Effective rate equations for the overdamped motion in fluctuating potentials

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We discuss physical and mathematical aspects of the overdamped motion of a Brownian particle in fluctuating potentials. It is shown that such a system can be described quantitatively by fluctuating rates if the potential fluctuations are slow compared to relaxation within the minima of the potential, and if the position of the minima does not fluctuate. Effective rates can be calculated; they describe the long-time dynamics of the system. Furthermore, we show the existence of a stationary solution of the Fokker-Planck equation that describes the motion within the fluctuating potential under some general conditions. We also show that a stationary solution of the rate equations with fluctuating rates exists.

DOI: 10.1103/PhysRevE.64.021106

PACS number(s): 05.40.-a, 02.50.Ey, 82.20.Fd, 87.16.Ac

I. INTRODUCTION

Thermally activated relaxation processes are an important mechanism for the dynamics of many physical, chemical, and biological systems. The diffusion of a Brownian particle in a potential with several minima and barriers serves as a paradigm for such relaxation processes. The minima of the potential represent stable or metastable states of the system, and the motion of the particle models the transition of the system from one state to the other and back. The most simple example in this class of models is the problem of diffusion over a single potential barrier pioneered by Kramers [1]. The dynamics of the diffusion over a barrier is dominated by a characteristic time scale, which is given by the mean first passage time for the escape out of the minimum of the potential. In a model with several minima and barriers, the corresponding time scales are given by the inverse rates for the transition from one state to another.

Often one uses kinetic rate equations to describe the dynamical properties of relaxation in chemical or biological systems. This description is clearly much simpler than the description of a Brownian particle in a potential. To calculate rate coefficients, it is often sufficient to know some of the properties of the potential energy surface, like barrier heights. In principle, if one knows the potential energy surface, both descriptions are equivalent. One can calculate static or dynamic properties of the system either using the Brownian particle in the potential or using the rate equations.

In many situations, the potential fluctuates due to some external fluctuations, chemical reactions, or oscillations. These fluctuations usually have a finite correlation time, which corresponds to a nonthermal noise. The most simple model with a fluctuating potential consists of a single fluctuating barrier, where the height of the barrier fluctuates between two different values. Doering and Gadoua [2] investigated such a simple model. They found a local minimum in the mean first passage time as a function of the barrier fluctuation rate. This effect has been called resonant activation and has been studied extensively [3-8].

In principle it should be possible to use a description with rate equations in the case with fluctuating barriers too. This picture has already been suggested by Bier and Astumian [4]. They showed that the long-time dynamics of a simple model with a dichotomously fluctuating linear ramp can be described by kinetic rate equations if the potential fluctuations are not too fast. A similar observation was made in [9], where it was shown that noise induced stability of a metastable state occurs in systems with a single fluctuating barrier.

In the present paper we consider potentials in d dimensions with several minima and barriers. Furthermore, we allow general potential fluctuations. For simplicity we assume that the potential fluctuations can be parametrized by a single Markovian noise process z(t). Often it is argued that the Markovian noise z(t) represents the cumulative effects of many weakly coupled environmental degrees of freedom. In that case the central limit theorem can be applied and z(t)becomes an Ornstein-Uhlenbeck process. On the other hand, there are many realistic situations where the fluctuation of the potential is triggered by a single (or few) environmental degree of freedom. The most simple case of a non-Gaussian fluctuating potential is a dichotomously fluctuating potential as discussed, e.g., in [2]. But in many applications, the potential fluctuates not only between two different values. Since all these different cases are of interest, we will investigate general noise processes. Furthermore, it is clear that there are always many possibilities to parametrize a fluctuating potential by a Markovian noise process. This is another reason why it is useful to consider general noise processes.

A main purpose of the present paper is to show under which conditions it is possible to pass from a model with a fluctuating potential to a set of rate equations with fluctuating rates. We will discuss general models in arbitrary dimensions and with potentials having several minima. The general picture we support by our calculations is that rate equations can be used if the potential fluctuations are slow compared to the relaxation within the minima of the potential and if the positions of the minima of the potential do not fluctuate. The long-time behavior of the model with a fluctuating potential is then similar to that obtained using rate equations. The fluctuation of the potential fluctuations may be fast compared to relaxation across the barriers of the potential.

After these remarks it is clear but important to mention that the applicability of rate equations depends also on the quantities one is interested in. As long as one is interested in stationary or quasistationary properties, rate equations yield good results. If one is interested in properties on time scales of the order of or smaller than the relaxation within the minima, rate equations cannot be applied.

Furthermore, we analyze systems with fluctuating rates in detail and derive effective rates (which do not fluctuate). The effective rates yield the relevant time scales for the long-time behavior of the system. They can be directly compared to the mean escape rates in the fluctuating potentials. Furthermore, they yield the stationary distribution, which can also be compared with results obtained directly by solving the Fokker-Planck equation with a fluctuating potential. The comparison is easily done for simple, one-dimensional systems with piecewise linear potentials. For such systems it is relatively easy to solve the Fokker-Planck equation directly.

The paper is organized as follows. In the next section we discuss how a kinetic rate equation with fluctuating rates can be obtained starting from a Fokker-Planck equation with a fluctuating potential, and in which situations one can expect it to be valid. We also show how one can derive effective rates from the kinetic rate equation with fluctuating rates. The effective rates yield the characteristic time scales for the escape out of the minima of the fluctuating potential. The general formula for the effective rates can be evaluated for special noise processes. This is done in Sec. III. The results are used to compare the description of a system by effective rates with the original Fokker-Planck equation. For our discussion the existence of a stationary solution of the Fokker-Planck equation with a fluctuating potential and the existence of a stationary solution of a kinetic rate equation with fluctuating rates is essential. Section III also contains the exact solution of the dichotomous two-state system and its derivation. Some of the exact results have already been mentioned without derivation in [9]. In Sec. IV we show that the stationary solutions of the rate equation with fluctuating rates and of the Fokker-Planck equation with a fluctuating potential exist under some general assumptions on the potential. This result is a generalization of a recent, similar result by Pechukas and Ankerhold [10]. Finally, in Sec. V, some possible applications are mentioned and conclusions of our results are drawn. As a specific example we discuss how our results can be used to interpret experimental findings for membrane proteins.

II. DERIVATION OF EFFECTIVE RATE EQUATIONS

A. Noise processes and the Fokker-Planck equation

The overdamped motion of a particle in a fluctuating potential is usually described by a Langevin equation,

$$\frac{d\vec{x}}{dt} = \vec{f}(\vec{x},t) + \sqrt{2T}\vec{\xi}(t).$$
(1)

We have chosen the units such that the friction constant is unity. $\xi(t)$ is a thermal (Gaussian white) noise, it satisfies $\langle \vec{\xi}(t) \rangle = 0$, $\langle \xi_a(t)\xi_b(t') \rangle = \delta_{ab}\delta(t-t')$. $\vec{f}(\vec{x},t) = -\nabla V(\vec{x},t)$ is the force of the fluctuating potential. We investigate situations where the potential tends to infinity for $|\vec{x}| \to \infty$ or where the potential is defined in a finite domain Ω with reflecting boundary conditions on the boundary $\partial \Omega$. Let us assume that the potential fluctuations can be parametrized by a single stochastic variable z(t), $V(\vec{x},t) = V(\vec{x},z(t))$. It is clear that such a representation is not unique. Instead of z(t)one can use any monotonic function of z(t) to parametrize V(x,t). Often one uses a simple linear ansatz V(x,z(t)) $= V_0(x) + z(t)\Delta V(x)$ with a dichotomous or an Ornstein-Uhlenbeck process z(t). As already mentioned in the Introduction, we want to investigate models for which the potential fluctuations represent the effects of few environmental degrees of freedom. Therefore we do not restrict ourselves to a linear ansatz or to special noise processes. The only assumption we make is that the noise process z(t) can be described using a Fokker-Planck equation for the probability distribution p(z,t),

$$\frac{\partial p(z,t)}{\partial t} = M_z p(z,t). \tag{2}$$

The right eigenfunctions of the generator M_z are denoted by $\phi_n(z)$, the eigenvalues by $-\lambda_n$,

$$M_z \phi_n(z) = -\lambda_n \phi_n(z), \qquad (3)$$

where $\lambda_0 = 0$, $\lambda_1 > 0$, $\lambda_n \ge \lambda_{n-1}$. $\phi_0(z)$ is the stationary distribution of the noise process z(t). We let $\phi_n(z) = g_n(z)\phi_0(z)$. Then one has

$$\int dz g_n(z) g_m(z) \phi_0(z) = \delta_{n,m} \,. \tag{4}$$

 $g_n(z)$ are orthogonal functions with respect to the weight function $\phi_0(z)$. For n=m this equation fixes the normalization of $\phi_n(z)$. The eigenvalue $\lambda_1 = \tau^{-1}$ determines the correlation time of the noise process z(t). We already mentioned that the parametrization of the fluctuating potential by a stochastic variable z(t) is not unique. One can always choose a different parametrization. The final results should of course be independent of the parametrization. If one changes the parametrization, the spectrum of M_z remains the same, but the right eigenfunctions are changed.

The motion of the overdamped particle in the fluctuating potential can be described using the joint probability distribution $\rho(\vec{x}, z, t)$ for $\vec{x}(t)$ and z(t). It obeys the Fokker-Planck equation

$$\frac{\partial \rho(\vec{x}, z, t)}{\partial t} = \vec{\nabla} \cdot [(\vec{\nabla} V(\vec{x}, z)) + T\vec{\nabla}] \rho(\vec{x}, z, t) + M_z \rho(\vec{x}, z, t).$$
(5)

 $V(\vec{x},z)$ must have the same properties as a function of \vec{x} as $V(\vec{x},t)$. It must tend to infinity sufficiently fast when $|\vec{x}|$ becomes large; or alternatively, we may define the Fokker-Planck equation on a finite, open domain Ω with reflecting boundary conditions on the boundary $\partial \Omega$. Furthermore we

assume that $V(\vec{x},z)$ is finite for any z out of the support of $\phi_0(z)$. $V(\vec{x},z)$ does not have any infinitely high potential wells inside Ω .

B. Rate equations

1. Rate equations for fixed potentials

Let us first briefly review the case of a potential that does not fluctuate. The Fokker-Planck operator has the form

$$L = \vec{\nabla} \cdot \left[\left(\vec{\nabla} V(\vec{x}) \right] + T \vec{\nabla} \right). \tag{6}$$

Let us assume that the potential V(x) has N distinct minima at the points $\vec{x_i}$, i = 0, ..., N-1 and that $V(\vec{x}) \rightarrow \infty$ for $|\vec{x}|$ $\rightarrow \infty$ sufficiently fast. L has one eigenvalue 0, the corresponding right eigenfunction is the stationary distribution $\propto \exp[-V(\vec{x})/T]$. -L is non-negative, i.e., all the other eigenvalues of -L are positive. Using a simple variational argument, one can show that -L has at least (N-1) eigenvalues that behave like $\exp(-c/T)$ for small T. Let us divide the entire space Ω in which the system is defined into subspaces Ω_i . $\vec{x}_i \in \Omega_j$ if and only if i=j, $\Omega_i \cap \Omega_j = \emptyset$ if $i \neq j$. $\bigcup_i \Omega_i$ = $\Omega \setminus (\bigcup_i \partial \Omega_i)$. $\partial \Omega_i$ denotes the boundary of Ω_i . The boundaries are chosen such that $n_i(x) \cdot \nabla V(x) = 0$ for $x \in \partial \Omega_i$, where $\vec{n}_i(\vec{x})$ is the normal vector to the boundary $\partial \Omega_i$ at \vec{x} . Such a division of Ω is always possible. Depending on the form of $V(\vec{x})$ it may not be unique. Let us now introduce subsets $\Omega_i \subset \Omega_i$, which do not contain a small region close to the boundary of Ω_i , but which do contain the positions x_i of the minima together with a sufficiently large region around them. Let $\tilde{\theta}_i(x)$ be a smooth function that is equal to unity if $\vec{x} \in \tilde{\Omega}_i$ and that vanishes if $\vec{x} \notin \Omega_i$. In $\Omega_i \setminus \tilde{\Omega}_i$ the function $\tilde{\theta}_i(\vec{x})$ drops smoothly from 1 to 0. We assume that sufficiently many derivatives of $\tilde{\theta}_i$ exist.

In a next step, we introduce $\hat{L} = \exp[V(\vec{x})/(2T)]L \exp[-V(\vec{x})/(2T)]$. \hat{L} is a Hermitian operator. It has the explicit form

$$\hat{L} = T\nabla^2 + \frac{1}{2}(\vec{\nabla}^2 V) - \frac{1}{4T}(\vec{\nabla} V)^2.$$
(7)

Clearly, \hat{L} has the same eigenvalues as L. exp $[-V(\vec{x})/(2T)]$ is the (un-normalized) eigenfunction of \hat{L} to eigenvalue 0. Let $\hat{g}_i(\vec{x}) = C_i \exp[-V(\vec{x})/(2T)]\tilde{\theta}_i(\vec{x})$, where C_i is chosen such that $\int d^d x |\hat{g}_i(\vec{x})|^2 = 1$. One has $\int d^d x \hat{g}_i(\vec{x}) \hat{g}_j(\vec{x}) = \delta_{i,j}$. Let $\hat{L}_{i,j} = \int d^d x \hat{g}_i(\vec{x}) \hat{L} \hat{g}_j(\vec{x})$. These quantities vanish for $i \neq j$. For i = j one obtains

$$\hat{L}_{ii} = \frac{T}{2} |C_i|^2 \int d^d x \vec{\nabla} \cdot [\exp\{-V(\vec{x})/T\} \nabla \tilde{\theta}_i(\vec{x})].$$
(8)

Since $\nabla \tilde{\theta}_i$ vanishes for $\vec{x} \in \tilde{\Omega}_i$, \hat{L}_{ii} contains a factor $\exp(-\Delta \hat{V}_i/T)$ where

$$\Delta \hat{V}_i = \min_{x \in \Omega_i \setminus \tilde{\Omega}_i} V(\vec{x}) - V(\vec{x}_i).$$
(9)

Let $\Delta \hat{V} = \min_i \Delta \hat{V}_i$. Then this shows that $-\hat{L}$ has N eigenvalues decay at least as fast as $\exp(-\Delta \hat{V}/T)$ to zero when T tends to zero.

Let us now consider L restricted to Ω_i but with reflecting boundary conditions on $\partial\Omega_i$. With these boundary conditions -L has one eigenvalue 0 and typically other eigenvalues that are much larger than $\exp(-\Delta \hat{V}/T)$. These eigenvalues describe the time scale for a decay within the potential well around the minimum at \vec{x}_i . On the original domain Ω one can therefore expect that -L has one eigenvalue 0 and N-1 eigenvalues that are of the order of $\exp(-c/T)$ and that all other eigenvalues of -L are much larger if T is sufficiently small. Let us denote the smallest N eigenvalues of -L by λ_i , $i=0,\ldots,N-1$, $\lambda_0=0$, and the corresponding right eigenfunctions by $g_i(\vec{x})$. Then, for time scales large compared to the intrawell relaxation times of the potential $V(\vec{x})$, we can make the ansatz

$$\rho(\vec{x},t) = \sum_{i=0}^{N-1} \rho_i(t) g_i(\vec{x})$$
(10)

for the solution of the Fokker-Planck equation $\partial \rho / \partial t = L\rho$. This ansatz should describe the long-time behavior of the system quite well. It is valid on time scales being large compared to the intrawell relaxation times in the different potential wells of the potential. $n_i(t) = \int_{x \in \Omega_i} d^d x \rho(\vec{x}, t)$ is the probability to find the particle in the region Ω_i , i.e., close to the minimum \vec{x}_i . Let $n_{ij} = \int_{x \in \Omega_i} d^d x g_j(\vec{x})$. One has $n_i(t) = \sum_j \rho_j(t) n_{ij}$. For the derivative with respect to time one obtains

$$\frac{dn_i}{dt} = \int_{x \in \Omega_i} d^d x L \rho = \sum_j \rho_j \int_{x \in \Omega_i} d^d x L g_j$$
$$= -\sum_j \lambda_j \rho_j n_{ij} = \sum_k r_{ik} n_k, \qquad (11)$$

where the matrix $R = (r_{ij})_{i,j=0,...,N-1}$ is given by $R = -N^{-1}\Lambda N$. *N* is the matrix $(n_{ij})_{i,j=0,...,N-1}$ and Λ is a diagonal matrix with the entries λ_j , j = 0, ..., N-1. The differential equation $dn_i/dt = \sum_k r_{ik}n_k$ is the rate equation for this problem. By construction, the eigenvalues of *R* are $-\lambda_i$.

Unfortunately, the definition of R is not useful for an explicit calculation. If one is able to solve the eigenvalue problem of L, there is no need to describe the problem by a rate equation. The rate equation is useful in situations, where it is not possible to obtain analytic results from the Fokker-Planck equation. Fortunately, approximate values for the rates r_{ij} often yield very good results. In the one-dimensional case the usual Kramers' rate is sufficiently good.

By construction, R is equivalent to the Fokker-Planck operator L projected onto the space of its eigenfunctions with N

largest eigenvalues. This means that dynamical properties on time scales corresponding to eigenvalues or larger can be well described by *R*. Properties on time scales $|\lambda_N|^{-1}$ or smaller cannot be calculated using *R*. λ_N sets the time scale for relaxation within the potential wells.

2. Rate equations for fluctuating potentials

Under which conditions is it possible to describe a problem with a fluctuating potential by rate equations? It is clear that one must be able to define n_i for a fluctuating potential. This is possible if the stationary distribution is strongly peaked at the minima \vec{x}_i . But in general one can have a situation, where the positions \vec{x}_i of the minima fluctuate as well. If this is the case, the probability distribution will not necessarily be peaked. Let us discuss that using a simple example: A one-dimensional dichotomously fluctuating potential.

For this example, z takes two values, ± 1 with equal probability. M_z has the eigenvalues 0 and $-\tau^{-1}$. In the eigenbasis of M_z the potential V(x,z) has the form

$$\begin{pmatrix} \langle V \rangle(x) & \Delta V(x) \\ \Delta V(x) & \langle V \rangle(x) \end{pmatrix},$$
 (12)

where $\langle V \rangle(x) = [V(x,+)+V(x,-)]/2$, $\Delta V(x) = [V(x,+)-V(x,-)]/2$. Let $f(x) = -d\langle V \rangle(x)/dx$, $\Delta f(x) = -d\Delta V(x)/dx$. Let us calculate the stationary probability $p_0(x) = \int dz \ \rho(x,z)$. The stationary joint probability distribution can be written as $\rho(x,z) = p_0(x) \phi_0(z) + p_1(x) \phi_1(z)$. Inserting this into the stationary Fokker-Planck equation yields

$$\frac{d}{dx} \left(f + z\Delta f - T\frac{d}{dx} \right) [p_0(x)\phi_0(z) + p_1(x)\phi_1(z)] + \tau^{-1}p_1(x)\phi_1(z) = 0.$$
(13)

This yields two equations for $p_0(x)$ and $p_1(x)$ corresponding to the two coefficients of $\phi_0(z)$ and $\phi_1(z)$, which both have to vanish. Eliminating $p_1(x)$ yields a single equation for $p_0(x)$, which can be written as

$$\left[1 + \tau \frac{d}{dx} \left(f - T \frac{d}{dx}\right)\right] (\Delta f)^{-1} \left(f - T \frac{d}{dx}\right) p_0(x) + \tau \frac{d}{dx} \Delta f p_0(x) = 0.$$
(14)

This equation can be derived directly from the Fokker-Planck equation. For T=0 it can be solved explicitly. One obtains

$$p_0(x) = C \frac{\Delta f}{f^2 - \Delta f^2} \exp\left(-\frac{1}{\tau} \int dx \frac{f}{f^2 - \Delta f^2}\right).$$
(15)

The integrand $f/(f^2 - \Delta f^2)$ diverges at the minima of $V(x, \pm)$. If the position of the minima does not change, $p_0(\vec{x})$ is a sum of δ functions located at the minima with appropriate

weights. With changing positions of the minima, the situation becomes different. In that case the particle moves between the minima. The most simple example to illustrate that is $V(x,\pm) = \frac{1}{2} \gamma_{\pm} (x-x_{\pm})^2$, $\gamma_{\pm} > 0$. Let $x_+ > x_-$. Then the support of $p_0(x)$ is the interval $[x_-, x_+]$. For $p_0(x)$ one obtains

$$p_{0}(x) = C \left[\frac{1}{\gamma_{+}(x_{+}-x)} + \frac{1}{\gamma_{-}(x-x_{-})} \right] \\ \times (x_{+}-x)^{1/(2\tau\gamma_{+})} (x-x_{-})^{1/(2\tau\gamma_{-})}$$
(16)

with some normalization constant *C*. The distribution function $p_0(x)$ is not peaked at x_{\pm} . It is clear that for small *T* >0 one obtains a similar result. This simple example shows that in a situation where the position of the minima depends on *z*, the stationary distribution of the system is not necessarily peaked at the minima of the potential.

Let us now assume that the potential V(x,z) has minima at $\vec{x_i}$ and that the positions $\vec{x_i}$ of the minima do not depend on z. For a fixed potential, rate equations can be used for time scales larger than the intrawell relaxation times. It is therefore clear that for potential fluctuations being faster than the intrawell relaxation times, one cannot use rate equations. If the potential fluctuations are slower than the intrawell relaxation times, and if the position of the minima does not fluctuate, one is in an adiabatic situation. The particle is able to follow the fluctuations of the potential as long as it stays close to a minimum. This means that one can calculate the rates for fixed z. The rates become functions of z. The system is thus described by rate equations with fluctuating rates,

$$\frac{dn_i}{dt} = \sum_j r_{ij}(z(t))n_j.$$
(17)

These rate equations are considerably simpler than the Fokker-Planck equation with a fluctuating potential. On the other hand, the validity of these equations is not entirely clear. We used only heuristic arguments to obtain the rate equations. In the following section we will compare results obtained from the rate equations with results from the Fokker-Planck equation. This allows us to show, in which range of parameters the rate equations can be used.

C. Elimination of z(t)

If one is interested in stationary properties or in the longtime behavior of the system, it is useful to derive effective rates for the transition from one state to another. The effective rates do not depend on z. They can be calculated by eliminating z(t). If the correlation time τ of z(t) is small, the elimination process is a standard elimination of a fast variable. To lowest order in τ the effective rates are simply given by the average of the rates $r_{ij}(z)$. To derive effective rates, let us transform the rate equation

$$\frac{d\vec{n}}{dt} = R\vec{n} \tag{18}$$

into a Fokker-Planck equation for the joint probability density $p(\vec{n},z,t)$,

$$\frac{\partial p(\vec{n},z,t)}{\partial t} = \vec{\nabla}_{\vec{n}} \cdot [R\vec{n} p(\vec{n},z,t)] + M_z p(\vec{n},z,t).$$
(19)

Since $\sum_i n_i = 1$ for all t, $p(\vec{n}, z, t)$ contains a factor $\delta(\sum_i n_i - 1)$. To calculate average values for \vec{n} , it is sufficient to consider the stationary case. Let

$$p(\vec{n},z) = p_0(\vec{n})\phi_0(z) + \bar{p}(\vec{n},z)$$
(20)

where

$$\int dz \, \bar{p}(\vec{n}, z) = 0, \qquad (21)$$

$$\int d^N n \, \bar{p}(\vec{n},z) = 0.$$
(22)

Let $\langle R \rangle = \int dz R(z) \phi_0(z)$. The stationary Fokker-Planck equation yields

$$\vec{\nabla}_{\vec{n}} \cdot [\langle R \rangle \vec{n} \, p_0(\vec{n})] = -\vec{\nabla}_{\vec{n}} \int dz \, R \vec{n} \, \vec{p}(\vec{n}, z) \tag{23}$$

and

$$\bar{p}(\vec{n},z) = -M_z^{-1} \vec{\nabla}_{\vec{n}} R \vec{n} \left[p_0(\vec{n}) \phi_0(z) + \bar{p}(\vec{n},z) \right].$$
(24)

Here, M_z^{-1} is the generalized inverse, it obeys $M_z^{-1}\phi_n(z) = -\lambda_n^{-1}\phi_n(z)$ for n > 0, $M_z^{-1}\phi_0(z) = 0$. Equation (23) yields

$$\langle R \rangle \vec{n} p_0(\vec{n}) = -\int dz \, R \vec{n} \, \vec{p}(\vec{n},z) + \vec{j}(\vec{n}) \tag{25}$$

with $\vec{\nabla}_{\vec{n}}\vec{j}(\vec{n}) = 0$. Let $\langle \vec{n} \rangle = \int d^N n \, \vec{n} p_0(\vec{n})$. Since

$$\int d^{N}n \,\vec{j}(\vec{n}) = \left\langle \frac{d\vec{n}}{dt} \right\rangle = 0 \tag{26}$$

one has

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$$R\rangle\langle\vec{n}\rangle = \sum_{l=0}^{\infty} \int dz \int d^{N}n(-1)^{l}R\vec{n}(M_{z}^{-1}\vec{\nabla}_{\vec{n}}R\vec{n})^{l+1}$$
$$\times p_{0}(\vec{n})\phi_{0}(z)$$
$$= -\sum_{l=0}^{\infty} \int dz \int d^{N}n(RM_{z}^{-1})^{l+1}R\vec{n}p_{0}(\vec{n})\phi_{0}(z).$$
(27)

The first line in Eq. (27) is obtained by solving Eq. (24) iteratively and inserting the result in Eq. (25). Integrating by parts yields the second line. Since R and M_z do not depend on \vec{n} , the integration over \vec{n} yields the average $\langle \vec{n} \rangle$. The sum over l gives finally

$$\int dz (1 - RM_z^{-1})^{-1} R\phi_0(z) \langle \vec{n} \rangle = 0.$$
 (28)

The entries of the matrix

$$R_{\rm eff} = \int dz (1 - RM_z^{-1})^{-1} R \phi_0(z)$$
 (29)

are the effective rates, the negative inverse of its diagonal matrix elements are the mean escape times out of the states *i*. The expansion in Eq. (27) can be viewed as a usual τ expansion. Expanding the right-hand side of Eq. (29) yields the τ expansion for the effective rates. But since $(1 - RM_z^{-1})$ is positive definite, the right-hand side of Eq. (29) is well defined for any τ . Whereas it is not possible to speak of an effective barrier for large τ [5], it is possible to obtain effective rates. But, depending on the noise process and on R(z), it cannot be excluded that the smallest nonvanishing eigenvalue of $R_{\rm eff}$ tends to zero in the limit $\tau \rightarrow \infty$. This may happen if the fluctuations of the barrier heights in the potential can be arbitrarily large, so that infinite barrier heights occur. In that case the solution for $\langle n \rangle$ has to be obtained from the solution for finite τ and the limit $\tau \rightarrow \infty$ has to be taken afterwards since Eq. (28) does not have a unique solution for $\tau = \infty$. As long as the eigenvalue 0 of R(z) is nondegenerate for all z in the support of $\phi_0(z)$, the matrix $R_{\rm eff}$ has a nondegenerate eigenvalue 0 and Eq. (28) has a unique solution. This is shown below in Sec. IV.

D. Applicability of effective rate equations

As we already discussed in Sec. II A, rate equations describe only properties on long-time scales. If one wants to investigate dynamical properties on characteristic time scales of the order of the typical relaxation within a potential well, one cannot use rate equations. If one wants to study a system with a fluctuating potential that fluctuates fast (i.e., one uses a time scale of the order of the typical relaxation within a potential well) one has to use the Fokker-Planck equation. This is true for equations with fluctuating rates as well as for the effective rate equations we derived in the last section.

Since in many applications details of the potential are unknown, one often uses rate equations. Especially in biological applications, e.g., for the description of membrane proteins, one usually uses rate equations. A typical example is the well-known paper by Petracchi et al. [11], who studied the effect of time dependent electric fields on membrane proteins experimentally and used rate equations to interpret their experimental findings. They were able to explain most of their results using rate equations and a very simple ansatz for the rates. In their experimental studies they found an interesting effect called phase anticipation. The rate equations on the other hand do not show this effect. They argued that phase anticipation occurs due to some special biological effect. From our discussion about the validity of rate equations it is clear that the effect of phase anticipation, which occurs on short time scales, cannot be described by rate equations. We will come back to this point in our conclusions.

In the following section we discuss some simple, onedimensional examples. The most simple case, the dichotomous two-state model, is exactly solvable. For this case and for more general examples we compare the solution of the rate equations with the solution of the Fokker-Planck equation. For all cases, rate equations yield very good results if the potential fluctuations are slow compared to the relaxation within the potential wells.

III. SOLUTIONS FOR SPECIAL CASES

A. The dichotomous two-state model

The most simple model with a fluctuating potential is a dichotomous two-state model. The potential $V(\vec{x},z)$ has two minima at \vec{x}_1 and \vec{x}_2 . z(t) is a dichotomous process. Such a model has two states and corresponding rates r_{ij} , i,j=1,2 for transitions between them. For this model we let $r_1:=r_{21}$ = $-r_{11}$, $r_2:=r_{12}=-r_{22}$. The rate equations are

$$\frac{dn_1}{dt} = -r_1 n_1 + r_2 n_2, \tag{30}$$

$$\frac{dn_2}{dt} = -r_2n_2 + r_1n_1. \tag{31}$$

Using the normalization condition $n_1 + n_2 = 1$ one obtains a single differential equation

$$\frac{dn_1}{dt} = -(r_1 + r_2)n_1 + r_2. \tag{32}$$

In the following the index of n_1 will be dropped. The Fokker-Planck equation for p(n,z,t) is

$$\frac{\partial p(n,z,t)}{\partial t} = \frac{\partial}{\partial n} [(r_1 + r_2)n - r_2] p(n,z,t) + M_z p(n,z,t).$$
(33)

In general it is possible to expand the stationary solution using the eigenfunctions of M_z ,

$$p(n,z) = \sum_{k} p_k(n)\phi_k(z).$$
(34)

For the matrix elements of $r_i(z)$ one has the representation

$$r_i(z)\phi_k(z) = \sum_l r_{ikl}\phi_l(z).$$
(35)

It is useful to define

$$r_{kl} := r_{1kl} + r_{2kl} \,. \tag{36}$$

The Fokker-Planck equation for the stationary distribution p(n,z) has now the form

$$\frac{\partial}{\partial n} [(r_1 + r_2)n - r_2]p(n,z) = \frac{\partial}{\partial n} \sum_{kl} (r_{kl}n - r_{2kl})p_k \phi_l$$
$$= \sum_l \lambda_l p_l \phi_l.$$
(37)

The two matrices $R_i = (r_{ikl})$ have the same eigenvectors and only positive eigenvalues. This must be true for $R = R_1$ $+R_2$ as well. In the dichotomous case R and R_2 are 2×2 matrices. The eigenvalues are r_{\pm} and $r_{2\pm}$. $r_i(z)$ takes the two values $r_{i\pm}$ with the probabilities p_{\pm} . Let n_{\pm} $= r_{2\pm}/r_{\pm}$. I assume $n_+ > n_-$. The explicit form is of R is

$$R = \begin{pmatrix} r_{00} & r_{01} \\ r_{10} & r_{11} \end{pmatrix} = \begin{pmatrix} p_{+}r_{+} + p_{-}r_{-} & \sqrt{p_{+}p_{-}(r_{+} - r_{-})} \\ \sqrt{p_{+}p_{-}(r_{+} - r_{-})} & p_{+}r_{-} + p_{-}r_{+} \end{pmatrix}$$
(38)

and similarly for R_2 . $p_0(n)$ vanishes for $n < n_-$ and for $n > n_+$. The two equations (37) are

$$(r_{00}n - r_{200})p_0 + (r_{10}n - r_{210})p_1 = 0, (39)$$

$$\frac{d}{dn} [(r_{01}n - r_{201})p_0 + (r_{11}n - r_{211})p_1] = \lambda_1 p_1.$$
(40)

The integration constant in Eq. (39) vanishes since the average of dn/dt vanishes. It may be used to eliminate p_1 , one obtains

$$\frac{d}{dn}Ap_0 = -\lambda_1 Bp_0, \tag{41}$$

where

$$A = \frac{f(n)}{r_{10}n - r_{210}},\tag{42}$$

$$B = \frac{r_{00}n - r_{200}}{r_{10}n - r_{210}}.$$
(43)

The solution of this equation is

$$p_0 = \frac{C}{A} \exp\left(-\lambda_1 \int \frac{B}{A} dn\right),\tag{44}$$

where C is a normalization constant and f(n) is given by

$$f(n) \coloneqq (r_{01}n - r_{201})(r_{10}n - r_{210}) - (r_{00}n - r_{200})(r_{11}n - r_{211})$$

= $-\det(Rn - R_2) = -\det(R)(n - n_+)(n - n_-).$ (45)

The integral in Eq. (44) can be calculated. The explicit solution is

$$p_0(n) = C(n - \tilde{n})(n - n_-)^{\alpha_- - 1}(n_+ - n)^{\alpha_+ - 1}, \quad (46)$$

where

$$\alpha_{-} = \frac{\lambda_{1} r_{00}(n_{0} - n_{-})}{r_{+} r_{-}(n_{+} - n_{-})}, \qquad (47)$$

$$\alpha_{+} = \frac{\lambda_{1} r_{00} (n_{+} - n_{0})}{r_{+} r_{-} (n_{+} - n_{-})}, \qquad (48)$$

$$\tilde{n} = \frac{r_{210}}{r_{10}} = \frac{r_{2+} - r_{2-}}{r_{+} - r_{-}}.$$
(49)

One can show that either $\tilde{n} < n_-$ or $\tilde{n} > n_+$, so that $p_0(n)$ is non-negative for $n \in [n_-, n_+]$ as it should be. For $\lambda_1 \rightarrow \infty$ one obtains the expected result $p_0(n) = \delta(n - n_0)$ where n_0 $=r_{200}/r_{00}$. For small λ_1 , $p_0(n)$ has algebraic singularities at the boundaries n_{\pm} . For small but finite λ_1 , the singularities remain integrable. In the limit $\lambda_1 \rightarrow 0$, $p_0(n)$ tends to a sum of two δ functions (with appropriate weights according to the dichotomous process) located at $n_{\pm} = r_{2\pm}/r_{\pm}$,

$$p_0(n) = p_+ \delta(n - n_+) + p_- \delta(n - n_-).$$
 (50)

The average value \overline{n} of n can be calculated explicitly. One obtains

$$\begin{split} \bar{n} &= \int_{n_{-}}^{n_{+}} np_{0}(n) dn \\ &= \frac{\int_{n_{-}}^{n_{+}} n(n-\tilde{n})(n-n_{-})^{\alpha_{-}-1}(n_{+}-n)^{\alpha_{+}-1} dn}{\int_{n_{-}}^{n_{+}} (n-\tilde{n})(n-n_{-})^{\alpha_{-}-1}(n_{+}-n)^{\alpha_{+}-1} dn} \\ &= \frac{n_{-}(n_{-}-\tilde{n})B(\alpha_{+},\alpha_{-}) + (2n_{-}-\tilde{n})(n_{+}-n_{-})B(\alpha_{+},\alpha_{-}+1) + (n_{+}-n_{-})^{2}B(\alpha_{+},\alpha_{-}+2)}{(n_{-}-\tilde{n})B(\alpha_{+},\alpha_{-}) + (n_{+}-n_{-})B(\alpha_{+},\alpha_{-}+1)} \\ &= \frac{n_{-}(n_{-}-\tilde{n}) + (2n_{-}-\tilde{n})(n_{0}-n_{-}) + (n_{+}-n_{-})(n_{0}-n_{-})\frac{(\alpha_{-}+1)}{(\alpha_{+}+\alpha_{-}+1)}}{n_{0}-\tilde{n}}. \end{split}$$
(51)

One can easily show that $\overline{n} = n_0$ in the limit $\lambda_1 \rightarrow \infty$, and $\overline{n} = p_+ n_+ + p_- n_-$ for $\lambda_1 = 0$. Finally one obtains

$$\bar{n} = n_0 + (p_+ n_+ + p_- n_- - n_0) \frac{\tau}{\bar{\tau} + \tau}, \qquad (52)$$

where $\tau = \lambda_1^{-1}$ and $\overline{\tau} = r_{00}/r_+r_- = p_+/r_- + p_-/r_+$. Thus, \overline{n} depends monotonously on τ . The expression (52) as well as the formula for $p_0(n)$, Eq. (46), have been mentioned without a detailed derivation in [9].

B. The general two-state model

For a general noise process z(t) the effective rates are

$$\bar{r}_i = \int dz [1 - r(z)M_z^{-1}]^{-1} r_i(z)\phi_0(z), \qquad (53)$$

where again $r_1 := r_{21} = -r_{11}$, $r_2 := r_{12} = -r_{22}$. This yields directly

$$\bar{n} = \frac{\int dz (1 - r(z)M_z^{-1})^{-1} r_2(z)\phi_0(z)}{\int dz (1 - r(z)M_z^{-1})^{-1} r(z)\phi_0(z)}.$$
(54)

Whether or not this expression can be calculated explicitly depends on the generator M_z of the noise process z(t). A case where the calculation is possible is a kangaroo process. A kangaroo process is a process where all nonvanishing eigenvalues of M_z are equal. In that case, Eq. (54) simplifies to

$$\bar{n} = \frac{\int dz [1 + \tau r(z)]^{-1} r_2(z) \phi_0(z)}{\int dz [1 + \tau r(z)]^{-1} r(z) \phi_0(z)}.$$
(55)

The derivation of the last result is not straightforward: One cannot simply replace M_z by $-\tau^{-1}$. But it turns out, that the result is the same, since additional factors in the denominator and in the numerator cancel each other. Equation (55) is valid for a dichotomous process as well and yields directly Eq. (52).

C. Comparison with the Fokker-Planck equation

As already mentioned, the validity of rate equations for fluctuating potentials is less clear than for fixed potentials. It is therefore important to compare results for rate equations with exact solutions of the Fokker-Planck equation. They can be obtained easily for one-dimensional piecewise linear potentials. This has been shown in detail for periodic potentials in [12,13]. In the present case the procedure is similar, let us therefore describe it only briefly. For this comparison, we restrict ourselves to the stationary distribution $p_0(x) = \int dz \rho(x,z)$ and to the mean escape times for the particle sitting in a minimum of the potential.

We first divide the domain Ω , on which the Fokker-Planck equation is defined, into a set of intervals I_i , $i = 1, \ldots, N_I$. We assume that the force $f(x,z) = -\partial V(x,z)/\partial x$ is constant on each interval I_i ,

$$f(x,z) = f_i(z) \quad \text{if } x \in I_i.$$
(56)

As a second step, we expand the stationary solution of the Fokker-Planck equation using the eigenbasis of M_z ,

$$\rho(x,z) = p_0(x)\phi_0(z) + \sum_{k>0} p'_k(x)\phi_k(z).$$
 (57)

On the interval I_i , the Fokker-Planck equation yields

$$f_{i}(z)p_{0}'(x)\phi_{0}(z) + Tp_{0}''(x)\phi_{0}(z) + \sum_{k>0} p_{k}''(x)f_{i}(z)\phi_{k}(z) + T\sum_{k>0} p_{k}'''(x)\phi_{k}(z) - \sum_{k>0} \lambda_{k}p_{k}'(x)\phi_{k}(z) = 0.$$
(58)

This equation can be integrated once with respect to x. Furthermore, expression of the form $f_i(z)\phi_k(z)$ can be expanded,

$$f_i(z)\phi_k(z) = \sum_{k'} f_{k,k'}^{(i)}\phi_{k'}(z).$$
 (59)

One then obtains

$$f_{0,0}^{(i)} p_0(x) + \sum_{k'>0} f_{k',0}^{(i)} p_{k'}'(x) + T p_0'(x) = 0, \qquad (60)$$

$$f_{0k}^{(i)}p_0(x) + \sum_{k'>0} f_{k',k}^{(i)}p_{k'}'(x) + Tp_k''(x) = \lambda_k p_k(x), \quad k > 0.$$
(61)

The coefficients $f_{k,k'}^{(i)}$ and therefore the final set of equations for $p_k(x)$ do not depend on the representation of the noise process z(t) as it should be. If the space of functions $\phi_k(z)$ has finite dimension N, Eqs. (60) and (61) form is a set of N differential equations with constant coefficients, which can be solved explicitely. One obtains a simple eigenvalue problem for an $N \times N$ matrix, which can be solved numerically. A general solution of this problem is a linear combination of the N different solutions of the eigenvalue problem. The remaining problem is thus to determine the $N_I N$ coefficients in these linear combinations. They are determined by the normalization of $p_0(x)$, by the continuity of $p_0(x)$, and by the continuity of $p_k(x)$ and $p'_k(x)$ for $k \ge 0$. This is a problem of solving $N_I N$ linear equations for $N_I N$ unknown coefficients, which can again be done easily numerically as long as $N_I N$ is not too large.

In Fig. 1 we show explicit results for a one-dimensional dichotomous two-state model. The potential $V(x,z) = \langle V \rangle(x) + z \Delta V(x)$ has two minima. z(t) is a dichotomous



FIG. 1. The mean occupancy in the left minimum of a piecewise linear, dichotomously fluctuating potential as a function of τ for various *T*. The parameters are *T*=0.2, 0.4, 0.6, 0.8, 1.0, the average force takes the values 10, 3/2, -3/2, 4/3, -4, -10; Δf takes the values 0, 1/4, -1/4, 2/3, -2, 0. The force jumps at *x*=-4, -2, 0, 3, 4. The maximal value of the mean occupancy decreases with increasing *T*. The dashed lines are the results from the effective rate equation.

process that takes the two values ± 1 . In this case N=2 and $N_I=6$. The plot shows the probability to find the particle in the left potential well as a function of τ . It can be compared to \overline{n} , which can be calculated from the effective rate equations. A similar comparison for a different, dichotomously fluctuating potential was shown in [9]. The results from the rate equation yield an accurate approximation for sufficiently large τ , as expected. The intrawell relaxation time can be estimated from the time the particle needs to reach a minimum of the potential if it starts close to the maximum. In our units it is somewhat larger than unity. The approximation becomes better for smaller temperature. The behavior is similar for other dichotomously fluctuating potentials with two minima.

For other stochastic processes z(t) one can draw similar conclusions. Let us consider as an example a potential $V(x,z) = \langle V \rangle(x) + z \Delta V(x)$, where now z(t) is a sum of N - 1 dichotomous processes, where each dichotomous process takes the values $\pm 1/\sqrt{N}$. This means that z(t) takes the values $(-N+2n+1)/\sqrt{N}$, $n=0,\ldots,N$ with the probability $p_n = 2^{-N} {N \choose n}$. In the limit $N \rightarrow \infty$ one obtains an Ornstein-Uhlenbeck process [14]. For $\langle V \rangle$ and ΔV we take the same values as in Fig. 1. Results for the mean occupancy are shown in Fig. 2 together with the corresponding results from the rate equation, obtained from Eq. (54). As for the dichotomous process, the results from the rate equation for not too small values of τ .

Similar results can be obtained for other noise processes as well. For kangaroo processes, the mean occupancy has been obtained explicitly, Eq. (55). The corresponding results for the mean occupancy are similar to what has been obtained so far. If one takes, for instance, kangaroo processes



FIG. 2. The mean occupancy in the left minimum for a fluctuating potential as in Fig. 1. f and Δf are the same as in Fig. 1, T = 0.4. The different curves are for sums of N-1 dichotomous processes, $N=2, \ldots, 6$. For increasing N the maximum of the curve is shifted to the left. The dashed lines are the results from the effective rate equation, the results for N=4,5,6 cannot be distinguished.

with the same $\phi_0(z)$ as for the sums of dichotomous processes in Fig. 2, the quantitative results are almost the same, they differ by less than 1%.

The results in Fig. 2 show that already for small N one obtains essentially the result for the Ornstein-Uhlenbeck process. The convergency to the Ornstein-Uhlenbeck process is very fast. For N>6 the curves lie on top of the curve for N=6. Clearly, the convergency depends on the parameters of the system. For smaller temperatures we observe a slower convergency.

It is possible to consider potentials with more than two minima as well. The corresponding rate equations describe systems with more than two states. For simple systems with few minima and, e.g., a dichotomously fluctuating potential, one can again compare the results from the Fokker-Planck equation with the results from the rate equations. The generic behavior of such systems is again described by the effective rates as long as τ is not too small. Whether or not Eq. (29) can be used to calculate the effective rates analytically depends on the structure of the problem. For complicated noise processes or for potentials with many minima, $R_{\rm eff}$ can only be obtained numerically.

IV. MATHEMATICAL ASPECTS

In the main part of this work we investigate the stationary solution of a Fokker-Planck equation (5) with a fluctuating potential, assuming that such a stationary solution exists. The aim of this section is to prove the existence of a stationary solution under certain assumptions that will be sufficient for our purpose. For the special case of a one-dimensional, dichotomously fluctuating potential, Pechukas and Ankerhold [10] proved the existence of a stationary solution under some special assumptions on the potential. The result of this section can be understood as a generalization of their results.

A. The discrete case

Let us first discuss the discrete case, which is very simple but shows some interesting aspects. Let us assume that a system has N states and that it can move from one state to another. The behavior is described by a rate equation

$$\frac{dp_i}{dt} = \sum_j r_{ij} p_j \,. \tag{62}$$

We now assume that the rates fluctuate between N_r different states $\alpha = 1, \ldots, N_r$. The fluctuation between these states is determined by a matrix $M = (m_{\alpha\beta})_{\alpha,\beta=1,\ldots,N_r}$ that contains the rates for the fluctuations. The probability $p_{i\alpha}$ of the fluctuating system to be in the state given by *i* and α is determined by

$$\frac{dp_{i\alpha}}{dt} = \sum_{j} r_{ij\alpha} p_{j\alpha} + \sum_{\beta} m_{\alpha\beta} p_{i\beta}.$$
(63)

The stationary solution of this equation, if it exists, must be the eigenvector of the matrix

$$\hat{R} = \begin{pmatrix} R_1 + m_{11}I & m_{12}I & m_{13}I & \cdots & m_{1N_r}I \\ m_{21}I & R_2 + m_{22}I & m_{23}I & \cdots & & \\ m_{31}I & m_{32}I & R_3 + m_{33}I & \ddots & \vdots \\ \vdots & & \ddots & \ddots & m_{N_r-1,N_r}I \\ m_{N_r}II & \cdots & & m_{N_r,N_r-1}I & R_{N_r} + m_{N_rN_r}I \end{pmatrix}$$
(64)

to eigenvalue 0. Here, *I* is the $N \times N$ unit matrix. The matrices R_{α} are rate matrices. As a consequence,

$$r_{ii\alpha} = -\sum_{j \neq i} r_{ji\alpha} \tag{65}$$

and $r_{ij\alpha} \ge 0$ for $i \ne j$. Furthermore, M is a rate matrix, $m_{\alpha\alpha} = -\sum_{\beta \ne \alpha} m_{\beta,\alpha}$, and $m_{\alpha,\beta} \ge 0$ for $\alpha \ne \beta$. Therefore, all the off-diagonal matrix elements of \hat{R} are non-negative. I assume that starting from some state of the system, any other state can be reached dynamically. This means that \hat{R} is irreducible.

Then, one can introduce a constant *c* that is larger than the modulus of any diagonal matrix element of \hat{R} . Let \hat{I} be the $N_r N \times N_r N$ unit matrix. Then the matrix $\hat{R} + c\hat{I}$ is irreducible and has only positive matrix elements. As a consequence, the Perron-Frobenius theorem applies. The eigenvalue with the largest modulus of $\hat{R} + c\hat{I}$ is nondegenerate, real, and positive, and the corresponding left and right eigenvectors have only non-negative entries. Now, the vector $(1,1,\ldots,1)$ is a left eigenvalue with the largest modulus of $\hat{R} + c\hat{I}$ with the eigenvalue *c*. Therefore *c* is the eigenvalue with the largest modulus of $\hat{R} + c\hat{I}$. This shows that 0 is the eigenvalue of \hat{R} with the largest real part and that its right eigenvalue has only non-negative entries. It is the stationary solution of Eq. (63). One can show that any solution of the rate equation (63) tends to the stationary solution for $t \rightarrow \infty$.

This proof is very simple, it is based on the fact that rates are non-negative and on the Perron-Frobenius theorem. Let us now take a look at the original Fokker-Planck equation (5). For simplicity, we discuss only the case where the potential fluctuates between N_r different states. z takes N_r different values. Let us assume that the Fokker-Planck equation is defined on a finite open domain Ω with reflecting boundary conditions on the boundary $\partial \Omega$. If one introduces some kind of coarse graining, i.e., a partition of Ω into N small parts Ω_i , it is then possible to apply the above result and a stationary solution exists. This holds for any partition of Ω into small parts. Therefore one should expect that a stationary solution of the Fokker-Planck equation (5) exists as well. We will show this in the next section using a different approach. One has to show that the Fokker-Planck operator has a unique right eigenstate with eigenvalue 0 and that any solution of the Fokker-Planck equations tends to that solution for $t \rightarrow \infty$. Since we are dealing with a differential operator and not with a matrix, it is not possible to use a Perron-Frobenius type argument.

B. The continuous case

Let us first show that the ratio of any two solutions of the Fokker-Planck equation tends to unity in the limit $t \rightarrow \infty$. As for the discrete case we assume that z takes N_r values z_{α} , $\alpha = 1, \ldots, N_r$ and we let

$$\rho(\vec{x}, z_{\alpha}, t) = p_{\alpha}(\vec{x}, t). \tag{66}$$

Furthermore, we again assume that the Fokker-Planck equation (5) is defined on a finite open domain Ω with reflecting boundary conditions on the boundary $\partial \Omega$. The Fokker-Planck equation can be written in the form

$$\frac{\partial p_{\alpha}}{\partial t} = \vec{\nabla} [\vec{\nabla} V_{\alpha}(\vec{x}) + T\vec{\nabla}] p_{\alpha} + \sum_{\beta} m_{\alpha\beta} p_{\beta}$$
$$= L_{\alpha} p_{\alpha} + \sum_{\beta} m_{\alpha\beta} p_{\beta}.$$
(67)

Let $p_{1\alpha}$, $p_{2\alpha}$ be two positive solutions of this equation. We define

$$H(t) = \sum_{\alpha} \int d^{d}x \, p_{1\alpha} \ln(p_{1\alpha}/p_{2\alpha}).$$
(68)

This function is often called Kullback information and has been introduced by Kullback [15,16]. By standard methods one can show that (i) $H(t) \ge 0$ and (ii) $dH/dt \le 0$. As a consequence H tends to zero for $t \rightarrow \infty$ and therefore $p_{1\alpha}/p_{2\alpha}$ tends to 1. The proof is similar to the one presented in chapter 6 of [17]. First, one shows

$$H(t) = \sum_{\alpha} \int d^d x \, p_{2\alpha} \left(\frac{p_{1\alpha}}{p_{2\alpha}} \ln \frac{p_{1\alpha}}{p_{2\alpha}} - \frac{p_{1\alpha}}{p_{2\alpha}} + 1 \right) \ge 0, \quad (69)$$

which follows from $r \ln r - r + 1 \ge 0$ for $r \ge 0$. Now, one can calculate

$$\frac{dH}{dt} = \sum_{\alpha} \int d^{d}x \left[\frac{dp_{1\alpha}}{dt} \ln(p_{1\alpha}/p_{2\alpha}) - (p_{1\alpha}/p_{2\alpha}) \frac{dp_{2\alpha}}{dt} \right]$$

$$= \sum_{\alpha} \int d^{d}x \left[(L_{\alpha}p_{1\alpha})\ln(p_{1\alpha}/p_{2\alpha}) + \sum_{\beta} m_{\alpha\beta}p_{\beta}\ln(p_{1\alpha}/p_{2\alpha}) - \frac{p_{1\alpha}}{p_{2\alpha}} \frac{dp_{2\alpha}}{dt} \right]$$

$$= \sum_{\alpha} \int d^{d}x \left[p_{1\alpha}L_{\alpha}^{\dagger}\ln(p_{1\alpha}/p_{2\alpha}) + \sum_{\beta} m_{\alpha\beta}p_{\beta}\ln(p_{1\alpha}/p_{2\alpha}) - \frac{p_{1\alpha}}{p_{2\alpha}} \frac{dp_{2\alpha}}{dt} \right]. \quad (70)$$

The first term may be evaluated using

$$L_{\alpha}^{\dagger}\ln(p_{1\alpha}/p_{2\alpha}) = \frac{p_{2\alpha}}{p_{1\alpha}}L_{\alpha}^{\dagger}(p_{1\alpha}/p_{2\alpha}) - T\frac{p_{2\alpha}^{2}}{p_{1\alpha}^{2}}[\vec{\nabla}(p_{1\alpha}/p_{2\alpha})]^{2}.$$
(71)

This yields

$$\frac{dH}{dt} = -T \sum_{\alpha} \int d^d x \, p_{1\alpha} [\vec{\nabla} \ln(p_{1\alpha}/p_{2\alpha})]^2 + \int d^d x \sum_{\alpha\beta} m_{\alpha\beta} \left(p_{1\beta} \ln \frac{p_{1\alpha}}{p_{2\alpha}} - p_{2\beta} \frac{p_{1\alpha}}{p_{2\alpha}} \right).$$
(72)

The first term is clearly negative or 0. To investigate the second term, let us introduce

$$\Lambda_{\alpha\beta} = \delta_{\alpha\beta} + \epsilon m_{\alpha\beta}, \qquad (73)$$

where ϵ should be sufficiently small ($\epsilon^{-1} > \max_{\alpha} |m_{\alpha\alpha}|$, I will take the limit $\epsilon \rightarrow 0$ below). Let $\hat{p}_{i\alpha} = \sum_{\beta} \Lambda_{\alpha\beta} p_{i\beta}$. One can show that

$$\int d^{d}x \sum_{\alpha} \left(p_{1\alpha} \ln \frac{p_{1\alpha}}{p_{2\alpha}} - \hat{p}_{1\alpha} \ln \frac{\hat{p}_{1\alpha}}{\hat{p}_{2\alpha}} \right)$$
$$= \int d^{d}x \sum_{\alpha\beta} \Lambda_{\alpha\beta} p_{1\beta} \ln \frac{p_{1\beta} \hat{p}_{2\alpha}}{p_{2\beta} \hat{p}_{1\alpha}}$$
$$\geq \int d^{d}x \sum_{\alpha\beta} \Lambda_{\alpha\beta} p_{1\beta} \left(1 - \frac{p_{2\beta} \hat{p}_{1\alpha}}{p_{1\beta} \hat{p}_{2\alpha}} \right) = 0. \quad (74)$$

Expanding the left-hand side in powers of ϵ , one obtains

$$\int d^d x \sum_{\alpha\beta} m_{\alpha\beta} \left(p_{1\beta} \ln \frac{p_{1\alpha}}{p_{2\alpha}} - p_{2\beta} \frac{p_{1\alpha}}{p_{2\alpha}} \right) \leq 0$$
(75)

and therefore $dH/dt \le 0$. The limit $p_{\alpha}(\vec{x}, t \to \infty)$ of a positive solution of the Fokker-Planck equation is positive. If the potential $V_{\alpha}(\vec{x})$ is bounded within the domain Ω , and if Ω is connected, this limit yields the unique stationary solution of the Fokker-Planck equation since the Fokker-Planck operator does not depend on time. This shows the existence of the stationary solution of the Fokker-Planck equation.

V. OUTLOOK AND CONCLUSIONS

One main result of this paper is that the overdamped motion of a Brownian particle in a fluctuating potential can be described by kinetic rate equations if (i) the potential fluctuations are slower than the relaxation of the particle within a minimum of the potential and if (ii) the positions of the minima of the potential do not fluctuate. For temperatures small compared to typical barrier heights of the potential, the quantitative agreement of the two different descriptions is very good if one calculates stationary or quasistationary properties. This shows that the long-time behavior of a Brownian particle in a fluctuating potential has universal properties and does not depend on the details of the potential but only on the rates for the transition over the various fluctuating barriers. This is important, since in many realistic situations details of the potential are not known. For instance, in the case of a cell surface receptor or some other protein in a cell membrane, one knows eventually something about the stable or metastable conformations of the protein, but (almost) nothing about the potential that describes the energy of the deformations of the protein. Even if the potential is known, it is often simpler to solve a kinetic rate equation instead of a Fokker-Planck equation.

Starting from the kinetic rate equation with fluctuating rates, we derived a formula for effective rates that holds for general noise processes. This formula yields directly information on the stationary and long-time properties of the system. It can be evaluated for many noise processes. The special case of a dichotomous two-state system can be solved exactly and one obtains the distribution functions for the occupancy of the two states.

As a byproduct, we could show the existence of a station-

ary solution of the Fokker-Planck equation with a fluctuating potential under some general conditions and thereby generalize a result published recently by Pechukas and Ankerhold [10].

As a direct consequence of our calculations, one can investigate where rate equations can yield good results and where not. This is important since in many applications, mainly in biological situations, rate equations are used for a direct comparison with experiments. For instance, dynamical properties of membrane proteins are often described and investigated by rate equations. An interesting and well-known example is the effect of time-dependent electric fields on such proteins. It has been discussed it detail by Astumian and Robertson [18] using rate equations. In such systems one observes strong amplifications of weak signals. Moss [19] suggested that the reason may be stochastic resonance. Kruglikov and Dertinger [20] gave a qualitative discussion supporting the occurence of stochastic resonance in such systems using a time-dependent potential, but a description using rate equations would have been possible as well. As already mentioned, Petracchi et al. [11] studied the effect of time-dependent electric fields on membrane proteins experimentally. As a specific example they used a K⁺ channel. They measured various stationary probabilities and dynamical quantities to describe the statistics of the transitions in these systems. As an interesting result they found a phase anticipation: One of the transitions occurs with a negative phase shift compared to the stimulus. They compared their experimental results with numerical results for a simple twostate Markov model, described by a rate equation. The rate equations do not show the phase anticipation. The authors argued that this effect has a biological origin and discussed various possible hypotheses to explain it.

From our discussion it is clear that rate equations do describe the stationary or long-time behavior of a system quite well. Indeed, the stationary or quasistationary properties calculated by Petracchi et al. [11] agree very well with their experimental findings. The phase anticipation is a dynamical effect that occurs on shorter time scales. Therefore one cannot expect that it can be obtained using rate equations. Our calculations suggest that this effect depends on the details of the potential and not only on the rates, which are determined mainly by the barrier heights. The main problem is that in the case of a specific membrane protein the potential that describes the dynamics is not known. Therefore it is not possible to model the K⁺ channel using a time-dependent potential. But it would be of general interest to investigate whether or not and under which conditions a Brownian particle in a time-dependent potential shows the described phase anticipation. This is clearly beyond the scope of the present paper.

ACKNOWLEDGMENT

I wish to thank H. Dertinger for various interesting discussions on the effect of weak electromagnetic fields on cells.

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