STRONG AND WEAK SEMICLASSICAL LIMIT FOR SOME ROUGH HAMILTONIANS

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ABSTRACT. We present several results concerning the semiclassical limit of the time dependent Schrödinger equation with potentials whose regularity doesn't guarantee the uniqueness of the underlying classical flow. Different topologies for the limit are considered and the situation where two bicharacteristics can be obtained out of the same initial point is emphasized.

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1. Introduction

Consider a wavefunction \( u^\varepsilon(x, t) : \mathbb{R}^n \times \mathbb{R} \to \mathbb{C} \) satisfying the Schrödinger equation,

\[
\tag{1.1}
i\varepsilon \frac{\partial}{\partial t} u^\varepsilon_t (x) = \left( -\frac{\varepsilon^2}{2} \Delta + V(x) \right) u^\varepsilon_t (x), \quad u^\varepsilon_t (x) = u^\varepsilon_0(x),
\]

or a mixed state \( D^\varepsilon_t \) satisfying the Heisenberg-von Neumann equation:

\[
\tag{1.2}
i\varepsilon \frac{\partial}{\partial t} D^\varepsilon_t = \left[ -\frac{\varepsilon^2}{2} \Delta + V, D^\varepsilon_t \right], \quad D^\varepsilon_t |_{t=0} = D^\varepsilon_0.
\]
The small parameter \( \varepsilon \) is called Planck’s constant in the context of quantum mechanics, and its physical meaning is that of the lengthscale of oscillation of the wavefunction in space.

Several techniques have been developed for the study of the semiclassical limit \( \varepsilon \to 0 \) of (1.1). One is based on the Wigner transform (WT) [9, 10, 13]. The WT \( W_\varepsilon^\varepsilon(x,k) \) of the wavefunction \( u_\varepsilon^\varepsilon \) is, defined in the standard semiclassical scaling,

\[
W_\varepsilon^\varepsilon(x,k) = W_\varepsilon^\varepsilon[u_\varepsilon^\varepsilon](x,k) = \int_{y \in \mathbb{R}^n} e^{-2\pi i y k} u_\varepsilon^\varepsilon(x + \varepsilon y / 2) \bar{u}_\varepsilon^\varepsilon(x - \varepsilon y / 2) dy.
\]

In that case \( W_\varepsilon^\varepsilon \) corresponds to a pure state. A Wigner function \( W_\varepsilon^\varepsilon \) need not be the Wigner transform of a wavefunction \( u_\varepsilon^\varepsilon \); when working with mixed states, i.e. with equation (1.2), the Wigner function can be

\[
W_\varepsilon^\varepsilon(x,k) = W_\varepsilon^\varepsilon[D_\varepsilon^\varepsilon](x,k) = \int_{y \in \mathbb{R}^n} e^{-2\pi i y k} K_\varepsilon^\varepsilon(x + \varepsilon y / 2, x - \varepsilon y / 2) dy
\]

where \( K_\varepsilon^\varepsilon(x,y) \) is the integral kernel of the mixed-state \( D_\varepsilon^\varepsilon \). It evolves in time under the well-known Wigner equation,

\[
\partial_t W_\varepsilon^\varepsilon(x,k) + 2\pi k \cdot \partial_x W_\varepsilon^\varepsilon(x,k) +
+ \frac{2}{\varepsilon} \text{Re} \left[ i \int e^{2\pi i S_x} \hat{V}(S) W_\varepsilon^\varepsilon(x,k - \frac{\varepsilon S}{2}) dS \right] = 0,
\]

\[W_{t=0}^\varepsilon = W_0^\varepsilon.\]

An important feature of the Wigner transform is that (under appropriate conditions) it allows for a very natural and compact description of the semiclassical limit: the WT has a physically meaningful limit as \( \varepsilon \) tends to zero, while in general the wavefunction \( u_\varepsilon^\varepsilon \) itself does not. The limit (in the weak-* topology of an appropriate algebra of test functions), called the Wigner measure,

\[
W_\varepsilon^\varepsilon \to W_0^0,
\]

satisfies the Liouville equation of classical mechanics

\[
\partial_t W_0^0 + 2\pi k \partial_x W_0^0 - \frac{1}{2\pi} \partial_x V \partial_k W_0^0 = 0,
\]

\[W_{t=0}^0 = W_0^0.\]

In general the Wigner measure \( W_0^0 \) is a non-negative measure on phase-space, associated with the family of initial data \( \{W_0^\varepsilon\} \). For the precise formulation of the result, see e.g. [13]. The limit problem (1.6) is well-posed in general for \( V(x) \in C^{1,1} \).

It was already observed in [13] that there are certain cases where semiclassical convergence to the Liouville equation could be shown, but the Liouville equation itself is not well posed – namely for \( V(x) \in C^1 \setminus C^{1,1} \). The result
then becomes as follows: the Wigner measure is a weak solution of the Liouville equation – but there is not enough information to select the correct weak solution. Families of initial data whose evolution oscillate between the several solutions of the Liouville equation were constructed in [13].

The one-dimensional case study considered in [13] was essentially equivalent to the following:

One considers the potential

\[ V(x) = -|x|^{1+\theta} \cdot \beta(x) \]

with \( \theta \in (0, 1) \), \( C^\infty \ni \beta = 1 \) on \([-1, +1]\), and \( \beta(x) = 0 \) for \(|x| \geq 2\). Then one can easily check that there exist, for \( t \) small enough, two families of bicharacteristics starting near the origin at \( t = 0 \), namely

\[ (X^\pm(t), P^\pm(t)) = (\pm c_0 t^\nu, \pm \frac{c_0 t^\nu}{2\pi}) \]

with \( \nu = \frac{2}{1-\theta} \) and \( c_0 = \left( \frac{(1-\theta)^2}{2} \right)^{1-\theta} \).

In this case it is possible to find semiclassical limits where one wavepacket (a \( \delta \)-function in classical phase-space) splits in two wavepackets. (In this particular example it happens immediately, i.e. \( \forall t > 0 \); different examples can be constructed where it happens after an \( \varepsilon \)-independent time \( t^* \)). The mass of each wavepacket depends on information which can be easily extracted from the quantum problem, but lost if we take the limit as usual in a straightforward way.

Weak semiclassical limits with rough potentials were also considered recently in [2] where the hypothesis on the potential is essentially, among other properties, that its gradient is in \( BV_{loc} \) (allowing for a Coulomb part to be added). It was proved there, roughly speaking, that the Wigner measure of the solutions of (1.1) at time \( t \), after averaging over initial data, tends weakly to the push-forward of its initial value by the flow associated to the underlying classical Hamiltonian system of ODEs by the Ambrosio-DiPerna-Lions theory [1, 7] (see also [8] for the same construction with density matrices). In particular this averaging condition forbids concentrating initial data, a fact consistent with the weak limit strategy and the fact that the flow is not defined everywhere.

Our results concern less general potentials than those treated in [2, 8] (still giving rise to ill-posed classical dynamics of course), but can deal with initial data concentrating to singular measures.

Other problems where the Wigner measure does not satisfy one well-posed Liouville equation have attracted attention recently. The difficulties can originate from lack of smoothness of the potential [11], or working on a torus instead of Euclidean space [14].
The first of our results deals with the case study from [13] mentioned earlier. Our second Theorem gives a general result concerning the approximation, in strong topology, of the Wigner function of the solution of (1.4) with the solution of a “smoothed” Liouville equation. Our third results goes back to the weak approximation, but with explicit remainder estimates.

2. Main results

We will use the space of test-functions $A$, from [13], namely

\[(2.1) \quad A = \{ \phi \in C(\mathbb{R}^{2n}) | \sup_x |F_{k \to K}[\phi(x,k)]|dK < \infty \},\]

equipped with the norm

\[(2.2) \quad ||\phi||_A = \int \sup_x |F_{k \to K}[\phi(x,k)]|dK.\]

Our first result will focus on the case study presented in the introduction; a family of initial data which concentrates at the origin at time $t = 0$ and splits into two separated wave packets at time $t > 0$. This is only an example among several others presented in section 5 and the more general results of the paper are given by the Theorems 2.2 and 2.3 below.

**Theorem 2.1.** Let $V$ given by (1.7) and

\[(2.3) \quad W^0_t(x,k) = \lambda^{\frac{7+3\theta}{9\theta}} w(\lambda^{\frac{1+\theta}{6}} x, \lambda^{\frac{1-\theta}{15}} k)\]

with $\lambda = \log(\frac{1}{\varepsilon})$, $w \in H^2 \cap L^\infty \cap L^1$, $\text{supp } w \subseteq \{|x|^2 + |k|^2 < 1\}$.

Then $\exists T > 0$ such that for all $t \in [0,T]$, the solution $W^\varepsilon_t$ of (1.5) converges in weak-$*$ sense in $A'$ to

\[(2.4) \quad W^\varepsilon_t = c_+ \delta_{(X^+(t),P^+(t))} + c_- \delta_{(X^-(t),P^-(t))},\]

with $(X^\pm(t), P^\pm(t))$ given by (1.8) and

\[c_\pm = \int_{x>0} w(x,k)dxdk.\]

In Remarque IV.3 of [13], one can find a construction of a family of normalized wave functions $u^+_\varepsilon$ concentrating at the origin, and whose Wigner functions, after extraction of a subsequence, converge to a Dirac mass centered on $(X^+(t), P^+(t))$. The same construction is obviously possible in order to get a family $u^-_\varepsilon$ whose Wigner functions, after extraction of a subsequence, converge to a Dirac mass centered on $(X^-(t), P^-(t))$. Defining $D^\varepsilon_0 = \frac{1}{\sqrt{2}}(|u^+_\varepsilon\rangle\langle u^+_\varepsilon| + |u^-_\varepsilon\rangle\langle u^-_\varepsilon|)$ one can check that the Wigner function of $D^\varepsilon_0$ will follow the conclusion of Theorem 2.1. However the present construction is different, more precise and is not obtained through compactness arguments: there is no need of extraction of subsequence, the scaling property is explicit (it is given by an implicit diagonal argument in a two scale sequence.
in [13]) and the choice of the initial datum is somewhat more general. The construction of [13] is based on the more-or-less explicit understanding of the flow around the singular point, while our approximation (i.e. Theorems 2.2 and 2.3 below) is built around a graceful degradation property (in a sense analogous to the use of the term, e.g., in [12]) of the flow (properties 1 and 2 in section 5), and applies automatically, whenever that property holds, to more general and less explicit potentials. Let us also remark the similarity with the long-time behaviour around a regular separatrix, studied extensively in [15].

We turn now to the more general results of this paper.

In the sequel we will consider potentials $V$ satisfies the following assumption:

**Assumption 1.**

\[
\int_{\mathbb{R}^n} |\hat{V}(S)| \frac{S^2}{1+S^2} dS < \infty,
\]

and moreover there are constants $C > 0$, $\theta \in (0, 1)$ such that for $m \in \{0, 1, 2\}$

\[
\forall 1 \leq a \leq b \leq +\infty : \int_{|S| \in (a,b)} |\hat{V}(S)||S|^m dS \leq \frac{C}{m-1-\theta} (b^{m-1-\theta} - a^{m-1-\theta})
\]

This is closely modeled after $V(x) = C|x|^{1+\theta}$; indeed it is easy to check that the aforementioned potential satisfies this condition. In section 5, we will also see some other relevant types of singularities (also generated from $|x|^{1+\theta}$ in some sense).

Denote

\[
\tilde{V}^\gamma(x) = \left(\frac{2}{\varepsilon^\gamma}\right)^{\frac{2}{\gamma}} \int e^{-\frac{2\pi}{\varepsilon^\gamma}|x-x'|^2} V(x')dx'.
\]

**Theorem 2.2.** Let us suppose that Assumption 1 holds, and there exist $T > 0$, $\delta \in (0, \frac{\theta}{2+\theta})$, $\gamma > \frac{2}{1+\theta}$ and $[W^\varepsilon_0] \in H^2(\mathbb{R}^{2n})$ such that

\[
|||W^\varepsilon_0||_{L^2} = o(1)||W^\varepsilon_0||_{L^2}
\]

and the solution of

\[
\partial_t \rho + 2\pi k \partial_x \rho - \frac{1}{2\pi} \partial_x \tilde{V}^\gamma \cdot \partial_x \rho = 0,
\]

with initial condition $\rho(t=0) = [W^\varepsilon_0]$ satisfies

\[
||\rho(t)||_{H^2} = O(\varepsilon^{-\delta}||W^\varepsilon_0||_{L^2})
\]

uniformly on $[0, T]$. 

Then the Wigner function $W_{}^\varepsilon_t$ of the solution of (1.4) satisfies, uniformly on $[0,T]$, 
\begin{equation}
||W_{}^\varepsilon_t - \rho_1^\varepsilon(t)||_{L^2} = O(\varepsilon \kappa ||W_0^\varepsilon||_{L^2} + ||W_0^\varepsilon - \rho_1^\varepsilon(t)||_{L^2}) = o(||W_0^\varepsilon||_{L^2}).
\end{equation}
where $\rho_1^\varepsilon$ is the solution of (2.8) with initial datum $\rho_1^\varepsilon(t = 0) = W_0^\varepsilon$ and
\[\kappa = \min\{\gamma \frac{1+\theta}{2} - 1, \frac{\theta}{2\pi\theta} - \delta\}\].

Our next result deals with the weak-$^\ast$ limit in $\varepsilon$ (which in general is not in $L^2$). It must be noted that the assumptions are strictly stronger than those of theorem 2.2, therefore the latter also holds.

**Theorem 2.3.** Assume that the hypotheses of Theorem 2.2 hold, let us suppose moreover that 
\begin{equation}
||[W_0^\varepsilon] - W_0^\varepsilon||_{L^2} = o(1),
\end{equation}
and

- $W_0^\varepsilon$ is compactly supported (uniformly in $\varepsilon$),
- $W_0^\varepsilon \in L^1 \cap L^2$,
- $W_0^\varepsilon$ is the Wigner function of a density matrix (positive operator of trace 1),
- $||W_0^\varepsilon||_{L^2} = o(\varepsilon^{-\kappa})$ where $\kappa$ is the number defined in Theorem 2.2.
- the total energy is uniformly bounded in $\varepsilon$,
\[\exists C_E : \int (2\pi^2 k^2 + V(x))W_0^\varepsilon(x,k)dk < C_E,\]
- the potential $V$ is in $W^{1,\infty} \cap H^1$, and bounded below.

Then, $\forall \phi \in \mathcal{A}$ and $\forall t \in [0,T]$, the Wigner function $W_{}^\varepsilon_t$ of the solution of (1.4) satisfies 
\begin{equation}
\langle W_{}^\varepsilon_t, \phi \rangle = \langle \rho_1^\varepsilon(t), \phi \rangle + E^\varepsilon,
\end{equation}
where again $\rho_1^\varepsilon$ satisfies (2.8) with initial datum $\rho_1^\varepsilon(t = 0) = W_0^\varepsilon$.

The error $E^\varepsilon$ can be estimated explicitly as follows: if $\mu = (\varepsilon^\kappa ||W_0^\varepsilon||_{L^2} + ||W_0^\varepsilon - \rho_1^\varepsilon(t)||_{L^2} + \sqrt{\varepsilon}||W_0^\varepsilon||_{H^1})^{\frac{1}{2}}$, 
\begin{equation}
E^\varepsilon = O\left(\mu||\phi||_{\mathcal{A}} + ||\langle W^\varepsilon, (\mathbb{1} - e^{-\varepsilon\Delta}) \phi \rangle ||\right).
\end{equation}

**Corollary 2.4.** If the assumptions of Theorem 2.3 are satisfied, the Wigner function $W^\varepsilon$ has a weak-$^\ast$ limit in $\mathcal{A}'$ if and only if $\rho_1^\varepsilon$ has a weak-$^\ast$ limit, and the two are equal, i.e.
\begin{equation}
\exists W_0^0 : \lim_{\varepsilon \to 0} \langle W_{}^\varepsilon_t, \phi \rangle = \langle W_0^0, \phi \rangle \forall \phi \in \mathcal{A}, \ t \in [0,T],
\end{equation}
if and only if 
\begin{equation}
\lim_{\varepsilon \to 0} \langle \rho_1^\varepsilon(t), \phi \rangle = \langle W_0^0, \phi \rangle \forall \phi \in \mathcal{A}, \ t \in [0,T].
\end{equation}
In other words, under the assumptions of theorem 2.3, one has to propagate solve equation (2.8) in time, in order to obtain $\rho_1(0)$ and then can take the limit $\epsilon \to 0$; but this no longer makes sense if the semiclassical limit is taken before time propagation. That is, unlike what happens with regular potentials, the semiclassical limit and the propagation in time by (a regularization of) the Liouville flow no longer commute.

\[ \rho_1(0) \xrightarrow{\epsilon 	o 0} \rho_1(t) \]

\[ W^0(0) \xrightarrow{\epsilon 	o 0} W^0(t) \]

\textbf{Figure 1.} In some cases, there is not enough information in $W^0(0)$ to determine $W^0(t)$.

\textbf{Organization of the paper:} Proofs of the main Theorems are in section 4. Checking the assumptions of the Theorems for various concrete problems is done in section 5.

3. Definitions and Notations

The Fourier transform is defined as

\begin{equation}
\hat{f}(k) = \mathcal{F}_{x \to k} [f(x)] = \int_{x \in \mathbb{R}^n} e^{-2\pi i k x} f(x) dx.
\end{equation}

For compactness, we will use the following notations:

\[ T^V_\epsilon W = \frac{2}{\epsilon} \text{Re} \left[ i \int e^{2\pi i S x} \hat{V}(S) W(x, k - \frac{\epsilon S}{2}) dS \right] =
\]

\[ = 2\mathcal{F}_{X,K \to x,k}^{-1} \left[ \int \hat{V}(S) \hat{W}(X - S, K) \frac{\sin(\pi \epsilon SK)}{\epsilon} dS \right],
\]

\[ T^V_0 W = -\frac{1}{2\pi} \partial_x \hat{V} \cdot \partial_x W =
\]

\[ = 2\pi \mathcal{F}_{X,K \to x,k}^{-1} \left[ \int \hat{V}(S) \hat{W}(X - S, K) S \cdot K dS \right].
\]

The Sobolev norms of order $m$ on phase-space will be defined as follows:

\[ ||f||_{W^m,p(\mathbb{R}^{2n})} = \sum_{|a| + |b| \leq m} ||\partial_x^a \partial_k^b f||_{L^p(\mathbb{R}^{2n})},
\]

where of course $a$ and $b$ are multi-indices of length $n$ each. Moreover, $H^m(\mathbb{R}^{2n}) = W^{m,2}(\mathbb{R}^{2n})$. 
Denote by $Φ$ the $ε$-dependent smoothing operator
\[ Φ : f(x) \mapsto \mathcal{F}_{a(x)}^{-1}[e^{-\frac{γ}{ε}σ^2a^2}\mathcal{F}_{z-a}[f(z)]]. \]
As a matter of notation we will use $\tilde{f} = Φf$. (The parameter $γ$ can be chosen freely, as long as $γ > \frac{2}{1+γ}$, in the context of the problem (1.5) with a potential satisfying Assumption 1. To see why the scaling is calibrated this way, one should consult the proof of lemma 4.1).

Denote $H(x, k, t)$ the Husimi function,
\[
H^ε(x, k) = \frac{\sqrt{2}}{ε^n} \int e^{-2π\frac{(x-x')^2}{ε^2} - 2π\frac{(k-k')^2}{ε^4}} W^ε(x', k', t)dx'dk' = (2ε)^{-\frac{n}{2}} \left( D^ε e^{-\frac{2πi}{ε}(x-x')} e^{-\frac{π}{ε}(x-x')^2} e^{\frac{2πi}{ε}(k-k')} e^{-\frac{π}{ε}(k-k')^2} \right)_{L^2},
\]
where of course $W^ε$ is the corresponding Wigner function, and $D^ε$ the corresponding operator,

\[
D^εu(x) = ε^n \int e^{2πi(x-y)W^ε(x + y/2, εk)}u(y)dydk.
\]
It is of crucial importance that $D^ε > 0 \Rightarrow H(x, k) > 0$.

4. Proofs

4.1. Proof of Theorem 2.1. Observe that $λ$ as used in the statement of theorem 2.1 is an effective upper bound; if we replace $λ$ by $λ' \leq λ$, $\frac{1}{γ'} = o(1)$, the theorem still holds.

Remark on notation: Here and in the sequel we use the notation $R = λ^{-\frac{1}{γ'}}$.

In lemma 5.3 it is shown essentially that Theorem 2.3 applies if $n = 1$, $V(x) = -|x|^{1+θ}$ in $\{|x| < 1\}$ (with a smooth cutoff outside of that, the specifics of which are irrelevant) and $W^*_0 = δ_x^{-1}δ_k^{-1}w(x, k)$, where $δ_k = R^{-\frac{1+θ}{2}}$, $δ_x = R^{-\frac{1+θ}{2}}$. (We avoid duplicating parts of the proof of lemma 5.3 here. Moreover, it is completely straightforward to see that the assumptions not checked explicitly there also hold). Therefore corollary 2.4 applies, i.e. it suffices to find the limit in $ε$ of $ρ^ε_1$.

With $[W^ε_0]$ defined as in equation (5.15), we have
\[
||W^ε_0 - [W^ε_0]| |_{L^1} \leq Cδ_x^{-1}δ_k^{-1} \int_{|x|<CR} |w(\frac{x}{δ_x}, \frac{k}{δ_k})|dxdk = \int_{|x|<CR} |w(x, k)|dxdk \leq Cδ_x^{-1}R = o(1).
\]
This means that $ρ^ε_2$ and $ρ^ε_1$ (the evolution under equation (2.8) of $[W^ε_0]$, $W^ε_0$ respectively) are interchangeable for our purposes, since
\[
\langle ρ^ε_1, ϕ \rangle = \langle ρ^ε_2, ϕ \rangle + \langle ρ^ε_1 - ρ^ε_2, ϕ \rangle
\]
and
\[ |\langle \rho_1^0 - \rho_2^0, \phi \rangle| \leq \|\rho_1^0 - \rho_2^0\|_{L^1} \|\phi\|_{L^\infty} \leq \|W_0^0 - [W_0^0]\|_{L^1} \|\phi\|_{\mathcal{A}}. \]

(We used the obvious bound \( \|\phi\|_{L^\infty} \leq \|\phi\|_{\mathcal{A}} = \int \sup_x |\mathcal{F}_{K \rightarrow K}[\phi(x,k)]| dK \)).

In particular, this means we can always work with \(|x| > CR\), since \(\rho_2^0\) stays outside of \(|x| < CR\) by construction. This allows certain ODEs we will use to be well-posed. Without loss of generality we will work for \(x > 0\), \(k > 0\) (working with \(k < 0\) makes no difference other than the opposite sign in the explicit expression for \(K(0)\) in equation (4.2), the case \(x < 0\) follows by symmetry). Denote by \(X(t), K(t)\) the unique solution of
\[
\dot{X}(t) = 2\pi K(t), \quad \dot{K}(t) = \frac{1+\theta}{2\pi} (X(t))^\theta,
\]
(4.2)
\[ X(0) = CR, \quad K(0) = |X(0)|^{\frac{1+\theta}{2}}, \]
i.e. the only branch of the level set of \(\{2\pi^2 k^2 + V(x)\} = 0\) in \(\{x > R \land k > 0\}\).

The interesting property of this trajectory is that, unlike what happens with regular potentials, it leaves zero in finite time (more generally, i.e. without restricting to \(\{x > R \land k > 0\}\), it reaches and leaves zero in finite time).

Claim: We will show that any characteristic of the Liouville equation (2.8) starting in \(\{x > R \land k > 0\} \cap \text{supp}[W_0^\varepsilon]\) converges to \(X(t), K(t)\) defined in equation (4.2).

Then it follows that
\[ \chi_{x>0}(x,k)[W_0^\varepsilon] \rightarrow \int \chi_{x>0}(x,k)[W_0^\varepsilon] dx dk \delta(x - X(t), k - K(t)). \]

Repeating the argument for \(x < 0\) gives a limit of two delta functions leaving zero along different trajectories.

Proof of the claim: Denote by \(X_1(t), K_1(t)\) the solution of
\[
\dot{X}_1(t) = 2\pi K_1(t), \quad \dot{K}_1(t) = -\frac{1}{2\pi} \partial_x \tilde{V}(X_1(t)),
\]
(4.3)
\[ X_1(0) = x_0, \quad K_1(0) = k_0, \]
where of course \((x_0, k_0) \in \{x > R \land k > 0\} \cap \text{supp}[W_0^\varepsilon]\) are as in the statement of the claim). Then equation (4.3) can be recast as
(4.4)
\[ \dot{X}_1(t) = 2\pi K_1(t), \quad \dot{K}_1(t) = -\frac{1}{2\pi} V_x(X_1(t)) - \frac{1}{2\pi} (\partial_x \tilde{V}(X_1(t)) - V_x(X_1(t))), \]
\[ X_1(0) = x_0, \quad K_1(0) = k_0. \]

Here we use the estimate \( |\partial_x \tilde{V}(X_1(t)) - V_x(X_1(t))| \leq CR^{\theta-1} \varepsilon^{\frac{\theta}{2}} = o(1) \) (see lemma 6.2 with \(V'(x)\) in the place of \(f\)). It is clear now that any such
characteristic converges to
\[
\dot{X}_2(t) = 2\pi K_2(t), \quad \dot{K}_2(t) = -\frac{1}{2\pi} V_x(X_2(t)),
\]
(4.5)
\[X_2(0) = x_0, \quad K_2(0) = k_0.\]
To conclude observe that trajectory \(X_2(t), K_2(t)\) is squeezed between two level sets of \(2\pi^2 k^2 + V(x)\) converging to each other (namely between \(\{2\pi^2 k^2 + V(x)\} = V(CR)\) and \(\{2\pi^2 k^2 + V(x)\} = V(R_1^{-}\)).

Essentially the same analysis applies to the example of lemma 5.4 as well.

4.2. Proof of Theorem 2.2. It will be helpful to recall the main objects we are going to use here. The WT for this problem, \(W = W_\varepsilon(x, k, t)\), satisfies the well-known Wigner equation (1.5) with initial data \(W_0^\varepsilon\). Moreover, \(W_1^\varepsilon\) and \(\rho_2^\varepsilon\) are defined as follows
\[
\left\{
\begin{array}{l}
\partial_t W_1^\varepsilon + 2\pi k \cdot \partial_x W_1^\varepsilon + T_\varepsilon V^\gamma W_1^\varepsilon = 0, \\
W_1^\varepsilon(t = 0) = W_0^\varepsilon,
\end{array}
\right.
\]
(4.6) and \(\rho_2^\varepsilon\),
\[
\left\{
\begin{array}{l}
\partial_t \rho_2^\varepsilon + 2\pi k \partial_x \rho_2^\varepsilon + T_0 V^\gamma \rho_2^\varepsilon = 0, \\
\rho_2^\varepsilon(t = 0) = [W_0^\varepsilon],
\end{array}
\right.
\]
(4.7)
Of course recall that \(\rho_1^\varepsilon\) was defined as the solution of equation (2.8) with initial data \(\rho_1^\varepsilon(t = 0) = W_0^\varepsilon\).

We partition the proof as follows:
\[
||W^\varepsilon - \rho_1^\varepsilon||_{L^2} \leq ||W^\varepsilon - W_1^\varepsilon||_{L^2} +
+||W_1^\varepsilon - \rho_2^\varepsilon||_{L^2} + ||\rho_1^\varepsilon - \rho_2^\varepsilon||_{L^2} = o(1)||W_0^\varepsilon||_{L^2}.
\]
(4.8)
We will use without further comments the elementary observation
\[
\left\|
\int f(x - s, k) g(s) ds \right\|_{L^2(\mathbb{R}^{2n})} \leq ||g||_{L^1(\mathbb{R}^n)} ||f||_{L^2(\mathbb{R}^n)}
\]
(4.9)

Lemma 4.1 (\(W \approx W_1\)). \(\forall t \in [0, T]\)
\[
||W_1^\varepsilon(t) - W^\varepsilon(t)||_{L^2} = O(\varepsilon^{1+\gamma} ||W_0^\varepsilon||_{L^2}).
\]
Proof: Denote
\[
h^\varepsilon = W^\varepsilon - W_1^\varepsilon.
\]
(4.10)
Obviously,
\begin{equation}
\partial_t h + 2\pi k \partial_x h + T^\gamma V h = -T^V - V h = -T^\gamma V h = -\gamma V h - \gamma V W \epsilon_1,
\end{equation}
\begin{equation}
h(x, k, 0) = 0.
\end{equation}

Since the Wigner equation has a bounded \( L^2 \) propagator, it suffices to bound in \( L^2 \) the rhs. Indeed, we have
\begin{equation}
\|T^V V W_1^\epsilon\|_{L^2} = C \|\hat{V}(S)(1 - e^{-\frac{\pi}{2} \gamma S^2})\hat{W}_1^\epsilon(X - S, K)\frac{\sin(\pi \epsilon S K)}{\epsilon} dS\|_{L^2} \leq
\end{equation}
\begin{equation}
\leq O(\epsilon^{-1})\|\hat{V}(S)(1 - e^{-\frac{\pi}{2} \gamma S^2})\|_{L^1}\|\hat{W}_1^\epsilon\|_{L^2} \leq
\end{equation}
\begin{equation}
\leq O(\epsilon^{-1})\|W_1^\epsilon\|_{L^2} \left[ \epsilon^\gamma \int_{|S| \leq \epsilon^{-\frac{\gamma}{2}}} |\hat{V}(S)| |S|^2 dS + \int_{|S| \geq \epsilon^{-\frac{\gamma}{2}}} |\hat{V}(S)| dS \right] =
\end{equation}
\begin{equation}
= O(\epsilon^{-1})\|W_1^\epsilon\|_{L^2} \left[ \epsilon^\gamma + \epsilon^\gamma \int_{\rho \in (1, \epsilon^{-\frac{\gamma}{2}})} \rho^0 \rho \rho^{-\theta} d\rho + \int_{\rho \geq \epsilon^{-\frac{\gamma}{2}}} \rho^{-2-\theta} d\rho \right] =
\end{equation}
\begin{equation}
= O(\epsilon^{-1})\|W_0^\epsilon\|_{L^2} \left[ \epsilon^\gamma + \epsilon^\gamma[-\frac{1}{2}(1-\theta)+1] + \epsilon^{\frac{1+\theta}{2}} \right].
\end{equation}

Asking that the remainder is small gives
\begin{equation}
\gamma > 1,
\end{equation}
\begin{equation}
\gamma[-\frac{1}{2}(1-\theta)+1] > 1 \iff \gamma > \frac{2}{1+\theta},
\end{equation}
\begin{equation}
\gamma^{\frac{1+\theta}{2}} > 1 \iff \gamma > \frac{2}{1+\theta}.
\end{equation}

So finally with the calibration \( \gamma > \frac{2}{1+\theta} \) the proof of lemma 4.1 is complete.

Moreover:

**Lemma 4.2** \( (W_1 \approx \rho_0^\epsilon), \forall t \in [0, T] \)
\begin{equation}
\|W_1^\epsilon(t) - \rho_0^\epsilon(t)\|_{L^2} = O(\epsilon^{\frac{\theta}{2}-\delta})\|W_0^\epsilon\|_{L^2}.
\end{equation}

**Proof:** It is straightforward to check that
\begin{equation}
\partial_t \hat{W}_1^\epsilon - 2\pi X \cdot \partial_K \hat{W}_1^\epsilon + 2 \int \hat{\gamma}(S) \hat{W}_1^\epsilon(X - S, K) \frac{\sin(\pi \epsilon S K)}{\epsilon} dS = 0,
\end{equation}
\begin{equation}
\partial_t \hat{\rho}_2^\epsilon - 2\pi X \cdot \partial_K \hat{\rho}_2^\epsilon + 2 \int \hat{\gamma}(S) \hat{\rho}_2^\epsilon(X - S, K) S \cdot K dS = 0,
\end{equation}
\begin{equation}
\hat{\rho}_2^\epsilon(0) = \rho_0^\epsilon.
\end{equation}
and therefore, if \( f = \hat{W}^{\varepsilon}_1 - \hat{\rho}_2 \),

\[
\partial_t f - 2\pi X \cdot \partial_K f + 2 \int \hat{V} (S) f (X - S, K) \frac{\sin(\pi \varepsilon S \cdot K)}{\varepsilon} dS = 2 \int \hat{V} (S) \hat{\rho}_2 (X - S, K) \left( 1 - \frac{\sin(\pi \varepsilon S \cdot K)}{\pi \varepsilon S \cdot K} \right) \pi S \cdot K dS,
\]

(4.16)

\[
f(x, k, 0) = W^{\varepsilon}_0 - [W^{\varepsilon}_0].
\]

Using the Duhamel formula,

(4.17)

\[
||f(t)||_{L^2} \leq ||W^{\varepsilon}_0 - [W^{\varepsilon}_0]||_{L^2} + 2T \sup_{t \in [0, T]} ||\int \hat{V} (S) \hat{\rho}_2 (X - S, K, t) \left( 1 - \frac{\sin(\pi \varepsilon S \cdot K)}{\pi \varepsilon S \cdot K} \right) \pi S \cdot K dS||_{L^2}.
\]

Recall that the first term above is (relatively) small by assumption; we will work out the other term:

To bound that, first of all observe that

(4.18)

\[
\left| \frac{\sin(\pi \varepsilon S \cdot K)}{\pi \varepsilon S \cdot K} - 1 \right| \leq C \min \{1, |\pi \varepsilon S \cdot K|\}
\]

and therefore for any \( b < 0 \)

(4.19)

\[
||\int \hat{V} (S) \hat{\rho}_2 (X - S, K) \left( 1 - \frac{\sin(\pi \varepsilon S \cdot K)}{\pi \varepsilon S \cdot K} \right) \pi S \cdot K dS||_{L^2} \leq
\]

\[
\leq C \varepsilon \left( \int_{0 < |S| \leq 1} ||\hat{V} (S)| \cdot |S|^2 |\hat{\rho}_2 (X - S, K)| \cdot |K|^2 dS||_{L^2} + \varepsilon^{1+2b} \int_{1 < |S| < \varepsilon^b} ||\hat{V} (S)| \hat{\rho}_2 (X - S, K)|K|^2 dS||_{L^2} + ||\int_{|S| > \varepsilon^b} \hat{V} (S) \hat{\rho}_2 (X - S, K)| SK| dS||_{L^2} \right) \leq
\]

\[
\leq C \left( \varepsilon^{1+2b} ||\hat{V} (S) \min \{|S|^2, 1\} ||_{L^1} ||\hat{\rho}_2||_{H^2} + \int_{|S| > \varepsilon^b} ||\hat{V} (S)| |S| dS||\hat{\rho}_2||_{H^1} \right) \leq
\]

\[
\leq C \left( \varepsilon^{1+2b} ||\hat{V} \frac{S^2}{1 + S^2} ||_{L^1} ||\hat{\rho}_2||_{H^2} + \int_{\rho = \varepsilon^b} \rho^{-1-\theta} d\rho ||\hat{\rho}_2||_{H^1} \right) =
\]

\[
= C \left( \varepsilon^{1+2b} ||\hat{\rho}_2||_{H^2} + \varepsilon^{-\theta b} ||\hat{\rho}_2||_{H^1} \leq C (\varepsilon^{1+2b} + \varepsilon^{-\theta b-\delta}) ||W^{\varepsilon}_0||_{L^2} \right)
\]
Calibrating the parameters is easy; we are given \( \theta \in (0, 1) \), \( \delta \) as in the statement of the theorem and we need to find \( b \in (-\frac{1}{2}, 0) \) so that

\[
1 + 2b - \delta > 0, \quad -\theta b - \delta > 0 \iff \delta - \frac{1}{2} < b < -\frac{\delta}{\theta}.
\]

It should now be clear that the constraint on \( \delta \) in statement of the theorem comes from the self-consistency check \(-\frac{1}{2} < \delta - \frac{1}{2} < -\frac{\delta}{\theta} < 0\). The proof of lemma 4.2 is complete.

Now the last step is to show that \( \rho^2_\epsilon \approx \rho_1^\epsilon \); but this follows by construction, since they satisfy the same equation (which has an \( L^2 \)-continuous propagator):

\[
||\rho^2_\epsilon(t) - \rho_1^\epsilon(t)||_{L^2} = ||W^\epsilon_0 - [W^\epsilon_0]||_{L^2}.
\]

The proof of Theorem 2.2 is complete.

4.3. Proof of Theorem 2.3. First of all observe that equation (2.10), under the assumptions of theorem 2.3 ||\( W^\epsilon_0 - [W^\epsilon_0] ||_{L^2} = o(1)\) and ||\( W^\epsilon_0 ||_{L^2} = o(\epsilon^{-\kappa})\), implies

\[
||W^\epsilon(t) - \rho_1^\epsilon(t)||_{L^2} = o(1).
\]

Therefore the result formally follows by

\[
\langle \rho_1^\epsilon(t), \phi \rangle = \langle W^\epsilon(t), \phi \rangle + \langle \rho_1^\epsilon(t) - W^\epsilon, \phi \rangle = \langle W^\epsilon(t), \phi \rangle + o(1).
\]

However for that computation to be correct we would need to have \( ||\phi||_{L^2} \leq ||\phi||_{A^\epsilon} \), which is not true. This section is mainly devoted to working around that. The Husimi function suppresses the spurious oscillations typically present in \( W^\epsilon \), and therefore is a useful tool in passing from \( L^2 \) asymptotics to \( L^1 \) asymptotics (see also [5] for a related construction).

**Theorem 4.3.** Suppose that the assumptions of Theorem 2.3 hold. Denote by \( W^\epsilon(x, k, t) \) the solution of the Wigner equation (1.5), and by \( H^\epsilon(x, k, t) \) the corresponding Husimi function.

Moreover, assume that there is a function \( \rho^\epsilon(x, k, t) \) such that

\[
\exists C_0 > 0 : \bigcup_{t \in [0, T]} \text{supp} \rho^\epsilon(x, k, t) \subseteq \{||(x, k)|| < C_0\}
\]

and

\[
\sup_{t \in [0, T]} ||H^\epsilon(t) - \rho^\epsilon(t)||_{L^2} = o(1).
\]

Then

\[
\sup_{t \in [0, T]} ||H^\epsilon(t) - \rho^\epsilon(t)||_{L^1} = o(1).
\]

**Remark:** The proof is essentially an adaptation of an argument from [5].
Proof: Before the main part of the proof observe that our assumptions on the potential imply that \( W^{\varepsilon}(x, k, t) \in L^1_{x, k} \forall t \), and

\[
\text{tr}(A^{\varepsilon}(t)) = \int W^{\varepsilon}(x, k, t) dx dk = 1.
\]

The first statement is elementary to prove, and the second well-known; see [5] for details. It is also straightforward that \( ||W^{\varepsilon}(t)||_{L^2} = ||W_0^{\varepsilon}||_{L^2} \) (it follows by the essential self-adjointness of the corresponding Schrödinger operator).

Denote \( h^{\varepsilon} = H^{\varepsilon} - \rho^{\varepsilon} \). Then, for any \( M > C_0 \),

\[
||h^{\varepsilon}||_{L^1} = \int |H^{\varepsilon} - \rho^{\varepsilon}| dx dk + \int |H^{\varepsilon} - \rho^{\varepsilon}| dx dk \leq
\]

\[
\leq \langle x_0, M \rangle \left( \left( ||(x, k)|, h^{\varepsilon}\right) + \int \left| H^{\varepsilon}\right| dx dk =
\right)
\]

\[
= CM^n ||h^{\varepsilon}||_{L^2} + \int \left| H^{\varepsilon}\right| dx dk.
\]

\( M = M(\varepsilon) \) will be calibrated so that \( \frac{1}{M} = o(1) \) and

\[
\sup_{t \in [0, T]} M^n(\varepsilon)||h^{\varepsilon}||_{L^2} = o(1).
\]

At this point we will need that \( V(x) > 0 \). However, since \( V \) is assumed to be bounded below, non-negativity follows without loss of generality by substituting \( V(x) \rightarrow V_2(x) = V(x) + ||V_-||_{L^\infty} \), where of course \( V_- \) is the negative part, \( V_-(x) = V(x)\chi_{\{V(x) < 0\}}(x) \).

Now, recalling that \( H(x, k) \geq 0 \), we have

\[
\int |H(x, k)| dx dk \leq M^{-2} \int |H(x, k)| dx dk \leq
\]

\[
\leq M^{-2} \int (k^2 + V(x))H(x, k) dx dk =
\]

\[
= M^{-2} \int (k^2 + V(x))W^{\varepsilon}(x, k, t) dx dk + M^{-2}C_\varepsilon \int W^{\varepsilon}(x, k) dx dk +
\]

\[
+ M^{-2} \int (V(x) - \left( \frac{\varepsilon}{2} \right)^2 \int e^{-\frac{2\pi |x-x'|^2}{\varepsilon}} V(x') dx') W^{\varepsilon}(x, k) dx dk \leq
\]

\[
\leq O(\varepsilon M^{-2}) +
\]

\[
+ M^{-2} ||(V(x) - \left( \frac{\varepsilon}{2} \right)^2 \int e^{-\frac{2\pi |x-x'|^2}{\varepsilon}} V(x') dx')||_{L^2} ||W^{\varepsilon}||_{L^2} \leq
\]

\[
\leq O(M^{-2}) (\varepsilon + \sqrt{\varepsilon}||V||_{H^1} ||W^{\varepsilon}||_{L^2}) = o(1).
\]

In this computation we used the energy conservation,

\[
\int (2\pi^2 k^2 + V(x))W^{\varepsilon}(x, k, t) dx dk = \text{tr} \left( (-\frac{\varepsilon}{2} \Delta + V(x))D^{\varepsilon}(t) \right) =
\]

\[
= \text{tr} \left( (-\frac{\varepsilon}{2} \Delta + V(x))D^{\varepsilon}(0) \right) = \int (2\pi^2 k^2 + V(x))W^{\varepsilon}(x, k, 0) dx dk \leq C_\varepsilon,
\]
and the obvious observation

$$\|V(x) - \left(\frac{2}{\varepsilon}\right)^n \int e^{-\frac{2\pi |x-x'|^2}{\varepsilon}} V(x') dx'\|_{L^2} = O(\sqrt{\varepsilon}\|V\|_{H^1})$$

(Also $\sqrt{\varepsilon}\|W_0^\varepsilon\|_{L^2} \leq \sqrt{\varepsilon}\|W_0^\varepsilon\|_{H^2} = o(1)$ virtue of equation (2.9)). The proof is complete.

**Conclusion of the proof of Theorem 2.3:** Now observe that the conditions of theorem 4.3 are satisfied with $\rho^\varepsilon = \rho_1^\varepsilon$.

Indeed, since $W_0^\varepsilon$ is of compact support, $\rho_1^\varepsilon(t)$ is of compact support uniformly in $\varepsilon$ (see lemma 6.5) for $t \in [0,T]$. Moreover, using again that $\|H^\varepsilon(t) - W^\varepsilon(t)\|_{L^2} \leq C\sqrt{\varepsilon}\|W^\varepsilon(t)\|_{H^1}$ (see e.g. lemma A.1 of [4] for a proof), it follows that

$$\|H^\varepsilon - \rho_1^\varepsilon\|_{L^2} \leq C\|W^\varepsilon - \rho_1^\varepsilon\|_{L^2} + \sqrt{\varepsilon}\|W^\varepsilon(t)\|_{H^1}.$$  

With that in mind we can calibrate more precisely the parameter $M = M(\varepsilon)$ which appeared earlier:

$$M = (\varepsilon^n\|W_0^\varepsilon\|_{L^2} + \|W_0^\varepsilon - [W_0^\varepsilon]\|_{L^2} + \sqrt{\varepsilon}\|W_0^\varepsilon\|_{H^1})^{-\frac{1}{2n}}.$$  

(Again we have $\sqrt{\varepsilon}\|W_0^\varepsilon\|_{H^1} \leq \sqrt{\varepsilon}\|W_0^\varepsilon\|_{H^2} = o(1)$ follows from equation (2.9)).

The result follows, making use of theorem 4.3, by virtue of

$$|\langle \rho_1^\varepsilon - W^\varepsilon, \phi \rangle| \leq |\langle \rho_1^\varepsilon - H^\varepsilon, \phi \rangle| + |\langle W^\varepsilon - H^\varepsilon, \phi \rangle| \leq$$

$$\leq \|\rho_1^\varepsilon - H^\varepsilon\|_{L^1} \|\phi\|_{L^\infty} + |\langle W^\varepsilon, \phi - \phi + \phi * \frac{\sqrt{2} e^{-\frac{2\pi |(x,k)|^2}{\varepsilon}}}{e^\frac{2\pi |(x,k)|^2}{\varepsilon}} \rangle| =$$

$$= O(M^{-2}) (1 + \sqrt{\varepsilon}\|W^\varepsilon\|_{L^2}) \|\phi\|_{A} + |\langle W^\varepsilon, \left(1 - e^{-\frac{2\pi \Delta}{\varepsilon}}\right) \phi \rangle|.$$  

It was shown in [13] that $W^\varepsilon$ and $H^\varepsilon$ have the same weak-* limit in $\mathcal{A}'$; $\langle W^\varepsilon, \left(1 - e^{-\frac{2\pi \Delta}{\varepsilon}}\right) \phi \rangle \to 0$ follows. One can estimate the rate of convergence if there is more smoothness information for the test-function $\phi$.

The proof is complete.

**5. Examples**

In this section we isolate a subset of the assumptions of Theorems 2.2, 2.3 and construct some non-trivial examples for which they are satisfied. Essentially these are “graceful degradation” properties for the flow. The remaining assumptions, which we don’t check explicitly, are much simpler to check (regularity of the potential and initial data, finite energy etc).
5.1. **Examples for Theorem 2.2.** Here we are concerned with Cauchy problems for the Liouville equation,

\[
\begin{aligned}
\partial_t \rho^\varepsilon_1 + 2\pi k \partial_x \rho^\varepsilon_1 - \frac{1}{2\pi} \partial_x \nabla \cdot \partial_x \rho^\varepsilon_1 &= 0, \\
\rho^\varepsilon_1(t=0) &= W^\varepsilon_0,
\end{aligned}
\]

which have the following

**Property 1.** There exist \( T > 0, \delta \in (0, \frac{\theta}{2+\theta}) \) and \([W^\varepsilon_0] \in H^2(\mathbb{R}^{2n})\) such that

\[(5.1) \quad \|[W^\varepsilon_0] - W^\varepsilon_0\|_{L^2} = o(1)\|W^\varepsilon_0\|_{L^2}\]

and the solution of

\[
\begin{aligned}
\partial_t \rho^\varepsilon_2 + 2\pi k \partial_x \rho^\varepsilon_2 - \frac{1}{2\pi} \partial_x \nabla \cdot \partial_x \rho^\varepsilon_2 &= 0, \\
\rho^\varepsilon_2(t=0) &= [W^\varepsilon_0],
\end{aligned}
\]

satisfies

\[(5.3) \quad \|\rho^\varepsilon_2(t)\|_{H^2} = O(\varepsilon^{-\delta}\|W^\varepsilon_0\|_{L^2}), \quad \text{uniformly for } t \in [0,T].\]

We treat \( \kappa, \gamma \) as given parameters; of course when checking the property in the context of Theorem 2.3 they are controlled by the statement of Theorem 2.2.

**Lemma 5.1.** Let \( \theta \in (0,1), \delta \in (0, \frac{\theta}{2+\theta}), n = 2, V(x) = -|x|^{1+\theta} \) with an appropriate smooth cutoff outside \( \{|x| < 2\} \). Assume moreover that there exists a function of compact support \( w \in H^2(\mathbb{R}^{4}) \cap L^\infty \), such that

\[(5.4) \quad W^\varepsilon_0(x,k) = \delta_x^{-2}\delta_k^{-2} w\left(\frac{x - (z_1, z_2)}{\delta_x}, \frac{k - (z_3, z_4)}{\delta_k}\right),\]

with \((z_1, z_2) = -L(z_3, z_4)\) (i.e. the wavepacket is “shot towards zero”). \( ||w||_{H^2}, ||w||_{L^\infty}, |z| \) have to be bounded uniformly in \( \varepsilon, ||w||_{H^2}, ||w||_{L^\infty}, |z| = O(1) \); other than that \( w \) and \( z \) can be allowed to depend on \( \varepsilon \).

The small parameters involved \( \delta_x, \delta_k, R = o(1) \) are calibrated as follows:

\[R = \left( \log \left( \frac{1}{\varepsilon} \right) \right)^{-\frac{1}{4}},\]

as in equation (6.26), and

\[\delta_x^{-1}(R + \delta_k) = o(1),\]

\[(5.5) \quad (\delta_x + \delta_k)^{-2}(\delta_k + R)^{-2} = O(\varepsilon^{-\frac{2}{7}}).\]

Then Property 1 is satisfied.

**Remarks:**
• In fact it follows from the lemma that property 1 holds for any \( \delta \in (0, \frac{\theta}{2}) \).

• The difficulties arise from the non-smootheness at zero, and the details of a smooth behaviour away from zero are irrelevant here. So we will work for a small enough time, before any trajectory starting in a neighbourhood of zero reaches the support of the cutoff function. Therefore, to keep the presentation simple, we will not introduce any explicit treatment of the cutoff function. This approach is followed in the sequel as well.

• It is obviously much easier to work if the wavepacket is not “shot towards zero”. However in that case it’s basically covered by the standard theory; the interesting cases are when in the limit a delta function passes over a set where the flow associated with the original potential \( V(x) \) doesn’t exist.

• An example of a scaling satisfying the above constraints is \( \delta_x = \sqrt{\delta_k} = \sqrt{R} \).

\[ \text{Proof:} \quad \text{Assume without loss of generality that } z = (0, -L, 0, 1). \text{ Consider } \phi : \mathbb{R} \to [0,1] \text{ to be a } C^\infty \text{ function such that} \]

\[ \phi(x) = 0, \quad |x| < \frac{1}{2}, \]

\[ \phi(x) = 1, \quad |x| > 1. \]

The modified initial data \( [W_0^\varepsilon] \) is selected as

\[ [W_0^\varepsilon](x, k) = W_0^\varepsilon(x, k)\phi \left( \frac{|x_1|}{R + 2L\delta_k} \right). \]

The key claim is that \( \rho_2^\varepsilon \) never touches \( \{|x| < R\} \). Indeed \( [W_0^\varepsilon] \) has two disjointly supported components. The calibration of the parameters and the cut-off is such that, if \( [W_0^\varepsilon] \) was propagated by a Liouville equation corresponding to \( V(x) = 0 \), each of the components would stay on either side of the strip \( \{|x_1| < R\} \) (see Figure 2). The idea is that since there is a potential driving the flow away from \( x = 0 \), the claim follows. (In fact it is slightly more complicated for this case, since the potential depends on \( x_2 \) as well: one easily gets an \( O(1) \) upper bound for the time it takes for the trajectories to reach \( x_2 = 0 \), and then checks that in this longer time, \( x_1 \) doesn’t have the time to change sign. The difference in an \( O(1) \) factor which we absorb in \( C \) below).

Therefore, in this problem we can take \( \tilde{V}_{\text{eff}}(x) = \tilde{V}(x)\phi \left( \frac{|x|}{R + C\delta_k} \right) \). Now for condition (5.3), making use of observation 6.6 (enlarging the cut-off area in fact strengthens observation 6.6, i.e. substituting \( R + 2L\delta_k \) for \( R \) is painless) and lemma 6.4 it follows that

\[ \|\rho_2^\varepsilon(t)\|_{H^2} \leq \|W_0^\varepsilon\|_{H^2} \varepsilon^{-\frac{1}{2}} \leq C\|W_0^\varepsilon\|_{H^2} \varepsilon^{-\frac{1}{2}}(\delta_k + R)^{-2} = \]

\[ = C\|W_0^\varepsilon\|_{L^2}(\delta_k + R)^{-2}(\delta_k + R)^{-2} \varepsilon^{-\frac{1}{2}}, \]
Figure 2. Any trajectory leaving the support of $W_0^\varepsilon$ has an initial velocity in a small cone around $k = (0, 1)$. If we were in free space, there wouldn’t be enough space for sufficient movement in $x_1$ to reach $\{|x_1| < R\}$. Therefore the solution is the sum of two components, supported on either side of $\{|x_1| < R\} = 0$. The presence of a repulsive potential in many cases can be easily factored in this construction.

and therefore – using the scaling of equation (5.5) – $\|\rho_2^\varepsilon(t)\|_{H^2} = O(\varepsilon^{-\delta}\|W_0^\varepsilon\|_{L^2})$.

Now for $W_0^\varepsilon - [W_0^\varepsilon]$:

$$\|W_0^\varepsilon - [W_0^\varepsilon]\|_{L^2}^2 \leq \delta_x^{-2n}\delta_k^{-2n} \int_{|x_1|<C(R+\delta_k)} |w(\frac{(x_1,x_2)}{\delta_x}, \frac{(k_1-1,k_2)}{\delta_k})|^2 dx dk =$$

$$= \delta_x^{-n}\delta_k^{-n} \int_{|x_1|<C^R+\delta_k} |w(x_1,x_2,k_1-1,k_2)|^2 dx dk \leq$$

$$\leq C\|W_0^\varepsilon\|_{L^2}^2 \|w\|_{L^\infty}^2 \delta_x^{-1}(R + \delta_k)$$

The same idea can be applied to different configurations:

Lemma 5.2. Consider the setup of lemma 5.1 with $z = (0, -2, 0, 1)$ and the only difference that

$$V(x) = -|x_1|^{1+\theta}\psi(x_1)\psi(x_2).$$

where $\psi$ is a smooth cutoff function, $\psi = 1 - \phi \in \mathcal{S}(\mathbb{R})$ ($\phi$ was defined in equation (5.6)).

Then Property 1 holds.
Proof: The geometry is essentially the same as before, and it is clear that the obvious adaptation of observation 6.6, i.e. the one with

$$\tilde{V}_{\text{eff}}(x) = \tilde{V}(x)\phi\left(\frac{x_1}{R + C\delta_k}\right)$$

holds.

Figure 3. The construction of lemma 5.1 holds, with obvious adjustments, to the example of lemma 5.2.

5.2. Examples for Theorem 2.3. In relation to Theorem 2.3, we introduce

Property 2. Given $V(x)$, $W_0^\varepsilon$ with $||W_0^\varepsilon||_{L^2} = o(\varepsilon^{-\kappa})$, there exist $T > 0$, $\delta \in (0, \frac{\theta}{\theta + \delta})$ and $[W_0^\varepsilon] \in H^2(\mathbb{R}^{2n})$ such that

$$||[W_0^\varepsilon] - W_0^\varepsilon||_{L^2} = o(1) \quad (5.11)$$

and the solution of

$$\left\{ \begin{array}{l}
\partial_t \rho_2^\varepsilon + 2\pi k \partial_x \rho_2^\varepsilon - \frac{1}{2\pi} \partial_x \tilde{V} \cdot \partial_x \rho_2^\varepsilon = 0, \\
\rho_2^\varepsilon(t = 0) = [W_0^\varepsilon],
\end{array} \right. \quad (5.12)$$

satisfies

$$||\rho_2^\varepsilon(t)||_{H^2} = O(\varepsilon^{-\delta}||W_0^\varepsilon||_{L^2}), \quad \text{uniformly for } t \in [0, T]. \quad (5.13)$$
Lemma 5.3. Let \( n = 1, \theta \in (0,1), V(x) = -|x|^{1+\theta} \) with an appropriate smooth cutoff outside \( \{|x| < 2\} \). Assume moreover that \( w(x,k) \in H^2 \cap L^\infty \cap L^1, \text{supp } w \subseteq \{|x|,|k| < 1\} \), and

\[
W_0^\varepsilon = \delta_x^{-1} \delta_k^{-1} w(\frac{x}{\delta_x}, \frac{k}{\delta_k}).
\]

Finally, set \( R = (\log \frac{1}{\varepsilon})^{-\frac{1}{\gamma}} \).

Then Property 2 is satisfied if \( \delta_k = C'R^{1+\theta}, \delta_x = C''R^{1-\theta} \).

Proof: We will cutoff a strip of the form \( \{|x| \leq CR\} \), and show that in fact this suffices.

As is illustrated in Figure 5.2, the preimage under the flow of \( \{|x| < R \wedge |k| < \delta_k\} \) is contained between the level sets \( \{2\pi^2 k^2 - \hat{V}(x) = \hat{V}(0)\} \) and \( \{2\pi^2 k^2 + \hat{V}(x) = \hat{V}(R)\} \) (in red).

A big help with the algebra will be to approximate \( \hat{V}(x) = -|x|^{1+\theta} + O(\varepsilon^{\frac{\gamma}{2}}) \). (\( \gamma \) is as in the statement of Theorem 2.2). Since all the other small parameters are (negative) powers of \( \log \frac{1}{\varepsilon} \), \( O(\varepsilon^{\frac{\gamma}{2}}) \) is negligible everywhere. (The justification is that without loss of generality we can localize the problem on a compact set, and then \( V(x) \in W^{1,\infty} \) uniformly in \( \varepsilon \). The conclusion follows by a standard observation on mollifiers, see lemma 6.2).

It is easily seen that the needed length is \( x_s \), defined by

\[
2\pi^2 \delta_k^2 + \hat{V}(x_s) = \hat{V}(R)
\]
or

\[
|x_s|^{1+\theta} = 2\pi^2 \delta_k^2 + R^{1+\theta} + O(\varepsilon^{\frac{\gamma}{2}}).
\]

Here is where \( \delta_k = C'R^{1+\theta} \) comes from; making that scaling we get \( x_s = ((1 + 2\pi^2)R^{1+\theta} + O(\varepsilon^{\frac{\gamma}{2}}))^{\frac{1}{1+\theta}} = C'R \).

So far we have ensured by construction that, if \( \delta_k = R^{1+\theta} \), (there is an \( O(1) \) constant \( C \) so that) if

\[
[W_0^\varepsilon](x,k) = W_0^\varepsilon(x,k) \phi\left(\frac{x}{CR}\right),
\]

its propagation under equation (5.12) will never enter \( \{|x| < R\} \).

Now for the approximation error:

\[
||W_0^\varepsilon - [W_0^\varepsilon]|^2_{L^2} \leq C\delta_x^{-2} \delta_k^{-2} \int_{|x|<CR} |w(\frac{x}{\delta_x}, \frac{k}{\delta_k})|^2 dx dk =
\]

\[
= C\delta_x^{-1} \delta_k^{-1} \int_{|x|<CR} |w(x,k)|^2 dx dk \leq C\delta_x^{-2} R^{\frac{1-\theta}{\gamma}}.
\]

This gives the constraint \( R^{\frac{1-\theta}{\gamma}} = o(\delta_x) \) for the error to be small, which is satisfied e.g. by our earlier selection \( \delta_x = CR^{\frac{1-\theta}{\gamma}} \).
Again, since all the small parameters are powers of \( \log(\frac{1}{\varepsilon}) \), \( ||W_0^\varepsilon||_{H^2} = o(\varepsilon^{-\frac{5}{2}}) \) follows automatically for any \( \delta \in (0, \frac{\theta}{2\pi}) \).

The proof is complete.

The following is a somewhat artificial example, but it highlights an interesting behaviour.

**Lemma 5.4.** Assume \( n = 2 \), \( V \) as in equation (5.10). Moreover, there is a function \( f \in \mathcal{S}(\mathbb{R}^4) \), \( \int f(x,k)dxdk = 1 \), such that

\[
\text{supp } f \subseteq \{|x| < 1 \} \times \{|k| < 1 \land k_1 > 0 \}.
\]

Denote by \( \tilde{f} \)

\[
\tilde{f}(x_1, x_2, k_1, k_2) = f(-x_1, x_2, -k_1, k_2)
\]

Now take

\[
W_0^\varepsilon(x, k) = \frac{1}{\delta_x \delta_k} \left(c_1 f \left( \frac{x_1 - 2R + x_2 + 2}{\delta_x}, \frac{k_1 + k_2 + 1}{\delta_k} \right) +
\right.
\]

\[
+ c_2 \tilde{f} \left( \frac{x_1 - 2R + x_2 + 2}{\delta_x}, \frac{k_1 + k_2 + 1}{\delta_k} \right).
\]

If \( \delta_k = C' \delta_x = C'' \sqrt{R} \), this problem has Property 2.
Proof: Two things should be obvious by construction: that $W^\varepsilon_0 \to \delta(x + (0, 2), k - (0, 1))$, and that the propagation of $W^\varepsilon_0$ under equation (2.8), i.e. $\rho^\varepsilon_1$, is never supported inside $\{|x_1| < R\}$. In other words, for this specifically constructed data, $W^\varepsilon_0 = [W^\varepsilon_0]$. The scaling of $\delta_x, \delta_k$ in this case controls simply the rate of concentration, and therefore $\|W^\varepsilon_0\|_{H^2}$. The result follows.

6. Auxiliary results

Lemma 6.1 (2nd order derivatives equations for the Liouville equation). Consider the Cauchy problem for the Liouville equation with potential $V(x)$ on $\mathbb{R}^n$,

$$ f_t + 2\pi k \nabla_x f - \frac{1}{2\pi} \nabla_x V(x) \nabla_k f = 0, $$

(6.1)

$$ f(t = 0) = f_0. $$

There are constants $C_1, C_2 > 0$ depending only on $n$ such that

$$ \|f(t)\|_{H^2} \leq C_1 e^{tC_2 \sup_{|\alpha| \leq 3} \|\partial_x V(x)\|_{L^\infty}} \|f_0\|_{H^2}. $$

Proof: The proof follows readily with the method of characteristics. It suffices to observe that if

$$ \dot{x}_i(t) = 2\pi k_i(t), \quad \dot{k}_i(t) = -\frac{1}{2\pi} \partial_{x_i} V(x(t)), $$

(6.2)

$$ z(t) = f(x(t), k(t)), \quad \dot{z}_j(t) = 0, $$

(6.3)

$$ z_{x_i}(t) = \partial_{x_i} f(x(t), k(t)), \quad z_{x_i x_j}(t) = \partial_{x_i x_j} f(x(t), k(t)), \quad z_{k_i}(t) = \partial_{k_i} f(x(t), k(t)), \quad z_{k_i k_j}(t) = \partial_{k_i k_j} f(x(t), k(t)), \quad z_{x_i k_j}(t) = \partial_{x_i k_j} f(x(t), k(t)), $$

it follows that

$$ \dot{z}_{x_i}(t) = \frac{1}{2\pi} \sum_m \partial_{x_i x_m} V(x(t)) z_{x_i k_m}(t), $$

(6.4)

$$ \dot{z}_{k_i}(t) = -2\pi z_{x_i}(t), $$

(6.5)

$$ \dot{z}_{k_i k_j}(t) = -2\pi (z_{x_i, x_j} + z_{x_i k_j}), $$

(6.6)

$$ \dot{z}_{k_i x_j}(t) = -2\pi z_{x_i x_j}(t) + \frac{1}{2\pi} \sum_m \partial_{x_m x_j} V(x(t)) z_{x_i k_m}(t). $$

(6.7)

The result now follows by the Gronwall inequality.
Lemma 6.2. Consider a function \( f \in W^{1,\infty} \), i.e. \( \sup_{|a| \leq 1} |a| \leq 1 \)
\( ||\partial_a f||_{L^\infty} < \infty \).

Then, if \( \bar{f} = \left( \frac{2}{\eta} \right)^{\frac{n}{2}} \int e^{-2\pi \frac{|x-x'|^2}{\eta}} f(x')dx' \), we have
\[
||f - \bar{f}||_{L^\infty} = O(||f||_{W^{1,\infty}} \sqrt{\eta}).
\]

Proof: Take any \( \zeta \in (0, \frac{1}{2}) \). Now we have
\[
|f(x) - \bar{f}(x)| = \left| \left( \frac{2}{\eta} \right)^{\frac{n}{2}} \int e^{-2\pi \frac{|x-x'|^2}{\eta}} [f(x - x') - f(x)]dx' \right| + O(\eta^\infty) \leq \]
\[
\leq C||f||_{W^{1,\infty}} \left( \frac{2}{\eta} \right)^{\frac{n}{2}} \int_{|x'| < 2^{\zeta - \frac{1}{2}}} e^{-2\pi \frac{|x'|^2}{\eta}} |x'|dx' + O(||f||_{L^\infty} \eta^\infty) = \]
\[
= O(\sqrt{\eta} 2^{\zeta})||f||_{W^{1,\infty}} \int_{|y| < 2^{\zeta}} e^{-2\pi |y|^2} |y|dy + O(||f||_{L^\infty} \eta^\infty) = \]
\[
= O(||f||_{W^{1,\infty}} \sqrt{\eta}).
\]

The proof is complete.

Remark: A sharper version is used in subsection 4.1, namely
\[
|f(x) - \bar{f}(x)| \leq C \left( \sup_{|a| = 1} || \partial_a f || \right) \left( \frac{2}{\eta} \right)^{\frac{n}{2}} \int_{|x'| < 2^{\zeta - \frac{1}{2}}} e^{-2\pi \frac{|x'|^2}{\eta}} |x'|dx' + \]
\[
+ O(||f||_{L^\infty} \eta^\infty).
\]

The following observations are used in Section 5:

Observation 6.3 (Locality of the Liouville equation). For a Liouville equation with initial data of compact support \( \rho_0 \) and a \( C^{1,1} \) potential, interchanging the potential with any one that coincides with it on the “path” of the solution,
\[
(6.12) \quad \mathbb{S} = \bigcup_{t \in [0,T]} \phi_t(\text{supp} \rho_0),
\]
does not change the solution.

Also,
Lemma 6.4 (An $H^2$ estimate). Consider equation (2.8) with initial data $\rho_0 = [W_0^\varepsilon] \in H^2$. Assume that there is a function $\tilde{V}_{eff}^\gamma$ such that
\begin{equation}
\tilde{V} = \tilde{V}_{eff}^\gamma \quad \text{on } S,
\end{equation}
where $S$ is defined as in (6.12), and
\begin{equation}
\sup_{|A| \leq 3} ||\partial_A^A \tilde{V}_{eff}(x)||_{L^\infty} = O\left( \log \left( \varepsilon^{-\frac{3}{2}} \right) \right).
\end{equation}
Then,
\begin{equation}
||\rho_0^\varepsilon||_{H^2} = O\left( \varepsilon^{-\frac{3}{2}} ||W_0^\varepsilon||_{H^2} \right).
\end{equation}

Proof: The proof consists of using $\tilde{V}_{eff}^\gamma$ in place of $\tilde{V}$, making use of observation 6.3, and then applying directly theorem 6.1.

We will also use the following

Lemma 6.5. It is easy to observe that, if $\nabla V \in L^\infty$ and $\text{supp} W_0^\varepsilon$ is compact, it follows that
\begin{equation}
\exists M > 0 : \forall t \in [0, T], \varepsilon > 0 : \bigcup_{t \in [0, T]} \text{supp } \rho_0^\varepsilon(t) \subseteq \{ |(x, k)| < M \}.
\end{equation}

Observation 6.6 ($\tilde{V}_{eff}^\gamma$ and $R$). Let $V(x) = -|x|^{1+\theta}$. $\tilde{V}$ is a mollified version, according to equation (2.6).

Consider $\phi : \mathbb{R}^n \to [0, 1]$ to be a $C^\infty$ function such that
\begin{equation}
\phi(x) = 0, \quad |x| < \frac{1}{2},
\end{equation}
\begin{equation}
\phi(x) = 1, \quad |x| > 1,
\end{equation}
and $||\phi||_{W^{3, \infty}} \leq 10$. It is clear that such a function exists; this is of course an arbitrary requirement, but one that allows us not to carry the cutoff function $\phi$ to other results.

Then, setting $R \geq (\log(\varepsilon^{-1}))^{-\frac{1}{3}}$, $R = o(1)$, it follows that
\begin{equation}
\tilde{V}_{eff}(x) := \phi\left( \frac{x}{R} \right) \tilde{V}(x),
\end{equation}
satisfies
\begin{equation}
\sup_{|A| \leq 3} ||\partial_A^A \tilde{V}_{eff}(x)||_{L^\infty} = O\left( \log \left( \varepsilon^{-\frac{3}{2}} \right) \right),
\end{equation}
while, of course,
\begin{equation}
\tilde{V}_{eff}^\gamma = \tilde{V} \quad \text{on } \mathbb{R}^n \setminus \{ |x| < R \}.
\end{equation}
Proof: Observe that $1 - \phi(x) = \psi(x) \in {\mathcal S}(\mathbb{R}^n)$. Moreover, making use of observations 6.5 and 6.3, we can restrict \( \tilde{V} \) to \{\(|x| < M\)\} without loss of generality,

\[
\tilde{V}(x) \mapsto \tilde{V}(x) \left( 1 - \phi \left( \frac{x}{2M} \right) \right).
\]

Denote

\[
\tilde{V}_{\text{eff}}(x) = \tilde{V}(x) \left( 1 - \frac{x}{2M} \right) \phi \left( \frac{x}{R} \right) =
\]

\[
\tilde{V}(x) \left( \frac{x}{2M} \right) - \psi \left( \frac{x}{R} \right).
\]

Since

\[
||\partial^A_x \tilde{f}||_{L^\infty} \leq ||\partial^A_x f||_{L^\infty},
\]

it suffices to work with

\[
\partial^A_x |x|^{1+\theta} \left( \psi \left( \frac{x}{2M} \right) - \psi \left( \frac{x}{R} \right) \right).
\]

Set \( x' = \frac{x}{R} \); then \( \partial_x = R^{-1} \partial_{x'} \). Now we have

\[
||\partial^A_x |x|^{1+\theta} \left( \psi \left( \frac{x}{2M} \right) - \psi \left( \frac{x}{R} \right) \right) ||_{L^\infty} =
\]

\[
= R^{\theta-2} ||\partial^A_x |x|^{1+\theta} \left( \psi \left( \frac{R}{2M} x \right) - \psi(x) \right) ||_{L^\infty} \leq
\]

\[
\leq R^{\theta-2} \sup_{|A| \leq 3} ||\partial^A_x x^{1+\theta} ||_{L^\infty([1/2, 2M])} \sup_{|A| \leq 3} ||\partial^A_x \left( \psi \left( \frac{R}{2M} x \right) - \psi(x) \right) ||_{L^\infty(\mathbb{R}^n)} \leq
\]

\[
\leq (2M)^{1+\theta} ||\phi||_{W^{3,\infty}} R^{-3}.
\]

Recall that the constant \( M \) is chosen so that any trajectory leaving the support of the initial data doesn’t not escape \{\(|x| < M\)\} for \( t \in [0, T] \). That is, to properly quantify it, one needs to consider initial data of compact support associated with the Liouville equation, and a time-scale \( T \). Assuming that the initial data is supported in \{\(|(x, k)| < R_0\)\}, it is easy to check that one can set \( M = R_0 + 1 + R_0 T + ||\nabla V(x)||_{L^\infty} T^2 \). So now for (6.19) to hold it suffices that

\[
\frac{2^{1+\theta} (1 + R_0 + R_0 T + ||\nabla V(x)||_{L^\infty} T^2)^{1+\theta}}{C^2} R^{-3} \leq \log(\varepsilon^{-1})
\]

The constant \( C \) comes from the \( O(\cdot) \) of equation (6.19); it is clear that it can be chosen so that the constraint finally becomes \( R \geq (\log(\varepsilon^{-1}))^{-\frac{1}{2}} \), which is satisfied by choosing

\[
R = (\log(\varepsilon^{-1}))^{-\frac{1}{2}}.
\]
REFERENCES


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