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ON THE SELECTION OF THE CLASSICAL LIMIT FOR POTENTIALS WITH BV DERIVATIVES

AGISSILAOS ATHANASSOULIS AND THIERRY PAUL

ABSTRACT. We consider the classical limit of the quantum evolution, with some rough potential, of wave packets concentrated near singular trajectories of the underlying dynamics. We prove that under appropriate conditions, even in the case of BV vector fields, the correct classical limit can be selected.

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1. INTRODUCTION

The fundamental equation of quantum mechanics is formulated either as the Schrödinger equation for a wavefunction,

$$(1.1) \quad \begin{aligned} i\varepsilon \frac{\partial}{\partial t} u^\varepsilon(t) &= \left(-\frac{\varepsilon^2}{2} \Delta + V(x) \right) u^\varepsilon(t), \\ u^\varepsilon(t=0) &= u_0^\varepsilon(x), \end{aligned}$$

or as the Heisenberg-von Neumann equation for a density matrix $D^\varepsilon(t)$,

$$(1.2) \quad \begin{aligned} i\varepsilon \frac{\partial}{\partial t} D^\varepsilon(t) &= \left[-\frac{\varepsilon^2}{2} \Delta + V, D^\varepsilon(t) \right], \\ D^\varepsilon(t=0) &= D_0^\varepsilon. \end{aligned}$$

The state of the quantum system is described by the operator D^ε (understood to be the orthonormal projector $|u^\varepsilon\rangle\langle u^\varepsilon|$ when working with the Schrödinger equation), evolving in time under the aforementioned equations. Indeed, under very general conditions on the potential V , known as Kato's conditions, this is a problem well posed for $u_0^\varepsilon \in L^2$ or D_0^ε being a Hilbert-Schmidt operator [9, 14, 15]. It is also well known that the positivity and trace of D_0^ε is preserved in time. The parameter ε is called Planck's constant, and when it is very small one usually expects the system to behave like a classical one.

An equivalent way to write equation (1.2) is in terms of the Wigner transform,

$$(1.3) \quad W^\varepsilon(x, k, t) = W^\varepsilon[D^\varepsilon(t)](x, k) = \int_{y \in \mathbb{R}^n} e^{-2\pi i y k} K^\varepsilon(x + \varepsilon \frac{y}{2}, x - \varepsilon \frac{y}{2}, t) dy$$

where for each time t $K^\varepsilon(x, y, t)$ is the integral kernel of the operator $D^\varepsilon(t)$. In the case of a *pure state* this yields

$$(1.4) \quad W^\varepsilon(x, k, t) = W^\varepsilon[u^\varepsilon(t)](x, k) = \int_{y \in \mathbb{R}^n} e^{-2\pi i y k} u^\varepsilon(x + \varepsilon \frac{y}{2}, t) \bar{u}^\varepsilon(x - \varepsilon \frac{y}{2}, t) dy.$$

Equation (1.2) implies [10]

$$(1.5) \quad \partial_t W^\varepsilon + 2\pi k \cdot \partial_x W^\varepsilon + i \int e^{-2\pi i S y} \frac{V(x + \frac{\varepsilon}{2} y) - V(x - \frac{\varepsilon}{2} y)}{\varepsilon} dy W^\varepsilon(x, k - S, t) dS = 0,$$

$$W^\varepsilon(t = 0) = W_0^\varepsilon = W^\varepsilon[D_0^\varepsilon].$$

We will often use the shorthand T_ε^V for the operator involving the potential (see Definitions and Notations in section 3 below).

The Wigner transform (WT) respects the structure of the problem, allowing the Wigner equation (1.5) to inherit properties (and in particular conservation laws) from (1.2), see e.g. lemma 5.4. Moreover, under appropriate conditions, it allows for a very natural and compact description of the semiclassical limit, $\varepsilon \rightarrow 0$. Indeed the WT has a physically meaningful limit as ε tends to zero, while in general the wavefunction u^ε itself (or the operator D^ε) does not. The limit (in the weak-* topology of an appropriate algebra of test functions), called the Wigner measure,

$$W^\varepsilon \rightharpoonup W^0,$$

satisfies the Liouville equation of classical mechanics

$$(1.6) \quad \partial_t W^0 + 2\pi k \partial_x W^0 - \frac{1}{2\pi} \partial_x V \cdot \partial_k W^0 = 0,$$

$$W^0(t = 0) = \lim_{\varepsilon \rightarrow 0} W_0^\varepsilon$$

if the potential is regular enough [10, 7].

For potentials with low regularity (less than $C^{1,1}$) two different problems naturally arise;

Problem 1: showing that the Wigner measure is a weak solution of (1.6); and, if possible,

Problem 2: constructing a selection principle that identifies the correct weak solution, since problem (1.6) can be ill-posed in that case.

Problem 1 is solved in [10] for $V \in C^1$ (here we will work for quite less regular, if somewhat specific, potentials). As is highlighted already for $V \in C^1 \setminus C^{1,1}$, there can be several weak solutions when the initial datum is a singular measure (i.e. supported on a set of measure zero, the most natural case being a single point). That is what we mean in the sequel when we say the quantum initial data W_0^ε is *concentrating* – concentrating to a point-supported measure $W_0^0 = \lim_{\varepsilon \rightarrow 0} W_0^\varepsilon$. Indeed, the Liouville equation can be made well-posed for much worse potentials than C^1 , but typically not for concentrating initial data [1, 5, 8, 11].

More recently, Problem 1 has been successfully handled in [2, 6] under BV condition of the vector field generated by the classical Hamiltonian, and for **non-concentrating** (in the sense defined above) and **small** (in operator norm) initial data. In that situation the solution of (1.5) tends weakly to the push forward of the limit of the initial datum by the so-called DiPerna-Lions-Bouchut-Ambrosio flow (defined a.e.). In [4] a new technique was formulated for handling semiclassical limits with concentrating initial data. Parts of it were used to work out problems with low-regularity potentials (also in strong topology), in [3]; see also [12] for a short review. In a nutshell, the technique of [3] applies to more general initial data than the results of [2, 6] (including in particular data **concentrating to a point-supported measure**), but less general potentials (roughly speaking, of $C^{1,\alpha}$ type).

In the present note we compute explicitly the Wigner measures for certain problems with potentials that are not in C^1 , but have measure valued second derivative (therefore in a sense having the typical flavour of BV vector fields), and initial data which concentrate on the points of singularity. We illustrate the idea behind this computation and discuss how this idea can be generalized to other problems with similar regularity, and, roughly speaking, isolated repulsive singularities. The main difference from [3] is that here we manage to work with other topologies, more appropriate for the problem, and this is crucial to achieving stronger results.

This paper is mainly methodological: we have a technique that can select the correct semiclassical limit in some cases, and at the moment it's not clear how to characterize in an intuitive, economical way the class of problems to which it applies. We illustrate the proof, point out some exotic constructions that can be achieved with it, and outline how to check its applicability to a given problem (and what kind of problems can one reasonably hope that it applies to). A more precise and conclusive characterization of its applicability will have to be given in a forthcoming publication.

2. MAIN RESULT: A CASE STUDY

Since the regularity condition on the potential does not insure unicity of the flow, we will compare the solution of the quantum problem (1.5) to the behaviour of the solution of the Liouville equation **with the same, ε -dependent, initial condition**. We will show that the limit of the solution of the problem (1.5) is the same as the limit of the solution of the Liouville equation with the quantum (i.e. ε -dependent) initial datum. In other words, the extra information needed for the selection of the weak solution that correctly captures the semiclassical limit is fully contained in the way in which the initial datum concentrates to a point-supported measure.

Theorem 2.1 (Explicit example). *Let $\psi : [0, \infty) \rightarrow [0, 1]$ be a C^∞ monotone cutoff function, $\psi(s) = 1$ for all $s < \frac{1}{2}$, $\psi(s) = 0$ for all $s > 1$. Now, for $x \in \mathbb{R}^2$, some $\theta \in [0, 1)$ set*

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

Moreover let

$$F_0^\varepsilon = \delta_x^{-2} \delta_k^{-2} w \left(\frac{x - (0, -X)}{\delta_x}, \frac{k - (0, K)}{\delta_k} \right)$$

for some $w \in H^2 \cap L^\infty \cap W^{1,1}$, $w(x, k) \geq 0$, $\int w dx dk = 1$, $\text{supp } w \subseteq \{x^2 + k^2 \leq 1\}$, $X > 1, K > 0$ and $\delta_k = \delta_x^2 = (-\log(\varepsilon^{\frac{1}{4}}))^{\frac{1}{2}}$ (see claim 4.1 for other possible scalings).

Let W^ε be the solution of

$$(2.1) \quad \begin{aligned} & \partial_t W^\varepsilon(x, k, t) + 2\pi k \cdot \partial_x W^\varepsilon(x, k, t) + \\ & + i \int e^{-2\pi i S y} \frac{V(x + \frac{\varepsilon}{2} y) - V(x - \frac{\varepsilon}{2} y)}{\varepsilon} dy W^\varepsilon(x, k - S, t) dS = 0, \end{aligned}$$

$$W^\varepsilon(t = 0) = F_0^\varepsilon,$$

and ρ^ε the solution of

$$(2.2) \quad \begin{aligned} & \partial_t \rho^\varepsilon + 2\pi k \cdot \partial_x \rho^\varepsilon - \frac{1}{2\pi} \partial_x V(x) \cdot \partial_k \rho^\varepsilon = 0, \\ & \rho^\varepsilon(t = 0) = F_0^\varepsilon. \end{aligned}$$

Then, for each $t \in [0, T]$

$$(2.3) \quad \lim_{\varepsilon \rightarrow 0} \langle W^\varepsilon(t) - \rho^\varepsilon(t), \phi \rangle = 0 \quad \forall \phi \in \mathcal{A} \cap L^2$$

(see Definition 3.1 below for the algebra \mathcal{A}).

In other words the semiclassical limit of (2.1) behaves exactly in the same way as the concentration limit for the Liouville equation (2.2).

Remark: One reason the result is phrased the way it is, is that approximation (2.3) can hold even when there is no unique semiclassical limit. If, for example we restrict the values of $\varepsilon \in \{\frac{1}{m}\}$ and we substitute the initial datum $F_0^\varepsilon = F_0^{\frac{1}{m}}$ by $F_0^{\frac{1}{m}}(x_1 - d^{\frac{1}{m}}, x_2, k_1, k_2)$ for $d^{\frac{1}{m}} = C(-1)^m (\log m)^{-\frac{1}{2}}$, then it is easy to check that the result applies, but $W^{\frac{1}{m}}(t)$ – and $\rho^{\frac{1}{m}}(t)$ with it – has two accumulation points, one scattered to the right of $\{x_1 = 0\}$ and the other to the left (see also next remark). The finding here is that, whatever the interaction with the singularity, it is *the same* for the quantum problem (2.1) and the classical problem (2.2). In other words, all the information needed to determine the interaction is contained in the initial datum.

Remark: We can use the theorem to compute in more detail the semiclassical limit. Indeed,

$$(2.4) \quad W^0(t) = c_+ \delta(\mathcal{X}^+(t), \mathcal{P}^+(t)) + c_- \delta(\mathcal{X}^-(t), \mathcal{P}^-(t)),$$

with $c_\pm = \int_{\pm x_1 > 0} w(x, k) dx dk$ and $(\mathcal{X}^\pm(t), \mathcal{P}^\pm(t))$ are the two diverging curves obtained by $(\mathcal{X}^\pm(t), \mathcal{P}^\pm(t)) = \lim_{\eta \rightarrow 0} (\mathcal{X}_\eta^\pm(t), \mathcal{P}_\eta^\pm(t))$,

$$\begin{aligned} \dot{\mathcal{X}}_\eta^\pm(t) &= 2\pi \mathcal{P}_\eta^\pm(t), & \dot{\mathcal{P}}_\eta^\pm(t) &= -\frac{1}{2\pi} \nabla V(\mathcal{X}_\eta^\pm(t)), \\ \mathcal{X}_\eta^\pm(t = 0) &= (\pm\eta, -X), & \mathcal{P}_\eta^\pm(t = 0) &= (0, K). \end{aligned}$$

This can in fact be used to construct quite exotic examples; in figure 1 one sees various trajectories (their projection on the \mathcal{X} plane, to be precise) for various values of the parameter K . (The value of the parameter X is not really interesting, since it merely determines the length of free motion before interaction with the potential starts). Observe moreover, that by an appropriate superposition (discrete

or continuous) of initial data as in theorem 2.1, one can construct an example where the outgoing waves are (discretely or continuously) distributed among almost any angle.

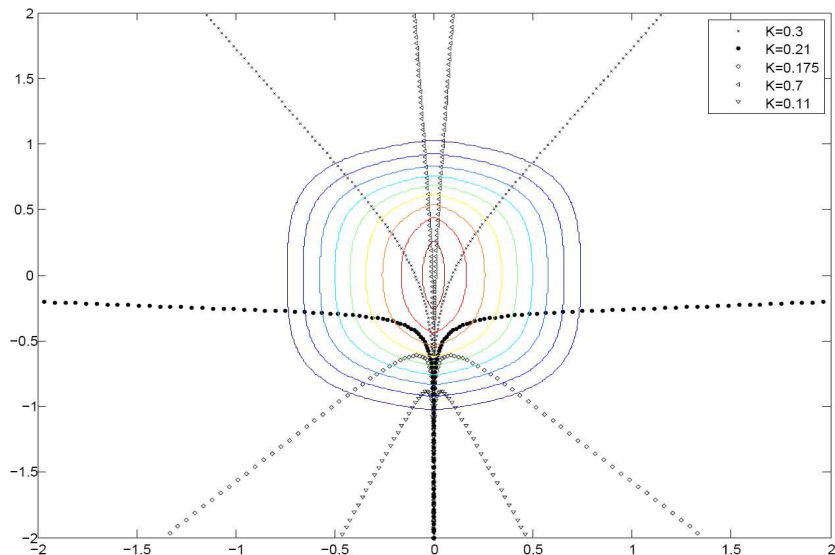


FIGURE 1. Pairs of trajectories $(\mathcal{X}^\pm, \mathcal{P}^\pm)$ corresponding to various values of the parameter K . The projection to the $(\mathcal{X}_1, \mathcal{X}_2)$ plane, as well as the contours of the potential $V(\mathcal{X}_1, \mathcal{X}_2)$ are shown.

This is a highly exotic example, and it may well turn out that its physical relevance is limited. However, it does point out a couple of points that could prove rather fruitful. First of all, the tools of semiclassical analysis can be extended and applied even to phenomena that are qualitatively very different from the smooth regime that they originate in. Such splitting of a particle and other possible exotic examples can motivate technical (and possibly numerical) experiments and refinements on existing methods and tools.

Moreover, this example points quite naturally to a more general (and, one expects, physical) situation: there is a “scattering process” happening “on the singularity”. That is, scattering with respect to local fast variables a neighbourhood of the singularity. In this particular case (of very slow concentration, since δ_x, δ_k are logarithmic in the semiclassical parameter) the scattering operator is “trivially” given by the concentration limit of the respective classical problem. That is, it is an operator fully determined and constructed by classical dynamics. In cases of faster concentration (e.g. pure states) one expects that an auxiliary *quantum problem* will need to be solved in a neighbourhood of the singularity to determine the “splitting”. In other words the irregular potential makes the corresponding flow to behave in $O(1)$ time-scales as a regular one would in *long times* (diverging trajectories), and this comparison can yield some very insightful finds [13].

Remark: This particular example rests on the exact alignment of the direction of propagation to the line of singularity (even a small rotation can, potentially, destroy it). As can be easily checked, if the potential is substituted by $V(x) = (1 - |x|^{1+\theta}) \psi(|x_1|)\psi(|x_2|)$ theorem 2.1 still holds, and then is robust to any perturbation of the initial data. Still, the only *nontrivial* case is when the particle passes *exactly* over the singularity, and the possibility of non-standard effects in that case.

At this point we must refer once again to [2]; a stable result, which covers “almost all initial data” – in which case the a.e. unique solution of the Liouville equation describes the semiclassical limit, and there is no room for exotic effects. In other words these interactions may be “atypical” in some sense, but if one tries to study them, they *have* to select initial data so that they fully interact.

Let us finally remark that an outline of how to check whether this technique applies to different problems is given in section 6.

3. DEFINITIONS AND NOTATIONS

The Fourier transform is defined as

$$(3.1) \quad \widehat{f}(k) = \mathcal{F}_{x \rightarrow k} [f(x)] = \int_{x \in \mathbb{R}^n} e^{-2\pi i k x} f(x) dx.$$

For compactness, we will use the following notations:

$$(3.2) \quad \mathcal{F}_2 W(x, K) = \int e^{-2\pi i k K} W(x, k) dk.$$

$$(3.3) \quad \begin{aligned} T_\varepsilon^V W &= i \int e^{-2\pi i S y} \frac{V(x + \frac{\varepsilon}{2} y) - V(x - \frac{\varepsilon}{2} y)}{\varepsilon} dy W(x, k - S) dS = \\ &= \mathcal{F}_{K \rightarrow k}^{-1} \left[\frac{V(x + \frac{\varepsilon}{2} K) - V(x - \frac{\varepsilon}{2} K)}{\varepsilon} \mathcal{F}_2 W(x, K) \right] = \mathcal{F}_{K \rightarrow k}^{-1} \left[\int_{s = -\frac{K}{2}}^{\frac{K}{2}} \nabla V(x + s) ds \mathcal{F}_2 W(x, K) \right], \end{aligned}$$

$$(3.4) \quad T_0^V W = -\frac{1}{2\pi} \partial_x V \cdot \partial_k W = \mathcal{F}_{K \rightarrow k}^{-1} [\nabla V(x) \cdot K \mathcal{F}_2 W(x, K)].$$

Definition 3.1. *The space \mathcal{A} is defined as the completion of the smooth functions of compact support $C_c^\infty(\mathbb{R}^{2n})$ under the norm*

$$\|f\|_{\mathcal{A}} = \int_K \sup_x |\mathcal{F}_2 f(x, K)| dK.$$

It follows that

$$\|f\|_{\mathcal{A}'} = \sup_K \int_x |\mathcal{F}_2 f(x, K)| dx.$$

Some basic properties that we will use are

$$\begin{aligned} \|f\|_{L^\infty} &\leq \|f\|_{\mathcal{A}} \\ \|f\|_{\mathcal{A}'} &\leq \|f\|_{L^1} \leq \|f\|_{\mathcal{M}} \end{aligned}$$

Definition 3.2. Every finite signed measure can be decomposed to positive and negative part, $\mu = \mu^+ - \mu^-$, for some finite non-negative measures μ^+, μ^- . We will denote the total variation of a signed measure

$$\|\mu\|_{\mathcal{M}} = \int_{\mathbb{R}^n} d\mu^+ + \int_{\mathbb{R}^n} d\mu^-.$$

Definition 3.3. Denote by $\Phi = \Phi^\varepsilon$ the smoothing operator

$$\Phi^\varepsilon : f(x) \mapsto \left(\frac{2}{\varepsilon}\right)^{\frac{n}{2}} \int e^{-2\pi \frac{(x-x')^2}{\varepsilon}} f(x') dx'.$$

When there is no danger of confusion, we will write

$$\tilde{f} := \Phi^\varepsilon f.$$

Definition 3.4. The Sobolev space $W^{m,p}$ is defined as the completion of the smooth functions of compact support $C_c^\infty(\mathbb{R}^n)$ under the norm

$$\|f\|_{W^{m,p}} = \sum_{|A| \leq m} \|\nabla^A f\|_{L^p}.$$

4. PROOF

4.1. Strategy of the proof. The intuition behind this proof is quite simple: the quantum initial datum, while concentrating to a point still is an L^2 function, hence the a.e. theory for its evolution under the Liouville equation applies. The concentration limit of problem (2.2) is therefore a natural candidate for the semiclassical limit. Of course to show a semiclassical approximation, a certain degree of smoothness in the potential and the Wigner function is needed, and in any frontal approach to this problem, such smoothness simply is not there.

A simple idea to try out is the following: can we cut-off the pieces of the Wigner function that approach too closely to the singularity? If we do that, can we meaningfully use an auxiliary function supported just far enough away from the singularity so that it preserves enough smoothness itself – as well as allowing one to cut-off the potential’s singularities? This was pretty much the program we followed before, in [3].

The new element here is that to strengthen that technique to potentials as bad as the ones we treat here, basically a much bigger piece of the initial datum would have to be cut-off, and we need to find a meaningful way in which such a cut-off introduces “small” errors. To do that, the positivity of the density matrix comes into play, and some very different considerations are needed.

Claim 4.1 holds all the compromises that need to be made, and is the conclusion of a lot trial and error. If one assumes it and move on in a first reading, the flow of the rest proof should provide a reasonable motivation for why these computations have to be just so.

4.2. Proof of theorem 2.1. To facilitate the presentation, let us introduce at this point a number of auxiliary functions (using the notations introduced in the

previous section):

$$(4.1) \quad \begin{aligned} \partial_t W^\varepsilon + 2\pi k \cdot \partial_x W^\varepsilon + T_\varepsilon^V W^\varepsilon &= 0, \\ W^\varepsilon(t=0) &= F_0^\varepsilon, \end{aligned}$$

$$(4.2) \quad \begin{aligned} \partial_t W_1^\varepsilon + 2\pi k \cdot \partial_x W_1^\varepsilon + T_\varepsilon^V W_1^\varepsilon &= 0, \\ W_2^\varepsilon(t=0) &= F_1^\varepsilon = \Phi F_0^\varepsilon, \end{aligned}$$

$$(4.3) \quad \begin{aligned} \partial_t W_2^\varepsilon + 2\pi k \cdot \partial_x W_2^\varepsilon + T_\varepsilon^V W_2^\varepsilon &= 0, \\ W_2^\varepsilon(t=0) &= F_2^\varepsilon = \Phi(1 - \psi(\frac{x}{R}))F_0^\varepsilon, \end{aligned}$$

$$(4.4) \quad \begin{aligned} \partial_t W_3^\varepsilon + 2\pi k \cdot \partial_x W_3^\varepsilon + T_\varepsilon^V W_3^\varepsilon &= 0, \\ W_3^\varepsilon(t=0) &= F_3^\varepsilon = (1 - \psi(\frac{x}{R}))F_0^\varepsilon, \end{aligned}$$

$$(4.5) \quad \begin{aligned} \partial_t \rho_1^\varepsilon + 2\pi k \cdot \partial_x \rho_1^\varepsilon + T_0^V \rho_1^\varepsilon &= 0, \\ \rho_1^\varepsilon(t=0) &= F_3^\varepsilon, \end{aligned}$$

$$(4.6) \quad \begin{aligned} \partial_t \rho^\varepsilon + 2\pi k \cdot \partial_x \rho^\varepsilon + T_0^V \rho^\varepsilon &= 0, \\ \rho^\varepsilon(t=0) &= F_0^\varepsilon. \end{aligned}$$

The function ψ is the same as in the statement of theorem 2.1, and R' will be set below.

Obviously, we are going to use these functions as stepping stones, passing from one to the other with the appropriate topology each time. Because this topology cannot be always the same, the end result is formulated as is, in weak sense: for all $\phi \in L^2 \cap \mathcal{A}$

$$\begin{aligned} |\langle W^\varepsilon(t) - \rho^\varepsilon(t), \phi \rangle| &\leq |\langle W^\varepsilon(t) - W_1^\varepsilon(t), \phi \rangle| + |\langle W_1^\varepsilon(t) - W_2^\varepsilon(t), \phi \rangle| + \\ &+ |\langle W_2^\varepsilon(t) - W_3^\varepsilon(t), \phi \rangle| + |\langle W_3^\varepsilon(t) - \rho_1^\varepsilon(t), \phi \rangle| + |\langle \rho_1^\varepsilon(t) - \rho^\varepsilon(t), \phi \rangle| \end{aligned}$$

The point is to collect the various constraints that would come from each of these building-block problems and satisfy them at the same time. That is essentially done in the following

Claim 4.1. *With appropriate calibration of δ_x, δ_k, R' , we have*

$$\|F_3^\varepsilon - F_0^\varepsilon\|_{L^1} = o(1),$$

$$\|\rho_1^\varepsilon\|_{L^\infty([0,T],H^2)} = o(\varepsilon^{-\frac{1}{2}}),$$

and

$$\|F_2^\varepsilon - F_3^\varepsilon\|_{L^2} = o(1).$$

Proof of the claim: The point is to make sure that $\rho_1^\varepsilon(t)$ doesn't pass through a neighbourhood of the set where the second and third derivatives of the potential are singular (or too large in any case); in this case the strip $\{|x| < R\}$, $R = (-\log(\varepsilon^{\frac{1}{4}}))^{\frac{1}{\theta-2}}$, for $t \in [0, T]$. In that case, using lemma 5.1 it follows readily that indeed

$$(4.7) \quad \|\rho_1^\varepsilon\|_{L^\infty([0,T],H^2)} = O(\varepsilon^{-\frac{1}{4}} \|F_3^\varepsilon\|_{H^2}).$$

(by recalling that in fact only the values of the potential's derivatives along the part of phase space that the solutions passes through matter).

To ensure that $\rho_1^\varepsilon(t)$ stays away from $\{|x| < R\}$, it turns out to be sufficient to exclude a somewhat larger strip $\{|x| < R'\}$. To see why, imagine firstly that we are in free space, $V = 0$. The largest possible momentum in the x_1 direction is δ_k ; in time T this can only cover a distance of $T\delta_k$; therefore $R' = R + T\delta_k$ should do it.

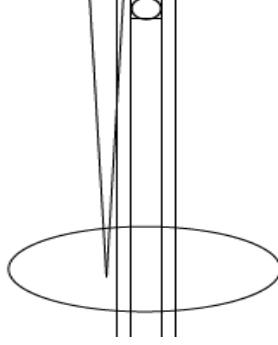


FIGURE 2. Any trajectory leaving the support of F_3^ε has an initial velocity in a small cone around $k = (0, K)$. If we were in free space, there wouldn't be enough time for sufficient movement in x_1 to reach $\{|x_1| < R\}$ for $t \in [0, T]$. Therefore the solution is the sum of two components, supported on either side of $\{|x_1| < R\}$. The time-scale T enters in the constants used, but the construction can be made for any fixed T .

Now we have to take into account the addition of the potential. The repulsive nature of the singularity (preserved by the cutoff $\psi(|x_1|)$) makes sure that any trajectory moving towards it will not be accelerated – but in fact slowed down; i.e. any movement towards the singularity will be *less* than $T\delta_k$. The presence of $\psi(|x_2|)$ serves only to slow somewhat movement in the x_2 direction (or even turn it back); but makes no difference whatsoever in the x_1 direction.

The timescale T is interesting when trajectories starting in the support of F_0^ε reach and leave the support of the potential; in any case $T = O(1)$.

Moreover, for the first part of the claim,

$$\begin{aligned} \|F_3^\varepsilon - F_0^\varepsilon\|_{L^1} &\leq \|(1 - \Phi^\varepsilon)\delta_x^{-2}\delta_k^{-2}w\left(\frac{x}{\delta_x}, \frac{k}{\delta_k}\right)\|_{L^1} + \\ &\quad + \|\psi\left(\frac{x}{R'}\right)\delta_x^{-2}\delta_k^{-2}w\left(\frac{x}{\delta_x}, \frac{k}{\delta_k}\right)\|_{L^1} \leq \\ &\leq \sqrt{\varepsilon}\|\delta_x^{-2}\delta_k^{-2}w\left(\frac{x}{\delta_x}, \frac{k}{\delta_k}\right)\|_{W^{1,1}} + \|w\|_{L^\infty}\|\psi\left(\frac{x\delta_x}{R'}\right)\|_{L^1} = \\ &= O\left(\frac{\sqrt{\varepsilon}}{\delta_x} + \frac{\sqrt{\varepsilon}}{\delta_k} + \frac{R'}{\delta_x}\right) \end{aligned}$$

Finally, for the third part,

$$\|F_2^\varepsilon - F_3^\varepsilon\|_{L^2} = \|(I - \Phi^\varepsilon)F_3^\varepsilon\|_{L^2} \leq \sqrt{\varepsilon}\|F_3^\varepsilon\|_{H^1} = O\left(\frac{\sqrt{\varepsilon}}{R}(\delta_x^{-1} + \delta_k^{-1})\delta_x^{-1}\delta_k^{-1}\right)$$

So collecting all the constraints, we have

$$\begin{aligned} \sqrt{\varepsilon} \ll \delta_x, \delta_k; R' \ll \delta_x &\Rightarrow \|F_3^\varepsilon - F_0^\varepsilon\|_{L^1} = o(1), \\ R' = (-\log(\varepsilon^{\frac{1}{4}}))^{\frac{1}{\theta-2}} + C\delta_k &\Rightarrow \|\rho_1^\varepsilon\|_{L^\infty([0,T],H^2)} = O(\varepsilon^{-\frac{1}{4}}\|F_3^\varepsilon\|_{H^2}), \\ (\delta_x^{-2} + \delta_k^{-2})\delta_x^{-1}\delta_k^{-1} \ll \varepsilon^{-\frac{1}{4}} &\Rightarrow \|F_3^\varepsilon\|_{H^2} = o(\varepsilon^{-\frac{1}{4}}), \\ \sqrt{\varepsilon} \ll R'\delta_x^2\delta_k, \sqrt{\varepsilon} \ll R'\delta_x\delta_k^2 &\Rightarrow \|F_2^\varepsilon - F_3^\varepsilon\|_{L^2} = o(1). \end{aligned}$$

A concrete scaling that makes all these constraints covalid is (recall that $R = (-\log(\varepsilon^{\frac{1}{4}}))^{\frac{1}{\theta-2}}$)

$$\delta_k = \delta_x^2 = R.$$

The proof of the claim is complete.

Remark: The claim apparently depends on the geometry of the problem; but really all we used was that the flow is repulsive away from the singularity. That is, that any trajectory would not go towards the line $\{x_1 = 0\}$ faster than it would on free space. Thus it makes no difference if e.g. a different potential – for which the same property holds – is used. Indeed if the potential was $V(x) = (1 - |x|^{1+\theta}) \psi(|x_1|)\psi(|x_2|)$ nothing needs to be changed.

Given this calibration, the proof proceeds in a very predictable fashion. Indeed, using the unitary propagation of the Wigner equation – theorem 5.4 – one observes that

$$\|F_0^\varepsilon - F_1^\varepsilon\|_{L^2} \leq \sqrt{\varepsilon}\|F_0^\varepsilon\|_{H^1} = O(\sqrt{\varepsilon}\delta_x\delta_k(\delta_x + \delta_k)) = o(1) \Rightarrow \|W_0^\varepsilon(t) - W_1^\varepsilon(t)\|_{L^2} = o(1) \forall t.$$

The estimate between W_1^ε and W_2^ε is one of the essential parts, and the technical innovation here. Denote by $U^\varepsilon(t)$ the propagator of the Wigner equation; we observe that

$$\begin{aligned} W_1^\varepsilon - W_2^\varepsilon &= U^\varepsilon(t)(F_1^\varepsilon - F_2^\varepsilon) = \\ &= U^\varepsilon(t) \left(\frac{2}{\varepsilon}\right)^2 \int e^{-2\pi\frac{(x-x')^2 + (k-k')^2}{\varepsilon}} \psi\left(\frac{x'}{R}\right)\delta_x^{-2}\delta_k^{-2} w\left(\frac{x' - (0,-2)}{\delta_x}, \frac{k' - (0,1)}{\delta_k}\right) dx' dk', \end{aligned}$$

the point being that $F_1^\varepsilon - F_2^\varepsilon$ is itself a density matrix – convolution of a positive measure with a coherent state; see lemma 5.2. Moreover, it is a *small* density matrix, and this is preserved in time; see lemma 5.5. Indeed it follows by the conservation of trace (and lemma 5.3 for density matrices) that

$$\|W_1^\varepsilon(t) - W_2^\varepsilon(t)\|_{\mathcal{A}'} = \|F_1^\varepsilon - F_2^\varepsilon\|_{\mathcal{A}'} \leq \|F_1^\varepsilon - F_2^\varepsilon\|_{L^1} = O\left(\frac{R'}{\delta_x}\right)$$

which was already calibrated to be $o(1)$ in the proof of claim 4.1. (Indeed following that calibration we have $\frac{R'}{\delta_x} = O(\sqrt{\delta_k} + \delta_x)$). Observe that $F_1^\varepsilon - F_2^\varepsilon$ is *not* $o(1)$ in L^2 sense, hence the introduction of the L^1 -like norm is necessary.

Moreover, by virtue of theorem 5.4 and claim 4.1

$$\|W_2^\varepsilon(t) - W_3^\varepsilon(t)\|_{L^2} = \|F_2^\varepsilon - F_3^\varepsilon\|_{L^2} = o(1).$$

The other non-trivial step is the quantum-classical dynamics comparison between W_3^ε and ρ_1^ε . Recall that $U^\varepsilon(t)$ the propagator of the Wigner equation, and denote

by $E(t)$ the propagator of the Liouville equation. Set $h^\varepsilon(t) = W_3^\varepsilon(t) - \rho_1^\varepsilon(t)$; then of course $h^\varepsilon(t) = \int_{\tau=0}^t U^\varepsilon(t)(T_\varepsilon^V - T_0^V)\rho_1^\varepsilon(\tau)d\tau$.

Using the Fourier transform in the k variable we have

$$\begin{aligned} & \left\| \int_{\tau=0}^t U(t)(T_\varepsilon^V - T_0^V)\rho_2^\varepsilon(\tau)d\tau \right\|_{L^2} \leq \\ & \leq T \sup_{t \in [0, T]} \left\| \left(\frac{V(x + \frac{\varepsilon K}{2}) - V(x - \frac{\varepsilon K}{2})}{\varepsilon} - \partial_x V(x) \cdot K \right) \mathcal{F}_2 \rho_2^\varepsilon(t) \right\|_{L^2} \leq \\ & \leq T \sup_{t \in [0, T]} \|\chi_{[0, \varepsilon^{-\frac{1}{2}}]}(|K|) \left(\frac{V(x + \frac{\varepsilon K}{2}) - V(x - \frac{\varepsilon K}{2})}{\varepsilon} - \partial_x V(x) \cdot K \right) \mathcal{F}_2 \rho_2^\varepsilon(t)\|_{L^2} + \\ & + T \sup_{t \in [0, T]} \|\chi_{(\varepsilon^{-\frac{1}{2}}, +\infty)}(|K|) \left(\frac{V(x + \frac{\varepsilon K}{2}) - V(x - \frac{\varepsilon K}{2})}{\varepsilon} - \partial_x V(x) \cdot K \right) \mathcal{F}_2 \rho_2^\varepsilon(t)\|_{L^2} \leq \\ & \leq \varepsilon \|\chi_{[0, \varepsilon^{-\frac{1}{2}}]}(|K|) \frac{|K|^2}{2} \sup_{\substack{\tau \in (-1, 1) \\ |A|=2}} |\partial_x^A V(x + \tau \frac{\varepsilon K}{2})| \mathcal{F}_2 \rho_2^\varepsilon(t)\|_{L^2} + O(\varepsilon^{\frac{1}{2}}) \| |K|^2 \mathcal{F}_2 \rho_2^\varepsilon(t) \|_{L^2}. \end{aligned}$$

For the first term we used a straightforward Taylor expansion, while for the complementary case we used the computation

$$\begin{aligned} & \left| \left(\frac{V(x + \frac{\varepsilon K}{2}) - V(x - \frac{\varepsilon K}{2})}{\varepsilon} - \partial_x V(x) \cdot K \right) \right| \leq \\ & \leq \left| \frac{1}{\varepsilon} \int_{-\frac{\varepsilon|K|}{2}}^{\frac{\varepsilon|K|}{2}} \left(\frac{\varepsilon|K|}{2} - t \right) \sum_{j,m} z_j z_m \partial_{x_j x_m} V(x + tz) dt \right| \leq C|K| \sum_{j,m} \left| \int_{-\frac{\varepsilon|K|}{2}}^{\frac{\varepsilon|K|}{2}} \partial_{x_j x_m} V(x + tz) dt \right| \leq \\ & \leq C|K| \sum_{A=2} \|\partial_x^A V\|_{\mathcal{M}}. \end{aligned}$$

The key to proceed is to make use of the well prepared initial datum F_3^ε : denote

$$S^\varepsilon = \{x \exists t \in [0, T], k \in \mathbb{R}^2 : \rho_1^\varepsilon(x, k, t) > 0\}$$

and

$$B^\varepsilon = \bigcup_{x_* \in S^\varepsilon} \{x \mid |x - x_*| < \sqrt{\varepsilon}\}.$$

Then

$$\chi_{[0, \varepsilon^{-\frac{1}{2}}]}(|K|) \sup_{\substack{\tau \in (-1, 1) \\ |A|=2}} |\partial_x^A V(x + \tau \frac{\varepsilon K}{2})| \leq \sup_{\substack{x \in B^\varepsilon \\ |A|=2}} |\partial_x^A V(x)|$$

and therefore

$$\begin{aligned} & \varepsilon \chi_{[0, \varepsilon^{-\frac{1}{2}}]}(|K|) \sup_{\substack{\tau \in (-1, 1) \\ |A|=2}} |\partial_x^A V(x + \tau \frac{\varepsilon K}{2})| \| |K|^2 \mathcal{F}_2 \rho_2^\varepsilon(t) \|_{L^2} \leq \\ & \leq \varepsilon \sup_{\substack{x \in B^\varepsilon \\ |A|=2}} |\partial_x^A V(x)| \sup_{t \in [0, T]} \|\rho_2^\varepsilon(t)\|_{H^2} \end{aligned}$$

finally yielding

$$\begin{aligned} & \left\| \int_{\tau=0}^t U(\tau)(T_\varepsilon^V - T_0^V)\rho_2^\varepsilon(\tau)d\tau \right\|_{L^2} = \\ & = O(\sqrt{\varepsilon}\|\rho_2^\varepsilon(t)\|_{H^2}(1 + \sqrt{\varepsilon} \sup_{\substack{x \in B^\varepsilon \\ |A|=2}} |\partial_x^A V(x)|)) = O(\varepsilon^{\frac{1}{4}}\|F_3^\varepsilon\|_{H^2}(1 + R^{\theta-1})) = o(1). \end{aligned}$$

This is possible precisely because our approximate initial datum F_3^ε stays away from a neighbourhood of the singular set $\{|x_1| = 0\}$ when pushed forward in time by the a.e. Ambrosio-Lions-Di Perna flow $E(t)$, as was checked in claim 4.1. (For a potential of the form $|x|$ the second derivatives are zero almost everywhere, and this creates the possibility of taking advantage of a very special structure. That's why we included the case $\theta > 0$ in the statement of the theorem, to show that in principle this technique work for a variety of localized repulsive singularities).

5. AUXILIARY RESULTS

Lemma 5.1 (2^{nd} order derivatives equations for the Liouville equation). *Consider the Cauchy problem for the Liouville equation with potential $V(x)$ on \mathbb{R}^n ,*

$$(5.1) \quad \begin{aligned} \partial_t f + 2\pi k \cdot \partial_x f - \frac{1}{2\pi} \partial_x V(x) \cdot \partial_k f &= 0, \\ f(t=0) &= f_0. \end{aligned}$$

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

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

Lemma 5.2 (Density matrices). *Let μ be a probability measure on \mathbb{R}^{2n} . Then*

$$W_0^\varepsilon(x, k) = \left(\frac{2}{\varepsilon}\right)^n \int e^{-2\pi \frac{(x-x')^2 + (k-k')^2}{\varepsilon}} d\mu(x', k')$$

is the Wigner function of a density matrix, i.e. the corresponding operator D^ε is a positive trace-class operator with $\text{tr}(D^\varepsilon) = 1$.

Proof: We know that $\text{tr}(D^\varepsilon) = \int W^\varepsilon dx dk = 1$. Let us now look at positivity. To that end, observe that the integral kernel

$$K^\varepsilon(x, y) = \left(\frac{2}{\varepsilon}\right)^{\frac{n}{2}} \int_{x_0, k_0} e^{2\pi i \frac{k_0}{\varepsilon}(x-y)} e^{-\frac{\pi}{\varepsilon} [(x-x_0)^2 + (y-x_0)^2]} d\mu(x_0, k_0)$$

is the kernel of a positive operator. Indeed:

$$\int K^\varepsilon(x, y) u(x) \bar{u}(y) dx dy = \left(\frac{2}{\varepsilon}\right)^{\frac{n}{2}} \int_{x_0, k_0} | \langle e^{-2\pi i \frac{k_0}{\varepsilon} x + \frac{\pi}{\varepsilon} (x-x_0)^2}, u \rangle |^2 d\mu(x_0, k_0) \geq 0.$$

The proof is complete by observing that the Wigner function corresponding to the kernel ρ^ε is

$$\int e^{-2\pi i k y} K^\varepsilon\left(x + \frac{\varepsilon y}{2}, x - \frac{\varepsilon y}{2}\right) dy = W_0^\varepsilon.$$

Lemma 5.3. *For any trace class operator D^ε , with corresponding Wigner function W^ε ,*

$$\|W^\varepsilon\|_{\mathcal{A}'} \leq \|D^\varepsilon\|_{tr} = tr(|D^\varepsilon|).$$

Moreover, if $D^\varepsilon \geq 0$

$$\|W^\varepsilon\|_{\mathcal{A}'} = \|D^\varepsilon\|_{tr} = tr(D^\varepsilon).$$

Proof. Let D^ε be a non-negative trace-class operator. Then, it admits a SVD expansion over orthonormal projectors

$$D^\varepsilon = \sum_m \lambda_m |u_m\rangle\langle u_m|,$$

where of course $\|\lambda_m\|_{l^1} = \|D^\varepsilon\|_{tr} = tr(D^\varepsilon)$, i.e. $\lambda_m \geq 0$, and $\langle u_m, u_l \rangle = \delta_{m,l}$. It follows that

$$W^\varepsilon = \sum_m \lambda_m W^\varepsilon[u_m] = \sum_m \lambda_m \int e^{-2\pi i k y} u_m(x + \frac{\varepsilon y}{2}) \bar{u}_m(x - \frac{\varepsilon y}{2}) dy$$

and, by straightforward substitution,

$$\begin{aligned} \|W^\varepsilon\|_{\mathcal{A}'} &= \sup_K \int_x |\sum_m \lambda_m u_m(x + \frac{\varepsilon K}{2}) \bar{u}_m(x - \frac{\varepsilon K}{2})| dx \leq \\ &\leq \sum_m \lambda_m \sup_K \int_x |u_m(x + \frac{\varepsilon K}{2}) \bar{u}_m(x - \frac{\varepsilon K}{2})| dx = \sum_m \lambda_m. \end{aligned}$$

Now observe that

$$\begin{aligned} \|W^\varepsilon\|_{\mathcal{A}'} &= \sup_K \int_x |\sum_m \lambda_m u_m(x + \frac{\varepsilon K}{2}) \bar{u}_m(x - \frac{\varepsilon K}{2})| dx \geq \\ &\geq \int_x \sum_m \lambda_m |u_m(x)|^2 dx = \sum_m \lambda_m. \end{aligned}$$

Hence, $D^\varepsilon \geq 0 \Rightarrow \|W^\varepsilon\|_{\mathcal{A}'} = tr(D^\varepsilon) = \|D^\varepsilon\|_{tr}$.

To conclude, observe that for $D^\varepsilon = D_+^\varepsilon - D_-^\varepsilon$

$$\|W^\varepsilon\|_{\mathcal{A}'} \leq \|W^\varepsilon[D_+^\varepsilon]\|_{\mathcal{A}'} + \|W^\varepsilon[D_-^\varepsilon]\|_{\mathcal{A}'} = tr(D_+^\varepsilon) + tr(D_-^\varepsilon) = \|D^\varepsilon\|_{tr}.$$

The proof is complete

The following theorem is well known, and follows from the correspondence between the density matrix and the Wigner function; see e.g. Theorem 2.1 of [11]:

Theorem 5.4 (L^2 regularity of the Wigner equation). *If the Schrödinger operator $-\frac{1}{2}\Delta + V(x)$ is essentially self-adjoint on $L^2(\mathbb{R}^n)$, then the corresponding Wigner equation preserves the L^2 norm, i.e. for $W_0^\varepsilon \in L^2(\mathbb{R}^{2n})$ there is a unique solution of*

$$(5.2) \quad \begin{aligned} \partial_t W^\varepsilon + 2\pi k \cdot \partial_x W^\varepsilon + \frac{2}{\varepsilon} Re \left[i \int e^{2\pi i S x} \hat{V}(S) W^\varepsilon(x, k - \frac{\varepsilon S}{2}) dS \right] &= 0, \\ W^\varepsilon(t=0) &= W_0^\varepsilon, \end{aligned}$$

and $\|W^\varepsilon(t)\|_{L^2(\mathbb{R}^{2n})} = \|W_0^\varepsilon\|_{L^2(\mathbb{R}^{2n})} \quad \forall t \in \mathbb{R}$.

Lemma 5.5 (Trace conservation). *Let W_0^ε be a Wigner function corresponding to a trace-class operator. Consider a potential V such that the corresponding Schrödinger operator is essentially self-adjoint. If by $W^\varepsilon(t)$ we denote its evolution in time under the corresponding Wigner equation, then*

$$tr(W_0^\varepsilon) = tr(W^\varepsilon(t)) \quad \forall t.$$

Proof. One easily checks that $W_0^\varepsilon = \sum_m \lambda_m |u_m\rangle \langle u_m| \Rightarrow W_0^\varepsilon = \sum_m \lambda_m |u_m(t)\rangle \langle u_m(t)|$,

where $u_m(t)$ is the evolution in time of $u_m \in L^2$ under the Schrödinger equation.

The result follows.

6. EXTENSIONS

As was mentioned earlier, this approach also applies if the potential is substituted by $V(x) = (1 - |x|^{1+\theta}) \psi(|x_1|)\psi(|x_2|)$. Other straightforward generalizations come by embedding the problem in higher-dimensional space (i.e. including more transverse dimensions).

Checking whether some version of this approach applies to a problem would start by building the counterpart of claim 4.1, and its second part in particular. More specifically:

- Set $N = \{(x, k) \mid \sup_{|A|=3} |\partial_x^A V(x)| > \zeta\}$ (ζ is a parameter of the order $-\log \varepsilon$ to some power). This is the “neighbourhood of the singularity” that we want to stay away from.

- Back-propagate it for the appropriate time scale, $M = \bigcup_{t \in [0, T]} \phi_{-t}(N)$. The

counterpart to F_0^3 should be chosen so that it does not enter M . M is not a small set; but we only need that cutting it off makes a small difference to F_0^ε , not that itself it is small. $F_0^\varepsilon(1 - \chi_M(x, k))$ is a starting point for F_3^ε ; a simpler cutoff might be preferable. Certainly though, if $\|F_0^\varepsilon(1 - \chi_M(x, k))\|_{L^1}$ is not small, then this approach does not apply.

- The previous step is only one constraint about at least how much we have to cut-off. If we cut-off any less (i.e. exclude a smaller set), we will enter too close to the singularity and the estimates will fail. An other constraint comes from having several H^2 bounds to check for our various approximate initial data. This means that if we exclude too small sets (e.g. too thin strips etc), the derivatives of the approximate data will be too large, and the estimates also fail. This is what makes us exclude a strip even when $\theta = 0$ in theorem 2.1, where the previous step would be fulfilled by taking N to be a line.

In other words, there are two different origins for constraints in this construction: on the one hand, we want to avoid the low-regularity region. The more we cut-off the better. On the other hand, we want the cut-offs to introduce small errors, including in various derivative-norms. This is a double sided constraint, as cutting off either a too thick or too thin strip will fail here. Of course the norms we can work in are constrained by the available conservation laws for the equations we work with, i.e. apparently have to be based on L^2 , and \mathcal{A}' as long as we keep track successfully of positivity. The fact that \mathcal{A}' scales like L^1 in concentrating data is crucial in allowing wider strips to be cut-off.

The fact that we will have to exclude domains measured in $-(\log \varepsilon)^{-1}$ (to some power) means that in general one can only work with slowly concentrating initial data. Moreover, the *repulsive* character of the singularity is used implicitly here, as it helps the set M not be too large, but, when projected on the hyperplane transverse to propagation, basically looks like the set N . It is not necessary per se, but a way to meaningfully control how much larger is the set M from N is needed.

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