

MoTP Exercises 2019/20

1.1 Stochastic Processes and Path Integrals

Exercise 1.1.1 (Stirling's approximation):

Use the Γ function definition:

$$\Gamma(n) \equiv \int_0^\infty x^{n-1} e^{-x} dx \quad n > 0 \quad \Gamma(n+1) = n! \quad (1.1)$$

together with the saddle point approximation to derive the result used in chapter 2 of Lecture Notes:

$$\ln n! = n \ln n - n + \frac{1}{2} \ln(2\pi n) + O\left(\frac{1}{n}\right) \quad (1.2)$$

Solution. To use the saddle-point approximation we need to rewrite (1.1) in the following form:

$$I(\lambda) = \int_S dx \exp\left(-\frac{F(x)}{\lambda}\right) \quad (1.3)$$

So that:

$$I(\lambda) \underset{\lambda \rightarrow 0}{\approx} \sqrt{2\pi\lambda} \left(\frac{\partial F}{\partial x}(x) \Big|_{x=x_0} \right)^{-1/2} \exp\left(-\frac{F(x_0)}{\lambda}\right) \quad (1.4)$$

where x_0 is a global minimum of $F(x)$.

First, we evaluate Γ at $n+1$, and express the integrand as a single exponential:

$$\Gamma(n+1) = n! = \int_0^{+\infty} dx x^n e^{-x} = \int_0^{+\infty} dx e^{-x+n \log x}$$

We want to collect a n in the exponential, and then define $\lambda = 1/n$, so that the saddle-point approximation $\lambda \rightarrow 0$ corresponds to the case of a large factorial

$n \rightarrow \infty$. To do this, we perform a change of variables $x \mapsto s$, so that $x = ns$, with $dx = n ds$:

$$\begin{aligned}\Gamma(n+1) &= \int_0^{+\infty} ds n \exp(-ns + n \log(ns)) = \\ &= n^{n+1} \int_0^{+\infty} ds \exp(n[\log s - s])\end{aligned}$$

In the last step we split the logarithm $n \log(ns) = n \log n + n \log s = n^n + n \log s$, extracted from the integral all terms not depending on s , and then collected the n as desired. Now, letting $\lambda = 1/n$ we have:

$$= n^{n+1} \int_0^{+\infty} ds \exp\left(\frac{\log s - s}{\lambda}\right)$$

which is in the desired form (??).

So, we compute the minimum of $F(s) = \log s - s$:

$$\begin{aligned}F'(s) &= \frac{d}{ds}(s - \log s) = 1 - \frac{1}{s} \stackrel{!}{=} 0 \Rightarrow s_0 = 1 \\ F''(s) &= \frac{1}{s^2} \Rightarrow F''(s_0) = 1 > 0\end{aligned}$$

And applying formula (1.4):

$$n! \underset{n \rightarrow \infty}{\approx} \sqrt{\frac{2\pi}{n}} \cdot 1 \cdot e^{-n} = \sqrt{2\pi} n^{n+\frac{1}{2}} e^{-n}$$

Finally, taking the logarithm leads to the result (1.2):

$$\log n! \underset{n \rightarrow \infty}{\approx} n \log n - n + \frac{1}{2} \log(2\pi n)$$

Exercise 1.1.2 (Random walk tends to a Gaussian):

Implement a numerical simulation to explicitly show how the solution of the ME for a 1-dimensional random walk with $p_{\pm} = 1/2$ tends to the Gaussian.

Exercise 1.1.3 (Non symmetrical motion):

Write the analogous of:

$$W(x, t + \epsilon) = \frac{1}{2}[W(x - l, t) + W(x + l, t)] \quad (1.5)$$

in the LN for the case with $p_+ = 1 - p_- \neq p_-$ and determine:

1. How they depend on l and ϵ in order to have a meaningful continuum limit

2. The resulting continuum equation and how to map it in the diffusion equation:

$$\partial_t W(x, t) = D \partial_x^2 W(x, t)$$

Solution. Consider a Brownian particle moving on a uniform lattice $\{x_i = i \cdot l\}_{i \in \mathbb{N}}$, making exactly one *step* at each *discrete instant* $\{t_n = n \cdot \epsilon\}_{n \in \mathbb{N}}$, with $l, \epsilon \in \mathbb{R}$ fixed. Denoting with p_+ the probability of a *step to the right*, and with p_- that of a *step to the left*, the Master Equation for the particle becomes:

$$W(x, t + \epsilon) = p_+ W(x - l, t) + p_- W(x + l, t) \quad (1.6)$$

1. We already derived (see 7/10) the expected position n at timestep n in that case:

$$\langle x \rangle_{t_n} = nl(p_+ - p_-) = t \frac{l}{\epsilon} (p_+ - p_-) \quad (1.7)$$

Intuitively, an unbalance $p_+ \neq p_-$ will result in a *preferred motion* proportional to that unbalance. Thus we can rewrite (1.7) as:

$$\langle x \rangle_{t_n} = vt \quad v = \frac{l}{\epsilon} (p_+ - p_-)$$

v is the *physical* parameter that needs to be fixed when performing the continuum limit. So, as $p_+ - p_- = 2p_+ - 1$ by normalization, we can find the desired relation between p_+ and v :

$$(2p_+ - 1) \frac{l}{\epsilon} \equiv v \Rightarrow p_+ = \frac{1}{2} \left[\frac{v\epsilon}{l} + 1 \right]$$

As before, we also need to fix $l^2/(2\epsilon) \equiv D$.

2. Expanding each term of (1.6) in a Taylor series we get:

$$\begin{aligned} \cancel{W(x, t)} + \epsilon \dot{W}(x, t) + \frac{\epsilon^2}{2} \ddot{W}(x, t) + O(\epsilon^3) &= p_+ \left[\cancel{W(x, t)} + lW'(x, t) + \frac{1}{2}l^2W''(x, t) + O(l^3) \right] \\ &+ p_- \left[\cancel{W(x, t)} - lW'(x, t) + \frac{1}{2}l^2W''(x, t) + O(l^3) \right] \end{aligned}$$

Using the normalization $p_+ + p_- = 1$ and dividing by ϵ leads to:

$$\dot{W}(x, t) + \frac{\epsilon}{2} \ddot{W}(x, t) + O(\epsilon^2) = (p_+ - p_-) \frac{l}{\epsilon} W'(x, t) + \frac{l^2}{2\epsilon} W''(x, t) + O\left(\frac{l^3}{\epsilon}\right)$$

In the continuum limit $l, \epsilon \rightarrow 0$, with fixed v and D , we get the diffusion equation:

$$\dot{W}(x, t) = vW'(x, t) + DW''(x, t)$$

which leads back to the usual diffusion equation if we set $v = 0$. Note that $p_+ = p_- \Rightarrow v = 0$, as it should be.

Exercise 1.1.4 (Multiple steps at once):

Write the analogous of:

$$W(x, t + \epsilon) = \frac{1}{2}[W(x - l, t) + W(x + l, t)]$$

for the case where the probability to make a step of length $sl \in \{\pm nl : n \in \mathbb{Z} \wedge n > 0\}$ is:

$$p(s) = \frac{1}{Z} \exp(-|s|\alpha)$$

where α is some fixed constant. Determine:

1. the normalization constant Z
2. what is the condition to have a meaningful continuum limit, discussing why the neglected terms do not contribute to such limit
3. which equation you get in the continuum limit

Exercise 1.1.5 (Expected values):

Use equation:

$$W(x, t) \equiv W(x, t | x_0, t_0) = \frac{1}{\sqrt{4\pi Dt}} \exp\left(-\frac{(x - x_0)^2}{4D(t - t_0)}\right) \quad t \geq t_0$$

to determine $\langle x \rangle_t$, $\langle x^2 \rangle_t$ and $\text{Var}_t(x)$.

Exercise 1.1.6 (Diffusion with boundaries):

Consider the diffusion equation:

$$\partial_t W(x, t) = D \partial_x^2 W(x, t)$$

in the domain $[0, \infty)$ instead of $(-\infty, \infty)$. To do that one needs the *boundary condition* (bc) that $W(x, t)$ has to satisfy at 0. Determine the bc for the following two cases and for each of them solve the diffusion equation with the initial condition $W(x, t = 0) = \delta(x - x_0)$ with $x_0 > 0$.

1. *Case of reflecting bc*: when the particle arrives at the origin it bounces back and remains in the domain. How is the flux of particles at 0?
2. *Case of absorbing bc*: when the particle arrives at the origin it is removed from the system (captured by a trap acting like a filter!) What is $W(x = 0, t)$ at all time t ? Notice that in this case we do not expect that the probability is conserved, i.e. we have instead a *Survival Probability*:

$$\mathcal{P}(t) \equiv \int_0^\infty W(x, t) dx$$

that decreases with t . Calculate it and determine its behavior in the two regimes $t \ll x_0^2/D$ and $t \gg x_0^2/D$. Why x_0^2/D is a relevant time scale? (Hint: use the fact that $e^{\pm ikx}$ are eigenfunctions of ∂_x^2 corresponding to the same eigenvalue and choose an appropriate linear combination of them so to satisfy the bc for the two cases. Be aware to ensure that the eigenfunctions so determined are orthonormal. Use also the fact that $\int_{\mathbb{R}} e^{iqx} dx = \delta(q)$)

1.2 Particles in a thermal bath

Exercise 1.2.1 (Harmonic oscillator with general initial condition):

The propagator for a stochastic harmonic oscillator is given by:

$$W(x, t|0, 0) = \sqrt{\frac{k}{2\pi D(1 - e^{-2kt})}} \exp\left(-\frac{k}{2D} \frac{x^2}{1 - e^{-2kt}}\right)$$

Derive the analogous result for $W(x, t|x_0, t_0)$.

Solution. Consider a particle of mass m , experiencing a drag force $F_d = -\gamma v$, an elastic force $F = -kx = -m\omega^2 x$ and thermal fluctuations with amplitude $\sqrt{2D}\gamma$. The equation of motion is given by:

$$m\ddot{x} = -\gamma\dot{x} - m\omega^2 x + \sqrt{2D}\gamma\xi$$

where $\xi(t)$ is a white noise *function*, meaning that $\langle \xi(t)\xi(t') \rangle = \delta(t-t')$ (infinite variance). Dividing by γ and taking the overdamped limit $m/\gamma \ll 0$ we can ignore the \ddot{x} term, leading to a first order SDE:

$$\dot{x} = -\underbrace{\frac{m\omega^2}{\gamma}}_k x + \sqrt{2D}x + \sqrt{2D}\xi \Rightarrow dx(t) = -kx(t) dt + \sqrt{2D} \underbrace{\xi dt}_{dB(t)} \quad (1.8)$$

Consider a time discretization $\{t_j\}_{j=1,\dots,n}$, with $t_n \equiv t$ and the usual notation $x(t_i) \equiv x_i$, $\Delta x_i \equiv x_i - x_{i-1}$. In the Ito prescription, equation (1.8) becomes:

$$\Delta x_i = -kx_{i-1}\Delta t_i + \sqrt{2D}\Delta B_i$$

The probability associated with a sequence $\{\Delta B_j\}_{j=1,\dots,n}$ independent increments is a product of gaussians:

$$\mathbb{P}(\{\Delta B_j\}_{j=1,\dots,n}) = \left(\prod_{i=1}^n \frac{d\Delta B_i}{\sqrt{2\pi\Delta t_i}} \right) \exp\left(-\sum_{i=1}^n \frac{\Delta B_i^2}{2\Delta t_i}\right) \quad (1.9)$$

From (1.9), we can find the probability of the path-increments $\{\Delta x_i\}_{i=1,\dots,n}$ with a change of random variable:

$$\Delta B_i = \frac{\Delta x_i + kx_{i-1}\Delta t_i}{\sqrt{2D}}$$

With the jacobian:

$$J = \det \left| \frac{\partial \{\Delta B_i\}}{\partial \{\Delta x_j\}} \right| = \left| \frac{\partial \{\Delta x_i\}}{\partial \{\Delta B_j\}} \right|^{-1} = \begin{vmatrix} \sqrt{2D} & 0 & \cdots & 0 \\ * & \sqrt{2D} & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ * & \cdots & * & \sqrt{2D} \end{vmatrix}^{-1} = (2D)^{-n/2}$$

The starred terms $*$ are generally non-zero (they are due to the presence of x_{i-1} in the Δx_i formula, which depends on ΔB_j with $j < i - 1$), but the matrix is still lower triangular, meaning that its determinant is just the product of the diagonal terms.

Performing the change of variables leads to:

$$\mathbb{P}(\{\Delta x_i\}_{i=1,\dots,n}) = \left(\prod_{i=1}^n \frac{d\Delta x_i}{\sqrt{4\pi D \Delta t_i}} \right) \exp \left(- \sum_{i=1}^n \frac{1}{2\Delta t_i} \left(\frac{\Delta x_i + kx_{i-1}\Delta t_i}{\sqrt{2D}} \right)^2 \right) \quad (1.10)$$

Taking the continuum limit $n \rightarrow \infty$:

$$dP \equiv \mathbb{P}(\{x(\tau)\}_{t_0 \leq \tau \leq t}) = \left(\prod_{\tau=t_0^+}^t \frac{dx(\tau)}{\sqrt{4\pi D d\tau}} \right) \exp \left(- \frac{1}{4D} \int_{t_0}^t (\dot{x} + kx)^2 d\tau \right)$$

We can finally consider the path integral for the propagator:

$$\begin{aligned} W(x_t, t | x_0, t_0) &= \langle \delta(x_t - x) \rangle_W = \int_{\mathbb{R}^T} \delta(x_t - x) dP = \\ &= \int_{\mathbb{R}^T} dx_W \delta(x_t - x) \exp \left(- \frac{1}{4D} \int_{t_0}^t (\dot{x} + kx)^2 d\tau \right) \end{aligned}$$

The quickest way to compute this integral is to use variational methods. So, consider the functional:

$$S[x(\tau)] = \int_{t_0}^t [\dot{x}(\tau) + kx(\tau)]^2 d\tau$$

The path $x_c(\tau)$ that stationarizes $S[x(\tau)]$ is the solution of the Euler-Lagrange equations:

$$\frac{d}{d\tau} \frac{\partial S}{\partial \dot{x}}(x_c) - \frac{\partial S}{\partial x}(x_c) = 0 \Rightarrow \ddot{x}_c(\tau) = k^2 x(\tau)$$

Leading to:

$$x_c(\tau) = Ae^{k\tau} + Be^{-k\tau}$$

The boundary conditions are $x_c(t_0) = x_0$ and $x_c(t) = x_t$ (this last one is given by the δ). So:

$$\begin{cases} x_0 = Ae^{kt_0} + Be^{-kt_0} \\ x_t = Ae^{kt} + Be^{-kt} \end{cases} \Rightarrow \begin{cases} A = \frac{x_t e^{kt} - x_0 e^{kt_0}}{e^{2kt} - e^{2kt_0}} \\ B = -Ae^{2kt_0} + x_0 e^{kt_0} \end{cases}$$

The integral is then:

$$W(x_t, t|x_0, t_0) = \Phi(t) \exp \left(-\frac{1}{4D} \int_{t_0}^t (\dot{x}_c + kx)^2 d\tau \right)$$

Note that:

$$\dot{x}_c + kx = 2kAe^{k\tau}$$

And so the integral becomes:

$$\int_{t_0}^t (2kAe^{k\tau})^2 d\tau = 2kA^2(e^{2kt} - e^{2kt_0}) = 2k \frac{[x_t e^{kt} - x_0 e^{kt_0}]^2}{e^{2kt} - e^{2kt_0}} = 2k \frac{[x_t - x_0 e^{-k(t-t_0)}]^2}{1 - e^{-2k(t-t_0)}}$$

To compute $\Phi(t)$ we impose the normalization:

$$\int_{\mathbb{R}} dx W(x, t|x_0, t_0) \stackrel{!}{=} 1 \Rightarrow \Phi(t) = \left[\int_{\mathbb{R}} dx \exp \left(-\frac{k}{2D} \frac{[x e^{kt} - x_0 e^{kt_0}]^2}{e^{2kt} - e^{2kt_0}} \right) \right]^{-1}$$

With the substitution $s = x e^{kt} - x_0 e^{kt_0}$ this is just a gaussian integral, evaluating to:

$$\Phi(t) = \sqrt{\frac{k}{2\pi D(1 - e^{-2k(t-t_0)})}}$$

And so the full propagator is:

$$W(x_t, t|x_0, t_0) = \sqrt{\frac{k}{2\pi D(1 - e^{-2k(t-t_0)})}} \exp \left(-\frac{k}{2D} \frac{[x_t - x_0 e^{-k(t-t_0)}]^2}{1 - e^{-2k(t-t_0)}} \right) \quad (1.11)$$

Exercise 1.2.2 (Stationary harmonic oscillator):

Derive the stationary solution $W^*(x)$ of the Fokker Planck equation for the harmonic oscillator, which obeys the following equation:

$$\partial_x [kxW^*(x) + D\partial_x W^*(x)] = 0$$

Explain the hypothesis underlying the derivation and the validity of the derived solution.

Solution. Recall the Fokker-Planck equation for the distribution $W(x, t)$ of a diffusing particle in a medium with diffusion coefficient $D(x, t)$, and in the presence of an external **conservative** force $F(x, t)$ with potential $V(x, t)$ and a drag force $F_d = -\gamma v$:

$$\frac{\partial}{\partial t} W(x, t) = -\frac{\partial}{\partial x} \left[f(x, t) W(x, t) - \frac{\partial}{\partial x} [D(x, t) W(x, t)] \right]$$

where:

$$f(x, t) = \frac{F_{\text{ext}}}{\gamma} = -\frac{1}{\gamma} \frac{\partial V}{\partial x}(x) \quad \gamma = 6\pi\eta a$$

At equilibrium, we expect a time independent solution $W^*(x)$, so that $\partial_t W^*(x) \equiv 0$. We assume, for simplicity, that $\gamma = 1$ and $D(x, t) \equiv D$ **constant**. Letting $F(x, t) = -kx$ be an elastic force, we arrive to:

$$0 = -\partial_x[-kxW^* - D\partial_x W^*] = kxW^*(x) + D\partial_x W^*(x)$$

This is a first order ODE that can be solved by separating variables:

$$\frac{d}{dx}W^* = -\frac{kx}{D}W^* \Rightarrow \frac{dW^*}{W^*} = -\frac{kx}{D}dx \Rightarrow W^*(x) = A \exp\left(-\frac{kx^2}{2D}\right)$$

To be valid, this solution must be **consistent** with the Boltzmann distribution:

$$W^*(x)_{\text{Boltz}} = \frac{1}{Z} \exp(-\beta V(x)) = \frac{1}{Z} \exp\left(-\beta \frac{kx^2}{2}\right)$$

Meaning that $D = 1/\beta = k_B T$.

Exercise 1.2.3 (Harmonic propagator with Fourier transforms):

Use Fourier transforms to derive the full time dependent propagator $W(x, t|x_0, t_0)$ from the FP equation of the harmonic oscillator:

$$\partial_t W(x, t|x_0, t_0) = \partial_x[kxW(x, t|x_0, t_0)] + D\partial_x^2 W(x, t|x_0, t_0) \quad (1.12)$$

Solution. The idea is to use the Fourier transform to *reduce* the equation to a simpler one, that can be hopefully solved.

First, we expand the first derivative:

$$\partial_t W(x, t|x_0, t_0) = kW(x, t|x_0, t_0) + kx\partial_x W(x, t|x_0, t_0) + D\partial_x^2 W(x, t|x_0, t_0)$$

For simplicity, let $W \equiv W(x, t|x_0, t_0)$. Its Fourier transform is given by:

$$\mathcal{F}[W](\omega) \equiv \tilde{W} = \int_{\mathbb{R}} dx e^{-i\omega x} W(x, t|x_0, t_0)$$

The Fourier transforms of the derivatives become:

$$\mathcal{F}[\partial_x W](\omega) = i\omega \tilde{W}; \quad \mathcal{F}[\partial_x^2 W](\omega) = (i\omega)^2 \tilde{W} = -\omega^2 \tilde{W}$$

(These formulas can be proven by repeated integration by parts). All that's left is to transform the remaining term:

$$\begin{aligned} \mathcal{F}[x\partial_x W](\omega) &= \int_{\mathbb{R}} dx e^{-i\omega x} x\partial_x W = \int_{\mathbb{R}} dx i\partial_\omega [e^{-i\omega x} \partial_x W] = i \frac{d}{d\omega} \underbrace{\int_{\mathbb{R}} dx e^{-i\omega x} \partial_x W}_{i\omega \tilde{W}} = \\ &= -\tilde{W} - \omega \partial_\omega \tilde{W} \end{aligned}$$

So (1.12) becomes:

$$\partial_t \tilde{W}(\omega, t) = k\tilde{W} - k\tilde{W} - k\omega \tilde{W}' - D\omega^2 \tilde{W} = -k\omega \tilde{W}'(\omega, t) - D\omega^2 \tilde{W}(\omega, t)$$

Rearranging:

$$k\omega \partial_w \tilde{W}(\omega, t) + \partial_t \tilde{W}(\omega, t) = -k\omega \tilde{W}'(\omega, t) \quad (1.13)$$

This is a *linear* first order partial differential equation. One way to solve it is by using the *method of characteristics*.

Method of characteristics. Consider a general *quasilinear* PDE:

$$a(x, y, z) \frac{\partial z}{\partial x} + b(x, y, z) \frac{\partial z}{\partial y} = c(x, y, z) \quad (1.14)$$

Quasilinear means that a and b can depend also on the dependent variable z , and not only on the independent variables x, y . A solution $z = z(x, y)$ is, geometrically, a *surface graph* immersed in \mathbb{R}^3 . Note that the normal at any point is the gradient of $f(x, y, z) = z(x, y) - z$, that is:

$$\nabla f(x, y, z) = \left(\frac{\partial z}{\partial x}(x, y), \frac{\partial z}{\partial y}(x, y), -1 \right)$$

Rearranging (1.14) we can rewrite it as a dot product:

$$\mathbf{v} \cdot \nabla f = 0 \quad \mathbf{v} = (a(x, y, z), b(x, y, z), c(x, y, z))^T \quad (1.15)$$

This means that, at any point (x, y, z) , the graph $f(x, y, z)$ is *tangent* to the vector field $\mathbf{v}: \mathbb{R}^3 \rightarrow \mathbb{R}^3$, $(x, y, z) \mapsto \mathbf{v}(x, y, z)$.

So, we can consider a set of parametric curves $t \mapsto (x(t), y(t), z(t))$, and *impose* the tangency condition:

$$\begin{cases} \frac{dx}{dt} = a(x, y, z) \\ \frac{dy}{dt} = b(x, y, z) \\ \frac{dz}{dt} = c(x, y, z) \end{cases}$$

This will result in 3 parametric equations in t . If we are able to solve one of the first two for t , we can substitute it and get the desired cartesian form $z = z(x, y)$.

Let $u(\omega, t)$ be a solution. Consider a parameterization $s \mapsto (\omega(s), t(s))$.

$$\frac{du}{ds}(\omega(s), t(s)) = \frac{\partial u}{\partial \omega} \frac{d\omega}{ds} + \frac{\partial u}{\partial t} \frac{dt}{ds} \quad (1.16)$$

By confronting (1.16) with (1.13) we get:

$$\frac{d\omega}{ds} = k\omega; \quad \frac{dt}{ds} = 1 \quad (1.17)$$

Note that now $d\omega/ds$ is exactly the left side of (1.13), so:

$$\frac{du}{ds}(\omega(s), t(s)) = -D\omega(s)^2 u(\omega(s), t(s)) \quad (1.18)$$

For the boundary condition, we suppose $W(x, 0) = \delta(x - x_0)$, meaning that:

$$\tilde{W}(\omega, 0) = \int_{\mathbb{R}} dx e^{-i\omega x} \delta(x - x_0) = e^{-i\omega x_0} = u(\omega, 0)$$

Let's fix $\omega = \omega_0$, and choose the parameterization so that $\omega(s = 0) = \omega_0$ and $t(s = 0) = 0$ (meaning that $u(s = 0) = u(\omega_0, 0)$). We can now solve (1.17):

$$\begin{cases} \frac{d\omega}{ds} = k\omega \\ \omega(0) = \omega_0 \end{cases} \Rightarrow \omega(s) = \omega_0 e^{ks} \quad \begin{cases} \frac{dt}{ds} = 1 \\ t(0) = 0 \end{cases} \Rightarrow t(s) = s \quad (1.19)$$

Finally we can substitute in (1.18) and solve it:

$$\frac{du}{ds} = -D\omega_0^2 e^{2ks} u \Rightarrow \ln |u| = -D\omega_0^2 \int e^{2ks} ds \Rightarrow u(s) = A \exp\left(-\frac{D\omega_0^2}{2k} e^{2ks}\right)$$

And imposing the boundary condition $u(0) = u(\omega_0, s)$ we get:

$$A = u(\omega_0, s) \exp\left(\frac{D\omega_0^2}{2k}\right) \Rightarrow u(\omega(s), t(s)) = u(\omega_0, s) \exp\left[\frac{D\omega_0^2}{2k}(1 - e^{2ks})\right]$$

Note that this solution is expressed as a function of the parameter s , and a starting point ω_0 . By expressing these two as a function of ω and t , we can recover the desired $u(\omega, t)$. To do this, we can simply invert the two solutions (1.19), obtaining:

$$\begin{cases} \omega_0 = \omega e^{-ks} \\ s = t \end{cases} \Rightarrow \begin{cases} \omega_0 = \omega e^{-kt} \\ s = t \end{cases}$$

So that:

$$\begin{aligned} u(\omega, t) \equiv \tilde{W}(\omega, t) &= \exp(-ix_0 \omega e^{-kt}) \exp\left[\frac{D\omega^2 e^{-2kt}}{2k}(1 - e^{2kt})\right] = \\ &= \exp(-ix_0 \omega e^{-kt}) \exp\left[-\frac{D\omega^2}{2k}(1 - e^{-2kt})\right] \end{aligned}$$

All that's left is to perform a Fourier anti-transform to obtain $W(\omega, t)$:

$$\begin{aligned}
W(\omega, t) &= \mathcal{F}^{-1}[\tilde{W}] = \frac{1}{2\pi} \int_{\mathbb{R}} d\omega e^{i\omega x} \exp(-ix_0\omega e^{-kt}) \exp\left[-\frac{D\omega^2}{2k}(1 - e^{-2kt})\right] = \\
&= \frac{1}{2\pi} \int_{\mathbb{R}} d\omega \exp\left(-\underbrace{\frac{D}{2k}(1 - e^{-2kt})}_{a} \omega^2 + i\underbrace{(x - x_0 e^{-kt})}_{b} \omega\right) = \\
&= \frac{1}{2\pi} \sqrt{\frac{\pi}{a}} \exp\left(\frac{b^2}{4a}\right) = \sqrt{\frac{k}{2\pi D(1 - e^{-2kt})}} \exp\left(-\frac{k}{2D} \frac{[x - x_0 e^{-kt}]^2}{1 - e^{-2kt}}\right)
\end{aligned}$$

Which is exactly the same solution found in (1.11).

Exercise 1.2.4 (Multidimensional Fokker-Planck):

Derive the multidimensional Fokker-Planck equation associated to the Langevin equation:

$$d\mathbf{x}^\alpha(t) = f^\alpha(\mathbf{x}(t), t) dt + \sqrt{2D_\alpha(\mathbf{x}(t), t)} dB^\alpha(t) \quad 1 \leq \alpha \leq n \quad (1.20)$$

Solution. We wish to derive from (1.20) a PDE involving the multi-dimensional pdf $W(\mathbf{x}, t)$. To do this, we consider an *ensemble* of paths generated by (1.20), from which we can compute average values, that we can compare with the analogues obtained using $W(\mathbf{x}, t)$, thus reaching the desired relation. First, we consider a generic non-anticipating *test function* $h(\mathbf{x}(t)): \mathbb{R}^n \rightarrow \mathbb{R}$ to be averaged. It's average is, by definition:

$$\langle h(\mathbf{x}(t)) \rangle = \int_{\mathbb{R}} d^n \mathbf{x} W(\mathbf{x}, t) h(\mathbf{x})$$

To construct the ODE, we need the time derivative:

$$\frac{d}{dt} \langle h(\mathbf{x}(t)) \rangle = \int_{\mathbb{R}} d^n \mathbf{x} \dot{W}(\mathbf{x}, t) h(\mathbf{x}) \quad (1.21)$$

We can construct this same derivative starting from (1.20). First consider the differential, i.e. the first order *change* of $h(\mathbf{x}(t))$ after a change of the argument $t \rightarrow t + dt$. We start by considering a change in $\mathbf{x} \rightarrow \mathbf{x} + d\mathbf{x}$, and then use (1.20) to express $d\mathbf{x}$ in terms of dt . Note that Ito's rules imply that $dx^\alpha dx^\beta = dt \delta_{\alpha\beta}$ which is linear in dt and needs to be considered - meaning that we need to expand the \mathbf{x} differential up to *second* order:

$$\begin{aligned}
dh(\mathbf{x}(t)) &= h(\mathbf{x}(t) + d\mathbf{x}(t)) - h(\mathbf{x}(t)) = \\
&= \underline{h(\mathbf{x}(t))} + \sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} dx^\alpha + \frac{1}{2} \sum_{\alpha, \beta=1}^n \frac{\partial^2 h(\mathbf{x})}{\partial x^\alpha \partial x^\beta} dx^\alpha dx^\beta - \underline{h(\mathbf{x}(t))} + O([dx]^3)
\end{aligned}$$

Note that:

$$\begin{aligned}
dx^\alpha dx^\beta &= (f^\alpha dt + \sqrt{2D_\alpha} dB^\alpha)(f^\beta dt + \sqrt{2D_\beta} dB^\beta) = \\
&= 2D_\alpha D_\beta dt \delta_{\alpha\beta} + O(dt^2) + O(dt dB) = 2D_\alpha^2 dt \delta_{\alpha\beta} + O(dt^{3/2})
\end{aligned}$$

And so:

$$\begin{aligned} dh(\mathbf{x}(t)) &= \sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} (f^\alpha dt + \sqrt{2D_\alpha} dB^\alpha) + \frac{1}{2} \sum_{\alpha=1}^n \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} 2D_\alpha dt = \\ &= dt \left[\sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} f^\alpha + D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right] + \sum_{\alpha=1}^n \sqrt{2D_\alpha} \frac{\partial h(\mathbf{x})}{\partial x^\alpha} dB^\alpha \end{aligned}$$

Taking the expected value:

$$\begin{aligned} d\langle h(\mathbf{x}(t)) \rangle &= \langle dt \left[\sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} f^\alpha + D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right] \rangle + \langle \sum_{\alpha=1}^n \sqrt{2D_\alpha} \frac{\partial h(\mathbf{x})}{\partial x^\alpha} dB^\alpha \rangle = \\ &\stackrel{(a)}{=} dt \langle \left[\sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} f^\alpha + D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right] \rangle + \sum_{\alpha=1}^n \langle \sqrt{2D_\alpha} \frac{\partial h(\mathbf{x})}{\partial x^\alpha} \rangle \underbrace{\langle dB^\alpha \rangle}_0 = \\ &= dt \langle \left[\sum_{\alpha=1}^n \frac{\partial h(\mathbf{x})}{\partial x^\alpha} f^\alpha + D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right] \rangle \end{aligned}$$

where in (a) we applied the linearity of the expected value, and then used the fact that h and D_α are non-anticipating, meaning that they are independent of dB^α , leading to a factorization.

Finally, dividing by dt and writing explicitly the averages leads to the desired time derivative:

$$\frac{d\langle h(\mathbf{x}(t)) \rangle}{dt} = \int_{\mathbb{R}^n} d^n \mathbf{x} W(\mathbf{x}, t) \left(\sum_{\alpha=1}^n f^\alpha \frac{\partial h(\mathbf{x})}{\partial x^\alpha} \right) + \int_{\mathbb{R}^n} d^n \mathbf{x} W(\mathbf{x}, t) \left(\sum_{\alpha=1}^n D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right)$$

With a repeated integration by parts we can *move* the derivatives on the $W(\mathbf{x}, t)$, allowing to factorize $h(\mathbf{x})$. This is done by exploiting the fact that $h(\mathbf{x})$ has compact support (as it is a test function), and so:

$$\begin{aligned} \int_{\mathbb{R}^n} d^n \mathbf{x} W(\mathbf{x}, t) \left(\sum_{\alpha=1}^n f^\alpha \frac{\partial h(\mathbf{x})}{\partial x^\alpha} \right) &= - \int_{\mathbb{R}^n} d^n \mathbf{x} h(\mathbf{x}) \left[\sum_{\alpha=1}^n \frac{\partial}{\partial x^\alpha} (W(\mathbf{x}, t) f^\alpha) \right] + \\ &\quad \int_{\mathbb{R}^{n-1}} d^{n-1} \mathbf{x} h(\mathbf{x}) W(\mathbf{x}, t) \sum_{\alpha=1}^n f^\alpha \Big|_{x^\alpha=-\infty}^{x^\alpha=+\infty} \end{aligned}$$

and the boundary term vanishes. A similar procedure holds for the second integral:

$$\int_{\mathbb{R}^n} d^n \mathbf{x} W(\mathbf{x}, t) \left(\sum_{\alpha=1}^n D_\alpha^2 \frac{\partial^2 h(\mathbf{x})}{\partial (x^\alpha)^2} \right) = \int_{\mathbb{R}^n} d^n \mathbf{x} h(\mathbf{x}) \sum_{\alpha=1}^n \frac{\partial^2}{\partial (x^\alpha)^2} [D_\alpha^2 W(\mathbf{x}, t)]$$

This leads to:

$$\begin{aligned} \frac{d\langle h(\mathbf{x}(t)) \rangle}{dt} &= - \int_{\mathbb{R}^n} d^n \mathbf{x} h(\mathbf{x}) \left[\sum_{\alpha=1}^n \frac{\partial}{\partial x^\alpha} (W(\mathbf{x}, t) f^\alpha) \right] + \\ &\quad \int_{\mathbb{R}^n} d^n \mathbf{x} h(\mathbf{x}) \sum_{\alpha=1}^n \frac{\partial^2}{\partial (x^\alpha)^2} [D_\alpha^2 W(\mathbf{x}, t)] \end{aligned} \quad (1.22)$$

Then we equate (1.21) and (1.22):

$$\int_{\mathbb{R}} d^n \mathbf{x} \dot{W}(\mathbf{x}, t) h(\mathbf{x}) = \int_{\mathbb{R}^n} d^n \mathbf{x} h(\mathbf{x}) \sum_{\alpha=1}^n \left[\frac{\partial^2}{\partial (x^\alpha)^2} [D_\alpha^2 W(\mathbf{x}, t)] - \frac{\partial}{\partial x^\alpha} (W(\mathbf{x}, t) f^\alpha) \right]$$

This equality holds for *any* $h(\mathbf{x})$, meaning that the integrands themselves (without the test function) must be everywhere equal:

$$\dot{W}(\mathbf{x}, t) = \sum_{\alpha=1}^n \left[\frac{\partial^2}{\partial (x^\alpha)^2} [D_\alpha^2 W(\mathbf{x}, t)] - \frac{\partial}{\partial x^\alpha} (W(\mathbf{x}, t) f^\alpha) \right]$$

If we suppose D_α to be independent of \mathbf{x} , we could rewrite this relation in a nicer vector form:

$$\dot{W}(\mathbf{x}, t) = \|\mathbf{D}\|^2 \nabla^2 W(\mathbf{x}, t) - \nabla \cdot (W(\mathbf{x}, t) \cdot \mathbf{f})$$

where $\mathbf{D} = (D_1, \dots, D_n)^T$.

Exercise 1.2.5 (Underdamped Wiener measure):

Derive the discretized Wiener measure for the underdamped Langevin equation:

$$m d\mathbf{v}(t) = (-\gamma \mathbf{v} + \mathbf{F}(\mathbf{r})) dt + \gamma \sqrt{2D} d\mathbf{B}$$

and discuss the formal continuum limit.

Solution. The equation can be rewritten as a system of two first order SDE:

$$\begin{cases} d\mathbf{x}(t) = \mathbf{v}(t) dt \\ d\mathbf{v}(t) = \left[-\frac{\gamma}{m} \mathbf{v}(t) + \mathbf{f}(\mathbf{r}) \right] dt + \frac{\gamma}{m} \sqrt{2D} d\mathbf{B} \end{cases}$$

with $\mathbf{f}(\mathbf{r}) = \mathbf{F}(\mathbf{r})/m$.

It is convenient to “symmetrize” the system, by adding an *independent* stochastic term in the first equation:

$$\begin{cases} d\mathbf{x}(t) = \mathbf{v}(t) dt + \sqrt{2\hat{D}} d\hat{\mathbf{B}} \\ d\mathbf{v}(t) = \left[-\frac{\gamma}{m} \mathbf{v}(t) + \mathbf{f}(\mathbf{r}) \right] dt + \frac{\gamma}{m} \sqrt{2D} d\mathbf{B} \end{cases}$$

In this way, we can write a joint pdf for both the position and velocity increments, and then take the limit $\hat{D} \rightarrow 0$. As the $d\mathbf{x}$ are *deterministic*, we expect them to follow a δ distribution.

Explicitly, we introduce a discretization $\{t_i\}_{i=0, \dots, n}$, with fixed endpoints $t_0 \equiv 0$, $t_n \equiv t$. Following Ito’s prescription, the equations become:

$$\begin{cases} \Delta \mathbf{x}_i = \mathbf{v}_{i-1} \Delta t + \sqrt{2\hat{D}} \Delta \hat{\mathbf{B}}_i \\ \Delta \mathbf{v}_i = \left[-\frac{\gamma}{m} \mathbf{v}_{i-1} + \mathbf{f}(\mathbf{x}_{i-1}) \right] \Delta t + \frac{\gamma}{m} \sqrt{2D} \Delta \mathbf{B}_i \end{cases} \quad (1.23)$$

With the usual notation $\mathbf{x}_j \equiv \mathbf{x}(t_j)$. The joint pdf for all the increments is:

$$\begin{aligned} dP(\Delta \mathbf{B}_1, \Delta \hat{\mathbf{B}}_1, \dots, \Delta \mathbf{B}_n, \Delta \hat{\mathbf{B}}_n) &= \left(\prod_{i=1}^n \frac{d^3 \Delta \mathbf{B}_i}{(2\pi \Delta t_i)^{3/2}} \frac{d^3 \Delta \hat{\mathbf{B}}_i}{(2\pi \Delta t_i)^{3/2}} \right) \cdot \\ &\cdot \exp \left(-\frac{1}{2} \sum_{i=1}^n \frac{\|\Delta \mathbf{B}_i\|^2 + \|\Delta \hat{\mathbf{B}}_i\|^2}{\Delta t_i} \right) \end{aligned}$$

where $d^3 \Delta \mathbf{B}_i \equiv \prod_{\alpha=1}^3 d\Delta B_i^\alpha$ (product of the differential of each component of the $d = 3$ vector $\Delta \mathbf{B}_i$).

To get the distribution of the position and velocity increments we perform a change of random variables, inverting (1.23):

$$\begin{aligned} \Delta \hat{\mathbf{B}}_i &= \frac{\Delta \mathbf{x}_i - \mathbf{v}_{i-1} \Delta t}{\sqrt{2\hat{D}}} \\ \Delta \mathbf{B}_i &= \frac{m}{\gamma \sqrt{2D}} \left(\Delta \mathbf{v}_i + \left[\frac{\gamma}{m} \mathbf{v}_{i-1} - \mathbf{f}(\mathbf{x}_{i-1}) \right] \Delta t_i \right) \end{aligned}$$

with jacobian:

$$\begin{aligned} \det \left| \frac{\partial \{\Delta \hat{B}_i^\alpha\}}{\partial \{\Delta x_j^\beta\}} \right| &= (2\hat{D})^{-3n/2} \\ \det \left| \frac{\partial \{\Delta B_i^\alpha\}}{\partial \{\Delta x_j^\beta\}} \right| &= \det \left| \frac{\partial \{\Delta x_j^\beta\}}{\partial \{\Delta B_i^\alpha\}} \right|^{-1} = \left(\frac{\gamma^2}{m^2} 2D \right)^{-3n/2} \end{aligned}$$

And so the final joint distribution is:

$$\begin{aligned} dP(\{\Delta \mathbf{x}_i, \Delta \mathbf{v}_i\}) &= \left(\prod_{i=1}^n \frac{d^3 \Delta \mathbf{x}_i}{(4\pi \hat{D} \Delta t_i)^{3/2}} \frac{d^3 \Delta \mathbf{v}_i}{(4\pi D \Delta t_i \gamma^2 / m^2)^{3/2}} \right) \cdot \\ &\cdot \exp \left(-\frac{m^2}{4D\gamma^2} \sum_{i=1}^n \left\| \frac{\Delta \mathbf{v}_i}{\Delta t_i} + \frac{\gamma}{m} \mathbf{v}_{i-1} - \mathbf{f}(\mathbf{x}_{i-1}) \right\|^2 \Delta t_i \right) \cdot \\ &\cdot \exp \left(-\frac{1}{4\hat{D}} \sum_{i=1}^n \left\| \frac{\Delta \mathbf{x}_i}{\Delta t_i} - \mathbf{v}_{i-1} \right\|^2 \Delta t_i \right) \end{aligned}$$

The highlighted terms become a $\delta(\Delta \mathbf{x}_i - \mathbf{v}_{i-1} \Delta t_i)$ in the limit $\hat{D} \rightarrow 0$ (according to the δ definition as the limit of a normalized gaussian with $\sigma \rightarrow 0$). Then, taking the continuum limit leads to:

$$\begin{aligned} dP(\{\mathbf{x}(\tau), \mathbf{v}(\tau)\}) &= \left(\prod_{\tau=0^+}^t d^3 \mathbf{x}(\tau) \frac{\delta^3(\dot{\mathbf{x}}(\tau) - \mathbf{v}(\tau))}{(d\tau)^3} \frac{d^3 \mathbf{v}(\tau)}{(4\pi D d\tau \gamma^2 / m^2)^{3/2}} \right) \cdot \\ &\cdot \exp \left(-\frac{m^2}{4D\gamma^2} \int_{0^+}^t d\tau \left\| \dot{\mathbf{v}}(\tau) + \frac{\gamma}{m} \mathbf{v}(\tau) - \mathbf{f}(\mathbf{x}(\tau)) \right\|^2 \right) \end{aligned}$$

Exercise 1.2.6 (Maxwell-Boltzmann consistency):

Verify that the Maxwell-Boltzmann distribution:

$$W^*(\mathbf{x}, \mathbf{v}) = \frac{1}{Z^*} \exp \left(-\beta \left[\frac{m \|\mathbf{v}\|^2}{2} + V(\mathbf{x}) \right] \right) \quad (1.24)$$

$$Z^* = \int_{\mathbb{R}^3} d^3 \mathbf{v} \int_{\mathcal{V}} d^3 \mathbf{x} \exp \left(-\beta \left[\frac{m \|\mathbf{v}\|^2}{2} + V(\mathbf{x}) \right] \right)$$

satisfies the Kramers equation:

$$0 = \nabla_{\mathbf{v}} \cdot \left[\left(\frac{\gamma \mathbf{v}}{m} - \frac{\mathbf{F}(\mathbf{x})}{m} \right) W(\mathbf{x}, \mathbf{v}) + \frac{\gamma^2 D}{m^2} \nabla_{\mathbf{v}} W(\mathbf{x}, \mathbf{v}) \right] - \nabla_{\mathbf{x}} \cdot (\mathbf{v} W(\mathbf{x}, \mathbf{v})) \quad (1.25)$$

if the noise amplitude D is given by the Einstein relation:

$$D = \frac{k_B T}{\gamma} = \frac{1}{\beta \gamma}$$

Solution. The idea is to just substitute (1.24) in (1.25). First we compute the relevant *blocks*:

$$\begin{aligned} \nabla_{\mathbf{v}} W^*(\mathbf{x}, \mathbf{v}) &= -\beta m \mathbf{v} W^*(\mathbf{x}, \mathbf{v}) \\ \nabla_{\mathbf{x}} \cdot (\mathbf{v} W^*) &\stackrel{(a)}{=} \mathbf{v} \cdot \nabla_{\mathbf{x}} W^* = -\mathbf{v} \cdot W^* \beta \nabla_{\mathbf{x}} V = \beta W^* \mathbf{v} \cdot \mathbf{F} \end{aligned}$$

where in (a) we used:

$$\nabla \cdot \mathbf{a} f(\mathbf{x}) = \sum_{i=1}^d \frac{\partial}{\partial x_i} [a_i f(\mathbf{x})] = \sum_{i=1}^d a_i \frac{\partial f}{\partial x_i}(\mathbf{x}) = \mathbf{a} \cdot \nabla f \quad \mathbf{a} \in \mathbb{R}^d \text{ constant}$$

Substituting in (1.25):

$$\begin{aligned} & \nabla_{\mathbf{v}} \cdot \left(\left[\frac{\gamma \mathbf{v}}{m} - \frac{\mathbf{F}(\mathbf{x})}{m} - \frac{\gamma^2 D}{m^2} \beta m \mathbf{v} \right] W^*(\mathbf{x}, \mathbf{v}) \right) - \beta W^* \mathbf{v} \cdot \mathbf{F} = \\ &= \nabla_{\mathbf{v}} \cdot \left[W^* \left(\frac{\gamma}{m} - \frac{\gamma^2 D \beta}{m} \right) \mathbf{v} \right] - \nabla_{\mathbf{v}} \cdot \frac{\mathbf{F}}{m} W^* - \beta W^* \mathbf{v} \cdot \mathbf{F} = \\ &= \nabla_{\mathbf{v}} \cdot \left[W^* \left(\frac{\gamma}{m} - \frac{\gamma^2 D \beta}{m} \right) \mathbf{v} \right] - \frac{\mathbf{F}}{m} \cdot \nabla_{\mathbf{v}} W^* - \beta W^* \mathbf{v} \cdot \mathbf{F} = \\ &= \nabla_{\mathbf{v}} \cdot \left[W^* \left(\frac{\gamma}{m} - \frac{\gamma^2 D \beta}{m} \right) \mathbf{v} \right] + \cancel{\frac{\mathbf{F} \cdot \mathbf{v}}{m} \beta m W^*} - \cancel{\beta W^* \mathbf{v} \cdot \mathbf{F}} \end{aligned}$$

Note that the first term vanishes when the expression highlighted in blue is 0, i.e. when:

$$\frac{\gamma}{m} - \frac{\gamma^2 D \beta}{m} = 0 \Leftrightarrow D = \frac{1}{\beta \gamma}$$

which is Einstein's relation for the diffusion coefficient.

Exercise 1.2.7:

Let $P_i(t)$ be the probability that a system is found in the (discrete) state i at time t . If $dt W_{ij}(t)$ represents the transition probability to go from state j to state i during the time interval $(t, t + dt)$, prove that the Master Equation governing the time evolution of the system is:

$$\dot{P}_i(t) = \sum_j (W_{ij}(t)P_j(t) - W_{ji}(t)P_i(t)) \equiv (H(t)P(t))_i$$

where $H_{ij}(t) = W_{ij}(t) - \delta_{ij} \sum_k W_{ki}(t)$.

1. If a_i is an observable quantity (not explicitly dependent on time) of the system when it is in state i , show that:

$$\frac{d\langle a \rangle_t}{dt} = \langle H^T a \rangle_t$$

where $\langle a \rangle_t = \sum_i P_i(t) a_i$

2. If the initial condition is $P_i(t_0) = \delta_{i,i_0}$, the corresponding solution of the Master Equation is called propagator and it will be denoted $P_{i,i_0}(t|t_0)$. Thus $P(t|t_0)$ is a matrix satisfying:

$$\frac{\partial P(t|t_0)}{\partial t} = H(t)P(t|t_0)$$

Show that:

$$\frac{\partial P(t|t_0)}{\partial t_0} = -P(t|t_0)H(t_0)$$

3. Assume now that the transition rates do not depend on time and that an equilibrium stationary state exists. A stationary state P^* satisfies the stationary condition $HP^* = 0$. An equilibrium stationary state, P^{eq} , besides to the stationary condition, satisfies also the so called *detailed balance* (DB) condition $W_{ij}P_j^{\text{eq}} = W_{ji}P_i^{\text{eq}}$ (explain what this means).

If S is the diagonal matrix $S_{ij} = \delta_{ij} \sqrt{P_i^{\text{eq}}}$ show that, as a consequence of the DB condition, the matrix $\hat{H} = S^{-1}HS$ is symmetric and semi-negative definite. Under the hypothesis that each state i can be reached through a path of non-zero transition rates from any state j show that the equilibrium state is unique.

Solution. Consider a uniform time discretization $\{t_n\}_{n \in \mathbb{N}}$, with $t_n - t_{n-1} \equiv \Delta t$. Suppose we know all the probabilities $\{P_j(t_n)\}$ of the system being in any state $j \in J$ at the present time t_n . The probability $P_{j \rightarrow i}$ of a particle transiting from $j \rightarrow i$ at time t_n is the product of the probability of the particle *being* initially at j ($P_j(t_n)$) and the transition probability $W_{ij}(t_n)\Delta t$:

$$P_{j \rightarrow i}(t_n) = W_{ij}(t_n)P_j(t_n)\Delta t$$

Then the probability of the system being in a certain state i at the next timestep t_{n+1} is just the total probability of the system arriving to i at t_{n+1} , that is:

$$P_i(t_{n+1}) = \sum_{j \in J} P_{j \rightarrow i}(t_n) = \sum_{j \in J} W_{ij}(t_n) P_j(t_n) \Delta t$$

We can split the sum to highlight the probability $P_{i \rightarrow i}$ of *remaining* in i , leading to:

$$P_i(t_{n+1}) = \sum_{j \in J \setminus \{i\}} W_{ij} P_j \Delta t + P_{i \rightarrow i}$$

The probability of remaining is just the probability of being in i and *not* transitioning to any other state from i :

$$P_{i \rightarrow i} = P_i \left(1 - \sum_{j \in J \setminus \{i\}} W_{ji} \Delta t \right)$$

Substituting back:

$$P_i(t_{n+1}) = \Delta t \sum_{j \in J \setminus \{i\}} (W_{ij} P_j - W_{ji} P_i) + P_i(t_n)$$

Rearranging and dividing by Δt leads to a Newton's different quotient, which becomes a time derivative in the continuum limit $\Delta t \rightarrow 0$:

$$\frac{P_i(t_n + \Delta t) - P_i(t_n)}{\Delta t} = \sum_{j \in J \setminus \{i\}} (W_{ij} P_j(t_n) - W_{ji} P_i(t_n)) \xrightarrow{\Delta t \rightarrow 0} \dot{P}_i = \sum_{j \in J \setminus \{i\}} (W_{ij} P_j - W_{ji} P_i) \quad (1.26)$$

Before continuing, we wish to rewrite \dot{P}_i as a *matrix multiplication*, i.e. in the form:

$$\dot{P}_i(t) = (H(t) \mathbf{P}(t))_i$$

for a certain $|J| \times |J|$ matrix H , with \mathbf{P} being the vector with the probabilities of each state $(P_j)_{j \in J}^T$. First, notice that we can extend the sum in (??) over the entire J , as the term where $j = i$ vanishes:

$$\dot{P}_i = \sum_{j \in J} (W_{ij} P_j - W_{ji} P_i)$$

Then we rewrite the second term as the following:

$$\sum_{j \in J} W_{ji} P_i = \sum_{\mathbf{k} \in J} W_{\mathbf{k}i} P_i = \sum_{j \in J} \sum_{\mathbf{k} \in J} W_{\mathbf{k}j} P_j \delta_{ij}$$

Now we can collect the P_j :

$$\dot{P}_i = \sum_{j \in J} \left(W_{ij} - \delta_{ij} \sum_{\mathbf{k} \in J} W_{\mathbf{k}j} \right) P_j = \sum_{j \in J} H_{ij} P_j = (H(t) \mathbf{P}(t))_i \quad (1.27)$$

with:

$$H_{ij}(t) = W_{ij}(t) - \delta_{ij} \sum_{k \in J} W_{kj}(t)$$

Note that $H_{ij}(t)$ differs from $W_{ij}(t)$ only on the diagonal elements, which are equal to (minus) the probability of *escape* from that state:

$$H_{jj} = W_{jj} - W_{jj} - \sum_{k \neq j} W_{kj} = - \sum_{k \neq j} W_{kj}$$

1. Let A be an observable of the system, assuming values a_i in each state i . At a fixed time t , the system state is described by the discrete probability distribution (or *probability mass function*) $P_i(t)$. So, the average of A at time t is:

$$\langle a \rangle_t = \sum_{i \in J} P_i(t) a_i$$

Suppose that a_i does not depend on time. Differentiating:

$$\begin{aligned} \frac{d\langle a \rangle_t}{dt} &= \sum_{i \in J} \dot{P}_i(t) a_i \stackrel{(1.27)}{=} \sum_{i \in J} \sum_{j \in J} H_{ij} P_j(t) a_i = \\ &= \sum_{j \in J} P_j(t) \left(\sum_{i \in J} H_{ij} a_i \right) = \sum_{j \in J} P_j(t) (H^T \mathbf{a})_j = \langle H^T \mathbf{a} \rangle_t \end{aligned}$$

where \mathbf{a} is the vector $(a_j)_{j \in J}^T$.

2. The propagator $P(i, t | i_0, t_0) \equiv P_{i, i_0}(t | t_0)$ is just the transition probability from an initial defined state i_0 at t_0 to a generic state i at t .

Consider now a uniform time discretization, and construct the desired time derivative:

$$\frac{\partial}{\partial t} P(i, t | i_0, t_0) = \lim_{\Delta t \rightarrow 0} \frac{P(i, t | i_0, t_0) - P(i, t | i_0, t_0 - \Delta t)}{\Delta t}$$

We choose this definition so that $t_0 - \Delta t < t_0 < t$, and we can apply the Chapman-Kolmogorov equation (that holds as the system is Markovian):

$$P(i, t | i_0, t_0 - \Delta t) = \sum_{i' \in J} P(i, t | i', t_0) P(i', t_0 | i_0, t_0 - \Delta t) \quad (1.28)$$

That is, the transition probability $(i_0, t_0 - \Delta t) \rightarrow (i, t)$ can be obtained by *splitting* the path into two steps $(i_0, t_0 - \Delta t) \rightarrow (i_0, t_0)$ and $(i_0, t_0) \rightarrow (i, t)$, multiplying the two transition probabilities, and summing over all the possible intermediate states i' .

Note now that $P(i', t_0 | i_0, t_0 - \Delta t)$ is a transition probability over *a single timestep*, and so can be computed using the transition probability matrix:

$$P(i', t_0 | i_0, t_0 - \Delta t) = W_{i' i_0} \Delta t$$

Substituting back in (1.28):

$$P(i, t|i_0, t_0 - \Delta t) = \sum_{i' \in J} P(i, t|i', t_0) W_{i'i_0}$$

As before, we highlight the case of the system remaining in the same state i_0 :

$$= \sum_{i' \neq i_0} P(i, t|i', t_0) W_{i'i_0} \Delta t + P(i, t|i_0, t_0) \left(1 - \sum_{k \neq i_0} W_{ki_0} \Delta t \right)$$

where the probability of remaining in i_0 is equal to the probability of *not* going to any other state k .

We can now construct the difference quotient:

$$\begin{aligned} \frac{P(i, t|i_0, t_0) - P(i, t|i_0, t_0 - \Delta t)}{\Delta t} &= - \sum_{i' \neq i_0} P(i, t|i', t_0) W_{i'i_0} + P(i, t|i_0, t_0) \sum_{k \neq i_0} W_{ki_0} = \\ &= - \sum_{i' \neq i_0} P(i, t|i', t_0) W_{i'i_0} + \sum_{i' \neq i_0} \delta_{i'i_0} P(i, t|i', t_0) \sum_{k \neq i_0} W_{ki_0} = \\ &= - \sum_{i' \neq i_0} P(i, t|i', t_0) \underbrace{\left[W_{i'i_0} - \delta_{i'i_0} \sum_{k \neq i_0} W_{ki_0} \right]}_{H_{i'i_0}} = \\ &= - \sum_{i' \neq j} P_{ii'}(t|t_0) H_{i'i_0}(t_0) \end{aligned}$$

And so, taking the continuum limit $\Delta t \rightarrow 0$:

$$\frac{\partial P_{ii_0}(t|t_0)}{\partial t_0} = - \sum_{i' \neq j} P_{ii'}(t|t_0) H_{i'i_0}(t_0) \Rightarrow \frac{\partial P(t|t_0)}{\partial t_0} = -P(t|t_0) H(t_0)$$

Where $P(t|t_0)$ is the $|J| \times |J|$ matrix with entries $P_{ij}(t|t_0)$.

3. The detailed balance (DB) condition is:

$$W_{ij} P_j^{\text{eq}} = W_{ji} P_i^{\text{eq}}$$

This means that the probability of a transition $j \rightarrow i$ is, at equilibrium, exactly the same as the probability of the inverse transition $i \rightarrow j$. In other words, every process that would change the state is *exactly balanced* by its inverse process.

We consider now a time-independent W_{ij} , and the diagonal matrix $S_{ij} = \delta_{ij} \sqrt{P_i^{\text{eq}}}$. Then:

$$\hat{H} = S^{-1} H S \Rightarrow \hat{H}_{ij} = \sum_{ks} \frac{1}{\sqrt{P_i^{\text{eq}}}} \delta_{ik} H_{ks} \delta_{sj} \sqrt{P_j^{\text{eq}}} = \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} H_{ij}$$

To prove that \hat{H} is symmetric, we need to show that the off-diagonal elements remain the same after inverting $j \leftrightarrow i$. That is:

$$(\hat{H}^T)_{ij} = \hat{H}_{ji} = \sqrt{\frac{P_i^{\text{eq}}}{P_j^{\text{eq}}}} H_{ji} \stackrel{?}{=} \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} H_{ij}$$

Recall that:

$$H_{ij} = W_{ij} - \delta_{ij} \sum_k W_{ki}$$

meaning that for $i \neq j$, $H_{ij} = W_{ij}$. So:

$$\sqrt{\frac{P_i^{\text{eq}}}{P_j^{\text{eq}}}} W_{ji} \stackrel{?}{=} \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} \Leftrightarrow W_{ji} P_i^{\text{eq}} \stackrel{?}{=} W_{ij} P_j^{\text{eq}}$$

And the latter is exactly the DB condition, and so DB implies \hat{H} symmetric.

To check if \hat{H} is negative definite, we need to show that:

$$\sum_{ij} x_i \hat{H}_{ij} x_j \leq 0 \quad \forall \mathbf{x} \in \mathbb{R}^{|J|} \setminus \{\mathbf{0}\}$$

Expanding:

$$\begin{aligned} \sum_{ij} x_i \hat{H}_{ij} x_j &= \sum_{ij} x_i \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} H_{ij} x_j = \sum_{ij} x_i \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} x_j - \sum_i x_i^2 \sqrt{\frac{P_i^{\text{eq}}}{P_i^{\text{eq}}}} \sum_k W_{ki} = \\ &= \sum_{ij} \left(\sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} x_i x_j - x_i^2 W_{ji} \right) \end{aligned}$$

As \hat{H} is symmetric:

$$\begin{aligned} \sum_{ij} x_i \hat{H}_{ij} x_j &= \sum_{ij} x_j \hat{H}_{ji} x_i = \sum_{ij} \left(\sqrt{\frac{P_i^{\text{eq}}}{P_j^{\text{eq}}}} W_{ji} x_i x_j - x_j^2 W_{ij} \right) = \\ &\stackrel{(a)}{=} \sum_{ij} \left(\sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} x_i x_j - x_j^2 W_{ij} \right) \end{aligned}$$

where in (a) we used (DB), or more precisely:

$$W_{ji} = W_{ij} \frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}$$

So the sum of the “two versions” of the product will be exactly two times the original sum:

$$\begin{aligned}
\sum_{ij} x_i \hat{H}_{ij} x_j &= \frac{1}{2} \sum_{ij} \left[x_i \hat{H}_{ij} x_j + \sum_{ij} x_j \hat{H}_{ji} x_i \right] = \\
&= \frac{1}{2} \sum_{ij} \left[2 \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} x_i x_j - x_i^2 W_{ji} - x_j^2 W_{ij} \right] = \\
&\stackrel{(\text{??})}{=} \frac{1}{2} \sum_{ij} \left[2 \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} W_{ij} x_i x_j - x_i^2 \frac{P_j^{\text{eq}}}{P_i^{\text{eq}}} W_{ij} - x_j^2 W_{ij} \right] = \\
&= -\frac{1}{2} \sum_{ij} \left[x_i \sqrt{\frac{P_j^{\text{eq}}}{P_i^{\text{eq}}}} + x_j \right]^2 W_{ij} \leq 0
\end{aligned}$$

As $W_{ij} \geq 0$.

All that's left is to show that the equilibrium state is unique, under the hypothesis that each state i can be reached through a path of non-zero transition rates from any state j .

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