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# UNIFORM RESOLVENT ESTIMATES FOR A NON-DISSIPATIVE HELMHOLTZ EQUATION 

Julien Royer

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# UNIFORM RESOLVENT ESTIMATES FOR A NON-DISSIPATIVE HELMHOLTZ EQUATION 

by Julien Royer


#### Abstract

We study the high frequency limit for a non-selfadjoint Helmholtz equation. This equation models the propagation of the electromagnetic field of a laser in an inhomogeneus material medium with non-constant absorption index. In this paper the absorption index can take negative values and we only use a damping condition on the classical limit of the problem. In this setting we first prove the absence of eigenvalue on the upper half-plane and close to an energy which satisfies this damping assumption. Then we generalize the resolvent estimates of Robert-Tamura and prove the limiting absorption principle. We finally study the semiclassical measures of the solution when the source term concentrates on a bounded submanifold of $\mathbb{R}^{n}$.


## 1. Introduction and statement of the main results

The purpose of this paper is to study on $\mathbb{R}^{n}, n \geqslant 1$, the high frequency limit for the Helmholtz equation in a non-dissipative setting. After rescaling, this equation can be written

$$
\begin{equation*}
\left(H_{h}-E\right) u_{h}=f_{h}, \quad \text { where } \quad H_{h}=-h^{2} \Delta+V_{1}(x)-i h V_{2}(x) . \tag{1.1}
\end{equation*}
$$

We recall that this equation models for instance the propagation of the electromagnetic field of a laser in an inhomogeneus material medium. In this setting

[^0]$V_{1}(x)-E$ is linked to the refraction index, $V_{2}(x)$ is the absorption index and $f_{h}$ is the source term. The parameter $h>0$ is proportional to the wavelength. In this paper we are interested in the asymptotic behavior of the solution $u_{h}$ when $h$ goes to 0 .

All along this paper, we assume that $V_{1}$ and $V_{2}$ are bounded and go to 0 at infinity. This implies in particular that the essential spectrum of $H_{h}$ on the Sobolev space $H^{2}\left(\mathbb{R}^{n}\right)$ is $\mathbb{R}_{+}$, as for the free Laplacian. Our purpose is to study the resolvent $\left(H_{h}-z\right)^{-1}$, where $h>0$ is small and $z \in \mathbb{C}_{+}=\{z \in \mathbb{C}: \operatorname{Im} z>0\}$ is close to $E \in \mathbb{R}_{+}^{*}$. We prove some estimates for this resolvent uniform in the spectral parameter $z$, in order to obtain the limiting absorption principle and then existence and uniqueness of an outgoing solution $u_{h}$ for (1.1). We also control the dependence in $h$ of these estimates, which gives an a priori estimate for the size of $u_{h}$ when $h$ goes to 0 . Note that it is not clear that the resolvent is well-defined. More precisely the operator $H_{h}$ may have eigenvalues in a strip of size $O(h)$ around the real axis. Therefore we first have to prove that it cannot happen where we study the resolvent.

Let $\delta>\frac{1}{2}$. In the self-adjoint case $\left(V_{2}=0\right)$ it is known that there exist a neighborhood $I$ of $E>0, h_{0}>0$ and $c \geqslant 0$ such that

$$
\begin{equation*}
\left.\forall h \in] 0, h_{0}\right], \sup _{\substack{\operatorname{Re} z \in I \\ \operatorname{Im} z \neq 0}}\left\|\langle x\rangle^{-\delta}\left(H_{1}^{h}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h} \tag{1.2}
\end{equation*}
$$

if and only if the energy $E$ is non-trapping (see (1.5) below). Here we denote by $H_{1}^{h}$ the self-adjoint Schrödinger operator $-h^{2} \Delta+V_{1}(x)$, by $\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)$ the space of bounded operators on $L^{2}\left(\mathbb{R}^{n}\right)$, and $\langle x\rangle=(1+|x|)^{\frac{1}{2}}$. D. Robert and H . Tamura [23] proved that the non-trapping condition is sufficient and X.P. Wang [31] proved its necessity. In fact, if the non-trapping condition is not satisfied then the norm in (1.2) is at least of size $|\ln h| / h$ (see [4]).

For this result and all along this paper the potential $V_{1}$ is assumed to be of long range: it is smooth and there exist constants $\rho>0$ and $c_{\alpha} \geqslant 0$ for $\alpha \in \mathbb{N}^{n}$ such that

$$
\begin{equation*}
\forall \alpha \in \mathbb{N}^{n}, \forall x \in \mathbb{R}^{n}, \quad\left|\partial^{\alpha} V_{1}(x)\right| \leqslant c_{\alpha}\langle x\rangle^{-\rho-|\alpha|} \tag{1.3}
\end{equation*}
$$

Let $p:(x, \xi) \mapsto \xi^{2}+V_{1}(x)$ be the semiclassical symbol of $H_{1}^{h}$ on $\mathbb{R}^{2 n} \simeq T^{*} \mathbb{R}^{n}$ and $\phi^{t}$ the corresponding Hamiltonian flow. For any $w \in \mathbb{R}^{2 n}, t \mapsto \phi^{t}(w)=$ $(X(t, w), \Xi(t, w))$ is the solution of the system

$$
\left\{\begin{array}{l}
\partial_{t} X(t, w)=2 \Xi(t, w),  \tag{1.4}\\
\partial_{t} \Xi(t, w)=-\nabla V_{1}(X(t, w)), \\
\phi^{0}(w)=w .
\end{array}\right.
$$

We recall that $E>0$ is said to be non-trapping if

$$
\begin{equation*}
\forall w \in p^{-1}(\{E\}), \quad|X(t, w)| \xrightarrow[t \rightarrow \pm \infty]{ }+\infty \tag{1.5}
\end{equation*}
$$

For $I \subset \mathbb{R}$, we introduce the following subsets of $p^{-1}(I)$ :

$$
\begin{aligned}
\Omega_{b}^{ \pm}(I) & =\left\{w \in p^{-1}(I): \sup _{t \geqslant 0}|X( \pm t, w)|<\infty\right\} \\
\Omega_{b}(I) & =\Omega_{b}^{-}(I) \cap \Omega_{b}^{+}(I) \\
\Omega_{\infty}^{ \pm}(I) & =\left\{w \in p^{-1}(I):|X( \pm t, w)| \underset{t \rightarrow+\infty}{ }+\infty\right\}
\end{aligned}
$$

In [26] we considered the dissipative case $V_{2} \geqslant 0$. We proved (1.2) for $\operatorname{Im} z>0$ under a damping assumption on trapped trajectories:

$$
\begin{equation*}
\forall w \in \Omega_{b}(\{E\}), \exists T \in \mathbb{R}, \quad V_{2}(X(T, w))>0 \tag{1.6}
\end{equation*}
$$

Notice that this generalizes the usual non-trapping condition: when $V_{2}=0$ then (1.6) is equivalent to (1.5).

To prove this result we developed a dissipative version of Mourre's theory [20], which we applied to the dissipative Schrödinger operator. For this we constructed an escape function as introduced by Ch. Gérard and A. Martinez [14], using the damping assumption to allow trapped trajectories. Note that L. Aloui and M. Khenissi also proved some resolvent estimates for a dissipative Schrödinger operator in [1]. They needed a similar assumption but used a different approach (see below).

We know that assumption (1.6) is both sufficient and necessary in the dissipative setting. Our purpose is now to relax the dissipative condition, allowing negative values for the absorption index $V_{2}$. In this case, the damping assumption need reformulating. The condition we are going to use in this paper is the following:

$$
\begin{equation*}
\forall w \in \Omega_{b}(\{E\}), \exists T>0, \quad \int_{0}^{T} V_{2}(X(t, w)) d t>0 \tag{1.7}
\end{equation*}
$$

This condition is in particular satisfied if $V_{2} \geqslant 0$ and (1.6) holds. From this point of view, the results we are going to prove here are stronger than those given in the dissipative setting. With Assumption (1.7) and compactness of $\Omega_{b}(\{E\})$, we can prove that for all $w \in \Omega_{b}(\{E\})$ we have

$$
\liminf _{T \rightarrow+\infty} \frac{1}{T} \int_{0}^{T} V_{2}(X(t, w)) d t>0
$$

which means that we have a damping condition "on average in time" on bounded trajectories. See Proposition 2.3 for a more precise statement. The relation between asymptotic spectral properties of the non-selfadjoint Helmholtz equation and the average of the absorption index on classical trajectories has already been studied on compact manifolds by J. Sjöstrand in [30].

In this setting we cannot use the dissipative version of Mourre's commutators method. We use the same approach as in [1] instead. The idea is due to G. Lebeau [19] and N. Burq [6]. It is a contradiction argument. We consider a family of functions which denies the result, a semiclassical mesure associated to this family and finally we prove that this measure is both zero and non-zero. This idea was used in [6] for a general self-adjoint and compactly supported perturbation of the Laplacian. In [17], Th. Jecko used the argument to give a new proof of (1.2) with a real-valued potential. The motivation was to give a proof which could be applied to matrix-valued operators. To allow long range potentials, the author introduced a bounded "escape function" which we use here. The method was then used in [8] for a potential with Coulomb singularities and in $[18,13,10]$ for a matrix-valued operator.

Let us now state the main results about the resolvent. An important difference with the dissipative case is that we do not know if the resolvent is well-defined, even on the upper half-plane $\mathbb{C}_{+}$. However, our operator $H_{h}$ is a relatively compact perturbation of the Laplacian, so according to Weyl's Theorem [21, §XIII.4], its essential spectrum is $\mathbb{R}_{+}$and it can only have isolated eigenvalues on $\mathbb{C}_{+}$. So in the results we state now, we first claim that $H_{h}$ has no eigenvalue in the considered region and then give an estimate for the resolvent. The first theorem is about spectral parameters whose imaginary parts are bigger than $\beta h$ for some $\beta>0$ :

Theorem 1.1. - Suppose $V_{2}$ is smooth with bounded derivatives and $V_{2}(x) \rightarrow 0$ when $|x| \rightarrow+\infty$. Let $E>0$ be an energy which satisfies the damping assumption (1.7) and $\beta>0$. Then there exist a neighborhood $I$ of $E$, $h_{0}>0$ and $c \geqslant 0$ such that for $\left.\left.h \in\right] 0, h_{0}\right]$ and

$$
z \in \mathbb{C}_{I, h \beta}=\{z \in \mathbb{C}: \operatorname{Re} z \in I, \operatorname{Im} z \geqslant h \beta\}
$$

the operator $\left(H_{h}-z\right): H^{2}\left(\mathbb{R}^{n}\right) \rightarrow L^{2}\left(\mathbb{R}^{n}\right)$ has a bounded inverse on $L^{2}\left(\mathbb{R}^{n}\right)$ and

$$
\left\|\left(H_{h}-z\right)^{-1}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h} .
$$

Theorem 1.1 is obvious when $H_{h}$ is self-adjoint or at least dissipative. We can take $c=\beta^{-1}$ in these cases. The statement remains easy in the nondissipative setting when $\beta>\left\|V_{2}\right\|_{\infty}$ (see (3.3)), and the point of the theorem is to prove it for any $\beta>0$. Notice that we can deduce an estimate of size $O\left(h^{-3}\right)$ for $\left(H_{h}-z\right)^{-1}$ as an operator in $\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right), H^{2}\left(\mathbb{R}^{n}\right)\right)$.

In a second step we study the resolvent up to the real axis, generalizing (1.2). Now $V_{2}$ has to be of short range, which means that there exist $\rho>0$ and constants $c_{\alpha} \geqslant 0$ for $\alpha \in \mathbb{N}$ such that

$$
\begin{equation*}
\forall \alpha \in \mathbb{N}^{n}, \forall x \in \mathbb{R}^{n}, \quad\left|\partial^{\alpha} V_{2}(x)\right| \leqslant c_{\alpha}\langle x\rangle^{-1-\rho-|\alpha|} \tag{1.8}
\end{equation*}
$$

Theorem 1.2. - Assume that $V_{2}$ is of short range. Let $E>0$ satisfy the damping assumption (1.7) and $\delta>\frac{1}{2}$. Then there exist a neighborhood $I$ of $E$, $h_{0}>0$ and $c \geqslant 0$ such that for $\left.\left.h \in\right] 0, h_{0}\right]$ and

$$
z \in \mathbb{C}_{I,+}=\{z \in \mathbb{C}: \operatorname{Re} z \in I, \operatorname{Im} z>0\}
$$

the operator $\left(H_{h}-z\right)$ has a bounded inverse on $L^{2}\left(\mathbb{R}^{n}\right)$ and

$$
\left\|\langle x\rangle^{-\delta}\left(H_{h}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h} .
$$

The proof of this theorem is inspired from [17]. In particular we use a bounded escape function at infinity to prove that the semiclassical measure we study is non-zero. But contrary to the selfadjoint case (this could also be done in the dissipative case) we cannot use this escape function to prove that this measure is supported in a compact subset of $\mathbb{R}^{2 n}$. We use instead the estimate for the outgoing solution of the Helmholtz equation in the incoming region proved in [24].

As in the dissipative case we have only given a result on the upper half-plane. Here we have no assumption about the sign of $V_{2}$, but there still is a damping condition in (1.7). The difference with the dissipative context is that we recover a symmetric situation under the stronger non-trapping condition (1.5), so that the result we have proved for $\operatorname{Im} z>0$ now holds when $\operatorname{Im} z<0$ :

Corollary 1.3. - Assume that $V_{2}$ is of short range. Let $E>0$ be a nontrapping energy and $\delta>\frac{1}{2}$. Then there exist a neighborhood $I$ of $E, h_{0}>0$ and $c \geqslant 0$ such that for $\left.h \in] 0, h_{0}\right], \operatorname{Re} z \in I$ and $\operatorname{Im} z \neq 0$ the operator $\left(H_{h}-z\right)$ has a bounded inverse on $L^{2}\left(\mathbb{R}^{n}\right)$ and

$$
\left\|\langle x\rangle^{-\delta}\left(H_{h}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h} .
$$

Once we have the uniform resolvent estimates, we can prove the limiting absorption principle. This question has been studied for a long range selfadjoint Schrödinger operator in [15] and [29]. It is proved that the equation $(H-E) u=f$ has a unique outgoing solution $u \in H_{\mathrm{loc}}^{2}\left(\mathbb{R}^{n}\right) \cap L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ when $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$, and this solution is given by the limit in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ of $(H-z)^{-1} f$ when $z \in \mathbb{C}_{+}$goes to $E$ (we do not work in the semiclassical limit here, so we only consider the case $h=1$ ). Here we denote by $L^{2, \delta}\left(\mathbb{R}^{n}\right)$ the weighted space $L^{2}\left(\langle x\rangle^{2 \delta} d x\right)$. An outgoing solution is a solution which satisfies a radiation condition of Sommerfeld type at infinity (see Definition 6.2). The strategy is to
prove first uniqueness of an outgoing solution. This is used to prove resolvent estimates, and then we can get the limiting absorption principle.

The result has been extended to the non-selfadjoint case in [28]. Y. Saito proves that when the potential has a short range imaginary part, the result can be extended where we have uniqueness of the outgoing solution, which no longer holds for any $E>0$.

Here we use the fact that we already have uniform resolvent estimates to obtain uniqueness of the outgoing solution and then the limiting absorption principle. With Theorem 1.2 we obtain the following result:

Theorem 1.4. - Let $V_{2}, E$ and $\delta$ be as in Theorem 1.2. Then there exist a neighborhood $I$ of $E$ and $h_{0}>0$ such that for any $\left.\left.h \in\right] 0, h_{0}\right], \lambda \in I$ and $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ the limit

$$
\lim _{\substack{z \overrightarrow{\mathbb{C}}_{I,+}^{\lambda} \\ z}}\left(H_{h}-z\right)^{-1} f
$$

exists in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ and defines the unique outgoing solution for the equation ( $\left.H_{h}-\lambda\right) u=f$ (see Definition 6.2).

Note that when the dissipative part $V_{2}$ is non-negative we can prove directly uniqueness of the outgoing solution and hence we can proceed as in the selfadjoint case, and $V_{2}$ can be of long range in this case (see Proposition 6.9).

Now that the outgoing solution

$$
u_{h}=\left(H_{h}-(E+i 0)\right)^{-1} f_{h}
$$

is well-defined in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ for any $\left.\left.h \in\right] 0, h_{0}\right]\left(h_{0}>0\right.$ being given by Theorem 1.4) and $f_{h} \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$, we can study its semiclassical measures when the source term $f_{h}$ concentrates on a submanifold of $\mathbb{R}^{n}$. We recall that a measure $\mu$ on the phase space $T^{*} \mathbb{R}^{n} \simeq \mathbb{R}^{2 n}$ is said to be a semiclassical measure for the family $\left(u_{h}\right)_{\left.h \in] 0, h_{0}\right]}$ if there exists a sequence $\left.\left.\left(h_{m}\right)_{m \in \mathbb{N}} \in\right] 0, h_{0}\right]^{\mathbb{N}}$ such that

$$
h_{m} \xrightarrow[m \rightarrow \infty]{ } 0 \quad \text { and } \quad \forall q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right), \quad\left\langle\mathrm{Op}_{h_{m}}^{w} u_{h_{m}}, u_{h_{m}}\right\rangle \underset{m \rightarrow \infty}{ } \int_{\mathbb{R}^{2 n}} q d \mu
$$

Here $\mathrm{Op}_{h}^{w}(q)$ denotes the Weyl $h$-quantization of the symbol $q$ :

$$
\mathrm{Op}_{h}^{w}(q) u(x)=\frac{1}{(2 \pi h)^{n}} \int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}} e^{\frac{i}{h}\langle x-y, \xi\rangle} q\left(\frac{x+y}{2}, \xi\right) u(y) d y d \xi
$$

We will also use the standard quantization:

$$
\mathrm{Op}_{h}(q) u(x)=\frac{1}{(2 \pi h)^{n}} \int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}} e^{\frac{i}{h}\langle x-y, \xi\rangle} q(x, \xi) u(y) d y d \xi
$$

The first paper about this problem is [2], where $f_{h}$ is assumed to concentrate on $\Gamma=\{0\}$ as $h$ goes to 0 (see also [7]). This was extended in [9, 33] to the case where $f_{h}$ concentrates on an affine subspace $\Gamma$ in $\mathbb{R}^{n}$. J.-F. Bony gave
in [3] another proof for the case $\Gamma=\{0\}$ using different assumptions and above all a different approach. We used this point of view in [27] to deal with the case where the absorption index is non-constant and $\Gamma$ is any bounded submanifold of dimension $d \in \llbracket 0, n-1 \rrbracket$ in $\mathbb{R}^{n}$. Trapped trajectories of energy $E$ for the classical flow were allowed under assumption (1.6). Our purpose here is to check that this result still holds -under Hypothesis (1.7)- when $V_{2}$ takes negative values. The proof is actually approximately the same. We will have to prove the estimates in the incoming region in this non-dissipative setting (see Theorem 4.1) and be careful with the fact that the semi-group generated by $H_{h}$ is no longer a contraction semi-group.

Throughout this paper we denote by $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ the set of smooth and compactly supported functions on $\mathbb{R}^{n}, C_{b}^{\infty}\left(\mathbb{R}^{n}\right)$ is the set of smooth functions whose derivatives are bounded and $\gamma\left(\mathbb{R}^{n}\right)$ is the Schwartz space of rapidly decaying functions. For $\delta \in \mathbb{R}$, we denote by $\delta\left(\langle x\rangle^{\delta}\right)$ the set of symbols $a \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ such that

$$
\forall \alpha, \beta \in \mathbb{N}^{n}, \exists c_{\alpha, \beta} \geqslant 0, \forall(x, \xi) \in \mathbb{R}^{2 n}, \quad\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi)\right| \leqslant c_{\alpha, \beta}\langle x\rangle^{\delta}
$$

We also denote by $\ell_{\delta}\left(\mathbb{R}^{2 n}\right)$ the set of symbols $a \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ such that

$$
\forall \alpha, \beta \in \mathbb{N}^{n}, \exists c_{\alpha, \beta} \geqslant 0, \forall(x, \xi) \in \mathbb{R}^{2 n}, \quad\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi)\right| \leqslant c_{\alpha, \beta}\langle x\rangle^{\delta-|\alpha|}
$$

For $R>0$ we denote by $B_{R}$ the open ball of radius $R$ in $\mathbb{R}^{n}$, by $B_{R}^{c}$ its complement in $\mathbb{R}^{n}$, and by $\mathbb{S}_{R}$ the sphere of radius $R$, endowed with the Lebesgue measure. We also set $B_{x}(R)=\{(x, \xi):|x|<R\} \subset \mathbb{R}^{2 n}$.

In Section 2 we recall some properties of the flow $\phi^{t}$ defined by (1.4) and discuss assumption (1.7). Section 3 is devoted to the proof of Theorem 1.1. Before giving a proof of Theorem 1.2 in Section 5, we state a non-dissipative version for the estimate in the incoming region of the outgoing solution for the Helmholtz equation (see Section 4). With the uniform resolvent estimates, we prove the limiting absorption principle in Section 6. We finally show in Section 7 that the result known about the semiclassical measure for the outgoing solution of (1.1) when the source term concentrates on a bounded submanifold of $\mathbb{R}^{n}$ remains valid in our non-dissipative setting.

## 2. More about classical dynamics and the weak damping assumption

The good properties of the flow at infinity come from the fact that for any $\nu>0$ there exists $\mathcal{R} \geqslant 0$ such that

$$
\begin{equation*}
\forall x \in \mathbb{R}^{n}, \quad|x| \geqslant \mathscr{R} \Longrightarrow\left|V_{1}(x)\right|+|x|\left|\nabla V_{1}(x)\right|<\nu \tag{2.1}
\end{equation*}
$$

This implies that far from the origin the refraction index has low influence on the flow and hence classical trajectories behave as in the free case.

In particular for $J=\left[E_{1}, E_{2}\right] \subset \mathbb{R}_{+}^{*}$ and $\mathscr{R}$ so large that (2.1) holds for some $\nu \in] 0, \frac{2 E_{1}}{3}\left[\right.$, we have $\partial_{t}^{2}|X(t, w)|^{2} \geqslant 8 E_{1}-12 \nu>0$ if $p(w) \in J$ and $|X(t, w)| \geqslant \mathcal{R}$. As a consequence a classical trajectory of energy $E \in J$ which leaves $B_{x}(\mathscr{R})$ cannot come back and goes to infinity. This implies that

$$
p^{-1}(J)=\Omega_{b}^{+}(J) \sqcup \Omega_{\infty}^{+}(J)=\Omega_{b}^{-}(J) \sqcup \Omega_{\infty}^{-}(J)=\Omega_{b}(J) \cup \Omega_{\infty}^{+}(J) \cup \Omega_{\infty}^{-}(J)
$$

Moreover $\Omega_{b}^{ \pm}(J)$ is closed in $\mathbb{R}^{2 n}$ and $\Omega_{b}(J) \subset B_{x}(\mathscr{R})$ is compact. If $\mathcal{B}_{ \pm}$is a bounded subset of $\Omega_{b}^{ \pm}(J)$, we choose $\mathscr{R}$ such that (2.1) holds for $\nu=\frac{E_{1}}{3}$ and $\mathcal{B}_{ \pm} \subset B_{x}(\mathcal{R})$ to prove that the set

$$
\left\{\phi^{ \pm t}(w), t \geqslant 0, w \in \mathcal{B}_{ \pm}\right\}
$$

is bounded in $\mathbb{R}^{2 n}$.
For $R \geqslant 0, d \geqslant 0$ and $\sigma \in[-1,1]$ we denote by

$$
\mathcal{Z}_{ \pm}(R, d, \sigma)=\left\{(x, \xi) \in \mathbb{R}^{2 n}:|x| \geqslant R,|\xi| \geqslant d \text { and } \pm\langle x, \xi\rangle \geqslant \pm \sigma|x||\xi|\right\}
$$

the incoming and outgoing regions. The proposition we prove now ensures that a trajectory starting outside some incoming region and far enough from the origin stays away from the influence of $V_{1}$ and hence has a "nice" behavior:

Proposition 2.1. - Let $E_{2} \geqslant E_{1}>0, J \subset\left[E_{1}, E_{2}\right]$ and $\sigma \in[0,1[$ be such that $\sigma^{2} E_{2}<E_{1}$. Then there exist $\mathcal{R}>0$ and $c_{0}>0$ such that

$$
\forall t \geqslant 0, \forall(x, \xi) \in \mathcal{Z}_{ \pm}(\mathcal{R}, 0, \mp \sigma) \cap p^{-1}(J), \quad|X( \pm t, x, \xi)| \geqslant c_{0}(t+|x|)
$$

Proof. - Let $\tilde{\mathscr{R}}$ be such that (2.1) holds for $\nu \in] 0, \frac{2 E_{1}}{3}[$ so small that

$$
\tilde{\nu}:=1-\sigma^{2} \frac{E_{2}+\nu}{E_{1}-\frac{3}{2} \nu}>0 .
$$

Let $\mathscr{R}$ be greater than $\tilde{\mathscr{R}} / \sqrt{\tilde{\nu}}$ and $(x, \xi) \in \mathcal{Z}_{ \pm}(\mathscr{R}, 0, \mp \sigma) \cap p^{-1}(J)$. Suppose there exists $t \geqslant 0$ such that $|X( \pm t, x, \xi)|<\tilde{\mathscr{R}}$ and let

$$
t_{0}=\inf \{t \geqslant 0,|X( \pm t, x, \xi)|<\tilde{\mathscr{R}}\} .
$$

Let $E_{3}=E_{1}-\frac{3}{2} \nu>0$. For $t \in\left[0, t_{0}\right]$ we have $\partial_{t}^{2}|X( \pm t, x, \xi)|^{2} \geqslant 8 E_{3}$ and hence

$$
|X( \pm t, x, \xi)|^{2} \geqslant|x|^{2}-4 t \sigma|x||\xi|+4 E_{3} t^{2} \geqslant|x|^{2}\left(1-\frac{\sigma^{2}|\xi|^{2}}{E_{3}}\right) \geqslant|x|^{2} \tilde{\nu}>\tilde{\mathscr{R}}^{2}
$$

This gives a contradiction when $t=t_{0}$, and proves that these inequalities actually hold for all $t \geqslant 0$. We also have

$$
\forall t \geqslant 2 \sigma|x||\xi| E_{3}, \quad|X( \pm t, x, \xi)|^{2}-|x|^{2}-2 E_{3} t^{2} \geqslant 2 t\left(E_{3} t-2 \sigma|x||\xi|\right) \geqslant 0
$$

which concludes the proof.

We now discuss assumption (1.7). We still denote by $V_{2}$ the function $(x, \xi) \mapsto V_{2}(x)$ on $\mathbb{R}^{2 n}$. We first remark that since $(x, \xi)$ and $(x,-\xi)$ are simultaneously in $\Omega_{b}(\{E\})$, Assumption (1.7) could be equivalently formulated looking at trajectories in the past:

$$
\begin{equation*}
(1.7) \quad \Longleftrightarrow \quad \forall w \in \Omega_{b}(\{E\}), \exists T \in \mathbb{R}, \quad V_{2}(X(-T, w))>0 \tag{2.2}
\end{equation*}
$$

Using compactness of $\Omega_{b}([E / 2,2 E])$, we can also check that assumption (1.7) is an open property:

Proposition 2.2. - If assumption (1.7) holds for some $E>0$, then it also holds for any $\lambda$ in some neighborhood of $E$.

We now prove that assumption (1.7) ensures that the absorption is positive on trapped trajectories "on average" in time, as mentioned in introduction.

Proposition 2.3. - Let $J \subset \mathbb{R}_{+}^{*}$ be such that assumption (1.7) holds for any $\lambda \in J$. Then for all compact $K \subset \Omega_{b}^{ \pm}(J)$ there exist $c_{0}>0$ and $C \geqslant 0$ such that

$$
\forall t \geqslant 0, \forall w \in K, \quad \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant c_{0} t-C .
$$

Remark 2.4. - This proposition implies that

$$
\forall w \in \Omega_{b}^{ \pm}(J), \quad \liminf _{t \rightarrow+\infty} \frac{1}{t} \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s>0
$$

The claim that the average in time of the absorption is positive on trapped trajectories becomes clear on periodic trajectories. If $w \in p^{-1}(J)$ and $T>0$ are such that $\phi^{T}(w)=w$, then

$$
\int_{0}^{T}\left(V_{2} \circ \phi^{t}\right)(w) d t>0
$$

To prove this we only have to apply (1.7) to $w_{0}=\phi^{t_{0}}(w) \in \Omega_{b}(\{E\})$, where $t_{0} \in[0, T]$ is the time for which $t \mapsto \int_{0}^{t} V_{2}(X(s, w)) d s$ reaches its maximum.

Proof of Proposition 2.3. - 1. Since $K$ is compact we can assume without loss of generality that $J$ is a compact subset of $\mathbb{R}_{+}^{*}$. Let $w \in \Omega_{b}(J)$. By assumption there exist $T, \gamma>0$ such that

$$
\int_{0}^{T}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant 2 \gamma
$$

Since the left-hand side is a continuous function of $w$, we can find a neighborhood $V_{w}$ of $w$ in $\mathbb{R}^{2 n}$ such that for all $v \in V_{w}$ we have

$$
\int_{0}^{T}\left(V_{2} \circ \phi^{ \pm s}\right)(v) d s \geqslant \gamma
$$

As $\Omega_{b}(J)$ is compact, it is covered by a finite number of such sets $V_{w}$. Hence we can find $T_{1}, \gamma_{1}>0$ such that

$$
\forall w \in \Omega_{b}(J), \exists t \in\left[0, T_{1}\right], \quad \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant \gamma_{1}
$$

2 . Let $\nu=\frac{2+T_{1} m_{-}}{\gamma_{1}}$ and $T_{2}=T_{1}(1+\nu)$, where $m_{-}:=-\inf _{\mathbb{R}^{n}} V_{2}$. Let $w \in \Omega_{b}(J)$. We set $t_{0}=0$ and for all $k \in \mathbb{N}$ we consider by induction $\left.\left.t_{k+1} \in\right] t_{k}, t_{k}+T_{1}\right]$ such that

$$
\int_{t_{k}}^{t_{k+1}}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant \gamma_{1}
$$

We necessarily have $t_{k+1} \geqslant t_{k}+\gamma_{1} /\left\|V_{2}\right\|_{\infty}$ for all $k \in \mathbb{N}$, and hence $t_{k} \rightarrow+\infty$ (if $V_{2}=0$ the statement of the proposition is empty). In particular any $t>0$ belongs to $\left.] t_{k}, t_{k}+T_{1}\right]$ for some $k \in \mathbb{N}$. Let $t \geqslant T_{2}$ and $N \in \mathbb{N}$ such that $\left.t \in] t_{N}, t_{N}+T_{1}\right]$. We have $N \geqslant \nu$ and hence

$$
\begin{aligned}
\int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s & \geqslant \sum_{k=0}^{N-1} \int_{t_{k}}^{t_{k+1}}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s+\int_{t_{N}}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \\
& \geqslant \nu \gamma_{1}-T_{1} m_{-}
\end{aligned}
$$

This proves that

$$
\forall w \in \Omega_{b}(J), \forall t \geqslant T_{2}, \quad \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant 2
$$

3. By continuity, there exists a neighborhood $\mathcal{U}$ of $\Omega_{b}(J)$ such that

$$
\forall w \in \mathcal{U}, \quad \int_{0}^{T_{2}}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant 1 .
$$

4. Now let $K$ be a compact subset of $\Omega_{b}^{ \pm}(J)$. We prove by contradiction that there exists $T_{K} \geqslant 0$ such that $\phi^{ \pm t}(w) \in \mathcal{U}$ for all $w \in K$ and $t \geqslant T_{K}$. If it is not the case we can find sequences $\left(t_{m}\right)_{m \in \mathbb{N}}$ and $\left(w_{m}\right)_{m \in \mathbb{N}}$ with $w_{m} \in K$ and $t_{m} \rightarrow+\infty$ such that $\phi^{ \pm t_{m}}\left(w_{m}\right) \notin \mathcal{U}$. Each term of the sequence $\left(\phi^{ \pm t_{m}}(w)\right)_{m \in \mathbb{N}}$ belongs to the bounded set $\bigcup_{t \geqslant 0} \phi^{ \pm t}(K)$, so after extracting a subsequence if necessary, we can assume that it converges to

$$
w_{\infty} \in \Omega_{b}^{ \pm}(J) \backslash \mathcal{U} \subset \Omega_{\infty}^{\mp}(\mathbb{R})
$$

Let $\mathcal{R}$ be such that (2.1) holds for $\nu=(\inf J) / 3>0$ and $K \cup\left\{w_{\infty}\right\} \subset B_{x}(\mathcal{R})$. There exists $T_{\infty} \geqslant 0$ such that $\left|X\left(\mp T_{\infty}, w_{\infty}\right)\right|>2 \mathcal{R}$. By continuity and properties of $\mathcal{R}$, we can find a neighborhood $\mathcal{V} \subset B_{x}(\mathscr{R})$ of $w_{\infty}$ such that $|X(\mp t, v)|>2 \mathcal{R}$ for all $v \in V$ and $t \geqslant T_{\infty}$. Hence for large $m$ we have $\phi^{\mp t_{m}}(\mathcal{V}) \cap K=\varnothing$, and in particular $\phi^{ \pm t_{m}}\left(w_{m}\right) \notin \mathcal{V}$. This gives a contradiction.
5. Let $w \in K, t \geqslant T_{K}$ and $N$ be the integer part of $\frac{t-T_{K}}{T_{2}}$. We have

$$
\begin{aligned}
& \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \\
& \geqslant
\end{aligned} \begin{aligned}
& T_{0} \\
&\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s+\sum_{k=0}^{N-1} \int_{T_{K}+k T_{2}}^{T_{K}+(k+1) T_{2}}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \\
& \geqslant \int_{T_{K}+N T_{2}}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \\
& \geqslant-T_{K} m_{-}+N-T_{2} m_{-} \\
& m_{-}+\frac{t-T_{K}}{T_{2}}-1-T_{2} m_{-}
\end{aligned}
$$

Since this integral is not less than $-T_{K} m_{-}$when $t \in\left[0, T_{K}\right]$, this gives the result with $c_{0}=1 / T_{2}$ and $C=1+\left(T_{K}+T_{2}\right) m_{-}+T_{K} / T_{2}$.

We are going to use in Section 7 a more precise result:
Proposition 2.5. - Let $\mathcal{R}>0$ and $J \subset \mathbb{R}_{+}^{*}$ such that assumption (1.7) holds for any $\lambda \in J$. Then for any compact subset $\tilde{K}$ of $p^{-1}(J)$ there exist $c_{0}, C>0$ such that

$$
\forall t \geqslant 0, \forall w \in \tilde{K}, \quad \int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant c_{0} t-C \quad \text { or } \quad|X( \pm t, w)| \geqslant \mathcal{R}
$$

If $K \subset \Omega_{b}^{ \pm}(J)$ this comes from Proposition 2.3, and if $K \subset \Omega_{\infty}^{ \pm}(J)$, the second conclusion holds for $t$ large enough, uniformly in $w \in K$, and the first conclusion is always true for finite times. The problem therefore comes from the boundary between $\Omega_{b}^{ \pm}(J)$ and $\Omega_{\infty}^{ \pm}(J)$.

Proof. - As above we may assume that $J$ is compact. Since the conclusion is stronger if $\mathcal{R}$ is taken larger, we may assume that (2.1) holds for $\nu=(\inf J) / 3$. Let $K=\tilde{K} \cap \Omega_{b}^{ \pm}(J)$. $K$ is a compact subset of $\Omega_{b}^{ \pm}(J)$. We use the notation introduced in the proof of Proposition 2.3. We know that there exists $T_{K} \geqslant 0$ such that $\phi^{ \pm t}(w) \in \mathcal{U}$ for all $w \in K$ and $t \geqslant T_{K}$. By continuity of the Hamiltonian flow, there exists a neighborhood $V$ of $K$ in $\tilde{K}$ such that $\phi^{ \pm T_{K}}(w) \in \mathcal{U}$ for all $w \in \mathcal{V}$. We now prove that there exists $T_{\mathscr{R}} \geqslant 0$ such that for all $w \in \mathcal{V} \backslash K \subset \Omega_{\infty}^{ \pm}(J)$ we can find $\tau_{w} \geqslant T_{K}$ which satisfies:
(2.3) $\forall t \in\left[T_{K}, \tau_{w}\right], \quad \phi^{ \pm t}(w) \in \mathcal{U} \quad$ and $\quad \forall t \geqslant \tau_{w}+T_{\mathscr{R}}, \quad|X( \pm t, w)| \geqslant \mathcal{R}$.

This means that even if we cannot say when a trajectory coming from $\mathcal{V} \backslash K$ will leave $B_{x}(\mathscr{R})$, we control the time it can stay in $B_{x}(\mathscr{R}) \backslash \mathcal{U}$. Assume that we
cannot find $T_{\mathscr{R}}$ such that (2.3) holds. Then there exists a sequence $\left(w_{m}\right)_{m \in \mathbb{N}}$ of elements in $\mathcal{V} \backslash K$ and times $t_{m} \geqslant T_{K}, \theta_{m} \geqslant m$ for $m \in \mathbb{N}$ such that

$$
\phi^{ \pm t_{m}}\left(w_{m}\right) \notin \mathcal{U} \quad \text { and } \quad\left|X\left( \pm\left(t_{m}+\theta_{m}\right), w_{m}\right)\right| \leqslant \mathcal{R} .
$$

After extracting a subsequence if necessary, we may assume that $w_{m}$ converges to $w_{\infty} \in \tilde{K}$. If $w_{\infty} \in \Omega_{\infty}^{ \pm}(J)$ then there exist $T_{\infty} \geqslant 0$ and a neighborhood $\mathcal{W}$ of $w_{\infty}$ such that $|X( \pm t, v)| \geqslant \mathcal{R}$ for all $t \geqslant T_{\infty}$ and $v \in \mathcal{W}$, which is impossible. This means that the limit $w_{\infty}$ actually belongs to $K$. If the sequence $\left(t_{m}\right)_{m \in \mathbb{N}}$ is bounded we can assume, after extraction, that $t_{m} \rightarrow t_{\infty} \geqslant T_{K}$, which cannot be true since we would have $\phi^{ \pm t_{m}}\left(w_{m}\right) \rightarrow \phi^{ \pm t}\left(w_{\infty}\right) \in \mathcal{U}$ and hence $\phi^{ \pm t_{m}}\left(w_{m}\right) \in \mathcal{U}$ for $m$ large enough. Extracting again a subsequence, we can assume that $t_{m} \rightarrow+\infty$. Let $v_{m}=\phi^{ \pm t_{m}}\left(w_{m}\right)$. The sequence $\left(v_{m}\right)_{m \in \mathbb{N}}$ is bounded so without loss of generality we can assume that it converges to some $v_{\infty} \in p^{-1}(J)$. Since $t_{m}, \theta_{m} \rightarrow \infty$ and the sequences $\left(\phi^{\mp t_{m}}\left(v_{m}\right)\right)_{m \in \mathbb{N}}$ and $\left(\phi^{ \pm \theta_{m}}\left(v_{m}\right)\right)_{m \in \mathbb{N}}$ are bounded, we obtain as before that $v_{\infty} \in \Omega_{b}(J)$, which gives a contradiction and hence proves (2.3). The complement $\tilde{K} \backslash \mathcal{V}$ is a compact subset of $\Omega_{\infty}^{ \pm}(J)$. Choosing $T_{\mathscr{R}}$ larger if necessary, we can assume that

$$
\forall w \in \tilde{K} \backslash \mathcal{V}, \forall t \geqslant T_{\mathscr{R}}, \quad|X( \pm t, w)| \geqslant \mathscr{R}
$$

As a consequence, given $w \in \tilde{K}$ and $t \geqslant 0$ such that $|X( \pm t, w)| \leqslant \mathcal{R}$, we have $\phi^{ \pm s}(w) \in \mathcal{U}$ for all $s \in\left[T_{K}, t-T_{\mathcal{R}}\right]$ and hence, as we did in the proof of Proposition 2.3, we obtain

$$
\int_{0}^{t}\left(V_{2} \circ \phi^{ \pm s}\right)(w) d s \geqslant-m_{-} T_{K}+\frac{t-T_{K}-T_{\mathscr{R}}}{T_{2}}-1-T_{2} m_{-}-T_{\mathscr{R}} m_{-}
$$

and conclude the proof.

## 3. Resolvent at distance of order $h$ from the real axis

In this section we give some general properties about the semiclassical measures we consider and prove Theorem 1.1.

Proposition 3.1. - Assume that $V_{2} \in C_{b}^{\infty}\left(\mathbb{R}^{n}\right)$. Let $\left(z_{m}\right)_{m \in \mathbb{N}} \in \mathbb{C}^{\mathbb{N}}$ and $\left.\left.\left(h_{m}\right)_{m \in \mathbb{N}} \in\right] 0,1\right]^{\mathbb{N}}$ be sequences such that
$h_{m} \xrightarrow[m \rightarrow \infty]{ } 0, \quad \lambda_{m}:=\operatorname{Re} z_{m} \xrightarrow[m \rightarrow \infty]{ } E>0 \quad$ and $\quad \beta_{m}:=h_{m}^{-1} \operatorname{Im} z_{m} \xrightarrow[m \rightarrow \infty]{ } \beta \in \mathbb{R}$.
Consider $\delta \geqslant 0$ and a sequence $\left(v_{m}\right)_{m \in \mathbb{N}} \in H^{2}\left(\mathbb{R}^{n}\right)^{\mathbb{N}}$ such that

$$
\left\|v_{m}\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}=1, \quad\left\|\left(H_{h_{m}}-z_{m}\right) v_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}=\underset{m \rightarrow \infty}{o}\left(h_{m}\right),
$$

and

$$
\begin{equation*}
\forall q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right), \quad\left\langle\mathrm{Op}_{h_{m}}^{w}(q) v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \underset{m \rightarrow \infty}{ } \int_{\mathbb{R}^{2 n}} q d \mu \tag{3.1}
\end{equation*}
$$

for some (non-negative) measure $\mu$ on $\mathbb{R}^{2 n}$. Then we have the following two properties.
(i) If $q \in \&\left(\langle x\rangle^{-2 \delta}\right)$ is supported outside $p^{-1}(J)$ for some neighborhood $J$ of $E$ we have

$$
\left\langle\mathrm{Op}_{h_{m}}^{w}(q) v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow \infty]{ } 0
$$

In particular $\mu$ is supported on $p^{-1}(\{E\})$ and for $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ we have

$$
\left\langle\chi v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow \infty]{ } \int_{\mathbb{R}^{2 n}} \chi(x) d \mu(x, \xi)
$$

(ii) For $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ and $t \geqslant 0$ we have

$$
\begin{equation*}
\int_{\mathbb{R}^{2 n}} q d \mu=\int_{\mathbb{R}^{2 n}}\left(q \circ \phi^{t}\right) \exp \left(-2 \int_{0}^{t}\left(V_{2}+\beta\right) \circ \phi^{t-s} d s\right) d \mu \tag{3.2}
\end{equation*}
$$

Proof. - (i) For $m$ large enough (such that $\lambda_{m} \in{ }^{\circ}$ ) we set:

$$
a_{m}(x, \xi)=\frac{q(x, \xi)\langle x\rangle^{2 \delta}}{p(x, \xi)-z_{m}}
$$

Since $q$ vanishes on $p^{-1}(J)$, we have $a_{m} \in \phi\left(\langle\xi\rangle^{-2}\right)$ uniformly in $m$. We can write

$$
\begin{aligned}
\left\langle\mathrm{Op}_{h_{m}}^{w}(q) v_{m}, v_{m}\right\rangle & \leqslant\left\|\mathrm{Op}_{h_{m}}^{w}(q) v_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)} \\
& \leqslant\left\|\mathrm{Op}_{h_{m}}^{w}\left(a_{m}\right)\left(H_{h_{m}}^{1}-z_{m}\right) v_{m}\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\underset{m \rightarrow \infty}{o}\left(h_{m}\right) \\
& \leqslant\left\|\mathrm{Op}_{h_{m}}^{w}\left(a_{m}\right)\left(H_{h_{m}}-z_{m}\right) v_{m}\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\underset{m \rightarrow \infty}{o}\left(h_{m}\right) \\
& \xrightarrow[m \rightarrow \infty]{o} 0
\end{aligned}
$$

which proves the first assertion. Applied with $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$, this proves that $\mu$ is supported on $p^{-1}(\{E\})$. Now let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and consider $\tilde{\chi} \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\mathbb{R}^{n} \times \operatorname{supp}(1-\tilde{\chi})$ does not intersect $p^{-1}(\{E\})$. Then

$$
\left\langle\mathrm{Op}_{h_{m}}^{w}(\chi(x)(1-\tilde{\chi}(\xi))) v_{m}, v_{m}\right\rangle \underset{m \rightarrow \infty}{ } 0
$$

and hence

$$
\begin{aligned}
\lim _{m \rightarrow \infty}\left\langle\chi(x) v_{m}, v_{m}\right\rangle & =\lim _{m \rightarrow \infty}\left\langle\mathrm{Op}_{h}^{w}(\chi(x) \tilde{\chi}(\xi)) v_{m}, v_{m}\right\rangle=\int_{\mathbb{R}^{2 n}} \chi(x) \tilde{\chi}(\xi) d \mu(x, \xi) \\
& =\int_{\mathbb{R}^{2 n}} \chi(x) d \mu(x, \xi)
\end{aligned}
$$

(ii) Let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ and $t \geqslant 0$. For $\tau \in[0, t]$ and $w \in \mathbb{R}^{2 n}$ we set

$$
q(\tau, w)=q\left(\phi^{t-\tau}(w)\right) \exp \left(-2 \int_{\tau}^{t}\left(V_{2}+\beta\right)\left(\phi^{s-\tau}(w)\right) d s\right)
$$

Since $\bigcup_{\tau \in[0, t]} \operatorname{supp}\left(q \circ \phi^{t-\tau}\right)$ is bounded in $\mathbb{R}^{2 n}$ we can use differentiation under the integral sign:

$$
\begin{aligned}
\frac{d}{d \tau} & \int_{\mathbb{R}^{2 n}} q(\tau) d \mu=\int_{\mathbb{R}^{2 n}} \frac{d}{d \tau} q(\tau) d \mu=\int_{\mathbb{R}^{2 n}}\left(2\left(V_{2}+\beta\right) q(\tau)-\{p, q(\tau)\}\right) d \mu \\
& =\lim _{m \rightarrow \infty}\left\langle\operatorname{Op}_{h_{m}}^{w}\left(2\left(V_{2}+\beta\right) q(\tau)-\{p, q(\tau)\}\right) v_{m}, v_{m}\right\rangle \\
& =\lim _{m \rightarrow \infty}\left\langle\left(2\left(V_{2}+\beta_{m}\right) \mathrm{Op}_{h_{m}}^{w}(q(\tau))-\frac{i}{h_{m}}\left[H_{h_{m}}^{1}, \mathrm{Op}_{h_{m}}^{w}(q(\tau))\right]\right) v_{m}, v_{m}\right\rangle \\
& =\lim _{m \rightarrow \infty} \frac{i}{h_{m}}\left\langle\left(\mathrm{Op}_{h_{m}}^{w}(q(\tau))\left(H_{h_{m}}-z_{m}\right)-\left(H_{h_{m}}-z_{m}\right)^{*} \mathrm{Op}_{h_{m}}^{w}(q(\tau))\right) v_{m}, v_{m}\right\rangle \\
& =0 .
\end{aligned}
$$

This gives statement (ii). Here we do not have to worry about decay properties of $v_{m}$ since we only work with compactly supported symbols.

We now turn to the proof of Theorem 1.1. We first remark that it is easy when $\beta>m_{-}=-\inf _{\mathbb{R}^{n}} V_{2}$ since

$$
H_{h}-z=\left(H_{h}-i h m_{-}\right)-\left(z-i h m_{-}\right)
$$

and $H_{h}-i h m_{-}$is maximal dissipative. This proves that if $\operatorname{Im} z>h m_{-}$the resolvent $\left(H_{h}-z\right)^{-1}$ is well-defined and

$$
\begin{equation*}
\left\|\left(H_{h}-z\right)^{-1}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{1}{\operatorname{Im} z-h m_{-}} . \tag{3.3}
\end{equation*}
$$

As said in introduction we proceed by contradiction to prove the general case. So we assume that we can find sequences $\left(v_{m}\right)_{m \in \mathbb{N}} \in H^{2}\left(\mathbb{R}^{n}\right)^{\mathbb{N}},\left(z_{m}\right)_{m \in \mathbb{N}} \in \mathbb{C}^{\mathbb{N}}$ and $\left.\left.\left(h_{m}\right)_{m \in \mathbb{N}} \in\right] 0,1\right]^{\mathbb{N}}$ such that

$$
h_{m} \rightarrow 0, \quad \lambda_{m}:=\operatorname{Re} z_{m} \rightarrow E, \quad \beta_{m}:=h_{m}^{-1} \operatorname{Im} z_{m} \geqslant \beta, \quad\left\|v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}=1
$$

and

$$
\left\|\left(H_{h_{m}}-z_{m}\right) v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}=\underset{m \rightarrow \infty}{o}\left(h_{m}\right) .
$$

We are going to prove that such a sequence $\left(v_{m}\right)_{m \in \mathbb{N}}$ cannot exist. First considering a sequence of eigenvectors, this will prove that for $h$ small enough, $H_{h}$ has no eigenvalue with real part close to $E$ and imaginary part larger than $h \beta$. But the essential spectrum of $H_{h}$ is $\mathbb{R}_{+}$, so if $z \in \mathbb{C}_{+}$belongs to the spectrum of $H_{h}$ it must be an (isolated) eigenvalue. Thus, the resolvent $\left(H_{h}-z\right)^{-1}$ is well-defined as a bounded operator from $L^{2}\left(\mathbb{R}^{n}\right)$ to $H^{2}\left(\mathbb{R}^{n}\right)$ when $h$ is small enough, $\operatorname{Im} z>h \beta$ and $\operatorname{Re} z$ is close to $E$. Applying again the argument now gives the estimate of Theorem 1.1.

Since we already have the result for large $\beta$, the sequence $\left(\beta_{m}\right)_{m \in \mathbb{N}}$ is necessarily bounded. After extracting a subsequence if necessary, we can assume
that $\beta_{m} \rightarrow \tilde{\beta} \geqslant \beta$. Since a bounded sequence in $L^{2}\left(\mathbb{R}^{n}\right)$ always has a semiclassical measure (see [5, 34]), we can assume after extracting another subsequence that (3.1) holds for some nonnegative Radon measure $\mu$ on $\mathbb{R}^{n}$. Our purpose is now to prove that $\mu$ is both zero and non-zero to get a contradiction.

Proposition 3.2. - The measure $\mu$ is non-zero.
Proof. - As $V_{2}$ goes to 0 at infinity, there exists $R \geqslant 0$ such that $V_{2}(x) \geqslant-\frac{\beta}{2}$ for all $x \in B_{R}^{c}$. We have

$$
\int_{\mathbb{R}^{n}}\left(V_{2}(x)+\beta_{m}\right)\left|v_{m}(x)\right|^{2} d x=-h_{m}^{-1} \operatorname{Im}\left\langle\left(H_{h_{m}}-z_{m}\right) v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow \infty]{ } 0
$$

and hence, for $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ equal to 1 on $B_{R}$,

$$
\begin{aligned}
\frac{\beta}{2} & =\frac{\beta}{2} \int_{\mathbb{R}^{n}}(1-\chi(x))\left|v_{m}(x)\right|^{2} d x+\frac{\beta}{2} \int_{\mathbb{R}^{n}} \chi(x)\left|v_{m}(x)\right|^{2} d x \\
& \leqslant \int_{\mathbb{R}^{n}}\left(V_{2}(x)+\beta_{m}\right)(1-\chi(x))\left|v_{m}(x)\right|^{2} d x+\frac{\beta}{2} \int_{\mathbb{R}^{n}} \chi(x)\left|v_{m}(x)\right|^{2} d x \\
& \leqslant \int_{\mathbb{R}^{n}}\left(\frac{\beta}{2}-V_{2}(x)-\beta_{m}\right) \chi(x)\left|v_{m}(x)\right|^{2} d x+\underset{m \rightarrow+\infty}{o}(1) \\
& \leqslant\left(\frac{\beta}{2}+m_{-}\right) \int_{\mathbb{R}^{n}} \chi(x)\left|v_{m}(x)\right|^{2} d x+\underset{m \rightarrow+\infty}{o}(1)
\end{aligned}
$$

This proves that

$$
\int_{\mathbb{R}^{2 n}} \chi(x) d \mu(x, \xi)=\lim _{m \rightarrow \infty} \int_{\mathbb{R}^{n}} \chi(x)\left|v_{m}(x)\right|^{2} d x \neq 0
$$

We now prove that $\mu$ is actually zero. Note that by Proposition 3.1 we already know that $\mu$ is supported on $p^{-1}(\{E\})$.

Proposition 3.3. - The total measure of $\mu$ is finite.
Proof. - Let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n},[0,1]\right)$. We have

$$
\left\|\mathrm{Op}_{h}^{w}(q)\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant C\|q\|_{L^{\infty}\left(\mathbb{R}^{2 n}\right)}+\underset{m \rightarrow \infty}{O}\left(\sqrt{h_{m}}\right)
$$

where $C$ only depends on the dimension $n$ (see for instance Theorem 5.1 in [34]), and hence:

$$
\int_{\mathbb{R}^{2 n}} q d \mu=\lim _{m \rightarrow \infty}\left\langle\operatorname{Op}_{h_{m}}^{w}(q) v_{m}, v_{m}\right\rangle \leqslant \limsup _{m \rightarrow \infty} C\|q\|_{L^{\infty}\left(\mathbb{R}^{2 n}\right)}\left\|v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \leqslant C
$$

Considering $q_{k}$ equal to 1 on the ball $B_{\mathbb{R}^{2 n}}(k)$ of radius $k \in \mathbb{N}$ in $\mathbb{R}^{2 n}$ proves that $\mu\left(B_{\mathbb{R}^{2 n}}(k)\right) \leqslant C$ for all $k \in \mathbb{N}$.

Proposition 3.4. - $\mu=0$ on $\Omega_{\infty}^{-}(\{E\})$.

Proof. - Let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n},[0,1]\right)$ supported in $\Omega_{\infty}^{-}\left(\mathbb{R}_{+}^{*}\right)$. There exists $T \geqslant 0$ such that for $w \in \operatorname{supp} q$ and $s \geqslant T$ we have $V_{2}(X(-s, w))+\beta \geqslant \frac{\beta}{2}$. Put

$$
C=\sup _{\operatorname{supp} q} \exp \left(-2 \int_{0}^{T}\left(V_{2}+\beta\right) \circ \phi^{-s} d s\right)
$$

According to (3.2) and Proposition 3.3, we have for all $t \geqslant T$ :

$$
\begin{aligned}
0 \leqslant \int_{\mathbb{R}^{2 n}} q d \mu & =\int_{\mathbb{R}^{2 n}}\left(q \circ \phi^{t}\right) \exp \left(-2 \int_{0}^{t}\left(V_{2}+\beta\right) \circ \phi^{t-s} d s\right) d \mu \\
& \leqslant \mu\left(\mathbb{R}^{2 n}\right) \sup _{\operatorname{supp} q} \exp \left(-2 \int_{0}^{t}\left(V_{2}+\beta\right) \circ \phi^{-s} d s\right) \\
& \leqslant C \mu\left(\mathbb{R}^{2 n}\right) \exp (-(t-T) \beta) \\
& \xrightarrow[t \rightarrow+\infty]{ } 0 .
\end{aligned}
$$

This implies that the integral of $q$ is zero and proves the proposition.
Proposition 3.5. - $\mu=0$ on $\Omega_{b}^{-}(\{E\})$ and hence on $\mathbb{R}^{2 n}$.
Proof. - We follow the idea of the previous proof, now using the absorption assumption on trapped trajectories. Let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n},[0,1]\right)$. Since $\mu$ is supported on $\Omega_{b}^{-}(\{E\})$ we have for all $t \geqslant 0$ :

$$
\begin{aligned}
0 \leqslant \int_{\mathbb{R}^{2 n}} q d \mu & =\int_{\mathbb{R}^{2 n}}\left(q \circ \phi^{t}\right) \exp \left(-2 \int_{0}^{t}\left(V_{2}+\beta\right) \circ \phi^{t-s} d s\right) d \mu \\
& \leqslant \mu\left(\mathbb{R}^{2 n}\right) \sup _{\Omega_{b}^{-}(\{E\}) \cap \operatorname{supp} q} \exp \left(-2 \int_{0}^{t}\left(V_{2}+\beta\right) \circ \phi^{-s} d s\right) .
\end{aligned}
$$

Now using Proposition 2.3, we can conclude that the integral of $q$ is zero.
Propositions 3.2 and 3.5 give the contradiction which proves Theorem 1.1. We remark that assumption (1.7) is stronger than necessary to prove Proposition 3.5 , since we did not use the fact that $\operatorname{Im} z_{m}$ is allways greater that $\beta$. We can actually prove the following result:

Corollary 3.6. - Let $E>0$ and $\beta>0$ such that

$$
\forall w \in \Omega_{b}(\{E\}), \exists T \geqslant 0, \quad \int_{0}^{T}\left(V_{2}(X(-s, w))+\beta\right) d s>0
$$

Then there exist a neighborhood $I$ of $E, h_{0}>0$ and $c \geqslant 0$ such that for $\left.h \in\right] 0, h_{0}$ ] the operator $H_{h}$ has no eigenvalue in $\mathbb{C}_{I, h \beta}$ and

$$
\forall z \in \mathbb{C}_{I, h \beta}, \quad\left\|\left(H_{h}-z\right)^{-1}\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h}
$$

As explained before Corollary 1.3, we obtain results on the lower half-plane under the more "symmetric" non-trapping condition:

Corollary 3.7. - Let $E>0$ be a non-trapping energy (see (1.5)) and $\beta>0$. Then there exist a neighborhood $I$ of $E, h_{0}>0$ and $c \geqslant 0$ such that for $\left.h \in] 0, h_{0}\right], \operatorname{Re} z \in I$ and $|\operatorname{Im} z| \geqslant h \beta$ the operator $\left(H_{h}-z\right)$ has a bounded inverse on $L^{2}\left(\mathbb{R}^{n}\right)$ and

$$
\left\|\left(H_{h}-z\right)^{-1}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h} .
$$

Proof. - We only have to apply Theorem 1.1 to $H_{h}$ and its adjoint $H_{h}^{*}$.

## 4. Estimate in the incoming region

In this section we prove an estimate for the outgoing solution of the Helmholtz equation in the incoming region.

Theorem 4.1. - Assume that $V_{2}$ is of short range (see (1.8)). Let $I \subset \mathbb{R}$, $h_{0}>0, \delta>\frac{1}{2}, c \geqslant 0, k \in \mathbb{N}$ and suppose that for $\left.\left.h \in\right] 0, h_{0}\right]$ and $z \in \mathbb{C}_{I,+}$ the resolvent $\left(H_{h}-z\right)^{-1}$ is well-defined and

$$
\left\|\langle x\rangle^{-\delta}\left(H_{h}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{c}{h^{k}} .
$$

Let $R_{1}>0, d>d_{1} \geqslant 0$ and $\left.\sigma, \sigma_{1} \in\right]-1,1\left[\right.$ such that $\sigma_{1}<\sigma$. Then there exists $R>R_{1}$ such that for $z \in \mathbb{C}_{I,+}, \omega \in \&_{0}\left(\mathbb{R}^{2 n}\right)$ supported outside $\mathcal{Z}_{-}\left(R_{1}, d_{1},-\sigma_{1}\right)$ and $\omega_{-} \in \wp_{0}\left(\mathbb{R}^{2 n}\right)$ supported in $Z_{-}(R, d,-\sigma)$ we have

$$
\left\|\langle x\rangle^{-\delta} \mathrm{Op}_{h}\left(\omega_{-}\right)\left(H_{h}-z\right)^{-1} \mathrm{Op}_{h}(\omega)\langle x\rangle^{-\delta}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)}=\underset{h \rightarrow 0}{O}\left(h^{\infty}\right),
$$

and the size of the rest is uniform in $z \in \mathbb{C}_{I,+}$. If the limiting absorption principle holds in $\mathcal{L}\left(L^{2, \delta}\left(\mathbb{R}^{n}\right), L^{2,-\delta}\left(\mathbb{R}^{n}\right)\right)$ for $\lambda \in I$ and $h>0$ small enough, then the estimate remains true for $\left(H_{h}-(\lambda+i 0)\right)^{-1}, \lambda \in I$. Moreover if $\operatorname{supp} \omega \subset B_{x}(r)$ for some $r>0$, then the weight $\langle x\rangle^{-\delta}$ can be replaced on both sides by $\langle x\rangle^{\beta}$ for any $\beta \in \mathbb{R}$.

REmark 4.2. - We are going to use this result for a small perturbation of a dissipative Schrödinger operator in order to prove Theorem 1.2 (see Proposition 5.3). Then, once Theorem 1.2 is proved, we can use Theorem 4.1 for the full non-dissipative Schrödinger operator $H_{h}$ we are interested in (see Section 7).

The result being stronger in this case, we can assume without loss of generality that $d_{1}$ is non-zero. The proof of this theorem follows that of the dissipative analog given in [27]. We recall the sketch of the proof for the reader convenience and refer to [24, 32, 27, 25] for more details.

Let $\left.d_{0} \in\right] 0, d_{1}\left[\right.$ and $\left.\sigma_{0} \in\right]-1, \sigma_{1}\left[\right.$. There exist $R_{0}>0$ and $\phi \in C^{\infty}\left(\mathbb{R}^{2 n}\right)$ such that

$$
\begin{equation*}
\forall(x, \xi) \in Z_{-}\left(R_{0}, d_{0},-\sigma_{0}\right), \quad\left|\nabla_{x} \phi(x, \xi)\right|^{2}+V_{1}(x)=|\xi|^{2} \tag{4.1}
\end{equation*}
$$

and, for some $\rho>0$ :

$$
\forall(x, \xi) \in \mathbb{R}^{2 n}, \forall \alpha, \beta \in \mathbb{N}^{n}, \quad\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta}(\phi(x, \xi)-\langle x, \xi\rangle)\right| \leqslant c_{\alpha, \beta}\langle x\rangle^{1-\rho-|\alpha|}
$$

(see [16]). As explained in [32], we can assume that the constants $c_{\alpha, \beta}>0$ for $\alpha, \beta \in \mathbb{N}^{n}$ are as small as we wish as long as we replace $\phi$ by

$$
\begin{equation*}
(x, \xi) \mapsto(\phi(x, \xi)-\langle x, \xi\rangle) \chi\left(\frac{x}{R}\right)+\langle x, \xi\rangle \tag{4.2}
\end{equation*}
$$

for $R \geqslant 2 R_{0}$ large enough and $\chi \in C^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\chi(x)=0$ for $|x| \leqslant 1 / 4$ and $\chi(x)=1$ for $|x| \geqslant 1 / 2$. In this case (4.1) remains valid on $Z_{-}\left(R / 2, d_{0},-\sigma_{0}\right)$.

For all $(x, \xi) \in \mathbb{R}^{2 n}$ we denote by $t \mapsto r(t, x, \xi) \in \mathbb{R}^{n}$ the solution of the problem

$$
\left\{\begin{array}{l}
\partial_{t} r(t, x, \xi)=\nabla_{x} \phi(r(t, x, \xi), \xi) \\
r(0, x, \xi)=x
\end{array}\right.
$$

We can check that this defines a smooth function on $\mathbb{R} \times \mathbb{R}^{2 n}$. If $R$ was chosen large enough in (4.2), then for $(x, \xi) \in \mathcal{Z}_{-}\left(0, d_{1},-\sigma_{1}\right)$ and $t \geqslant 0$ we have

$$
|r(-t, x, \xi)| \geqslant c_{0}(t+|x|)
$$

for some $c_{0}>0$. Moreover for $\alpha, \beta \in \mathbb{N}^{n}$ with $|\alpha|+|\beta| \geqslant 1$ there exists $c_{\alpha, \beta} \geqslant 0$ such that

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} r(-t, x, \xi)\right| \leqslant c_{\alpha, \beta}(t+\langle x\rangle)\langle x\rangle^{-|\alpha|} .
$$

For $t \geqslant 0$ and $(x, \xi) \in Z_{-}\left(0, d_{1},-\sigma_{1}\right)$ we now set

$$
F(t, x, \xi)=\Delta_{x} \phi(r(t, x, \xi), \xi)-V_{2}(r(t, x, \xi))
$$

and define on $Z_{-}\left(0, d_{1},-\sigma_{1}\right)$ the symbols

$$
a_{0}(x, \xi)=\exp \left(-\int_{0}^{\infty} F(-2 s, x, \xi) d s\right)
$$

and, for $j \geqslant 1$ :

$$
a_{j}(x, \xi)=i \int_{0}^{+\infty} \Delta_{x} a_{j-1}(r(-2 \tau, x, \xi), \xi) \exp \left(-\int_{0}^{\tau} F(-2 s, x, \xi) d s\right) d \tau
$$

These functions are solutions of the transport equations

$$
2 \nabla_{x} a_{0} \cdot \nabla_{x} \phi+a_{0} \Delta_{x} \phi-a_{0} V_{2}=0
$$

and, for $j \geqslant 1$ :

$$
2 \nabla_{x} a_{j} \cdot \nabla_{x} \phi+a_{j} \Delta_{x} \phi-a_{j} V_{2}-i \Delta_{x} a_{j-1}=0
$$

Moreover $a_{j}$ decays as a function of $\S_{-j}\left(\mathbb{R}^{2 n}\right)$ and there exists $c_{0}>0$ such that

$$
\forall(x, \xi) \in \mathcal{Z}_{-}\left(0, d_{1},-\sigma_{1}\right), \quad\left|a_{0}(x, \xi)\right| \geqslant c_{0}
$$

Note that $V_{2}$ has to be of short range here but the sign does not matter.
Since we work on $\mathcal{Z}_{-}\left(0, d_{1},-\sigma_{1}\right)$, we now introduce a cut-off function as follows. We choose $R_{2}$ and $R_{3}$ such that $\max \left(R_{1}, R / 2\right)<R_{2}<R_{3}<R, d_{2}$ and $d_{3}$ such that $d_{1}<d_{2}<d_{3}<d$ and finally $\sigma_{2}$ and $\sigma_{3}$ such that $\sigma_{1}<\sigma_{2}<\sigma_{3}<\sigma$. Then we consider $\chi_{1}, \chi_{2}, \chi_{3} \in C^{\infty}(\mathbb{R},[0,1])$ such that $\chi_{1}(s)=0$ for $s \leqslant R_{2}$, $\chi_{1}(s)=1$ for $s \geqslant R_{3}, \chi_{2}(s)=0$ for $s \leqslant d_{2}, \chi_{2}(s)=1$ for $s \geqslant d_{3}, \chi_{3}(s)=0$ for $s \leqslant \sigma_{2}$ and $\chi_{3}(s)=1$ for $s \geqslant \sigma_{3}$. We fix $N \in \mathbb{N}$. Let us define

$$
a(h)=\sum_{j=0}^{N} h^{j} a_{j} \quad \text { and } \quad b(x, \xi, h)=\chi_{1}(|x|) \chi_{2}(|\xi|) \chi_{3}\left(-\frac{x \cdot \xi}{|x||\xi|}\right) a(x, \xi, h) .
$$

We also consider
$p(h)=\frac{i}{h}\left(\left|\nabla_{x} \phi\right|^{2}+V_{1}-|\xi|^{2}\right) b(h)+\left(2 \nabla_{x} b(h) \cdot \nabla_{x} \phi+b(h) \Delta_{x} \phi-b(h) V_{2}\right)-i h \Delta_{x} b(h)$.
The symbols $b(h)$ and $p(h)$ are supported in $\mathcal{Z}_{-}\left(R_{2}, d_{2},-\sigma_{2}\right)$ and for $\alpha, \beta \in \mathbb{N}^{n}$ there exists a constant $c_{\alpha, \beta}>0$ such that for $\left.\left.h \in\right] 0,1\right]$ we have

$$
\forall(x, \xi) \in Z_{-}\left(R_{2}, d_{2},-\sigma_{2}\right), \quad\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} b(x, \xi, h)\right|+\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} p(x, \xi, h)\right| \leqslant c_{\alpha, \beta}\langle x\rangle^{-|\alpha|}
$$

and

$$
\forall(x, \xi) \in Z_{-}\left(R_{3}, d_{3},-\sigma_{3}\right), \quad\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} p(x, \xi, h)\right| \leqslant c_{\alpha, \beta} h^{N+1}\langle x\rangle^{-2-N-|\alpha|}
$$

If $R$ is chosen large enough, $\left.R_{5} \in\right] R_{3}, R\left[, d_{5} \in\right] d_{3}, d\left[\right.$ and $\left.\sigma_{5} \in\right] \sigma_{3}, \sigma[$, then we can construct (see [32, Lemma 4.5]) a symbol $e(h)=\sum_{j=0}^{N} h^{j} e_{j}$ such that $e_{j} \in \&_{-j}\left(\mathbb{R}^{2 n}\right)$ is supported in $Z_{-}\left(R_{5}, d_{5},-\sigma_{5}\right)$ for all $j \in \llbracket 0, N \rrbracket$ and

$$
I_{h}(e(h), \phi) I_{h}(b(h), \phi)^{*}=\mathrm{Op}_{h}\left(\omega_{-}\right)+h^{N+1} \mathrm{Op}_{h}(r(h)),
$$

where $r(h) \in \delta_{-N}\left(\mathbb{R}^{2 n}\right)$ uniformly in $\left.\left.h \in\right] 0,1\right]$ and for $u \in \&\left(\mathbb{R}^{n}\right)$ we have set

$$
I_{h}(b, \phi) u(x)=\frac{1}{(2 \pi h)^{n}} \int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}} e^{\frac{i}{h}(\phi(x, \xi)-\langle y, \xi\rangle)} b(x, \xi) u(y) d \xi d y
$$

For any $t \geqslant 0$ we have

$$
I_{h}(b(h), \phi)^{*} U_{h}(t)=U_{0}^{h}(t) I_{h}(b(h), \phi)^{*}-\int_{0}^{t} U_{0}^{h}(s) I_{h}(p, \phi)^{*} U_{h}(t-s) d s
$$

and hence

$$
\begin{aligned}
\mathrm{Op}_{h}\left(\omega_{-}\right) U_{h}(t)= & -h^{N+1} \mathrm{Op}_{h}(r(h)) U_{h}(t)+I_{h}(e(h), \phi) U_{0}^{h}(t) I_{h}(b(h), \phi)^{*} \\
& -\int_{0}^{t} I_{h}(e(h), \phi) U_{0}^{h}(s) I_{h}(p(h), \phi) U_{h}(t-s) d s
\end{aligned}
$$

Contrary to the dissipative case, we cannot write

$$
\left(H_{h}-z\right)^{-1}=\frac{i}{h} \int_{0}^{\infty} e^{-\frac{i t}{h}\left(H_{h}-z\right)} d t=\frac{i}{h} \int_{0}^{\infty} e^{\frac{i t}{h} z} U_{h}(t) d t
$$

for any $z \in \mathbb{C}_{+}$, but this is at least possible when $\operatorname{Im} z>h\left\|V_{2}\right\|_{\infty}$. For such a $z$, we obtain from the previous equality:

$$
\begin{aligned}
&\langle x\rangle^{-\delta} \mathrm{Op}_{h}\left(\omega_{-}\right)\left(H_{h}-z\right)^{-1} \mathrm{Op}_{h}(\omega)\langle x\rangle^{-\delta} \\
&=-h^{N+1}\langle x\rangle^{-\delta} \mathrm{Op}_{h}(r(h))\left(H_{h}-z\right)^{-1} \mathrm{Op}_{h}(\omega)\langle x\rangle^{-\delta} \\
& \quad+\frac{i}{h}\langle x\rangle^{-\delta} \int_{0}^{\infty} e^{\frac{i t}{h} z} I_{h}(e(h), \phi) U_{0}^{h}(t) I_{h}(b(h), \phi)^{*} \mathrm{Op}_{h}(\omega)\langle x\rangle^{-\delta} d t \\
&-\langle x\rangle^{-\delta} \int_{0}^{\infty} e^{\frac{i s}{h} z} I_{h}(e(h), \phi) U_{0}^{h}(s) I_{h}(p(h), \phi)^{*}\left(H_{h}-z\right)^{-1} \mathrm{Op}_{h}(\omega)\langle x\rangle^{-\delta} d s .
\end{aligned}
$$

This equality is proved for $\operatorname{Im} z>\left\|V_{2}\right\|_{\infty}$ but the two integrands decay with time uniformly in $z \in \mathbb{C}_{I,+}$. Each term is holomorphic on $\mathbb{C}_{I,+}$, so for any $h \in] 0, h_{0}$ ] this equality remains valid on $\mathbb{C}_{I,+}$ by unique continuation. Then it only remains to chose $N$ large enough and estimate each term of the right-hand side to conclude (we use the assumption on resolvent estimates here). This can be done as in the dissipative case, using Propositions 3.2 and 3.5 of [32].

## 5. Uniform resolvent estimates

We now prove the uniform resolvent estimates for the non-dissipative Schrödinger operator up to the real axis. In order to use Theorem 4.1 we assume that $V_{2}$ is of short range. But we expect Theorem 1.2 to be true under a weaker assumption on $V_{2}$, so we are going to give the other arguments only assuming that $V_{2} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ is of long range with a short range negative part: there exist $\rho>0, C \geqslant 0$ and constants $c_{\alpha}$ for $\alpha \in \mathbb{R}^{n}$ such that for all $x \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
V_{2}(x) \geqslant-C\langle x\rangle^{-1-\rho} \quad \text { and } \quad \forall \alpha \in \mathbb{N}^{n},\left|\partial^{\alpha} V_{2}(x)\right| \leqslant c_{\alpha}\langle x\rangle^{-|\alpha|-\rho} \tag{5.1}
\end{equation*}
$$

As for Theorem 1.1, we proceed by contradiction. We suppose that Theorem 1.2 is wrong and consider sequences $\left(v_{m}\right)_{m \in \mathbb{N}} \in H^{2}\left(\mathbb{R}^{n}\right)^{\mathbb{N}},\left(z_{m}\right)_{m \in \mathbb{N}} \in \mathbb{C}^{\mathbb{N}}$

[^1]and $\left.\left.\left(h_{m}\right)_{m \in \mathbb{N}} \in\right] 0,1\right]^{\mathbb{N}}$ such that if we set $\lambda_{m}=\operatorname{Re} z_{m}$ and $\beta_{m}=h_{m}^{-1} \operatorname{Im} z_{m}$ then for some $\delta \in] \frac{1}{2}, \frac{1+\rho}{2}[$ we have
$$
h_{m} \rightarrow 0, \quad \lambda_{m} \rightarrow E, \quad 0<\beta_{m} \rightarrow 0, \quad\left\|v_{m}\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}=1
$$
and
$$
\left\|\left(H_{h_{m}}-z_{m}\right) v_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}=\underset{m \rightarrow \infty}{o}\left(h_{m}\right) .
$$

We remark that $v_{m}$ is assumed to be in $H^{2}\left(\mathbb{R}^{n}\right)$ for all $m \in \mathbb{N}$, but is only uniformly bounded in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$.

As for Theorem 1.1, we show that such a sequence $\left(v_{m}\right)_{m \in \mathbb{N}}$ cannot exist, which proves first that for $h$ small enough, $H_{h}$ has no eigenvalue with real part close to $E$ and positive imaginary part, and then the estimate of Theorem 1.2.

If there exists a subsequence $\left(m_{k}\right)_{k \in \mathbb{N}}$ such that $\beta_{m_{k}} \geqslant \beta>0$ for all $k \in \mathbb{N}$, then we obtain a contradiction with Theorem 1.1. Therefore we can assume that

$$
\beta_{m} \xrightarrow[m \rightarrow \infty]{ } 0
$$

After extracting a subsequence if necessary, we can assume that (3.1) holds for some non-negative Radon measure $\mu$. We already know that $\mu$ is supported in $p^{-1}(\{E\})$. In order to get a contradiction, we prove that $\mu=0$ and $\mu \neq 0$.

Let $W_{2}=V_{2}+2 C\langle x\rangle^{-1-\rho} \geqslant C\langle x\rangle^{-1-\rho}$, the constant $C$ being given by (5.1). We first prove that $\mu \neq 0$. The proof relies on the existence of an escape function in the sense of [17]:

Proposition 5.1. - Let $E>0$. There exist $f \in C_{b}^{\infty}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$, $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ and $\tilde{\chi} \in C_{0}^{\infty}(\mathbb{R},[0,1])$ equal to 1 in a neighborhood of $E$ such that

$$
\forall(x, \xi) \in \mathbb{R}^{2 n}, \quad\{p, f\}(x, \xi)=(1-\chi(x)) \tilde{\chi}(p(x, \xi))\langle x\rangle^{-2 \delta}
$$

The proof of this proposition is postponed to Appendix A.

Proposition 5.2. - The measure $\mu$ is non-zero.

Proof. - 1. Let $\theta \in C_{0}^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ be supported in $B_{2}$ and equal to 1 on $B_{1}$. For $R>0$ we set $\theta_{R}(x)=\theta\left(\frac{x}{R}\right)$. If there exists $R>0$ such that

$$
\begin{equation*}
\int_{\mathbb{R}^{2 n}} \theta_{R}(x)\langle x\rangle^{-1-\rho} d \mu(x, \xi) \neq 0 \tag{5.2}
\end{equation*}
$$

then the proposition is proved. Otherwise, according to Proposition 3.1, we have for any $R>0$

$$
\begin{aligned}
& \limsup _{m \rightarrow \infty}\left|\left\langle\langle x\rangle^{-1-\rho} v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}\right| \\
& \leqslant \limsup _{m \rightarrow \infty}\left|\left\langle\theta_{R}\langle x\rangle^{-1-\rho} v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}\right| \\
& \quad+\limsup _{m \rightarrow \infty}\left|\left\langle\left(1-\theta_{R}\right)\langle x\rangle^{-1-\rho} v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}\right| \\
& \leqslant\left\|\langle x\rangle^{2 \delta-1-\rho}\left(1-\theta_{R}\right)\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)} \leqslant\langle R\rangle^{2 \delta-1-\rho} .
\end{aligned}
$$

This proves that

$$
\begin{equation*}
\left\langle\langle x\rangle^{-1-\rho} v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow+\infty]{ } 0, \tag{5.3}
\end{equation*}
$$

and hence:

$$
\begin{align*}
& \beta_{m}\left\|v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}+\left\|\sqrt{W_{2}} v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}  \tag{5.4}\\
&=-h_{m}^{-1} \operatorname{Im}\left\langle\left(H_{h_{m}}-z_{m}\right) v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \\
&+2 C\left\langle\langle x\rangle^{-1-\rho} v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow+\infty]{ } 0 .
\end{align*}
$$

Since both terms of the left-hand side are non-negative, this means that each goes to 0 as $m$ goes to $+\infty$. Moreover $\sqrt{W_{2}}$ is smooth and all its derivatives of order at least 1 belong to $\delta\left(\langle x\rangle^{-\delta}\right)$, so for any $f \in C_{b}^{\infty}\left(\mathbb{R}^{2 n}\right)$ we have (5.5)

$$
\left\|\sqrt{W_{2}} \mathrm{Op}_{h}^{w}(f) v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\|\mathrm{Op}_{h}^{w}(f) \sqrt{W_{2}} v_{m}\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}+\underset{m \rightarrow \infty}{O}\left(h_{m}\right) \underset{m \rightarrow \infty}{ } 0
$$

2. Let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n},[0,1]\right), \tilde{\chi} \in C_{0}^{\infty}(\mathbb{R},[0,1])$ and $f \in C_{b}^{\infty}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$ be given by Proposition 5.1. For $m \in \mathbb{N}$ we have:

$$
\begin{aligned}
1= & \left\langle\langle x\rangle^{-2 \delta} v_{m}, v_{m}\right\rangle \\
= & \left\langle\langle x\rangle^{-2 \delta} \chi(x) v_{m}, v_{m}\right\rangle+\left\langle\mathrm{Op}_{h_{m}}^{w}\left(\langle x\rangle^{-2 \delta}(1-\chi(x))((1-\tilde{\chi}) \circ p)\right) v_{m}, v_{m}\right\rangle \\
& +\left\langle\mathrm{Op}_{h_{m}}^{w}\left(\langle x\rangle^{-2 \delta}(1-\chi(x))(\tilde{\chi} \circ p)\right) v_{m}, v_{m}\right\rangle .
\end{aligned}
$$

According to Proposition 3.1, the second term goes to 0 as $m$ goes to $+\infty$. We now show that this also holds for the third term in order to prove that

$$
\int_{\mathbb{R}^{2 n}}\langle x\rangle^{-2 \delta} \chi(x) d \mu(x, \xi)=\lim _{m \rightarrow \infty}\left\langle\langle x\rangle^{-2 \delta} \chi(x) v_{m}, v_{m}\right\rangle=1 \neq 0 .
$$

3. We have

$$
\frac{i}{h_{m}}\left[H_{h_{m}}^{1}, \mathrm{Op}_{h_{m}}^{w}(f)\right]=\mathrm{Op}_{h_{m}}^{w}(\{p, f\})+h_{m}^{2} \mathrm{Op}_{h_{m}}^{w}\left(r_{3}\left(h_{m}\right)\right),
$$

where $r_{3}(h) \in \phi\left(\langle x\rangle^{-1-\rho}\right)$ uniformly in $\left.\left.h \in\right] 0,1\right]$, and hence

$$
\begin{aligned}
&\left\langle\mathrm{Op}_{h}^{w}\left((1-\chi(x))\langle x\rangle^{-2 \delta}(\tilde{\chi} \circ p)\right) v_{m}, v_{m}\right\rangle=\left\langle\mathrm{Op}_{h_{m}}^{w}(\{p, f\}) v_{m}, v_{m}\right\rangle \\
&= \frac{i}{h_{m}}\left\langle\left[H_{1}^{h}, \mathrm{Op}_{h_{m}}^{w}(f)\right] v_{m}, v_{m}\right\rangle+\underset{m \rightarrow \infty}{O}\left(h_{m}^{2}\right)
\end{aligned}
$$

Since $v_{m} \in H^{2}\left(\mathbb{R}^{n}\right)$ for all $m \in \mathbb{N}$ and according to (5.3)-(5.5) we have

$$
\begin{align*}
& \left\langle\mathrm{Op}_{h}^{w}\left((\tilde{\chi} \circ p)(1-\chi(x))\langle x\rangle^{-2 \delta}\right) v_{m}, v_{m}\right\rangle  \tag{5.6}\\
& =\frac{i}{h_{m}}\left\langle\left(\left(H_{h_{m}}-z_{m}\right)^{*} \operatorname{Op}_{h_{m}}^{w}(f)-\mathrm{Op}_{h_{m}}^{w}(f)\left(H_{h_{m}}-z_{m}\right)\right) v_{m}, v_{m}\right\rangle+\underset{m \rightarrow \infty}{o}(1) \\
& \underset{m \rightarrow \infty}{ } 0
\end{align*}
$$

which concludes the proof.
Let

$$
H_{h}^{2}=-h^{2} \Delta+V_{1}(x)-i h W_{2}(x) .
$$

The Schrödinger operator $H_{h}^{2}$ is dissipative and its dissipative part is positive on trapped trajectories of energy $E$, so there exist a neighborhood $I$ of $E$, $h_{0}>0$ and $C_{2}>0$ such that for $\left.\left.h \in\right] 0, h_{0}\right]$ and $z \in \mathbb{C}_{I,+}$ we have

$$
\left\|\langle x\rangle^{-\delta}\left(H_{h}^{2}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{C_{2}}{h}
$$

(see [26]). Since $2 \delta<1+\rho$ we can write $2 C\langle x\rangle^{-1-\rho}=W_{3}+W_{4}$ where $W_{4} \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and

$$
\forall x \in \mathbb{R}^{n}, \quad\langle x\rangle^{2 \delta}\left|W_{3}(x)\right| \leqslant \frac{1}{2 C_{2}}
$$

Put

$$
H_{h}^{3}=-h^{2} \Delta+V_{1}(x)-i h W_{2}(x)+i h W_{3}(x)=H_{h}^{2}+i h W_{3}(x) .
$$

Let $z \in \mathbb{C}_{I,+}$. By a standard perturbation argument, we know that for $\left.\left.h \in\right] 0, h_{0}\right]$ the resolvent $\left(H_{h}^{3}-z\right)^{-1}$ is well-defined and

$$
\begin{equation*}
\left\|\langle x\rangle^{-\delta}\left(H_{h}^{3}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{2 C_{2}}{h} . \tag{5.7}
\end{equation*}
$$

As a result we can apply Theorem 4.1 with $H_{h}^{3}$ on $\mathbb{C}_{I,+}$ (we recall that $V_{2}$ has to be of short range here). We use this result to prove that the semiclassical measure $\mu$ is supported outside $\Omega_{\infty}^{-}(\{E\})$ :

Proposition 5.3. - $\mu=0$ on $\Omega_{\infty}^{-}(\{E\})$.

Proof. - Let $J \subset I \cap] E / 2,2 E[$ be a neighborhood of $E$. We first check that $\mu=0$ in the incoming region $\mathcal{Z}_{-}(R, 0,-1 / 2)$ for some $R$ large enough. Let $R_{1}$ be such that supp $\left.W_{4} \subset B_{R_{1}}, d \in\right] 0, \sqrt{E / 2}\left[\right.$ and $\sigma=\frac{1}{2}$. Let $R$ be given by Theorem 4.1 applied to $H_{h}^{3}$, and finally $\omega_{-} \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ supported in $Z_{-}(R, d,-1 / 2)$. For $m$ large enough the operator $\left(H_{h_{m}}^{3}-z_{m}\right)$ has a bounded inverse, so we can write

$$
\begin{aligned}
\operatorname{Op}_{h_{m}}^{w}\left(\omega_{-}\right) v_{m}= & \operatorname{Op}_{h_{m}}^{w}\left(\omega_{-}\right)\left(H_{h_{m}}^{3}-z_{m}\right)^{-1}\left(H_{h_{m}}-z_{m}\right) v_{m} \\
& -i h_{m} \mathrm{Op}_{h_{m}}^{w}\left(\omega_{-}\right)\left(H_{h_{m}}^{3}-z_{m}\right)^{-1} W_{4} v_{m} .
\end{aligned}
$$

According to (5.7) and Theorem 4.1 we obtain

$$
\begin{align*}
& \left\|\mathrm{Op}_{h_{m}}^{w}\left(\omega_{-}\right) v_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}  \tag{5.8}\\
& =\left\|\langle x\rangle^{\delta} \mathrm{Op}_{h}^{w}\left(\omega_{-}\right)\langle x\rangle^{\delta}\right\|\left\|\langle x\rangle^{-\delta}\left(H_{h_{m}}^{3}-z_{m}\right)^{-1}\langle x\rangle^{-\delta}\right\|\left\|\left(H_{h_{m}}-z_{m}\right) v_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)} \\
& \quad+h_{m}\left\|\langle x\rangle^{\delta} \mathrm{Op}_{h}^{w}\left(\omega_{-}\right)\left(H_{h_{m}}^{3}-z_{m}\right)^{-1} W_{4}\langle x\rangle^{\delta}\right\|\left\|v_{m}\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}
\end{align*}
$$

$$
\xrightarrow[m \rightarrow \infty]{ } 0
$$

This proves that

$$
\left\langle\mathrm{Op}_{h_{m}}^{w}\left(\omega_{-}\right) v_{m}, v_{m}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \xrightarrow[m \rightarrow \infty]{ } 0
$$

and hence $\mu=0$ on $Z_{-}(R, d,-1 / 2)$. Now let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ supported in $\Omega_{\infty}^{-}(J)$. For $t \geqslant 0$ large enough we have $\phi^{-t}(\operatorname{supp} q) \subset Z_{-}(R, d,-1 / 2)$. According to Proposition 3.1 we obtain

$$
\int_{\mathbb{R}^{2 n}} q d \mu=\int_{\mathbb{R}^{2 n}}\left(q \circ \phi^{t}\right) \exp \left(-2 \int_{0}^{t} V_{2} \circ \phi^{t-s} d s\right) d \mu=0
$$

which proves that $\mu=0$ on $\Omega_{\infty}^{-}(J)$.
When (5.2) holds for some $R>0$, then Proposition 5.2 is proved but not (5.6). So we cannot use it to show that $\mu$ is zero at infinity as in the self-adjoint case (see [17]).

Proposition 5.4. - $\mu=0$ on $\Omega_{b}^{-}(\{E\})$.
Proof. - This proposition is proved as Proposition 3.5, even though the total measure of $\mu$ is no longer necessarily finite. Let $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n},[0,1]\right)$. We know that $\mu$ is supported in $\Omega_{b}^{-}(\{E\})$, so according to Proposition 3.1 we can write

$$
\int_{\mathbb{R}^{2 n}} q d \mu=\int_{\mathbb{R}^{2 n}} \mathbb{1}_{\Omega_{b}^{-}(\{E\})}\left(q \circ \phi^{t}\right) \exp \left(-2 \int_{0}^{t} V_{2} \circ \phi^{t-s} d s\right) d \mu
$$

Since the set

$$
\bigcup_{t \geqslant 0} \phi^{-t}\left(\operatorname{supp} q \cap \Omega_{b}^{-}(\{E\})\right)
$$

is bounded, there exists $c \geqslant 0$ such that for all $t \geqslant 0$

$$
\int_{\mathbb{R}^{2 n}} q d \mu \leqslant c \sup _{\operatorname{supp} q \cap \Omega_{b}^{-}(\{E\})} \exp \left(-2 \int_{0}^{t} V_{2} \circ \phi^{-s} d s\right)
$$

Then we can conclude with Proposition 2.3.

## 6. Limiting Absorption Principle

After having proved resolvent estimates on $\mathbb{C}_{I,+}$, we can show the limiting absorption principle and study existence and uniqueness of an outgoing solution for (1.1). Before giving more precise statements, we introduce some notation. Let

$$
\mathbb{C}_{++}=\{\zeta \in \mathbb{C}: \operatorname{Re} \zeta>0, \operatorname{Im} \zeta \geqslant 0\}
$$

For $u \in H^{1}\left(\mathbb{R}^{n}\right)$ we set on $\mathbb{R}^{n} \backslash\{0\}$ :

$$
\partial_{r} u=\frac{x \cdot \nabla u}{|x|}, \quad \mathscr{D}_{r} u=\partial_{r} u+\frac{n-1}{2|x|} u \quad \text { and } \quad \nabla_{\perp} u=\nabla u-\frac{x \partial_{r} u}{|x|} .
$$

This defines functions in $L^{2}\left(\mathbb{R}^{n} \backslash \mathcal{V}\right)$ for any neighborhood $V$ of 0 . We are going to use the following standard properties of these operators:
Proposition 6.1. - Let $u, v \in H^{1}\left(\mathbb{R}^{n}\right)$ and assume that $u$ is compactly supported outside a neighborhood of 0 . We have

$$
\left\langle\mathscr{D}_{r} u, v\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=-\left\langle u, \mathscr{D}_{r} v\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

and

$$
\left\langle\partial_{r} u, \partial_{r} v\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle\mathscr{D}_{r} u, \mathscr{D}_{r} v\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}+\left\langle\frac{(n-1)(n-3)}{4|x|^{2}} u, v\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

If moreover $u$ belongs to $H^{2}\left(\mathbb{R}^{n}\right)$ we also have

$$
\partial_{r} \nabla_{\perp} u(x)=-\frac{1}{|x|} \nabla_{\perp} u(x)+\nabla_{\perp} \partial_{r} u(x)
$$

In this section we consider a Schrödinger operator

$$
H=-\Delta+V_{1}(x)-i V_{2}(x)
$$

with domain $\mathscr{D}(H)=H^{2}\left(\mathbb{R}^{n}\right)$. $V_{1}$ and $V_{2}$ are bounded and real-valued. We assume that

$$
\begin{equation*}
V_{1}-i V_{2}=W_{1}-i W_{2}+W_{3}, \tag{6.1}
\end{equation*}
$$

where:
(i) $W_{1}$ and $W_{2}$ are differentiable,
(ii) for all $x \in \mathbb{R}^{n}$ we have $W_{1}(x) \in \mathbb{R}, W_{2}(x) \geqslant 0$ and $W_{3}(x) \in \mathbb{C}$,
(iii) there exist $\rho \in] 0,1]$ and $c \geqslant 0$ such that for all $x \in \mathbb{R}^{n}$ we have
$\left|W_{1}(x)\right|+W_{2}(x) \leqslant c\langle x\rangle^{-\rho} \quad$ and $\quad\left|\nabla W_{1}(x)\right|+\left|\nabla W_{2}(x)\right|+\left|W_{3}(x)\right| \leqslant c\langle x\rangle^{-1-\rho}$.
For $x \in B_{1}^{c}$ we also set

$$
\tilde{W}_{3}(x)=W_{3}(x)+\frac{(n-1)(n-3)}{4|x|^{2}} .
$$

Définition 6.2. - Let $\zeta \in \mathbb{C}_{++}, f \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{n}\right)$, and suppose that $u \in H_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right)$ is a solution for the equation

$$
\begin{equation*}
\left(H-\zeta^{2}\right) u=f \tag{6.3}
\end{equation*}
$$

Then we say that $u$ is an outgoing solution for (6.3) if there exists $\delta>\frac{1}{2}$ such that $\left(\mathscr{D}_{r}-i \zeta\right) u \in L^{2, \delta-1}\left(B_{1}^{c}\right)$.

Let $\delta \in] \frac{1}{2}, \frac{1}{2}+\frac{\rho}{4}$ [ be fixed for all this section, $\rho$ being given by (6.2). Let $K$ be a compact subset of $\mathbb{C}_{++}$such that $K=\overline{K \cap \mathbb{C}_{+}}$. We set $K^{*}=K \cap \mathbb{C}_{+}$.

Proposition 6.3. - Assume the resolvent $\left(H-\zeta^{2}\right)^{-1}$ is defined for all $\zeta \in K^{*}$ and the equation $\left(H-\zeta^{2}\right) u=0$ has no non-trivial outgoing solution when $\zeta \in K$.

Let $\lambda \in K \cap \mathbb{R}_{+}^{*}$ and $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$. Then $\left(H-\zeta^{2}\right)^{-1} f$ converges in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ to the unique outgoing solution for the equation $\left(H-\lambda^{2}\right) u=f$ when $\zeta \in K^{*}$ goes to $\lambda$.

Moreover, there exists a constant $C \geqslant 0$ such that for any $\zeta \in K$ and $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$, if we denote by $u$ the unique outgoing solution for the equation $\left(H-\zeta^{2}\right) u=f$ then we have for all $R>0$ the following estimates:

$$
\begin{equation*}
\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\left\|\left(\mathscr{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(B_{1}^{c}\right)}+R^{\delta-\frac{1}{2}}\|u\|_{L^{2,-\delta}\left(B_{R}^{c}\right)} \leqslant C\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)} \tag{6.4}
\end{equation*}
$$

We denote by $\left(H-\left(\lambda^{2}+i 0\right)\right)^{-1} f$ the unique outgoing solution $u \in H_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right) \cap$ $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ for the equation $\left(H-\lambda^{2}\right) u=f$. To prove this proposition, we study the behavior of $\left(H-\zeta^{2}\right)^{-1} f, \zeta \in K^{*}$, at infinity. In a compact subset of $\mathbb{R}^{n}$, the estimates we need are given by interior regularity (see for instance [12, § 6.3.1]):

Proposition 6.4. - Let $z \in \mathbb{C}$ and $R \geqslant 0$. There exists $C \geqslant 0$ such that if $f \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{n}\right)$ and $u \in H_{\mathrm{loc}}^{2}\left(\mathbb{R}^{n}\right)$ is a solution for the equation $(H-z) u=f$ then we have

$$
\|u\|_{H^{2}\left(B_{R}\right)} \leqslant C\left(\|f\|_{B_{R+1}}+\|u\|_{B_{R+1}}\right) .
$$

Moreover $C$ is uniform for $(z, R)$ in a compact subset of $\mathbb{C} \times \mathbb{R}_{+}$.

The difficulty is to give some estimates of $\left(H-\zeta^{2}\right)^{-1} f$ uniform when $\zeta$ approaches $\mathbb{R}_{+}^{*}$. For some fixed $\zeta \in K^{*}$ we have the following lemma:

Lemma 6.5. - Let $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right), \zeta \in K^{*}$ and $u=\left(H-\zeta^{2}\right)^{-1} f$. Then $u$ and $\nabla u$ belong to $L^{2, \delta}\left(\mathbb{R}^{n}\right)$.

Proof. - Let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ be supported in $B_{2}$ and equal to 1 on $B_{1}$. For any $R>0$ we set $\chi_{R}: x \mapsto \chi\left(\frac{x}{R}\right)$. Then we have to check that [ $\chi_{R}\langle x\rangle^{\delta},\left(H-\zeta^{2}\right)^{-1}$ ] is bounded uniformly in $R>0$. We prove similarly that $\chi_{R}\langle x\rangle^{\delta} \partial_{j} u$ is in $L^{2}\left(\mathbb{R}^{n}\right)$ uniformly in $R>0$ for all $j \in \llbracket 1, n \rrbracket$.

The self-adjoint version for the first estimate of the following result is Lemma 4.1 in [29]:

Lemma 6.6. - There exists $C$ such that for $\zeta=\zeta_{1}+i \zeta_{2} \in K^{*}, f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ and $u=\left(H-\zeta^{2}\right)^{-1} f$ we have

$$
\zeta_{2}\|u\|_{L^{2,1-\delta}\left(\mathbb{R}^{n}\right)} \leqslant C\left(\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\left\|\left(\mathscr{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)}+\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}\right)
$$

and

$$
\begin{aligned}
& \| \\
& \left\|\sqrt{W_{2}} u\right\|_{L^{2, \frac{1}{2}-\delta}\left(\mathbb{R}^{n}\right)}^{2} \\
& \quad \leqslant C\left(\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\left\|\left(\mathscr{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)}+\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}\right)\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}
\end{aligned}
$$

Proof. - Let $\chi \in C^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ be equal to 1 outside $B_{3}$ and equal to 0 on $B_{2}$. For $R \geqslant 1$ we set $\chi_{R}: x \mapsto \chi(x)-\chi(x / R)$. Let $\left.\left.\alpha \in\right]-\infty, 1-\delta\right]$, $\zeta=\zeta_{1}+i \zeta_{2} \in K^{*}, f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ and $u=\left(H-\zeta^{2}\right)^{-1} f \in H^{2}\left(\mathbb{R}^{n}\right)$. For all $R \geqslant 1$ we can write

$$
\left\langle\left(H-\zeta^{2}\right) u, \chi_{R}^{2}(x)(1+|x|)^{2 \alpha} u\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle f, \chi_{R}^{2}(x)(1+|x|)^{2 \alpha} u\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} .
$$

Taking the imaginary part in this equality gives

$$
\begin{aligned}
& \left\|\chi_{R} \sqrt{W_{2}} u\right\|_{L^{2, \alpha}\left(\mathbb{R}^{n}\right)}^{2}+2 \zeta_{1} \zeta_{2}\left\|\chi_{R} u\right\|_{L^{2, \alpha}\left(\mathbb{R}^{n}\right)}^{2} \\
& \leqslant\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}\|u\|_{L^{2,2 \alpha-\delta}\left(\mathbb{R}^{n}\right)}+\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}\left\|W_{3} u\right\|_{L^{2,2 \alpha+\delta}\left(\mathbb{R}^{n}\right)} \\
& \quad+\left\|\chi_{R} \partial_{r} u\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}\left\|\left(2(1+|x|)^{2 \alpha} \partial_{r} \chi_{R}+2 \alpha(1+|x|)^{2 \alpha-1} \chi_{R}\right) u\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)} \\
& \leqslant c\left(\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}+\left\|\left(\mathcal{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n} \backslash B_{2}\right)}+\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}\right)\|u\|_{L^{2,2 \alpha+\delta-1}\left(\mathbb{R}^{n}\right)},
\end{aligned}
$$

where $c$ does not depend on $R$. Then we can let $R$ go to infinity in the left-hand side. Proposition 6.4 and this inequality with $\alpha=1-\delta$ give the first estimate. We take $\alpha=\frac{1}{2}-\delta$ to obtain the second.

For all $\zeta \in K^{*}$ and $x \in \mathbb{R}^{n}$ the complex number $\zeta^{2}+i W_{2}(x)$ belongs to $\mathbb{C}_{+}$ and hence has a unique square root $\zeta_{W}(x, \zeta)=\zeta_{1}(x, \zeta)+i \zeta_{2}(x, \zeta)$ in $\mathbb{C}_{++}$. Moreover there exists $C>0$ such that for all $\zeta=\zeta_{1}+i \zeta_{2} \in K^{*}$ and $x \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
C^{-1}\left(2 \zeta_{1} \zeta_{2}+W_{2}(x)\right) \leqslant \zeta_{2}(x, z) \leqslant C\left(2 \zeta_{1} \zeta_{2}+W_{2}(x)\right) \tag{6.5}
\end{equation*}
$$

and

$$
\left|\nabla_{x} \zeta_{W}(x, \zeta)\right| \leqslant C\langle x\rangle^{-1-\rho},
$$

where $\rho>0$ is given by (6.2). For $\zeta \in K^{*}$ we set

$$
\mathscr{D}_{r}^{\zeta}=\mathscr{D}_{r}-i \zeta_{W}(x, \zeta) .
$$

Proposition 6.7. - There exists $C$ such that for $\zeta=\zeta_{1}+i \zeta_{2} \in K^{*}$, $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ and $u=\left(H-\zeta^{2}\right)^{-1} f$ we have

$$
\left\|\left(\mathscr{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)}+\left\|\nabla_{\perp} u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)} \leqslant C\left(\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}\right) .
$$

The proof is inspired from the self-adjoint version given in [29], to which we refer for some additional details. The new difficulty here is to take into account the long-range imaginary part $W_{2}$ of the potential.

Proof. - 1. Since

$$
\left|\zeta_{W}(x, \zeta)-\zeta\right|=\frac{W_{2}(x)}{\left|\zeta_{W}(x, \zeta)+\zeta\right|} \leqslant c W_{2}(x)
$$

we have

$$
\left\|\left(\mathscr{D}_{r}-i \zeta\right) u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)} \leqslant\left\|\mathscr{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)}+c\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}
$$

where $c \geqslant 0$ denotes different constants which do not depend on $\zeta \in K^{*}$ or $f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$. It is therefore enough to prove the proposition with $\left(\mathscr{D}_{r}-i \zeta\right)$ replaced by $\mathscr{D}_{r}^{\zeta}$.
2. We consider $\chi: x \mapsto \tilde{\chi}(|x|)$, where $\tilde{\chi} \in C^{\infty}(\mathbb{R},[0,1])$ is non-decreasing, equal to 0 on $]-\infty, 1]$ and equal to 1 on $\left[2,+\infty\left[\right.\right.$. For $R>1$ we set $\chi_{R}(x)=$ $\tilde{\chi}(|x|)-\tilde{\chi}(|x|-R)$. Let $\zeta=\zeta_{1}+i \zeta_{2} \in K^{*}, f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ and $u=\left(H-\zeta^{2}\right)^{-1} f$. It is enough to consider the case $\mathcal{N}_{u}<\mathcal{N}_{d}$, where

$$
\mathcal{N}_{u}=\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}
$$

and

$$
\mathcal{N}_{d}=\max \left(\left\|\chi \mathscr{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)},\left\|\chi \nabla_{\perp} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)}\right)
$$

According to Proposition 6.1 we can write for any $R>1$ :

$$
\begin{aligned}
\operatorname{Re} & \left\langle f, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
= & \operatorname{Re}\left\langle\left(H-\zeta^{2}\right) u, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
= & \operatorname{Re}\left\langle\mathscr{D}_{r} u, \mathscr{D}_{r} \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle+\operatorname{Re}\left\langle\nabla_{\perp} u, \nabla_{\perp} \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
& +\operatorname{Re}\left\langle W_{1}(x) u+\tilde{W}_{3}(x) u-\zeta_{W}(x, \zeta)^{2} u, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
= & A_{1}(R)+A_{2}(R)+A_{3}(R)+A_{4}(R)+A_{5}(R) .
\end{aligned}
$$

3. We have

$$
\begin{aligned}
A_{1}(R)+A_{5}(R)= & \operatorname{Re}\left\langle\mathscr{D}_{r}^{\zeta} u, \mathscr{D}_{r} \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
& +\operatorname{Re}\left\langle-i\left(\partial_{r} \zeta_{W}(x, \zeta)+\zeta_{W}(x, \zeta) \mathscr{D}_{r}^{\zeta}\right) u, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle
\end{aligned}
$$

According to Lemma 6.5, $\mathscr{D}_{r}^{\zeta} u$ belongs to $L^{2, \delta}\left(B_{1}^{c}\right)$ and hence

$$
\begin{aligned}
\liminf _{R \rightarrow \infty} \operatorname{Re}\left\langle\mathscr{D}_{r}^{\zeta} u\right. & \left., \mathscr{D}_{r} \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
& =\frac{1}{2} \liminf _{R \rightarrow \infty}\left\langle\mathscr{D}_{r}^{\zeta} u, \partial_{r}\left(\chi_{R}^{2}(1+|x|)^{2 \delta-1}\right) \mathscr{D}_{r}^{\zeta} u\right\rangle \\
& \geqslant\left(\delta-\frac{1}{2}\right)\left\|\chi \mathscr{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)}^{2}-c\left\|\mathscr{D}_{r}^{\zeta} u\right\|_{L^{2}\left(B_{2} \cap B_{1}^{c}\right)}^{2}
\end{aligned}
$$

Moreover

$$
\begin{gathered}
\liminf _{R \rightarrow \infty} \operatorname{Re}\left\langle-i\left(\partial_{r} \zeta_{W}(x, \zeta)+\zeta_{W}(x, \zeta) \mathscr{D}_{r}^{\zeta}\right) u, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
\quad \geqslant \liminf _{R \rightarrow \infty} \operatorname{Re}\left\langle-i \partial_{r} \zeta_{W}(x, \zeta) u, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \\
\quad \geqslant-c\|u\|_{L^{2,-1-\rho+\delta}\left(\mathbb{R}^{n}\right)}\left\|\chi \mathscr{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)}
\end{gathered}
$$

According to Proposition 6.4, these two estimates give

$$
\begin{equation*}
\liminf _{R \rightarrow \infty}\left(A_{1}(R)+A_{5}(R)\right) \geqslant\left(\delta-\frac{1}{2}\right)\left\|\chi \mathcal{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(\mathbb{R}^{n}\right)}-c \mathcal{N}_{d} \mathcal{N}_{u} \tag{6.6}
\end{equation*}
$$

With the same kind of argument and using the last property of Proposition 6.1 we prove that

$$
\liminf _{R \rightarrow \infty} A_{2}(R) \geqslant\left(\frac{3}{2}-\delta\right)\left\|\chi \nabla_{\perp} u\right\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}^{2}-c \mathcal{N}_{u} \mathcal{N}_{d}
$$

Using (6.5) and Lemma 6.6 we show that

$$
\liminf _{R \rightarrow \infty} A_{3}(R) \geqslant-c \mathcal{N}_{u} \mathcal{N}_{d}
$$

We have the same estimate for $A_{4}(R)$ and

$$
\forall R>1, \quad \operatorname{Re}\left\langle f, \chi_{R}^{2}(1+|x|)^{2 \delta-1} \mathscr{D}_{r}^{\zeta} u\right\rangle \leqslant c \mathcal{N}_{u} \mathcal{N}_{d}
$$

so we finally have

$$
\left(\frac{3}{2}-\delta\right)\left\|\nabla_{\perp} u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)}^{2}+\left(\delta-\frac{1}{2}\right)\left\|\chi \mathscr{D}_{r}^{\zeta} u\right\|_{L^{2, \delta-1}\left(B_{2}^{c}\right)}^{2} \leqslant c \mathcal{N}_{u} \mathcal{N}_{d}
$$

which, together with Proposition 6.4, concludes the proof.
As in the self-adjoint case we also use the following estimate:
Proposition 6.8. - There exists $C$ such that for $R \geqslant 1, \zeta \in K^{*}, f \in L^{2, \delta}\left(\mathbb{R}^{n}\right)$ and $u=\left(H-\zeta^{2}\right)^{-1} f$ we have

$$
\|u\|_{L^{2,-\delta}\left(B_{R}^{c}\right)}^{2} \leqslant C R^{1-2 \delta}\left(\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}^{2}+\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}^{2}\right) .
$$

Proof. - For all $r>0$ we have
(6.7) $\left|\left(\mathscr{D}_{r}-i \zeta\right) u\right|_{\mathbb{S}_{r}}^{2}$
$=\left|\mathscr{D}_{r} u+\zeta_{2} u\right|_{\mathbb{S}_{r}}^{2}+\zeta_{1}^{2}|u|_{\mathbb{S}_{r}}^{2}+2 \zeta_{1}\left\langle V_{2} u, u\right\rangle_{B_{r}}+4 \zeta_{1}^{2} \zeta_{2}\|u\|_{B_{r}}^{2}+2 \zeta_{1} \operatorname{Im}\langle f, u\rangle_{B_{r}}$.
In particular

$$
\zeta_{1}^{2}|u|_{\mathbb{S}_{r}}^{2} \leqslant\left|\left(\mathscr{D}_{r}-i \zeta\right) u\right|_{\mathbb{S}_{r}}^{2}+2 \zeta_{1}\|f\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}+c\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}^{2}
$$

We multiply by $(1+r)^{-2 \delta}$ and integrate from $R$ to $+\infty$ to prove the proposition.

With all these estimates and uniqueness of the outgoing solution for the equation $\left(H-\lambda^{2}\right) u=0$ we can now conclude that the limiting absorption principle holds as in the self-adjoint case.

Proof of Proposition 6.3. - The first step is to prove that if we have the estimates (6.4) on $K^{*}$, then we have the limiting absorption principle. We refer to [29, Lemma 2.6] and recall briefly the idea. If $\zeta_{m} \rightarrow \lambda$ and $u_{m}=\left(H-\zeta_{m}^{2}\right)^{-1} f$, then a subsequence of $\left(u_{m}\right)_{m \in \mathbb{N}}$ converges to some $u$ in $L_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right)$ according to Proposition 6.4. We obtain convergence in $L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ according to the last estimate of (6.4). Using again Proposition 6.4 we also have convergence in $H_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right)$. The limit $u$ is necessarily an outgoing solution for the equation $\left(H-\lambda^{2}\right) u=f$. And since we have uniqueness for such a solution, we actually have convergence of the whole sequence.

To apply this result, we still have to check the first estimate of (6.4). This is a contradiction argument: we assume that we can find sequences $\left(f_{m}\right)_{m \in \mathbb{N}} \in L^{2, \delta}\left(\mathbb{R}^{n}\right)^{\mathbb{N}}$ and $\left(z_{m}\right)_{m \in \mathbb{N}} \in\left(K^{*}\right)^{\mathbb{N}}$ such that

$$
z_{m} \xrightarrow[m \rightarrow \infty]{ } \lambda \in K \cap \mathbb{R}_{+}^{*}, \quad\left\|\left(H-z_{m}\right)^{-1} f_{m}\right\|_{L^{2,-\delta}}=1
$$

and

$$
\left\|f_{m}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}=\underset{m \rightarrow \infty}{o}(1) .
$$

All the estimates of (6.4) hold for this sequence and hence $\left(H-z_{m}\right)^{-1} f_{m}$ converges to an outgoing solution for the equation $(H-\lambda) u=0$, which must be zero. This gives a contradiction (see the proof of Theorem 2.3 in [29]).

For the dissipative Schrödinger operator everything holds as in the selfadjoint case on the upper half-plane:

Proposition 6.9. - If $V_{2} \geqslant 0$ then the assumptions and hence the conclusions of Proposition 6.3 hold for any compact subset $K$ of $\mathbb{C}_{++}$.

Proof. - If $V_{2} \geqslant 0$ then the operator $H$ with domain $H^{2}\left(\mathbb{R}^{n}\right)$ is maximal dissipative, so its resolvent $(H-z)^{-1}$ is at least well-defined for $\operatorname{Im} z>0$. Moreover, if $u \in H_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right) \cap L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ is an outgoing solution for the equation $\left(H-\zeta^{2}\right) u=0$ for some $\zeta \in K$ then, according to (6.7), $\zeta \in \mathbb{R}$ and $u$ vanishes on the support of $V_{2}$ and hence is an outgoing solution for the equation $\left(H_{1}-\zeta^{2}\right) u=0$, where $H_{1}=-\Delta+V_{1}(x)$ is the self-adjoint part of $H$. We know from [29, §3] that such a solution must be zero.

In the non-dissipative case we do not have such a systematic result. In order to prove Theorem 1.4 we use the fact that for $h>0$ small enough and $z \in \mathbb{C}_{I,+}$ we already have existence and uniform estimates for the resolvent $\left(H_{h}-z\right)^{-1}$.

Proof of Theorem 1.4. - Let $\theta_{0} \geqslant 0$ be such that $V_{2}(x)+\theta_{0}\langle x\rangle^{-1-\rho} \geqslant 0$ for all $x \in \mathbb{R}^{n}$. There exist a neighborhood $I$ of $E, h_{0}>0$ and $C>0$ such that for all $\left.h \in] 0, h_{0}\right], z \in \mathbb{C}_{I,+}$ and $\theta \in\left[0, \theta_{0}\right]$ we have

$$
\left\|\langle x\rangle^{-\delta}\left(H_{h}^{\theta}-z\right)^{-1}\langle x\rangle^{-\delta}\right\|_{\underline{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \frac{C}{h}
$$

where $H_{h}^{\theta}=H_{h}-i h \theta\langle x\rangle^{-1-\rho}$. Indeed, Theorem 1.2 gives such an estimate for any fixed $\tilde{\theta} \in\left[0, \theta_{0}\right]$, and the perturbation argument already used to prove (5.7) gives an estimate uniform for $\theta$ in a neighborhood of $\tilde{\theta}$. Now let $\theta_{1} \in\left[0, \theta_{0}\right]$ and assume that $u=0$ is the unique outgoing solution for the equation $\left(H_{h}^{\theta}-z\right) u=0$ when $\theta \in\left[\theta_{1}, \theta_{0}\right], z \in \mathbb{C}_{I,+} \cup I$ and $\left.\left.h \in\right] 0, h_{0}\right]$. Let $\left.\left.h \in\right] 0, h_{0}\right]$, $z \in \mathbb{C}_{I,+} \cup I, \theta \in\left[\theta_{1}-C / 2, \theta_{1}\right]$ and let $u \in H_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right) \cap L^{2,-\delta}\left(\mathbb{R}^{n}\right)$ be an outgoing solution for the equation $\left(H_{h}^{\theta}-z\right) u=0$. We have

$$
\left(H_{h}^{\theta_{1}}-z\right) u=-i h\left(\theta_{1}-\theta\right)\langle x\rangle^{-1-\rho} u \in L^{2, \delta}\left(\mathbb{R}^{n}\right)
$$

According to Proposition 6.3, the outgoing solution for the equation $\left(H_{h}^{\theta_{1}}-z\right) u=f$ is given by the limiting absorption principle if $z \in I$ and hence in any case we have

$$
\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)} \leqslant \frac{C}{h}\left\|h\left(\theta_{1}-\theta\right)\langle x\rangle^{-1-\rho} u\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)} \leqslant \frac{1}{2}\|u\|_{L^{2,-\delta}\left(\mathbb{R}^{n}\right)}
$$

This proves that $u=0$.

## 7. Semiclassical Measure

We study in this section the semiclassical measures for the outgoing solution of (1.1) when the source term $f_{h}$ concentrates on a bounded submanifold of $\mathbb{R}^{n}$, $V_{1}$ is of long range and $V_{2}$ is of short range. We adapt to this setting the proof given in [27] for the dissipative case.

Let $\Gamma$ be a (bounded) submanifold of dimension $d \in \llbracket 0, n-1 \rrbracket$ in $\mathbb{R}^{n}$. We consider $A \in C_{0}^{\infty}(\Gamma), S \in \gamma\left(\mathbb{R}^{n}\right)$ and define

$$
\begin{equation*}
f_{h}(x)=\int_{\Gamma} A(z) S\left(\frac{x-z}{h}\right) d \sigma_{\Gamma}(z) \tag{7.1}
\end{equation*}
$$

where $\sigma_{\Gamma}$ is the Lebesgue mesure on $\Gamma$. We can check that $f_{h}$ is microlocalized on $N \Gamma$ and $\left\|f_{h}\right\|_{L^{2, \delta}\left(\mathbb{R}^{n}\right)}=O(\sqrt{h})$ for any $\delta>\frac{1}{2}$. Let $E>0$ be an energy which satisfies assumption (1.7). We assume that

$$
\begin{equation*}
\forall z \in \Gamma, \quad V_{1}(z)<E . \tag{7.2}
\end{equation*}
$$

Let

$$
\begin{aligned}
N_{E} \Gamma & =N \Gamma \cap p^{-1}(\{E\}) \\
& =\left\{(z, \xi) \in \Gamma \times \mathbb{R}^{n}: \xi \perp T_{z} \Gamma \text { and }|\xi|^{2}=E-V_{1}(z)\right\} \subset \mathbb{R}^{2 n}
\end{aligned}
$$

where $T_{z} \Gamma$ is the tangent space of $\Gamma$ at point $z$. Assumption (7.2) ensures that $N_{E} \Gamma$ is a submanifold of dimension $n-1$ in $\mathbb{R}^{2 n}$. The Riemannian structure $g$ on $N_{E} \Gamma$ is defined as follows. For $(z, \xi) \in N_{E} \Gamma$ and $(Z, \Xi)$, $(\tilde{Z}, \tilde{\Xi}) \in T_{(z, \xi)} N_{E} \Gamma \subset \mathbb{R}^{2 n}$ we set

$$
g_{(z, \xi)}((Z, \Xi),(\tilde{Z}, \tilde{\Xi}))=\langle Z, \tilde{Z}\rangle_{\mathbb{R}^{n}}+\left\langle\Xi_{\perp}, \tilde{\Xi}_{\perp}\right\rangle_{\mathbb{R}^{n}}
$$

where $\Xi_{\perp}, \tilde{\Xi}_{\perp}$ are the orthogonal projections of $\Xi, \tilde{\Xi} \in \mathbb{R}^{n}$ on $\left(T_{z} \Gamma \oplus \mathbb{R} \xi\right)^{\perp}=$ $T_{\xi}\left(\left(\{z\} \times \mathbb{R}^{n}\right) \cap N_{E} \Gamma\right)$. We denote by $\sigma_{N_{E} \Gamma}$ the canonical measure on $N_{E} \Gamma$ given by $g$, and assume that

$$
\begin{equation*}
\sigma_{N_{E} \Gamma}\left(\left\{(z, \xi) \in N_{E} \Gamma: \exists t>0, \phi^{t}(z, \xi) \in N_{E} \Gamma\right\}\right)=0 . \tag{7.3}
\end{equation*}
$$

With all these assumptions we can prove the following result:
Theorem 7.1. - Let $f_{h}$ be given by (7.1) and $u_{h}$ be the outgoing solution for the Helmholtz Equation (1.1). Let assumptions (1.3), (1.7), (1.8), (7.2) and (7.3) be fulfilled.
(i) There exists a non-negative Radon measure $\mu$ on $\mathbb{R}^{2 n}$ such that

$$
\begin{equation*}
\forall q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right), \quad\left\langle\mathrm{Op}_{h}^{w}(q) u_{h}, u_{h}\right\rangle \underset{h \rightarrow 0}{ } \int_{\mathbb{R}^{2 n}} q d \mu \tag{7.4}
\end{equation*}
$$

(ii) This measure is characterized by the following three properties: a. $\mu$ is supported in $p^{-1}(\{E\})$.

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b. For any $\sigma \in] 0,1[$ there exists $R \geqslant 0$ such that $\mu=0$ in the incoming region $Z_{-}(R, 0,-\sigma)$.
c. $\mu$ is solution of the Liouville equation

$$
\begin{equation*}
\{p, \mu\}+2 V_{2} \mu=\pi(2 \pi)^{d-n}|A(z)|^{2}|\xi|^{-1}|\hat{S}(\xi)|^{2} \sigma_{N_{E} \Gamma} \tag{7.5}
\end{equation*}
$$

(iii) These three properties imply that for $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ we have

$$
\begin{equation*}
\int_{\mathbb{R}^{2 n}} q d \mu \tag{7.6}
\end{equation*}
$$

$$
=\int_{0}^{+\infty} \int_{N_{E} \Gamma} \pi(2 \pi)^{d-n}|A(z)|^{2}|\xi|^{-1}|\hat{S}(\xi)|^{2} q\left(\phi^{t}(z, \xi)\right) e^{-2 \int_{0}^{t} V_{2}(X(s, z, \xi)) d s} d \sigma_{N_{E} \Gamma}(z, \xi) d t
$$

Note that as in $[3,27]$ we can let $E$ depend on $h: E_{h}=E_{0}+h E_{1}+o(h) \in \mathbb{C}_{+}$, where $E_{0}>0$ satisfies assumption (1.7) and $\operatorname{Im} E_{h} \geqslant 0$. Then $V_{2}$ has to be replaced by $V_{2}+\operatorname{Im} E_{1}$ in (7.5) and (7.6).

We recall the sketch of the proof, discuss differences with the dissipative case and refer to [3, 27] for details (see also [25]). We first remark that the limit (7.4) is zero when $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ is supported outside $p^{-1}(\{E\})$. For $q$ supported close to $p^{-1}(\{E\})$ the idea is to replace the resolvent which defines $u_{h}$ by the integral over finite times of the propagator. More precisely, for $T \geqslant 0$ and $h \in] 0,1]$ we set

$$
u_{h}^{T}=\frac{i}{h} \int_{0}^{\infty} \chi_{T}(t) e^{-\frac{i t}{h}\left(H_{h}-E\right)} f_{h}
$$

where $\chi_{T}(t)=\chi(t-T)$ and $\chi \in C^{\infty}(\mathbb{R},[0,1])$ is equal to 1 in a neighborhood of ] $-\infty, 0$ ] and equal to 0 on $\left[\tau_{0},+\infty\left[\right.\right.$ for some $\tau_{0}>0$ small enough (see [27]). Then we can study separately the contribution of different times. For small times we proceed exactly as in the dissipative case, using a W.K.B. approximation of the propagator (see below). For intermediate times, and then to prove that $u_{h}^{T}$ is in some sense a good approximation of $u_{h}$ for large $T$ and small $h$, we need a non-selfadjoint version of Egorov's Theorem.

According to (3.3) and Hille-Yosida Theorem (see for instance Theorem II.3.5 in [11]) we know that $H_{h}$ generates a continuous semi-group, which we denote by $U_{h}(t)$, and

$$
\forall t \geqslant 0, \quad\left\|U_{h}(t)\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant e^{t m_{-}}
$$

(we recall that $\left.m_{-}=-\inf V_{2}\right)$. Let $W_{2}, \tilde{W}_{2} \in C_{b}^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}\right), W=W_{2}+\tilde{W}_{2}$ and, for $t \in \mathbb{R}$ :

$$
U_{2}^{h}(t)=e^{-\frac{i t}{h}\left(H_{1}^{h}-i h W_{2}\right)} \quad \text { and } \quad \tilde{U}_{2}^{h}(t)=e^{-\frac{i t}{h}\left(H_{1}^{h}-i h \tilde{W}_{2}\right)} .
$$

The usual Egorov's Theorem extends without modification to the nondissipative case:

Theorem 7.2. - Let $a \in C_{b}^{\infty}\left(\mathbb{R}^{2 n}\right)$. There exists a family of symbols $\alpha_{j}(t)$ for $j \in \mathbb{N}$ and $t \in \mathbb{R}$ such that:
(i) For all $t \in \mathbb{R}, N \in \mathbb{N}$ and $h \in] 0,1]$ we have

$$
U_{2}^{h}(t)^{*} \mathrm{Op}_{h}^{w}(a) \tilde{U}_{2}^{h}(t)=\sum_{j=0}^{N} h^{j} \mathrm{Op}_{h}^{w}\left(\alpha_{j}(t)\right)+h^{N+1} R_{N}(t, h),
$$

where $R_{N}(t, h)$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ uniformly in $\left.\left.h \in\right] 0,1\right]$ and $t \in[0, T]$ for any $T \geqslant 0$.
(ii) We have

$$
\alpha_{0}(t)=\left(a \circ \phi^{t}\right) e^{-\int_{0}^{t} W \circ \phi^{s} d s} .
$$

(iii) For $t \in \mathbb{R}$ and $j \in \mathbb{N}$ we have

$$
\operatorname{supp} \alpha_{j}(t) \subset \phi^{-t}(\operatorname{supp} a)
$$

Let $w \in \mathbb{R}^{2 n}, T \geqslant 0$ and $0<t_{w, 1}<\cdots<t_{w, K_{w}^{T}} \leqslant T+\tau_{0}$ be the times between 0 and $T+\tau_{0}$ for which $\phi^{-t_{w, k}}(w) \in N_{E} \Gamma_{0}=N_{E} \Gamma \cap\left(\operatorname{supp} A \times \mathbb{R}^{n}\right)$. For $\tau_{w}>0$ small enough we consider $\chi_{w} \in C_{0}^{\infty}(] 0,2 \tau_{w}[)$ equal to 1 in a neighborhood of $\tau_{w}$. For $k \in \llbracket 1, K_{w}^{T} \rrbracket$ we prove that in $L^{2}\left(\mathbb{R}^{n}\right)$

$$
\frac{i}{h} \int_{0}^{\infty} \chi_{T}(t) \chi_{w}\left(t-t_{w, k}+\tau_{w}\right) e^{-\frac{i t}{h}\left(H_{h}-E\right)} f_{h}=B_{w, k}^{T}(h)+\underset{h \rightarrow 0}{O}(\sqrt{h})
$$

where $B_{w, k}^{T}(h)$ is a Lagrangian distribution of Lagrangian submanifold

$$
\Lambda_{w, k}^{T}=\left\{\phi^{t}(z, \xi),(z, \xi) \in N_{E} \Gamma, t \in\right] t_{w, k}-\tau_{w}, t_{w, k}+\tau_{w}[ \} .
$$

This means that there exist $N \in \mathbb{N}, b_{w, k}^{T} \in C_{0}^{\infty}\left(\mathbb{R}^{n+N}\right)$ and a non-degenerate phase function $\psi \in C_{b}^{\infty}\left(\mathbb{R}^{n+N}, \mathbb{R}\right)$ (if $\nabla_{\theta} \psi(x, \theta)=0$ for some $(x, \theta) \in \mathbb{R}^{n+N}$, then the $N$ linear forms $d_{(x, \theta)} \partial_{\theta_{i}} \psi: \mathbb{R}^{n+N} \rightarrow \mathbb{R}, 1 \leqslant i \leqslant N$, are linearly independent) such that

$$
B_{w, k}^{T}(h)=\frac{1}{(2 \pi h)^{\frac{N}{2}}} \int_{\mathbb{R}^{N}} e^{\frac{i}{h} \psi(x, \theta)} b(x, \theta) d \theta
$$

and

$$
\left\{\left(x, \nabla_{x} \psi(x, \theta)\right) \text { for }(x, \theta) \in \mathbb{R}^{n+N} \text { such that } \nabla_{\theta} \psi(x, \theta)=0\right\} \subset \Lambda_{w, k}^{T}
$$

(this replaces what is said in [27]). This is proved by direct computations when $t_{w, k}$ is replaced by $\tau_{w}$ (and $N=0$ in this case), and then we use the fact that for any $t>0$ the propagator $U_{h}(t)$ can be seen as a Fourier Integral Operator and maps, up to $O\left(h^{\infty}\right)$, a Lagrangian distribution of submanifold $\Lambda$ to some Lagrangian distribution of submanifold $\phi^{t}(\Lambda)$. We know that for such
a Lagrangian distribution there exists a smooth and non-negative function $\nu_{w, k}^{T}$ on $\Lambda_{w, k}^{T}$ such that

$$
\forall q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right), \quad\left\langle\mathrm{Op}_{h}^{w}(q) B_{w, k}^{T}(h), B_{w, k}^{T}(h)\right\rangle \underset{h \rightarrow 0}{\longrightarrow} \int_{\Lambda_{w, k}^{T}} q(\tilde{w}) \nu(\tilde{w}) d \sigma_{\Lambda_{w, k}^{T}}(\tilde{w})
$$

where $\sigma_{\Lambda_{w, k}^{T}}$ is the Lebesgue measure on $\Lambda_{w, k}^{T}$. According to Egorov's Theorem, times far from 0 and $t_{w, k}^{T}\left(k \in \llbracket 1, K_{w}^{T} \rrbracket\right)$ do not give any contribution around $w$ at the limit $h \rightarrow 0$, so we can prove that (7.4) holds for some measure $\mu_{T}$ if $u_{h}$ is replaced by $u_{h}^{T}$.

It remains to study the contribution of large times. In [27, Prop. 2.3] we used the fact that the damping factor $\exp \left(-\int_{0}^{t} V_{2} \circ \phi^{s} d s\right)$ is a non-increasing function of $t$, which is no longer the case. We use Proposition 2.5 instead:

Proposition 7.3. - Let $J$ be a neighborhood of $E$ such that assumption (1.7) holds for all $\lambda \in J$. Let $K_{1}$ and $K_{2}$ be compact subsets of $p^{-1}(J)$ and $\varepsilon>0$. Then there exists $T_{0} \geqslant 0$ such that for $q_{1}, q_{2} \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ respectively supported in $K_{1}$ and $K_{2}$ we have

$$
\forall T \geqslant T_{0}, \quad \limsup _{h \rightarrow 0}\left\|\mathrm{Op}_{h}^{w}\left(q_{1}\right) U_{h}(T) \mathrm{Op}_{h}^{w}\left(q_{2}\right)\right\|_{\mathcal{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \varepsilon\left\|q_{1}\right\|_{\infty}\left\|q_{2}\right\|_{\infty}
$$

Proof. - Let $q_{1}, q_{2} \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ be respectively supported in $K_{1}$ and $K_{2}$. Let $t \mapsto U_{1}^{h}(t)$ denote the unitary group generated by the self-adjoint part $H_{1}^{h}$ of $H_{h}$. According to Egorov's Theorem we have for all $T \geqslant 0$ :

$$
\begin{aligned}
& \| \mathrm{Op}_{h}^{w}\left(q_{1}\right) \\
& \quad U_{h}(T) \mathrm{Op}_{h}^{w}\left(q_{2}\right) \|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \\
& \quad=\left\|U_{1}^{h}(T)^{*} \operatorname{Op}_{h}^{w}\left(q_{1}\right) U_{h}(T) \mathrm{Op}_{h}^{w}\left(q_{2}\right)\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \\
& \quad=\left\|\operatorname{Op}_{h}^{w}\left(q_{2}\left(q_{1} \circ \phi^{T}\right) e^{-\int_{0}^{T} V_{2} \circ \phi^{s} d s}\right)\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)}+\underset{h \rightarrow 0}{O}(h) \\
& \quad \leqslant C \sup _{w \in \mathbb{R}^{2 n}}\left|q_{2}(w) q_{1}\left(\phi^{T}(w)\right) e^{-\int_{0}^{T} V_{2}(X(s, w)) d s}\right|+\underset{h \rightarrow 0}{O}(\sqrt{h}),
\end{aligned}
$$

where the size of the rest depends on $T, q_{1}$ and $q_{2}$. The constant $C$ only depends on the dimension $n$. According to Proposition 2.5 there exists $T_{0}$ such that for $T \geqslant T_{0}$ and $w \in K_{2}$ we have

$$
C e^{-\int_{0}^{T} V_{2}(X(s, w)) d s} \leqslant \varepsilon \quad \text { or } \quad \phi^{T}(w) \notin K_{1} .
$$

Therefore we have for all $T \geqslant T_{0}$ :

$$
\left\|\mathrm{Op}_{h}^{w}\left(q_{1}\right) U_{h}(T) \mathrm{Op}_{h}^{w}\left(q_{2}\right)\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)} \leqslant \varepsilon\left\|q_{1}\right\|_{\infty}\left\|q_{2}\right\|_{\infty}+\underset{h \rightarrow 0}{O}(\sqrt{h})
$$

It only remains to take the limit $h \rightarrow 0$ for fixed $T, q_{1}$ and $q_{2}$ to conclude.

For the rest of the proof we proceed as in the dissipative case. We only have to be careful with the proof of Lemma 5.4 in [27] since the resolvent $\left(H_{h}-z\right)^{-1}$ cannot be written as the integral of the propagator over positive times for all $z \in \mathbb{C}_{+}$. However, for $h>0$ small enough and $z \in \mathbb{C}_{+}$close to $E$, $\left(H_{h}-z\right)^{-1} f_{h}$ is well-defined and belongs to $H^{2}\left(\mathbb{R}^{n}\right)$. Therefore we can write

$$
\begin{aligned}
& \mathrm{Op}_{h}^{w}(q)\left(H_{h}-z\right)^{-1} f_{h}-\mathrm{Op}_{h}^{w}(q) e^{-\frac{i T}{h}\left(H_{h}-z\right)}\left(H_{h}-z\right)^{-1} f_{h} \\
&=-\int_{0}^{T} \mathrm{Op}_{h}^{w}(q) \frac{d}{d t} e^{-\frac{i t}{h}\left(H_{h}-z\right)}\left(H_{h}-z\right)^{-1} f_{h} d t \\
&= \frac{i}{h} \int_{0}^{\infty} \chi_{T}(t) \mathrm{Op}_{h}^{w}(q) e^{-\frac{i t}{h}\left(H_{h}-z\right)} f_{h} d t \\
&-\frac{i}{h} \int_{0}^{\infty} \chi(t) \mathrm{Op}_{h}^{w}(q) e^{-\frac{i T}{h}\left(H_{h}-z\right)} e^{-\frac{i t}{h}\left(H_{h}-z\right)} f_{h} d t .
\end{aligned}
$$

Note also that since $\left\|U_{h}(T)\right\|_{\mathscr{L}\left(L^{2}\left(\mathbb{R}^{n}\right)\right)}$ is no longer estimated uniformly by 1 , some rests depend on $T$ in the proof of this lemma. This is not a real problem since we take the limit $h \rightarrow 0$ for fixed $T$. We finally obtain that for any compact subset $K$ of $p^{-1}(J)$ and $\varepsilon>0$ there exists $T_{0}$ such that for $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ supported in $K$ we have

$$
\forall T \geqslant T_{0}, \quad \limsup _{h \rightarrow 0}\left|\left\langle\mathrm{Op}_{h}^{w}(q) u_{h}, u_{h}\right\rangle-\left\langle\mathrm{Op}_{h}^{w}(q) u_{h}^{T}, u_{h}^{T}\right\rangle\right| \leqslant \varepsilon\|q\|_{\infty}
$$

With this estimate we can check that for all $q \in C_{0}^{\infty}\left(\mathbb{R}^{2 n}\right)$ the function $T \mapsto \int q d \mu_{T}$ has a limit when $T \rightarrow+\infty$, that this limit defines a non-negative measure $\mu$ on $\mathbb{R}^{2 n}$, and finally that (7.4) holds for this measure. All the properties of $\mu$ stated in Theorem 7.1 are proved as in the dissipative case (in particular we use Theorem 4.1 to prove (ii) b.).

## Appendix A

## Construction of an escape function

In this appendix we prove Proposition 5.1. A similar result (with an inequality) is proved in [17]. The purpose was to give a proof which could be extended for matrix-valued operators. The version we give here is much more convenient in our context.

Let $J=] \frac{E}{2}, 2 E[, \sigma \in] 0, \frac{1}{2}[$ and $\mathcal{R}$ given by Proposition 2.1. We set

$$
Z_{J, \pm}=Z_{ \pm}(\mathcal{R}, 0, \mp \sigma) \cap p^{-1}(J) .
$$

Proposition A.1. - If $\mathscr{R}$ was chosen large enough, then for $\alpha, \beta \in \mathbb{N}^{n}$ such that $|\alpha|+|\beta| \geqslant 1$ there exists $c_{\alpha, \beta}$ such that for $t \geqslant 0$ and $(x, \xi) \in \mathcal{Z}_{J, \pm}$ we have

$$
\begin{equation*}
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} \phi^{ \pm t}(x, \xi)\right| \leqslant c_{\alpha, \beta}\langle t\rangle \tag{A.1}
\end{equation*}
$$

We know (see for instance Lemma IV. 9 in [22]) that the derivatives of the flow $\phi^{t}$ are uniformly bounded as long as $t$ stays in a bounded subset of $\mathbb{R}$, but may grow exponentially fast with time. The purpose of this proposition is to check that if we only look at the flow far from the origin (where it is "almost free") then we recover a growth of size $O(t)$ as in the free case $(x, \xi) \mapsto(x+2 t \xi, \xi)$.

Proof. - 1. We prove the proposition for $(x, \xi) \in Z_{J,+}$, the case $(x, \xi) \in Z_{J,-}$ being analogous. Let

$$
A(t, x, \xi)=\left(\begin{array}{cc}
J_{x} X(t, x, \xi) & J_{\xi} X(t, x, \xi) \\
J_{x} \Xi(t, x, \xi) & J_{\xi} \Xi(t, x, \xi)
\end{array}\right) \in M_{2 n}(\mathbb{R})
$$

where for instance $J_{\xi} X$ denotes the partial jacobian matrix of $X$ with respect to $\xi$. Suppose that

$$
\limsup _{t \rightarrow+\infty} \frac{\|A(t)\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)}}{t}=+\infty
$$

Differentiating the system (1.4) with respect to $x$ and then with respect to $\xi$, we see that

$$
\partial_{t} A(t, x, \xi)=B(t, x, \xi) \cdot A(t, x, \xi)
$$

where

$$
B(t, x, \xi)=\left(\begin{array}{cc}
0 & 2 \mathrm{I}_{n} \\
- \text { Hess } V_{1}(X(t, x, \xi)) & 0
\end{array}\right) \in M_{2 n}(\mathbb{R})
$$

and hence:
$\partial_{t}^{2} A(t, x, \xi)=\partial_{t} B(t, x, \xi) \cdot A(t, x, \xi)+B(t, x, \xi)^{2} \cdot A(t, x, \xi)=: C(t, x, \xi) \cdot A(t, x, \xi)$
According to Proposition 2.1, there exists $c_{0}>0$ such that for $(x, \xi) \in Z_{J,+}$ and $t \geqslant 0$ we have $|X(t, x, \xi)| \geqslant c_{0}(|x|+t)$, so

$$
\|C(t, x, \xi)\| \leqslant\left\|\partial_{t} B(t, x, \xi)\right\|+\left\|B(t, x, \xi)^{2}\right\| \leqslant c(|x|+t)^{-2-\rho}
$$

where $c$ depends neither on $(x, \xi) \in \mathcal{Z}_{J,+}$ nor on $t \geqslant 0$. For $m \in \mathbb{N}$ let

$$
t_{m}=\inf \left\{t \geqslant 1:\|A(t, x, \xi)\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)} \geqslant m t\right\}
$$

Since the derivatives of $\phi^{t}$ are in $L^{\infty}\left(\mathbb{R}^{2 n}\right)$ uniformly for $t$ in a compact subset of $\mathbb{R}$, we have $t_{m} \rightarrow+\infty$. According to Taylor's formula we have

$$
\begin{aligned}
& \frac{\left\|A\left(t_{m}\right)\right\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)}}{t_{m}} \\
& \leqslant \frac{\|A(0)\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)}}{t_{m}}+\left\|\partial_{t} A(0)\right\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)} \\
&+\frac{1}{t_{m}} \int_{0}^{t_{m}}\left(t_{m}-s\right) c(\mathscr{R}+s)^{-2-\rho}\|A(s)\|_{L^{\infty}\left(Z_{J,+}, M_{2 n}(\mathbb{R})\right)} d s \\
& \leqslant c+c \int_{0}^{t_{m}} \frac{t_{m}-s}{t_{m}}(\mathscr{R}+s)^{-2-\rho} s m d s \\
& \leqslant c+c m \int_{0}^{t_{m}}(\mathscr{R}+s)^{-1-\rho} d s \\
& \leqslant c+c m \mathcal{R}^{-\rho}
\end{aligned}
$$

where the different constants $c$ depend neither on $m \in \mathbb{N}$ nor on the choice of $\mathscr{R}$. If $\mathscr{R}$ was chosen so large that $c \mathscr{R}^{-\rho} \leqslant \frac{1}{4}$, then the right-hand side is less than $m / 2$ for large $m$, which gives a contradiction. The case $|\alpha|+|\beta|=1$ is proved.
2. We now proceed by induction on $|\alpha|+|\beta|$. Let $\alpha, \beta \in \mathbb{N}^{n}$ be such that $|\alpha|+|\beta| \geqslant 2$ and assume that the result is proved for any derivative of order less than $|\alpha|+|\beta|$. For $j \in \llbracket 1, n \rrbracket$, the differential operator $\partial_{x}^{\alpha} \partial_{\xi}^{\beta}$ applied to (1.4) gives

$$
\left\{\begin{array}{l}
\partial_{t} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} X_{j}(t, x, \xi)=2 \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{j}(t, x, \xi) \\
\partial_{t} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{j}(t, x, \xi)=-\sum_{l=1}^{n}\left(\partial_{x_{j}} \partial_{x_{l}} V_{1}\right)(X(t, x, \xi)) \partial_{x}^{\alpha} \partial_{\xi}^{\beta} X_{l}(t, x, \xi)+b_{\alpha, \beta, j}(t, x, \xi)
\end{array}\right.
$$

where $b_{\alpha, \beta, j}$ is a sum of terms of the form

$$
-\left(\partial_{x_{j}} \partial^{\nu} V_{1}\right)(X(t, x, \xi)) \prod_{k=1}^{|\nu|} \partial_{x}^{\alpha_{k}} \partial_{\xi}^{\beta_{k}} X_{j_{k}}(t, x, \xi)
$$

where $|\nu| \geqslant 2, \sum_{k=1}^{|\nu|} \alpha_{k}=\alpha, \sum_{k=1}^{|\nu|} \beta_{k}=\beta$ and for $k \in \llbracket 1,|\nu| \rrbracket: j_{k} \in \llbracket 1, n \rrbracket$ and $\left|\alpha_{k}\right|+\left|\beta_{k}\right| \geqslant 1$. In particular for all $k$ we have $\left|\alpha_{k}\right|+\left|\beta_{k}\right|<|\alpha|+|\beta|$, so each term is estimated by

$$
\begin{aligned}
& \left|\left(\partial_{x_{j}} \partial^{\nu} V_{1}\right)(X(t, x, \xi)) \prod_{k=1}^{|\nu|} \partial_{x}^{\alpha_{k}} \partial_{\xi}^{\beta_{k}} X_{j_{k}}(t, x, \xi)\right| \\
& \quad \leqslant c(|x|+t)^{-1-\rho-|\nu|}\langle t\rangle^{\nu} \leqslant c(|x|+t)^{-1-\rho}
\end{aligned}
$$

where $c$ depends neither on $t \geqslant 0$ nor on $(x, \xi) \in Z_{J,+}$, and hence

$$
\left|b_{\alpha, \beta, j}(t, x, \xi)\right| \leqslant c(|x|+t)^{-1-\rho} .
$$

We also have

$$
\left|\partial_{t} b_{\alpha, \beta, j}(t, x, \xi)\right| \leqslant c(|x|+t)^{-1-\rho}
$$

If we set
$A_{\alpha, \beta}(t, x, \xi)=\left(\begin{array}{c}\partial_{x}^{\alpha} \partial_{\xi}^{\beta} X_{1}(t, x, \xi) \\ \vdots \\ \partial_{x}^{\alpha} \partial_{\xi}^{\beta} X_{n}(t, x, \xi) \\ \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{1}(t, x, \xi) \\ \vdots \\ \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{n}(t, x, \xi)\end{array}\right) \in \mathbb{R}^{2 n} \quad$ and $\quad D_{\alpha, \beta}(t, x, \xi)=\left(\begin{array}{c}0 \\ \vdots \\ 0 \\ b_{\alpha, \beta, 1}(t, x, \xi) \\ \vdots \\ b_{\alpha, \beta, n}(t, x, \xi)\end{array}\right) \in \mathbb{R}^{2 n}$,
we have

$$
\partial_{t} A_{\alpha, \beta}(t, x, \xi)=B(t, x, \xi) \cdot A_{\alpha, \beta}(t, x, \xi)+D_{\alpha, \beta}(t, x, \xi)
$$

and
$\partial_{t}^{2} A_{\alpha, \beta}(t, x, \xi)=C(t, x, \xi) \cdot A_{\alpha, \beta}(t, x, \xi)+B(t, x, \xi) \cdot D_{\alpha, \beta}(t, x, \xi)+\partial_{t} D_{\alpha, \beta}(t, x, \xi)$.
Then we can conclude as above. Note that the matrix $C$ is the same, and hence the choice of $\mathscr{R}$ does not depend on $(\alpha, \beta)$.

The estimate of size $O(t)$ for $X$ is what was expected since this is indeed what we have in the free case, but we can improve the result for $\Xi$ :

Corollary A.2. - For any $\alpha, \beta \in \mathbb{N}^{n}$ there exists $c_{\alpha, \beta} \geqslant 0$ such that for $t \geqslant 0$ and $(x, \xi) \in \mathcal{Z}_{J, \pm}$ we have

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi( \pm t, x, \xi)\right| \leqslant c_{\alpha, \beta} .
$$

Proof. - Let $\alpha, \beta \in \mathbb{N}^{n}$. We have proved that

$$
\begin{aligned}
& \partial_{t} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{j}( \pm t, x, \xi) \\
& \quad=\mp \sum_{l=1}^{n}\left(\partial_{x_{j}} \partial_{x_{l}} V_{1}\right)(X( \pm t, x, \xi)) \partial_{x}^{\alpha} \partial_{\xi}^{\beta} X_{l}( \pm t, x, \xi)+\underset{t \rightarrow+\infty}{O}\left(t^{-1-\rho}\right)
\end{aligned}
$$

where the rest is uniform in $(x, \xi) \in Z_{J, \pm}$. With the estimates we now have on the derivatives of $X$ this means that there exists $c_{\alpha, \beta} \geqslant 0$ such that for all $(x, \xi) \in Z_{J, \pm}$ and $t \geqslant 0$ we have

$$
\left|\partial_{t} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} \Xi_{j}( \pm t, x, \xi)\right| \leqslant c_{\alpha, \beta}\langle t\rangle^{-1-\rho}
$$

It only remains to integrate in time to conclude.

Corollary A.3. - For any $\alpha, \beta \in \mathbb{N}^{n}$ there exists $c_{\alpha, \beta} \geqslant 0$ such that for $t \geqslant 0$ and $(x, \xi) \in \mathcal{Z}_{J, \pm}$ we have

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} \frac{X( \pm t, x, \xi) \cdot \Xi( \pm t, x, \xi)}{|X( \pm t, x, \xi)||\Xi( \pm t, x, \xi)|}\right| \leqslant c_{\alpha, \beta} .
$$

Corollary A.4. - Let $\delta>\frac{1}{2}$. Then for any $\alpha, \beta \in \mathbb{N}^{n}$ there exists $c_{\alpha, \beta} \geqslant 0$ such that for all $t \geqslant 0$ and $(x, \xi) \in \mathcal{Z}_{J, \pm}$ we have

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta}\langle X( \pm t, x, \xi)\rangle^{-2 \delta}\right| \leqslant c_{\alpha, \beta}(|x|+t)^{-2 \delta}
$$

Proof. - Let $\alpha, \beta \in \mathbb{N}^{n}$. We remark that $\partial_{x}^{\alpha} \partial_{\xi}^{\beta}\langle X( \pm t, x, \xi)\rangle^{-2 \delta}$ is a sum of terms of the form

$$
c_{K}(X( \pm t, x, \xi))\langle X( \pm t, x, \xi)\rangle^{-2 \delta-K} \prod_{k=1}^{K} \partial_{x}^{\alpha_{k}} \partial_{\xi}^{\beta_{k}} X( \pm t, x, \xi)
$$

where $K \in \llbracket 1,|\alpha|+|\beta| \rrbracket, c_{K}(x)=\langle x\rangle^{2 \delta+K} \frac{d^{K}}{d x^{K}}\langle x\rangle^{-2 \delta}$ is bounded, $\alpha=\sum_{k=1}^{K} \alpha_{k}$ and $\beta=\sum_{k=1}^{K} \beta_{k}$. Since $\langle X( \pm t, x, \xi)\rangle^{-2 \delta-K} \leqslant c(|x|+t)^{-2 d-K}$ for some $c \geqslant 0$, we only have to apply Proposition A. 1 to conclude.

## Now we can prove Proposition 5.1 :

Proof. - Let $\tilde{\chi} \in C_{0}^{\infty}(\mathbb{R})$ be supported in $J$ and equal to 1 in a neighborhood of $E$. Let $\chi_{+}, \chi_{-} \in C^{\infty}(\mathbb{R})$ such that supp $\left.\chi_{+} \subset\right]-\sigma,+\infty\left[\right.$, supp $\left.\chi_{-} \subset\right]-\infty, \sigma[$ and $\chi_{+}+\chi_{-}=1$ on $\mathbb{R}$. Let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ be equal to 1 on $B_{\mathcal{R}+1}$. Let

$$
g_{ \pm}:(x, \xi) \mapsto \chi_{ \pm}\left(\frac{x \cdot \xi}{|x||\xi|}\right)(1-\chi(x)) \tilde{\chi}(p(x, \xi))\langle x\rangle^{-2 \delta}
$$

and, for $w \in \mathbb{R}^{2 n}$ :

$$
f_{ \pm}(w)= \pm \int_{0}^{+\infty} g_{ \pm}\left(\phi^{\mp t}(w)\right) d t
$$

In particular the functions $g_{ \pm}$are supported in $Z_{J, \pm}$. Let $w=(x, \xi) \in \mathbb{R}^{2 n}$. There exists $T_{w} \geqslant 0$ such that $\phi^{ \pm}\left(Z_{J, \pm}\right) \cap B_{x}(2|x|)=\varnothing$ for all $t \geqslant T_{w}$ and hence

$$
\forall v \in B_{x}(2|x|), \forall t \geqslant T_{w}, \quad g_{ \pm}\left(\phi^{\mp t}(v)\right)=0
$$

According to the regularity theorems under the integral sign the functions $f_{+}$ and $f_{-}$are smooth around $w$. Since this holds for any $w \in \mathbb{R}^{2 n}, f_{+}$and $f_{-}$are smooth on $\mathbb{R}^{2 n}$. Moreover their derivatives along the flow $\phi^{t}$ are given by

$$
\left\{p, f_{ \pm}\right\}= \pm \int_{0}^{+\infty}\left\{p, g_{ \pm} \circ \phi^{\mp t}\right\} d t=g_{ \pm}
$$

We now check that all the derivatives of $f_{ \pm}$are bounded. According to the estimates proved above, for $\alpha, \beta \in \mathbb{N}^{n}$ there exists $c_{\alpha, \beta} \geqslant 0$ such that for $(x, \xi) \in Z_{J, \pm}$ and $t \geqslant 0$ we have

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta}\left(g_{ \pm} \circ \phi^{ \pm t}\right)(x, \xi)\right| \leqslant c_{\alpha, \beta}(|x|+t)^{-2 \delta}
$$

Let $w \in \mathbb{R}^{2 n}$ such that $\phi^{\mp t}(w) \in Z_{J, \pm}$ for some $t \geqslant 0$ (otherwise the derivatives of $f_{ \pm}$vanish at $w$ ). Let $t_{0}$ be the maximum of such times $t$. We have

$$
\begin{aligned}
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} f_{ \pm}(w)\right| & \leqslant \int_{0}^{t_{0}}\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta}\left(g_{ \pm} \circ \phi^{\mp t}\right)(w)\right| d t \\
& =\int_{0}^{t_{0}}\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta}\left(g_{ \pm} \circ \phi^{ \pm t}\right)\left(\phi^{\mp t_{0}}(w)\right)\right| d t \\
& \leqslant c_{\alpha, \beta} \int_{0}^{+\infty}(\mathscr{R}+t)^{-2 \delta} d t .
\end{aligned}
$$

This means that $f_{+}$and $f_{-}$belong to $C_{b}^{\infty}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$. It only remains to set $f=f_{+}+f_{-}$to conclude.

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    Julien Royer, Institut de mathématiques de Toulouse, 118, route de Narbonne, 31062
    Toulouse Cédex 9, France - E-mail: julien.royer@math.univ-toulouse.fr 2010 Mathematics Subject Classification. - 35J10, 47A10, 47A55, 47B44, 47G30, 81Q20.
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