# INVERSE SCATTERING AT FIXED ENERGY ON SURFACES WITH EUCLIDEAN ENDS 

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#### Abstract

On a fixed Riemann surface ( $M_{0}, g_{0}$ ) with $N$ Euclidean ends and genus $g$, we show that, under a topological condition, the scattering matrix $S_{V}(\lambda)$ at frequency $\lambda>0$ for the operator $\Delta+V$ determines the potential $V$ if $V \in C^{1, \alpha}\left(M_{0}\right) \cap e^{-\gamma d\left(\cdot, z_{0}\right)^{j}} L^{\infty}\left(M_{0}\right)$ for all $\gamma>0$ and for some $j \in\{1,2\}$, where $d\left(z, z_{0}\right)$ denotes the distance from $z$ to a fixed point $z_{0} \in M_{0}$. The topological condition is given by $N \geq \max (2 g+1,2)$ for $j=1$ and by $N \geq g+1$ if $j=2$. In $\mathbb{R}^{2}$ this implies that the operator $S_{V}(\lambda)$ determines any $C^{1, \alpha}$ potential $V$ such that $V(z)=O\left(e^{-\gamma|z|^{2}}\right)$ for all $\gamma>0$.


## 1. Introduction

The purpose of this paper is to prove an inverse scattering result at fixed frequency $\lambda>0$ in dimension 2 . The typical question one can ask is to show that the scattering matrix $S_{V}(\lambda)$ for the Schrödinger operator $\Delta+V$ determines the potential. This is known to be false if $V$ is only assumed to be Schwartz, by the example of Grinevich-Novikov [6], but it is also known to be true for exponentially decaying potentials (i.e. $V \in e^{-\gamma|z|} L^{\infty}\left(\mathbb{R}^{2}\right)$ for some $\gamma>0$ ) with norm smaller than a constant depending on the frequency $\lambda$, see Novikov [15]. For other partial results we refer to [2], [10], [19], [20], [21]. The determinacy of $V$ from $S_{V}(\lambda)$ when $V$ is compactly supported, without any smallness assumption on the norm, follows from the recent work of Bukhgeim [1] on the inverse boundary problem after a standard reduction to the Dirichlet-to-Neumann operator on a large sphere (see [25] for this reduction).

In dimensions $n \geq 3$, it is proved in Novikov [16] (see also [3] for the case of magnetic Schrödinger operators) that the scattering matrix at a fixed frequency $\lambda$ determines an exponentially decaying potential. When $V$ is compactly supported this also follows directly from the result by Sylvester-Uhlmann [22] on the inverse boundary problem, by reducing to the Dirichlet-to-Neumann operator on a large sphere. Melrose [14] gave a direct proof of the last result based on the methods of [22], and this proof was extended to exponentially decaying potentials in [26] and to the magnetic case in [17]. In the geometric scattering setting, [11, 12] reconstruct the asymptotic expansion of a potential or metrics from the scattering operator at fixed frequency on asymptotically Euclidean/hyperbolic manifolds. Further results of this type are given in [27, 28].

The method for proving the determinacy of $V$ from $S_{V}(\lambda)$ in $[14,26]$ is based on the construction of complex geometric optics solutions $u(z)=e^{\rho \cdot z}(1+r(\rho, z))$ of $\left(\Delta+V-\lambda^{2}\right) u=$ 0 with $\rho \in \mathbb{C}^{n}, z \in \mathbb{R}^{n}$, and the density of the oscillating scattering solutions $u_{\mathrm{sc}}(z)=$ $\int_{S^{n-1}} \Phi_{V}(\lambda, z, \omega) f(\omega) d \omega$ within those complex geometric optics solutions, where $\Phi_{V}(\lambda, z, \omega)=$ $e^{i \lambda \omega \cdot z}+e^{-i \lambda \omega . z}|z|^{-\frac{1}{2}(n-1)} a(\lambda, z, \omega)$ are the perturbed plane wave solutions (here $\omega \in S^{n-1}$ and $a \in L^{\infty}$ ). Unlike when $n \geq 3$, the problem in dimension 2 is that the set of complex geometrical optics solutions of this type is not large enough to show that the Fourier transform of $V_{1}-V_{2}$ is 0 .

The real novelty in the recent work of Bukhgeim [1] in dimension 2 is the construction of new complex geometric optics solutions (at least on a bounded domain $\Omega \subset \mathbb{C}$ ) of $\left(\Delta+V_{i}\right) u_{i}=0$ of the form $u_{1}=e^{\Phi / h}\left(1+r_{1}(h)\right)$ and $u_{2}=e^{-\Phi / h}\left(1+r_{2}(h)\right)$ with $0<h \ll 1$ where $\Phi$ is a holomorphic function in $\mathbb{C}$ with a unique non-degenerate critical point at a fixed $z_{0} \in \mathbb{C}$ (for instance $\left.\Phi(z)=\left(z-z_{0}\right)^{2}\right)$, and $\left\|r_{j}(h)\right\|_{L^{p}}$ is small as $h \rightarrow 0$ for $p>1$. These solutions allow to use stationary phase at $z_{0}$ to get

$$
\int_{\Omega}\left(V_{1}-V_{2}\right) u_{1} \overline{u_{2}}=C\left(V_{1}\left(z_{0}\right)-V_{2}\left(z_{0}\right)\right) h+o(h), \quad C \neq 0
$$

as $h \rightarrow 0$ and thus, if the Dirichlet-to-Neumann operators on $\partial \Omega$ are the same, then $V_{1}\left(z_{0}\right)=$ $V_{2}\left(z_{0}\right)$.

One of the problems to extend this to inverse scattering is that a holomorphic function in $\mathbb{C}$ with a non-degenerate critical point needs to grow at least quadratically at infinity, which would somehow force to consider potentials $V$ having Gaussian decay. On the other hand, if we allow the function to be meromorphic with simple poles, then we can construct such functions, having a single critical point at any given point $p$, for instance by considering $\Phi(z)=(z-p)^{2} / z$. Of course, with such $\Phi$ we then need to work on $\mathbb{C} \backslash\{0\}$, which is conformal to a surface with no hole but with 2 Euclidean ends, and $\Phi$ has linear growth in the ends. In general, on a surface with genus $g$ and $N$ Euclidean ends, we can use the Riemann-Roch theorem to construct holomorphic functions with linear or quadratic growth in the ends, the dimension of the space of such functions depending on $g, N$.

In the present work, we apply this idea to obtain an inverse scattering result for $\Delta_{g_{0}}+V$ on a fixed Riemann surface ( $M_{0}, g_{0}$ ) with Euclidean ends, under some topological condition on $M_{0}$ and some decay condition on $V$.
Theorem 1.1. Let $\left(M_{0}, g_{0}\right)$ be a non-compact Riemann surface with genus $g$ and $N$ ends isometric to $\mathbb{R}^{2} \backslash\{|z| \leq 1\}$ with metric $|d z|^{2}$. Let $V_{1}$ and $V_{2}$ be two potentials in $C^{1, \alpha}\left(M_{0}\right)$ with $\alpha>0$, and such that $S_{V_{1}}(\lambda)=S_{V_{2}}(\lambda)$ for some $\lambda>0$. Let $d\left(z, z_{0}\right)$ denote the distance between $z$ and a fixed point $z_{0} \in M_{0}$.
(i) If $N \geq \max (2 g+1,2)$ and $V_{i} \in e^{-\gamma d\left(\cdot, z_{0}\right)} L^{\infty}\left(M_{0}\right)$ for all $\gamma>0$, then $V_{1}=V_{2}$.
(ii) If $N \geq g+1$ and $V_{i} \in e^{-\gamma d\left(\cdot, z_{0}\right)^{2}} L^{\infty}\left(M_{0}\right)$ for all $\gamma>0$, then $V_{1}=V_{2}$.

In $\mathbb{R}^{2}$, where $g=0$ and $N=1$, we have an immediate corollary:
Corollary 1.2. Let $\lambda>0$ and let $V_{1}, V_{2} \in C^{1, \alpha}\left(\mathbb{R}^{2}\right) \cap e^{-\gamma|z|^{2}} L^{\infty}\left(\mathbb{R}^{2}\right)$ for all $\gamma>0$. If the scattering matrices satisfy $S_{V_{1}}(\lambda)=S_{V_{2}}(\lambda)$, then $V_{1}=V_{2}$.

This is an improvement on the result of Bukhgeim [1] which shows identifiability for compactly supported functions, and in a certain sense on the result of Novikov [15] since it is assumed there that the potential has to be of small $L^{\infty}$ norm.

The structure of the paper is as follows. In Section 2 we employ the Riemann-Roch theorem and a transversality argument to construct Morse holomorphic functions on ( $M_{0}, g_{0}$ ) with linear or quadratic growth in the ends. Section 3 considers Carleman estimates with harmonic weights on ( $M_{0}, g_{0}$ ), where suitable convexification and weights at the ends are required since the surface is non compact. Complex geometrical optics solutions are constructed in Section 4. Section 5 discusses direct scattering theory on surfaces with Euclidean ends and contains the proof that scattering solutions are dense in the set of suitable solutions, and Section 6 gives the proof of Theorem 1.1. Finally, there is an appendix discussing a Paley-Wiener type result for functions with Gaussian decay which is needed to prove density of scattering solutions.

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## 2. Holomorphic Morse functions on a surface with Euclidean ends

2.1. Riemann surfaces with Euclidean ends. Let $\left(M_{0}, g_{0}\right)$ be a non-compact connected smooth Riemannian surface with $N$ ends $E_{1}, \ldots, E_{N}$ which are Euclidean, i.e. isometric to $\mathbb{C} \backslash\{|z| \leq 1\}$ with metric $|d z|^{2}$. By using a complex inversion $z \rightarrow 1 / z$, each end is also isometric to a pointed disk

$$
E_{i} \simeq\{|z| \leq 1, z \neq 0\} \text { with metric } \frac{|d z|^{2}}{|z|^{4}}
$$

thus conformal to the Euclidean metric on the pointed disk. The surface $M_{0}$ can then be compactified by adding the points corresponding to $z=0$ in each pointed disk corresponding to an end $E_{i}$, we obtain a closed Riemann surface $M$ with a natural complex structure induced by that of $M_{0}$, or equivalently a smooth conformal class on $M$ induced by that of $M_{0}$. Another way of thinking is to say that $M_{0}$ is the closed Riemann surface $M$ with $N$ points $e_{1}, \ldots, e_{N}$ removed. The Riemann surface $M$ has holomorphic charts $z_{\alpha}: U_{\alpha} \rightarrow \mathbb{C}$ and we will denote by $z_{1}, \ldots z_{N}$ the complex coordinates corresponding to the ends of $M_{0}$, or equivalently to the neighbourhoods of the points $e_{i}$. The Hodge star operator $\star$ acts on the cotangent bundle $T^{*} M$, its eigenvalues are $\pm i$ and the respective eigenspaces $T_{1,0}^{*} M:=\operatorname{ker}(\star+i \mathrm{Id})$ and $T_{0,1}^{*} M:=\operatorname{ker}(\star-i \mathrm{Id})$ are sub-bundles of the complexified cotangent bundle $\mathbb{C} T^{*} M$ and the splitting $\mathbb{C} T^{*} M=T_{1,0}^{*} M \oplus T_{0,1}^{*} M$ holds as complex vector spaces. Since $\star$ is conformally invariant on 1 -forms on $M$, the complex structure depends only on the conformal class of $g$. In holomorphic coordinates $z=x+i y$ in a chart $U_{\alpha}$, one has $\star(u d x+v d y)=-v d x+u d y$ and

$$
\left.T_{1,0}^{*} M\right|_{U_{\alpha}} \simeq \mathbb{C} d z,\left.\quad T_{0,1}^{*} M\right|_{U_{\alpha}} \simeq \mathbb{C} d \bar{z}
$$

where $d z=d x+i d y$ and $d \bar{z}=d x-i d y$. We define the natural projections induced by the splitting of $\mathbb{C} T^{*} M$

$$
\pi_{1,0}: \mathbb{C} T^{*} M \rightarrow T_{1,0}^{*} M, \quad \pi_{0,1}: \mathbb{C} T^{*} M \rightarrow T_{0,1}^{*} M
$$

The exterior derivative $d$ defines the de Rham complex $0 \rightarrow \Lambda^{0} \rightarrow \Lambda^{1} \rightarrow \Lambda^{2} \rightarrow 0$ where $\Lambda^{k}:=$ $\Lambda^{k} T^{*} M$ denotes the real bundle of $k$-forms on $M$. Let us denote $\mathbb{C} \Lambda^{k}$ the complexification of $\Lambda^{k}$, then the $\partial$ and $\bar{\partial}$ operators can be defined as differential operators $\partial: \mathbb{C} \Lambda^{0} \rightarrow T_{1,0}^{*} M$ and $\bar{\partial}: \mathbb{C} \Lambda 0 \rightarrow T_{0,1}^{*} M$ by

$$
\partial f:=\pi_{1,0} d f, \quad \bar{\partial} f:=\pi_{0,1} d f
$$

they satisfy $d=\partial+\bar{\partial}$ and are expressed in holomorphic coordinates by

$$
\partial f=\partial_{z} f d z, \quad \bar{\partial} f=\partial_{\bar{z}} f d \bar{z}
$$

with $\partial_{z}:=\frac{1}{2}\left(\partial_{x}-i \partial_{y}\right)$ and $\partial_{\bar{z}}:=\frac{1}{2}\left(\partial_{x}+i \partial_{y}\right)$. Similarly, one can define the $\partial$ and $\bar{\partial}$ operators from $\mathbb{C} \Lambda^{1}$ to $\mathbb{C} \Lambda^{2}$ by setting

$$
\partial\left(\omega_{1,0}+\omega_{0,1}\right):=d \omega_{0,1}, \quad \bar{\partial}\left(\omega_{1,0}+\omega_{0,1}\right):=d \omega_{1,0}
$$

if $\omega_{0,1} \in T_{0,1}^{*} M$ and $\omega_{1,0} \in T_{1,0}^{*} M$. In coordinates this is simply

$$
\partial(u d z+v d \bar{z})=\partial v \wedge d \bar{z}, \quad \bar{\partial}(u d z+v d \bar{z})=\bar{\partial} u \wedge d z
$$

If $g$ is a metric on $M$ whose conformal class induces the complex structure $T_{1,0}^{*} M$, there is a natural operator, the Laplacian acting on functions and defined by

$$
\Delta f:=-2 i \star \bar{\partial} \partial f=d^{*} d
$$

where $d^{*}$ is the adjoint of $d$ through the metric $g$ and $\star$ is the Hodge star operator mapping $\Lambda^{2}$ to $\Lambda^{0}$ and induced by $g$ as well.
2.2. Holomorphic functions. We are going to construct Carleman weights given by holomorphic functions on $M_{0}$ which grow at most linearly or quadratically in the ends. We will use the Riemann-Roch theorem, following ideas of [7], however, the difference in the present case is that we have very little freedom to construct these holomorphic functions, simply because there is just a finite dimensional space of such functions by Riemann-Roch. For the convenience of the reader, and to fix notations, we recall the usual Riemann-Roch index theorem (see Farkas-Kra [5] for more details). A divisor $D$ on $M$ is an element

$$
D=\left(\left(p_{1}, n_{1}\right), \ldots,\left(p_{k}, n_{k}\right)\right) \in(M \times \mathbb{Z})^{k}, \text { where } k \in \mathbb{N}
$$

which will also be denoted $D=\prod_{i=1}^{k} p_{i}^{n_{i}}$ or $D=\prod_{p \in M} p^{\alpha(p)}$ where $\alpha(p)=0$ for all $p$ except $\alpha\left(p_{i}\right)=n_{i}$. The inverse divisor of $D$ is defined to be $D^{-1}:=\prod_{p \in M} p^{-\alpha(p)}$ and the degree of the divisor $D$ is defined by $\operatorname{deg}(D):=\sum_{i=1}^{k} n_{i}=\sum_{p \in M} \alpha(p)$. A non-zero meromorphic function on $M$ is said to have divisor $D$ if $(f):=\prod_{p \in M} p^{\operatorname{ord}(p)}$ is equal to $D$, where $\operatorname{ord}(p)$ denotes the order of $p$ as a pole or zero of $f$ (with positive sign convention for zeros). Notice that in this case we have $\operatorname{deg}(f)=0$. For divisors $D^{\prime}=\prod_{p \in M} p^{\alpha^{\prime}(p)}$ and $D=\prod_{p \in M} p^{\alpha(p)}$, we say that $D^{\prime} \geq D$ if $\alpha^{\prime}(p) \geq \alpha(p)$ for all $p \in M$. The same exact notions apply for meromorphic 1-forms on $M$. Then we define for a divisor $D$

$$
\begin{aligned}
r(D) & :=\operatorname{dim}(\{f \text { meromorphic function on } M ;(f) \geq D\} \cup\{0\}) \\
i(D) & :=\operatorname{dim}(\{u \text { meromorphic } 1 \text { form on } M ;(u) \geq D\} \cup\{0\})
\end{aligned}
$$

The Riemann-Roch theorem states the following identity: for any divisor $D$ on the closed Riemann surface $M$ of genus $g$,

$$
\begin{equation*}
r\left(D^{-1}\right)=i(D)+\operatorname{deg}(D)-g+1 \tag{1}
\end{equation*}
$$

Notice also that for any divisor $D$ with $\operatorname{deg}(D)>0$, one has $r(D)=0$ since $\operatorname{deg}(f)=0$ for all $f$ meromorphic. By [5, Th. p70], let $D$ be a divisor, then for any non-zero meromorphic 1-form $\omega$ on $M$, one has

$$
\begin{equation*}
i(D)=r\left(D(\omega)^{-1}\right) \tag{2}
\end{equation*}
$$

which is thus independent of $\omega$. For instance, if $D=1$, we know that the only holomorphic function on $M$ is 1 and one has $1=r(1)=r\left((\omega)^{-1}\right)-g+1$ and thus $r\left((\omega)^{-1}\right)=g$ if $\omega$ is a non-zero meromorphic 1 form. Now if $D=(\omega)$, we obtain again from (1)

$$
g=r\left((\omega)^{-1}\right)=2-g+\operatorname{deg}((\omega))
$$

which gives $\operatorname{deg}((\omega))=2(g-1)$ for any non-zero meromorphic 1-form $\omega$. In particular, if $D$ is a divisor such that $\operatorname{deg}(D)>2(g-1)$, then we get $\operatorname{deg}\left(D(\omega)^{-1}\right)=\operatorname{deg}(D)-2(g-1)>0$ and thus $i(D)=r\left(D(\omega)^{-1}\right)=0$, which implies by (1)

$$
\begin{equation*}
\operatorname{deg}(D)>2(g-1) \Longrightarrow r\left(D^{-1}\right)=\operatorname{deg}(D)-g+1 \geq g \tag{3}
\end{equation*}
$$

Now we deduce the

Lemma 2.1. Let $e_{1}, \ldots, e_{N}$ be distinct points on a closed Riemann surface $M$ with genus $g$, and let $z_{0}$ be another point of $M \backslash\left\{e_{1}, \ldots, e_{N}\right\}$. If $N \geq \max (2 g+1,2)$, the following hold true:
(i) there exists a meromorphic function $f$ on $M$ with at most simple poles, all contained in $\left\{e_{1}, \ldots, e_{N}\right\}$, such that $\partial f\left(z_{0}\right) \neq 0$,
(ii) there exists a meromorphic function $h$ on $M$ with at most simple poles, all contained in $\left\{e_{1}, \ldots, e_{N}\right\}$, such that $z_{0}$ is a zero of order at least 2 of $h$.

Proof. Let first $g \geq 1$, so that $N \geq 2 g+1$. By the discussion before the Lemma, we know that there are at least $g+2$ linearly independent (over $\mathbb{C}$ ) meromorphic functions $f_{0}, \ldots, f_{g+1}$ on $M$ with at most simple poles, all contained in $\left\{e_{1}, \ldots, e_{2 g+1}\right\}$. Without loss of generality, one can set $f_{0}=1$ and by linear combinations we can assume that $f_{1}\left(z_{0}\right)=\cdots=f_{g+1}\left(z_{0}\right)=0$. Now consider the divisor $D_{j}=e_{1} \ldots e_{2 g+1} z_{0}^{-j}$ for $j=1,2$, with degree $\operatorname{deg}\left(D_{j}\right)=2 g+1-j$, then by the Riemann-Roch formula (more precisely (3))

$$
r\left(D_{j}^{-1}\right)=g+2-j
$$

Thus, since $r\left(D_{1}^{-1}\right)>r\left(D_{2}^{-1}\right)=g$ and using the assumption that $g \geq 1$, we deduce that there is a function in $\operatorname{span}\left(f_{1}, \ldots, f_{g+1}\right)$ which has a zero of order 2 at $z_{0}$ and a function which has a zero of order exactly 1 at $z_{0}$. The same method clearly works if $g=0$ by taking two points $e_{1}, e_{2}$ instead of just $e_{1}$.

If we allow double poles instead of simple poles, the proof of Lemma 2.1 shows the
Lemma 2.2. Let $e_{1}, \ldots, e_{N}$ be distinct points on a closed Riemann surface $M$ with genus $g$, and let $z_{0}$ be another point of $M \backslash\left\{e_{1}, \ldots, e_{N}\right\}$. If $N \geq g+1$, then the following hold true: (i) there exists a meromorphic function $f$ on $M$ with at most double poles, all contained in $\left\{e_{1}, \ldots, e_{N}\right\}$, such that $\partial f\left(z_{0}\right) \neq 0$,
(ii) there exists a meromorphic function $h$ on $M$ with at most double poles, all contained in $\left\{e_{1}, \ldots, e_{N}\right\}$, such that $z_{0}$ is a zero of order at least 2 of $h$.
2.3. Morse holomorphic functions with prescribed critical points. We follow in this section the arguments used in [7] to construct holomorphic functions with non-degenerate critical points (i.e. Morse holomorphic functions) on the surface $M_{0}$ with genus $g$ and $N$ ends, such that these functions have at most linear growth (resp. quadratic growth) in the ends if $N \geq \max (2 g+1,2)$ (resp. if $N \geq g+1)$. We let $\mathcal{H}$ be the complex vector space spanned by the meromorphic functions on $M$ with divisors larger or equal to $e_{1}^{-1} \ldots e_{N}^{-1}$ (resp. by $e_{1}^{-2} \ldots e_{N}^{-2}$ ) if we work with functions having linear growth (resp. quadratic growth), where $e_{1}, \ldots e_{N} \in M$ are points corresponding to the ends of $M_{0}$ as explained in Section 2. Note that $\mathcal{H}$ is a complex vector space of complex dimension greater or equal to $N-g+1$ (resp. $2 N-g+1$ ) for the $e_{1}^{-1} \ldots e_{N}^{-1}$ divisor (resp. the $e_{1}^{-2} \ldots e_{N}^{-2}$ divisor). We will also consider the real vector space $H$ spanned by the real parts and imaginary parts of functions in $\mathcal{H}$, this is a real vector space which admits a Lebesgue measure. We now prove the following

Lemma 2.3. The set of functions $u \in H$ which are not Morse in $M_{0}$ has measure 0 in $H$, in particular its complement is dense in $H$.

Proof. We use an argument very similar to that used by Uhlenbeck [24]. We start by defining $m: M_{0} \times H \rightarrow T^{*} M_{0}$ by $(p, u) \mapsto(p, d u(p)) \in T_{p}^{*} M_{0}$. This is clearly a smooth map, linear in the second variable, moreover $m_{u}:=m(., u)=(\cdot, d u(\cdot))$ is smooth on $M_{0}$. The map $u$ is a

Morse function if and only if $m_{u}$ is transverse to the zero section, denoted $T_{0}^{*} M_{0}$, of $T^{*} M_{0}$, i.e. if

$$
\operatorname{Image}\left(D_{p} m_{u}\right)+T_{m_{u}(p)}\left(T_{0}^{*} M_{0}\right)=T_{m_{u}(p)}\left(T^{*} M_{0}\right), \quad \forall p \in M_{0} \text { such that } m_{u}(p)=(p, 0)
$$

This is equivalent to the fact that the Hessian of $u$ at critical points is non-degenerate (see for instance Lemma 2.8 of [24]). We recall the following transversality result, the proof of which is contained in [24, Th.2] by replacing Sard-Smale theorem by the usual finite dimensional Sard theorem:

Theorem 2.4. Let $m: X \times H \rightarrow W$ be a $C^{k}$ map and $X, W$ be smooth manifolds and $H$ a finite dimensional vector space, if $W^{\prime} \subset W$ is a submanifold such that $k>\max (1, \operatorname{dim} X-$ $\left.\operatorname{dim} W+\operatorname{dim} W^{\prime}\right)$, then the transversality of the map $m$ to $W^{\prime}$ implies that the complement of the set $\left\{u \in H ; m_{u}\right.$ is transverse to $\left.W^{\prime}\right\}$ in $H$ has Lebesgue measure 0 .

We want to apply this result with $X:=M_{0}, W:=T^{*} M_{0}$ and $W^{\prime}:=T_{0}^{*} M_{0}$, and with the map $m$ as defined above. We have thus proved our Lemma if one can show that $m$ is transverse to $W^{\prime}$. Let $(p, u)$ such that $m(p, u)=(p, 0) \in W^{\prime}$. Then identifying $T_{(p, 0)}\left(T^{*} M_{0}\right)$ with $T_{p} M_{0} \oplus T_{p}^{*} M_{0}$, one has

$$
D m_{(p, u)}(z, v)=\left(z, d v(p)+\operatorname{Hess}_{p}(u) z\right)
$$

where $\operatorname{Hess}_{p}(u)$ is the Hessian of $u$ at the point $p$, viewed as a linear map from $T_{p} M_{0}$ to $T_{p}^{*} M_{0}$ (note that this is different from the covariant Hessian defined by the Levi-Civita connection). To prove that $m$ is transverse to $W^{\prime}$ we need to show that $(z, v) \rightarrow\left(z, d v(p)+\operatorname{Hess}_{p}(u) z\right)$ is onto from $T_{p} M_{0} \oplus H$ to $T_{p} M_{0} \oplus T_{p}^{*} M_{0}$, which is realized if the map $v \rightarrow d v(p)$ from $H$ to $T_{p}^{*} M_{0}$ is onto. But from Lemma 2.1, we know that there exists a meromorphic function $f$ with real part $v=\operatorname{Re}(f) \in H$ such that $v(p)=0$ and $d v(p) \neq 0$ as an element of $T_{p}^{*} M_{0}$. We can then take $v_{1}:=v$ and $v_{2}:=\operatorname{Im}(f)$, which are functions of $H$ such that $d v_{1}(p)$ and $d v_{2}(p)$ are linearly independent in $T_{p}^{*} M_{0}$ by the Cauchy-Riemann equation $\bar{\partial} f=0$. This shows our claim and ends the proof by using Theorem 2.4.

In particular, by the Cauchy-Riemann equation, this Lemma implies that the set of Morse functions in $\mathcal{H}$ is dense in $\mathcal{H}$. We deduce

Proposition 2.1. There exists a dense set of points $p$ in $M_{0}$ such that there exists a Morse holomorphic function $f \in \mathcal{H}$ on $M_{0}$ which has a critical point at $p$.

Proof. Let $p$ be a point of $M_{0}$ and let $u$ be a holomorphic function with a zero of order at least 2 at $p$, the existence is ensured by Lemma 2.1. Let $B(p, \eta)$ be a any small ball of radius $\eta>0$ near $p$, then by Lemma 2.3, for any $\epsilon>0$, we can approach $u$ by a holomorphic Morse function $u_{\epsilon} \in \mathcal{H}_{\epsilon}$ which is at distance less than $\epsilon$ of $u$ in a fixed norm on the finite dimensional space $\mathcal{H}$. Rouché's theorem for $\partial_{z} u_{\epsilon}$ and $\partial_{z} u$ (which are viewed as functions locally near $p$ ) implies that $\partial_{z} u_{\epsilon}$ has at least one zero of order exactly 1 in $B(p, \eta)$ if $\epsilon$ is chosen small enough. Thus there is a Morse function in $\mathcal{H}$ with a critical point arbitrarily close to $p$.

Remark 2.5. In the case where the surface $M$ has genus 0 and $N$ ends, we have an explicit formula for the function in Proposition 2.1: indeed $M_{0}$ is conformal to $\mathbb{C} \backslash\left\{e_{1}, \ldots, e_{N-1}\right\}$ for some $e_{i} \in \mathbb{C}$ - i.e. the Riemann sphere minus $N$ points - then the function $f(z)=$ $\left(z-z_{0}\right)^{2} /\left(z-e_{1}\right)$ with $z_{0} \notin\left\{e_{1}, \ldots, e_{N-1}\right\}$ has $z_{0}$ for unique critical point in $\mathbb{C} \backslash\left\{e_{1}, \ldots, e_{N-1}\right\}$ and it is non-degenerate.

We end this section by the following Lemmas which will be used for the amplitude of the complex geometric optics solutions but not for the phase.
Lemma 2.6. For any $p_{0}, p_{1}, \ldots p_{n} \in M_{0}$ some points of $M_{0}$ and $L \in \mathbb{N}$, then there exists a function a(z) holomorphic on $M_{0}$ which vanishes to order $L$ at all $p_{j}$ for $j=1, \ldots, n$ and such that $a\left(p_{0}\right) \neq 0$. Moreover $a(z)$ can be chosen to have at most polynomial growth in the ends, i.e. $|a(z)| \leq C|z|^{J}$ for some $J \in \mathbb{N}$.
Proof. It suffices to find on $M$ some meromorphic function with divisor greater or equal to $D:=e_{1}^{-J} \ldots e_{N}^{-J} p_{1}^{L} \ldots p_{n}^{L}$ but not greater or equal to $D p_{0}$ and this is insured by RiemannRoch theorem as long as $J N-n L \geq 2 g$ since then $r(D)=-g+1+J N-n L$ and $r\left(D p_{0}\right)=$ $-g+J N-n L$.

Lemma 2.7. Let $\left\{p_{0}, p_{1}, . ., p_{n}\right\} \subset M_{0}$ be a set of $n+1$ disjoint points. Let $c_{0}, c_{1}, \ldots, c_{K} \in \mathbb{C}$, $L \in \mathbb{N}$, and let $z$ be a complex coordinate near $p_{0}$ such that $p_{0}=\{z=0\}$. Then there exists a holomorphic function $f$ on $M_{0}$ with zeros of order at least $L$ at each $p_{j}$, such that $f(z)=c_{0}+c_{1} z+\ldots+c_{K} z^{K}+O\left(|z|^{K+1}\right)$ in the coordinate $z$. Moreover $f$ can be chosen so that there is $J \in \mathbb{N}$ such that, in the ends, $\left|\partial_{z}^{\ell} f(z)\right|=O\left(|z|^{J}\right)$ for all $\ell \in \mathbb{N}_{0}$.

Proof. The proof goes along the same lines as in Lemma 2.6. By induction on $K$ and linear combinations, it suffices to prove it for $c_{0}=\cdots=c_{K-1}=0$. As in the proof of Lemma 2.6, if $J$ is taken large enough, there exists a function with divisor greater or equal to $D:=e_{1}^{-J} \ldots e_{N}^{-J} p_{0}^{K-1} p_{1}^{L} \ldots p_{n}^{L}$ but not greater or equal to $D p_{0}$. Then it suffices to multiply this function by $c_{K}$ times the inverse of the coefficient of $z^{K}$ in its Taylor expansion at $z=0$.
2.4. Laplacian on weighted spaces. Let $x$ be a smooth positive function on $M_{0}$, which is equal to $|z|^{-1}$ for $|z|>r_{0}$ in the ends $E_{i} \simeq\{z \in \mathbb{C} ;|z|>1\}$, where $r_{0}$ is a large fixed number. We now show that the Laplacian $\Delta_{g_{0}}$ on a surface with Euclidean ends has a right inverse on the weighted spaces $x^{-J} L^{2}\left(M_{0}\right)$ for $J \notin \mathbb{N}$ positive.

Lemma 2.8. For any $J>-1$ which is not an integer, there exists a continuous operator $G$ mapping $x^{-J} L^{2}\left(M_{0}\right)$ to $x^{-J-2} L^{2}\left(M_{0}\right)$ such that $\Delta_{g_{0}} G=\mathrm{Id}$.
Proof. Let $g_{b}:=x^{2} g_{0}$ be a metric conformal to $g_{0}$. The metric $g_{b}$ in the ends can be written $g_{b}=d x^{2} / x^{2}+d \theta_{S^{1}}^{2}$ by using radial coordinates $x=|z|^{-1}, \theta=z /|z| \in S^{1}$, this is thus a b-metric in the sense of Melrose [13], giving the surface a geometry of surface with cylindrical ends. Let us define for $m \in \mathbb{N}_{0}$

$$
H_{b}^{m}\left(M_{0}\right):=\left\{u \in L^{2}\left(M_{0} ; \operatorname{dvol}_{g_{b}}\right) ;\left(x \partial_{x}\right)^{j} \partial_{\theta}^{k} u \in L^{2}\left(M_{0} ; \operatorname{dvol}_{g_{b}}\right) \text { for all } j+k \leq m\right\} .
$$

The Laplacian has the form $\Delta_{g_{b}}=-\left(x \partial_{x}\right)^{2}+\Delta_{S^{1}}$ in the ends, and the indicial roots of $\Delta_{g_{b}}$ in the sense of Section 5.2 of [13] are given by the complex numbers $\lambda$ such that $x^{-i \lambda} \Delta_{g_{b}} x^{i \lambda}$ is not invertible as an operator acting on the circle $S_{\theta}^{1}$. Thus the indicial roots are the solutions of $\lambda^{2}+k^{2}=0$ where $k^{2}$ runs over the eigenvalues of $\Delta_{S^{1}}$, that is, $k \in \mathbb{Z}$. The roots are simple at $\pm i k \in i \mathbb{Z} \backslash\{0\}$ and 0 is a double root. In Theorem 5.60 of [13], Melrose proves that $\Delta_{g_{b}}$ is Fredholm on $x^{a} H_{b}^{2}\left(M_{0}\right)$ if and only if $-a$ is not the imaginary part of some indicial root, that is here $a \notin \mathbb{Z}$. For $J>0$, the kernel of $\Delta_{g_{b}}$ on the space $x^{J} H_{b}^{2}\left(M_{0}\right)$ is clearly trivial by an energy estimate. Thus $\Delta_{g_{b}}: x^{-J} H_{b}^{0}\left(M_{0}\right) \rightarrow x^{-J} H_{b}^{-2}\left(M_{0}\right)$ is surjective for $J>0$ and $J \notin \mathbb{Z}$, and the same then holds for $\Delta_{g_{b}}: x^{-J} H_{b}^{2}\left(M_{0}\right) \rightarrow x^{-J} H_{b}^{0}\left(M_{0}\right)$ by elliptic regularity.

Now we can use Proposition 5.64 of [13], which asserts, for all positive $J \notin \mathbb{Z}$, the existence of a pseudodifferential operator $G_{b}$ mapping continuously $x^{-J} H_{b}^{0}\left(M_{0}\right)$ to $x^{-J} H_{b}^{2}\left(M_{0}\right)$ such that $\Delta_{g_{b}} G_{b}=$ Id. Thus if we set $G=G_{b} x^{-2}$, we have $\Delta_{g_{0}} G=\mathrm{Id}$ and $G$ maps continuously $x^{-J+1} L^{2}\left(M_{0}\right)$ to $x^{-J-1} L^{2}\left(M_{0}\right)$ (note that $L^{2}\left(M_{0}\right)=x H_{b}^{0}\left(M_{0}\right)$ ).

## 3. Carleman Estimate for Harmonic Weights with Critical Points

3.1. The linear weight case. In this section, we prove a Carleman estimate using harmonic weights with non-degenerate critical points, in a way similar to [7]. Here however we need to work on a non compact surface and with weighted spaces. We first consider a Morse holomorphic function $\Phi \in \mathcal{H}$ obtained from Proposition 2.1 with the condition that $\Phi$ has linear growth in the ends, which corresponds to the case where $V \in e^{-\gamma / x} L^{\infty}\left(M_{0}\right)$ for all $\gamma>0$. The Carleman weight will be the harmonic function $\varphi:=\operatorname{Re}(\Phi)$. We let $x$ be a positive smooth function on $M_{0}$ such that $x=|z|^{-1}$ in the complex charts $\{z \in \mathbb{C} ;|z|>1\} \simeq E_{i}$ covering the end $E_{i}$.

Let $\delta \in(0,1)$ be small and let us take $\varphi_{0} \in x^{-\alpha} L^{2}\left(M_{0}\right)$ a solution of $\Delta_{g_{0}} \varphi_{0}=x^{2-\delta}$, a solution exists by Proposition 2.8 if $\alpha>1+\delta$. Actually, by using Proposition 5.61 of [13], if we choose $\alpha<2$, then it is easy to see that $\varphi_{0}$ is smooth on $M_{0}$ and has polyhomogeneous expansion as $|z| \rightarrow \infty$, with leading asymptotic in the end $E_{i}$ given by $\varphi_{0}=-x^{-\delta} / \delta^{2}+c_{i} \log (x)+d_{i}+O(x)$ for some $c_{i}, d_{i}$ which are smooth functions in $S^{1}$. For $\epsilon>0$ small, we define the convexified weight $\varphi_{\epsilon}:=\varphi-\frac{h}{\epsilon} \varphi_{0}$.

We recall from the proof of Proposition 3.1 in [7] the following estimate which is valid in any compact set $K \subset M_{0}$ : for all $w \in C_{0}^{\infty}(K)$, we have

$$
\begin{equation*}
\frac{C}{\epsilon}\left(\frac{1}{h}\|w\|_{L^{2}}^{2}+\frac{1}{h^{2}}\|w|d \varphi|\|_{L^{2}}^{2}+\frac{1}{h^{2}}\left\|w \mid d \varphi_{\epsilon}\right\|_{L^{2}}^{2}+\|d w\|_{L^{2}(K)}^{2}\right) \leq\left\|e^{\varphi_{\epsilon} / h} \Delta_{g} e^{-\varphi_{\epsilon} / h} w\right\|_{L^{2}}^{2} \tag{4}
\end{equation*}
$$

where $C$ depends on $K$ but not on $h$ and $\epsilon$.
So for functions supported in the end $E_{i}$, it clearly suffices to obtain a Carleman estimate in $E_{i} \simeq \mathbb{R}^{2} \backslash\{|z| \leq 1\}$ by using the Euclidean coordinate $z$ of the end.
Proposition 3.1. Let $\delta \in(0,1)$, and $\varphi_{\epsilon}$ as above, then there exists $C>0$ such that for all $\epsilon \gg h>0$ small enough, and all $u \in C_{0}^{\infty}\left(E_{i}\right)$

$$
h^{2}\left\|e^{\varphi_{\epsilon} / h}\left(\Delta-\lambda^{2}\right) e^{-\varphi_{\epsilon} / h} u\right\|_{L^{2}}^{2} \geq \frac{C}{\epsilon}\left(\left\|x^{1-\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{2}\left\|x^{1-\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right) .
$$

Proof. The metric $g_{0}$ can be extended to $\mathbb{R}^{2}$ to be the Euclidean metric and we shall denote by $\Delta$ the flat positive Laplacian on $\mathbb{R}^{2}$. Let us write $P:=\Delta_{g_{0}}-\lambda^{2}$, then the operator $P_{h}:=h^{2} e^{\varphi_{\epsilon} / h} P e^{-\varphi_{\epsilon} / h}$ is given by

$$
P_{h}=h^{2} \Delta-\left|d \varphi_{\epsilon}\right|^{2}+2 h \nabla \varphi_{\epsilon} \cdot \nabla-h \Delta \varphi_{\epsilon}-h^{2} \lambda^{2}
$$

following the notation of [4, Chap. 4.3], it is a semiclassical operator in $S^{0}\left(\langle\xi\rangle^{2}\right)$ with semiclassical full Weyl symbol

$$
\sigma\left(P_{h}\right):=|\xi|^{2}-\left|d \varphi_{\epsilon}\right|^{2}-h^{2} \lambda^{2}+2 i\left\langle d \varphi_{\epsilon}, \xi\right\rangle=a+i b .
$$

We can define $A:=\left(P_{h}+P_{h}^{*}\right) / 2=h^{2} \Delta-\left|d \varphi_{\epsilon}\right|^{2}-h^{2} \lambda^{2}$ and $B:=\left(P_{h}-P_{h}^{*}\right) / 2 i=-2 i h \nabla \varphi_{\epsilon} . \nabla+$ $i h \Delta \varphi_{\epsilon}$ which have respective semiclassical full symbols $a$ and $b$, i.e. $A=\operatorname{Op}_{h}(a)$ and $B=$
$\mathrm{Op}_{h}(b)$ for the Weyl quantization. Notice that $A, B$ are symmetric operators, thus for all $u \in C_{0}^{\infty}\left(E_{i}\right)$

$$
\begin{equation*}
\|(A+i B) u\|^{2}=\left\langle\left(A^{2}+B^{2}+i[A, B]\right) u, u\right\rangle \tag{5}
\end{equation*}
$$

It is easy to check that the operator $i h^{-1}[A, B]$ is a semiclassical differential operator in $S^{0}\left(\langle\xi\rangle^{2}\right)$ with full semiclassical symbol

$$
\begin{equation*}
\{a, b\}(\xi)=4\left(D^{2} \varphi_{\epsilon}\left(d \varphi_{\epsilon}, d \varphi_{\epsilon}\right)+D^{2} \varphi_{\epsilon}(\xi, \xi)\right) \tag{6}
\end{equation*}
$$

Let us now decompose the Hessian of $\varphi_{\epsilon}$ in the basis $\left(d \varphi_{\epsilon}, \theta\right)$ where $\theta$ is a covector orthogonal to $d \varphi_{\epsilon}$ and of norm $\left|d \varphi_{\epsilon}\right|$. This yields coordinates $\xi=\xi_{0} d \varphi_{\epsilon}+\xi_{1} \theta$ and there exist smooth functions $M, N, K$ so that

$$
D^{2} \varphi_{\epsilon}(\xi, \xi)=\left|d \varphi_{\epsilon}\right|^{2}\left(M \xi_{0}^{2}+N \xi_{1}^{2}+2 K \xi_{0} \xi_{1}\right)
$$

Notice that $\varphi_{\epsilon}$ has a polyhomogeneous expansion at infinity of the form

$$
\varphi_{\epsilon}(z)=\gamma \cdot z+\frac{h}{\epsilon} \frac{r^{\delta}}{\delta^{2}}+c_{1} \log (r)+c_{2}+c_{3} r^{-1}+O\left(r^{-2}\right)
$$

where $r=|z|, \omega=z / r, \gamma=\left(\gamma_{1}, \gamma_{2}\right) \in \mathbb{R}^{2}$ and $c_{i}$ are some smooth functions on $S^{1}$ depending on $h$; in particular we have

$$
d \varphi_{\epsilon}=\gamma_{1} d z_{1}+\gamma_{2} d z_{2}+O\left(r^{-1+\delta}\right), \quad \partial_{z}^{\alpha} \partial_{\bar{z}}^{\beta} \varphi_{\epsilon}(z)=O\left(r^{-2+\delta}\right) \quad \text { for all } \alpha+\beta \geq 2
$$

which implies that $M, N, K \in r^{-2+\delta} L^{\infty}\left(E_{i}\right)$. Then one can write

$$
\begin{aligned}
\{a, b\} & =4\left|d \varphi_{\epsilon}\right|^{2}\left(M+M \xi_{0}^{2}+N \xi_{1}^{2}+2 K \xi_{0} \xi_{1}\right) \\
& =4\left(N\left(a+h^{2} \lambda^{2}\right)+\left((M-N) \xi_{0}+2 K \xi_{1}\right) b / 2+(N+M)\left|d \varphi_{\epsilon}\right|^{2}\right)
\end{aligned}
$$

and since $M+N=\operatorname{Tr}\left(D^{2} \varphi_{\epsilon}\right)=-\Delta \varphi_{\epsilon}=h \Delta \varphi_{0} / \epsilon$ we obtain

$$
\begin{gather*}
\{a, b\}=4\left|d \varphi_{\epsilon}\right|^{2}\left(c(z)\left(a+h^{2} \lambda^{2}\right)+\ell(z, \xi) b+\frac{h}{\epsilon} r^{-2+\delta}\right) \\
c(z)=\frac{N}{\left|d \varphi_{\epsilon}\right|^{2}}, \quad \ell(z, \xi)=\frac{(M-N) \xi_{0}+2 K \xi_{1}}{2\left|d \varphi_{\epsilon}\right|^{2}} \tag{7}
\end{gather*}
$$

Now, we take a smooth extension of $\left|d \varphi_{\epsilon}\right|^{2}, a(z, \xi), \ell(z, \xi)$ and $r$ to $z \in \mathbb{R}^{2}$, this can done for instance by extending $r$ as a smooth positive function on $\mathbb{R}^{2}$ and then extending $d \varphi$ and $d \varphi_{0}$ to smooth non vanishing 1 -forms on $\mathbb{R}^{2}$ (not necessarily exact) so that $\left|d \varphi_{\epsilon}\right|^{2}$ is smooth positive (for small $h$ ) and polynomial in $h$ and $a, \ell$ are of the same form as in $\{|z|>1\}$. Let us define the symbol and quantized differential operator on $\mathbb{R}^{2}$

$$
e:=4\left|d \varphi_{\epsilon}\right|^{2}\left(c(z)\left(a+h^{2} \lambda^{2}\right)+\ell(z, \xi) b\right), \quad E:=\mathrm{Op}_{h}(e)
$$

and write

$$
\begin{align*}
& i h^{-1} r^{1-\frac{\delta}{2}}[A, B] r^{1-\frac{\delta}{2}}=h F+r^{1-\frac{\delta}{2}} E r^{1-\frac{\delta}{2}}-\frac{h}{\epsilon}\left(A^{2}+B^{2}\right) \\
& \text { with } F:=h^{-1} r^{1-\frac{\delta}{2}}\left(i h^{-1}[A, B]-E\right) r^{1-\frac{\delta}{2}}+\frac{1}{\epsilon}\left(A^{2}+B^{2}\right) \tag{8}
\end{align*}
$$

We deduce from (6) and (7) the following

Lemma 3.2. The operator $F$ is a semiclassical differential operator in the class $S^{0}\left(\langle\xi\rangle^{4}\right)$ with semiclassical principal symbol

$$
\sigma(F)(\xi)=\frac{4|d \varphi|^{2}}{\epsilon}+\frac{1}{\epsilon}\left(|\xi|^{2}-|d \varphi|^{2}\right)^{2}+\frac{4}{\epsilon}(\langle\xi, d \varphi\rangle)^{2} .
$$

By the semiclassical Gårding estimate, we obtain the
Corollary 3.3. The operator $F$ of Lemma 3.2 is such that there is a constant $C$ so that

$$
\langle F u, u\rangle \geq \frac{C}{\epsilon}\left(\|u\|_{L^{2}}^{2}+h^{2}\|d u\|_{L^{2}}^{2}\right) .
$$

Proof. It suffices to use that $\sigma(F)(\xi) \geq \frac{C^{\prime}}{\epsilon}\left(1+|\xi|^{4}\right)$ for some $C^{\prime}>0$ and use the semiclassical Gårding estimate.

So by writing $\langle i[A, B] u, u\rangle=\left\langle i r^{1-\frac{\delta}{2}}[A, B] r^{1-\frac{\delta}{2}} r^{-1+\frac{\delta}{2}} u, r^{-1+\frac{\delta}{2}} u\right\rangle$ in (5) and using (8) and Corollary 3.3 , we obtain that there exists $C>0$ such that for all $u \in C_{0}^{\infty}\left(E_{i}\right)$

$$
\begin{align*}
\left\|P_{h} u\right\|_{L^{2}}^{2} \geq & \left\langle\left(A^{2}+B^{2}\right) u, u\right\rangle+\frac{C h^{2}}{\epsilon}\left(\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{2}\left\|r^{-1+\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right)+h\langle E u, u\rangle \\
& -\frac{h^{2}}{\epsilon}\left(\left\|A\left(r^{-1+\frac{\delta}{2}} u\right)\right\|_{L^{2}}^{2}+\left\|B\left(r^{-1+\frac{\delta}{2}} u\right)\right\|_{L^{2}}^{2}\right) \tag{9}
\end{align*}
$$

We observe that $h^{-1}\left[A, r^{-1+\frac{\delta}{2}}\right] r^{1+\frac{\delta}{2}} \in S^{0}(\langle\xi\rangle)$ and $h^{-1}\left[B, r^{-1+\frac{\delta}{2}}\right] r^{1+\frac{\delta}{2}} \in h S^{0}(1)$, and thus
$\left.\left\|A\left(r^{-1+\frac{\delta}{2}} u\right)\right\|_{L^{2}}^{2}+\left\|B\left(r^{-1+\frac{\delta}{2}} u\right)\right\|_{L^{2}}^{2}\right) \leq C^{\prime}\left(\|A u\|_{L^{2}}^{2}+\|B u\|_{L^{2}}^{2}+h^{2}\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{4}\left\|r^{-1+\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right)$
for some $C^{\prime}>0$. Taking $h$ small, this implies with (9) that there exists a new constant $C>0$ such that

$$
\begin{equation*}
\left\|P_{h} u\right\|_{L^{2}}^{2} \geq \frac{1}{2}\left\langle\left(A^{2}+B^{2}\right) u, u\right\rangle+\frac{C h^{2}}{\epsilon}\left(\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{2}\left\|r^{-1+\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right)+h\langle E u, u\rangle \tag{10}
\end{equation*}
$$

It remains to deal with $h\langle E u, u\rangle$ : we first write $E=4\left|d \varphi_{\epsilon}\right|^{2}\left(c(z)\left(A+h^{2} \lambda^{2}\right)+\mathrm{Op}_{h}(\ell) B\right)+$ $h r^{-1+\frac{\delta}{2}} S r^{-1+\frac{\delta}{2}}$ where $S$ is a semiclassical differential operator in the class $S^{0}(\langle\xi\rangle)$ by the decay estimates on $c(z), \ell(z, \xi)$ as $z \rightarrow \infty$, then by Cauchy-Schwartz (and with $L:=\mathrm{Op}_{h}(\ell)$ )

$$
\begin{aligned}
|\langle h E u, u\rangle| & \leq C h\left(\|A u\|_{L^{2}}+h^{2}\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}+h\left\|S r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}\right)\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}+C h\|B u\|_{L^{2}}\|L u\|_{L^{2}} \\
& \leq \frac{1}{4}\|A u\|_{L^{2}}^{2}+h^{2}\left\|S r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+C h^{2}\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+\frac{1}{4}\|B u\|_{L^{2}}^{2}+C h^{2}\|L u\|_{L^{2}}^{2}
\end{aligned}
$$

where $C$ is a constant independent of $h, \epsilon$ but may change from line to line. Now we observe that $L r^{1-\frac{\delta}{2}}$ and $S$ are in $S^{0}(\langle\xi\rangle)$ and thus

$$
\left\|S r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+\|L u\|_{L^{2}}^{2} \leq C\left(\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{2}\left\|r^{-1+\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right)
$$

which by (10) implies that there exists $C>0$ such that for all $\epsilon \gg h>0$ with $\epsilon$ small enough

$$
\left\|P_{h} u\right\|_{L^{2}}^{2} \geq \frac{C h^{2}}{\epsilon}\left(\left\|r^{-1+\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+h^{2}\left\|r^{-1+\frac{\delta}{2}} d u\right\|_{L^{2}}^{2}\right)
$$

for all $u \in C_{0}^{\infty}\left(E_{i}\right)$. The proof is complete.

Combining now Proposition 3.1 and (4), we obtain

Proposition 3.4. Let $\left(M_{0}, g_{0}\right)$ be a Riemann surface with Euclidean ends with $x$ a boundary defining function of the radial compactification $\bar{M}_{0}$ and let $\varphi_{\epsilon}=\varphi-\frac{h}{\epsilon} \varphi_{0}$ where $\varphi$ is a harmonic function with non-degenerate critical points and linear growth on $M_{0}$ and $\varphi_{0}$ satisfies $\Delta_{g_{0}} \varphi_{0}=x^{2-\delta}$ as above. Then for all $V \in x^{1-\frac{\delta}{2}} L^{\infty}\left(M_{0}\right)$ there exists an $h_{0}>0, \epsilon_{0}$ and $C>0$ such that for all $0<h<h_{0}, h \ll \epsilon<\epsilon_{0}$ and $u \in C_{0}^{\infty}\left(M_{0}\right)$, we have

$$
\begin{equation*}
\frac{1}{h}\left\|x^{1-\frac{\delta}{2}} u\right\|_{L^{2}}^{2}+\frac{1}{h^{2}}\left\|x^{1-\frac{\delta}{2}} u|d \varphi|\right\|_{L^{2}}^{2}+\left\|x^{1-\frac{\delta}{2}} d u\right\|_{L^{2}}^{2} \leq C \epsilon\left\|e^{\varphi_{\epsilon} / h}\left(\Delta_{g}+V-\lambda^{2}\right) e^{-\varphi_{\epsilon} / h} u\right\|_{L^{2}}^{2} \tag{11}
\end{equation*}
$$

Proof. As in the proof of Proposition 3.1 in [7], by taking $\epsilon$ small enough, we see that the combination of (4) and Proposition 3.1 shows that for any $w \in C_{0}^{\infty}\left(M_{0}\right)$,

$$
\begin{gathered}
\frac{C}{\epsilon}\left(\frac{1}{h}\left\|x^{1-\frac{\delta}{2}} w\right\|_{L^{2}}^{2}+\frac{1}{h^{2}}\left\|x^{1-\frac{\delta}{2}} w|d \varphi|\right\|_{L^{2}}^{2}+\frac{1}{h^{2}}\left\|x^{1-\frac{\delta}{2}} w\left|d \varphi_{\epsilon}\right|\right\|_{L^{2}}^{2}+\left\|x^{1-\frac{\delta}{2}} d w\right\|_{L^{2}}^{2}\right) \\
\leq\left\|e^{\frac{\varphi_{\epsilon}}{h}}\left(\Delta-\lambda^{2}\right) e^{-\frac{\varphi_{\epsilon}}{h}} w\right\|_{L^{2}}^{2}
\end{gathered}
$$

which ends the proof.
3.2. The quadratic weight case for surfaces. In this section, $\varphi$ has quadratic growth at infinity, which corresponds to the case where $V \in e^{-\gamma / x^{2}} L^{\infty}$ for all $\gamma>0$. The proof when $\varphi$ has quadratic growth at infinity is even simpler than the linear growth case. We define $\varphi_{0} \in x^{-2} L^{\infty}$ to be a solution of $\Delta_{g_{0}} \varphi_{0}=1$, this is possible by Lemma 2.8 and one easily obtains from Proposition 5.61 of [13] that $\varphi_{0}=-x^{-2} / 4+O\left(x^{-1}\right)$ as $x \rightarrow 0$. We let $\varphi_{\epsilon}:=\varphi-\frac{h}{\epsilon} \varphi_{0}$ which satisfies $\Delta_{g_{0}} \varphi_{\epsilon} / h=-1 / \epsilon$.

If $K \subset M_{0}$ is a compact set, the Carleman estimate (4) in $K$ is satisfied by Proposition 3.1 of [7], it then remains to get the estimate in the ends $E_{1}, \ldots, E_{N}$. But the exact same proof as in Lemma 3.1 and Lemma 3.2 of [7] gives directly that for any $w \in C_{0}^{\infty}\left(E_{i}\right)$

$$
\begin{equation*}
\frac{C}{\epsilon}\left(\frac{1}{h}\|w\|_{L^{2}}^{2}+\frac{1}{h^{2}}\|w|d \varphi|\|_{L^{2}}^{2}+\frac{1}{h^{2}}\left\|w\left|d \varphi_{\epsilon}\right|\right\|_{L^{2}}^{2}+\|d w\|_{L^{2}}^{2}\right) \leq\left\|e^{\varphi_{\epsilon} / h} \Delta_{g_{0}} e^{-\varphi_{\epsilon} / h} w\right\|_{L^{2}}^{2} \tag{12}
\end{equation*}
$$

for some $C>0$ independent of $\epsilon, h$ and it suffices to glue the estimates in $K$ and in the ends $E_{i}$ as in Proposition 3.1 of [7], to obtain (12) for any $w \in C_{0}^{\infty}\left(M_{0}\right)$. Then by using triangle inequality

$$
\left\|e^{\varphi_{\epsilon} / h}\left(\Delta_{g_{0}}+V-\lambda^{2}\right) e^{-\varphi_{\epsilon} / h} u\right\|_{L^{2}} \leq\left\|e^{\varphi_{\epsilon} / h} \Delta_{g_{0}} e^{-\varphi_{\epsilon} / h} u\right\|_{L^{2}}+C\|u\|_{L^{2}}
$$

for some $C$ depending on $\lambda,\|V\|_{L^{\infty}}$, we see that the $V-\lambda^{2}$ term can be absorbed by the left hand side of (12) and we finally deduce

Proposition 3.5. Let $\left(M_{0}, g_{0}\right)$ be a Riemann surface with Euclidean ends and let $\varphi_{\epsilon}=$ $\varphi-\frac{h}{\epsilon} \varphi_{0}$ where $\varphi$ is a harmonic function with non-degenerate critical points and quadratic growth on $M_{0}$ and $\varphi_{0}$ satisfies $\Delta_{g_{0}} \varphi_{0}=1$ with $\varphi_{0} \in x^{-2} L^{\infty}\left(M_{0}\right)$. Then for all $V \in L^{\infty}$ there exists an $h_{0}>0, \epsilon_{0}$ and $C>0$ such that for all $0<h<h_{0}, h \ll \epsilon<\epsilon_{0}$ and $u \in C_{0}^{\infty}\left(M_{0}\right)$

$$
\frac{C}{\epsilon}\left(\frac{1}{h}\|u\|_{L^{2}}^{2}+\frac{1}{h^{2}}\|u \mid d \varphi\|_{L^{2}}^{2}+\|d u\|_{L^{2}}^{2}\right) \leq\left\|e^{\varphi_{\epsilon} / h}\left(\Delta_{g_{0}}+V-\lambda^{2}\right) e^{-\varphi_{\epsilon} / h} u\right\|_{L^{2}}^{2}
$$

The main difference with the linear weight case is that one can use a convexification which has quadratic growth at infinity which allows to absorb the $\lambda^{2}$ term, while it was not the case for the linearly growing weights.

## 4. Complex Geometric Optics on a Riemann Surface with Euclidean ends

As in $[1,9,7]$, the method for identifying the potential at a point $p$ is to construct complex geometric optic solutions depending on a small parameter $h>0$, with phase a Morse holomorphic function with a non-degenerate critical point at $p$, and then to apply the stationary phase method. Here, in addition, we need the phase to be of linear growth at infinity if $V \in e^{-\gamma / x} L^{\infty}$ for all $\gamma>0$ while the phase has to be of quadratic growth at infinity if $V \in e^{-\gamma / x^{2}} L^{\infty}$ for all $\gamma>0$.

We shall now assume that $M_{0}$ is a non-compact surface with genus $g$ with $N$ ends equipped with a metric $g_{0}$ which is Euclidean in the ends, and $V$ is a $C^{1, \alpha}$ function in $M_{0}$. Moreover, if $V \in e^{-\gamma / x} L^{\infty}$ for all $\gamma>0$, we ask that $N \geq \max (2 g+1,2)$ while if $V \in e^{-\gamma / x^{2}} L^{\infty}$ for all $\gamma>0$, we assume that $N \geq g+1$. As above, let us use a smooth positive function $x$ which is equal to 1 in a large compact set of $M_{0}$ and is equal to $x=|z|^{-1}$ in the regions $|z|>r_{0}$ of the ends $E_{i} \simeq\{z \in \mathbb{C} ;|z|>1\}$, where $r_{0}$ is a fixed large number. This function is a boundary defining function of the radial compactification of $M_{0}$ in the sense of Melrose [13]. To construct the complex geometric optics solutions, we will need to work with the weighted spaces $x^{-\alpha} L^{2}\left(M_{0}\right)$ where $\alpha \in \mathbb{R}_{+}$.

Let $\mathcal{H}$ be the finite dimensional complex vector space defined in the beginning of Section 2.3. Choose $p \in M_{0}$ such that there exists a Morse holomorphic function $\Phi=\varphi+i \psi \in \mathcal{H}$ on $M_{0}$, with a critical point at $p$; there is a dense set of such points by Proposition 2.1. The purpose of this section is to construct solutions $u$ on $M_{0}$ of $\left(\Delta-\lambda^{2}+V\right) u=0$ of the form

$$
\begin{equation*}
u=e^{\Phi / h}\left(a+r_{1}+r_{2}\right) \tag{13}
\end{equation*}
$$

for $h>0$ small, where $a \in x^{-J+1} L^{2}$ with $J \in \mathbb{R}_{+} \backslash \mathbb{N}$ is a holomorphic function on $M_{0}$, obtained by Lemma 2.6, such that $a(p) \neq 0$ and $a$ vanishing to order $L$ (for some fixed large $L)$ at all other critical points of $\Phi$, and finally $r_{1}, r_{2}$ will be remainder terms which are small as $h \rightarrow 0$ and have particular properties near the critical points of $\Phi$. More precisely, $e^{\varphi_{0} / \epsilon} r_{2}$ will be a $o_{L^{2}}(h)$ and $r_{1}$ will be a $O_{x^{-J} L^{2}}(h)$ but with an explicit expression, which can be used to obtain sufficient information in order to apply the stationary phase method.
4.0.1. Construction of $r_{1}$. We want to construct $r_{1}=O_{x^{-J} L^{2}}(h)$ which satisfies

$$
e^{-\Phi / h}\left(\Delta_{g_{0}}-\lambda^{2}+V\right) e^{\Phi / h}\left(a+r_{1}\right)=O_{x^{-J} L^{2}}(h)
$$

for some large $J \in \mathbb{R}_{+} \backslash \mathbb{N}$ so that $a \in x^{-J+1} L^{2}$.
Let $G$ be the operator of Lemma 2.8, mapping continuously $x^{-J+1} L^{2}\left(M_{0}\right)$ to $x^{-J-1} L^{2}\left(M_{0}\right)$. Then clearly $\bar{\partial} \partial G=\frac{i}{2} \star^{-1}$ when acting on $x^{-J+1} L^{2}$, here $\star^{-1}$ is the inverse of $\star$ mapping functions to 2 -forms. First, we will search for $r_{1}$ satisfying

$$
\begin{equation*}
e^{-2 i \psi / h} \partial e^{2 i \psi / h} r_{1}=-\partial G\left(a\left(V-\lambda^{2}\right)\right)+\omega+O_{x^{-J} H^{1}}(h) \tag{14}
\end{equation*}
$$

with $\omega \in x^{-J} L^{2}\left(M_{0}\right)$ a holomorphic 1-form on $M_{0}$ and $\left\|r_{1}\right\|_{x^{-J} L^{2}}=O(h)$. Indeed, using the fact that $\Phi$ is holomorphic we have

$$
e^{-\Phi / h} \Delta_{g_{0}} e^{\Phi / h}=-2 i \star \bar{\partial} e^{-\Phi / h} \partial e^{\Phi / h}=-2 i \star \bar{\partial} e^{-\frac{1}{h}(\Phi-\bar{\Phi})} \partial e^{\frac{1}{h}(\Phi-\bar{\Phi})}=-2 i \star \bar{\partial} e^{-2 i \psi / h} \partial e^{2 i \psi / h}
$$

and applying $-2 i \star \bar{\partial}$ to (14), this gives

$$
e^{-\Phi / h}\left(\Delta_{g_{0}}+V\right) e^{\Phi / h} r_{1}=-a\left(V-\lambda^{2}\right)+O_{x^{-J} L^{2}}(h)
$$

Writing $-\partial G\left(a\left(V-\lambda^{2}\right)\right)=: c(z) d z$ in local complex coordinates, $c(z)$ is $C^{2, \alpha}$ by elliptic regularity and we have $2 i \partial_{\bar{z}} c(z)=a\left(V-\lambda^{2}\right)$, therefore $\partial_{z} \partial_{\bar{z}} c\left(p^{\prime}\right)=\partial_{\bar{z}}^{2} c\left(p^{\prime}\right)=0$ at each critical point $p^{\prime} \neq p$ by construction of the function $a$. Therefore, we deduce that at each critical point $p^{\prime} \neq p, c(z)$ has Taylor series expansion $\sum_{j=0}^{2} c_{j} z^{j}+O\left(|z|^{2+\alpha}\right)$. That is, all the lower order terms of the Taylor expansion of $c(z)$ around $p^{\prime}$ are polynomials of $z$ only. By Lemma 2.7, and possibly by taking $J$ larger, there exists a holomorphic function $f \in x^{-J} L^{2}$ such that $\omega:=\partial f$ has Taylor expansion equal to that of $\partial G\left(a\left(V-\lambda^{2}\right)\right)$ at all critical points $p^{\prime} \neq p$ of $\Phi$. We deduce that, if $b:=-\partial G\left(a\left(V-\lambda^{2}\right)\right)+\omega=b(z) d z$, we have

$$
\begin{array}{rc}
\left|\partial_{\bar{z}}^{m} \partial_{z}^{\ell} b(z)\right|= & O\left(|z|^{2+\alpha-\ell-m}\right), \\
|b(z)|=O(|z|), & \text { for } \ell+m \leq 2, \text { at critical points } p^{\prime} \neq p  \tag{15}\\
\text { if } p^{\prime}=p
\end{array}
$$

Now, we let $\chi_{1} \in C_{0}^{\infty}\left(M_{0}\right)$ be a cutoff function supported in a small neighbourhood $U_{p}$ of the critical point $p$ and identically 1 near $p$, and $\chi \in C_{0}^{\infty}\left(M_{0}\right)$ is defined similarly with $\chi=1$ on the support of $\chi_{1}$. We will construct $r_{1}$ to be a sum $r_{1}=r_{11}+h r_{12}$ where $r_{11}$ is a compactly supported approximate solution of (14) near the critical point $p$ of $\Phi$ and $r_{12}$ is correction term supported away from $p$. We define locally in complex coordinates centered at $p$ and containing the support of $\chi$

$$
\begin{equation*}
r_{11}:=\chi e^{-2 i \psi / h} R\left(e^{2 i \psi / h} \chi_{1} b\right) \tag{16}
\end{equation*}
$$

where $R f(z):=-(2 \pi i)^{-1} \int_{\mathbb{R}^{2}} \frac{1}{\bar{z}-\bar{\xi}} f d \bar{\xi} \wedge d \xi$ for $f \in L^{\infty}$ compactly supported is the classical Cauchy operator inverting locally $\partial_{z}$ ( $r_{11}$ is extended by 0 outside the neighbourhood of $p$ ). The function $r_{11}$ is in $C^{3, \alpha}\left(M_{0}\right)$ and we have

$$
\begin{gather*}
e^{-2 i \psi / h} \partial\left(e^{2 i \psi / h} r_{11}\right)=\chi_{1}\left(-\partial G\left(a\left(V-\lambda^{2}\right)\right)+\omega\right)+\eta \\
\text { with } \eta:=e^{-2 i \psi / h} R\left(e^{2 i \psi / h} \chi_{1} b\right) \partial \chi \tag{17}
\end{gather*}
$$

We then construct $r_{12}$ by observing that $b$ vanishes to order $2+\alpha$ at critical points of $\Phi$ other than $p$ (from (15)), and $\partial \chi=0$ in a neighbourhood of any critical point of $\psi$, so we can find $r_{12}$ satisfying

$$
\begin{equation*}
2 i r_{12} \partial \psi=\left(1-\chi_{1}\right) b \tag{18}
\end{equation*}
$$

This is possible since both $\partial \psi$ and the right hand side are valued in $T_{1,0}^{*} M_{0}$ and $\partial \psi$ has finitely many isolated zeroes on $M_{0}: r_{12}$ is then a function which is in $C^{2, \alpha}\left(M_{0} \backslash P\right)$ where $P:=\left\{p_{1}, \ldots, p_{n}\right\}$ is the set of critical points other than $p$, it extends to a function in $C^{1, \alpha}\left(M_{0}\right)$ and it satisfies in local complex coordinates $z$ at each $p_{j}$

$$
\left|\partial_{\bar{z}}^{\beta} \partial_{z}^{\gamma} r_{12}(z)\right| \leq C|z|^{1+\alpha-\beta-\gamma}, \quad \beta+\gamma \leq 2
$$

by using also the fact that $\partial \psi$ can be locally be considered as a smooth function with a zero of order 1 at each $p_{j}$. Moreover $b \in x^{-J} H^{2}\left(M_{0}\right)$ thus $r_{1} \in x^{-J} H^{2}\left(M_{0}\right)$ and we have

$$
e^{-2 i \psi / h} \partial\left(e^{2 i \psi / h} r_{1}\right)=b+h \partial r_{12}+\eta=-\partial G\left(a\left(V-\lambda^{2}\right)\right)+\omega+h \partial r_{12}+\eta
$$

Lemma 4.1. The following estimates hold true

$$
\begin{gathered}
\|\eta\|_{H^{2}\left(M_{0}\right)}=O(|\log h|), \quad\|\eta\|_{H^{1}\left(M_{0}\right)} \leq O(h|\log h|), \quad\left\|x^{J} \partial r_{12}\right\|_{H^{1}\left(M_{0}\right)}=O(1) \\
\left\|x^{J} r_{1}\right\|_{L^{2}}=O(h), \quad\left\|x^{J}\left(r_{1}-h \widetilde{r}_{12}\right)\right\|_{L^{2}}=o(h)
\end{gathered}
$$

where $\widetilde{r}_{12}$ solves $2 i \widetilde{r}_{12} \partial \psi=b$.

Proof. The proof is exactly the same as the proof of Lemma 4.2 in [8], except that one needs to add the weight $x^{J}$ to have bounded integrals.

As a direct consequence, we have
Corollary 4.2. With $r_{1}=r_{11}+h r_{12}$, there exists $J>0$ such that

$$
\left\|e^{-\Phi / h}\left(\Delta_{g_{0}}-\lambda^{2}+V\right) e^{\Phi / h}\left(a+r_{1}\right)\right\|_{x^{-J} L^{2}\left(M_{0}\right)}=O(h|\log h|)
$$

4.0.2. Construction of $r_{2}$. In this section, we complete the construction of the complex geometric optic solutions. We deal with the general case of surfaces and we shall show the following

Proposition 4.1. If $\varphi_{0}$ is the subharmonic function constructed in Section 3, then for $\epsilon$ small enough there exist solutions to $\left(\Delta_{g_{0}}-\lambda^{2}+V\right) u=0$ of the form $u=e^{\Phi / h}\left(a+r_{1}+r_{2}\right)$ with $r_{1}=r_{11}+h r_{12}$ constructed in the previous section and $r_{2} \in e^{-\varphi_{0} / \epsilon} L^{2}$ satisfying $\left\|e^{\varphi_{0} / \epsilon} r_{2}\right\|_{L^{2}} \leq$ $C h^{3 / 2}|\log h|$.

This is a consequence of the following Lemma (which follows from the Carleman estimate obtained in Section 3 above)

Lemma 4.3. Let $\delta \in(0,1), V \in x^{1-\frac{\delta}{2}} L^{\infty}\left(M_{0}\right)$, and $\varphi_{\epsilon}=\varphi-\frac{h}{\epsilon} \varphi_{0}$ a weight with linear growth at infinity as in Proposition 3.4. For all $f \in L^{2}\left(M_{0}\right)$ and all $h>0$ small enough, there exists a solution $v \in L^{2}\left(M_{0}\right)$ to the equation

$$
\begin{equation*}
e^{-\varphi_{\epsilon} / h}\left(\Delta_{g}-\lambda^{2}+V\right) e^{\varphi_{\epsilon} / h} v=x^{1-\frac{\delta}{2}} f \tag{19}
\end{equation*}
$$

satisfying

$$
\|v\|_{L^{2}\left(M_{0}\right)} \leq C h^{\frac{1}{2}}\|f\|_{L^{2}\left(M_{0}\right)}
$$

If $\varphi_{\epsilon}$ has quadratic growth at infinity, the same result is true when $V \in L^{\infty}\left(M_{0}\right)$ but $x^{1-\frac{\delta}{2}} f$ can be replaced by $f \in L^{2}$ in (19).

Proof. The proof is based on a duality argument. Let $P_{h}:=e^{\varphi_{\epsilon} / h}\left(\Delta_{g}-\lambda^{2}+V\right) e^{-\varphi_{\epsilon} / h}$ and for all $h>0$ the real vector space $\mathcal{A}:=\left\{u \in x^{-1+\frac{\delta}{2}} H^{1}\left(M_{0}\right) ; P_{h} u \in L^{2}\left(M_{0}\right)\right\}$ equipped with the real scalar product

$$
(u, w)_{\mathcal{A}}:=\left\langle P_{h} u, P_{h} w\right\rangle_{L^{2}} .
$$

By the Carleman estimate of Proposition 3.4, the space $\mathcal{A}$ is a Hilbert space equipped with the scalar product above if $h<h_{0}$, and thus the linear functional $L: w \rightarrow \int_{M_{0}} x^{1-\frac{\delta}{2}} f w \mathrm{dvol}_{g_{0}}$ on $\mathcal{A}$ is continuous with norm bounded by $C h^{\frac{1}{2}}\|f\|_{L^{2}}$ by Proposition 3.4, and by Riesz theorem there is an element $u \in \mathcal{A}$ such that $(., u)_{\mathcal{A}}=L$ and with norm bounded by the norm of $L$. It remains to take $v:=P_{h} u$ which solves $P_{h}^{*} v=x^{1-\frac{\delta}{2}} f$ where $P_{h}^{*}=e^{-\varphi_{\epsilon} / h}\left(\Delta_{g}-\lambda^{2}+V\right) e^{\varphi_{\epsilon} / h}$ is the adjoint of $P_{h}$ and $v$ satisfies the desired norm estimate. The proof when the weight $\varphi_{\epsilon}$ has quadratic growth at infinity is the same, but improves slightly due to the Carleman estimate of Proposition 3.5.

Proof of Proposition 4.1. We first solve the equation

$$
\left(\Delta+V-\lambda^{2}\right) e^{\varphi_{\epsilon} / h} \widetilde{r}_{2}=x^{1-\frac{\delta}{2}} e^{\varphi_{\epsilon} / h}\left(x^{-1+\frac{\delta}{2}} e^{-\varphi_{\epsilon} / h}\left(\Delta+V-\lambda^{2}\right) e^{\Phi / h}\left(a+r_{1}\right)\right)
$$

by using Lemma 4.3 and the fact that for $J$ large, there is $C>0$ such that for all $h<h_{0}$

$$
\left\|x^{-1+\frac{\delta}{2}} e^{-\varphi_{\epsilon} / h}\left(\Delta+V-\lambda^{2}\right) e^{\Phi / h}\left(a+r_{1}\right)\right\|_{L^{2}} \leq C\left\|x^{J} e^{-\Phi / h}\left(\Delta-\lambda^{2}+V\right) e^{\Phi / h}\left(a+r_{1}\right)\right\|_{L^{2}}
$$

since $x^{-J-1} e^{\varphi_{0} / \epsilon} \in L^{\infty}\left(M_{0}\right)$ for all $J$ (recall that $\varphi_{0} \sim-x^{-\delta} / \delta^{2}$ as $x \rightarrow 0$ ). But now the right hand side is bounded by $O(h|\log h|)$ according to Corollary 4.2 , therefore we set $r_{2}:=-e^{-i \psi / h-\varphi_{0} / \epsilon} \widetilde{r}_{2}$ which satisfies $\left(\Delta_{g_{0}}-\lambda^{2}+V\right) e^{\Phi / h}\left(a+r_{1}+r_{2}\right)=0$ and, by Lemma 4.3, the norm estimate $\left\|e^{\varphi_{0} / \epsilon} r_{2}\right\|_{L^{2}} \leq O\left(h^{3 / 2}|\log h|\right)$.

## 5. Scattering on surface with Euclidean ends

Let $\left(M_{0}, g_{0}\right)$ be a surface with Euclidean ends and $V \in e^{-\gamma / x} L^{\infty}\left(M_{0}\right)$ for some $\gamma$. The scattering theory in this setting is described for instance in Melrose [14], here we will follow this presentation (see also Section 3 in Uhlmann-Vasy [26] for the $\mathbb{R}^{n}$ case). First, using standard methods in scattering theory, we define the resolvent on the continuous spectrum as follows

Lemma 5.1. The resolvent $R_{V}(\lambda):=\left(\Delta_{g_{0}}+V-\lambda^{2}\right)^{-1}$ admits a meromorphic extension from $\{\operatorname{Im}(\lambda)<0\}$ to $\{\operatorname{Im}(\lambda) \leq A, \operatorname{Re}(\lambda) \neq 0\}$, as a family of operators mapping $e^{-\gamma / x} L^{2}\left(M_{0}\right)$ to $e^{\gamma / x} L^{2}\left(M_{0}\right)$ for any $\gamma>A$. Moreover, for $\lambda \in \mathbb{R} \backslash\{0\}$ not a pole, $R_{V}(\lambda)$ maps continuously $x^{\alpha} L^{2}$ to $x^{-\alpha} L^{2}$ for any $\alpha>1 / 2$.
Proof. The statement is known for $V=0$ and $M_{0}=\mathbb{R}^{2}$ by using the explicit formula of the resolvent convolution kernel on $\mathbb{R}^{2}$ in terms of Hankel functions (see for instance [14]), we shall denote $R_{0}(\lambda)$ this continued resolvent. More precisely, for all $A>0$, the operator $R_{0}(\lambda)$ continues analytically from $\{\operatorname{Im}(\lambda)<0\}$ to $\{\operatorname{Im}(\lambda) \leq A, \operatorname{Re}(\lambda) \neq 0\}$ as a family of bounded operators mapping $e^{-\gamma / x} L^{2}$ to $e^{\gamma / x} L^{2}$ for any $\gamma>A$. Now we can set $\chi \in C_{0}^{\infty}\left(M_{0}\right)$ such that $1-\chi$ is supported in the ends $E_{i}$, and let $\chi_{0}, \chi_{1} \in C_{0}^{\infty}\left(M_{0}\right)$ such that $\left(1-\chi_{0}\right)=1$ on the support of $(1-\chi)$ and $\chi_{1}=1$ on the support of $\chi$. Let $\lambda_{0} \in-i \mathbb{R}_{+}$with $i \lambda_{0} \gg 0$, then the resolvent $R_{0}\left(\lambda_{0}\right)$ is well defined from $L^{2}\left(M_{0}\right)$ to $H^{2}\left(M_{0}\right)$ since the Laplacian is essentially self-adjoint [23, Proposition 8.2.4], and we have a parametrix

$$
E(\lambda):=\left(1-\chi_{0}\right) R_{0}(\lambda)(1-\chi)+\chi_{1} R_{0}\left(\lambda_{0}\right) \chi
$$

which satisfies

$$
\begin{gathered}
\left(\Delta_{g_{0}}-\lambda^{2}+V\right) E(\lambda)=1+K(\lambda) \\
K(\lambda):=\left(\left[\Delta_{g_{0}}, \chi_{1}\right]-\left(\lambda^{2}-\lambda_{0}^{2}\right) \chi_{1}\right) R_{0}\left(\lambda_{0}\right) \chi-\left[\Delta_{g_{0}}, \chi_{0}\right] R_{0}(\lambda)(1-\chi)+V E(\lambda)
\end{gathered}
$$

where here we use the notation $R_{0}(\lambda)$ for an integral kernel on $M_{0}$, which in the charts $\left\{z \in \mathbb{R}^{2} ;|z|>1\right\}$ corresponding the ends $E_{1}, \ldots E_{N}$, is given by the integral kernel of $\left(\Delta_{\mathbb{R}^{2}}-\right.$ $\left.\lambda^{2}\right)^{-1}$. Using the explicit expression of the convolution kernel of $R_{0}(\lambda)$ in the ends (see for instance Section 1.5 of [14]) and the decay assumption on $V$, it is direct to see that for $\operatorname{Im}(\lambda)<A, \operatorname{Re}(\lambda) \neq 0$, the map $\lambda \mapsto K(\lambda)$ a is compact analytic family of bounded operators from $e^{-\gamma / x} L^{2}$ to $e^{-\gamma / x} L^{2}$ for any $\gamma>A$. Moreover $1+K\left(\lambda_{0}\right)$ is invertible since $\left\|K\left(\lambda_{0}\right)\right\|_{L^{2} \rightarrow L^{2}} \leq 1 / 2$ if $i \lambda_{0}$ is large enough. Then by analytic Fredholm theory, the resolvent $R_{V}(\lambda)$ has an meromorphic extension to $\operatorname{Im}(\lambda)<A, \operatorname{Re}(\lambda) \neq 0$ as a bounded operator from $e^{-\gamma / x} L^{2}$ to $e^{\gamma / x} L^{2}$ if $\gamma>A$, given by

$$
R_{V}(\lambda)=E(\lambda)(1+K(\lambda))^{-1}
$$

Now $(1+K(\lambda))^{-1}=1+Q(\lambda)$ for some $Q(\lambda)=-K(\lambda)(1+K(\lambda))^{-1}$ mapping $e^{-\gamma / x} L^{2}$ to itself for any $\gamma>A$, which proves the mapping properties of $R_{V}(\lambda)$ on exponential weighted spaces. For the mapping properties on $\{\operatorname{Re}(\lambda)=0\}$, a similar argument works.

A corollary of this Lemma is the mapping property
Corollary 5.2. For $\lambda \in \mathbb{R} \backslash\{0\}$ not a pole of $R_{V}(\lambda)$, and $f \in e^{-\gamma / x} L^{\infty}$ for some $\gamma>0$, then there exists $v \in C^{\infty}\left(\partial \bar{M}_{0}\right)$ such that

$$
R_{V}(\lambda) f-x^{\frac{1}{2}} e^{-i \lambda / x} v \in L^{2}
$$

Proof. Using the expression $R_{V}(\lambda)=E(\lambda)(1+Q(\lambda))$ of the proof of Lemma 5.1, it suffices to know the mapping property of $E(\lambda)$ on $e^{-\gamma / x} L^{2}$, but since outside a compact set (i.e. in the ends) $E(\lambda)$ is given by the free resolvent on $\mathbb{R}^{2}$, this amounts to proving the statement in $\mathbb{R}^{2}$, which is well-known: for instance, this is proved for $f \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$ in Section 1.7 [14] but the proof extends easily to $f \in e^{-\gamma / x} L^{\infty}\left(\mathbb{R}^{2}\right)$ since the only used assumption on $f$ for applying a stationary phase argument is actually that the Fourier transform $\hat{f}(z)$ has a holomorphic extension in a complex neighbourhood of $\mathbb{R}^{2}$.

We also have a boundary pairing, the proof of which is exactly the same as [14, Lemma 2.2] (see also Proposition 3.1 of [26]).

Lemma 5.3. For $\lambda>0$ and $V \in e^{-\gamma / x} L^{\infty}\left(M_{0}\right)$, if $u_{ \pm} \in x^{-\alpha} L^{2}\left(M_{0}\right)$ for some $\alpha>1 / 2$ and $\left(\Delta_{g_{0}}-\lambda^{2}+V\right) u_{ \pm} \in x^{\alpha} L^{2}\left(M_{0}\right)$ with

$$
u_{+}-x^{\frac{1}{2}} e^{i \lambda / x} f_{++}-x^{\frac{1}{2}} e^{-i \lambda / x} f_{+-} \in L^{2}, \quad u_{-}-x^{\frac{1}{2}} e^{i \lambda / x} f_{-+}-x^{\frac{1}{2}} e^{-i \lambda / x} f_{--} \in L^{2}
$$

for some $f_{ \pm \pm} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$, then

$$
\left\langle u_{+},\left(\Delta_{g_{0}}+V-\lambda^{2}\right) u_{-}\right\rangle-\left\langle\left(\Delta_{g_{0}}+V-\lambda^{2}\right) u_{+}, u_{-}\right\rangle=2 i \lambda \int_{\partial \bar{M}_{0}}\left(f_{++} \overline{f_{-+}}-f_{+-} \overline{f_{--}}\right)
$$

where the volume form on $\partial \bar{M}_{0} \simeq \sqcup_{i=1}^{N} S^{1}$ is induced by the metric $\left.x^{2} g\right|_{T \partial \bar{M}_{0}}$.
As a corollary, the same exact arguments as in Sections 2.2 to 2.5 in [14] show ${ }^{1}$
Corollary 5.4. The operator $R_{V}(\lambda)$ is analytic on $\lambda \in \mathbb{R} \backslash\{0\}$ as a bounded operator from $x^{\alpha} L^{2}$ to $x^{-\alpha} L^{2}$ if $\alpha>1 / 2$.

In $\mathbb{R}^{2}$ there is a Poisson operator $P_{0}(\lambda)$ mapping $C^{\infty}\left(S^{1}\right)$ to $x^{-\alpha} L^{2}\left(\mathbb{R}^{2}\right)$ for $\alpha>1 / 2$, which satisfies that for any $f_{+} \in C^{\infty}\left(S^{1}\right)$ there exists $f_{-} \in C^{\infty}\left(S^{1}\right)$ such that

$$
P_{0}(\lambda) f_{+}-x^{\frac{1}{2}} e^{i \lambda / x} f_{+}-x^{\frac{1}{2}} e^{-i \lambda / x} f_{-} \in L^{2}, \quad\left(\Delta-\lambda^{2}\right) P_{0}(\lambda) f_{+}=0
$$

We can therefore define in our case a similar Poisson operator $P_{V}(\lambda)$ mapping $C^{\infty}\left(\partial \bar{M}_{0}\right)$ to $x^{-\alpha} L^{2}$ for $\alpha>1 / 2$, by

$$
\begin{equation*}
P_{V}(\lambda) f_{+}:=(1-\chi) P_{0}(\lambda) f_{+}-R_{V}(\lambda)\left(\Delta_{g_{0}}+V-\lambda^{2}\right)(1-\chi) P_{0}(\lambda) f_{+} \tag{20}
\end{equation*}
$$

where $1-\chi \in C^{\infty}\left(M_{0}\right)$ equals 1 in the ends $E_{i}$ and $P_{0}(\lambda)$ denotes here the Schwartz kernel of the Poisson operator on $\mathbb{R}^{2}$ pulled back to each of the Euclidean ends $E_{i}$ of $M_{0}$ in the

[^0]obvious way. Then, since $\left(\Delta_{g_{0}}+V-\lambda^{2}\right)(1-\chi) P_{0}(\lambda) f_{+} \in e^{-\gamma / x} L^{2}$ for all $\gamma>0$, it suffices to use Corollaries 5.2 and 5.4 to see that it defines an analytic Poisson operator $P_{V}(\lambda)$ on $\lambda \in \mathbb{R} \backslash\{0\}$ satisfying that for all $f_{+} \in C^{\infty}\left(\partial \overline{M_{0}}\right)$, there exists $f_{-} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$ such that
\[

$$
\begin{equation*}
P_{V}(\lambda) f_{+}-x^{\frac{1}{2}} e^{i \lambda / x} f_{+}-x^{\frac{1}{2}} e^{-i \lambda / x} f_{-} \in L^{2}, \quad\left(\Delta+V-\lambda^{2}\right) P_{V}(\lambda) f_{+}=0 \tag{21}
\end{equation*}
$$

\]

Moreover, it is easily seen to be the unique solution of (21): indeed, if two such solutions exist then the difference is a solution $u$ with asymptotic $x^{\frac{1}{2}} e^{-i \lambda / x} f_{-}+L^{2}$ for some $f_{-} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$, but applying Lemma 5.3 with $u_{-}=u_{+}=u$ shows that $f_{-}=0$, thus $u \in L^{2}$, which implies $u=0$ by Corollary 5.4.
Definition 5.5. The scattering matrix $S_{V}(\lambda): C^{\infty}\left(\partial \bar{M}_{0}\right) \rightarrow C^{\infty}\left(\partial \bar{M}_{0}\right)$ for $\lambda \in \mathbb{R} \backslash\{0\}$ is defined to be the map $S_{V}(\lambda) f_{+}:=f_{-}$where $f_{-}$is given by the asymptotic

$$
P_{V}(\lambda) f_{+}=x^{\frac{1}{2}} e^{i \lambda / x} f_{+}+x^{\frac{1}{2}} e^{-i \lambda / x} f_{-}+g, \quad \text { with } g \in L^{2}
$$

We remark that, using Lemma 5.3 and the uniqueness of the Poisson operator, one easily deduces for $\lambda \in \mathbb{R} \backslash\{0\}$

$$
\begin{equation*}
S_{V}(\lambda)^{*}=S_{V}(-\lambda)=S_{V}(\lambda)^{-1} \tag{22}
\end{equation*}
$$

where the scalar product on $L^{2}\left(\partial \bar{M}_{0}\right)$ is induced by the metric $\left.x^{2} g_{0}\right|_{T \partial \bar{M}_{0}}$.
We can now state a density result similar to Proposition 3.3 of [26]:
Proposition 5.6. If $V \in e^{-\gamma_{0} / x} L^{\infty}\left(M_{0}\right)$ (resp. $V \in e^{-\gamma_{0} / x^{2}} L^{\infty}\left(M_{0}\right)$ ) for some $\gamma_{0}>0$, and $\lambda \in \mathbb{R} \backslash\{0\}$, then for any $0<\gamma<\gamma^{\prime}<\gamma_{0}$ the set

$$
\mathcal{F}:=\left\{P_{V}(\lambda) f_{+} ; f_{+} \in C^{\infty}\left(\partial \bar{M}_{0}\right)\right\}
$$

is dense in the null space of $\Delta_{g_{0}}+V-\lambda^{2}$ in $e^{\gamma / x} L^{2}\left(M_{0}\right)$ for the topology of $e^{\gamma^{\prime} / x} L^{2}\left(M_{0}\right)$ (resp. in $e^{\gamma / x^{2}} L^{2}\left(M_{0}\right)$ for the topology of $e^{\gamma^{\prime} / x^{2}} L^{2}\left(M_{0}\right)$ ).
Proof. First assume $V \in e^{-\gamma_{0} / x} L^{\infty}\left(M_{0}\right)$. Let $w \in e^{-\gamma^{\prime} / x} L^{2}$ be orthogonal to $\mathcal{F}$, and set $u_{-}:=$ $R_{V}(\lambda) w$ and $u_{+}=P_{V}(\lambda) f_{++}$for some $f_{++} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$. Then, define $f_{--} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$ by $R_{V}(\lambda) w-x^{\frac{1}{2}} e^{-i \lambda / x} f_{--} \in L^{2}$, and from Lemma 5.3 we obtain $\left\langle f_{+-}, f_{--}\right\rangle=0$ since $\left\langle w, P_{V}(\lambda) f_{++}\right\rangle=0$ by assumption. Since $f_{+-}=S_{V}(\lambda) f_{++}$is arbitrary, then $f_{--}=0$ and $u_{-} \in L^{2}$. In particular, from the parametrix constructed in the proof of Lemma 5.1

$$
R_{V}(\lambda) w-\left(1-\chi_{0}\right) R_{0}(\lambda)(1-\chi)(1+Q(\lambda)) w \in L^{2}
$$

with $(1+Q(\lambda)) w \in e^{-\gamma^{\prime} / x} L^{2}$. Since in each end, $R_{0}(\lambda)$ is the integral kernel of the free resolvent of the Euclidean Laplacian on $\mathbb{R}^{2}$ and $\left(1-\chi_{0}\right)$ and $(1-\chi)$ are supported in the ends, we can view the term $\left(1-\chi_{0}\right) R_{0}(\lambda)(1-\chi)(1+Q(\lambda)) w$ as a disjoint sum (over the ends) of functions on $\mathbb{R}^{2}$ of the form

$$
\begin{equation*}
\left(1-\chi_{0}(z)\right) \frac{1}{(2 \pi)^{2}} \int_{\mathbb{R}^{2}} e^{i z \xi}\left(\xi^{2}-\lambda^{2}-i 0\right)^{-1} \hat{f}(\xi) d \xi \tag{23}
\end{equation*}
$$

where in each end $E_{i}, f=(1-\chi)(1+Q(\lambda)) w \in e^{-\gamma^{\prime} / x} L^{2}\left(E_{i}\right)$ can be considered as a function in $e^{-\gamma^{\prime}|z|} L^{2}\left(\mathbb{R}^{2}\right)$. By the Paley-Wiener theorem, $\hat{f}$ is holomorphic in a strip $U=\left\{|\operatorname{Im}(\xi)|<\gamma^{\prime}\right\}$ with bound $\sup _{\eta \leq \gamma}\|\hat{f}(\cdot+i \eta)\|_{L^{2}\left(\mathbb{R}^{2}\right)}<\infty$ for all $\gamma<\gamma^{\prime}$, so the fact that (23) is in $L^{2}$ implies that $\hat{f}$ vanishes at the real sphere $\left\{\xi \in \mathbb{R}^{2} ; \xi^{2}=\lambda^{2}\right\}$, and thus there exists $h$ holomorphic in $U$ such that $\hat{f}(\xi)=\left(\xi^{2}-\lambda^{2}\right) h(\xi)$ (see e.g. the proof of Lemma 2.5 in [17]), and satisfying the same types of $L^{2}$ estimates as $\hat{f}$ in $U$ on lines $\operatorname{Im}(\xi)=$ cst. By the Paley-Wiener theorem
again, we deduce that (23) is in $e^{-\gamma|z|} L^{2}$ and thus $R_{V}(\lambda) w \in e^{-\gamma / x} L^{2}\left(M_{0}\right)$ for any $\gamma<\gamma^{\prime}$. Then if $v \in e^{\gamma / x} L^{2}\left(M_{0}\right)$ and $\left(\Delta_{g_{0}}+V-\lambda^{2}\right) v=0$, one has by integration by parts

$$
0=\left\langle R_{V}(\lambda) w,\left(\Delta_{g_{0}}+V-\lambda^{2}\right) v\right\rangle=\langle w, v\rangle
$$

which ends the proof in the case $V \in e^{-\gamma_{0} / x} L^{\infty}\left(M_{0}\right)$. The quadratic decay case $V \in$ $e^{-\gamma_{0} / x^{2}} L^{\infty}\left(M_{0}\right)$ is exactly similar but instead of Paley-Wiener theorem, we use Corollary 7.3 and the inclusions $e^{-\gamma^{\prime} / x^{2}} L^{2} \subset e^{-\gamma^{\prime \prime} / x^{2}} L^{1} \cap e^{-\gamma^{\prime \prime} / x^{2}} L^{2}$ and $e^{-\gamma^{\prime} / x^{2}} L^{\infty} \subset e^{-\gamma / x^{2}} L^{2}$ for all $\gamma<\gamma^{\prime \prime}<\gamma^{\prime}$.

## 6. Identifying the potential

6.1. The case of a surface. On a Riemann surface $\left(M_{0}, g_{0}\right)$ with $N$ Euclidean ends and genus $g$, we assume that $V_{1}, V_{2} \in C^{1, \alpha}\left(M_{0}\right)$ are two real valued potentials such that the respective scattering operators $S_{V_{1}}(\lambda)$ and $S_{V_{2}}(\lambda)$ agree for a fixed $\lambda>0$. We also assume that for all $\gamma>0$

$$
V_{1}, V_{2} \in \begin{cases}e^{-\gamma / x} L^{\infty}\left(M_{0}\right) & \text { if } N \geq \max (2 g+1,2) \\ e^{-\gamma / x^{2}} L^{\infty}\left(M_{0}\right) & \text { if } N \geq g+1\end{cases}
$$

By considering the asymptotics of $u_{1}:=P_{V_{1}}(\lambda) f_{1}$ and $P_{V_{2}}(-\lambda) f_{2}$ for $f_{i} \in C^{\infty}\left(\partial \bar{M}_{0}\right)$ we easily have by integration by parts that

$$
\begin{align*}
\int_{M_{0}}\left(V_{1}-V_{2}\right) u_{1} \overline{u_{2}} \operatorname{dvol}_{g_{0}} & =-2 i \lambda \int_{\partial \bar{M}_{0}} S_{V_{1}}(\lambda) f_{1} \cdot \overline{f_{2}}-f_{1} \cdot \overline{S_{V_{2}}(-\lambda) f_{2}} \\
& =-2 i \lambda \int_{\partial \bar{M}_{0}}\left(S_{V_{1}}(\lambda)-S_{V_{2}}(\lambda)\right) f_{1} \cdot \overline{f_{2}}=0 \tag{24}
\end{align*}
$$

by using (22). From Proposition 5.6, this implies by density that, if $V \in e^{-\gamma / x} L^{\infty}$ (resp. $V \in e^{-\gamma / x^{2}} L^{\infty}$ for all $\gamma>0$ ), then for all solutions $u_{i}$ of ( $\left.\Delta_{g_{0}}+V_{i}-\lambda^{2}\right) u_{i}=0$ in $e^{\gamma^{\prime} / x} L^{2}\left(M_{0}\right)$ (resp. $u_{i} \in e^{\gamma^{\prime} / x^{2}} L^{2}\left(M_{0}\right)$ ) for some $\gamma^{\prime}>0$, we have

$$
\begin{equation*}
\int_{M_{0}}\left(V_{1}-V_{2}\right) u_{1} \overline{u_{2}} \operatorname{dvol}_{g_{0}}=0 \tag{25}
\end{equation*}
$$

We shall now use our complex geometric optics solutions as special solutions in the weighted space $e^{-\gamma^{\prime} / h x} L^{2}\left(M_{0}\right)$ (resp. $\left.e^{-\gamma^{\prime} / h x^{2}} L^{2}\left(M_{0}\right)\right)$ for some $\gamma^{\prime}>0$ if $V \in e^{-\gamma / x} L^{\infty}$ (resp. $V \in$ $e^{-\gamma / x^{2}} L^{\infty}$ ) for all $\gamma>0$.

Let $p \in M_{0}$ be such that, using Proposition 2.1, we can choose a holomorphic Morse function $\Phi=\varphi+i \psi$ with linear or quadratic growth on $M_{0}$ (depending on the topological assumption), with a critical point at $p$. Then for the complex geometric optics solutions $u_{1}, u_{2}$ with phase $\Phi$ constructed in Section 4, the identity (25) holds true. We will then deduce the
Proposition 6.1. Let $\lambda \in(0, \infty)$ and assume that $S_{V_{1}}(\lambda)=S_{V_{2}}(\lambda)$, then $V_{1}(p)=V_{2}(p)$.
Proof. Let $u_{1}$ and $u_{2}$ be solutions on $M_{0}$ to

$$
\left(\Delta_{g}+V_{j}-\lambda^{2}\right) u_{j}=0
$$

constructed in Section 4 with phase $\Phi$ for $u_{1}$ and $-\Phi$ for $u_{2}$, thus of the form

$$
u_{1}=e^{\Phi / h}\left(a+r_{1}^{1}+r_{2}^{1}\right), \quad u_{2}=e^{-\Phi / h}\left(a+r_{1}^{2}+r_{2}^{2}\right) .
$$

We have the identity

$$
\int_{M_{0}} u_{1}\left(V_{1}-V_{2}\right) \overline{u_{2}} \operatorname{dvol}_{g_{0}}=0
$$

Then by using the estimates in Lemma 4.1 and Proposition 4.1 we have, as $h \rightarrow 0$,

$$
\int_{M_{0}} e^{2 i \psi / h}|a|^{2}\left(V_{1}-V_{2}\right) \operatorname{dvol}_{g_{0}}+h \int_{M_{0}} e^{2 i \psi / h}\left(\bar{a} \widetilde{r}_{12}^{1}+a \overline{\widetilde{r}_{12}^{2}}\right)\left(V_{1}-V_{2}\right) \operatorname{dvol}_{g_{0}}+o(h)=0
$$

where $\widetilde{r}_{12}^{j} \in L^{\infty}\left(M_{0}\right)$ are defined in Lemma 4.1, with the superscript $j$ refering to the solution for the potential $V_{j}$; in particular these functions $\widetilde{r}_{12}^{j}$ are independent of $h$.

By splitting $V_{i}(\cdot)=\left(V_{i}(\cdot)-V_{i}(p)\right)+V_{i}(p)$ and using the $C^{1, \alpha}$ regularity assumption on $V_{i}$, one can use stationary phase for the $V_{i}(p)$ term and integration by parts to gain a power of $h$ for the $V_{i}(\cdot)-V_{i}(p)$ term (see the proof of Lemma 5.4 in [8] for details) to deduce

$$
\int_{M_{0}} e^{2 i \psi / h}|a|^{2}\left(V_{1}-V_{2}\right) \operatorname{dvol}_{g_{0}}=C h\left(V_{1}(p)-V_{2}(p)\right)+o(h)
$$

for some $C \neq 0$. Therefore,

$$
C h\left(V_{1}(p)-V_{2}(p)\right)+h \int_{M_{0}} e^{2 i \psi / h}\left(\bar{a} \widetilde{r}_{12}^{1}+a \overline{\widetilde{r}_{12}^{2}}\right)\left(V_{1}-V_{2}\right) \operatorname{dvol}_{g_{0}}=o(h)
$$

Now to deal with the middle terms, it suffices to apply a Riemann-Lebesgue type argument like Lemma 5.3 of [8] to deduce that it is a $o(h)$. The argument is simply to approximate the amplitude in the $L^{1}\left(M_{0}\right)$ norm by a smooth compactly supported function and then use stationary phase to deal with the smooth function. We have thus proved that $V_{1}(p)=V_{2}(p)$ by taking $h \rightarrow 0$.

## 7. Appendix

To obtain mapping properties of the resolvent of $\Delta_{\mathbb{R}^{2}}$ acting on functions with Gaussian decay, we shall give two Lemmas on Fourier transforms of functions with Gaussian decay.
Lemma 7.1. Let $f(z) \in e^{-\gamma|z|^{2}} L^{2}\left(\mathbb{R}^{2}\right)$ for some $\gamma>0$. Then the Fourier transform $\hat{f}(\xi)$ extends analytically to $\mathbb{C}^{2}$ and for all $\xi, \eta \in \mathbb{R}^{2}$,

$$
\|\hat{f}(\xi+i \eta)\|_{L^{2}\left(\mathbb{R}^{2}, d \xi\right)} \leq 2 \pi e^{\frac{|\eta|^{2}}{4 \gamma}}\left\|e^{\gamma|z|^{2}} f\right\|_{L^{2}\left(\mathbb{R}^{2}\right)}
$$

If $f(z) \in e^{-\gamma|z|^{2}} L^{1}\left(\mathbb{R}^{2}\right)$ for some $\gamma>0$ then

$$
\sup _{\xi \in \mathbb{R}^{2}}|\hat{f}(\xi+i \eta)| \leq e^{\frac{|\eta|^{2}}{4 \gamma}}\left\|e^{\gamma|z|^{2}} f\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}
$$

Proof. The first statement is clear. For the bound, we write

$$
\hat{f}(\xi+i \eta)=e^{\frac{|\eta|^{2}}{4 \gamma}} \int_{\mathbb{R}^{2}} e^{-i \xi \cdot z} e^{-\gamma\left|z-\frac{\eta}{2 \gamma}\right|^{2}} e^{\gamma|z|^{2}} f(z) d z=e^{\frac{|\eta|^{2}}{4 \gamma}} \mathcal{F}_{z \rightarrow \xi}\left(e^{-\gamma\left|z-\frac{\eta}{2 \gamma}\right|^{2}} e^{\gamma|z|^{2}} f(z)\right)
$$

But the function $e^{-\gamma\left|z-\frac{\eta}{2 \gamma}\right|^{2}} e^{\gamma|z|^{2}} f(z)$ is in $L^{2}\left(\mathbb{R}^{2}, d z\right)$ and its norm is bounded uniformly by $\left\|e^{\gamma|z|^{2}} f\right\|_{L^{2}}$, thus it suffices to use the Plancherel theorem to obtain the desired bound. The $L^{\infty}$ bound is similar.

Lemma 7.2. Let $F(\xi+i \eta)$ be a complex analytic function on $\mathbb{R}^{2}+i \mathbb{R}^{2}=\mathbb{C}^{2}$ such that there is $C>0$ and $\gamma>0$ with

$$
\|F(\xi+i \eta)\|_{L^{2}\left(\mathbb{R}^{2}, d \xi\right)} \leq C e^{\frac{|\eta|^{2}}{4 \gamma}} \text { and } \sup _{\xi \in \mathbb{R}^{2}}|F(\xi+i \eta)| \leq C e^{\frac{|\eta|^{2}}{4 \gamma}}
$$

If $F$ vanishes on the real submanifold $\left\{|\xi|^{2}=\lambda^{2}\right\}$, then $\mathcal{F}_{\xi \rightarrow z}^{-1}\left(\frac{F(\xi)}{|\xi|^{2}-\lambda^{2}}\right) \in e^{-\gamma|z|^{2}} L^{\infty}\left(\mathbb{R}^{2}\right)$.
Proof. First by analyticity of $F$, one has that $F$ vanishes on the complex hypersurface $M_{\lambda}:=\left\{\zeta \in \mathbb{C}^{2} ; \zeta . \zeta=\lambda^{2}\right\}$ (see for instance the proof of Lemma 2.5 of [17]), and in particular $G(\zeta)=F(\zeta) /\left(\zeta . \zeta-\lambda^{2}\right)$ is an analytic function on $\mathbb{C}^{2}$. We will first prove that for each $\eta \in \mathbb{R}^{2}$, $G(\xi+i \eta) \in L^{1}\left(\mathbb{R}^{2}, d \xi\right) \cap L^{\infty}\left(\mathbb{R}^{2}, d \xi\right)$ and

$$
\begin{equation*}
\|G(\xi+i \eta)\|_{L^{1}\left(\mathbb{R}^{2}, d \xi\right)} \leq C e^{\frac{\mid \eta \eta^{2}}{4 \gamma}} \tag{26}
\end{equation*}
$$

If $|\eta| \leq 2$ we choose the disc $B:=\left\{\xi \in \mathbb{R}^{2} ;|\xi|^{2}<2\left(4+\lambda^{2}\right)\right\}$ and let $\zeta:=\xi+i \eta$. Then $\|G(\xi+i \eta)\|_{L^{1}(B, d \xi)}$ and $\left\|\left(\zeta . \zeta-\lambda^{2}\right)^{-1}\right\|_{L^{2}\left(\mathbb{R}^{2} \backslash B, d \xi\right)}$ are uniformly bounded for $|\eta| \leq 2$, and we obtain by Cauchy-Schwarz that (26) holds for $|\eta| \leq 2$. For the case $|\eta|>2$ we define $U_{\eta}:=\left\{\xi \in \mathbb{R}^{2} ;\left|\zeta . \zeta-\lambda^{2}\right|>|\eta|\right\}$ and note that

$$
\begin{aligned}
& \sup _{|\eta|>2}\left\|\left(\zeta . \zeta-\lambda^{2}\right)^{-1}\right\|_{L^{1}\left(\mathbb{R}^{2} \backslash U_{\eta}, d \xi\right)}<\infty \\
& \sup _{|\eta|>2}\left\|\left(\zeta . \zeta-\lambda^{2}\right)^{-1}\right\|_{L^{2}\left(U_{\eta}, d \xi\right)}<\infty
\end{aligned}
$$

These results follow by decomposing the integration sets to parts where one can change coordinates $\xi_{1}+i \xi_{2}$ to $\tilde{\xi}_{1}+i \tilde{\xi}_{2}:=\zeta . \zeta-\lambda^{2}$, and by evaluating simple integrals. Then (26) follows from Cauchy-Schwarz and the estimates for $F$.

Let $\eta=2 \gamma z$, we use a contour deformation from $\mathbb{R}^{2}$ to $2 i \gamma z+\mathbb{R}^{2}$ in $\mathbb{C}^{2}$,

$$
\int_{\mathbb{R}^{2}} e^{i z \cdot \xi} G(\xi) d \xi=\int_{\mathbb{R}^{2}} e^{i z .(\xi+2 i \gamma z)} G(\xi+2 i \gamma z) d \xi
$$

which is justified by the fact that $G(\xi+i \eta) \in L^{1}\left(\mathbb{R}^{2} \times K, d \xi d \eta\right)$ for any compact set $K$ in $\mathbb{R}^{2}$ by the uniform bound (26). Now using (26) again shows that

$$
\left|\int_{\mathbb{R}^{2}} e^{i z \cdot \xi} G(\xi) d \xi\right| \leq C e^{-\gamma|z|^{2}}
$$

which ends the proof.
Corollary 7.3. Let $f(z) \in e^{-\gamma|z|^{2}} L^{2}\left(\mathbb{R}^{2}\right) \cap e^{-\gamma|z|^{2}} L^{1}\left(\mathbb{R}^{2}\right)$ for some $\gamma>0$. Assume that its Fourier transform $\hat{f}(\xi)$ vanishes on the sphere $\{|\xi|=|\lambda|\}$, then one has

$$
\mathcal{F}_{\xi \rightarrow z}^{-1}\left(\frac{\hat{f}(\xi)}{|\xi|^{2}-\lambda^{2}}\right) \in e^{-\gamma|z|^{2}} L^{\infty}\left(\mathbb{R}^{2}\right)
$$

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[^0]:    ${ }^{1}$ In [14], a unique continuation is used for Schwartz solutions of $\left(\Delta+V-\lambda^{2}\right) u=0$ when $V$ is a compactly supported potential on $\mathbb{R}^{n}$ but the same result is also true in our setting, this is a consequence of a standard Carleman estimate.

