Week 5 January 31-February 4: Metropolis Algoritm and Markov Chains, Importance Sampling, Fokker-Planck and Langevin equations

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Feb 2, 2022

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Overview of week 5, January 31-February 4

Topics

- Markov Chain Monte Carlo
- Metropolis-Hastings sampling and Importance Sampling
- Fokker-Planck and Langevin equations

Teaching Material, videos and written material

- Overview video on Metropolis algoritm
- Lecture notes here.
- ► See also Lectures from FYS3150/4150 on the Metropolis Algorithm

Basics of the Metropolis Algorithm

The Metropolis et al. algorithm was invented by Metropolis et. a 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

The basic of the Metropolis Algorithm

- ▶ 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 $\mathcal{P}_i^{(n\to\infty)}\to p_i$ 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.

More on the Metropolis

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:

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

Metropolis algorithm, setting it up

Since the probability of making some transition must be 1, $\sum_j T_{i o j} = 1$, and Eq. (1) becomes

$$\mathcal{P}_{i}^{(n)} = \mathcal{P}_{i}^{(n-1)} + \sum_{i} \left[\mathcal{P}_{j}^{(n-1)} T_{j \to i} A_{j \to i} - \mathcal{P}_{i}^{(n-1)} T_{i \to j} A_{i \to j} \right] . \quad (2)$$

Metropolis continues

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

$$\sum_{i} [p_{j} T_{j \to i} A_{j \to i} - p_{i} T_{i \to j} A_{i \to j}] = 0,$$
(3)

Detailed Balance

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

$$\frac{A_{j\to i}}{A_{i\to j}} = \frac{p_i T_{i\to j}}{p_j T_{j\to i}}.$$
 (4)

More on Detailed Balance

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

$$A_{j\to i} = \min\left(1, \frac{p_i T_{i\to j}}{p_j T_{j\to i}}\right). \tag{5}$$

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.¹ 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.

 $^{^{1}}$ The penalty function method uses just such a factor to compensate for p_{i} that are evaluated stochastically and are therefore noisy.

Dynamical Equation

The dynamical equation can be written as

$$\mathcal{P}_i^{(n)} = \sum_j M_{ij} \mathcal{P}_j^{(n-1)} \tag{6}$$

with the matrix M given by

$$M_{ij} = \delta_{ij} \left[1 - \sum_{k} T_{i \to k} A_{i \to k} \right] + T_{j \to i} A_{j \to i}. \tag{7}$$

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

Interpreting 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.

Even a defective matrix has at least one left and right eigenvector for each eigenvalue. An example of a defective matrix is

$$\begin{bmatrix} 0 & 1 \\ 0 & 0 \end{bmatrix},$$

with two zero eigenvalues, only one right eigenvector

$$\begin{bmatrix} 1 \\ 0 \end{bmatrix}$$

and only one left eigenvector (0 1).

Gershgorin bounds and Metropolis

The Gershgorin bounds for the eigenvalues can be derived by multiplying on the left with the eigenvector with the maximum and minimum eigenvalues,

$$\sum_{i} \psi_{i}^{\max} M_{ij} = \lambda_{\max} \psi_{j}^{\max}$$

$$\sum_{i} \psi_{i}^{\min} M_{ij} = \lambda_{\min} \psi_{j}^{\min}$$
(8)

Normalizing the Eigenvectors

Next we choose the normalization of these eigenvectors so that the largest element (or one of the equally largest elements) has value 1. Let's call this element k, and we can therefore bound the magnitude of the other elements to be less than or equal to 1. This leads to the inequalities, using the property that $M_{ij} \geq 0$,

$$\sum_{i} M_{ik} \le \lambda_{\max}$$

$$M_{kk} - \sum_{i \ne k} M_{ik} \ge \lambda_{\min}$$
(9)

where the equality from the maximum will occur only if the eigenvector takes the value 1 for all values of i where $M_{ik} \neq 0$, and the equality for the minimum will occur only if the eigenvector takes the value -1 for all values of $i \neq k$ where $M_{ik} \neq 0$.

More Metropolis analysis

That the maximum eigenvalue is 1 follows immediately from the property that $\sum_{i} M_{ik} = 1$. Similarly the minimum eigenvalue can be -1, but only if $M_{kk} = 0$ and the magnitude of all the other elements ψ_i^{\min} of the eigenvector that multiply nonzero elements M_{ik} are -1.

Let's first see what the properties of M must be to eliminate any -1 eigenvalues. To have a -1 eigenvalue, the left eigenvector must contain only ± 1 and 0 values. Taking in turn each ± 1 value as the maximum, so that it corresponds to the index k, the nonzero M_{ik} values must correspond to i index values of the eigenvector which have opposite sign elements. That is, the M matrix must break up ... back to set A. In particular, there can be no rejections of these moves in the cycle since the -1 eigenvalue requires $M_{kk}=0$. To guarantee no eigenvalues with eigenvalue -1, we simply have to

into sets of states that always make transitions from set A to set B make sure that there are no cycles among states. Notice that this is generally trivial since such cycles cannot have any rejections at any stage. An example of such a cycle is sampling a noninteracting Ising spin. If the transition is taken to flip the spin, and the energy

Final Considerations I

Next we need to make sure that there is only one left eigenvector with eigenvalue 1. To get an eigenvalue 1, the left eigenvector must be constructed from only ones and zeroes. It is straightforward to see that a vector made up of ones and zeroes can only be an eigenvector with eigenvalue 1 if the matrix element $M_{ii}=0$ for all cases where $\psi_i\neq\psi_i$. That is we can choose an index i and take $\psi_i = 1$. We require all elements ψ_i where $M_{ii} \neq 0$ to also have the value 1. Continuing we then require all elements $\psi_{\ell} M_{i\ell}$ to have value 1. Only if the matrix M can be put into block diagonal form can there be more than one choice for the left eigenvector with eigenvalue 1. We therefore require that the transition matrix not be in block diagonal form. This simply means that we must choose the transition probability so that we can get from any allowed state to any other in a series of transitions.

Final Considerations II

Finally, we note that for a defective matrix, with more eigenvalues than independent eigenvectors for eigenvalue 1, the left and right eigenvectors of eigenvalue 1 would be orthogonal. Here the left eigenvector is all 1 except for states that can never be reached, and the right eigenvector is $p_i > 0$ except for states that give zero probability. We already require that we can reach all states that contribute to p_i . Therefore the left and right eigenvectors with eigenvalue 1 do not correspond to a defective sector of the matrix and they are unique. The Metropolis algorithm therefore converges exponentially to the desired distribution.

Final Considerations III

The requirements for the transition $T_{i\rightarrow i}$ are

- ➤ A series of transitions must let us to get from any allowed state to any other by a finite series of transitions.
- ▶ The transitions cannot be grouped into sets of states, A, B, C, such that transitions from A go to B, B to C etc and finally back to A. With condition (a) satisfied, this condition will always be satisfied if either $T_{i\rightarrow i} \neq 0$ or there are some rejected moves.

Importance Sampling: Overview of what needs to be coded

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(x, t) obeys (as an approximation) the so-called Fokker-Planck equation

$$\frac{\partial P}{\partial t} = \sum_{i} D \frac{\partial}{\partial x_{i}} \left(\frac{\partial}{\partial x_{i}} - F_{i} \right) P(x, t),$$

where 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,

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

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

$$\frac{\partial^2 P}{\partial x_i^2} = P \frac{\partial g}{\partial P} \left(\frac{\partial P}{\partial x_i} \right)^2 + P g \frac{\partial^2 P}{\partial x_i^2} + g \left(\frac{\partial P}{\partial x_i} \right)^2.$$

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

$$\mathsf{F} = 2\frac{1}{\Psi_{\mathcal{T}}} \nabla \Psi_{\mathcal{T}},$$

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,

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

Code example for the interacting case with importance sampling

We are now ready to implement importance sampling. This is done here for the two-electron case with the Coulomb interaction, as in the previous example. We have two variational parameters α and β .

After the set up of files

```
# Common imports
import os
# Where to save the figures and data files
PROJECT ROOT DIR = "Results"
FIGURE_ID = "Results/FigureFiles"
DATA_ID = "Results/VMCQdotImportance"
if not os.path.exists(PROJECT_ROOT_DIR):
    os.mkdir(PROJECT ROOT DIR)
if not os.path.exists(FIGURE_ID):
    os.makedirs(FIGURE_ID)
if not os.path.exists(DATA_ID):
    os.makedirs(DATA_ID)
def image_path(fig_id):
   return os.path.join(FIGURE_ID, fig_id)
```

The full code is this link. Here we include only the parts pertaining to the computation of the quantum force and the Metropolis update. The program is a modification of our previous c++ program discussed previously. Here we display only the part from the *vmcsolver.cpp* file. Note the usage of the function *GaussianDeviate*

```
void VMCSolver::runMonteCarloIntegration()
  r0ld = zeros<mat>(nParticles, nDimensions);
  rNew = zeros<mat>(nParticles, nDimensions);
  QForceOld = zeros<mat>(nParticles, nDimensions);
  QForceNew = 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++) {</pre>
    for(int i = 0 \cdot i < nDimensions \cdot i++)
```

```
for(int cycle = 0; cycle < nCycles; cycle++) {</pre>
 // Store the current value of the wave function
 waveFunctionOld = waveFunction(rOld);
 QuantumForce(r0ld, QForce0ld); QForce0ld = QForce0ld*h/waveFunctio
  // New position to test
 for(int i = 0; i < nParticles; i++) {</pre>
    for(int j = 0; j < nDimensions; j++) {</pre>
      rNew(i,j) = rOld(i,j) + GaussianDeviate(&idum)*sqrt(timestep)+
    // for the other particles we need to set the position to the o
    // we move only one particle at the time
    for (int k = 0; k < nParticles; k++) {
      if ( k != i) {
        for (int j=0; j < nDimensions; j++) {
          rNew(k,j) = rOld(k,j);
```

```
// loop over Monte Carlo cycles
   // Recalculate the value of the wave function and the quantum for
   waveFunctionNew = waveFunction(rNew);
   QuantumForce(rNew,QForceNew) = QForceNew*h/waveFunctionNew;
    // we compute the log of the ratio of the greens functions to b
    // Metropolis-Hastings algorithm
   GreensFunction = 0.0;
   for (int j=0; j < nDimensions; j++) {</pre>
      GreensFunction += 0.5*(QForceOld(i,j)+QForceNew(i,j))*
        (D*timestep*0.5*(QForceOld(i,j)-QForceNew(i,j))-rNew(i,j)+rO
   GreensFunction = exp(GreensFunction);
    // The Metropolis test is performed by moving one particle at th
    if(ran2(&idum) <= GreensFunction*(waveFunctionNew*waveFunctionNe
      for(int j = 0; j < nDimensions; j++) {</pre>
        rOld(i,j) = rNew(i,j);
        QForceOld(i,j) = QForceNew(i,j);
        waveFunctionOld = waveFunctionNew;
   } else {
      for(int j = 0; j < nDimensions; j++) {</pre>
        rNew(i,j) = rOld(i,j);
        QForceNew(i,j) = QForceOld(i,j);
```

Note numerical derivatives

```
double VMCSolver::QuantumForce(const mat &r, mat &QForce)
   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
   double kineticEnergy = 0;
   for(int i = 0; i < nParticles; i++) {</pre>
        for(int j = 0; j < nDimensions; j++) {</pre>
            rPlus(i,j) += h;
            rMinus(i,j) -= h;
            waveFunctionMinus = waveFunction(rMinus);
            waveFunctionPlus = waveFunction(rPlus);
            QForce(i,j) = (waveFunctionPlus-waveFunctionMinus);
            rPlus(i,j) = r(i,j);
            rMinus(i,j) = r(i,j);
```

The general derivative formula of the Jastrow factor is (the subscript \mathcal{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)

$$R \equiv \frac{\Psi_T^{new}}{\Psi_T^{old}} = \frac{\Psi_{OB}^{new}}{\Psi_{OB}^{old}} \frac{\Psi_C^{new}}{\Psi_C^{old}}$$

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

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

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 $g = g(r_{ij})$ as well.

In our algorithm we will move one particle at the time, say the *kth*-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

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

For the Pade-Jastrow form

$$R_C = rac{\Psi_C^{
m new}}{\Psi_C^{
m cur}} = rac{\exp U_{new}}{\exp U_{cur}} = \exp \Delta U,$$

where

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

One needs to develop a special algorithm that runs only through the elements of the upper triangular matrix 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

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

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,

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

where

$$r_{ij} = |r_j - r_i| = (x_j - x_i)e_1 + (y_j - y_i)e_2 + (z_j - z_i)e_3$$

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

$$\Psi_C = \prod_{i < i} \exp f(r_{ij}),$$

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

$$\frac{\nabla_k^2 \Psi_C}{\Psi_C} = \sum_{ij \neq k} \frac{(r_k - r_i)(r_k - r_j)}{r_{ki}r_{kj}} f'(r_{ki}) f'(r_{kj}) + \sum_{i \neq k} \left(f''(r_{kj}) + \frac{2}{r_{kj}} f'(r_{kj}) \right)$$

Use the C++ random class for random number generations

```
// Initialize the seed and call the Mersienne algo
std::random_device rd;
std::mt19937_64 gen(rd());
// Set up the uniform distribution for x \in [[0, 1]
std::uniform_real_distribution<double> UniformNumberGenerator(0.0,1.
std::normal_distribution<double> Normaldistribution(0.0,1.0);
```

Use the C++ random class for RNGs, the Mersenne twister class

Finding the new position for importance sampling

```
for (int cycles = 1; cycles <= NumberMCsamples; cycles++){
   // new position
  for (int i = 0; i < NumberParticles; i++) {
    for (int j = 0; j < Dimension; j++) {
        // gaussian deviate to compute new positions using a given timent NewPosition(i,j) = OldPosition(i,j) + Normaldistribution(gen)*
}</pre>
```

Use the C++ random class for RNGs, the Metropolis test

Using the uniform distribution for the Metropolis test

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 $W(y' \to y)$ 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.

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)

$$W(0 \rightarrow 1) \exp{-(\epsilon_0/kT)} = W(1 \rightarrow 0) \exp{-(\epsilon_1/kT)}$$

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

$$w_i(t) = \sum_j W(j \to i)w_j(t=0).$$

The continuous analog to $w_i(0)$ is

$$w(x) \rightarrow \delta(x)$$
,

where we now have generalized the one-dimensional position x to a generic-dimensional vector x. The Kroenecker δ function is replaced by the δ distribution function $\delta(x)$ at t=0.

The transition from a state j to a state i is now replaced by a transition to a state with position y from a state with position 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

$$w(y, t + \Delta t) = \int W(y, t + \Delta t | x, t) w(x, t) dx,$$

and after m time steps we have

$$w(y, t + m\Delta t) = \int W(y, t + m\Delta t | x, t) w(x, t) dx.$$

When equilibrium is reached we have

$$w(y) = \int W(y|x,t)w(x)dx,$$

that is no time-dependence. Note our change of notation for W

We can solve the equation for w(y,t) by making a Fourier transform to momentum space. The PDF w(x,t) is related to its Fourier transform $\tilde{w}(k,t)$ through

$$w(x,t) = \int_{-\infty}^{\infty} dk \exp(ikx) \tilde{w}(k,t),$$

and using the definition of the δ -function

$$\delta(\mathsf{x}) = \frac{1}{2\pi} \int_{-\infty}^{\infty} d\mathsf{k} \exp{(i\mathsf{k}\mathsf{x})},$$

we see that

$$\tilde{w}(\mathsf{k},0)=1/2\pi.$$

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

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

With the Fourier transform we obtain

$$w(x,t) = \int_{-\infty}^{\infty} dk \exp\left[ikx\right] \frac{1}{2\pi} \exp\left[-(Dk^2t)\right] = \frac{1}{\sqrt{4\pi Dt}} \exp\left[-(x^2/4Dt)\right]$$

with the normalization condition

$$\int_{-\infty}^{\infty} w(x,t)dx = 1.$$

The solution represents the probability of finding our random walker at position x at time t if the initial distribution was placed at 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 $n \to \infty$ converges to the normal distribution. It is then possible to show that

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

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 $w(x_0, t_0)$, w(x', t') and w(x, t). We have then

$$w(x,t) = \int_{-\infty}^{\infty} W(x.t|x'.t')w(x',t')dx',$$

and

$$w(x,t) = \int_{-\infty}^{\infty} W(x.t|x_0.t_0)w(x_0,t_0)dx_0,$$

and

$$w(x',t') = \int_{-\infty}^{\infty} W(x'.t'|x_0,t_0)w(x_0,t_0)dx_0.$$

We can combine these equations and arrive at the famous Einstein-Smoluchenski-Kolmogorov-Chapman (ESKC) relation

$$W(xt|x_0t_0) = \int_{-\infty}^{\infty} W(x,t|x',t')W(x',t'|x_0,t_0)dx'.$$

We can replace the spatial dependence with a dependence upon say the velocity (or momentum), that is we have

$$W(v, t|v_0, t_0) = \int_{-\infty}^{\infty} W(v, t|v', t')W(v', t'|v_0, t_0)dx'.$$

We will now derive the Fokker-Planck equation. We start from the ESKC equation

$$W(x, t|x_0, t_0) = \int_{-\infty}^{\infty} W(x, t|x', t') W(x', t'|x_0, t_0) dx'.$$

Define $s=t'-t_0$, $\tau=t-t'$ and $t-t_0=s+\tau$. We have then

$$W(\mathsf{x},s+\tau|\mathsf{x}_0)=\int_{-\infty}^{\infty}W(\mathsf{x},\tau|\mathsf{x}')W(\mathsf{x}',s|\mathsf{x}_0)d\mathsf{x}'.$$

Assume now that τ is very small so that we can make an expansion in terms of a small step xi, with $x' = x - \xi$, that is

$$W(x, s|x_0) + \frac{\partial W}{\partial s}\tau + O(\tau^2) = \int_{-\infty}^{\infty} W(x, \tau|x - \xi)W(x - \xi, s|x_0)dx'.$$

We assume that $W(x, \tau | x - \xi)$ takes non-negligible values only when ξ is small. This is just another way of stating the Master equation!!

We say thus that 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

$$W(x,\tau|x-\xi)W(x-\xi,s|x_0) = \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} \left[W(x+\xi,\tau|x)W(x,s|x_0) \right].$$

We can then rewrite the ESKC equation as

$$\frac{\partial W}{\partial s}\tau = -W(x, s|x_0) + \sum_{n=0}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial x^n} \left[W(x, s|x_0) \int_{-\infty}^{\infty} \xi^n W(x + \xi, \tau|x) dx \right]$$

We have neglected higher powers of τ and have used that for n=0 we get simply $W(x,s|x_0)$ due to normalization.

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We have neglected higher powers of τ and have used that for n=0 we get simply $W(x,s|x_0)$ due to normalization.

We simplify the above by introducing the moments

$$M_n = \frac{1}{\tau} \int_{-\infty}^{\infty} \xi^n W(x + \xi, \tau | x) d\xi = \frac{\langle [\Delta x(\tau)]^n \rangle}{\tau},$$

resulting in

$$\frac{\partial W(\mathsf{x}, s|\mathsf{x}_0)}{\partial s} = \sum_{n=1}^{\infty} \frac{(-\xi)^n}{n!} \frac{\partial^n}{\partial \mathsf{x}^n} \left[W(\mathsf{x}, s|\mathsf{x}_0) M_n \right].$$

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

$$\frac{\partial W(\mathsf{x}, \mathsf{s}|\mathsf{x}_0)}{\partial \mathsf{s}} = -\frac{\partial M_1 W(\mathsf{x}, \mathsf{s}|\mathsf{x}_0)}{\partial \mathsf{x}} + \frac{1}{2} \frac{\partial^2 M_2 W(\mathsf{x}, \mathsf{s}|\mathsf{x}_0)}{\partial \mathsf{x}^2}.$$

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 F(t). The equations of motion are

$$ightharpoonup rac{d\mathbf{r}}{dt} = \mathbf{v} \text{ and }$$

Langevin equation

From hydrodynamics we know that the friction constant $\boldsymbol{\xi}$ is given by

$$\xi = 6\pi \eta a/m$$

where η is the viscosity of the solvent and ${\bf a}$ is the radius of the particle .

Solving the second equation in the previous slide we get

$$\mathsf{v}(t) = \mathsf{v}_0 e^{-\xi t} + \int_0^t d\tau e^{-\xi (t-\tau)} \mathsf{F}(\tau).$$

Langevin equation

If we want to get some useful information out of this, we have to average over all possible realizations of F(t), with the initial velocity as a condition. A useful quantity for example is

$$\langle \mathsf{v}(t) \cdot \mathsf{v}(t) \rangle_{\mathsf{v}_0} = \mathsf{v}_0^{-\xi 2t} + 2 \int_0^t d\tau e^{-\xi (2t - \tau)} \mathsf{v}_0 \cdot \langle \mathsf{F}(\tau) \rangle_{\mathsf{v}_0}$$
$$+ \int_0^t d\tau' \int_0^t d\tau e^{-\xi (2t - \tau - \tau')} \langle \mathsf{F}(\tau) \cdot \mathsf{F}(\tau') \rangle_{\mathsf{v}_0}.$$

Importance sampling, Fokker-Planck and Langevin equations 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

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

and

$$\langle \mathsf{F}(t) \cdot \mathsf{F}(t') \rangle_{\mathsf{v}_0} = C_{\mathsf{v}_0} \delta(t - t').$$

We omit the subscript v_0 , when the quantity of interest turns out to be independent of v_0 . Using the last three equations we get

$$\langle \mathsf{v}(t) \cdot \mathsf{v}(t) \rangle_{\mathsf{v}_0} = v_0^2 e^{-2\xi t} + \frac{C_{\mathsf{v}_0}}{2\xi} (1 - e^{-2\xi t}).$$

For large t this should be equal to 3kT/m, from which it follows that

$$\langle \mathsf{F}(t) \cdot \mathsf{F}(t') \rangle = 6 \frac{kT}{m} \xi \delta(t - t').$$

Langevin equation

Integrating

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

we get

$$\mathsf{r}(t) = \mathsf{r}_0 + \mathsf{v}_0 \frac{1}{\xi} (1 - e^{-\xi t}) + \int_0^t d au \int_0^ au au' e^{-\xi(au - au')} \mathsf{F}(au'),$$

from which we calculate the mean square displacement

$$\langle (\mathsf{r}(t) - \mathsf{r}_0)^2 \rangle_{\mathsf{v}_0} = \frac{\mathsf{v}_0^2}{\xi} (1 - e^{-\xi t})^2 + \frac{3kT}{m\xi^2} (2\xi t - 3 + 4e^{-\xi t} - e^{-2\xi t}).$$

Langevin equation

For very large t this becomes

$$\langle (\mathsf{r}(t) - \mathsf{r}_0)^2 \rangle = \frac{6kT}{m\xi}t$$

from which we get the Einstein relation

$$D=\frac{kT}{m\xi}$$

where we have used $\langle (r(t) - r_0)^2 \rangle = 6Dt$.