# A SEMICLASSICAL APPROACH TO THE KRAMERS-SMOLUCHOWSKI EQUATION 

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#### Abstract

We consider the Kramers-Smoluchowski equation at a low temperature regime and show how semiclassical techniques developed for the study of the Witten Laplacian and Fokker-Planck equation provide quantitative results. This equation comes from molecular dynamics and temperature plays the role of a semiclassical paramater. The presentation is self-contained in the one dimensional case, with pointers to the recent paper [15] for results needed in higher dimensions. One purpose of this note is to provide a simple introduction to semiclassical methods in this context.


## 1. Introduction

The Kramers-Smoluchowski equation describes the time evolution of the probability density of a particle undergoing a Brownian motion under the influence of a chemical potential - see [1] for the background and references. Mathematical treatments in the low temperature regime have been provided by Peletier et al [16] using $\Gamma$-convergence, by Herrmann-Niethammer [11] using Wasserstein gradient flows and by Evans-Tabrizian [5].

The purpose of this note is to explain how precise quantitative results can be obtained using semiclassical methods developed by, among others, Bovier, Gayrard, Helffer, Hérau, Hitrik, Klein, Nier and Sjöstrand [2, 7, 8, 9, 10] for the study of spectral asymptotics for Witten Laplacians [17] and for Fokker-Planck operators. The semiclassical parameter $h$ is the (low) temperature. This approach is much closer in spirit to the heuristic arguments in the physics literature $[6,13]$ and the main point is that the Kramers-Smoluchowski equation is the heat equation for the Witten Laplacian acting on functions. Here we give a self-contained presentation of the one dimensional case and explain how the recent paper by the first author [15] can be used to obtain results in higher dimensions.

Let $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a smooth function. Consider the corresponding KramersSmoluchowski equation:

$$
\left\{\begin{array}{c}
\partial_{t} \rho=\partial_{x} \cdot\left(\partial_{x} \rho+\epsilon^{-2} \rho \partial_{x} \varphi\right)  \tag{1.1}\\
\rho_{\mid t=0}=\rho_{0}
\end{array}\right.
$$

where $\epsilon \in(0,1]$ denotes the temperature of the system and will be the small asymptotic parameter. Assume that there exists $C>0$ and a compact $K \subset \mathbb{R}^{d}$ such that for all


Figure 1. A one dimensional potential with interesting KramersSmoluchowski dynamics.
$x \in \mathbb{R}^{d} \backslash K$, we have

$$
\begin{equation*}
|\partial \varphi(x)| \geq \frac{1}{C}, \quad\left|\partial_{x_{i} x_{j}}^{2} \varphi\right| \leq C|\partial \varphi|^{2}, \quad \varphi(x) \geq C|x| \tag{1.2}
\end{equation*}
$$

Suppose additionally that $\varphi$ is a Morse function, that is, $\varphi$ has isolated and nondegenerate critical points. Then, thanks to the above assumptions the set $\mathcal{U}$ of critical points of $\varphi$ is finite. For $p=0, \ldots, d$, we denote by $\mathcal{U}^{(p)}$ the set of critical points of index $p$. Denote

$$
\begin{equation*}
\varphi_{0}:=\inf _{x \in \mathbb{R}^{d}} \varphi(x)=\inf _{\mathbf{m} \in \mathcal{U}^{(0)}} \varphi(\mathbf{m}) \quad \text { and } \quad \sigma_{1}:=\sup _{\mathbf{s} \in \mathcal{U}^{(1)}} \varphi(\mathbf{s}) . \tag{1.3}
\end{equation*}
$$

Thanks to (1.2), the sublevel set of $\sigma_{1}$ is decomposed in finitely many connected components $E_{1}, \ldots, E_{N}$ :

$$
\begin{equation*}
\left\{x \in \mathbb{R}^{d}, \varphi(x)<\sigma_{1}\right\}=\bigsqcup_{n=1}^{N} E_{n} \tag{1.4}
\end{equation*}
$$

We assume that

$$
\begin{equation*}
\inf _{x \in E_{n}} \varphi(x)=\varphi_{0}, \quad \forall n=1, \ldots, N, \quad \text { and } \quad \varphi(\mathbf{s})=\sigma_{1}, \quad \forall \mathbf{s} \in \mathcal{U}^{(1)} \tag{1.5}
\end{equation*}
$$

which corresponds to the situation where $\varphi$ admits $N$ wells of the same height. In order to avoid heavy notation, we also assume that for $n=1, \ldots, N$ the minimum of $\varphi$ on $E_{n}$ is attained in a single point that we denote by $\mathbf{m}_{n}$.

The associated Arrhenius number, $S=\sigma_{1}-\varphi_{0}$, governs the long time dynamics of (1.1). That is made quantitative in Theorems 1 below. More general assumptions can be made as will be clear from the proofs. We restrict ourselves to the case in which the asymptotics are cleanest.

To state the simplest result let us assume that $d=1$ and that the second derivative of $\varphi$ is constant on the sets $\mathcal{U}^{(0)}$ and $\mathcal{U}^{(1)}$ :

$$
\begin{equation*}
\varphi^{\prime \prime}(\mathbf{m})=\mu, \quad \forall \mathbf{m} \in \mathcal{U}^{(0)} \quad \text { and } \quad \varphi^{\prime \prime}(\mathbf{s})=-\nu, \quad \forall \mathbf{s} \in \mathcal{U}^{(1)} \tag{1.6}
\end{equation*}
$$

for some $\mu, \nu>0$. The potential then looks like the one shown in Fig.1. We introduce the matrix

$$
A_{0}=\frac{\kappa}{\pi}\left(\begin{array}{rrrrrrrr}
1 & -1 & 0 & 0 & \ldots & \ldots & \ldots & 0  \tag{1.7}\\
-1 & 2 & -1 & 0 & \ldots & \ldots & \ldots & 0 \\
0 & -1 & 2 & -1 & 0 & \ldots & \ldots & 0 \\
\vdots & 0 & -1 & 2 & \ddots & \ddots & \ldots & 0 \\
\vdots & \vdots & \ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
\vdots & \vdots & \ddots & \ddots & \ddots & \ddots & -1 & 0 \\
0 & \vdots & \ddots & \ddots & \ddots & -1 & 2 & -1 \\
0 & 0 & \ldots & \ldots & \ldots & 0 & -1 & 1
\end{array}\right)
$$

with $\kappa=\sqrt{\mu \nu}$. This matrix is positive semi-definite with a simple eigenvalue at 0 .
Theorem 1. Suppose that $d=1$ and $\varphi$ satisfies (1.2), (1.5) and (1.6). Suppose that

$$
\begin{equation*}
\rho_{0}=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{2}}\left(\sum_{n=1}^{N} \beta_{n} \mathbb{1}_{E_{n}}+r_{\epsilon}\right) e^{-\varphi / \epsilon^{2}}, \quad \lim _{\epsilon \rightarrow 0}\left\|r_{\epsilon}\right\|_{L^{\infty}}=0, \quad \beta \in \mathbb{R}^{N} \tag{1.8}
\end{equation*}
$$

then the solution to (1.1) satisfies, uniformly for $\tau \geq 0$,

$$
\begin{equation*}
\rho\left(2 \epsilon^{2} e^{S / \epsilon^{2}} \tau, x\right) \rightarrow \sum_{n=1}^{N} \alpha_{n}(\tau) \delta_{\mathbf{m}_{n}}(x), \quad \epsilon \rightarrow 0 \tag{1.9}
\end{equation*}
$$

in the sense of distributions in $x$, where $S=\sigma_{1}-\varphi_{0}$ and where $\alpha(\tau)=\left(\alpha_{1}, \ldots, \alpha_{n}\right)(\tau)$ solves

$$
\begin{equation*}
\partial_{\tau} \alpha=-A_{0} \alpha, \quad \alpha(0)=\beta \tag{1.10}
\end{equation*}
$$

with $A_{0}$ given by (1.7).
The above result is a generalization of Theorem 2.5 in [5] where the case of a doublewell is considered and estimates are uniform on compact time intervals only. We remark that the equation considered in [5] has also an additional transverse variable (varying slowly). A development of the methods presented in this note would also allow having such variables. Since our goal is to explain general ideas in a simple setting we do not address this issue here.

A higher dimensional version of Theorem 1 is given in Theorem 3 in §3. In this higher dimensional setting, the matrix $A_{0}$ becomes a graph Laplacian for a graph obtained by taking minima as vertices and saddle points as edges. The same graph Laplacian was used by Landim et al [14] in the context of a discrete model of the Kramers-Smoluchowski equation.

Using methods of [5] and [2], Theorem 3 was also proved by Seo-Tabrizian [12], but as the other previous papers, without uniformity in time (that is, with convergence uniform for $t \in[0, T])$.

Here, Theorem 1 is a consequence of a more precise asymptotic formula given in Theorem 2 formulated using the Witten Laplacian. Provided that certain topological
assumptions are satisfied (see $[15, \S 1.1, \S 1.2]$ ) an analogue of Theorem 1 in higher dimensions is immediate - see $\S 3$ for geometrically interesting examples.

The need for the new results of [15] comes from the fact that in the papers on the lowlying eigenvalues of the Witten Laplacian $[2,7,8,9,10]$ the authors make assumptions on the relative positions of minima and of saddle points. These assumptions mean that the Arrhenius numbers are distinct and hence potentials for which the KramersSmoluchowski dynamics (1.9) is interesting are excluded. With this motivation the general case was studied in [15] and to explain how the results of that paper can be used in higher dimensions we give a self-contained presentation in dimension one.

We remark that we need specially prepared initial data (1.8) to obtain results valid for all times. Also, $E_{n}$ 's in the statement can be replaced by any interval in $E_{n}$ containing the minimum. Theorem 2 also shows that a weaker result is valid for any $L^{2}$ data: suppose that $\rho_{0} \in L_{\varphi}^{2}:=L^{2}\left(e^{\varphi(x) / \epsilon^{2}} d x\right)$ and that

$$
\beta_{n}:=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \int_{E_{n}} \rho_{0}(x) d x .
$$

Then, uniformly for $\tau \geq 0$,

$$
\begin{align*}
& \rho(t, x)=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \sum_{n=1}^{N} \alpha_{n}\left(\left(2 \epsilon^{2}\right)^{-1} e^{-S / \epsilon^{2}} t\right) \mathbb{1}_{E_{n}}(x) e^{-\varphi / \epsilon^{2}}+r_{\epsilon}(t, x),  \tag{1.11}\\
& \left\|r_{\epsilon}(t)\right\|_{L^{1}(d x)} \leq C\left(\epsilon^{\frac{5}{2}}+\epsilon^{\frac{1}{2}} e^{-t \epsilon^{2}}\right)\left\|\rho_{0}\right\|_{L_{\varphi}^{2}}, \quad L_{\varphi}^{2}:=L^{2}\left(\mathbb{R}, e^{\varphi(x) / \epsilon^{2}} d x\right) .
\end{align*}
$$

where $\alpha$ solves (1.10). The proof of (1.11) is given at the end of $\S 2.6$.
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## 2. Dimension one

In this section we assume that the dimension is equal to $d=1$. That allows to present self-contained proofs which indicate the strategy for higher dimension.

Ordering the sets $E_{n}$ such that $\mathbf{m}_{1}<\mathbf{m}_{2}<\ldots<\mathbf{m}_{N}$ it follows that for all $n=$ $1, \ldots N-1 \bar{E}_{n} \cap \bar{E}_{n+1}=\left\{\mathbf{s}_{n}\right\}$ is a maximum and we assume additionally that there exists $\mu_{n}, \nu_{k}>0$ such that for $n=1, \ldots, N$ and $k=1, \ldots, N-1$,

$$
\begin{equation*}
\varphi^{\prime \prime}\left(\mathbf{m}_{n}\right)=\mu_{n} \quad \text { and } \quad \varphi^{\prime \prime}\left(\mathbf{s}_{k}\right)=-\nu_{k} . \tag{2.1}
\end{equation*}
$$

Using this notation we define a symmetric $N \times N$ matrix: $A_{0}=\left(a_{i j}\right)_{1 \leq i, j \leq N}$, where (with the convention that $\nu_{0}=\nu_{N}=0$ )

$$
\begin{equation*}
a_{i i}=\pi^{-1} \mu_{j}^{\frac{1}{2}}\left(\nu_{j-1}^{\frac{1}{2}}+\nu_{j}^{\frac{1}{2}}\right), \quad a_{i, i+1}=-\pi^{-1} \nu_{i}^{\frac{1}{2}} \mu_{i}^{\frac{1}{4}} \mu_{i+1}^{\frac{1}{4}}, \quad 1 \leq i \leq N-1, \tag{2.2}
\end{equation*}
$$

and $a_{i, i+k}=0$, for $k>1, a_{i j}=a_{j i}$. The matrix $A_{0}$ is symmetric positive and the eigenvalue 0 has multiplicity 1 . When $\mu_{j}$ 's and $\nu_{j}$ 's are all equal our matrix takes the particularly simple form (1.7).

First, observe that we can assume without loss of generality that $\varphi_{0}=0$. Define the operator appearing on the right hand side of (1.1) by

$$
P:=\partial_{x} \cdot\left(\partial_{x}+\epsilon^{-2} \partial_{x} \varphi\right)
$$

and denote

$$
h=2 \epsilon^{2} .
$$

Then, considering $e^{ \pm \varphi / h}$ as a multiplication operator,

$$
P=\partial_{x} \circ\left(\partial_{x}+2 h^{-1} \partial_{x} \varphi\right)=\partial_{x} \circ e^{-2 \varphi / h} \circ \partial_{x} \circ e^{2 \varphi / h}
$$

and

$$
e^{\varphi / h} \circ P \circ e^{-\varphi / h}=-h^{-2} \Delta_{\varphi}, \quad \Delta_{\varphi}:=-h^{2} \Delta+\left|\partial_{x} \varphi\right|^{2}-h \Delta \varphi .
$$

Hence, $\rho$ is solution of (1.1) if $u(t, x):=e^{\varphi(x) / h} \rho\left(h^{2} t, x\right)$ is a solution of

$$
\begin{equation*}
\partial_{t} u=-\Delta_{\varphi} u, \quad u_{\mid t=0}=u_{0}:=\rho_{0} e^{\varphi / h} . \tag{2.3}
\end{equation*}
$$

In order to state our result for this equation, we denote

$$
\begin{equation*}
\psi_{n}(x):=c_{n}(h) h^{-\frac{1}{4}} \mathbb{1}_{E_{n}}(x) e^{-\left(\varphi-\varphi_{0}\right)(x) / h}, \quad \forall n=1, \ldots, N, \tag{2.4}
\end{equation*}
$$

where $c_{n}(h)$ is a normalization constant such that $\left\|\psi_{n}\right\|_{L^{2}}=1$. The method of steepest descent shows that

$$
\begin{equation*}
c_{n}(h) \sim \sum_{k=0}^{\infty} c_{n, k} h^{k}, \quad c_{n, 0}=\left(\mu_{n} / \pi\right)^{\frac{1}{4}}, \quad \forall n=1, \ldots N . \tag{2.5}
\end{equation*}
$$

We then define a map $\Psi: \mathbb{R}^{N} \rightarrow L^{2}$ by

$$
\begin{equation*}
\Psi(\beta):=\sum_{n=1}^{N} \beta_{n} \psi_{n}, \quad \forall \beta=\left(\beta_{1}, \ldots, \beta_{N}\right) \in \mathbb{R}^{N} . \tag{2.6}
\end{equation*}
$$

The following theorem describes the dynamic of the above equation when $h \rightarrow 0$.
Theorem 2. There exists $C>0$ and $h_{0}>0$ such that for all $\beta \in \mathbb{R}^{N}$ and all $0<h<h_{0}$, we have

$$
\begin{equation*}
\left\|e^{-t \Delta_{\varphi}} \Psi(\beta)-\Psi\left(e^{-t \nu_{h} A} \beta\right)\right\|_{L^{2}} \leq C e^{-\frac{1}{C h}}|\beta|, \quad \forall t \geq 0 \tag{2.7}
\end{equation*}
$$

where $\nu_{h}=h e^{-2 S / h}, S=\sigma_{1}-\varphi_{0}$, and $A=A(h)$ is a real symmetric positive matrix having a classical expansion $A \sim \sum_{k=0}^{\infty} h^{k} A_{k}$ with $A_{0}$ given by (2.2). In addition,

$$
\begin{equation*}
\left\|e^{-t \Delta_{\varphi}} \Psi(\beta)-\Psi\left(e^{-t \nu_{h} A_{0}} \beta\right)\right\|_{L^{2}} \leq C h|\beta| \tag{2.8}
\end{equation*}
$$

uniformly with respect to $t \geq 0$.
We first show how

Theorem 2 implies Theorem 1. First recall that we assume here $\mu_{n}=\mu$ for all $n=$ $1, \ldots N$ and $\nu_{k}=\nu$ for all $k=1, \ldots N-1$. Suppose that $\rho$ is the solution to (1.1) with $\rho_{0}$ as in Theorem 1. Then $u(t, x):=e^{\varphi(x) / h} \rho\left(h^{2} t, x\right)$ is a solution of (2.3), that is $u(t)=e^{-t \Delta_{\varphi}} u_{0}$ with

$$
\begin{align*}
u_{0} & =\rho_{0} e^{\varphi / 2 \epsilon^{2}}=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{2}}\left(\sum_{n=1}^{N} \beta_{n} \mathbb{1}_{E_{n}}+r_{\epsilon}\right) e^{-\varphi / 2 \epsilon^{2}} \\
& =\left(\frac{\mu}{\pi h}\right)^{\frac{1}{2}}\left(\sum_{n=1}^{N} \beta_{n} \mathbb{1}_{E_{n}}+r_{h}\right) e^{-\varphi / h} \tag{2.9}
\end{align*}
$$

Since, $c_{n}(h)=(\mu / \pi)^{\frac{1}{4}}+\mathcal{O}(h)$ it follows that

$$
u_{0}=(\mu / \pi h)^{\frac{1}{4}} \Psi(\beta)+\tilde{r}_{h}, \quad \tilde{r}_{h}=\left(\mathcal{O}\left(h^{\frac{1}{2}}\right)+h^{-\frac{1}{2}} r_{h}\right) e^{-\varphi / h}
$$

Since $h^{-\frac{1}{2}} e^{-\varphi / h}=\mathcal{O}_{L^{1}}(1)$, we have $\tilde{r}_{h} \rightarrow 0$ in $L^{1}$ when $h \rightarrow 0$. Hence, it follows from (2.8) (Theorem 2) that

$$
\begin{aligned}
\rho\left(h^{2} t, x\right) & =e^{-\varphi(x) / h} u(t, x)=e^{-\varphi(x) / h} e^{-t \Delta_{\varphi}}\left((\mu / \pi h)^{\frac{1}{4}} \Psi(\beta)+\tilde{r}_{h}\right) \\
& =e^{-\varphi(x) / h}\left((\mu / \pi h)^{\frac{1}{4}} \Psi\left(e^{-t \nu_{h} A_{0}} \beta\right)+e^{-t \Delta_{\varphi}} \tilde{r}_{h}+\mathcal{O}_{L^{2}}(h)\right)
\end{aligned}
$$

With the new time variable $s=t \nu_{h}$, we obtain

$$
\begin{equation*}
\rho\left(s h e^{2 S / h}, x\right)=e^{-\varphi(x) / h}\left((\mu / \pi h)^{\frac{1}{4}} \Psi\left(e^{-s A_{0}} \beta\right)+e^{-t \Delta_{\varphi}} \tilde{r}_{h}+\mathcal{O}_{L^{2}}(h)\right) \tag{2.10}
\end{equation*}
$$

and denoting $\alpha(s)=e^{-s A_{0}} \beta$, we get

$$
e^{-\varphi(x) / h}(\mu / \pi h)^{\frac{1}{4}} \Psi\left(e^{-s A_{0}} \beta\right)=(\mu / \pi)^{\frac{1}{4}} \sum_{n=1}^{N} \alpha_{n}(t) h^{-\frac{1}{2}} c_{n}(h) \chi_{n}(x) e^{-2 \varphi(x) / h} .
$$

On the other hand, $h^{-\frac{1}{2}} \chi_{n}(x) e^{-2 \varphi(x) / h} \longrightarrow(\pi / \mu)^{\frac{1}{2}} \delta_{x=m_{n}}$, as $h \rightarrow 0$, in the sense of distributions. Since, $c_{n}(h)=(\mu / \pi)^{\frac{1}{4}}+\mathcal{O}(h)$, it follows that

$$
\begin{equation*}
e^{-\varphi / h}(\mu / \pi h)^{\frac{1}{4}} \Psi\left(e^{-s A_{0}} \beta\right) \longrightarrow \sum_{n=1}^{N} \alpha_{n}(t) \delta_{x=m_{n}} \tag{2.11}
\end{equation*}
$$

when $h \rightarrow 0$. Moreover, since $e^{-t \Delta_{\varphi}}$ is bounded by 1 on $L^{2}$, then

$$
\left\|h^{-\frac{1}{2}} e^{-\varphi / h} e^{-t \Delta_{\varphi}}\left(r_{h} e^{-\varphi / h}\right)\right\|_{L^{1}} \leq\left\|r_{h}\right\|_{L^{\infty}}\left\|h^{-\frac{1}{4}} e^{-\varphi / h}\right\|_{L^{2}}^{2} \leq C\left\|r_{h}\right\|_{L^{\infty}}
$$

and recalling that $r_{h} \rightarrow 0$ in $L^{\infty}$, we see that

$$
\begin{equation*}
e^{-\varphi(x) / h}\left(e^{-t \Delta_{\varphi}} \tilde{r}_{h}+\mathcal{O}_{L^{2}}(h)\right) \longrightarrow 0 \tag{2.12}
\end{equation*}
$$

in the sense of distributions. Inserting (2.11) and (2.12) into (2.10) and recalling that $h=2 \epsilon^{2}$, we obtain (1.9).
2.1. Witten Laplacian in dimension one. The Witten Laplacian is particularly simple in dimension one but one can already observe features which play a crucial role in general study. For more information we refer to $[4, \S 11.1]$ and [7].

We first consider $\Delta_{\varphi}$ acting on $C_{c}^{\infty}(\mathbb{R})$ and recall a supersymmetric structure which is the starting point of our analysis:

$$
\begin{equation*}
\Delta_{\varphi}=d_{\varphi}^{*} \circ d_{\varphi} \tag{2.13}
\end{equation*}
$$

with $d_{\varphi}=e^{-\varphi / h} \circ h \partial_{x} \circ e^{\varphi / h}=h \partial_{x}+\partial_{x} \varphi$ and $d_{\varphi}^{*}=-h \partial_{x}+\partial_{x} \varphi=-d_{-\varphi}$. From this square structure, it is clear that $\Delta_{\varphi}$ is non negative and that we can use the Friedrichs extension to define a self-adjoint operator $\Delta_{\varphi}$ with domain denoted $D\left(\Delta_{\varphi}\right)$. Moreover, it follows from (1.2) that there exists $c_{0}, h_{0}>0$ such that for $0<h<h_{0}$,

$$
\begin{equation*}
\sigma_{\mathrm{ess}}\left(\Delta_{\varphi}\right) \subset\left[c_{0},+\infty\right) \tag{2.14}
\end{equation*}
$$

Therefore, $\sigma\left(\Delta_{\varphi}\right) \cap\left[0, c_{0}\right)$ consists of eigenvalues of finite multiplicity and with no accumulation points except possibly $c_{0}$.

The following proposition gives a preliminary description of the low-lying eigenvalues.

Proposition 1. There exist $\varepsilon_{0}, h_{0}>0$ such that for any $h \in\left(0, h_{0}\right], \Delta_{ \pm \varphi}$ has exactly $N$ eigenvalues $0 \leq \lambda_{1}^{ \pm} \leq \lambda_{2}^{ \pm} \ldots \leq \lambda_{N}^{ \pm}$in the interval $\left[0, \epsilon_{0} h\right]$. Moreover, for any $\epsilon>0$ there exists $C$ such that

$$
\begin{equation*}
\lambda_{n}^{ \pm}(h) \leq C e^{-(S-\epsilon) / h} \tag{2.15}
\end{equation*}
$$

where $S=\sigma_{1}-\varphi_{0}$.
Remark. The proof applies to any $\varphi$ which satisfies the first two inequalities in (1.2). If one assumes additionally that $\varphi(x) \geq C|x|$ for $|x|$ large, then $e^{-\varphi / h} \in D\left(\Delta_{\varphi}\right)$. Since $d_{\varphi}\left(e^{-\varphi / h}\right)=0$ it follows that $\lambda_{0}^{+}=0$.

Proof. This is proved in [4, Theorem 11.1] with $h^{\frac{3}{2}}$ in place of $\epsilon_{0} h$. The proof applies in any dimension and we present it in that greater generality for $\varphi$ satisfying

$$
|\partial \varphi(x)| \geq \frac{1}{C}, \quad\left|\partial_{x_{i} x_{j}}^{2} \varphi\right| \leq C|\partial \varphi|^{2} .
$$

The fact that there exists at least $N$ eigenvalues in the interval $\left[0, C e^{-(S-\epsilon) / h}\right]$ is a direct consequence of the existence of $N$ linearly independent quasi-modes - see Lemma 2 and (2.31) below.

To show that $N$ is the exact number of eigenvalues in $\left[0, \epsilon_{0} h\right)$ it suffices to find a $N$ dimensional vector space $V$ and $\varepsilon_{0}>0$ such that the operator $\Delta_{\varphi}$ is bounded from below by $\varepsilon_{0} h$ on $V^{\perp}$ - see for instance [18, Theorem C.15].

To find $V$ we introduce a family of harmonic oscillators associated to minima $\mathbf{m} \in$ $\mathcal{U}^{(0)}$ and obtained by replacing $\varphi$ by its harmonic approximation in the expression for $\Delta_{\varphi}$ :

$$
H_{\mathbf{m}}:=-h^{2} \Delta+\left|\varphi^{\prime \prime}(\mathbf{m})(x-\mathbf{m})\right|^{2}-h \Delta \varphi(\mathbf{m}), \quad \mathbf{m} \in \mathcal{U}^{(0)} .
$$

The spectrum of this operator is known explicitly, see [7, Sect 2.1] with the simple eigenvalue 0 at the bottom. We denote by $e_{\mathrm{m}}$ the normalized eigenfunction, $H_{\mathrm{m}} e_{\mathrm{m}}=0$. The other eigenvalues of $H_{\mathbf{m}}$ are bounded from below by $c_{0} h$ for some $c_{0}>0$.

Let $\chi \in C_{c}^{\infty}\left(\mathbb{R}^{d} ;[0,1]\right)$ be equal to 1 near 0 and satisfy $\left(1-\chi^{2}\right)^{\frac{1}{2}} \in C^{\infty}\left(\mathbb{R}^{d}\right)$. We define $\chi_{\mathbf{m}}(x)=\chi((x-\mathbf{m}) / \sqrt{M h})$ where $M>0$ will be chosen later. For $h$ small enough, the functions $\chi_{\mathrm{m}}$ have disjoint supports and hence the function $\chi_{\infty}$ defined by $1-\chi_{\infty}^{2}=\sum_{\mathbf{m} \in \mathcal{U}^{(0)}} \chi_{\mathbf{m}}^{2}$ is smooth. We define the $N$-dimensional vector space

$$
V=\operatorname{span}\left\{\chi_{\mathbf{m}} e_{\mathbf{m}}, \mathbf{m} \in \mathcal{U}^{(0)}\right\}
$$

The proof is completed if we show that there exist $\varepsilon_{0}, h_{0}>0$ such that

$$
\begin{equation*}
\left.\left.\left\langle\Delta_{\varphi} u, u\right\rangle \geq \varepsilon_{0} h\|u\|^{2}, \quad \forall u \in V^{\perp} \cap D\left(\Delta_{\varphi}\right), \quad \forall h \in\right] 0, h_{0}\right] . \tag{2.16}
\end{equation*}
$$

To establish (2.16) we use the following localization formula the verification of which is left to the reader (see [4, Theorem 3.2]):

$$
\Delta_{\varphi}=\sum_{\mathbf{m} \in \mathcal{U}^{(0)} \cup\{\infty\}} \chi_{\mathbf{m}} \circ \Delta_{\varphi} \circ \chi_{\mathbf{m}}-h^{2} \sum_{\mathbf{m} \in \mathcal{U}^{(0)} \cup\{\infty\}}\left|\nabla \chi_{\mathbf{m}}\right|^{2} .
$$

Since, $\nabla \chi_{\mathbf{m}}=\mathcal{O}\left((M h)^{-\frac{1}{2}}\right)$, this implies, for $u \in D\left(\Delta_{\varphi}\right)$, that

$$
\begin{equation*}
\left\langle\Delta_{\varphi} u, u\right\rangle=\left\langle\Delta_{\varphi} \chi_{\infty} u, \chi_{\infty} u\right\rangle+\sum_{\mathbf{m} \in \mathcal{U}^{(0)}}\left\langle\Delta_{\varphi} \chi_{\mathbf{m}} u, \chi_{\mathbf{m}} u\right\rangle+\mathcal{O}\left(h M^{-1}\|u\|^{2}\right) . \tag{2.17}
\end{equation*}
$$

On the support of $\chi_{\infty}$ we have $|\nabla \varphi|^{2}-h \Delta \varphi \geq(1-\mathcal{O}(h))|\nabla \varphi|^{2} \geq c_{1} M h$ for some $c_{1}>0$, and hence

$$
\begin{equation*}
\left\langle\Delta_{\varphi} \chi_{\infty} u, \chi_{\infty} u\right\rangle \geq M c_{1} h\left\|\chi_{\infty} u\right\|^{2} \tag{2.18}
\end{equation*}
$$

On the other hand, near any $\mathbf{m} \in \mathcal{U}^{(0)},|\nabla \varphi(x)|^{2}=\left|\varphi^{\prime \prime}(\mathbf{m})(x-\mathbf{m})\right|^{2}+\mathcal{O}\left(|x-\mathbf{m}|^{3}\right)$ and $\varphi^{\prime \prime}(x)=\varphi^{\prime \prime}(\mathbf{m})+\mathcal{O}(|x-\mathbf{m}|)$. Since on the support of $\chi_{\mathbf{m}}$ we have $|x-\mathbf{m}| \leq \sqrt{M h}$, it follows that

$$
\begin{equation*}
\left\langle\Delta_{\varphi} \chi_{\mathbf{m}} u, \chi_{\mathbf{m}} u\right\rangle=\left\langle H_{\mathbf{m}} \chi_{\mathbf{m}} u, \chi_{\mathbf{m}} u\right\rangle+\mathcal{O}\left((M h)^{\frac{3}{2}}\right) \tag{2.19}
\end{equation*}
$$

We now assume that $u \in D\left(\Delta_{\varphi}\right)$ is orthogonal to $\chi_{\mathbf{m}} e_{\mathbf{m}}$ for all $\mathbf{m}$. Then $\chi_{\mathbf{m}} u$ is orthogonal to $e_{\mathbf{m}}$. Since the spectral gap of $H_{\mathbf{m}}$ is bounded from below by $c_{0} h$, (2.19) shows that

$$
\begin{equation*}
\left\langle\Delta_{\varphi} \chi_{\mathbf{m}} u, \chi_{\mathbf{m}} u\right\rangle \geq c_{0} h\left\|\chi_{\mathbf{m}} u\right\|^{2}+\mathcal{O}\left((M h)^{\frac{3}{2}}\|u\|^{2}\right), \quad \forall \mathbf{m} \in \mathcal{U}^{(0)} \tag{2.20}
\end{equation*}
$$

Combining this with (2.17), (2.18) and (2.20) gives

$$
\begin{aligned}
\left\langle\Delta_{\varphi} u, u\right\rangle & \geq c_{0} h \sum_{\mathbf{m} \in \mathcal{U}^{(0)} \cup\{\infty\}}\left\|\chi_{\mathbf{m}} u\right\|^{2}+\mathcal{O}\left(h M^{-1}\|u\|^{2}\right)+\mathcal{O}\left((M h)^{\frac{3}{2}}\|u\|^{2}\right) \\
& \geq c_{0} h\|u\|^{2}+\mathcal{O}\left(h M^{-1}\|u\|^{2}\right)+\mathcal{O}\left((M h)^{\frac{3}{2}}\|u\|^{2}\right)
\end{aligned}
$$

Taking $M$ large enough completes the proof of (2.16).
We denote by $E^{(0)}$ the subspace spanned by eigenfunctions of these low lying eigenvalues and by

$$
\begin{equation*}
\Pi^{(0)}:=\mathbb{1}_{\left[0, \varepsilon_{0} h\right]}\left(\Delta_{\varphi}\right) \tag{2.21}
\end{equation*}
$$

the spectral projection onto $E^{(0)}$. This projector is expressed by the standard contour integral

$$
\begin{equation*}
\Pi^{(0)}=\frac{1}{2 \pi i} \int_{\partial B\left(0, \delta \varepsilon_{0} h\right)}\left(z-\Delta_{\varphi}\right)^{-1} d z \tag{2.22}
\end{equation*}
$$

In our analysis, we will also need the operator $\Delta_{-\varphi}$, noting that in dimension one $\Delta_{-\varphi}$ is the Witten Laplacian on 1-forms. Since $-\varphi$ has exactly $N-1$ minima (given by the $N-1$ maxima of $\varphi$ ), it follows from Proposition 1 that there exists $\varepsilon_{1}>0$ such that $\Delta_{-\varphi}$ has $N-1$ eigenvalues in $\left[0, \varepsilon_{1} h\right]$ and that these eigenvalues are actually exponentially small. Observe that because of the condition $\varphi(x) \geq C|x|$ at infinity, the function $e^{\varphi / h}$ is not square integrable. Consequently, unlike in the case of $\Delta_{\varphi}$, we cannot conclude that the lowest eigenvalue is equal to 0

We denote by $E^{(1)}$ the subspace spanned by eigenfunctions of these low-lying eigenfunctions of $\Delta_{-\varphi}$ and by $\Pi^{(1)}$ the corresponding projector onto $E^{(1)}$,

$$
\begin{equation*}
\Pi^{(1)}=\mathbb{1}_{\left[0, \varepsilon_{1} h\right]}\left(\Delta_{-\varphi}\right) . \tag{2.23}
\end{equation*}
$$

Similarly to (2.22), we have

$$
\begin{equation*}
\Pi^{(1)}=\frac{1}{2 \pi i} \int_{\partial B\left(0, \delta \varepsilon_{1} h\right)}\left(z-\Delta_{-\varphi}\right)^{-1} d z \tag{2.24}
\end{equation*}
$$

for any $0<\delta<1$.
2.2. Supersymmetry. The key point in the analysis is the following intertwining relations which follows directly from (2.13)

$$
\begin{equation*}
\Delta_{-\varphi} \circ d_{\varphi}=d_{\varphi} \circ \Delta_{\varphi} \tag{2.25}
\end{equation*}
$$

and its adjoint relation

$$
\begin{equation*}
d_{\varphi}^{*} \circ \Delta_{-\varphi}=\Delta_{\varphi} \circ d_{\varphi}^{*} . \tag{2.26}
\end{equation*}
$$

From these relations we deduce that $d_{\varphi}\left(E^{(0)}\right) \subset E^{(1)}$ and $d_{\varphi}^{*}\left(E^{(1)}\right) \subset E^{(0)}$. Indeed, suppose that $\Delta_{\varphi} u=\lambda u$, with $u \neq 0$ and $\lambda \in\left[0, \varepsilon_{0} h\right]$. Then, we see from (2.25) that

$$
\Delta_{-\varphi}\left(d_{\varphi} u\right)=d_{\varphi}\left(\Delta_{\varphi} u\right)=\lambda d_{\varphi} u
$$

Therefore, either $d_{\varphi} u$ is null and obviously belongs to $E^{(1)}$ or $d_{\varphi} u \neq 0$ and hence $d_{\varphi} u$ is an eigenvector of $\Delta_{-\varphi}$ associated to $\lambda \in\left[0, \varepsilon_{0} h\right]$. This proves the first statement. The inclusion $d_{\varphi}^{*}\left(E^{(1)}\right) \subset E^{(0)}$ is obtained by similar arguments.

By definition, the operator $\Delta_{\varphi}$ maps $E^{(0)}$ into itself and we can consider its restriction to $E^{(0)}$. From the above discussion we know also that $d_{\varphi}\left(E^{(0)}\right) \subset E^{(1)}$ and $d_{\varphi}^{*}\left(E^{(1)}\right) \subset$ $E^{(0)}$. Hence we consider $\mathcal{L}=\left(d_{\varphi}\right)_{\mid E^{(0)} \rightarrow E^{(1)}}$ and $\mathcal{L}^{*}=\left(d_{\varphi}^{*}\right)_{\mid E^{(1)} \rightarrow E^{(0)}}$. When restricted to $E^{(0)}$, the structure equation (2.13) becomes

$$
\begin{equation*}
\mathcal{M}=\mathcal{L}^{*} \mathcal{L} \quad \text { with } \quad \mathcal{M}:=\left.\Delta_{\varphi}\right|_{E^{(0)}}, \quad \mathcal{L}:=\left.\left(d_{\varphi}\right)\right|_{E^{(0)} \rightarrow E^{(1)}} \tag{2.27}
\end{equation*}
$$

2.3. Quasi-modes for $\Delta_{\varphi}$. Let $\delta_{0}=\inf \left\{\operatorname{diam}\left(E_{n}\right), n=1, \ldots, N\right\}$ and let $\epsilon>0$ be small with respect to $\delta_{0}$. For all $n=1, \ldots, N$, let $\chi_{n}$ be smooth cut-off functions such that

$$
\left\{\begin{array}{c}
0 \leq \chi_{n} \leq 1  \tag{2.28}\\
\operatorname{supp}\left(\chi_{n}\right) \subset\left\{x \in E_{n}, \varphi(x) \leq \sigma_{1}-\epsilon\right\} \\
\chi_{n}=1 \text { on }\left\{x \in E_{n}, \varphi(x) \leq \sigma_{1}-2 \epsilon\right\}
\end{array}\right.
$$

where $\epsilon>0$ will be chosen small (in particular much smaller than $\delta_{0}$ in (2.32)). Consider now the family of approximated eigenfunctions defined by

$$
\begin{equation*}
f_{n}^{(0)}(x)=h^{-\frac{1}{4}} c_{n}(h) \chi_{n}(x) e^{-\varphi(x) / h}, \quad\left\|f_{n}^{(0)}\right\|_{L^{2}}=1 \tag{2.29}
\end{equation*}
$$

where $c_{n}(h)=\varphi^{\prime \prime}\left(m_{n}\right)^{\frac{1}{4}} \pi^{-\frac{1}{4}}+\mathcal{O}(h)$. We introduce the projection of these quasi-modes onto the eigenspace space $E^{(0)}$ :

$$
\begin{equation*}
g_{n}^{(0)}:=\Pi^{(0)} f_{n}^{(0)} \tag{2.30}
\end{equation*}
$$

Lemma 2. The approximate eigenfunctions defined by (2.29) satisfy

$$
\left\langle f_{n}^{(0)}, f_{m}^{(0)}\right\rangle=\delta_{n, m}, \quad \forall n, m=1, \ldots, N
$$

and

$$
d_{\varphi} f_{n}^{(0)}=\mathcal{O}_{L^{2}}\left(e^{-(S-\epsilon) / h}\right), \quad g_{n}^{(0)}-f_{n}^{(0)}=\mathcal{O}_{L^{2}}\left(e^{-\left(S-\epsilon^{\prime}\right) / h}\right)
$$

for any $\epsilon^{\prime}>\epsilon$.
Proof. The first statement is a direct consequence of the support properties of the cut-off functions $\chi_{n}$ and the choice of the normalizing constant. To see the second estimate, recall that $d_{\varphi} e^{-\varphi / h}=0$. Hence

$$
d_{\varphi} f_{n}^{(0)}(x)=h^{\frac{3}{4}} c_{n}(h) \chi_{n}^{\prime}(x) e^{-\varphi(x) / h}
$$

Moreover, thanks to (2.28), there exists $c>0$ such that for $\epsilon>0$ small enough we have $\varphi(x) \geq S-\epsilon$ for $x \in \operatorname{supp}\left(\chi_{n}^{\prime}\right)$. Combining these two facts gives estimates on $d_{\varphi} f_{n}^{(0)}$.

We now prove the estimate on $g_{n}^{(0)}-f_{n}^{(0)}$. We first observe that

$$
\Delta_{\varphi} f_{n}^{(0)}=d_{\varphi}^{*} d_{\varphi} f_{n}^{(0)}=h^{\frac{3}{4}} c_{n}(h) d_{\varphi}^{*}\left(\chi_{n}^{\prime} e^{-\varphi / h}\right)=h^{\frac{3}{4}} c_{n}(h)\left(-h \chi_{n}^{\prime \prime}+2 \partial_{x} \varphi \chi_{n}^{\prime}\right) e^{-\varphi / h}
$$

and the same argument as before shows that

$$
\begin{equation*}
\Delta_{\varphi} f_{n}^{(0)}=\mathcal{O}_{L^{2}}\left(e^{-(S-\epsilon) / h}\right) \tag{2.31}
\end{equation*}
$$

From (2.22) and Cauchy formula, it follows that

$$
\begin{aligned}
g_{n}^{(0)}-f_{n}^{(0)} & =\Pi^{(0)} f_{n}^{(0)}-f_{n}^{(0)}=\frac{1}{2 \pi i} \int_{\gamma}\left(z-\Delta_{\varphi}\right)^{-1} f_{n}^{(0)} d z-\frac{1}{2 \pi i} \int_{\gamma} z^{-1} f_{n}^{(0)} d z \\
& =\frac{1}{2 \pi i} \int_{\gamma}\left(z-\Delta_{\varphi}\right)^{-1} z^{-1} \Delta_{\varphi} f_{n}^{(0)} d z
\end{aligned}
$$

with $\gamma=\partial B\left(0, \delta \epsilon_{0} h\right), 0<\delta<1$. Since $\Delta_{\varphi}$ is selfadjoint and $\sigma\left(\Delta_{\varphi}\right) \cap\left[0, \epsilon_{0} h\right] \subset$ [ $\left.0, e^{-1 / C h}\right]$, we have for $\alpha>0$ small enough

$$
\left\|\left(z-\Delta_{\varphi}\right)^{-1}\right\|=\mathcal{O}\left(h^{-1}\right)
$$

uniformly for $z \in \gamma$. Using (2.31), we get $\left\|\left(z-\Delta_{\varphi}\right)^{-1} z^{-1} \Delta_{\varphi} f_{k}^{(0)}\right\|=\mathcal{O}\left(h^{-2} e^{-(S-\epsilon) / h}\right)$, and, after integration, $\left\|g_{k}^{(0)}-f_{k}^{(0)}\right\|=\mathcal{O}\left(h^{-1} e^{-(S-\epsilon) / h}\right)=\mathcal{O}\left(e^{-\left(S-\epsilon^{\prime}\right) / h}\right)$, for any $\epsilon^{\prime}>$ $\epsilon$.
2.4. Quasi-modes for $\Delta_{-\varphi}$. Since, $\varphi$ and $-\varphi$ share similar properties, the construction of the preceding section produces quasi-modes for $\Delta_{-\varphi}$. Eventually we will only need quasi-modes localized near the maxima $\mathbf{s}_{k}$. Hence, let $\theta_{k} \in C_{c}^{\infty}(\mathbb{R} ;[0,1])$ satisfy

$$
\begin{equation*}
\operatorname{supp} \theta_{k} \subset\left\{\left|x-\mathbf{s}_{k}\right| \leq \delta_{0}\right\}, \quad \theta_{k}=1 \text { on }\left\{\left|x-\mathbf{s}_{k}\right| \leq \frac{\delta_{0}}{2}\right\} . \tag{2.32}
\end{equation*}
$$

We take $\epsilon$ in the definition (2.28) small enough then for all $k=1, \ldots N-1$, we have

$$
\begin{equation*}
\theta_{k} \chi_{k}^{\prime}=\chi_{k,+}^{\prime} \text { and } \theta_{k} \chi_{k+1}^{\prime}=\chi_{k+1,-}^{\prime} \tag{2.33}
\end{equation*}
$$

where $\chi_{k, \pm}$ are the smooth functions defined by

$$
\chi_{k,+}(x)=\left\{\begin{array}{cl}
\chi_{k}(x) & \text { if } x \geq m_{k}, \\
1 & \text { if } x<m_{k},
\end{array} \quad \chi_{k,-}(x)=\left\{\begin{array}{cl}
\chi_{k}(x) & \text { if } x \leq m_{k} \\
1 & \text { if } x>m_{k}
\end{array}\right.\right.
$$

Moreover, we also have $\theta_{k} \theta_{l}=0$ for all $k \neq l$. The family of quasi-modes associated to these cut-off functions is given by

$$
\begin{equation*}
f_{k}^{(1)}(x):=h^{-\frac{1}{4}} d_{k}(h) \theta_{k}(x) e^{(\varphi(x)-S) / h}, \quad\left\|f_{k}^{(1)}\right\|_{L^{2}}=1 \tag{2.34}
\end{equation*}
$$

where $d_{k}(h)=\left|\varphi^{\prime \prime}\left(s_{k}\right)\right|^{\frac{1}{4}} \pi^{-\frac{1}{4}}+\mathcal{O}(h)$ is the normalizing constant. Again, we introduce the projection of these quasi-modes onto the eigenspace $E^{(1)}$ :

$$
\begin{equation*}
g_{k}^{(1)}(x):=\Pi^{(1)} f_{k}^{(1)} \tag{2.35}
\end{equation*}
$$

Lemma 3. There exists $\alpha>0$ independant of $\epsilon$ such that the following hold true

$$
\begin{array}{cl}
\left\langle f_{k}^{(1)}, f_{l}^{(1)}\right\rangle=\delta_{k, l}, & \forall k, l=1, \ldots, N-1, \\
d_{\varphi}^{*} f_{k}^{(1)}=\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right), \quad g_{k}^{(1)}-f_{k}^{(1)}=\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)
\end{array}
$$

Proof. The proof follows the same lines as the proof of Lemma 2.
2.5. Computation of the operator $\mathcal{L}$. In this section we represent $\mathcal{L}$ in a suitable basis. For that we first observe that the bases $\left(g_{n}^{(0)}\right)$ and $\left(g_{k}^{(1)}\right)$ are quasi-orthonormal. Indeed, thanks to Lemmas 2 and 3, we have

$$
\left\langle g_{n}^{(0)}, g_{m}^{(0)}\right\rangle=\delta_{n, m}+\mathcal{O}\left(e^{-\alpha / h}\right), \quad \forall n, m=1, \ldots, N
$$

and

$$
\left\langle g_{k}^{(1)}, g_{l}^{(1)}\right\rangle=\delta_{k, l}+\mathcal{O}\left(e^{-\alpha / h}\right), \quad \forall k, l=1, \ldots, N-1
$$

for some $\alpha>0$. We then obtain orthonormal bases of $E^{(0)}$ and $E^{(1)}$ :

$$
\begin{gather*}
\quad\left(g_{n}^{(0)}\right)_{1 \leq n \leq N} \xrightarrow{\text { Gramm-Schmidt process }}\left(e_{n}^{(0)}\right)_{1 \leq n \leq N}, \\
\left(g_{k}^{(1)}\right)_{1 \leq k \leq N-1} \xrightarrow{\text { Gramm-Schmidt process }}\left(e_{k}^{(1)}\right)_{1 \leq n \leq N-1} . \tag{2.36}
\end{gather*}
$$

It follows from the approximate orthonormality above that the change of basis matrix $P_{j}$ from $\left(g_{n}^{(j)}\right)$ to $\left(e_{n}^{(j)}\right)$ satisfies

$$
\begin{equation*}
P_{j}=I+\mathcal{O}\left(e^{-\alpha / h}\right) \tag{2.37}
\end{equation*}
$$

for $j=0,1$. To describe the matrix of $\mathcal{L}$ in the bases $\left(e_{n}^{(0)}\right)$ and $\left(e_{k}^{(1)}\right)$ we introduce a $N-1 \times N$ matrix $\hat{L}=\left(\hat{\ell}_{i j}\right)$ defined by

$$
\begin{equation*}
\hat{\ell}_{i j}=\left\langle f_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle \tag{2.38}
\end{equation*}
$$

We claim that the matrices $L$ and $\hat{L}$ are very close. To see that we give a precise expansion of $\hat{L}$ :

Lemma 4. The matrix $\hat{L}$ defined by (2.38) is given by $\hat{L}=(h / \pi)^{\frac{1}{2}} e^{-S / h} \bar{L}$ where $\bar{L}$ admits a classical expansion $\bar{L} \sim \sum_{k=0}^{\infty} h^{k} L_{k}$ with

$$
L_{0}=\left(\begin{array}{cccccc}
-\nu_{1}^{\frac{1}{4}} \mu_{1}^{\frac{1}{4}} & \nu_{1}^{\frac{1}{4}} \mu_{2}^{\frac{1}{4}} & 0 & 0 & \ldots & 0  \tag{2.39}\\
0 & -\nu_{2}^{\frac{1}{4}} \mu_{2}^{\frac{1}{4}} & \nu_{2}^{\frac{1}{4}} \mu_{3}^{\frac{1}{4}} & 0 & \ldots & 0 \\
\vdots & \ddots & \ddots & \ddots & \ddots & \vdots \\
\vdots & \ddots & \ddots & \ddots & \ddots & 0 \\
0 & 0 & \ldots & 0 & -\nu_{n-1}^{\frac{1}{4}} \mu_{n-1}^{\frac{1}{4}} & \nu_{n-1}^{\frac{1}{4}} \mu_{n}^{\frac{1}{4}}
\end{array}\right)
$$

Proof. From (2.29) and (2.34), we have

$$
\begin{aligned}
\hat{\ell}_{i j} & =\left\langle f_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle=h^{-\frac{1}{2}} d_{i}(h) c_{j}(h) \int_{\mathbb{R}} \theta_{i}(x) e^{(\varphi(x)-S) / h} d_{\varphi}\left(\chi_{j}(x) e^{-\varphi(x) / h}\right) d x \\
& =h^{\frac{1}{2}} d_{i}(h) c_{j}(h) e^{-S / h} \int_{\mathbb{R}} \theta_{i}(x) \chi_{j}^{\prime}(x) d x
\end{aligned}
$$

Moreover, since $\operatorname{supp} \theta_{i} \cap \operatorname{supp} \chi_{j}=\emptyset$ except for $j=i$ or $j=i+1$, it follows from (2.33) that

$$
\begin{equation*}
\int_{\mathbb{R}} \theta_{i}(x) \chi_{j}^{\prime}(x) d x=\delta_{i, j} \int_{\mathbb{R}} \chi_{i,+}^{\prime}(x) d x+\delta_{i+1, j} \int_{\mathbb{R}} \chi_{i,-}^{\prime}(x) d x=-\delta_{i, j}+\delta_{i+1, j} \tag{2.40}
\end{equation*}
$$

On the other hand, we recall that $d_{i}(h)$ and $c_{j}(h)$ both have a classical expansion. Together with the above equality, this shows that $\hat{L}$ has the required form and it remains to prove the formula giving $L_{0}$. To that end we observe that

$$
d_{i}(h) c_{j}(h)=\pi^{-\frac{1}{2}}\left(\left(\left|\varphi^{\prime \prime}\left(s_{i}\right)\right| \varphi^{\prime \prime}\left(m_{j}\right)\right)^{\frac{1}{4}}+\mathcal{O}(h)\right)=\mu_{j}^{\frac{1}{4}} \nu_{i}^{\frac{1}{4}} \pi^{-\frac{1}{2}}+\mathcal{O}(h)
$$

in the notation of (2.1). Combining this with (2.40) we obtain

$$
\hat{\ell}_{i j}=h^{\frac{1}{2}} \pi^{-\frac{1}{2}} e^{-S / h} \mu_{j}^{\frac{1}{4}} \nu_{i}^{\frac{1}{4}}\left(-\delta_{i, j}+\delta_{i+1, j}+\mathcal{O}(h)\right)
$$

which gives (2.39).
Lemma 5. Let $L$ be the matrix of $\mathcal{L}$ in the basis obtained in (2.36). There exists $\alpha^{\prime}>0$ such that $L=\hat{L}+\mathcal{O}\left(e^{-\left(S+\alpha^{\prime}\right) / h}\right)$, where $\hat{L}$ is defined by (2.38) and is described in Lemma 4.

Proof. It follows from (2.37) that

$$
\begin{equation*}
L=\left(I+\mathcal{O}\left(e^{-\alpha / h}\right)\right) \tilde{L}\left(I+\mathcal{O}\left(e^{-\alpha / h}\right)\right) \tag{2.41}
\end{equation*}
$$

where $\tilde{L}=\left(\tilde{\ell}_{i, j}\right)$ with $\tilde{\ell}_{i, j}=\left\langle g_{i}^{(1)}, d_{\varphi} g_{j}^{(0)}\right\rangle$. Moreover, (2.25) implies that $\Pi^{(1)} d_{\varphi}=$ $d_{\varphi} \Pi^{(0)}$. Using this identity and the fact that $\Pi^{(0)}, \Pi^{(1)}$ are orthogonal projections, we have

$$
\begin{aligned}
\left\langle g_{i}^{(1)}, d_{\varphi} g_{j}^{(0)}\right\rangle & =\left\langle g_{i}^{(1)}, d_{\varphi} \Pi^{(0)} f_{j}^{(0)}\right\rangle=\left\langle g_{i}^{(1)}, \Pi^{(1)} d_{\varphi} f_{j}^{(0)}\right\rangle=\left\langle g_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle \\
& =\left\langle f_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle+\left\langle g_{i}^{(1)}-f_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle
\end{aligned}
$$

But from Lemmas 2, 3 and the Cauchy-Schwarz inequality we get

$$
\left|\left\langle g_{i}^{(1)}-f_{i}^{(1)}, d_{\varphi} f_{j}^{(0)}\right\rangle\right| \leq C e^{-\left(\alpha+S-\epsilon^{\prime}\right) / h}
$$

Since $\alpha$ is independent of $\epsilon^{\prime}$ which can be chosen as small as we want, it follows that there exists $\alpha^{\prime}>0$ such that $\tilde{\ell}_{i, j}=\hat{\ell}_{i, j}+\mathcal{O}\left(e^{-\left(S+\alpha^{\prime}\right) / h}\right)$. Combining this estimate, (2.41) and the fact that $\hat{\ell}_{i, j}=\mathcal{O}\left(e^{-S / h}\right)$, we get the announced result.

It is now easy to describe $\mathcal{M}$ as a matrix:
Lemma 6. Let $M$ be the matrix representation of $\mathcal{M}$ in the basis $\left(e_{n}^{(0)}\right)$. Then

$$
M=h e^{-2 S / h} A
$$

where $A$ is symmetric positive with a classical expansion $A \sim \sum_{k=0}^{\infty} h^{k} A_{k}$ with $A_{0}$ given by (2.2).

Proof. By definition, $M=L^{*} L$ and it follows from Lemma 4 and 5 that

$$
L^{*} L=\left(\hat{L}+\mathcal{O}\left(e^{-\left(S+\alpha^{\prime}\right) / h}\right)\right)^{*}\left(\hat{L}+\mathcal{O}\left(e^{-\left(S+\alpha^{\prime}\right) / h}\right)\right)=h e^{-2 S / h}\left(\bar{L}^{*} \bar{L}+\mathcal{O}\left(e^{-\alpha^{\prime} / h}\right)\right)
$$

Then, $A:=h^{-1} e^{2 S / h} L^{*} L$ is clearly positive and admits a classical expansion since $\bar{L}$ does. Moreover, the leading term of this expansion is $\bar{L}_{0}^{*} \bar{L}_{0}$ and a simple computation shows that $\bar{L}_{0}^{*} \bar{L}_{0}=A_{0}$, where $A_{0}$ is given by (2.2).

Remark. Innocent as this lemma might seem, the supersymmetric structure, that is writing $-\left.\Delta_{\varphi}\right|_{E^{(0)}}$ using $d_{\varphi}$, is very useful here.

Lemma 7. Denote by $\mu_{1}(h) \leq \ldots \leq \mu_{k}(h)$ the eigenvalues of $A(h)$. Then,

$$
\mu_{0}(h)=0 \text { and } \mu_{k}(h)=\mu_{k}^{0}+\mathcal{O}(h), \quad \forall k \geq 2
$$

where $0=\mu_{1}^{0}<\mu_{2}^{0} \leq \mu_{3}^{0} \leq \ldots \leq \mu_{N}^{0}$ denote the eigenvalues of $A_{0}$. Moreover, a normalized eigenvector associated to $\mu_{1}^{0}$ is $\xi^{0}=N^{-\frac{1}{2}}(1, \ldots, 1)$ and there exists a normalized vector $\xi(h) \in \operatorname{ker}(A(h))$, such that

$$
\begin{equation*}
\xi(h)=\xi_{0}+\mathcal{O}(h) \tag{2.42}
\end{equation*}
$$

Proof. Many of the statements of this lemma are immediate consequence of Lemma 6. We emphasize the fact that 0 belongs to $\sigma(A)$ since $0 \in \sigma\left(\Delta_{\varphi}\right)$. The fact that $\xi^{0}$ is in the kernel of $A_{0}$ is a simple computation. Eventually, for any $\xi \in \operatorname{ker}(A(h))$, we have

$$
\xi-\left\langle\xi, \xi_{0}\right\rangle \xi_{0}=\frac{1}{2 i \pi}\left(\int_{\gamma} z^{-1} \xi d z-\int_{\gamma}\left(A_{0}-z\right)^{-1} \xi d z\right)=\frac{1}{2 i \pi} \int_{\gamma}\left(A_{0}-z\right)^{-1} z^{-1} A_{0} \xi d z
$$

where $\gamma$ is a small path around 0 in $\mathbb{C}$. Since $A_{0} \xi=\mathcal{O}(h)$ we obtain (2.42).
2.6. Proof of Theorem 2. Let $u$ be solution of (2.3) with $u_{0}=\Psi(\beta),|\beta| \leq 1$ (see (2.4),(2.6) and (2.29) for definitions of $\psi_{n}, \Psi$ and $f_{n}^{(0)}$, respectively). Then,

$$
\begin{align*}
u & =e^{-t \Delta_{\varphi}} \Pi^{(0)} u_{0}+e^{-t \Delta_{\varphi}} \widehat{\Pi}^{(0)} u_{0} \\
& =e^{-t \mathcal{M}} \Pi^{(0)} u_{0}+e^{-t \Delta_{\varphi}} \widehat{\Pi}^{(0)} u_{0}, \quad \widehat{\Pi}^{(0)}:=I-\Pi^{(0)} . \tag{2.43}
\end{align*}
$$

Since $\mathbb{1}_{E_{n}}-\chi_{n}$ is supported near $\left\{s_{n-1}, s_{n}\right\}$, then $\psi_{n}-f_{n}^{(0)}=\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)$, for all $n=1, \ldots, N$, and it follows that

$$
u_{0}=\bar{u}_{0}+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right), \quad \bar{u}_{0}:=\sum_{n=1}^{N} \beta_{n} f_{n}^{(0)} .
$$

Then, using Lemma 2 and (2.37), we get $u_{0}=\tilde{u}_{0}+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)$ with $\tilde{u}_{0}:=\sum_{n=1}^{N} \beta_{n} e_{n}^{(0)}$, where $e_{n}^{(0)}$ is the orthonormal basis of $E^{(0)}$ given by (2.36). Since $\Pi^{(0)} e_{n}^{(0)}=e_{n}^{(0)}$ and $\widehat{\Pi}^{(0)} e_{n}^{(0)}=0$, we have

$$
u(t)=e^{-t \mathcal{M}} \tilde{u}_{0}+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)
$$

If $M$ is the matrix of the operator $\mathcal{M}$ in the basis $\left(e_{n}^{(0)}\right)$ then

$$
u(t)=\sum_{n=1}^{N}\left(e^{-t M} \beta\right)_{n} e_{n}^{(0)}+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)
$$

Going back from $e_{n}^{(0)}$ to $\psi_{n}$ as above, we see that

$$
\begin{equation*}
u(t)=\sum_{n=1}^{N}\left(e^{-t M} \beta\right)_{n} \psi_{n}+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right)=\Psi\left(e^{-t M} \beta\right)+\mathcal{O}_{L^{2}}\left(e^{-\alpha / h}\right) \tag{2.44}
\end{equation*}
$$

and the proof of (2.7) (main statement in Theorem 2) is complete. We now prove (2.8). Since the linear map $\psi: \mathbb{C}^{N} \rightarrow L^{2}(d x)$ is bounded uniformly with respect to $h$, and thanks to (2.7), the proof reduces to showing (after time rescaling) that there exists $C>0$ such that for all $\beta \in \mathbb{R}^{N}$,

$$
\begin{equation*}
\left|e^{-\tau A}-e^{-\tau A_{0}}\right| \leq C h, \quad \forall \tau \geq 0 \tag{2.45}
\end{equation*}
$$

Since, by Lemma $7, A$ and $A_{0}$ both have 0 as a simple eigenvalue with the approximate eigenvector given by $(1, \cdots, 1)$, we see that for any norm on $\mathbb{C}^{N}$,

$$
\begin{aligned}
\left|e^{-\tau A}-e^{-\tau A_{0}}\right| & \leq\left|e^{-\tau A_{0}}\right|_{\{(1, \cdots, 1)\}^{\perp}}\left|I-e^{-\tau \mathcal{O}(h)}\right|_{\ell^{2} \rightarrow \ell^{2}}+C h \\
& \leq C e^{-c \tau} \tau h+C h=\mathcal{O}(h) .
\end{aligned}
$$

which is exactly (2.45).
We now prove one of the consequences of Theorem 2.
Proof of (1.11). We have seen in the preceding proof that $e_{n}^{(0)}-\psi_{n}=\mathcal{O}_{L^{2}}\left(e^{-C / \epsilon^{2}}\right)$ and since

$$
\left\|\psi_{n}-\left(\mu / 2 \pi \epsilon^{2}\right)^{\frac{1}{4}} \mathbb{1}_{E_{n}} e^{-\varphi / 2 \epsilon^{2}}\right\|_{L^{2}}=\mathcal{O}\left(\epsilon^{2}\right)
$$

it follows that $\Pi^{(0)} u_{0}=\psi(\beta)+\mathcal{O}\left(\epsilon^{2}\left\|u_{0}\right\|_{L^{2}}\right)$ with $\beta \in \mathbb{C}^{N}$ given by

$$
\beta_{n}=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \int_{E_{n}} u_{0}(x) e^{-\varphi(x) / 2 \epsilon^{2}} d x=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \int_{E_{n}} \rho_{0}(x) d x .
$$

Applying 2 (second part of Theorem 2) with $h=2 \epsilon^{2}$ gives

$$
\begin{equation*}
e^{-t \Delta_{\varphi}} \Pi^{(0)} u_{0}=\sum_{n=1}^{N}\left(e^{-t \nu_{h} A_{0}} \beta\right)_{n}\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \mathbb{1}_{E_{n}} e^{-\varphi / 2 \epsilon^{2}}+\mathcal{O}_{L^{2}}\left(\epsilon^{2}\right)\left\|u_{0}\right\|_{L^{2}} \tag{2.46}
\end{equation*}
$$

On the other hand, Proposition 1 shows that

$$
\begin{equation*}
e^{-t \Delta_{\varphi}}\left(I-\Pi^{(0)}\right) u_{0}=\mathcal{O}_{L^{2}}\left(e^{-t \epsilon^{2} / C}\right)\left\|u_{0}\right\|_{L^{2}} \tag{2.47}
\end{equation*}
$$

Since $\rho\left(h^{2} t\right)=e^{-\varphi / h} u(t),(2.46)$ and (2.47) yield

$$
\begin{equation*}
\rho\left(2 \epsilon^{2} e^{S / \epsilon^{2}} \tau\right)=\sum_{n=1}^{N}\left(e^{-\tau A_{0}} \beta\right)_{n}\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{4}} \mathbb{1}_{E_{n}} e^{-\varphi / \epsilon^{2}}+r_{\epsilon}(\tau) \tag{2.48}
\end{equation*}
$$

with

$$
r_{\epsilon}(\tau)=e^{-\varphi / 2 \epsilon^{2}}\left(\mathcal{O}_{L^{2}}\left(e^{-c \tau e^{S / \epsilon^{2}}}\right)+\mathcal{O}_{L^{2}}\left(\epsilon^{2}\right)\right)\left\|\rho_{0}\right\|_{L_{\varphi}^{2}}
$$

By Cauchy-Schwartz it follows that $\left\|r_{\epsilon}(\tau)\right\|_{L^{1}} \leq C\left(\epsilon^{\frac{5}{2}}+e^{-c \tau e^{S / \epsilon^{2}}}\right)\left\|\rho_{0}\right\|_{L_{\varphi}^{2}}$.

## 3. A higher dimensional example

The same principles apply when the wells may have different height and in higher dimensions. In both cases there are interesting combinatorial and topological (when $d>1)$ complications and we refer to $[15, \S 1.1, \S 1.2]$ for a presentation and references. To illustrate this we give a higher dimensional result in a simplified setting.

Suppose that $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is a smooth Morse function satisfying (1.2) and denote by $\mathcal{U}^{(j)}$ the finite sets of critical points of index $j, n_{j}:=\left|\mathcal{U}^{(j)}\right|$. We assume that (1.5) holds and write $S:=\sigma_{1}-\varphi_{0}$. In the notation of (1.4) we have $n_{0}=N$ and we also assume


Figure 2. Left: The sublevel set $\left\{\varphi<\sigma_{1}\right\}$ (dashed region) associated to a potential $\varphi$ (1.5). The x's represent local minima, the o's, local maxima. Right: The graph associated to the potential on the left.
that each $E_{n}$ contains exactly one minimum. Hence we can label the components by the minima:

$$
\forall n=1, \ldots, N, \quad \exists!\mathbf{m} \in E_{n}, \min _{x \in E_{n}} \varphi(x)=\varphi(\mathbf{m})
$$

and we denote $E(\mathbf{m}):=E_{n}$. Since $\varphi$ is a Morse function,

$$
\begin{gather*}
\forall \mathbf{m}, \mathbf{m}^{\prime} \in \mathcal{U}^{(0)}, \quad \mathbf{m} \neq \mathbf{m}^{\prime} \Longrightarrow \bar{E}(\mathbf{m}) \cap \bar{E}\left(\mathbf{m}^{\prime}\right) \subset \mathcal{U}^{(1)} \\
\forall \mathbf{s} \in \mathcal{U}^{(1)}, \quad \exists!\mathbf{m}, \mathbf{m}^{\prime} \in \mathcal{U}^{(0)}, \quad \mathbf{s} \in \bar{E}(\mathbf{m}) \cap \bar{E}\left(\mathbf{m}^{\prime}\right) \tag{3.1}
\end{gather*}
$$

To simplify the presentation we make an addition assumption

$$
\begin{equation*}
\forall \mathbf{m}, \mathbf{m}^{\prime} \in \mathcal{U}^{(0)}, \quad \mathbf{m} \neq \mathbf{m}^{\prime} \Longrightarrow\left|\bar{E}(\mathbf{m}) \cap \bar{E}\left(\mathbf{m}^{\prime}\right)\right| \leq 1 \tag{3.2}
\end{equation*}
$$

Under these assumptions, the set $\mathcal{U}^{(0)} \times \mathcal{U}^{(1)}$ defines a graph $\mathcal{G}$. The elements of $\mathcal{U}^{(0)}$ are the vertices of $\mathcal{G}$ and elements of $\mathcal{U}^{(1)}$ are the edges of $\mathcal{G}$ : $\mathbf{s} \in \mathcal{U}^{(1)}$ is an edge between $\mathbf{m}$ and $\mathbf{m}^{\prime}$ in $\mathcal{U}^{(0)}$ if $s \in \bar{E}(\mathbf{m}) \cap E\left(\mathbf{m}^{\prime}\right)$ - see Fig. 2 for an example.

The same graph has been constructed in [14] for a certain discrete model of the Kramers-Smoluchowski equation.

We now introduce the discrete Laplace operator on $\mathcal{G}, M_{\mathcal{G}}$ - see [3] for the background and results about $M_{\mathcal{G}}$. If the degree $d(\mathbf{m})$ is defined as the number of edges at the vertex $\mathbf{m}, M_{\mathcal{G}}$ is given by the matrix $\left(a_{\mathbf{m}, \mathbf{m}^{\prime}}\right)_{\mathbf{m}, \mathbf{m}^{\prime} \in \mathcal{U}^{(0)}}$ :

$$
a_{\mathbf{m}, \mathbf{m}^{\prime}}=\left\{\begin{array}{cl}
d(\mathbf{m}), & \mathbf{m}=\mathbf{m}^{\prime}  \tag{3.3}\\
-1 & \mathbf{m} \neq \mathbf{m}^{\prime}, \quad \bar{E}(\mathbf{m}) \cap \bar{E}\left(\mathbf{m}^{\prime}\right) \neq \emptyset \\
0 & \text { otherwise }
\end{array}\right.
$$

Among basic properties of the matrix $M_{\mathcal{G}}$, we recall:

- it has a square structure $M_{\mathcal{G}}=\mathcal{L}^{*} \mathcal{L}$, where $\mathcal{L}$ is the transpose of the incidence matrix of any oriented version of the graph $\mathcal{G}$. In particular, $M_{\mathcal{G}}$ is symmetric positive.
- thanks to (1.5) and [15, Proposition B.1], the graph $\mathcal{G}$ is connected.
- 0 is a simple eigenvalue of $M_{\mathcal{G}}$

We make one more assumption which is a higher dimensional analogue of the hypothesis in Theorem 1: there exist $\mu, \nu>0$ such that

$$
\begin{align*}
& \operatorname{det} \varphi^{\prime \prime}(\mathbf{m})=\mu, \forall \mathbf{m} \in \mathcal{U}^{(0)} \\
& \frac{\lambda_{1}(\mathbf{s})^{2}}{\operatorname{det} \varphi^{\prime \prime}(\mathbf{s})}=-\nu, \quad \forall \mathbf{s} \in \mathcal{U}^{(1)} \tag{3.4}
\end{align*}
$$

where $\lambda_{1}(\mathbf{s})$ is the unique negative eigenvalue of $\varphi^{\prime \prime}(\mathbf{s})$. Assumptions (3.2) and (3.4) can be easily removed. Without (3.4) the graph $\mathcal{G}$ is replaced by a weighted graph with a weight function depending explicitly of the values of $\varphi^{\prime \prime}$ at critical points. Removing (3.2) leads to multigraphs in which there may be several edges between two vertices. This can be also handled easily.

Assumption (1.5) however is more fundamental and removing it results in major complications. We refer to [15] for results in that situation. Here we restrict ourselves to making the following
Remark. Under the assumption (1.5) the proof presented in the one dimensional case applies with relatively simple modifications. The serious difference lies in the description of $E^{(1)}$, the eigenspace of $\Delta_{\varphi}$ on one-forms, in terms of exponentially accurate quasi-modes (in one dimension it was easily done using Lemma 3). That description is however provided by Helffer-Sjöstrand in the self contained Section 2.2 of [9] - see Theorem 2.5 there. The computation of (2.38) becomes more involved and is based on the method of stationary phase - see Helffer-Klein-Nier [8, Proof of Proposition 6.4].

The analogue of Theorem 1 is
Theorem 3. Suppose that $\varphi$ satisfies (1.2),(1.5),(3.2) and (3.4). If

$$
\begin{equation*}
\rho_{0}=\left(\frac{\mu}{2 \pi \epsilon^{2}}\right)^{\frac{1}{2}}\left(\sum_{n=1}^{N} \beta_{n} \mathbb{1}_{E_{n}}+r_{\epsilon}\right) e^{-\varphi / \epsilon^{2}}, \quad \lim _{\epsilon \rightarrow 0}\left\|r_{\epsilon}\right\|_{L^{\infty}}=0, \quad \beta \in \mathbb{R}^{N} \tag{3.5}
\end{equation*}
$$

then the solution to (1.1) satisfies, uniformly for $\tau \geq 0$,

$$
\begin{equation*}
\rho\left(2 \epsilon^{2} e^{S / \epsilon^{2}} \tau, x\right) \rightarrow \sum_{n=1}^{N} \alpha_{n}(\tau) \delta_{m_{n}}(x), \quad \epsilon \rightarrow 0 \tag{3.6}
\end{equation*}
$$

in the sense of distributions in $x$, where $\alpha(t)=\left(\alpha_{1}, \ldots, \alpha_{n}\right)(\tau)$ solves

$$
\begin{equation*}
\partial_{\tau} \alpha=-\kappa M_{\mathcal{G}} \alpha, \quad \alpha(0)=\beta, \tag{3.7}
\end{equation*}
$$

with $M_{\mathcal{G}}$ is given by (3.3) and $\kappa=\pi^{-1} \mu^{\frac{1}{2}} \nu^{\frac{1}{2}}$ with $\mu$ and $\nu$ in (3.4).
We also have the analogue of (1.11) for any initial data.
As in the one dimensional case this theorem is a consequence of a more precise theorem formulated using the localized states

$$
\begin{equation*}
\psi_{n}(x)=c_{n}(h) h^{-\frac{d}{4}} \mathbb{1}_{E_{n}}(x) e^{-\left(\varphi-\varphi_{0}\right)(x) / h} \tag{3.8}
\end{equation*}
$$

where $c_{n}(h)$ is a normalization constant such that $\left\|\psi_{n}\right\|_{L^{2}}=1$. We then define a map $\Psi: \mathbb{R}^{N} \rightarrow L^{2}\left(\mathbb{R}^{d}\right)$ by

$$
\begin{equation*}
\Psi(\beta)=\sum_{n=1}^{N} \beta_{n} \psi_{n}, \quad \forall \beta=\left(\beta_{1}, \ldots, \beta_{N}\right) \in \mathbb{R}^{N} \tag{3.9}
\end{equation*}
$$

We have the following analogue of Theorem 2.
Theorem 4. Suppose $\varphi$ satisfies (1.2),(1.5),(3.2) and (3.4). There exists $C>0$ and $h_{0}>0$ such that for all $\beta \in \mathbb{R}^{N}$ and all $0<h<h_{0}$, we have

$$
\left\|e^{-t \Delta_{\varphi}} \Psi(\beta)-\Psi\left(e^{-t \kappa \nu_{h} A} \beta\right)\right\|_{L^{2}} \leq C e^{-1 / C h}, \quad t \geq 0
$$

where $\nu_{h}=h e^{-2 S / h}, \kappa=\pi^{-1} \mu^{\frac{1}{2}} \nu^{\frac{1}{2}}$ and $A=A(h)$ is a real symmetric positive matrix having a classical expansion $A \sim \sum_{k=0}^{\infty} h^{k} A_{k}$ and $A_{0}=M_{\mathcal{G}}$ with $M_{\mathcal{G}}$ the Laplace matrix defined by (3.3).


Figure 3. A two dimensional potential which is a cyclic analogue of the potential shown in Fig.1: the corresponding matrix describing the Kramer-Smoluchowski evolution is given by (3.10). It should compared to the matrix (1.7) for the potential in Fig.1. The corresponding cyclic graph is shown on the right.

We conclude by one example $[15, \S 6.3]$ for which the graph $\mathcal{G}$ is elementary. We assume that $d=2, \varphi$ has a maximum at $x=0$, there are $N$ minima, $m_{n}, N$ saddle points, $s_{n}$, and that (1.5) holds - see Fig.3. We assume also that

$$
\operatorname{det} \varphi^{\prime \prime}\left(\mathbf{m}_{n}\right)=\mu>0, \quad \frac{\lambda_{1}\left(\mathbf{s}_{n}\right)}{\lambda_{2}\left(\mathbf{s}_{n}\right)}=-\nu<0
$$

where for $\mathbf{s} \in \mathcal{U}^{(1)}, \lambda_{1}(\mathbf{s})>0>\lambda_{2}(\mathbf{s})$ denote the two eigenvalues of $\varphi^{\prime \prime}(\mathbf{s})$.

Then assumptions of Theorem 4 are satisfied. The graph $\mathcal{G}$ associated to $\varphi$ is the cyclic graph with $N$ vertices and the corresponding Laplacian is given by

$$
\mathcal{A}_{\mathcal{G}}=\left(\begin{array}{rrrrrrrr}
2 & -1 & 0 & 0 & \ldots & \ldots & \ldots & -1  \tag{3.10}\\
-1 & 2 & -1 & 0 & \ldots & \ldots & \ldots & 0 \\
0 & -1 & 2 & -1 & 0 & \ldots & \ldots & 0 \\
\vdots & 0 & -1 & 2 & \ddots & \ddots & \ldots & 0 \\
\vdots & \vdots & \ddots & \ddots & \ddots & \ddots & \ddots & 0 \\
\vdots & \vdots & \ddots & \ddots & \ddots & \ddots & -1 & 0 \\
0 & \vdots & \ddots & \ddots & \ddots & -1 & 2 & -1 \\
1 & 0 & \ldots & \ldots & \ldots & 0 & -1 & 2
\end{array}\right)
$$

## References

1. N. Berglund, Kramers' law: validity, derivations and generalisations, Markov Process. Related Fields 19 (2013), no. 3, 459-490.
2. A. Bovier, V. Gayrard, and M. Klein, Metastability in reversible diffusion processes. II. Precise asymptotics for small eigenvalues, J. Eur. Math. Soc. 7 (2005), no. 1, 69-99.
3. D. M. Cvetković, M. Doob, and H. Sachs, Spectra of graphs, third ed., Johann Ambrosius Barth, Heidelberg, 1995, Theory and applications.
4. H. Cycon, R. Froese, W. Kirsch, and B. Simon, Schrödinger operators with application to quantum mechanics and global geometry, study ed., Texts and Monographs in Physics, Springer-Verlag, 1987.
5. L. C. Evans and P. Tabrizian, Asymptotics for scaled Kramers-Smoluchowski equations, SIAM J. Math. Anal. 48 (2016), no. 4, 2944-2961.
6. A.H. Eyring, The activated complex in chemical reactions, J. Chem. Phys. (1935), no. 3, 107-115.
7. B. Helffer, Semi-classical analysis for the Schrödinger operator and applications, Lecture Notes in Mathematics, vol. 1336, Springer-Verlag, 1988.
8. B. Helffer, M. Klein, and F. Nier, Quantitative analysis of metastability in reversible diffusion processes via a Witten complex approach, Mat. Contemp. 26 (2004), 41-85.
9. B. Helffer and J. Sjöstrand, Puits multiples en mécanique semi-classique. IV. Étude du complexe de Witten, Comm. Partial Differential Equations 10 (1985), no. 3, 245-340.
10. F. Hérau, M. Hitrik, and J. Sjöstrand, Tunnel effect and symmetries for Kramers-Fokker-Planck type operators, J. Inst. Math. Jussieu 10 (2011), no. 3, 567-634.
11. M. Herrmann and B. Niethammer, Kramers' formula for chemical reactions in the context of Wasserstein gradient flows, Commun. Math. Sci. 9 (2011), no. 2, 623-635.
12. S. Insuk and P. Tabrizian, Kramers-Smoluchowski equation for general potentials, preprint.
13. H. A. Kramers, Brownian motion in a field of force and the diffusion model of chemical reactions, Physica 7 (1940), 284-304.
14. C. Landim, R. Misturini, and K. Tsunoda, Metastability of reversible random walks in potential fields, J. Stat. Phys. 160 (2015), no. 6, 1449-1482.
15. L. Michel, On small eigenvalues of Witten Laplacian, arXiv:1702.01837.
16. M. A. Peletier, G. Savaré, and M. Veneroni, Chemical reactions as $\Gamma$-limit of diffusion, SIAM Rev. 54 (2012), no. 2, 327-352.
17. E. Witten, Supersymmetry and Morse theory, J. Differential Geom. 17 (1982), no. 4, 661-692 (1983).
18. M. Zworski, Semiclassical analysis, Graduate Studies in Mathematics, vol. 138, American Mathematical Society, 2012.

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