Carlo methods

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Quantum Monte Carlo Motivation

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.

Quantum Monte Carlo Motivation

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.

Quantum Monte Carlo Motivation

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.

Quantum Monte Carlo Motivation

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.

Quantum Monte Carlo Motivation

- Construct first a trial wave function $\psi_T(R, \boldsymbol{\alpha})$, for a many-body system consisting of N particles located at positions

$R = (R_1, \ldots, R_N)$. The trial wave function depends on α variational parameters $\boldsymbol{\alpha} = (\alpha_1, \ldots, \alpha_M)$.

- Then we evaluate the expectation value of the hamiltonian H

$$E[H] = \langle H \rangle = \frac{\int dR \Psi_T^*(R, \boldsymbol{\alpha}) H(R) \Psi_T(R, \boldsymbol{\alpha})}{\int dR \Psi_T^*(R, \boldsymbol{\alpha}) \Psi_T(R, \boldsymbol{\alpha})}.$$

- Thereafter we vary α according to some minimization algorithm and return to the first step.

Quantum Monte Carlo Motivation

Basic steps

Choose a trial wave function $\psi_T(R)$.

$$P(R) = \frac{|\psi_T(R)|^2}{\int |\psi_T(R)|^2\, dR}.$$

This is our new probability distribution function (PDF). The approximation to the expectation value of the Hamiltonian is now

$$E[H(\alpha)] = \frac{\int dR \Psi_T^*(R,\alpha) H(R) \Psi_T(R,\alpha)}{\int dR \Psi_T^*(R,\alpha)\Psi_T(R,\alpha)}.$$

Quantum Monte Carlo Motivation

Define a new quantity

$$E_L(R,\boldsymbol{\alpha}) = \frac{1}{\psi_T(R,\boldsymbol{\alpha})} H \psi_T(R,\boldsymbol{\alpha}),$$

called the local energy, which, together with our trial PDF yields

$$E[H(\boldsymbol{\alpha})] = \int P(R) E_L(R) dR \approx \frac{1}{N} \sum_{i=1}^{N} P(R_i,\boldsymbol{\alpha}) E_L(R_i,\boldsymbol{\alpha})$$

with N being the number of Monte Carlo samples.

Quantum Monte Carlo

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 $R_p = R + r * step$ 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

$$-\frac{\hbar^2}{2m}\frac{\partial^2 u(r)}{\partial r^2} - \left(\frac{ke^2}{r} - \frac{\hbar^2 l(l+1)}{2mr^2}\right) u(r) = Eu(r),$$

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

$$u_T^\alpha(\rho) = \alpha\rho e^{-\alpha\rho}.$$

Quantum Monte Carlo: hydrogen atom

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	$\sigma/\sqrt{N}$
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

Quantum Monte Carlo: hydrogen atom

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 hamiltonan 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

$$\langle H^n \rangle = \frac{\int dR \Psi_T^*(R) H^n(R) \Psi_T(R)}{\int dR \Psi_T^*(R)\Psi_T(R)} = \text{constant} \times \frac{\int dR \Psi_T^*(R)\Psi_T(R)}{\int dR \Psi_T^*(R)\Psi_T(R)} = \text{co}$$

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.

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 $r_{12} = |r_1 - r_2|$.

Quantum Monte Carlo: the helium atom

The hamiltonian becomes then

$$\hat{H} = -\frac{\hbar^2 \nabla_1^2}{2m} - \frac{\hbar^2 \nabla_2^2}{2m} - \frac{2ke^2}{r_1} - \frac{2ke^2}{r_2} + \frac{ke^2}{r_{12}},$$

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.

Quantum Monte Carlo: the helium atom

Choice of trial wave function for Helium: Assume $r_1 \to 0$.

$$E_L(R) = \frac{1}{\psi_T(R)} H \psi_T(R) = \frac{1}{\psi_T(R)} \left(-\frac{1}{2} \nabla_1^2 - \frac{Z}{r_1} \right) \psi_T(R) + \text{finite terms.}$$

$$E_L(R) = \frac{1}{\mathcal{R}_T(r_1)} \left(-\frac{1}{2} \frac{d^2}{dr_1^2} - \frac{1}{r_1} \frac{d}{dr_1} - \frac{Z}{r_1} \right) \mathcal{R}_T(r_1) + \text{finite terms}$$

For small values of r_1, the terms which dominate are

$$\lim_{r_1 \to 0} E_L(R) = \frac{1}{\mathcal{R}_T(r_1)} \left(-\frac{1}{r_1} \frac{d}{dr_1} - \frac{Z}{r_1} \right) \mathcal{R}_T(r_1),$$

since the second derivative does not diverge due to the finiteness of Ψ at the origin.

Quantum Monte Carlo: the helium atom

This results in

$$\frac{1}{\mathcal{R}_T(r_1)} \frac{d\mathcal{R}_T(r_1)}{dr_1} = -Z,$$

and

$$\mathcal{R}_T(r_1) \propto e^{-Zr_1}.$$

A similar condition applies to electron 2 as well. For orbital momenta $l > 0$ we have

$$\frac{1}{\mathcal{R}_T(r)} \frac{d\mathcal{R}_T(r)}{dr} = -\frac{Z}{l+1}.$$

Similarly, studying the case $r_{12} \to 0$ 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

$$\psi_T(R) = \phi(r_1)\phi(r_2)\ldots\phi(r_N) \prod f(r_{ij}),$$

The first attempt at solving the helium atom

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

$$\langle E_{L1} \rangle = \alpha^2 - 2\alpha \left(Z - \frac{5}{16} \right)$$

The first attempt at solving the Helium atom

With closed form formulae we can speed up the computation of the correlation. In our case we write it as

$$\Psi_C = \exp \left\{ \sum_{i<j} \frac{a r_{ij}}{1 + \beta r_{ij}} \right\},$$

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*.

The first attempt at solving the Helium atom

We can test this by computing the local energy for our helium wave function

$$\psi_T(r_1, r_2) = \exp\left(-\alpha(r_1 + r_2)\right)\exp\left(\frac{r_{12}}{2(1 + \beta r_{12})}\right),$$

with α and β as variational parameters.
The local energy is for this case

$$E_{L2} = E_{L1} + \frac{1}{2(1 + \beta r_{12})^2}\left\{\frac{\alpha(r_1 + r_2)}{r_{12}}(1 - \frac{r_1 r_2}{r_1 r_2}) - \frac{1}{2(1 + \beta r_{12})^2} - \frac{2}{r_{12}} + \right.$$

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.

The first attempt at solving the Helium atom

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
http://docs.sympy.org/latest/index.html. 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

$$\phi_{2s}(r) = (Zr - 2)\exp{-(\frac{1}{2}Zr)},$$

with $r^2 = x^2 + y^2 + z^2$.

```
from sympy import symbols, diff, exp, sqrt
x, y, z, Z = symbols('x y z Z')
r = sqrt(x*x + y*y + z*z)
r
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

The first attempt at solving the Helium atom

We can improve our output by factorizing and substituting expressions

```
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))
```

The first attempt at solving the Helium atom

We can in turn look at second derivatives

```
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)
# Collect the Z values
(diff(diff(phi, x), x) + diff(diff(phi, y), y) +diff(diff(phi, z), z))
# Factorize also the r**2 terms
(diff(diff(phi, x), x) + diff(diff(phi, y), y) + diff(diff(phi, z), z)
print printing.ccode((diff(diff(phi, x), x) + diff(diff(phi, y), y) +
```

With some practice this allows one to be able to check one's own calculation and translate automatically into code lines.

The first attempt at solving the Helium atom

The c++ code with a VMC Solver class, main program first

```
#include "vmcsolver.h"
#include <iostream>
using namespace std;

int main()
{
    VMCSolver *solver = new VMCSolver();
    solver->runMonteCarloIntegration();
    return 0;
}
```

The first attempt at solving the Helium atom

The c++ code with a VMC Solver class, the VMCSolver header file

```
#ifndef VMCSOLVER_H
#define VMCSOLVER_H
#include <armadillo>
using namespace arma;
class VMCSolver
{
public:
    VMCSolver();
    void runMonteCarloIntegration();

private:
    double waveFunction(const mat &r);
    double localEnergy(const mat &r);
    int nDimensions;
    int charge;
    double stepLength;
    int nParticles;
    double h;
    double h2;
    long idum;
    double alpha;
    int nCycles;
    mat rOld;
    mat rNew;
};
#endif // VMCSOLVER_H
```

Slide 1

The c++ code with a VMC Solver class, VMCSolver codes, initialize

```cpp
#include "vmcsolver.h"
#include "lib.h"
#include <armadillo>
#include <iostream>
using namespace arma;
using namespace std;

VMCSolver::VMCSolver() :
    nDimensions(3),
    charge(2),
    stepLength(1.0),
    nParticles(2),
    h(0.001),
    h2(1000000),
    idum(-1),
    alpha(0.5*charge),
    nCycles(1000000)
{
}
```

Slide 2

The c++ code with a VMC Solver class, VMCSolver codes

```cpp
void VMCSolver::runMonteCarloIntegration()
{
    rOld = zeros<mat>(nParticles, nDimensions);
    rNew = zeros<mat>(nParticles, nDimensions);
    double waveFunctionOld = 0;
    double waveFunctionNew = 0;
    double energySum = 0;
    double energySquaredSum = 0;
    double deltaE;
    // initial trial positions
    for(int i = 0; i < nParticles; i++) {
        for(int j = 0; j < nDimensions; j++) {
            rOld(i,j) = stepLength * (ran2(&idum) - 0.5);
        }
    }
    rNew = rOld;
    // loop over Monte Carlo cycles
    for(int cycle = 0; cycle < nCycles; cycle++) {
        // Store the current value of the wave function
        waveFunctionOld = waveFunction(rOld);
        // New position to test
        for(int i = 0; i < nParticles; i++) {
            for(int j = 0; j < nDimensions; j++) {
                rNew(i,j) = rOld(i,j) + stepLength*(ran2(&idum) - 0.5)
            }
            // Recalculate the value of the wave function
```

Slide 3

The c++ code with a VMC Solver class, VMCSolver codes

```cpp
double VMCSolver::localEnergy(const mat &r)
{
    mat rPlus = zeros<mat>(nParticles, nDimensions);
    mat rMinus = zeros<mat>(nParticles, nDimensions);
    rPlus = rMinus = r;
    double waveFunctionMinus = 0;
    double waveFunctionPlus = 0;
    double waveFunctionCurrent = waveFunction(r);
    // Kinetic energy, brute force derivations
    double kineticEnergy = 0;
    for(int i = 0; i < nParticles; i++) {
        for(int j = 0; j < nDimensions; j++) {
            rPlus(i,j) += h;
            rMinus(i,j) -= h;
            waveFunctionMinus = waveFunction(rMinus);
            waveFunctionPlus = waveFunction(rPlus);
            kineticEnergy -= (waveFunctionMinus + waveFunctionPlus - 2
            rPlus(i,j) = r(i,j);
            rMinus(i,j) = r(i,j);
        }
    }
    kineticEnergy = 0.5 * h2 * kineticEnergy / waveFunctionCurrent;
    // Potential energy
    double potentialEnergy = 0;
    double rSingleParticle = 0;
    for(int i = 0; i < nParticles; i++) {
```

Slide 4

The c++ code with a VMC Solver class, VMCSolver codes

```cpp
double VMCSolver::waveFunction(const mat &r)
{
    double argument = 0;
    for(int i = 0; i < nParticles; i++) {
        double rSingleParticle = 0;
        for(int j = 0; j < nDimensions; j++) {
            rSingleParticle += r(i,j) * r(i,j);
        }
        argument += sqrt(rSingleParticle);
    }
    return exp(-argument * alpha);
}
```

Slide 5

The c++ code with a VMC Solver class, the VMCSolver header file

```cpp
#include <armadillo>
#include <iostream>
using namespace arma;
using namespace std;
double ran2(long *);

class VMCSolver
{
public:
    VMCSolver();
    void runMonteCarloIntegration();

private:
    double waveFunction(const mat &r);
    double localEnergy(const mat &r);
    int nDimensions;
    int charge;
    double stepLength;
    int nParticles;
    double h;
    double h2;
    long idum;
    double alpha;
    int nCycles;
    mat rOld;
    mat rNew;
```

Slide 6

Exercises for first lab session, Thursday 22

- If you have never used git, Qt, armadillo etc, get familiar with them, see the guides at the official UiO website of the course.
- Study the simple program in the folder https://github.com/CompPhysics/ComputationalPhysics2/tree/gh-pages/doc/pub/vmc/programs/
- Implement the closed form expression for the local energy and the so-called quantum force
- Convince yourself that the closed form expressions are correct. Check both wave functions
- Implement the closed form expression for the local energy and compare with a code where the second derivatives are computed numerically.

The Metropolis algorithm

The Metropolis algorithm , see http://scitation.aip.org/content/aip/journal/jcp/21/6/10.1063/1.1699114 (see also the FYS3150 lectures http://www.uio.no/studier/emner/matnat/fys/FYS3150/h14/index.html) 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 $\mathcal{P}_i^{(n)}$ 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 $T_{i \to j}$.
- Accept the new state j with probability $A_{i \to j}$ and use it as the next sample. With probability $1 - A_{i \to j}$ the move is rejected and the original state i is used again as a sample.

The Metropolis algorithm

We wish to derive the required properties of T and A such that $\mathcal{P}_i^{(n \to \infty)} \to p_i$ 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 Metropolis algorithm

The dynamical equation for $\mathcal{P}_i^{(n)}$ 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:

$$\mathcal{P}_i^{(n)} = \sum_j \left[\mathcal{P}_j^{(n-1)} T_{j \to i} A_{j \to i} + \mathcal{P}_i^{(n-1)} T_{i \to j} \left(1 - A_{i \to j} \right) \right].$$

Since the probability of making some transition must be 1, $\sum_j T_{i \to j} = 1$, and the above equation becomes

$$\mathcal{P}_i^{(n)} = \mathcal{P}_i^{(n-1)} + \sum_j \left[\mathcal{P}_j^{(n-1)} T_{j \to i} A_{j \to i} - \mathcal{P}_i^{(n-1)} T_{i \to j} A_{i \to j} \right].$$

The Metropolis algorithm

For large n we require that $\mathcal{P}_i^{(n \to \infty)} = p_i$, the desired probability distribution. Taking this limit, gives the balance requirement

$$\sum_j \left[p_j T_{j \to i} A_{j \to i} - p_i T_{i \to j} A_{i \to j} \right] = 0.$$

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

$$\frac{A_{j \to i}}{A_{i \to j}} = \frac{p_i T_{i \to j}}{p_j T_{j \to i}}.$$

The Metropolis algorithm

The Metropolis choice is to maximize the A values, that is

$$A_{j \to i} = \min \left(1, \frac{p_i T_{i \to j}}{p_j T_{j \to i}} \right).$$

Other choices are possible, but they all correspond to multilplying $A_{i \to j}$ and $A_{j \to i}$ by the same constant smaller than unity.[a]

[a] The penalty function method uses just such a factor to compensate for p_i that are evaluated stochastically and are therefore noisy.

The Metropolis algorithm

Having chosen the acceptance probabilities, we have guaranteed that if the $\mathcal{P}_i^{(n)}$ 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

$$\mathcal{P}_i^{(n)} = \sum_j M_{ij} \mathcal{P}_j^{(n-1)}$$

with the matrix M given by

$$M_{ij} = \delta_{ij} \left[1 - \sum_k T_{i \to k} A_{i \to k} \right] + T_{j \to i} A_{j \to i}.$$

Summing over i shows that $\sum_i M_{ij} = 1$, and since $\sum_k T_{i \to k} = 1$, and $A_{i \to k} \le 1$, the elements of the matrix satisfy $M_{ij} \ge 0$. The matrix M is therefore a stochastic matrix.

The Metropolis algorithm

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.

Why blocking?

Statistical analysis

- Monte Carlo simulations can be treated as *computer experiments*
- The results can be analysed with the same statistical tools as we would use analysing experimental data.
- As in all experiments, we are looking for expectation values and an estimate of how accurate they are, i.e., possible sources for errors.

Why blocking?

Statistical analysis

- As in other experiments, Monte Carlo experiments have two classes of errors:
  - Statistical errors
  - Systematical errors
- Statistical errors can be estimated using standard tools from statistics
- Systematical errors are method specific and must be treated differently from case to case. (In VMC a common source is the step length or time step in importance sampling)

Statistics and blocking

The *probability distribution function (PDF)* is a function $p(x)$ on the domain which, in the discrete case, gives us the probability or relative frequency with which these values of X occur:

$$p(x) = \text{prob}(X = x)$$

In the continuous case, the PDF does not directly depict the actual probability. Instead we define the probability for the stochastic variable to assume any value on an infinitesimal interval around x to be $p(x)dx$. The continuous function $p(x)$ then gives us the *density* of the probability rather than the probability itself. The probability for a stochastic variable to assume any value on a non-infinitesimal interval $[a, b]$ is then just the integral:

$$\text{prob}(a \leq X \leq b) = \int_a^b p(x)dx$$

Qualitatively speaking, a stochastic variable represents the values of numbers chosen as if by chance from some specified PDF so that

Statistics and blocking

Also of interest to us is the *cumulative probability distribution function (CDF)*, $P(x)$, which is just the probability for a stochastic variable X to assume any value less than x:

$$P(x) = \text{Prob}(X \leq x) = \int_{-\infty}^x p(x')dx'$$

The relation between a CDF and its corresponding PDF is then:

$$p(x) = \frac{d}{dx}P(x)$$

Statistics and blocking

A particularly useful class of special expectation values are the *moments*. The n-th moment of the PDF p is defined as follows:

$$\langle x^n \rangle \equiv \int x^n p(x)\, dx$$

The zero-th moment $\langle 1 \rangle$ is just the normalization condition of p. The first moment, $\langle x \rangle$, is called the *mean* of p and often denoted by the letter μ:

$$\langle x \rangle = \mu \equiv \int x p(x)\, dx$$

A special version of the moments is the set of *central moments*, the n-th central moment defined as:

$$\langle (x - \langle x \rangle)^n \rangle \equiv \int (x - \langle x \rangle)^n p(x)\, dx$$

The zero-th and first central moments are both trivial, equal 1 and 0, respectively. But the second central moment, known as the *variance* of p, is of particular interest. For the stochastic variable X, the variance is denoted as σ_X^2 or $\text{var}(X)$:

$$
\begin{aligned}
\sigma_X^2 \;=\; \text{var}(X) &= \langle (x - \langle x \rangle)^2 \rangle = \int (x - \langle x \rangle)^2 p(x)\, dx \quad (1)\\
&= \int \left(x^2 - 2x\langle x \rangle + \langle x \rangle^2 \right) p(x)\, dx \quad (2)\\
&= \langle x^2 \rangle - 2\langle x \rangle\langle x \rangle + \langle x \rangle^2 \quad (3)\\
&= \langle x^2 \rangle - \langle x \rangle^2 \quad (4)
\end{aligned}
$$

The square root of the variance, $\sigma = \sqrt{\langle (x - \langle x \rangle)^2 \rangle}$ is called the

---

Another important quantity is the so called covariance, a variant of the above defined variance. Consider again the set $\{X_i\}$ of n stochastic variables (not necessarily uncorrelated) with the multivariate PDF $P(x_1, \ldots, x_n)$. The *covariance* of two of the stochastic variables, X_i and X_j, is defined as follows:

$$
\begin{aligned}
\text{cov}(X_i, X_j) &\equiv \langle (x_i - \langle x_i \rangle)(x_j - \langle x_j \rangle) \rangle\\
&= \int \cdots \int (x_i - \langle x_i \rangle)(x_j - \langle x_j \rangle)\, P(x_1, \ldots, x_n)\, dx_1 \ldots dx_n \quad (5)
\end{aligned}
$$

with

$$\langle x_i \rangle = \int \cdots \int x_i\, P(x_1, \ldots, x_n)\, dx_1 \ldots dx_n$$

If we consider the above covariance as a matrix $C_{ij} = \text{cov}(X_i, X_j)$, then the diagonal elements are just the familiar variances, $C_{ii} = \text{cov}(X_i, X_i) = \text{var}(X_i)$. It turns out that all the off-diagonal elements are zero if the stochastic variables are uncorrelated. This is easy to show, keeping in mind the linearity of the expectation value. Consider the stochastic variables X_i and X_j, $(i \neq j)$:

$$
\begin{aligned}
\text{cov}(X_i, X_j) &= \langle (x_i - \langle x_i \rangle)(x_j - \langle x_j \rangle) \rangle \quad (6)\\
&= \langle x_i x_j - x_i \langle x_j \rangle - \langle x_i \rangle x_j + \langle x_i \rangle \langle x_j \rangle \rangle \quad (7)\\
&= \langle x_i x_j \rangle - \langle x_i \langle x_j \rangle \rangle - \langle \langle x_i \rangle x_j \rangle + \langle \langle x_i \rangle \langle x_j \rangle \rangle \quad (8)\\
&= \langle x_i x_j \rangle - \langle x_i \rangle \langle x_j \rangle - \langle x_i \rangle \langle x_j \rangle + \langle x_i \rangle \langle x_j \rangle \quad (9)\\
&= \langle x_i x_j \rangle - \langle x_i \rangle \langle x_j \rangle \quad (10)
\end{aligned}
$$

If X_i and X_j are independent, we get $\langle x_i x_j \rangle = \langle x_i \rangle \langle x_j \rangle$, resulting in $\text{cov}(X_i, X_j) = 0$ $(i \neq j)$.

Also useful for us is the covariance of linear combinations of stochastic variables. Let $\{X_i\}$ and $\{Y_i\}$ be two sets of stochastic variables. Let also $\{a_i\}$ and $\{b_i\}$ be two sets of scalars. Consider the linear combination:

$$U = \sum_i a_i X_i \qquad V = \sum_j b_j Y_j$$

By the linearity of the expectation value

$$\text{cov}(U, V) = \sum_{i,j} a_i b_j \text{cov}(X_i, Y_j)$$

Now, since the variance is just $\text{var}(X_i) = \text{cov}(X_i, X_i)$, we get the variance of the linear combination $U = \sum_i a_i X_i$:

$$\text{var}(U) = \sum_{i,j} a_i a_j \text{cov}(X_i, X_j) \quad (11)$$

And in the special case when the stochastic variables are uncorrelated, the off-diagonal elements of the covariance are as we know zero, resulting in:

$$\text{var}(U) = \sum_i a_i^2 \text{cov}(X_i, X_i) = \sum_i a_i^2 \text{var}(X_i)$$

$$\text{var}\left(\sum_i a_i X_i\right) = \sum_i a_i^2 \text{var}(X_i)$$

which will become very useful in our study of the error in the mean value of a set of measurements.

A *stochastic process* is a process that produces sequentially a chain of values:

$$\{x_1, x_2, \ldots x_k, \ldots\}.$$

We will call these values our *measurements* and the entire set as our measured *sample*. The action of measuring all the elements of a sample we will call a stochastic *experiment* since, operationally, they are often associated with results of empirical observation of some physical or mathematical phenomena; precisely an experiment. We assume that these values are distributed according to some PDF $p_X(x)$, where X is just the formal symbol for the stochastic variable whose PDF is $p_X(x)$. Instead of trying to determine the full distribution p we are often only interested in finding the few lowest moments, like the mean μ_X and the variance σ_X.

In practical situations a sample is always of finite size. Let that size be n. The expectation value of a sample, the *sample mean*, is then defined as follows:

$$\bar{x}_n \equiv \frac{1}{n}\sum_{k=1}^{n} x_k$$

The *sample variance* is:

$$\mathrm{var}(x) \equiv \frac{1}{n}\sum_{k=1}^{n}(x_k - \bar{x}_n)^2$$

its square root being the *standard deviation of the sample*. The *sample covariance* is:

$$\mathrm{cov}(x) \equiv \frac{1}{n}\sum_{kl}(x_k - \bar{x}_n)(x_l - \bar{x}_n)$$

Note that the sample variance is the sample covariance without the cross terms. In a similar manner as the covariance in Eq. (5) is a measure of the correlation between two stochastic variables, the above defined sample covariance is a measure of the sequential correlation between succeeding measurements of a sample. These quantities, being known experimental values, differ significantly from and must not be confused with the similarly named quantities for stochastic variables, mean μ_X, variance $\mathrm{var}(X)$ and covariance $\mathrm{cov}(X, Y)$.

The law of large numbers states that as the size of our sample grows to infinity, the sample mean approaches the true mean μ_X of the chosen PDF:

$$\lim_{n\to\infty} \bar{x}_n = \mu_X$$

The sample mean $\bar{x}_n$ works therefore as an estimate of the true mean μ_X.
What we need to find out is how good an approximation $\bar{x}_n$ is to μ_X. In any stochastic measurement, an estimated mean is of no use to us without a measure of its error. A quantity that tells us how well we can reproduce it in another experiment. We are therefore interested in the PDF of the sample mean itself. Its standard deviation will be a measure of the spread of sample means, and we will simply call it the *error* of the sample mean, or just sample error, and denote it by err_X. In practice, we will only be able to produce an *estimate* of the sample error since the exact value would require the knowledge of the true PDFs behind, which we usually do not have.

The straight forward brute force way of estimating the sample error is simply by producing a number of samples, and treating the mean of each as a measurement. The standard deviation of these means will then be an estimate of the original sample error. If we are unable to produce more than one sample, we can split it up sequentially into smaller ones, treating each in the same way as above. This procedure is known as *blocking* and will be given more attention shortly. At this point it is worth while exploring more indirect methods of estimation that will help us understand some important underlying principles of correlational effects.

Let us first take a look at what happens to the sample error as the size of the sample grows. In a sample, each of the measurements x_i can be associated with its own stochastic variable X_i. The stochastic variable $\overline{X}_n$ for the sample mean $\bar{x}_n$ is then just a linear combination, already familiar to us:

$$\overline{X}_n = \frac{1}{n}\sum_{i=1}^{n} X_i$$

All the coefficients are just equal $1/n$. The PDF of $\overline{X}_n$, denoted by $p_{\overline{X}_n}(x)$ is the desired PDF of the sample means.

The probability density of obtaining a sample mean $\bar{x}_n$ is the product of probabilities of obtaining arbitrary values $x_1, x_2, \ldots, x_n$ with the constraint that the mean of the set $\{x_i\}$ is $\bar{x}_n$:

$$p_{\overline{X}_n}(x) = \int p_X(x_1)\cdots\int p_X(x_n)\,\delta\!\left(x - \frac{x_1 + x_2 + \cdots + x_n}{n}\right)dx_n\cdots dx_1$$

And in particular we are interested in its variance $\mathrm{var}(\overline{X}_n)$.

It is generally not possible to express $p_{\overline{X}_n}(x)$ in a closed form given an arbitrary PDF p_X and a number n. But for the limit $n \to \infty$ it is possible to make an approximation. The very important result is called *the central limit theorem*. It tells us that as n goes to infinity, $p_{\overline{X}_n}(x)$ approaches a Gaussian distribution whose mean and variance equal the true mean and variance, μ_X and σ_X^2, respectively:

$$\lim_{n \to \infty} p_{\overline{X}_n}(x) = \left(\frac{n}{2\pi \mathrm{var}(X)}\right)^{1/2} e^{-\frac{n(x-\bar{x}_n)^2}{2\mathrm{var}(X)}} \qquad (12)$$

The desired variance $\mathrm{var}(\overline{X}_n)$, i.e. the sample error squared err_X^2, is given by:

$$\mathrm{err}_X^2 = \mathrm{var}(\overline{X}_n) = \frac{1}{n^2}\sum_{ij} \mathrm{cov}(X_i, X_j) \qquad (13)$$

We see now that in order to calculate the exact error of the sample with the above expression, we would need the true means μ_{X_i} of the stochastic variables X_i. To calculate these requires that we know the true multivariate PDF of all the X_i. But this PDF is unknown to us, we have only got the measurements of one sample. The best we can do is to let the sample itself be an estimate of the PDF of each of the X_i, estimating all properties of X_i through the measurements of the sample.

Our estimate of μ_{X_i} is then the sample mean $\bar{x}$ itself, in accordance with the the central limit theorem:

$$\mu_{X_i} = \langle x_i \rangle \approx \frac{1}{n}\sum_{k=1}^n x_k = \bar{x}$$

Using $\bar{x}$ in place of μ_{X_i} we can give an *estimate* of the covariance in Eq. (13)

$$\mathrm{cov}(X_i, X_j) = \langle(x_i - \langle x_i\rangle)(x_j - \langle x_j\rangle)\rangle \approx \langle(x_i - \bar{x})(x_j - \bar{x})\rangle,$$

resulting in

$$\frac{1}{n}\sum_l\left(\frac{1}{n}\sum_k(x_k - \bar{x}_n)(x_l - \bar{x}_n)\right) = \frac{1}{n}\frac{1}{n}\sum_{kl}(x_k - \bar{x}_n)(x_l - \bar{x}_n) = \frac{1}{n}\mathrm{cov}(x)$$

By the same procedure we can use the sample variance as an estimate of the variance of any of the stochastic variables X_i

$$\mathrm{var}(X_i) = \langle x_i - \langle x_i\rangle\rangle \approx \langle x_i - \bar{x}_n\rangle,$$

which is approximated as

$$\mathrm{var}(X_i) \approx \frac{1}{n}\sum_{k=1}^n(x_k - \bar{x}_n) = \mathrm{var}(x) \qquad (14)$$

Now we can calculate an estimate of the error err_X of the sample mean $\bar{x}_n$:

$$\begin{aligned}\mathrm{err}_X^2 &= \frac{1}{n^2}\sum_{ij}\mathrm{cov}(X_i, X_j)\\ &\approx \frac{1}{n^2}\sum_{ij}\frac{1}{n}\mathrm{cov}(x) = \frac{1}{n^2}n^2\frac{1}{n}\mathrm{cov}(x)\\ &\approx \frac{1}{n}\mathrm{cov}(x) \qquad (15)\end{aligned}$$

In the special case that the measurements of the sample are uncorrelated (equivalently the stochastic variables X_i are uncorrelated) we have that the off-diagonal elements of the covariance are zero. This gives the following estimate of the sample error:

$$\mathrm{err}_X^2 = \frac{1}{n^2}\sum_{ij}\mathrm{cov}(X_i, X_j) = \frac{1}{n^2}\sum_i\mathrm{var}(X_i),$$

resulting in

$$\mathrm{err}_X^2 \approx \frac{1}{n^2}\sum_i\mathrm{var}(x) = \frac{1}{n}\mathrm{var}(x) \qquad (16)$$

where in the second step we have used Eq. (14). The error of the sample is then just its standard deviation divided by the square root of the number of measurements the sample contains. This is a very useful formula which is easy to compute. It acts as a first approximation to the error, but in numerical experiments, we cannot overlook the always present correlations.

For computational purposes one usually splits up the estimate of err_X^2, given by Eq. (15), into two parts

$$\mathrm{err}_X^2 = \frac{1}{n}\mathrm{var}(x) + \frac{1}{n}(\mathrm{cov}(x) - \mathrm{var}(x)),$$

which equals

$$\frac{1}{n^2}\sum_{k=1}^n(x_k - \bar{x}_n)^2 + \frac{2}{n^2}\sum_{k<l}(x_k - \bar{x}_n)(x_l - \bar{x}_n) \qquad (17)$$

The first term is the same as the error in the uncorrelated case, Eq. (16). This means that the second term accounts for the error correction due to correlation between the measurements. For uncorrelated measurements this second term is zero.

Computationally the uncorrelated first term is much easier to treat efficiently than the second.

$$\text{var}(x) = \frac{1}{n}\sum_{k=1}^{n}(x_k - \bar{x}_n)^2 = \left(\frac{1}{n}\sum_{k=1}^{n}x_k^2\right) - \bar{x}_n^2$$

We just accumulate separately the values x^2 and x for every measurement x we receive. The correlation term, though, has to be calculated at the end of the experiment since we need all the measurements to calculate the cross terms. Therefore, all measurements have to be stored throughout the experiment.

Let us analyze the problem by splitting up the correlation term into partial sums of the form:

$$f_d = \frac{1}{n-d}\sum_{k=1}^{n-d}(x_k - \bar{x}_n)(x_{k+d} - \bar{x}_n)$$

The correlation term of the error can now be rewritten in terms of f_d

$$\frac{2}{n}\sum_{k<l}(x_k - \bar{x}_n)(x_l - \bar{x}_n) = 2\sum_{d=1}^{n-1}f_d$$

The value of f_d reflects the correlation between measurements separated by the distance d in the sample samples. Notice that for $d = 0$, f is just the sample variance, $\text{var}(x)$. If we divide f_d by $\text{var}(x)$, we arrive at the so called *autocorrelation function*

$$\kappa_d = \frac{f_d}{\text{var}(x)}$$

The sample error (see eq. (17)) can now be written in terms of the autocorrelation function:

$$
\begin{aligned}
\text{err}_X^2 &= \frac{1}{n}\text{var}(x) + \frac{2}{n}\cdot\text{var}(x)\sum_{d=1}^{n-1}\frac{f_d}{\text{var}(x)} \\
&= \left(1 + 2\sum_{d=1}^{n-1}\kappa_d\right)\frac{1}{n}\text{var}(x) \\
&= \frac{\tau}{n}\cdot\text{var}(x) \qquad (18)
\end{aligned}
$$

and we see that err_X can be expressed in terms the uncorrelated sample variance times a correction factor τ which accounts for the correlation between measurements. We call this correction factor the *autocorrelation time*:

$$\tau = 1 + 2\sum_{d=1}^{n-1}\kappa_d \qquad (19)$$

For a correlation free experiment, τ equals 1. From the point of view of eq. (18) we can interpret a sequential correlation as an effective reduction of the number of measurements by a factor τ. The effective number of measurements becomes:

$$n_{\text{eff}} = \frac{n}{\tau}$$

To neglect the autocorrelation time τ will always cause our simple uncorrelated estimate of $\text{err}_X^2 \approx \text{var}(x)/n$ to be less than the true sample error. The estimate of the error will be too *good*. On the other hand, the calculation of the full autocorrelation time poses an efficiency problem if the set of measurements is very large.

The so-called time-displacement autocorrelation $\phi(t)$ for a quantity $\mathcal{M}$ is given by

$$\phi(t) = \int dt'\left[\mathcal{M}(t') - \langle\mathcal{M}\rangle\right]\left[\mathcal{M}(t'+t) - \langle\mathcal{M}\rangle\right],$$

which can be rewritten as

$$\phi(t) = \int dt'\left[\mathcal{M}(t')\mathcal{M}(t'+t) - \langle\mathcal{M}\rangle^2\right],$$

where $\langle\mathcal{M}\rangle$ is the average value and $\mathcal{M}(t)$ its instantaneous value. We can discretize this function as follows, where we used our set of computed values $\mathcal{M}(t)$ for a set of discretized times (our Monte Carlo cycles corresponding to moving all electrons?)

$$\phi(t) = \frac{1}{t_{\max} - t}\sum_{t'=0}^{t_{\max}-t}\mathcal{M}(t')\mathcal{M}(t'+t) - \frac{1}{t_{\max}-t}\sum_{t'=0}^{t_{\max}-t}\mathcal{M}(t')\times\frac{1}{t_{\max}-t}$$

One should be careful with times close to $t_{\max}$, the upper limit of the sums becomes small and we end up integrating over a rather small time interval. This means that the statistical error in $\phi(t)$ due to the random nature of the fluctuations in $\mathcal{M}(t)$ can become large. One should therefore choose $t \ll t_{\max}$.
Note that the variable $\mathcal{M}$ can be any expectation values of interest. The time-correlation function gives a measure of the correlation between the various values of the variable at a time t' and a time $t' + t$. If we multiply the values of $\mathcal{M}$ at these two different times, we will get a positive contribution if they are fluctuating in the same direction, or a negative value if they fluctuate in the opposite direction. If we then integrate over time, or use the discretized version of, the time correlation function $\phi(t)$ should take a non-zero value if the fluctuations are correlated, else it should gradually go to zero. For times a long way apart the different values of $\mathcal{M}$ are most likely uncorrelated and $\phi(t)$ should be zero.

Time Auto-correlation Function

We can derive the correlation time by observing that our Metropolis algorithm is based on a random walk in the space of all possible spin configurations. Our probability distribution function $\hat{w}(t)$ after a given number of time steps t could be written as

$$\hat{w}(t) = \hat{W}^t \hat{w}(0),$$

with $\hat{w}(0)$ the distribution at $t = 0$ and $\hat{W}$ representing the transition probability matrix. We can always expand $\hat{w}(0)$ in terms of the right eigenvectors of $\hat{v}$ of $\hat{W}$ as

$$\hat{w}(0) = \sum_i \alpha_i \hat{v}_i,$$

resulting in

$$\hat{w}(t) = \hat{W}^t \hat{w}(0) = \hat{W}^t \sum_i \alpha_i \hat{v}_i = \sum_i \lambda_i^t \alpha_i \hat{v}_i,$$

with λ_i the ith eigenvalue corresponding to the eigenvector $\hat{v}_i$.

Time Auto-correlation Function

If we assume that λ_0 is the largest eigenvector we see that in the limit $t \to \infty$, $\hat{w}(t)$ becomes proportional to the corresponding eigenvector $\hat{v}_0$. This is our steady state or final distribution.

We can relate this property to an observable like the mean energy. With the probabilty $\hat{w}(t)$ (which in our case is the squared trial wave function) we can write the expectation values as

$$\langle \mathcal{M}(t) \rangle = \sum_\mu \hat{w}(t)_\mu \mathcal{M}_\mu,$$

or as the scalar of a vector product

$$\langle \mathcal{M}(t) \rangle = \hat{w}(t) \mathbf{m},$$

with $\mathbf{m}$ being the vector whose elements are the values of $\mathcal{M}_\mu$ in its various microstates μ.

Time Auto-correlation Function

We rewrite this relation as

$$\langle \mathcal{M}(t) \rangle = \hat{w}(t) \mathbf{m} = \sum_i \lambda_i^t \alpha_i \hat{v}_i \mathbf{m}_i.$$

If we define $m_i = \hat{v}_i \mathbf{m}_i$ as the expectation value of $\mathcal{M}$ in the ith eigenstate we can rewrite the last equation as

$$\langle \mathcal{M}(t) \rangle = \sum_i \lambda_i^t \alpha_i m_i.$$

Since we have that in the limit $t \to \infty$ the mean value is dominated by the largest eigenvalue λ_0, we can rewrite the last equation as

$$\langle \mathcal{M}(t) \rangle = \langle \mathcal{M}(\infty) \rangle + \sum_{i \neq 0} \lambda_i^t \alpha_i m_i.$$

We define the quantity

$$\tau_i = -\frac{1}{\ldots}.$$

Time Auto-correlation Function

The quantities τ_i are the correlation times for the system. They control also the auto-correlation function discussed above. The longest correlation time is obviously given by the second largest eigenvalue τ_1, which normally defines the correlation time discussed above. For large times, this is the only correlation time that survives. If higher eigenvalues of the transition matrix are well separated from λ_1 and we simulate long enough, τ_1 may well define the correlation time. In other cases we may not be able to extract a reliable result for τ_1. Coming back to the time correlation function $\phi(t)$ we can present a more general definition in terms of the mean magnetizations $\langle \mathcal{M}(t) \rangle$. Recalling that the mean value is equal to $\langle \mathcal{M}(\infty) \rangle$ we arrive at the expectation values

$$\phi(t) = \langle \mathcal{M}(0) - \mathcal{M}(\infty) \rangle \langle \mathcal{M}(t) - \mathcal{M}(\infty) \rangle,$$

resulting in

$$\phi(t) = \sum_{i,j \neq 0} m_i \alpha_i m_j \alpha_j e^{-t/\tau_i},$$

Correlation Time

If the correlation function decays exponentially

$$\phi(t) \sim \exp\left(-t/\tau\right)$$

then the exponential correlation time can be computed as the average

$$\tau_{\text{exp}} = -\langle \frac{t}{\log\left|\frac{\phi(t)}{\phi(0)}\right|} \rangle.$$

If the decay is exponential, then

$$\int_0^\infty dt \phi(t) = \int_0^\infty dt \phi(0) \exp\left(-t/\tau\right) = \tau \phi(0),$$

which suggests another measure of correlation

$$\tau_{\text{int}} = \sum_k \frac{\phi(k)}{\phi(0)},$$

called the integrated correlation time.

What is blocking?

Blocking

- Say that we have a set of samples from a Monte Carlo experiment
- Assuming (wrongly) that our samples are uncorrelated our best estimate of the standard deviation of the mean $\langle \mathcal{M} \rangle$ is given by

$$\sigma = \sqrt{\frac{1}{n} \left(\langle \mathcal{M}^2 \rangle - \langle \mathcal{M} \rangle^2 \right)}$$

- If the samples are correlated we can rewrite our results to show that

$$\sigma = \sqrt{\frac{1 + 2\tau/\Delta t}{n} \left(\langle \mathcal{M}^2 \rangle - \langle \mathcal{M} \rangle^2 \right)}$$

where τ is the correlation time (the time between a sample and the next uncorrelated sample) and Δt is time between each sample

What is blocking?

Blocking

- If $\Delta t \gg \tau$ our first estimate of σ still holds
- Much more common that $\Delta t < \tau$
- In the method of data blocking we divide the sequence of samples into blocks
- We then take the mean $\langle \mathcal{M}_i \rangle$ of block $i = 1 \ldots n_{blocks}$ to calculate the total mean and variance
- The size of each block must be so large that sample j of block i is not correlated with sample j of block $i + 1$
- The correlation time τ would be a good choice

What is blocking?

Blocking

- Problem: We don't know τ or it is too expensive to compute
- Solution: Make a plot of std. dev. as a function of blocksize
- The estimate of std. dev. of correlated data is too low $\rightarrow$ the error will increase with increasing block size until the blocks are uncorrelated, where we reach a plateau
- When the std. dev. stops increasing the blocks are uncorrelated

Implementation

Main ideas

- Do a Monte Carlo simulation, storing all samples to file
- Do the statistical analysis on this file, independently of your Monte Carlo program
- Read the file into an array
- Loop over various block sizes
- For each block size n_b, loop over the array in steps of n_b taking the mean of elements $i n_b, \ldots, (i+1) n_b$
- Take the mean and variance of the resulting array
- Write the results for each block size to file for later analysis

Importance sampling

Main ideas

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, see http://www.elsevier.com/books/stochastic-processes-in-physics-and-chemistry/van-kampen/978-0-444-52965-7 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)

$$\frac{\partial P}{\partial t} = D \frac{\partial}{\partial x} \left(\frac{\partial}{\partial x} - F \right) P(x, t),$$

Importance sampling

Main ideas

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

$$y = x + DF(x)\Delta t + \xi,$$

where ξ is gaussian random variable and Δt is a chosen time step.

Importance sampling

Main ideas

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 ith 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,

$$\frac{\partial^2 P}{\partial \mathbf{x_i}^2} = P \frac{\partial}{\partial \mathbf{x_i}} \mathbf{F_i} + \mathbf{F_i} \frac{\partial}{\partial \mathbf{x_i}} P.$$

Importance sampling

Main ideas

The drift vector should be of the form $\mathbf{F} = g(\mathbf{x})\frac{\partial P}{\partial \mathbf{x}}$. Then,

$$\frac{\partial^2 P}{\partial \mathbf{x}_i^2} = P\frac{\partial g}{\partial P}\left(\frac{\partial P}{\partial \mathbf{x}_i}\right)^2 + Pg\frac{\partial^2 P}{\partial \mathbf{x}_i^2} + g\left(\frac{\partial P}{\partial \mathbf{x}_i}\right)^2.$$

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

$$\boxed{\mathbf{F} = 2\frac{1}{\Psi_T}\nabla\Psi_T,} \tag{20}$$

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.

Importance sampling

Main ideas

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, see http://scitation.aip.org/content/aip/journal/jcp/21/6/10.1063/1.1699114 and URL "http://biomet.oxfordjournals.org/content/57/1/97.abstract",

$$q(y, x) = \frac{G(x, y, \Delta t)|\Psi_T(y)|^2}{G(y, x, \Delta t)|\Psi_T(x)|^2}$$

Importance sampling, program elements

```
for (variate=1; variate <= max_variations; variate++){
    // initialisations of variational parameters and energies
    beta += 0.1;
    energy = energy2 = delta_e = 0.0;
    //  initial trial position, note calling with beta
    for (i = 0; i < number_particles; i++) {
        for ( j=0; j < dimension; j++) {
            r_old[i][j] = gaussian_deviate(&idum)*sqrt(timestep);
        }
    }
    wfold = wave_function(r_old, beta);
    quantum_force(r_old, qforce_old, beta, wfold);
```

Importance sampling, program elements

```
    // loop over monte carlo cycles
    for (cycles = 1; cycles <= number_cycles; cycles++){
        // new position
        for (i = 0; i < number_particles; i++) {
            for ( j=0; j < dimension; j++) {
                // gaussian deviate to compute new positions using a given timestep
                r_new[i][j] = r_old[i][j] + gaussian_deviate(&idum)*sqrt(timestep)+q
                // we move only one particle  at the time
                for (k = 0; k < number_particles; k++) {
                    if ( k != i) {
                        for ( j=0; j < dimension; j++) {
                            r_new[k][j] = r_old[k][j];
                        }
                    }
                }
            }
```

Importance sampling, program elements

```
            wfnew = wave_function(r_new, beta);
            quantum_force(r_new, qforce_new, beta, wfnew);
            // we compute the log of the ratio of the greens functions to be used
            // Metropolis-Hastings algorithm
            greensfunction = 0.0;
            for ( j=0; j < dimension; j++) {
                greensfunction += 0.5*(qforce_old[i][j]+qforce_new[i][j])*
                    (D*timestep*0.5*(qforce_old[i][j]-qforce_new[i][j])-r_new[i][j]+r_
            }
            greensfunction = exp(greensfunction);
```

Importance sampling, program elements

```
            // The Metropolis test is performed by moving one particle  at the t
            if(ran2(&idum) <= greensfunction*wfnew*wfnew/wfold/wfold ) {
                for ( j=0; j < dimension; j++) {
                    r_old[i][j] = r_new[i][j];
                    qforce_old[i][j] = qforce_new[i][j];
                }
                wfold = wfnew;
                .....
```

```
void  quantum_force(double **r, double **qforce, double beta, double w
{
    int i, j;
    double wfminus, wfplus;
    double **r_plus, **r_minus;

    r_plus = (double **) matrix( number_particles, dimension, sizeof(dou
    r_minus = (double **) matrix( number_particles, dimension, sizeof(do
    for (i = 0; i < number_particles; i++) {
        for ( j=0; j < dimension; j++) {
            r_plus[i][j] = r_minus[i][j] = r[i][j];
        }
    }
    ...
```

```
// compute the first derivative
for (i = 0; i < number_particles; i++) {
    for (j = 0; j < dimension; j++) {
        r_plus[i][j] = r[i][j]+h;
        r_minus[i][j] = r[i][j]-h;
        wfminus = wave_function(r_minus, beta);
        wfplus  = wave_function(r_plus, beta);
        qforce[i][j] = (wfplus-wfminus)*2.0/wf/(2*h);
        r_plus[i][j] = r[i][j];
        r_minus[i][j] = r[i][j];
    }
}
```

The general derivative formula of the Jastrow factor is

$$\frac{1}{\Psi_C}\frac{\partial \Psi_C}{\partial x_k} = \sum_{i=1}^{k-1}\frac{\partial g_{ik}}{\partial x_k} + \sum_{i=k+1}^{N}\frac{\partial g_{ki}}{\partial x_k}$$

However, with our

$$\Psi_C = \prod_{i<j} g(r_{ij}) = \exp\left\{\sum_{i<j}\frac{ar_{ij}}{1+\beta r_{ij}}\right\},$$

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.

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

$$R \equiv \frac{\Psi_T^{new}}{\Psi_T^{cur}} = \underbrace{\frac{\Psi_D^{new}}{\Psi_D^{cur}}}_{R_{SD}}\underbrace{\frac{\Psi_C^{new}}{\Psi_C^{cur}}}_{R_C}.$$

Here Ψ_D is our Slater determinant while Ψ_C is our correlation function. We need to optimize $\nabla\Psi_T/\Psi_T$ 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.

$$\frac{\nabla\Psi}{\Psi} = \frac{\nabla(\Psi_D\Psi_C)}{\Psi_D\Psi_C} = \frac{\Psi_C\nabla\Psi_D + \Psi_D\nabla\Psi_C}{\Psi_D\Psi_C} = \frac{\nabla\Psi_D}{\Psi_D} + \frac{\nabla\Psi_C}{\Psi_C}$$

Main ideas

The expectation value of the kinetic energy expressed in atomic units for electron i is

$$\langle\hat{K}_i\rangle = -\frac{1}{2}\frac{\langle\Psi|\nabla_i^2|\Psi\rangle}{\langle\Psi|\Psi\rangle},$$

$$\hat{K}_i = -\frac{1}{2}\frac{\nabla_i^2\Psi}{\Psi}.$$

Main ideas

$$\begin{aligned}
\frac{\nabla^2\Psi}{\Psi} &= \frac{\nabla^2(\Psi_D\Psi_C)}{\Psi_D\Psi_C} = \frac{\nabla\cdot[\nabla(\Psi_D\Psi_C)]}{\Psi_D\Psi_C} = \frac{\nabla\cdot[\Psi_C\nabla\Psi_D + \Psi_D\nabla\Psi_C]}{\Psi_D\Psi_C} \\
&= \frac{\nabla\Psi_C\cdot\nabla\Psi_D + \Psi_C\nabla^2\Psi_D + \nabla\Psi_D\cdot\nabla\Psi_C + \Psi_D\nabla^2\Psi_C}{\Psi_D\Psi_C}
\end{aligned}$$

$$\tag{21}$$

$$\frac{\nabla^2\Psi}{\Psi} = \frac{\nabla^2\Psi_D}{\Psi_D} + \frac{\nabla^2\Psi_C}{\Psi_C} + 2\frac{\nabla\Psi_D}{\Psi_D}\cdot\frac{\nabla\Psi_C}{\Psi_C} \tag{22}$$

Importance sampling

Main ideas

We define 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 $r_{ij} = |\mathbf{r}_i - \mathbf{r}_j| = \sqrt{(x_i - x_j)^2 + (y_i - y_j)^2 + (z_i - z_j)^2}$ for three dimensions and $r_{ij} = |\mathbf{r}_i - \mathbf{r}_j| = \sqrt{(x_i - x_j)^2 + (y_i - y_j)^2}$ for two dimensions.

In our particular case we have

$$\Psi_C = \prod_{i<j} g(r_{ij}) = \exp\left\{\sum_{i<j} f(r_{ij})\right\} = \exp\left\{\sum_{i<j} \frac{a r_{ij}}{1 + \beta r_{ij}}\right\},$$

Importance sampling

Main ideas

The total number of different relative distances r_{ij} is $N(N-1)/2$. In a matrix storage format, the set forms 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 kth-particle . Keep this in mind in the discussion to come.

Importance sampling

Main ideas

$$R_C = \frac{\Psi_C^{\text{new}}}{\Psi_C^{\text{cur}}} = \prod_{i=1}^{k-1} \frac{g_{ik}^{\text{new}}}{g_{ik}^{\text{cur}}} \prod_{i=k+1}^{N} \frac{g_{ki}^{\text{new}}}{g_{ki}^{\text{cur}}}.$$

For the Pade-Jastrow form

$$R_C = \frac{\Psi_C^{\text{new}}}{\Psi_C^{\text{cur}}} = \frac{e^{U_{new}}}{e^{U_{cur}}} = e^{\Delta U},$$

where

$$\Delta U = \sum_{i=1}^{k-1} \left(f_{ik}^{\text{new}} - f_{ik}^{\text{cur}}\right) + \sum_{i=k+1}^{N} \left(f_{ki}^{\text{new}} - f_{ki}^{\text{cur}}\right)$$

Importance sampling

Main ideas

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.

Importance sampling

Main ideas

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 $\mathcal{O}(N)$.

Importance sampling

Main ideas

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

$$\mathbf{r}_{ij} = |\mathbf{r}_j - \mathbf{r}_i| = (x_j - x_i)\vec{e}_1 + (y_j - y_i)\vec{e}_2 + (z_j - z_i)\vec{e}_3$$

is the relative distance.

When the correlation function is the *linear Pade-Jastrow* we have

$$f_{ij} = \frac{ar_{ij}}{(1 + \beta r_{ij})},$$

which yields the closed form expression

$$\frac{\partial f_{ij}}{\partial r_{ij}} = \frac{a}{(1 + \beta r_{ij})^2}.$$

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

$$\Psi_C = \prod_{i<j}\exp f(r_{ij}) = \exp\left\{\sum_{i<j}\frac{ar_{ij}}{1 + \beta r_{ij}}\right\},$$

and it is easy to see that for particle k we have

$$\frac{\nabla_k^2\Psi_C}{\Psi_C} = \sum_{ij\neq k}\frac{(\mathbf{r}_k - \mathbf{r}_i)(\mathbf{r}_k - \mathbf{r}_j)}{r_{ki}r_{kj}}f'(r_{ki})f'(r_{kj}) + \sum_{j\neq k}\left(f''(r_{kj}) + \frac{2}{r_{kj}}f'(r_{kj})\right)$$

Importance sampling

Using
$$f(r_{ij}) = \frac{ar_{ij}}{1 + \beta r_{ij}},$$

and $g'(r_{kj}) = dg(r_{kj})/dr_{kj}$ and $g''(r_{kj}) = d^2 g(r_{kj})/dr_{kj}^2$ we find that for particle k we have

$$\frac{\nabla_k^2 \Psi_C}{\Psi_C} = \sum_{ij \neq k} \frac{(\mathbf{r}_k - \mathbf{r}_i)(\mathbf{r}_k - \mathbf{r}_j)}{r_{ki} r_{kj}} \frac{a}{(1 + \beta r_{ki})^2} \frac{a}{(1 + \beta r_{kj})^2} + \sum_{j \neq k} \left(\frac{2a}{r_{kj}(1 + \beta r_k} \right.$$

Importance sampling

For the correlation part

$$\Psi_C = \prod_{i<j} g(r_{ij}) = \exp\left\{ \sum_{i<j} \frac{ar_{ij}}{1 + \beta r_{ij}} \right\},$$

we need to take into account whether electrons have equal or opposite spins since we have to obey the electron-electron cusp condition as well. When the electrons have equal spins

$$a = 1/4,$$

while for opposite spins (like the ground state in Helium)

$$a = 1/2$$

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 $[x, x + dx]$.

Importance sampling, Fokker-Planck and Langevin equations

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.

Importance sampling, Fokker-Planck and Langevin equations

For many physical systems initial distributions of a stochastic variable y tend to equilibrium distributions: $w(y, t) \to w_0(y)$ as $t \to \infty$. In equilibrium detailed balance constrains the transition rates

$$W(y \to y')w(y) = W(y' \to y)w_0(y),$$

where $W(y' \to y)$ 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 $W(y' \to y)$ and the reverse $W(y \to y')$ may be very different.

Importance sampling, Fokker-Planck and Langevin equations

Consider, for instance, a simple system that has only two energy levels $\epsilon_0 = 0$ and $\epsilon_1 = \Delta E$.
For a system governed by the Boltzmann distribution we find (the partition function has been taken out)

$$W(0 \to 1) \exp -\epsilon_0/kT = W(1 \to 0) \exp -\epsilon_1/kT$$

We get then

$$\frac{W(1 \to 0)}{W(0 \to 1)} = \exp -\Delta E/kT,$$

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 $\Delta t = \epsilon$ 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

$$\tilde{w}(\mathbf{k}, t) = \tilde{w}(\mathbf{k}, 0) \exp{\left[-(D\mathbf{k}^2 t)\right]} = \frac{1}{2\pi} \exp{\left[-(D\mathbf{k}^2 t)\right]}.$$

With the Fourier transform we obtain

$$w(\mathbf{x}, t) = \int_{-\infty}^{\infty} d\mathbf{k} \exp{[i\mathbf{k}\mathbf{x}]} \frac{1}{2\pi} \exp{\left[-(D\mathbf{k}^2 t)\right]} = \frac{1}{\sqrt{4\pi Dt}} \exp{\left[-(\mathbf{x}^2/4Dt)\right]}$$

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

$$W(il - jl, n\epsilon) \rightarrow W(\mathbf{y}, t + \Delta t | \mathbf{x}, t) = \frac{1}{\sqrt{4\pi D\Delta t}} \exp{\left[-((\mathbf{y} - \mathbf{x})^2/4D\Delta t)\right]},$$

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

$$W(\mathbf{x}, \tau|\mathbf{x}-\xi)W(\mathbf{x}-\xi, s|\mathbf{x}_0) = \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!}\frac{\partial^n}{\partial x^n}\left[W(\mathbf{x}+\xi, \tau|\mathbf{x})W(\mathbf{x}, s|\mathbf{x}_0)\right].$$

We can then rewrite the ESKC equation as

$$\frac{\partial W}{\partial s}\tau = -W(\mathbf{x}, s|\mathbf{x}_0) + \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!}\frac{\partial^n}{\partial x^n}\left[W(\mathbf{x}, s|\mathbf{x}_0)\int_{-\infty}^{\infty}\xi^n W(\mathbf{x}+\xi, \tau|\mathbf{x})d\right.$$

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

$$\frac{\partial W(\mathbf{x}, s | \mathbf{x}_0)}{\partial s} \tau = -W(\mathbf{x}, s | \mathbf{x}_0) + \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} \left[W(\mathbf{x}, s | \mathbf{x}_0) \int_{-\infty}^{\infty} \xi^n W(\mathbf{x} + \right.$$

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 $\tau \to 0$ we assume that $\langle [\Delta x(\tau)]^n \rangle \to 0$ 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).

Langevin equation

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
- $\frac{d\mathbf{v}}{dt} = -\xi \mathbf{v} + \mathbf{F}$.

Langevin equation

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

$$\mathbf{v}(t) = \mathbf{v_0}e^{-\xi t} + \int_0^t d\tau e^{-\xi(t-\tau)}\mathbf{F}(\tau).$$

Langevin equation

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

$$\langle \mathbf{v}(t) \cdot \mathbf{v}(t)\rangle_{\mathbf{v_0}} = v_0^{-\xi 2t} + 2\int_0^t d\tau e^{-\xi(2t-\tau)}\mathbf{v_0} \cdot \langle\mathbf{F}(\tau)\rangle_{\mathbf{v_0}}$$

$$+ \int_0^t d\tau' \int_0^t d\tau e^{-\xi(2t-\tau-\tau')}\langle\mathbf{F}(\tau) \cdot \mathbf{F}(\tau')\rangle_{\mathbf{v_0}}.$$

Langevin equation

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

$$\langle\mathbf{F}(t)\rangle = 0,$$

$$\langle\mathbf{F}(t) \cdot \mathbf{F}(t')\rangle_{\mathbf{v_0}} = C_{\mathbf{v_0}}\delta(t - t').$$

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 .

Langevin equation

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

$$\langle(\mathbf{r}(t) - \mathbf{r_0})^2\rangle_{\mathbf{v_0}} = \frac{v_0^2}{\xi}(1 - e^{-\xi t})^2 + \frac{3kT}{m\xi^2}(2\xi t - 3 + 4e^{-\xi t} - e^{-2\xi t}).$$

Langevin equation

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 $\langle(\mathbf{r}(t) - \mathbf{r_0})^2\rangle = 6Dt$.