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## RESONANCES FOR SCHRODINGER OPERATORS ON INPINITE CYLINDERS AND OTHER PRODUCTS

# RESONANCES FOR SCHRÖDINGER OPERATORS ON INFINITE CYLINDERS AND OTHER PRODUCTS 

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#### Abstract

We study the resonances of Schrödinger operators on the infinite product $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$, where $d$ is odd, $\mathbb{S}^{1}$ is the unit circle, and the potential $V$ lies in $L_{c}^{\infty}(X)$. This paper shows that at high energy, resonances of the Schrödinger operator $-\Delta+V$ on $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$ which are near the continuous spectrum are approximated by the resonances of $-\Delta+V_{0}$ on $X$, where the potential $V_{0}$ is given by averaging $V$ over the unit circle. These resonances are, in turn, given in terms of the resonances of a Schrödinger operator on $\mathbb{R}^{d}$ which lie in a bounded set. If the potential is smooth, we obtain improved localization of the resonances, particularly in the case of simple, rank 1 poles of the corresponding scattering resolvent on $\mathbb{R}^{d}$. In that case, we obtain the leading order correction for the location of the corresponding high-energy resonances. In addition to direct results about the location of resonances, we show that at high energies away from the resonances, the resolvent of the model operator $-\Delta+V_{0}$ on $X$ approximates that of $-\Delta+V$ on $X$. If $d=1$, in certain cases this implies the existence of an asymptotic expansion of solutions of the wave equation. Again for the special case of $d=1$, we obtain a resonant rigidity type result for the zero potential among all real-valued smooth potentials.


## 1. Introduction

We study the Schrödinger operator $-\Delta+V$ on the manifold $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$ with the product metric, where $d$ is odd, $\mathbb{S}^{1}$ is the unit circle, and $V \in L_{c}^{\infty}(X)$. In the special case $d=1, X$ is the infinite cylinder $\mathbb{R} \times \mathbb{S}^{1}$. We show that in the large energy limit, resonances near the continuous spectrum are well approximated by those of $-\Delta+V_{0}$, where $V_{0}$ is the average of $V$ over $\mathbb{S}^{1}: V_{0}(x)=\frac{1}{2 \pi} \int_{0}^{2 \pi} V(x, \theta) d \theta$. By a separation of variables argument, these, in turn, are determined by the low energy resonances of the Schrödinger operator $-\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}$ on $\mathbb{R}^{d}$. In the case of smooth potentials $V$, for simple rank 1 poles of the (scattering) resolvent of $-\sum_{j=0}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}$, we find the leading-order corrections to the location of the corresponding poles of the resolvent of $-\Delta+V$ on $X$. Among other things, this allows us to prove that no other smooth real-valued potential on $\mathbb{R} \times \mathbb{S}^{1}$ has the same resonances as the zero potential. For potentials with $V_{0} \equiv 0$, we show the existence of large resonance-free regions. When $d=1$ and $V \in C_{c}^{\infty}(X ; \mathbb{R})$, under certain hypotheses on the potential $V_{0}$ we are able to give an asymptotic expansion of solutions of the wave equation. For the case of $d=1$ we study a simple example of a nontrivial potential $V$ with $V_{0} \equiv 0$ and locate some of the corresponding resonances. Some of these results are reminiscent of Drouot's results [2018] for rapidly oscillating potentials on $\mathbb{R}^{d}$.

[^0]Let $\Delta \leq 0$ denote the Laplacian on $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$ with the product metric. For $V \in L_{c}^{\infty}(X)$ the Schrödinger operator $-\Delta+V$ has continuous spectrum $[0, \infty)$, with multiplicity which increases at each threshold $j^{2}$, for $j \in \mathbb{N}_{0}$. For $\operatorname{Im} \zeta>0$, set $R_{V}(\zeta)=\left(-\Delta+V-\zeta^{2}\right)^{-1}$. This (scattering) resolvent has a meromorphic continuation to $\hat{Z}$, the minimal Riemann surface for which $\tau_{l}(\zeta) \stackrel{\text { def }}{=}\left(\zeta^{2}-l^{2}\right)^{1 / 2}$ is a single-valued analytic function for each $l \in \mathbb{N}_{0}$. The resonances are poles of the resolvent $R_{V}(\zeta)$. We refer to the portion of $\hat{Z}$ for which $\operatorname{Im} \tau_{l}(\zeta)>0$ for all $l \in \mathbb{N}_{0}$ as the physical space. In this set $R_{V}$ is a bounded operator on $L^{2}(X)$, away from a discrete set of points which correspond to (square roots of) eigenvalues. For $l \in \mathbb{N}_{0}$ and $\rho>0$, denote by $B_{l}(\rho)$ the connected component of $\left\{\zeta \in \hat{Z}:\left|\tau_{l}(\zeta)\right|<\rho\right\}$ which nontrivially intersects both the physical space and the set $\left\{\zeta \in \hat{Z}: \operatorname{Re} \tau_{0}(\zeta)>0\right\}$. Using as the coordinate $\tau_{l}(\zeta), B_{l}(\rho)$ is identified with the disk of radius $\rho$ in the complex plane, centered at the origin, and this identification is compatible with the complex structure of $\hat{Z} \upharpoonright_{B_{l}(\rho)}$ if $\rho<\sqrt{2 l-1}$. The point $\tau_{l}(\zeta)=0$ in $B_{l}(\rho)$ corresponds to the $l$-th threshold. We study the resonances of $-\Delta+V$ in $B_{l}(\rho)$, or $B_{l}(\alpha \log l)$, as $l \rightarrow \infty$. Results of Section 6 show that these are the high-energy resonances "near" the continuous spectrum which have $\operatorname{Re} \tau_{0}>0$.

For a function $V \in L_{c}^{\infty}(X)$ and $m \in \mathbb{Z}$ define

$$
V_{m}(x)=\frac{1}{2 \pi} \int_{0}^{2 \pi} V(x, \theta) e^{-i m \theta} d \theta
$$

so that $V(x, \theta)=\sum_{m=-\infty}^{\infty} V_{m}(x) e^{i m \theta}$. The minimal assumption on a potential $V$ in most of this paper will be that

$$
\begin{equation*}
V \in L_{c}^{\infty}(X) \quad \text { and } \quad\left\|V_{m}\right\|_{L^{\infty}}=O\left(|m|^{-\delta}\right) \quad \text { for some } \delta \text { with } 0<\delta \leq \frac{1}{2} \tag{1-1}
\end{equation*}
$$

Note that this imposes an assumption on $\delta$ as well, which we shall include when we invoke hypothesis (1-1). We use the notation $\Delta_{0}=\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}$ for the Laplacian on $\mathbb{R}^{d}$,

$$
\begin{equation*}
R_{V_{0}, 0}(\lambda)=\left(-\Delta_{0}+V_{0}-\lambda^{2}\right)^{-1}, \quad \text { if } \operatorname{Im} \lambda>0 \tag{1-2}
\end{equation*}
$$

with the same notation for its meromorphic continuation to the complex plane-see Section 3A. The poles of $R_{V_{0}, 0}$ in $\mathbb{C}$ are the resonances of $-\Delta_{0}+V_{0}$. The multiplicity $m_{V_{0}, 0}\left(\lambda_{0}\right)$ of a resonance of $-\Delta_{0}+V_{0}$ at $\lambda_{0}$ is given by the dimension of the range of the singular part of the resolvent at $\lambda_{0}$; this is discussed further in Section 4.
Theorem 1.1. Let $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$, d odd, and let $V \in L_{c}^{\infty}(X)$ satisfy $\left\|V_{m}\right\|_{L^{\infty}}=O\left(|m|^{-\delta}\right)$ for some $\delta$ with $0<\delta \leq \frac{1}{2}$. Suppose $\lambda_{0} \in \mathbb{C}, \lambda_{0} \neq 0$, is a resonance of $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$, of multiplicity $m_{V_{0}, 0}\left(\lambda_{0}\right)$. Let $\rho \in \mathbb{R}, \rho>\left|\lambda_{0}\right|$. Then there are $C_{0}>0, L>0$ so that for $l>L, l \in \mathbb{N}$ there are exactly $2 m_{V_{0}, 0}\left(\lambda_{0}\right)$ resonances, when counted with multiplicity, of $-\Delta+V$ in the set

$$
\left\{\zeta \in B_{l}(\rho):\left|\tau_{l}(\zeta)-\lambda_{0}\right|<C_{0} l^{-\delta /\left(m_{V_{0}, 0}\left(\lambda_{0}\right)\right)}\right\}
$$

Here, and elsewhere in the paper, the apparent "doubling" of the number of poles (when counted with multiplicity) on $X$ as compared with those on $\mathbb{R}^{d}$ is due to the fact that for $l \in \mathbb{N}, l^{2}$ is an eigenvalue of $-d^{2} / d \theta^{2}$ on $\mathbb{S}^{1}$ of multiplicity two. This can be seen immediately in the simplest case, $V \equiv V_{0}$, by separating variables.


Figure 1. A schematic showing resonances of $-\Delta+V$ in $B_{l}(\rho)$, pictured in the $\tau_{l}$-coordinate. Each red $\times$ indicates a single resonance of even multiplicity or a cluster of resonances. The hatched region indicates the portion of $B_{l}(\rho)$ which lies in the physical space. By comparing Figure 2, Section 3B one can see how this fits in the larger picture.

In this paper we refer to any pole of the resolvent as a resonance, including those which correspond to eigenvalues. The second part of Theorem 1.2, for which $V$ is assumed to be smooth, implies an improved localization of the resonances for smooth potentials.

The next theorem shows that, other than possible poles near the threshold, the poles as described above are all the poles in $B_{l}(\rho)$ for sufficiently large $l$.

Theorem 1.2. Let $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$, d odd, and suppose $V$ satisfies the hypothesis (1-1). Choose $\rho>0$ so that if $\lambda_{j}$ is a pole of $R_{V_{0}, 0}(\lambda)$, then $\left|\lambda_{j}\right| \neq \rho$. Set

$$
\Lambda_{\rho}=\left\{\lambda_{j} \in \mathbb{C}:\left|\lambda_{j}\right|<\rho \text { and } \lambda_{j} \text { is a pole of } R_{V_{0}, 0}(\lambda)\right\} .
$$

Let $\epsilon^{\prime}>0$ be so that $\epsilon^{\prime}<\min \left\{\left|\lambda_{j}\right|: \lambda_{j} \in \Lambda_{\rho}, \lambda_{j} \neq 0\right\}$. Then there are $\widetilde{C}, L>0$ so that for $l>L, l \in \mathbb{N}$, there are no resonances of $-\Delta+V$ in

$$
\left\{\zeta \in B_{l}(\rho):\left|\tau_{l}(\zeta)\right|>\epsilon^{\prime} \text { and }\left|\tau_{l}(\zeta)-\lambda_{j}\right|>\widetilde{C} l^{-\delta / m_{V_{0}, 0}\left(\lambda_{j}\right)} \text { for all } \lambda_{j} \in \Lambda_{\rho}\right\}
$$

Moreover, if $V$ is smooth for perhaps larger $L$ and $\widetilde{C}$, for $l>L$ there are no resonances in

$$
\left\{\zeta \in B_{l}(\rho):\left|\tau_{l}(\zeta)\right|>\epsilon^{\prime} \text { and }\left|\tau_{l}(\zeta)-\lambda_{j}\right|>\widetilde{C} l^{-2 /\left(m_{V_{0}, 0}\left(\lambda_{j}\right)\right)} \text { for all } \lambda_{j} \in \Lambda_{\rho}\right\}
$$

In addition, if $R_{V_{0}, 0}(\lambda)$ is analytic in a neighborhood of the origin, then there are no poles in $B_{l}\left(\epsilon^{\prime}\right)$ for $l$ sufficiently large.

We comment that smoothness of the potential $V$ is more than is needed for the second part of Theorem 1.2. It would suffice to have $V \in C^{k}(X)$, for some $k$ sufficiently large. In order to simplify the proofs, we have not tracked the value of $k$ which is needed.

To help visualize these theorems, we include Figure 1, which is a schematic showing the resonances of $-\Delta+V$ in $B_{l}(\rho)$ for large $l$, using the $\tau_{l}$-coordinate. This schematic is familiar from odd-dimensional scattering theory; that this should be so is a consequence of Theorems $1.1-1.3$. One difference is that in this diagram, the only portion of $B_{l}(\rho)$ which lies in the physical space is the portion which is in the
first quadrant, indicated by hatching. Another is that each $x$ indicates either a single resonance of even multiplicity, or a cluster of resonances. See Figure 2 to see how $B_{l}$ fits in a larger context.

For Schrödinger operators on $\mathbb{R}^{d}$, the behavior of the singularities of the resolvent at the origin is delicate. For example, notions of multiplicity of a resonance which agree at points away from the origin may differ at the origin; see [Dyatlov and Zworski 2019, Theorem 2.8]. These same sorts of issues arise at thresholds in the case under study here, and accounts for the fact that this next theorem, which concerns resonances very near the thresholds, is weaker than the previous ones.

Theorem 1.3. Let $V$ satisfy (1-1) and suppose the resolvent of $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$ has a pole at 0 of order $r>0$, and multiplicity $m_{V_{0}, 0}(0)$ Then there are $C, L>0$ so that $-\Delta+V$ on $X$ has at least $2 m_{V_{0}, 0}(0)$ resonances, when counted with multiplicity, in $B_{l}\left(\mathrm{Cl}^{-\delta / r}\right)$ when $l>L, l \in \mathbb{N}$. Moreover, there is an $\epsilon>0$ so that $-\Delta+V$ has no poles in $B_{l}(\epsilon) \backslash B_{l}\left(C l^{-\delta / r}\right)$ when $l>L$. If $V \in C_{c}^{\infty}(X)$, then this can be improved to show that there is a $C_{1}>0$ so that $-\Delta+V$ has no poles in $B_{l}(\epsilon) \backslash B_{l}\left(C_{1} l^{-2 / r}\right)$ when $l>L$. Moreover, under the hypothesis (1-1), if $r=1$ there are exactly $2 m_{V_{0}, 0}(0)$ resonances of $-\Delta+V$ in $B_{l}\left(C l^{-\delta}\right)$ for $l>L$.

Suppose for the moment that $V_{0}$ is real-valued. In this case, it is well known that if $d=1$ the order of the pole of the resolvent of $-d^{2} / d x^{2}+V_{0}$ at 0 cannot exceed 1 , and if it is 1 , then $m_{V_{0}, 0}(0)=1$ [Dyatlov and Zworski 2019, Theorem 2.7]. If $d \geq 3$ is odd, then the order of the pole of the resolvent of $-\Delta_{0}+V_{0}$ at 0 cannot exceed 2 [Dyatlov and Zworski 2019, Lemma 3.16]. For general $V$ and $r$, the order of the pole at 0 can be bounded from above in terms of $m_{V_{0}, 0}(0)$, and in the case $d=1, m_{V_{0}, 0}(0)$ can be bounded above by $r$.

It is of particular interest to understand poles of the resolvent $R_{V}$ near the physical region. In Section 6 we show that there are large regions near the physical region that contain no resonances. A consequence of those results is that large energy resonances near the continuous spectrum and having $\operatorname{Re} \tau_{0}(\zeta)>0$ are contained in regions of the form $B_{l}(\rho)$, where $\rho$ depends on how near the continuous spectrum we wish to look. In Section 6 we further justify our focus on the resonances in sets $B_{l}(\rho)$.

Theorems 1.1-1.3 combined with results of Section 6 yield the following corollary. Here $d \hat{z}$ is a distance on $\hat{Z}$, defined in Section 6. The boundary of the physical region corresponds to the continuous spectrum. In the corollary, we use $\left\{\zeta_{j}^{b}\right\}$ to denote a sequence of points in $\hat{Z}$, to distinguish them from $\zeta_{l}$ which is used elsewhere to denote a particular mapping from an open subset of the complex plane into $\hat{Z}$.

Corollary 1.4. Let $V \in L_{c}^{\infty}(X ; \mathbb{R})$ satisfy (1-1). Then $R_{V}(\zeta)$ has a sequence $\left\{\zeta_{j}^{b}\right\}_{j=1}^{\infty}$ of poles satisfying both $\left|\tau_{0}\left(\zeta_{j}^{b}\right)\right| \rightarrow \infty$ as $j \rightarrow \infty$ and $d_{\hat{Z}}\left(\zeta_{j}^{b}\right.$, physical region $) \rightarrow 0$ as $j \rightarrow \infty$ if and only if $R_{V_{0}, 0}(\lambda)$ has at least one pole in $i[0, \infty)$.

In particular, if $d=1$, by [Reed and Simon 1978, Theorem XIII.110] if $\int_{X} V \leq 0$ then $R_{V}(\zeta)$ has such a sequence of poles approaching the physical space. In contrast, if $V_{0}(x) \geq 0$ for all $x$ and $V_{0}$ is nontrivial, $R_{V}(\zeta)$ does not have such a sequence of poles. Note that for any fixed $k_{0} \in \mathbb{N}$, we have $\left|\tau_{0}\left(\zeta_{j}^{b}\right)\right| \rightarrow \infty$ as $j \rightarrow \infty$ if and only if $\left|\tau_{k_{0}}\left(\zeta_{j}^{b}\right)\right| \rightarrow \infty$ as $j \rightarrow \infty$. We remark that we could prove an analog of Corollary 1.4 for complex-valued potentials as well.

If we enlarge the region centered at the threshold $l^{2}$ with increasing $l$, we have less fine localization of the resonances, see Theorem 7.1. However, when $V_{0}$, the average of the potential, is identically zero, we can get a larger resonance-free region. The difference in the next result for $d=1$ and $d \geq 3$ is due to the fact that the resolvent of $-d^{2} / d x^{2}$ on $\mathbb{R}$ has a pole at the origin, but that of $-\Delta_{0}$ on $\mathbb{R}^{d}$ for $d \geq 3$ odd does not.

Theorem 1.5. Let $V \in L_{c}^{\infty}(X)$ satisfy (1-1), and suppose $V_{0} \equiv 0$. If $d=1$ there are $\alpha, c_{0}>0$ so that for $l \in \mathbb{N}$ sufficiently large there are no resonances of $-\Delta+V$ in the $\operatorname{set}\left\{\zeta \in B_{l}(\alpha \log l):\left|\tau_{l}(\zeta)\right|>c_{0} / l^{\delta}\right\}$. If $d \geq 3$ is odd, there is an $\alpha>0$ so that for $l$ sufficiently large there are no resonances of $-\Delta+V$ in the set $B_{l}(\alpha \log l)$.

There is a sense in which this theorem is sharp; see Proposition 12.6 for a computation for the case $d=1$ with the potential $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$, where $\chi_{I_{0}}$ is the characteristic function of the interval $[-1,1]$.

We can find the leading correction term for high-energy resonances of $-\Delta+V$ which correspond to simple resonances of $-\Delta_{0}+V_{0}$. In the next theorem, $\nabla_{0}$ is the gradient on $\mathbb{R}^{d}$, so that

$$
\nabla_{0} f=\left(\frac{\partial}{\partial x_{1}} f, \frac{\partial}{\partial x_{2}} f, \ldots, \frac{\partial}{\partial x_{d}} f\right)
$$

Theorem 1.6. Let $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$, d odd, $V \in C_{c}^{\infty}(X)$, and suppose $\lambda_{0} \in \mathbb{C}$ is a simple pole of the scattering resolvent $R_{V_{0}, 0}$ of $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$, and that the residue of $R_{V_{0}, 0}$ at $\lambda_{0}$ has rank 1 . Suppose for any $\chi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
\chi\left(R_{V_{0}, 0}(\lambda)-\frac{i}{\lambda-\lambda_{0}} u \otimes u\right) \chi \tag{1-3}
\end{equation*}
$$

is analytic near $\lambda=\lambda_{0}$. Let $\rho>\left|\lambda_{0}\right|$. Then there are $\epsilon, L>0$ so that for $l>L$ there are exactly two poles of $R_{V}(\zeta)$, when counted with multiplicity, in $\left\{\zeta \in B_{l}\left(\left|\lambda_{0}\right|+1\right):\left|\tau_{l}(\zeta)-\lambda_{0}\right|<\epsilon\right\}$, and each pole of $R_{V}(\zeta)$ in this set satisfies

$$
\tau_{l}(\zeta)=\lambda_{0}-\frac{i}{4 l^{2}} \sum_{k \neq 0} \frac{1}{k^{2}} \int_{\mathbb{R}}\left(k^{2} V_{-k} V_{k} u^{2}+\left(\nabla_{0} V_{-k} \cdot \nabla_{0} V_{k}\right) u^{2}\right)(x) d x+O\left(l^{-3}\right)
$$

We note that the normalization of the singularity in (1-3) is chosen so that if $V$ is real-valued and $\lambda_{0} \in i[0, \infty)$, then $u$ is real-valued. There is some further discussion of $u$ at the beginning of Section 10 . Proposition 12.3 shows that the leading correction may be rather different for a nonsmooth potential by considering the special case of the potential on $\mathbb{R} \times \mathbb{S}^{1}$ given by $V(x, \theta)=2 \cos \theta \chi_{I_{0}}(x)$, where $\chi_{I_{0}}$ is the characteristic function of the interval $[-1,1]$. As for Theorem 1.2, the proof of Theorem 1.6 only needs $V$ to be $C^{k}$ for some $k$ sufficiently large. Since Theorem 1.8 requires smoothness of the potential only for an application of Theorem 1.6, the same is true for it. Again, we have chosen not to track this value in the interest of simplifying proofs.

If $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d} ; \mathbb{R}\right)$ and the operator $-\Delta_{0}+V_{0}$ on $L^{2}\left(\mathbb{R}^{d}\right)$ has a simple negative eigenvalue $-\beta^{2}$, then this negative eigenvalue corresponds to a simple pole of $R_{V_{0}, 0}$ on the positive imaginary axis at $i|\beta|$, and the residue has rank 1. By Theorem 1.1 (or Corollary 1.4), in this case $R_{V}$ has a sequence of poles approaching the physical space. If $V \in C_{c}^{\infty}(X ; \mathbb{R})$, the poles approach the physical space very rapidly.

Theorem 1.7. Suppose $V \in C_{c}^{\infty}(X ; \mathbb{R})$ and $\lambda_{0} \in \mathbb{C}$ is a simple pole of $R_{V_{0}}(\lambda)$ with $\operatorname{Re} \lambda_{0}=0$, with residue of $R_{V_{0}}$ at $\lambda_{0}$ having rank 1. Then there is an $\epsilon>0$ so that if $\left\{\zeta_{l}^{b}\right\}_{l=L}^{\infty} \subset \hat{Z}$ is a sequence of poles of $R_{V}$ with $\zeta_{l}^{\text {b }} \in B_{l}\left(\left|\lambda_{0}\right|+1\right)$ and $\left|\tau_{l}\left(\zeta_{l}^{b}\right)-\lambda_{0}\right|<\epsilon$, then $\operatorname{Re} \tau_{l}\left(\zeta_{l}^{b}\right)=O\left(l^{-N}\right)$ for any $N$. In particular, this implies that if $\operatorname{Im} \lambda_{0}>0$, then $d_{\hat{Z}}\left(\zeta_{l}^{b}\right.$, physical region $)=O\left(l^{-N}\right)$.

Proposition 12.3 demonstrates the necessity of assuming some regularity of the potential, at least for $d=1$, by studying the resonances very near the $l$-th threshold for a certain real-valued potential with a jump singularity. These resonances in $B_{l}(1)$ arise from the pole of $R_{0,0}(\lambda)$ at $\lambda_{0}=0$. They have

$$
\left|\tau_{l}\left(\zeta_{l}^{b}\right)\right|=O\left(l^{-3 / 2}\right)
$$

and, for a subsequence of $l$ 's tending to infinity,

$$
\left|\operatorname{Im}\left(\tau_{l}\left(\zeta_{l}^{\mathrm{b}}\right)\right)\right|>\frac{1}{10} l^{-3 / 2}
$$

This paper was initially motivated by the case $d=1$, as $\mathbb{R} \times \mathbb{S}^{1}$ provides a particularly simple example of a manifold with infinite cylindrical ends and as such provides a testing ground for studying resonances for Schrödinger operators on such manifolds. Most of the proofs of the preceding theorems are essentially identical for any odd dimension of the factor $\mathbb{R}^{d}$, so we have included the more general results. However, Theorems 1.8 and 1.9 are particular to the $d=1$ case.

As a corollary of Theorems 1.1, 1.3, and 1.6, we get in the case $d=1$ a uniqueness-type result for the zero potential among smooth real-valued potentials.

Theorem 1.8. Let $V \in C_{c}^{\infty}\left(\mathbb{R} \times \mathbb{S}^{1} ; \mathbb{R}\right)$. Suppose for each $\rho>0$ there is a sequence

$$
\left\{l_{j}\right\}_{j=1}^{\infty}=\left\{l_{j}(\rho)\right\}_{j=1}^{\infty} \subset \mathbb{N}
$$

with $l_{j} \rightarrow \infty$ when $j \rightarrow \infty$ so that in $B_{l_{j}}(\rho)$ the resonances of $-\Delta+V$ and $-\Delta$ on $X=\mathbb{R} \times \mathbb{S}^{1}$ are the same. Then $V \equiv 0$.

This result is false if we omit the hypothesis that $V$ is real-valued. For example, for $V_{1} \in C_{c}^{\infty}(\mathbb{R})$ set $V(x, \theta)=V_{1}(x) e^{i \theta}$. Then the operators $-\Delta+V$ and $-\Delta$ have the same resonances; see [Autin 2011] or [Christiansen 2004, Section 4]. This example can be easily generalized.

As part of our study of the distribution of resonances, we prove that, in a suitable sense, near the physical region of $\hat{Z}, R_{V}$ is well approximated by $R_{V_{0}}$ away from the poles of $R_{V_{0}}$; see Proposition 5.4 and Lemma 5.5. In the case $d=1$, this and results of [Christiansen and Datchev 2022] give a wave expansion; see Theorem 1.9.

Let $X=\mathbb{R} \times \mathbb{S}^{1}, V \in C_{c}^{\infty}(X ; \mathbb{R})$, and suppose $-\Delta+V$ has finitely many eigenvalues $\mu_{1}, \mu_{2}, \ldots, \mu_{J}$, repeated with multiplicity, with associated orthonormal eigenfunctions $\left\{\eta_{j}\right\}$, so that $(-\Delta+V) \eta_{j}=\mu_{j} \eta_{j}$. Let $u$ satisfy

$$
\begin{gather*}
\frac{\partial^{2}}{\partial t^{2}} u-\Delta u+V u=0  \tag{1-4}\\
\left(u, u_{t}\right) \upharpoonright_{t=0}=\left(f_{1}, f_{2}\right) \in C_{c}^{\infty}(X) \times C_{c}^{\infty}(X)
\end{gather*}
$$

Theorem 1.9. Let $X=\mathbb{R} \times \mathbb{S}^{1}$ and $V, f_{1}, f_{2} \in C_{c}^{\infty}(X)$, with $V$ real-valued, and suppose $-d^{2} / d x^{2}+V_{0}$ on $\mathbb{R}$ has no negative eigenvalues and no resonance at 0 . Let $u$ be the solution of $(1-4)$ on $[0, \infty) \times X$. Then for each $k_{0} \in \mathbb{N}$ we can write $u(t)=u_{e}(t)+u_{\mathrm{thr}, k_{0}}(t)+u_{r, k_{0}}(t)$, where

$$
\left.\left.\left.\begin{array}{rl}
u_{e}(t, x, \theta)= & \sum_{\substack{\mu_{j} \in \sigma_{p}(-\Delta+V) \\
\mu_{j} \neq 0}} \eta_{j}(x, \theta)\left(\cos \left(\left(\mu_{j}\right)^{1 / 2} t\right)\left\langle f_{1}, \eta_{j}\right\rangle\right.
\end{array}\right) \frac{\sin \left(\left(\mu_{j}\right)^{1 / 2} t\right)}{\left(\mu_{j}\right)^{1 / 2}}\left\langle f_{2}, \eta_{j}\right\rangle\right)\right)
$$

and

$$
u_{\mathrm{thr}, k_{0}}(t, x, \theta)=b_{0,0,+}(x, \theta)+\sum_{k=0}^{k_{0}-1} t^{-1 / 2-k} \sum_{j=1}^{\infty}\left(e^{i t j} b_{j, k,+}(x, \theta)+e^{-i t j} b_{j, k,-}(x, \theta)\right)
$$

for some $b_{j, k, \pm} \in\langle x\rangle^{1 / 2+2 k+\epsilon} L^{2}(X)$. For any $\chi \in C_{c}^{\infty}(X)$ there is a constant $C$ so that

$$
\sum_{j=1}^{\infty}\left\|\chi b_{j, k, \pm}\right\|_{L^{2}(X)}<C, \quad k=0,1,2, \ldots, k_{0}-1
$$

and

$$
\left\|\chi u_{r, k_{0}}(t)\right\|_{L^{2}(X)} \leq C t^{-k_{0}} \quad \text { for } t \text { sufficiently large. }
$$

The assumption that $-d^{2} / d x^{2}+V_{0}$ on $\mathbb{R}$ has no negative eigenvalues and no resonance at 0 means, by Theorem 1.2, that $R_{V}$ has at most finitely many poles on the boundary of the physical space. In particular, this means at most finitely many eigenvalues of $-\Delta+V$, so that the sum in $u_{e}$ is finite. Further, there are at most finitely many poles at thresholds, and this implies via results of [Christiansen and Datchev 2022] that at most finitely many of the $b_{j, 0, \pm}$ are nonzero.

If $-d^{2} / d x^{2}+V_{0}$ on $\mathbb{R}$ has one or more negative eigenvalues, it seems plausible that there is an asymptotic expansion of solutions of the wave equation on compact sets. Since in this case by Theorem 1.7 the resolvent $R_{V}$ may have a sequence of poles rapidly approaching, but not lying in, the continuous spectrum, such an expansion would need to take these into account and is more complicated - see for example [Tang and Zworski 2000] for an expansion in a Euclidean scattering setting with resonances approaching the continuous spectrum. In our setting proving the existence of such an expansion may use techniques similar to those of [Christiansen and Datchev 2022] but does not follow directly from the results of that work. Proving this is outside the scope of this paper.

In this paper we have, for simplicity, limited ourselves to the case of Schrödinger operators on $\mathbb{R}^{d} \times \mathbb{S}$. However, many of our results for $L^{\infty}$ potentials hold as well for Schrödinger operators with Dirichlet or Neumann boundary conditions on $\mathbb{R}^{d-1} \times(0, \infty) \times \mathbb{S}$ or on $\mathbb{R}^{d} \times(0, \pi)$.

1A. Relation to previous work. This paper was inspired in part by two different sets of papers. The first are papers which study eigenvalues and resonances of Schrödinger operators on $\mathbb{R}^{d}$ with rapidly oscillating potentials, and includes [Borisov 2006; Borisov and Gadylshin 2006; Duchêne and Weinstein 2011; Duchêne et al. 2014; 2015; Dimassi 2016; Drouot 2018]. Of these the most closely related to this
paper is that of Drouot [2018], which studies the distribution of resonances of Schrödinger operators $-\Delta_{0}+V_{\epsilon}$ on $\mathbb{R}^{d}$ with $d$ odd. Here

$$
V_{\epsilon}(x)=V_{0}(x)+\sum_{k \in \mathbb{Z}^{d}, k \neq 0} V_{k}(x) e^{i k \cdot x / \epsilon}, \quad x \in \mathbb{R}^{d}
$$

Drouot shows in quantitative ways that in the limit $\epsilon \downarrow 0$, resonances of $-\Delta_{0}+V_{\epsilon}$ near the continuous spectrum are well approximated by those of $-\Delta_{0}+V_{0}$. In addition, he proves some refinements related, for example, to the leading order correction of the positions of the resonances. Theorems 1.1, 1.2, 1.3, 1.5, and 1.6, as well as some computations in Section 12, are inspired by results in [Drouot 2018]. However, the proofs are quite different. In part, this is because the different setting requires different techniques. Additionally, Drouot's results come mainly from studying regularized determinants. While this has the potential of producing in some instances more refined results than we obtain here, it requires a substantial amount of technical work. We have chosen instead to mostly avoid determinants, or to work only with determinants of operators of the type $I+F$, where $F$ is finite rank. Instead, we use an operator Rouché theorem of Gohberg and Sigal [1971]. In some places this may allow for sharper results than could be obtained by using a regularized determinant. We note in addition that in the setting of [Drouot 2018], the resonances lie on the complex plane, while for us, the resonances lie on a Riemann surface which is a countable but infinite cover of the complex plane, with infinitely many branch points. This means that some of the techniques used in [Drouot 2018] cannot be applied here.

A less direct source of inspiration is work done on the distribution of eigenvalues of the Schrödinger operator $-\Delta_{\mathbb{S}^{n}}+W$ on the sphere $\mathbb{S}^{n}$ (and certain other compact manifolds), $n \geq 2$; see for example [Weinstein 1977; Widom 1979]. In this setting, eigenvalues of the Schrödinger operator occur in bands. Roughly speaking, these authors show that a suitable average of the potential $W$ can be used to obtain information about the distribution of high-energy eigenvalues of the Schrödinger operator within these bands. This average is over closed geodesics, rather than over all of $\mathbb{S}^{n}$. Of course, our function $V_{0}(x)$ is the average of the potential $V$ over the cross section of $\mathbb{S}^{1}$, the unique closed geodesic on $\mathbb{S}^{1}$.

This paper was originally motivated by the $d=1$ case, which gives $X=\mathbb{R} \times \mathbb{S}^{1}$, a manifold with an infinite cylindrical end. The spectral and scattering theory of manifolds with infinite cylindrical ends has been studied in, for example, [Goldstein 1974; Guillopé 1989; Melrose 1993]. There is a large literature studying the existence of eigenvalues and, in certain settings, the locations of resonances for such manifolds and the related problems of waveguides which have a "one-dimensional infinity" as our $d=1$ case does; see, e.g., [Levitin and Marletta 2008] or the monograph [Exner and Kovařík 2015]. This monograph also includes some results for manifolds with "higher-dimensional infinity". Many of these results focus on low-energy eigenvalues or resonances. We mention the papers [Christiansen 2002; 2004; Christiansen and Datchev 2021; Christiansen and Zworski 1995; Parnovski 1995; Edward 2002] which are more directly connected with high-energy behavior.

1B. Comments regarding other product manifolds. This paper studies only Schrödinger operators on $\mathbb{R}^{d} \times \mathbb{S}^{1}$, where $d$ is odd. Here we comment on why we require that $d$ be odd and on the choice of $\mathbb{S}^{1}$ for the second factor.

For Euclidean scattering, e.g., for the Schrödinger operator $-\Delta_{\mathbb{R}^{d}}+V_{\mathbb{R}^{d}}$ on $\mathbb{R}^{d}$ with $V_{\mathbb{R}^{d}} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, the space to which the resolvent continues is determined by the dimension: for odd $d$ the meromorphic continuation is to the complex plane, and for even $d$ the meromorphic continuation is to $\Lambda$, the logarithmic cover of $\mathbb{C} \backslash\{0\}$. This means that certain questions related to the distribution of resonances are more difficult in even dimensional Euclidean scattering than in odd dimensional Euclidean scattering. For the problem we consider here, the Riemann surface on which the resonances live is a bit involved to describe when $d$ is odd; see Section 3B. The Riemann surface when $d$ is even is much more complicated, requiring as its building block $\Lambda$ rather than $\mathbb{C}$. It is, however, clear that some of our results, appropriately interpreted, hold if $d$ is even as long as we stay away from thresholds. In the interest of clarity we do not pursue this here.

Next we turn to the choice of the factor $\mathbb{S}^{1}$. There are three things that make this an especially nice choice:
(1) The spacing between the distinct eigenvalues grows as the eigenvalues grow.
(2) Upon averaging in $\mathbb{S}^{1}$, we get a model operator that we understand fairly well.
(3) There is a choice of eigenfunctions of the Laplacian on $\mathbb{S}^{1}$ so that a product of two eigenfunctions is again an eigenfunction: $e^{i j \theta} e^{i k \theta}=e^{i(j+k) \theta}$.

Not all of our results require this last property. In view of [Weinstein 1977; Widom 1979], it would be natural to think of replacing $\mathbb{S}^{1}$ with $\mathbb{S}^{m}$. Of course, the spacing of distinct eigenvalues of the Laplacian on $\mathbb{S}^{m}$ is similar to that for $\mathbb{S}^{1}$. However, when using a factor $\mathbb{S}^{m}$ with $m>1$, obtaining a model operator is much more complicated, and it seems any results for general potentials would likely be substantially weaker.

1C. Ideas from the proofs. Our starting point for the study of resonances of $-\Delta+V$ is an identification of the resonances with the points $\zeta$ for which the operator $I+\left(V-V_{0}\right) R_{V_{0}}(\zeta) \chi$ has nontrivial null space. Here $R_{W}(\zeta)$ is the meromorphic continuation of the resolvent of $-\Delta+W$, and $\chi \in L_{c}^{\infty}(X)$ satisfies $\chi V=V$ and is, for convenience, chosen independent of $\theta$. By separating variables, we can understand $R_{V_{0}}$ in terms of the resolvent of $-\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}(x)$ on $\mathbb{R}^{d}$.

We use two well-known and related properties of the resolvent of $-\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}(x)$ on $\mathbb{R}^{d}$. One is the estimate

$$
\left\|\tilde{\chi}\left(-\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}-(\lambda+i 0)^{2}\right)^{-1} \tilde{\chi}\right\|=O\left(|\lambda|^{-1}\right)
$$

as $\lambda \rightarrow \infty$ for $\lambda \in \mathbb{R}$ and $\tilde{\chi} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. The second is the existence of a logarithmic resonance-free neighborhood of the real axis.

An immediate consequence of this second fact and the fact that the distance between thresholds of our operator $-\Delta+V$ on $X$ increases at high energy is that if $V=V_{0}$, at high energy near the thresholds the resonances of $-\Delta+V_{0}$ on $X$ are determined by low-energy resonances of $-\sum_{j=1}^{d} \partial^{2} / \partial x_{j}^{2}+V_{0}$ on $\mathbb{R}^{d}$. Moreover, using these facts and an operator Rouché theorem of Gohberg and Sigal [1971], we are able to show that at high energy near the thresholds the zeros of $I+\left(V-V_{0}\right) R_{V_{0}} \chi$ are approximated by the
poles of $\chi R_{V_{0}} \chi$. These ideas underlie the proofs of the $L^{\infty}$ results of Theorems 1.1-1.3 and 1.5. They are also central to the proofs of the smooth versions of these results and of Theorem 1.6, although these proofs require some additional study of the resolvent of $-\Delta+V_{0}$ when $V_{0}$ is smooth.

1D. Organization. In Section 3 we recall some results from Euclidean scattering and show that the resolvent of $-\Delta+V$ on $X$ has a meromorphic continuation to $\hat{Z}$. (We note that this latter is known; see Section 3 for references.) We define the multiplicity of a pole of the resolvent, and give several useful identities involving it in Section 4. In addition, this section introduces some notation and results related to the operator Rouché theorem of Gohberg and Sigal [1971]. With these preliminaries we prove Theorems 1.1 and 1.2 in the case of an $L^{\infty}$ potential $V$, using results from [Gohberg and Sigal 1971]. Section 6 contains more discussion of the Riemann surface $\hat{Z}$ and shows the existence of resonance-free regions which are, at high energy, near the physical region and away from thresholds. This provides the missing pieces of the proof of Corollary 1.4. Combining these with the resolvent estimates of Section 5 and results of [Christiansen and Datchev 2022] proves Theorem 1.9.

Section 8 contains preliminary computations which are needed to refine our results for smooth potentials. The smooth case of Theorem 1.2 is proved with techniques similar to that of the $L^{\infty}$ result, but using in addition results of Section 8.

In Section 10 we prove Theorems 1.6 and 1.7. We do this using Fredholm determinants, but determinants of the form $\operatorname{det}(I+F)$, where $F$ is a finite-rank operator. Theorem 1.8 follows rather directly from the earlier results. Finally, in Section 12, in the case $d=1$ we give approximations of some of the high-energy resonances for a particularly simple potential which has $V_{0} \equiv 0$ and which is not smooth.

## 2. Notation and conventions

On $X=\mathbb{R}^{d} \times \mathbb{S}^{1}$ we use the coordinates $(x, \theta)$ or $\left(x^{\prime}, \theta^{\prime}\right)$, with $x, x^{\prime} \in \mathbb{R}^{d}$ and $\theta, \theta^{\prime} \in[0,2 \pi)$.
Throughout the paper, $V \in L_{c}^{\infty}(X)$ and $l \in \mathbb{N}_{0}$, and the dimension $d$ of $\mathbb{R}^{d}$ is odd. We use $C$ to stand for a positive constant, the value of which may change without comment.

Suppose $A$ and $B$ are linear operators on domains in $L^{2}\left(\mathbb{R}^{d}\right)$ and $L^{2}\left(\mathbb{S}^{1}\right)$, respectively, and are given by

$$
(A f)(x)=\int_{\mathbb{R}^{d}} A\left(x, x^{\prime}\right) f\left(x^{\prime}\right) d x^{\prime} \quad \text { and } \quad(B g)(\theta)=\int_{0}^{2 \pi} B\left(\theta, \theta^{\prime}\right) g\left(\theta^{\prime}\right) d \theta
$$

Then $A$ and $B$ give rise to linear operators on domains in $L^{2}(X)$, which we again denote by $A$ and $B$, and which are given by

$$
(A h)(x, \theta)=\int_{\mathbb{R}^{d}} A\left(x, x^{\prime}\right) h\left(x^{\prime}, \theta\right) d x^{\prime} \quad \text { and } \quad(B h)(x, \theta)=\int_{\mathbb{R}^{d}} B\left(\theta, \theta^{\prime}\right) h\left(x, \theta^{\prime}\right) d \theta^{\prime}
$$

For $f, g \in L^{2}\left(\mathbb{R}^{d}\right)$, the operator $f \otimes g: L^{2}\left(\mathbb{R}^{d}\right) \rightarrow L^{2}\left(\mathbb{R}^{d}\right)$ is defined via

$$
((f \otimes g) h)(x)=f(x) \int_{\mathbb{R}^{d}} g\left(x^{\prime}\right) h\left(x^{\prime}\right) d x^{\prime}
$$

If $f, g \in L^{2}(X)$, the operator $f \otimes g$ on $L^{2}(X)$ is defined analogously.

We list some repeatedly used notation for the convenience of the reader:

- The Laplacians on $\mathbb{R}^{d}$ and $X$ are given, respectively, by

$$
\Delta_{0}=\sum_{j=1}^{d} \frac{\partial^{2}}{\partial x_{j}^{2}} \quad \text { and } \quad \Delta=\sum_{j=1}^{d} \frac{\partial^{2}}{\partial x_{j}^{2}}+\frac{\partial^{2}}{\partial \theta^{2}}
$$

- $V_{m}(x)=\frac{1}{2 \pi} \int_{0}^{2 \pi} V(x, \theta) e^{-i m \theta} d \theta$ for $m \in \mathbb{Z}$.
- $V^{\#}=V^{\#}(x, \theta)=V(x, \theta)-V_{0}(x)$.
- $B_{l}(\rho)$ and $D_{l}\left(\lambda_{0}, \rho\right)$ are open sets in $\hat{Z}$, defined in Sections 1 and 5, respectively.
- $R_{V}$ is the (scattering) resolvent of $-\Delta+V$ on $X$; see Section 3B.
- $R_{V_{0}, 0}$ is the (scattering) resolvent of $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$; see Section 3 .
- $m_{V}\left(\zeta_{0}\right)$ is the multiplicity of $\zeta_{0} \in \hat{Z}$ as a pole of $R_{V}$; see (4-1).
- $m_{V_{0}, 0}\left(\lambda_{0}\right)$ is the multiplicity of $\lambda_{0} \in \mathbb{C}$ as a pole of $R_{V_{0}, 0}$; see (4-2).
- $\zeta_{l}:\{z \in \mathbb{C}:|z|<\sqrt{2 l-1}\} \rightarrow B_{l}(\sqrt{2 l-1}) \subset \hat{Z}$ is the (local) inverse of

$$
B_{l}(\sqrt{2 l-1}) \ni \zeta \mapsto \tau_{l}(\zeta) \in\{z \in \mathbb{C}:|z|<\sqrt{2 l-1}\} \subset \mathbb{C} .
$$

## 3. Odd-dimensional Euclidean scattering and continuation of the resolvent

We begin by fixing notation and recalling some well-known facts from Euclidean scattering theory. We then use these to give a self-contained proof that the resolvent of $-\Delta+V$ on $X$ has a meromorphic continuation to $\hat{Z}$.

3A. The Euclidean resolvent. Let $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right), d$ odd, and set

$$
R_{V_{0}, 0}(\lambda)=\left(-\Delta_{0}+V_{0}-\lambda^{2}\right)^{-1}: L^{2}\left(\mathbb{R}^{d}\right) \rightarrow L^{2}\left(\mathbb{R}^{d}\right)
$$

when $\operatorname{Im} \lambda>0$. The 0 in the second place in the subscript is to help us think of this as a model operator, as we shall see. We shall later use the explicit expression for the resolvent as an integral when $d=1$, $f \in L^{2}(\mathbb{R})$, and $\operatorname{Im} \lambda>0$ :

$$
\begin{equation*}
\left(R_{0,0}(\lambda) f\right)(x)=\frac{i}{2 \lambda} \int e^{i \lambda\left|x-x^{\prime}\right|} f\left(x^{\prime}\right) d x^{\prime} \quad \text { for } d=1 \tag{3-1}
\end{equation*}
$$

From this we can see immediately that if $\chi \in L_{c}^{\infty}(\mathbb{R})$, then $\chi R_{0,0}(\lambda) \chi$ has a meromorphic continuation to $\mathbb{C} \backslash\{0\}$. The same is true when $d \geq 3$ is odd: if $\chi \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, then $\chi R_{0,0}(\lambda) \chi$ has an analytic continuation to the complex plane, see [Dyatlov and Zworski 2019, Theorem 3.3]. In higher dimensions, the Schwartz kernel is given in terms of a Hankel function. It is well known, see [Dyatlov and Zworski 2019, Theorem 3.8], that if $V_{0}, \chi \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, then $\chi R_{V_{0}, 0}(\lambda) \chi$ has a meromorphic continuation to the complex plane. Alternatively, restricting the domain and enlarging the range, $R_{V_{0}, 0}(\lambda): L_{c}^{2}\left(\mathbb{R}^{d}\right) \rightarrow H_{\mathrm{loc}}^{2}\left(\mathbb{R}^{d}\right)$ has a meromorphic continuation to $\mathbb{C}$.

The following lemma is well known, but we include it for completeness, as it is crucial for our arguments.

Lemma 3.1. Let $V_{0}, \chi \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. Then there are constants $C_{0}, C_{1}>0$ so that $\chi R_{V_{0}, 0}(\lambda) \chi$ is analytic in $\left\{\lambda \in \mathbb{C}:|\operatorname{Re} \lambda|>C_{0}, \operatorname{Im} \lambda>-C_{1} \log (1+|\operatorname{Re} \lambda|)\right\}$. Moreover, in this region $\left\|\chi R_{V_{0}, 0}(\lambda) \chi\right\|=O\left(|\lambda|^{-1}\right)$. Proof. Without loss of generality, we may assume $\chi V_{0}=V_{0}$. Then

$$
\chi R_{V_{0}, 0}(\lambda) \chi=\chi R_{0,0}(\lambda) \chi\left(I+V_{0} R_{0,0}(\lambda) \chi\right)^{-1} .
$$

Since from (3-1) when $d=1$ or [Dyatlov and Zworski 2019, Theorem 3.1] when $d \geq 3$, there is a $C>0$ so that

$$
\left\|V R_{0,0}(\lambda) \chi\right\| \leq C e^{C(\operatorname{Im} \lambda)_{-}} /|\lambda|
$$

where $(\operatorname{Im} \lambda)_{-}=\max (0,-\operatorname{Im} \lambda)$; the lemma follows immediately.
3B. The resolvent of $-\Delta+V$ on $X$ and the Riemann surface $\hat{Z}$. Recall that when $d=1, X$ is a manifold with infinite cylindrical ends. For a manifold with infinite cylindrical ends, the space to which the resolvent of a Schrödinger operator continues is determined by the distinct eigenvalues of the Laplacian on the cross-section of the end(s). Here that means $\left\{j^{2}\right\}_{j \in \mathbb{N}_{0}}$, since this is the set of (distinct) eigenvalues of $-d^{2} / d \theta^{2}$ on $\mathbb{S}^{1}$. As we show below, the resolvent for $-\Delta+V$ on $\mathbb{R}^{d} \times \mathbb{S}^{1}$ has a meromorphic continuation to the same space as that of the resolvent of $-\Delta+V$ on $\mathbb{R} \times \mathbb{S}^{1}$, provided $d$ is odd.

For $j \in \mathbb{N}_{0}$ and $\zeta \in \mathbb{C}, \operatorname{Im} \zeta>0$, set

$$
\tau_{j}(\zeta) \stackrel{\text { def }}{=}\left(\zeta^{2}-j^{2}\right)^{1 / 2}
$$

with $\operatorname{Im} \tau_{j}(\zeta)>0$. Set $\tau_{-j}(\zeta)=\tau_{j}(\zeta)$ if $j \in \mathbb{N}$.
The Riemann surface $\hat{Z}$ is defined to be the minimal Riemann surface on which, for each $j \in \mathbb{N}_{0}, \tau_{j}$ is a single-valued analytic function on $\hat{Z}$. We briefly describe its construction and some of its properties. Note that $\tau_{0}(\zeta)=\zeta$ for $\zeta$ in the upper half-plane, and this has, of course, an analytic continuation to $\mathbb{C}$. Now $\tau_{1}(\zeta)=\tau_{-1}(\zeta)$ is an analytic function of $\zeta \in \mathbb{C} \backslash((-\infty, 1] \cup[1, \infty))$, and there is a minimal Riemann surface $\hat{Z}_{1}$ so that $\tau_{1}$ extends analytically to $\hat{Z}_{1}$. This is a double cover of $\mathbb{C}$, ramified at the points $\pm 1$. This process can be repeated for each $j \in \mathbb{N}$, resulting in a minimal Riemann surface $\hat{Z}$ on which $\tau_{j}$ is analytic for each $j \in \mathbb{N}_{0}$. We define a projection $p: \hat{Z} \rightarrow \mathbb{C}$ as follows. For $\zeta$ in the physical space, identified with the upper half-plane, $p(\zeta)=\zeta$, and $p$ is in general the analytic continuation of this function. Then $\hat{Z}$ has infinitely many ramification points which project under $p$ to $j \in \mathbb{Z} \backslash\{0\}$. We call the set $\left\{\zeta \in \hat{Z}: \operatorname{Im} \tau_{j}(\zeta)>0\right.$ for all $\left.j \in \mathbb{N}_{0}\right\}$ the physical space, or physical region. For further discussion of this Riemann surface; see [Melrose 1993, Section 6.6].

We shall say that a point $\zeta_{0} \in \hat{Z}$ corresponds to a threshold if $\tau_{0}\left(\zeta_{0}\right) \in \mathbb{Z}$. Note that with this definition, all the ramification points of $\hat{Z}$ correspond to thresholds. In addition, the set of points corresponding to thresholds includes those points projecting to 0 . These might naturally also be considered ramification points of $\hat{Z}$, as in some sense by choosing $\zeta^{2}$ to originally be our spectral parameter we have already made the cuts corresponding to the zero threshold.

In order to separate variables below, we introduce the orthogonal projections $\mathcal{P}_{k}: L^{2}(X) \rightarrow L^{2}(X)$ defined for $k \in \mathbb{Z}$ by

$$
\begin{aligned}
& \left(\mathcal{P}_{k} f\right)(x, \theta)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f\left(x, \theta^{\prime}\right)\left(e^{i k\left(\theta-\theta^{\prime}\right)}+e^{-i k\left(\theta-\theta^{\prime}\right)}\right) d \theta^{\prime} \quad \text { if } k \in \mathbb{N} \\
& \left(\mathcal{P}_{0} f\right)(x, \theta)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f\left(x, \theta^{\prime}\right) d \theta^{\prime}
\end{aligned}
$$

We shall use these throughout the paper.
Let $V \in L_{c}^{\infty}(X)$. For $\zeta \in \mathbb{C}$ with $\operatorname{Im} \zeta>0$, set $R_{V}(\zeta)=\left(-\Delta+V-\zeta^{2}\right)^{-1}$. Consider first the special case where $V \in L_{c}^{\infty}(X)$ is independent of $\theta$. Then $V=V_{0}$, and we can think of $V_{0}$ as an element of $L_{c}^{\infty}(X)$ or of $L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. In this special case we can separate variables to obtain

$$
\begin{equation*}
R_{V_{0}}(\zeta)=\sum_{k=0}^{\infty} R_{V_{0}, 0}\left(\tau_{k}(\zeta)\right) \mathcal{P}_{k} \tag{3-2}
\end{equation*}
$$

The explicit expression (3-2) for $R_{V_{0}}$ using separation of variables shows that if $\chi \in L_{c}^{\infty}(X)$, then $\chi R_{V_{0}} \chi$ and $R_{V_{0}}: L_{c}^{2}(X) \rightarrow H_{\mathrm{loc}}^{2}(X)$ have meromorphic continuations to $\hat{Z}$. In fact, the same is true for $\chi R_{V} \chi$ and $R_{V}$ for general $V \in L_{c}^{\infty}(X)$. This is well known, at least when $d=1$, see [Goldstein 1974; Guillopé 1989; Melrose 1993], though we sketch a proof below, valid for all odd $d$.

If $\zeta \in \mathbb{C}, \operatorname{Im} \zeta>0$, then

$$
\left(-\Delta+V-\zeta^{2}\right) R_{0}(\zeta)=I+V R_{0}(\zeta)
$$

Multiplying by a function $\chi \in L_{c}^{\infty}(X)$ with $\chi V=V$,

$$
\left(-\Delta+V-\zeta^{2}\right) R_{0}(\zeta) \chi=\chi\left(I+V R_{0}(\zeta) \chi\right)
$$

implying that

$$
\begin{equation*}
\chi R_{0}(\zeta) \chi=\chi R_{V}(\zeta) \chi\left(I+V R_{0}(\zeta) \chi\right) \tag{3-3}
\end{equation*}
$$

or

$$
\begin{equation*}
\chi R_{V}(\zeta) \chi=\chi R_{0}(\zeta) \chi\left(I+V R_{0}(\zeta) \chi\right)^{-1} \tag{3-4}
\end{equation*}
$$

Using $I-V R_{0}(\zeta) \chi\left(I+V R_{0}(\zeta) \chi\right)^{-1}=\left(I+V R_{0}(\zeta) \chi\right)^{-1}$ and (3-4) yields

$$
\begin{equation*}
\left(I+V R_{0}(\zeta) \chi\right)^{-1}=I-V R_{V}(\zeta) \chi \tag{3-5}
\end{equation*}
$$

compare [Dyatlov and Zworski 2019, (2.2.15)-(2.2.16)]. Likewise, writing

$$
\begin{equation*}
V^{\#} \stackrel{\text { def }}{=} V-V_{0} \tag{3-6}
\end{equation*}
$$

we find, making the additional hypothesis that $\chi V^{\#}=V^{\#}$,

$$
\begin{equation*}
\chi R_{V_{0}}(\zeta) \chi=\chi R_{V}(\zeta) \chi\left(I+V^{\#} R_{V_{0}}(\zeta) \chi\right) \quad \text { and } \quad\left(I+V^{\#} R_{V_{0}}(\zeta) \chi\right)^{-1}=I-V^{\#} R_{V}(\zeta) \chi \tag{3-7}
\end{equation*}
$$

Each of these is helpful. Since $V R_{0}(\zeta) \chi: L^{2}(X) \rightarrow L^{2}(X)$ is compact and has a meromorphic extension to $\hat{Z}$, and $I+V R_{0}(\zeta) \chi$ is invertible for $\zeta$ in the physical space with $\operatorname{Im} \zeta$ sufficiently large, meromorphic Fredholm theory ensures that $\left(I+V R_{0}(\zeta) \chi\right)^{-1}$ is a meromorphic operator-valued function on $\hat{Z}$, and each


Figure 2. On the left, $B_{l}(\rho)$ in the $\tau_{l}$-coordinate; on the right, a portion of $B_{l}(\rho)$ in the $w=\left(\tau_{0}(\zeta)\right)^{2}$-coordinate for larger context. In the $w$ diagram, $(-\Delta+V-w)^{-1}$ is bounded in the upper half-plane and the red dots on the horizontal axis indicate thresholds. The hatching denotes the portion of $B_{l}(\rho)$ in the physical region; the shaded region indicates the rest which is visible in the $w$ plane diagram.
of (3-3)-(3-5) and (3-7) holds on all of $\hat{Z}$. Moreover, writing $I+V R_{0}=\left(I+V R_{0}(I-\chi)\right)\left(I+V R_{0} \chi\right)$ and noting that $\left(I+V R_{0}(I-\chi)\right)^{-1}=I-V R_{0}(I-\chi)$, this shows that

$$
\left.R_{V}(\zeta)=R_{0}(\zeta)\left(I+V R_{0}(\zeta) \chi\right)\right)^{-1}\left(I-V R_{0}(\zeta)(I-\chi)\right): L_{c}^{2}(X) \rightarrow H_{\mathrm{loc}}^{2}(X)
$$

has a meromorphic continuation to $\hat{Z}$.
We note from (3-2) that $R_{V_{0}}$ is bounded on $L^{2}(X)$ when $\zeta$ is in the physical space and is away from a discrete set of poles (corresponding to eigenvalues). The same is true of $R_{V}$.

Throughout this paper we shall mainly work with subsets of $B_{l}(\sqrt{2 l-1}) \subset \hat{Z}$, for $l \in \mathbb{N}$. We recall $B_{l}(\rho)$ is defined to be the connected component of $\left\{\zeta \in \hat{Z}:\left|\tau_{l}(\zeta)\right|<\rho\right\}$ which has nonempty intersection with both the physical space and the portion of $\hat{Z}$ with $\operatorname{Re} \tau_{0}(\zeta)>0$. The choice of $\sqrt{2 l-1}$ in $B_{l}(\sqrt{2 l-1})$ is made because then (for $l \geq 1) B_{l}(\sqrt{2 l-1})$ contains only a single point of $\hat{Z}$ corresponding to a threshold, the one associated with the eigenvalue $l^{2}$ of $-d^{2} / d \theta^{2}$ on $\mathbb{S}^{1}$. If $\epsilon>0$, then $z=\tau_{l}(\zeta)$ gives the complex structure of $\left.\hat{Z}\right|_{B_{l}(\sqrt{2 l-1}-\epsilon)}$, and $B_{l}(\sqrt{2 l-1}-\epsilon)$ is naturally identified with a disk $B_{\mathbb{C}}(\sqrt{2 l-1}-\epsilon)$ of radius $\sqrt{2 l-1}-\epsilon$ in $\mathbb{C}$, centered at the origin. In this coordinate $z$, we have that $z=0$ corresponds to the threshold $l^{2}$ and the intersection of $B_{\mathbb{C}}(\sqrt{2 l-1}-\epsilon)$ with the first quadrant corresponds to a region in physical space, and so has $\operatorname{Im} \tau_{k}>0$ for all $k \in \mathbb{N}_{0}$. If $z$ lies in the intersection of $B_{\mathbb{C}}(\sqrt{2 l-1}-\epsilon)$ with the fourth quadrant, then $\operatorname{Im} \tau_{k}(\zeta(z))<0$ for $0 \leq k \leq l$ and $\operatorname{Im} \tau_{k}(\zeta(z))>0$ for $k>l$. On the other hand, if $z$ lies in the intersection of $B_{\mathbb{C}}(\sqrt{2 l-1}-\epsilon)$ with the second quadrant, then $\operatorname{Im} \tau_{k}(\zeta(z))<0$ for $0 \leq k \leq l-1$ and $\operatorname{Im} \tau_{k}(\zeta(z))>0$ for $k \geq l$. Figure 2 shows a schematic of $B_{l}(\rho)$ and, for context, the portion of $B_{l}(\rho)$ which is visible in the $w=\left(\tau_{0}(\zeta)\right)^{2}$ plane. We note that while we have used the spectral parameter $\zeta^{2}$ in the definition of $R_{V}(\zeta)$ to be consistent with the usual odd-dimensional Euclidean scattering resolvent, the diagram on the right in Figure 2 uses as spectral parameter $w=\left(\tau_{0}(\zeta)\right)^{2}$ to make a more easily digested diagram. To put the diagram in context, think of $(-\Delta+V-w)^{-1}$ as having meromorphic continuation from the upper half-plane to $\left\{w \in \mathbb{C} \backslash\left(\bigcup_{j=0}^{\infty}\left(j^{2}+i(-\infty, 0]\right)\right)\right\}$ (which can, of course, be identified with a subset of $\hat{Z}$ ).

On the open set $B_{l}(\sqrt{2 l-1}-\epsilon), z=\tau_{l}(\zeta)$ is a coordinate compatible with the complex structure of $\hat{Z}$. Thus it is natural to use $\tau_{l}$ as a local coordinate. We write

$$
\zeta_{l}:\{z \in \mathbb{C}:|z|<\sqrt{2 l-1}-\epsilon\} \rightarrow B_{l}(\sqrt{2 l-1}-\epsilon) \subset \hat{Z}
$$

as the function satisfying

$$
\zeta_{l}\left(\tau_{l}(\zeta)\right)=\zeta \quad \text { for all } \zeta \in B_{l}(\sqrt{2 l-1}-\epsilon)
$$

We note that if $\zeta \in B_{l}(\sqrt{2 l-1}-\epsilon)$, then $\operatorname{Re} \tau_{j}(\zeta)>0$ if $0 \leq j<l$, and $\operatorname{Im} \tau_{j}(\zeta)>0$ if $j>l$.
The next lemma follows easily from (3-2) and Lemma 3.1, but is fundamental to many of the results of this paper.

Lemma 3.2. Let $V_{0} \in L_{c}^{\infty}(\mathbb{R}), \alpha>0$, and $\chi \in L_{c}^{\infty}(X)$. Then for $l$ sufficiently large, uniformly for $\zeta \in B_{l}(\alpha \log l)$, we have $\left\|\chi\left(I-\mathcal{P}_{l}\right) R_{V_{0}}(\zeta) \chi\right\|=O\left(l^{-1 / 2}\right)$.

Proof. Set $\tau_{l}=z$ and $|z|<\alpha \log l$. Then using the identity

$$
\tau_{k}^{2}=\tau_{l}^{2}+l^{2}-k^{2}
$$

for $l$ sufficiently large, $\left|\tau_{k}\left(\zeta_{l}(z)\right)\right|>\sqrt{l}$ for $k \in \mathbb{N}_{0}, k \neq l$. Moreover, $\operatorname{Im} \tau_{k}\left(\zeta_{l}(z)\right)>0$ if $k>l$, and $\left|\operatorname{Im} \tau_{k}\left(\zeta_{l}(z)\right)\right|=O(1)$ if $k<l$. Then the lemma follows from Lemma 3.1 and the representation of $R_{V_{0}, 0}$ given by (3-2).

## 4. Multiplicities of poles and results of [Gohberg and Sigal 1971]

For an operator $A$ depending meromorphically on $\zeta \in \mathbb{C}$ or $\zeta \in \hat{Z}$, let $\Xi\left(A, \zeta_{0}\right)$ denote the principal part of the Laurent expansion of $A$ at $\zeta_{0}$. For $V \in L_{c}^{\infty}(X)$ and $\zeta_{0} \in \hat{Z}$, define

$$
\begin{equation*}
m_{V}\left(\zeta_{0}\right) \stackrel{\text { def }}{=} \operatorname{rank} \Xi\left(R_{V}, \zeta_{0}\right)\left(L_{c}^{2}(X)\right) \tag{4-1}
\end{equation*}
$$

Suppose $\chi \in L_{c}^{\infty}(X)$ satisfies $\chi V=V$ (and, if $V \equiv 0, \chi$ is nontrivial). Then it follows from an expansion of $R_{V}$ at its singularities as in [Dyatlov and Zworski 2019, Theorems 2.5, 2.7, 3.9, 3.17] and a unique continuation result, e.g., [Jerison and Kenig 1985, Remark 6.7], that $m_{V}\left(\zeta_{0}\right)=\operatorname{rank} \Xi\left(\chi R_{V} \chi, \zeta_{0}\right)$. Note that if $R_{V}$ is analytic at $\zeta_{0}$, then $m_{V}\left(\zeta_{0}\right)=0$.

If $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and $\lambda_{0} \in \mathbb{C}$ we define

$$
\begin{equation*}
m_{V_{0}, 0}\left(\lambda_{0}\right) \stackrel{\text { def }}{=} \operatorname{rank} \Xi\left(R_{V_{0}, 0}, \lambda_{0}\right)\left(L_{c}^{2}\left(\mathbb{R}^{d}\right)\right) \tag{4-2}
\end{equation*}
$$

Again, the second 0 in the subscript is meant to help us think of this as corresponding to a model. As for $m_{V}$, if $\chi \in L_{c}^{\infty}(\mathbb{R})$ satisfies $\chi V=V$ (and $\chi$ is nontrivial if $V_{0} \equiv 0$ ), then $m_{V_{0}, 0}\left(\lambda_{0}\right)=\operatorname{rank} \Xi\left(\chi R_{V_{0}, 0} \chi, \lambda_{0}\right)$.

We recall some definitions and results of [Gohberg and Sigal 1971], adapted to our setting.
Let $A$ be a bounded linear operator on a complex Hilbert space $\mathcal{H}$, depending meromorphically on $z \in \Omega \subset \mathbb{C}$, where $\Omega$ is a domain. Near a point $z_{0} \in \Omega$, we have $A(z)=\sum_{j=-n}^{\infty}\left(z-z_{0}\right)^{j} A_{j}$. If the operators $A_{-1}, \ldots, A_{-n}$ are finite rank, then we say $A$ is finitely meromorphic at $z_{0}$. If $A$ is finitely meromorphic at each $z_{0} \in \Omega$, then $A$ is finitely meromorphic on $\Omega$. Suppose that $A$ is a compact operator on $\mathcal{H}, A$ is
finitely meromorphic on $\Omega$, and $I+A\left(z_{1}\right)$ is invertible for some $z_{1} \in \Omega$. Then by the meromorphic Fredholm theorem, $(I+A(z))^{-1}$ is finitely meromorphic on $\Omega$.

Suppose $A$ is a finitely meromorphic operator on a domain $\Omega$, with $(I+A)^{-1}$ also finitely meromorphic on $\Omega$. Below we denote the derivative of $A$ with respect to $z$ by $\dot{A}$. Then for $z_{0} \in \Omega$, define

$$
M\left(I+A, z_{0}\right) \stackrel{\text { def }}{=} \frac{1}{2 \pi i} \operatorname{tr} \int_{\gamma_{z_{0}}} \dot{A}(z)(I+A(z))^{-1} d z
$$

where $\gamma_{z_{0}}$ is a positively oriented circle, centered at $z_{0}$ with radius $\epsilon$. Here we choose $\epsilon$ small enough that $\left\{\left|z-z_{0}\right| \leq \epsilon\right\} \subset \Omega$ and neither $A$ nor $(I+A)^{-1}$ has poles in the set $\left\{z: 0<\left|z-z_{0}\right| \leq \epsilon\right\}$.

Our definition of finitely meromorphic is local, so it makes sense on domains in $\hat{Z}$ as well, using a local coordinate compatible with the complex structure of $\hat{Z}$. Likewise, we can define $M\left(I+A, \zeta_{0}\right)$ for such operators. (This requires the choice of a circle small enough that it has in its interior at most one ramification point of $\hat{Z}$.)

We will say the linear operator $A$ on the Hilbert space $\mathcal{H}$ satisfies hypotheses (H1) on a domain $\Omega \subset \mathbb{C}$ if $A$ is a finitely meromorphic, compact operator defined on $\Omega$, and $I+A$ is invertible for at least one point in $\Omega$ and hence has a finitely meromorphic inverse in $\Omega$.

The following lemma is a direct consequence of [Gohberg and Sigal 1971, Proposition 5].
Lemma 4.1 [Gohberg and Sigal 1971, Proposition 5]. Suppose $A, B: \mathcal{H} \rightarrow \mathcal{H}$ satisfy hypotheses (H1), and suppose $B$ and $(I+B)^{-1}$ are analytic on $\Omega$. Then for $z_{0} \in \Omega$,

$$
M\left(I+A, z_{0}\right)=M\left((I+A)(I+B), z_{0}\right)
$$

Let $T: L^{2}(X) \rightarrow L^{2}(X)$ be a bounded linear operator. We shall repeatedly make use of the straightforward identities

$$
\begin{equation*}
I+T \mathcal{P}_{l}=\left(I+\mathcal{P}_{l} T \mathcal{P}_{l}\right)\left(I+\left(I-\mathcal{P}_{l}\right) T \mathcal{P}_{l}\right) \quad \text { and } \quad\left(I+\left(I-\mathcal{P}_{l}\right) T \mathcal{P}_{l}\right)^{-1}=I-\left(I-\mathcal{P}_{l}\right) T \mathcal{P}_{l} \tag{4-3}
\end{equation*}
$$

Lemma 4.2. Let $A: L^{2}(X) \rightarrow L^{2}(X)$ satisfy hypotheses $(\mathrm{H} 1)$ on a domain $\Omega$. Then for $z_{0} \in \Omega$,

$$
M\left(I+A \mathcal{P}_{l}, z_{0}\right)=M\left(I+\mathcal{P}_{l} A \mathcal{P}_{l}, z_{0}\right)
$$

Proof. Using (4-3) implies that

$$
\begin{equation*}
M\left(I+A \mathcal{P}_{l}, z_{0}\right)=\frac{1}{2 \pi i} \operatorname{tr} \int_{\gamma_{z_{0}}} \dot{A}(z) \mathcal{P}_{l}\left(I+A(z) \mathcal{P}_{l}\right)^{-1} d z=\frac{1}{2 \pi i} \operatorname{tr} \int_{\gamma_{z_{0}}} \dot{A}(z) \mathcal{P}_{l}\left(I+\mathcal{P}_{l} A(z) \mathcal{P}_{l}\right)^{-1} d z \tag{4-4}
\end{equation*}
$$

where $\gamma_{z_{0}}$ is a small circle centered at $z_{0}$ as in the definition of $M\left(I+A, z_{0}\right)$.
Because $\mathcal{P}_{l}$ is a projection, using the cyclicity of the $\operatorname{trace}, \operatorname{tr}\left(B \mathcal{P}_{l}\right)=\operatorname{tr}\left(\mathcal{P}_{l} B \mathcal{P}_{l}\right)$ for a trace class operator $B: L^{2}(X) \rightarrow L^{2}(X)$. Using this in (4-4) gives

$$
M\left(I+A \mathcal{P}_{l}, z_{0}\right)=\frac{1}{2 \pi i} \operatorname{tr} \int_{\gamma_{z_{0}}} \mathcal{P}_{l} \dot{A}(z) \mathcal{P}_{l}\left(I+\mathcal{P}_{l} A(z) \mathcal{P}_{l}\right)^{-1} d z=M\left(I+\mathcal{P}_{l} A \mathcal{P}_{l}, z_{0}\right)
$$

The following proposition is a variant of a well-known result in the study of resonances of Schrödinger operators on $\mathbb{R}^{d}$; compare, e.g., [Dyatlov and Zworski 2019, Theorem 3.15].

Proposition 4.3. Suppose $V \in L_{c}^{\infty}(X)$ is nontrivial, and let $\chi \in L_{c}^{\infty}(X)$ satisfy $\chi V=V$. Then the operator $R_{V}(\zeta)$ has a pole at $\zeta_{0} \in \hat{Z}$ if and only if the operator $I+V R_{0}(\zeta) \chi$ has nontrivial null space at $\zeta_{0}$. Moreover, if $\zeta_{0}$ does not correspond to a threshold, then

$$
m_{V}\left(\zeta_{0}\right)=M\left(I+V R_{0} \chi, \zeta_{0}\right)
$$

Proof. A proof follows by essentially the same method as [Dyatlov and Zworski 2019, Theorem 3.15].
We recall the notation $V^{\#}=V-V_{0}$. Another useful identity is the following.
Lemma 4.4. Let $\chi \in L_{c}^{\infty}(X)$ satisfy $\chi V=V$ and $\chi V_{0}=V_{0}$. Then for $\zeta_{0} \in \hat{Z}$ so that $\zeta_{0}$ does not correspond to a threshold, we have

$$
m_{V}\left(\zeta_{0}\right)=M\left(I+V^{\#} R_{V_{0}} \chi, \zeta_{0}\right)+m_{V_{0}}\left(\zeta_{0}\right)
$$

Proof. We first note that

$$
\begin{equation*}
I+V R_{0} \chi=\left(I+V^{\#} R_{0} \chi\left(I+V_{0} R_{0} \chi\right)^{-1}\right)\left(I+V_{0} R_{0} \chi\right)=\left(I+V^{\#} R_{V_{0}} \chi\right)\left(I+V_{0} R_{0} \chi\right) \tag{4-5}
\end{equation*}
$$

Thus using Proposition 4.3 and [Gohberg and Sigal 1971, Theorem 5.2] gives

$$
\begin{aligned}
m_{V}\left(\zeta_{0}\right)=M\left(I+V R_{0} \chi, \zeta_{0}\right) & =M\left(I+V^{\#} R_{V_{0}} \chi, \zeta_{0}\right)+M\left(I+V_{0} R_{0} \chi, \zeta_{0}\right) \\
& =M\left(I+V^{\#} R_{V_{0}} \chi, \zeta_{0}\right)+m_{V_{0}}\left(\zeta_{0}\right)
\end{aligned}
$$

Lemma 4.5. Suppose $V$, $\chi \in L_{c}^{\infty}(X)$, with $\chi V=V$, and $\chi$ is independent of $\theta$. Let $\alpha>0$. Then there is an $L>0$ so that for $l>L$

$$
M\left(I+V R_{0} \chi, \zeta_{0}\right)=M\left(I+\mathcal{P}_{l}\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V R_{0} \mathcal{P}_{l} \chi, \zeta_{0}\right)
$$

for any $\zeta_{0} \in B_{l}(\alpha \log l)$.
Proof. We begin by writing

$$
I+V R_{0} \chi=\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)\left(I+\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V R_{0} \mathcal{P}_{l} \chi\right)
$$

and noting that since by Lemma $3.2\left\|V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right\|=O\left(l^{-1 / 2}\right)$ uniformly on $B_{l}(\alpha \log l)$ there is an $L>0$ so that for $l>L,\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1}$ is analytic on $B_{l}(\alpha \log l)$. Thus for these $l$ by Lemma 4.1 $M\left(I+V R_{0} \chi, \zeta_{0}\right)=M\left(I+\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V R_{0} \mathcal{P}_{l} \chi, \zeta_{0}\right)$ for any $\zeta_{0} \in B_{l}(\alpha \log l)$. An application of Lemma 4.2 completes the proof.

Lemma 4.6. Let $V, \chi \in L_{c}^{\infty}(X)$, with $V$ satisfying (1-1), $\chi V=V$, and $\chi$ independent of $\theta$. Set $A_{l, V}=\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1}$ and $B_{l, V}=V R_{0} \mathcal{P}_{l} \chi$. Let $K \subset \mathbb{C}$ be a compact set such that $R_{V_{0}, 0}$ is analytic on $K$, and suppose $0 \notin K$ if $d=1$. Choose $\rho>0$ so that $K \subset\{\lambda \in \mathbb{C}:|\lambda|<\rho\}$, and set $K_{l}=\left\{\zeta \in B_{l}(\rho): \tau_{l}(\zeta) \in K\right\}$. Then for sufficiently large $l$,

$$
\begin{equation*}
\left\|\mathcal{P}_{l}\left(A_{l, V} B_{l, V}-A_{l, V_{0}} B_{l, V_{0}}\right)\right\|=O\left(l^{-\delta}\right) \tag{4-6}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|\left(I+\mathcal{P}_{l} A_{l, V_{0}} B_{l, V_{0}}\right)^{-1} \mathcal{P}_{l}\left(A_{l, V} B_{l, V}-A_{l, V_{0}} B_{l, V_{0}}\right)\right\|=O\left(l^{-\delta}\right) \tag{4-7}
\end{equation*}
$$

uniformly for $\zeta \in K_{l}$.

Proof. We write

$$
\begin{equation*}
\mathcal{P}_{l}\left(A_{l, V} B_{l, V}-A_{l, V_{0}} B_{l, V_{0}}\right)=\mathcal{P}_{l}\left(A_{l, V}-A_{l, V_{0}}\right) B_{l, V}+\mathcal{P}_{l} A_{l, V_{0}}\left(B_{l, V}-B_{l, V_{0}}\right) \tag{4-8}
\end{equation*}
$$

By Lemma 3.2, $\left\|A_{l, V}-I\right\|=O\left(l^{-1 / 2}\right)$ and $\left\|A_{l, V_{0}}-I\right\|=O\left(l^{-1 / 2}\right)$ uniformly on $B_{l}(\rho)$, so that the first term on the left-hand side is $O\left(l^{-1 / 2}\right)$. Moreover,

$$
\mathcal{P}_{l} A_{l, V_{0}}\left(B_{l, V}-B_{l, V_{0}}\right)=A_{l, V_{0}} \mathcal{P}_{l}\left(B_{l, V}-B_{l, V_{0}}\right)=A_{l, V_{0}} \mathcal{P}_{l} V^{\#} R_{0} \mathcal{P}_{l}
$$

and $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-\delta}\right)$ by our assumption on $V$. Hence the norm of the second term on the right-hand side of (4-8) is $O\left(l^{-\delta}\right)$. This proves (4-6).

On $K_{l}$,

$$
\begin{equation*}
I+\mathcal{P}_{l} A_{l, V_{0}} B_{l, V_{0}}=I+\mathcal{P}_{l} B_{l, V_{0}}+O\left(l^{-1 / 2}\right)=I+\mathcal{P}_{l} V_{0} R_{0} \chi+O\left(l^{-1 / 2}\right) \tag{4-9}
\end{equation*}
$$

But

$$
\left(I+\mathcal{P}_{l} V_{0} R_{0} \chi\right)^{-1}=I-\mathcal{P}_{l}+\left(I-V_{0} R_{V_{0}, 0}\left(\tau_{l}\right) \chi\right) \mathcal{P}_{l}=I-\mathcal{P}_{l}+T \mathcal{P}_{l}
$$

where $T$ is given by $T=\left(I+V_{0} R_{0,0}\left(\tau_{l}\right) \chi\right)^{-1}=I-V_{0} R_{V_{0}, 0}\left(\tau_{l}\right) \chi$. By our choice of $K$, we have that $T$ is uniformly bounded for $\tau_{l} \in K$ or for $\zeta \in K_{l}$, and hence $\left(I+\mathcal{P}_{l} V_{0} R_{0} \chi\right)^{-1}$ is bounded on $K_{l}$. Using (4-9), this shows $\left(I+\mathcal{P}_{l} A_{l, V_{0}} B_{l, V_{0}}\right)^{-1}$ is bounded on $K_{l}$, and thus, by (4-6), we get (4-7).

## 5. A resolvent estimate and localizing the resonances in the $L^{\infty}$ case: Proofs of Theorems 1.1, 1.2, and 1.3

In this section we prove Theorems 1.1-1.3 in the case of an $L^{\infty}$ potential $V$, providing a high-energy localization of the resonances in sets $B_{l}(\rho)$. We also prove Proposition 5.4 and Lemma 5.5, which show that the resolvent for the potential $V_{0}$ is, at high energies, a good approximation of the resolvent for the potential $V$ away from poles.

We shall use notation for a disk in the $\tau_{l}$-coordinate in $B_{l}(\rho)$. For $\lambda_{0} \in \mathbb{C}$ and $r_{0}>0$, set $\rho=\left|\lambda_{0}\right|+r_{0}+1$, and define, for $2 l>\rho^{2}+1, D_{l}\left(\lambda_{0}, r_{0}\right) \subset B_{l}(\rho) \subset \hat{Z}$ by

$$
D_{l}\left(\lambda_{0}, r_{0}\right) \stackrel{\text { def }}{=}\left\{\zeta \in B_{l}(\rho):\left|\tau_{l}(\zeta)-\lambda_{0}\right|<r_{0}\right\}
$$

A preliminary step in the proof of Theorem 1.1 is the following proposition, which provides an initial localization of the resonances.
Proposition 5.1. Let $V \in L_{c}^{\infty}(X)$ satisfy (1-1). Suppose $\lambda_{0} \in \mathbb{C}, \lambda_{0} \neq 0$ is a resonance of $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$, of multiplicity $m_{V_{0}, 0}\left(\lambda_{0}\right)$. Then there are $L, \epsilon>0$ so that

$$
\sum_{\substack{\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right) \\ m_{V}(\zeta)>0}} m_{V}(\zeta)=2 m_{V_{0}, 0}\left(\lambda_{0}\right)
$$

when $l>L$.
Proof. Choose $\epsilon>0$ so that $R_{V_{0}, 0}(\lambda)$ is analytic on $0<\left|\lambda-\lambda_{0}\right| \leq \epsilon$ and $\epsilon<\left|\lambda_{0}\right|$. By our expression (3-2) for $R_{V_{0}}$, using separation of variables and Lemmas 3.1 and 3.2, $m_{V_{0}}\left(\zeta_{l}\left(\lambda_{0}\right)\right)=2 m_{V_{0}, 0}\left(\lambda_{0}\right)$ for $l$ sufficiently large. The factor of 2 on the right comes from the fact that the range of $\mathcal{P}_{l}$ (as an operator on $L^{2}\left(\mathbb{S}^{1}\right)$ ) has
rank 2 for $l>0$. Choose $\chi \in L_{c}^{\infty}(X)$ independent of $\theta$ so that $\chi V=V$. From Proposition 4.3 and our choice of $\epsilon$, for $l$ sufficiently large,

$$
m_{V_{0}}\left(\zeta_{l}\left(\lambda_{0}\right)\right)=M\left(I+V_{0} R_{0} \chi, \zeta_{l}\left(\lambda_{0}\right)\right)=\sum_{\substack{\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right) \\ M\left(I+V_{0} R_{0} \chi, \zeta\right) \neq 0}} M\left(I+V_{0} R_{0} \chi, \zeta\right)
$$

Lemma 4.5 implies that if $W=V_{0}$ or $W=V$,

$$
\begin{equation*}
M\left(I+W R_{0} \chi, \zeta^{\prime}\right)=M\left(I+\mathcal{P}_{l}\left(I+W R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} W R_{0} \mathcal{P}_{l} \chi, \zeta^{\prime}\right) \quad \text { for } \zeta^{\prime} \in D_{l}\left(\lambda_{0}, \epsilon\right) \tag{5-1}
\end{equation*}
$$

if $l$ is sufficiently large.
By Lemma 4.6 and an operator Rouché theorem [Gohberg and Sigal 1971, Theorem 2.2], for $l$ sufficiently large,

$$
\begin{gather*}
\sum_{\substack{\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right) \\
M\left(I+\mathcal{P}_{l}\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V R_{0} \mathcal{P}_{l} \chi, \zeta\right) \neq 0}} M\left(I+\mathcal{P}_{l}\left(I+V R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V R_{0} \mathcal{P}_{l} \chi, \zeta\right) \\
=\sum_{\substack{\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right) \\
M\left(I+\mathcal{P}_{l}\left(I+V_{0} R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V_{0} R_{0} \mathcal{P}_{l} \chi, \zeta\right) \neq 0}} M\left(I+\mathcal{P}_{l}\left(I+V_{0} R_{0}\left(I-\mathcal{P}_{l}\right) \chi\right)^{-1} V_{0} R_{0} \mathcal{P}_{l} \chi, \zeta\right)
\end{gather*}
$$

Combining (5-1) (with $W=V$ and with $W=V_{0}$ ), (5-2), and another application of Proposition 4.3, this time with $V$, proves the proposition.

5A. Proofs of Theorems 1.1 and 1.2 for $V \in L_{\boldsymbol{c}}^{\infty}(X)$. Theorem 1.1 follows from combining the result of Theorem 1.2 for $L^{\infty}$ potentials and Proposition 5.1. In this section we prove Theorem 1.2 for $L^{\infty}$ potentials $V$.

Recall by the definition of $\Xi\left(R_{V_{0}, 0}, \lambda_{0}\right)$, if $\lambda_{0} \in \mathbb{C}$ is a pole of $R_{V_{0}, 0}$, then $R_{V_{0}, 0}(\lambda)-\Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right)$ is analytic at $\lambda_{0}$. Define

$$
\begin{equation*}
R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right) \stackrel{\text { def }}{=} R_{V_{0}}(\zeta)-\Xi\left(R_{V_{0}}, \zeta_{l}\left(\lambda_{0}\right)\right) \tag{5-3}
\end{equation*}
$$

For $l$ sufficiently large, by (3-2) and Lemma 3.2

$$
\begin{equation*}
R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)=R_{V_{0}}(\zeta)-\left.\Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \tag{5-4}
\end{equation*}
$$

Note that if $R_{V_{0}}$ is analytic at $\zeta_{l}\left(\lambda_{0}\right)$, then $R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)=R_{V_{0}}(\zeta)$.
Lemma 5.2. Suppose $V, \chi \in L_{c}^{\infty}(X)$ and $V$ satisfies (1-1). Let $\lambda_{0} \in \mathbb{C}$ and $R_{V_{0}}^{\mathrm{reg}}=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$ be the operator defined in (5-3). If $R_{V_{0}, 0}(\lambda)$ is analytic for $0<\left|\lambda-\lambda_{0}\right| \leq \epsilon$, then for $l$ sufficiently large,

$$
V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi=V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right) \chi
$$

is analytic on $\bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$, and as $l \rightarrow \infty$ the estimate $\left\|\chi R_{V_{0}}^{\mathrm{reg}}(\zeta) V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right\|=O\left(l^{-\delta}\right)$ holds uniformly for $\zeta \in \bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$.

Proof. Without loss of generality we can assume $\chi$ is independent of $\theta$ and $\chi V=V$. By (3-2) and Lemma 3.2, for $l$ sufficiently large, $R_{V_{0}}^{\text {reg }}(\zeta)$ is analytic and bounded in $\bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$. We write

$$
\begin{align*}
& \chi R_{V_{0}}^{\mathrm{reg}}(\zeta) V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi \\
&=\chi R_{V_{0}}^{\mathrm{reg}} \chi\left(I-\mathcal{P}_{l}\right) V^{\#} R_{V_{0}}^{\mathrm{reg}} \chi+\chi R_{V_{0}}^{\mathrm{reg}} \chi \mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}} \chi\left(I-\mathcal{P}_{l}\right)+\chi R_{V_{0}}^{\mathrm{reg}} \chi \mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}} \chi \mathcal{P}_{l} . \tag{5-5}
\end{align*}
$$

Now for $\zeta \in \bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$ and $l$ sufficiently large, $\left\|\chi R_{V_{0}}^{\mathrm{reg}} \chi\left(I-\mathcal{P}_{l}\right)\right\|=O\left(l^{-1 / 2}\right)$ uniformly in $\bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$. Since $\left\|V_{m}\right\|=O\left(|m|^{-\delta}\right)$ we have $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-\delta}\right)$, and so

$$
\left\|\mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}} \chi \mathcal{P}_{l}\right\|=\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l} R_{V_{0}}^{\mathrm{reg}} \chi \mathcal{P}_{l}\right\|=O\left(l^{-\delta}\right)
$$

A related lemma which we also need is the following.
Lemma 5.3. Let $V, \chi \in L_{c}^{\infty}(X)$ with $V$ satisfying (1-1). Let $K \subset \mathbb{C}$ be a compact set on which $R_{V_{0}, 0}$ is analytic and suppose $K \subset\{\lambda \in \mathbb{C}:|\lambda|<\rho\}$. Set $K_{l} \stackrel{\text { def }}{=}\left\{\zeta \in B_{l}(\rho): \tau_{l}(\zeta) \in K\right\} \subset \hat{Z}$. Then for $l$ sufficiently large, $\left\|\chi R_{V_{0}} V^{\#} R_{V_{0}} \chi\right\|=O\left(l^{-\delta}\right)$ uniformly on $K_{l}$.

Proof. This lemma can be proved by mimicking the proof of Lemma 5.2. Alternatively, it can be proved by covering $K_{l}$ with a finite number of neighborhoods on which Lemma 5.2 holds.

Proof of Theorem 1.2 for $V \in L_{c}^{\infty}(X)$. We shall use the identities (3-7). Thus poles of $R_{V}$ in $B_{l}(\rho)$ are the values of $\zeta \in B_{l}(\rho)$ such that $I+V^{\#} R_{V_{0}}(\zeta) \chi$ is not invertible. Here $\chi \in C_{c}^{\infty}(X)$ satisfies $\chi V=V$ and is independent of $\theta$.
(1) For each $\lambda_{j} \in \Lambda_{\rho}, \lambda_{j} \neq 0$, let $\epsilon_{j}>0$ be as guaranteed by Proposition 5.1, so that there are exactly $2 m_{V_{0}, 0}\left(\lambda_{0}\right)$ resonances (counted with multiplicity) of $-\Delta+V$ in $D_{l}\left(\lambda_{j}, \epsilon_{j}\right)$ for $l$ sufficiently large. Set

$$
\begin{aligned}
K & =\left\{\lambda \in \mathbb{C}: \epsilon^{\prime} \leq|\lambda| \leq \rho \text { and }\left|\lambda-\lambda_{j}\right| \geq \epsilon_{j} \text { for all } \lambda_{j} \in \Lambda_{\rho}\right\} \\
K_{l} & =\left\{\zeta \in B_{l}(\rho+1): \tau_{l}(\zeta) \in K\right\}=\bar{B}_{l}(\rho) \backslash\left(D_{l}\left(0, \epsilon^{\prime}\right) \bigcup_{\lambda_{j} \in \Lambda_{\rho}} D_{l}\left(\lambda_{j}, \epsilon_{j}\right)\right)
\end{aligned}
$$

By an application of Lemma 5.3, for $l$ sufficiently large, $I+V^{\#} R_{V_{0}}(\zeta) \chi$ is invertible by its Neumann series on $K_{l}$. Thus by (3-7) $R_{V}$ has no poles on $K_{l}$ for $l$ sufficiently large.
(2) Now we work on $D_{l}\left(\lambda_{j}, \epsilon_{j}\right)$ and set $R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; l, \lambda_{j}\right)$, so that

$$
R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}(\zeta)-\left.\Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{j}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}
$$

for $l$ sufficiently large. By our choice of $\epsilon_{j}$ this is analytic on $\bar{D}_{l}\left(\lambda_{j}, \epsilon_{j}\right)$ for large enough $l$. Then by Lemma $5.2 I+V^{\#} R_{V_{0}}^{\text {reg }}(\zeta) \chi$ is invertible in $D_{l}\left(\lambda_{j}, \epsilon_{j}\right)$, with

$$
\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1}=I-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi+O_{L^{2}(X) \rightarrow L^{2}(X)}\left(l^{-\delta}\right)
$$

for $\zeta \in \bar{D}_{l}\left(\lambda_{j}, \epsilon_{j}\right)$. Thus on $\bar{D}_{l}\left(\lambda_{j}, \epsilon_{j}\right)$,

$$
\begin{equation*}
I+V^{\#} R_{V_{0}} \chi=\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)\left(I+\left.\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi\right) \tag{5-6}
\end{equation*}
$$

By (5-6) and (4-3), $I+V^{\#} R_{V_{0}} \chi$ is invertible at a point $\zeta \in D_{l}\left(\lambda_{j}, \epsilon_{j}\right)$ if and only if

$$
I+\left.\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi
$$

is invertible at $\zeta$. There is a $C_{j}$ so that $\left\|\chi \Xi\left(R_{V_{0}, 0}, \lambda_{j}\right) \chi\right\| \leq C_{j}\left|\lambda-\lambda_{j}\right|^{-m_{V_{0}, 0}\left(\lambda_{j}\right)}$ on $\left\{\lambda \in \mathbb{C}:\left|\lambda-\lambda_{j}\right| \leq \epsilon_{j}\right\}$; see [Dyatlov and Zworski 2019, Theorems 2.5 and 3.9]. Thus on $D_{l}\left(\lambda_{j}, \epsilon_{j}\right)$, using Lemma 5.2,

$$
\begin{aligned}
\| \mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right. & )\left.^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}\left(\lambda, \lambda_{0}\right)\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi \| \\
& =\left\|\left.\sum_{j=0}^{\infty} \mathcal{P}_{l}\left(-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\right)^{j} V^{\#} \Xi\left(R_{V_{0}, 0}\left(\lambda, \lambda_{0}\right)\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi\right\| \\
& \leq\left\|\left.\mathcal{P}_{l}\left(I-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\right) V^{\#} \Xi\left(R_{V_{0}, 0}\left(\lambda, \lambda_{0}\right)\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi\right\|+C_{j}^{\prime} l^{-\delta}\left|\tau_{l}(\zeta)-\lambda_{j}\right|^{-m_{V_{0}, 0}\left(\lambda_{j}\right)}
\end{aligned}
$$

Now we use Lemma 3.2, $\left\|V_{m}\right\|_{L^{\infty}}=O\left(|m|^{-\delta}\right)$, and the fact that $\mathcal{P}_{l}$ commutes with $R_{V_{0}, 0}$ so that

$$
\left\|\mathcal{P}_{l}\left(I-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\right) V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-\delta}\right)
$$

on $\bar{D}_{l}\left(\lambda_{j}, \epsilon_{j}\right)$. Thus there is a (new) $C_{j}^{\prime}$ so that

$$
\left\|\left.\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}\left(\lambda, \lambda_{0}\right)\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi\right\| \leq C_{j}^{\prime} l^{-\delta}\left|\tau_{l}(\zeta)-\lambda_{j}\right|^{-m_{V_{0}, 0}\left(\lambda_{j}\right)}
$$

on $\bar{D}_{l}\left(\lambda_{j}, \epsilon_{j}\right)$. Therefore

$$
I+\left.\mathcal{P}_{l}\left(I+R_{V_{0}}^{\mathrm{reg}}(\zeta)\right)^{-1} \Xi\left(R_{V_{0}, 0}\left(\lambda, \lambda_{0}\right)\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l} \chi
$$

is invertible in this region if $\left|\tau_{l}(\zeta)-\lambda_{j}\right| \geq C_{j} l^{-\delta / m_{V_{0}, 0}\left(\lambda_{j}\right)}$, where we can take $C_{j}=\left(2 C_{j}^{\prime}\right)^{1 / m_{V_{0}, 0}\left(\lambda_{j}\right)}$. Taking $\widetilde{C}=\max _{\lambda_{j} \in \Lambda_{\rho}} C_{j}$ finishes the proof of Theorem 1.2 away from $\tau_{l}=0$.
(3) If $R_{V_{0}, 0}(\lambda)$ does not have a pole at the origin, then there is a $\delta>0$ so that for $l$ sufficiently large, $R_{V_{0}}(\zeta)$ is analytic in $\bar{B}_{l}(\delta)$. Thus by Lemma 5.3, for $l$ sufficiently large, $R_{V}(\zeta)$ is analytic in $\bar{B}_{l}(\delta)$.

5B. Approximating the resolvent $\boldsymbol{R}_{V}$. In a sense made precise below in Proposition 5.4 and Lemma 5.5, at high energies $R_{V_{0}}$ approximates $R_{V}$ well away from resonances. The first result is useful for neighborhoods of thresholds.

Proposition 5.4. Let $V, \chi \in L_{c}^{\infty}(X)$, with $V$ satisfying (1-1). Let $K \subset \mathbb{C}$ be a compact set on which $R_{V_{0}, 0}$ is analytic and suppose $K \subset\{\lambda \in \mathbb{C}:|\lambda|<\rho\}$. Define $K_{l} \stackrel{\text { def }}{=}\left\{\zeta \in B_{l}(\rho): \tau_{l}(\zeta) \in K\right\} \subset \hat{Z}$. Then for $l$ sufficiently large, $R_{V}$ is analytic on $K_{l}$. Moreover, if $\chi \in L_{c}^{\infty}(X)$, then $\left\|\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi\right\|=O\left(l^{-\delta}\right)$ uniformly for $\zeta \in K_{l}$.
Proof. Without loss of generality we may assume $\chi$ is independent of $\theta$ and satisfies $\chi V=V$. Then $\chi R_{V_{0}} \chi=\chi R_{V} \chi\left(I+V^{\#} R_{V_{0}} \chi\right)$. Since by Lemma $5.3\left\|\left(V^{\#} R_{V_{0}} \chi\right)^{2}\right\| \leq \frac{1}{2}$ on $K_{l}$ for $l$ sufficiently large, $I+V^{\#} R_{V_{0}} \chi$ is invertible as $\left(I+V^{\#} R_{V_{0}} \chi\right)^{-1}=\sum_{j=0}^{\infty}\left(-V^{\#} R_{V_{0}} \chi\right)^{j}$, and thus $R_{V}$ is analytic on $K_{l}$. Moreover,

$$
\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi=\chi \sum_{j=1}^{\infty} R_{V_{0}}(\zeta)\left(-V^{\#} R_{V_{0}}(\zeta) \chi\right)^{j}
$$

By applying Lemma 5.3 twice, this becomes

$$
\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi=-\chi R_{V_{0}}(\zeta) V^{\#} R_{V_{0}}(\zeta) \chi+O_{L^{2} \rightarrow L^{2}}\left(l^{-\delta}\right)=O_{L^{2} \rightarrow L^{2}}\left(l^{-\delta}\right)
$$

A similar result with a similar proof is the following lemma. The points $\zeta \in \hat{Z}$ considered in this lemma lie on the boundary of the physical space, but are away from the thresholds.

Lemma 5.5. Let $V, \chi \in L_{c}^{\infty}(X)$, with $V$ satisfying (1-1). Then there are constants $M, L>0$ so that if $l>L, \zeta \in B_{l}(\sqrt{2 l-1}), \tau_{l}(\zeta) \in i[0, \infty)$, and $M<\frac{\tau_{l}(\zeta)}{i}<\sqrt{2 l-1}-\frac{M}{\sqrt{l}}$,

$$
\begin{equation*}
\text { then }\left\|\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi\right\|=O\left(l^{-\delta}\right) \tag{5-7}
\end{equation*}
$$

Likewise, there are constants $M_{1}, L_{1}>0$ so that

$$
\begin{array}{r}
\text { if } l>L_{1}, \zeta \in B_{l}(\sqrt{2 l-1}), \tau_{l}(\zeta) \in[0, \infty), \text { and } M_{1}<\tau_{l}(\zeta)<\sqrt{2 l-1}-\frac{M_{1}}{\sqrt{l}} \\
\text { then }\left\|\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi\right\|=O\left(l^{-\delta}\right) . \tag{5-8}
\end{array}
$$

Proof. This proof is very similar to the proof of Proposition 5.4. We outline the proof of the first statement only, as the proof of the second is analogous.

Without loss of generality, we may assume $\chi$ is independent of $\theta$ and satisfies $\chi V=V$.
We next note that if $\zeta \in B_{l}(\sqrt{2 l-1})$, then for $l>3$ either $\left|\tau_{l}(\zeta)\right|>\frac{1}{4} \sqrt{2 l-1}$ or $\left|\tau_{l-1}(\zeta)\right|>\frac{1}{4} \sqrt{2 l-1}$ or both are true. In either case, if $\tau_{l}(\zeta) \in i[0, \infty)$, then there is a $c_{0}>0$ so that $\left|\tau_{j}(\zeta)\right|>c_{0} l^{1 / 2}$ for $j \neq l, l-1$. Moreover, again with $\tau_{l}(\zeta) \in i[0, \infty), \operatorname{Im} \tau_{j}(\zeta)>0$ if $j>l$ and $\operatorname{Im} \tau_{j}(\zeta)=0$ if $0 \leq j<l$.

Suppose $\zeta \in B_{l}(\sqrt{2 l-1}), \tau_{l}(\zeta) \in i[0, \infty)$, and $\left|\tau_{l}(\zeta)\right|>\frac{1}{4} \sqrt{2 l-1}$. Then using Lemma 3.1 and (3-2) we see that

$$
\left\|\chi R_{V_{0}}(\zeta) \chi\left(I-\mathcal{P}_{l-1}\right)\right\|=O\left(l^{-1 / 2}\right)
$$

By Lemma 3.1 there is a $C>0$ so that if $\lambda \in \mathbb{R},|\lambda|>C$, then

$$
\left\|V^{\#}\right\|_{L^{\infty}}\left\|\chi R_{V_{0}, 0}(\lambda) \chi\right\| \leq \frac{1}{2}
$$

Choose $M>C+1$; then if $\tau_{l}(\zeta) \in i[0, \infty)$ with

$$
\frac{\tau_{l}(\zeta)}{i}<\sqrt{2 l-1}-\frac{M}{\sqrt{l}}
$$

for $l$ sufficiently large $\left|\tau_{l-1}(\zeta)\right|>C$. Now we restrict ourselves to $\tau_{l}(\zeta) \in i[0, \infty)$ with

$$
\frac{1}{4} \sqrt{2 l-1}<\frac{\tau_{l}(\zeta)}{i}<\sqrt{2 l-1}-\frac{M}{\sqrt{l}}
$$

Since $\left\|\mathcal{P}_{l-1} V^{\#} \mathcal{P}_{l-1}\right\|=O\left(l^{-\delta}\right)$ by our assumption on $\left\|V_{m}\right\|_{L^{\infty}}$,

$$
\left\|\chi R_{V_{0}}(\zeta) \mathcal{P}_{l-1} V^{\#} R_{V_{0}}(\zeta) \mathcal{P}_{l-1} \chi\right\|=O\left(l^{-\delta}\right)
$$

and we can follow the proof of Lemma 5.2 to show that $\left\|\chi R_{V_{0}}(\zeta) V^{\#} R_{V_{0}}(\zeta) \chi\right\|=O\left(l^{-\delta}\right)$. Then

$$
\begin{aligned}
\left\|\chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi\right\| & =\left\|\chi R_{V_{0}}(\zeta) \chi\left(\left(I+V^{\#} R_{V_{0}}(\zeta) \chi\right)^{-1}-I\right)\right\| \\
& =\left\|\chi R_{V_{0}}(\zeta) V^{\#} R_{V_{0}}(\zeta) \chi\right\|+O\left(l^{-\delta}\right)=O\left(l^{-\delta}\right)
\end{aligned}
$$

proving the lemma when $\tau_{l}(\zeta) \in i[0, \infty)$ with $\frac{1}{4} \sqrt{2 l-1}<\frac{1}{i} \tau_{l}(\zeta)<\sqrt{2 l-1}-M / \sqrt{l}$. A similar argument, singling out $\mathcal{P}_{l}$ rather than $\mathcal{P}_{l-1}$, handles the case when $\tau_{l}(\zeta) \in i[0, \infty)$ with $\frac{1}{4} \sqrt{2 l-1}<\left|\tau_{l-1}(\zeta)\right|$.

5C. Proof of Theorem 1.3. Theorem 1.3 concerns poles of $R_{V}$ arising as perturbations of threshold poles of $R_{V_{0}}(\zeta)$. Using separation of variables as in (3-2), these threshold poles, in turn, correspond to a pole of $R_{V_{0}, 0}(\lambda)$ at $\lambda=0$.

We begin with a lemma about poles of $R_{V_{0}}(\lambda)$ at the origin. This result is well known if $V_{0}$ is real-valued.

Lemma 5.6. Suppose $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, and near $\lambda=0$

$$
\begin{equation*}
R_{V_{0}, 0}(\lambda)=\sum_{k=1}^{k_{0}} \frac{1}{\lambda^{k}} A_{k}+A(\lambda) \tag{5-9}
\end{equation*}
$$

where $A$ is analytic in a neighborhood of the origin. Then $m_{V_{0}, 0}(0)=\max _{0 \leq t \leq 1} \operatorname{rank}\left(A_{1}+t A_{2}\right)$.
Since $A_{1}, A_{2}$ are finite-rank, the rank of $A_{1}+t A_{2}$ is equal to its maximum for all but a finite number of values of $t$ in $[0,1]$.

Proof. Using the expansion (5-9) and the identity $\left(-\Delta_{0}+V_{0}-\lambda^{2}\right) R_{V_{0}, 0}(\lambda)=I$ shows that for $k>0$, $\left(-\Delta_{0}+V_{0}\right) A_{k}=A_{k+2}$, where we use the convention $A_{k+2}=0$ if $k+2>k_{0}$. Just as in [Dyatlov and Zworski 2019, Theorem 2.5], one can use this and the fact that $-\Delta_{0}+V_{0}$ commutes with $R_{V_{0}, 0}$ to show that for $j \in \mathbb{N}$, $\operatorname{Ran}\left(A_{2 j}\right) \subset \operatorname{Ran}\left(A_{2}\right)$ and $\operatorname{Ran}\left(A_{2 j+1}\right) \subset \operatorname{Ran}\left(A_{1}\right)$. Here $\operatorname{Ran}\left(A_{k}\right)$ denotes the range of the operator $A_{k}$ on $L_{c}^{2}\left(\mathbb{R}^{d}\right)$. Since $m_{V_{0}, 0}(0)=\operatorname{dim}\left(\bigcup_{k=1}^{k_{0}} \operatorname{Ran}\left(A_{k}\right)\right)$, this shows $m_{V_{0}, 0}(0)=\operatorname{dim}\left(\operatorname{Ran} A_{1} \cup \operatorname{Ran} A_{2}\right)$. But

$$
\operatorname{dim}\left(\operatorname{Ran} A_{1} \cup \operatorname{Ran} A_{2}\right)=\max _{t \in[0,1]} \operatorname{dim} \operatorname{Ran}\left(A_{1}+t A_{2}\right)=\max _{t \in[0,1]} \operatorname{rank}\left(A_{1}+t A_{2}\right)
$$

proving the lemma.
Lemma 5.7. Let $V \in L_{c}^{\infty}(X)$ satisfy (1-1). Let $\epsilon>0$ be chosen so that $R_{V_{0}, 0}(\lambda)$ has no poles in $\{\lambda \in \mathbb{C}: 0<|\lambda|<2 \epsilon\}$, and let $\gamma_{l} \subset B_{l}(2 \epsilon) \subset \hat{Z}$ be the curve $\left\{\left|\tau_{l}\right|=\epsilon\right\}$ with positive orientation. Then for $t \in[0,1]$ and $l$ sufficiently large,

$$
\operatorname{rank} \int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) R_{V}(\zeta) d \tau_{l}(\zeta) \geq \operatorname{rank} \int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) R_{V_{0}}(\zeta) d \tau_{l}(\zeta)
$$

Proof. We assume $V^{\#}$ is nontrivial, since otherwise there is nothing to prove.
We first point out that if $R_{V_{0}, 0}(\lambda)=\sum_{k=1}^{k_{0}} \lambda^{-k} A_{k}+A(\lambda)$, with $A(\lambda)$ analytic near $\lambda=0$, then for $l$ sufficiently large

$$
\int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) R_{V_{0}}(\zeta) d \tau_{l}(\zeta)=\int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) R_{V_{0}, 0}\left(\tau_{l}(\zeta)\right) \mathcal{P}_{l} d \tau_{l}(\zeta)=2 \pi i\left(A_{1}+t A_{2}\right) \mathcal{P}_{l}
$$

Let $\chi \in L_{c}^{\infty}(X)$ satisfy $\chi V=V$, with $\chi$ independent of $\theta$. Using Proposition 5.4 , for $l$ sufficiently large,

$$
\left\|\int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) \chi\left(R_{V}(\zeta)-R_{V_{0}}(\zeta)\right) \chi d \tau_{l}(\zeta)\right\|=O\left(l^{-\delta}\right)
$$

Thus

$$
\left\|\int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) \chi R_{V}(\zeta) \chi d \tau_{l}(\zeta)-2 \pi i \chi\left(A_{1}+t A_{2}\right) \mathcal{P}_{l} \chi\right\|=O\left(l^{-\delta}\right)
$$

and this implies that for $l$ sufficiently large,

$$
\begin{equation*}
\operatorname{rank} \int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) \chi R_{V}(\zeta) \chi d \tau_{l}(\zeta) \geq 2 \operatorname{rank}\left(\chi\left(A_{1}+t A_{2}\right) \chi\right) \tag{5-10}
\end{equation*}
$$

But since $\left(-\Delta_{0}+V_{0}\right)^{k_{0}} A_{j}=0$ for $j=1$, 2, a unique continuation theorem, e.g., [Jerison and Kenig 1985], ensures that $\operatorname{rank}\left(A_{1}+t A_{2}\right)=\operatorname{rank}\left(\chi\left(A_{1}+t A_{2}\right) \chi\right)$, and similarly

$$
\operatorname{rank} \int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) \chi R_{V}(\zeta) \chi d \tau_{l}(\zeta)=\operatorname{rank} \int_{\gamma_{l}}\left(1+t \tau_{l}(\zeta)\right) R_{V}(\zeta) d \tau_{l}(\zeta)
$$

Lemma 5.8. Let $V_{0}, \chi \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, with $\chi V_{0}=V_{0}$. Suppose $R_{V_{0}}(\lambda)$ has a pole of order 1 at the origin. Then for $l$ sufficiently large, $2\left(m_{V_{0}, 0}(0)-m_{0,0}(0)\right)=M\left(I+V_{0} R_{0} \chi, \zeta_{l}(0)\right)$.

Proof. We note here that the requirement that $l$ is sufficiently large is to ensure that, using (3-2), any poles of $R_{V_{0}}$ at $\zeta_{l}(0)$ arise from poles of $R_{V_{0}}$ at the origin. Then via separation of variables it suffices to show that

$$
m_{V_{0}, 0}(0)-m_{0,0}(0)=M\left(I+V_{0} R_{0,0}(\lambda) \chi, 0\right)
$$

For $d=1$, then $m_{V_{0}, 0}(0)=1$ and if $V_{0}$ is real-valued, this follows immediately from [Dyatlov and Zworski $2019,(2.2 .31)]$. For complex-valued $V_{0}$, the proof is similar, if one uses the assumption that $R_{V_{0}}$ has a simple pole at the origin. When $d \geq 3$ is odd, the lemma follows as in the proof of [Dyatlov and Zworski 2019, Theorem 3.15]. In each case, the assumption that the pole is of order 1 is important.

Lemma 5.9. Let $V \in L_{c}^{\infty}(X)$ satisfy (1-1). Let $\epsilon>0$ be chosen so that $R_{V_{0}, 0}(\lambda)$ has no poles in $\{\lambda \in \mathbb{C}: 0<|\lambda|<2 \epsilon\}$. Suppose $R_{V_{0}, 0}(\lambda)$ has a pole of order 1 at the origin, with residue of rank $m_{V_{0}, 0}(0)$. Then for $l$ sufficiently large,

$$
\sum_{\substack{\zeta \in D_{l}(\epsilon) \\ m_{V}(\zeta) \neq 0}} m_{V}(\zeta) \leq 2 m_{V_{0}, 0}(0)
$$

Proof. Let $\chi \in L_{c}^{\infty}(X)$ be independent of $\theta$ and satisfy $\chi V=V$. We first claim that for any $\zeta_{0} \in \hat{Z}$,

$$
\begin{equation*}
m_{V}\left(\zeta_{0}\right) \leq M\left(I+V R_{0} \chi, \zeta_{0}\right)+m_{0}\left(\zeta_{0}\right) \tag{5-11}
\end{equation*}
$$

If $\zeta_{0}$ does not correspond to a threshold, then $m_{0}\left(\zeta_{0}\right)=0$ and this follows from the stronger Proposition 4.3. If $\zeta_{0}$ does correspond to a threshold, this follows from a simplified adaptation of the proof of [Dyatlov and Zworski 2019, Theorem 3.15].

Arguing as in the proof of Proposition 5.1, using Lemmas 4.5 and 4.6 and an operator Rouché theorem [Gohberg and Sigal 1971, Theorem 2.2], for $l$ sufficiently large,

$$
\begin{equation*}
\sum_{\substack{\zeta \in B_{l}(\epsilon) \\ M\left(I+V R_{0} \chi, \zeta\right) \neq 0}} M\left(I+V R_{0} \chi, \zeta\right)=\sum_{\substack{\zeta \in B_{l}(\epsilon) \\ M\left(I+V_{0} R_{0} \chi, \zeta\right) \neq 0}} M\left(I+V_{0} R_{0} \chi, \zeta\right)=M\left(I+V_{0} R_{0} \chi, \zeta_{l}(0)\right) . \tag{5-12}
\end{equation*}
$$

But by our assumptions and Lemma 5.8, $M\left(I+V_{0} R_{0} \chi, \zeta_{l}(0)\right)=2\left(m_{V_{0}, 0}(0)-m_{0,0}(0)\right)$ for $l$ sufficiently large. Using this, (5-12), and applying (5-11) completes the proof.

Proof of Theorem 1.3 under the assumption $\left\|V_{m}\right\|_{L^{\infty}}=O\left(|m|^{-\delta}\right)$. Let $\epsilon>0$ be as in the statement of Lemma 5.7. By applying Lemmas 5.6 and 5.7, we see that for $l$ sufficiently large,

$$
\sum_{\substack{\zeta \in B_{l}(\epsilon) \\ m_{V}(\zeta) \neq 0}} m_{V}(\zeta) \geq \sum_{\substack{\zeta \in B_{l}(\epsilon) \\ m_{V_{0}}(\zeta) \neq 0}} m_{V_{0}}(\zeta)=2 m_{V_{0}, 0}(0)
$$

Thus for $l$ sufficiently large $R_{V}$ has at least $2 m_{V_{0}, 0}(0)$ poles in $B_{l}(\epsilon)$. If $R_{V_{0}, 0}(\lambda)$ has a simple pole at the origin, then applying in addition Lemma 5.9 we see that $R_{V}$ has at exactly $2 m_{V_{0}, 0}(0)$ poles in $B_{l}(\epsilon)$.

To finish the proof of the theorem for the $L^{\infty}$ case we need to refine the estimate on the location of the resonances in $B_{l}(\epsilon)$. We do this by showing that there is a $C>0$ so that there are no resonances in $B_{l}(\epsilon) \backslash B_{l}\left(C l^{-\delta / r}\right)$ for $l$ sufficiently large. This follows almost exactly the proof of Theorem 1.2, point 2, with $\lambda_{j}$ replaced by 0 . The difference here is that the bound on the singular part of $\chi R_{V_{0}} \chi$ at the origin is given by $\left\|\chi \Xi\left(R_{V_{0}}, 0\right) \chi\right\| \leq C|\lambda|^{-r}$; that is, $m_{V_{0}, 0}\left(\lambda_{j}\right)$ is replaced by $r$ rather than $m_{V_{0}, 0}(0)$. Having made this minor adaptation, the remainder of the proof follows without change.

## 6. Resonance-free regions, poles of $R_{V}$ and $R_{\bar{V}}$, and the proofs of Corollary 1.4 and Theorem 1.9

Thus far we have focused on resonances in the sets $B_{l}(\rho)$, for $l$ large. In this section we justify this by showing that the high-energy resonances near the physical space which also have $\operatorname{Re} \tau_{0}(\zeta)>0$ lie in $B_{l}(\rho)$, for $\rho$ sufficiently large. We do this by showing the existence of large resonance-free regions in $B_{l}(\sqrt{2 l-1})$. We discuss $\hat{Z}$ further, focusing on the region near the physical space. We describe the relationship between the resolvents $R_{V}$ and $R_{\bar{V}}$, where $\bar{V}$ is the complex conjugate of $V$; see Lemma 6.2. This lemma shows that we can understand the poles of $R_{V}$ which are near the physical space and have $\operatorname{Re} \tau_{0}(\zeta)<0$ by understanding the poles of $R_{\bar{V}}$ which are near the physical space and have $\operatorname{Re} \tau_{0}(\zeta)>0$.

Lemma 6.1. Let $V \in L_{c}^{\infty}(X)$. Then for any $0<\gamma<1$ there are $M_{+}, c_{+}>0$ so that the region

$$
U_{l}^{+} \stackrel{\text { def }}{=}\left\{\zeta \in B_{l}(\sqrt{2 l-1}): M_{+}<\operatorname{Re}\left(\tau_{l}(\zeta)\right)<\gamma \sqrt{2 l}, \operatorname{Im} \tau_{l}(\zeta)>-c_{+} \log \operatorname{Re}\left(\tau_{l}(\zeta)\right)\right\}
$$

contains no poles of $R_{V}$ forl sufficiently large. Likewise, for any $\alpha>0$ and $0<\gamma<1$, there is a constant $M_{-}>0$ so that

$$
U_{l}^{-} \stackrel{\text { def }}{=}\left\{\zeta \in B_{l}(\sqrt{2 l-1}): M_{-}<\operatorname{Im}\left(\tau_{l}(\zeta)\right)<\gamma \sqrt{2 l}, \operatorname{Re} \tau_{l}(\zeta)>-\alpha\right\}
$$

contains no poles of $R_{V}$ for $l$ sufficiently large.
The region $U_{l}^{+}$is reminiscent of the logarithmic resonance-free regions familiar from potential scattering on $\mathbb{R}^{d}$. We note that there is substantial overlap between $U_{l}^{+}$and $U_{l+1}^{-}$.
Proof. Let $\chi \in L_{c}^{\infty}(X)$ be independent of $\theta$ and satisfy $\chi V=V$ and $0 \leq \chi \leq 1$. To prove the lemma, we use $\chi R_{V}(\zeta) \chi=\chi R_{0}(\zeta)\left(I+V R_{0}(\zeta) \chi\right)^{-1}$ and the representation (3-2) via separation of variables.

From (3-2) and the estimate $\left\|\chi R_{0,0}(\lambda) \chi\right\| \leq C e^{(C \operatorname{Im} \lambda)_{-}}-|\lambda|$, there are constants $C_{1}, C_{2}$ so that

$$
\left\|V R_{0}(\zeta) \chi\right\| \leq \sup _{j \in \mathbb{N}_{0}}\left(\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{j}(\zeta)\right)_{-}}}{\left|\tau_{j}(\zeta)\right|}\right)
$$

First consider $U_{l}^{+}$. Set $c_{+}=1 / C_{2}-\delta_{+}$, where $\delta_{+}>0, \delta_{+}<1 / C_{2}$, and take $M_{+}>\left(2 C_{1}\right)^{1 /\left(\delta_{+} C_{2}\right)}$. Then if $\zeta \in U_{l}^{+}$,

$$
\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{l}(\zeta)\right)-}}{\left|\tau_{l}(\zeta)\right|}<\frac{1}{2}
$$

If $j<l$ and $\zeta \in U_{l}^{+}$, then $\left|\tau_{j}(\zeta)\right| \geq\left|\tau_{l}(\zeta)\right|$ and a computation shows

$$
\frac{e^{C_{2}\left(\operatorname{Im} \tau_{j}(\zeta)\right)_{-}}}{\left|\tau_{j}(\zeta)\right|}<\frac{e^{C_{2}\left(\operatorname{Im} \tau_{l}(\zeta)\right)_{-}}}{\left|\tau_{l}(\zeta)\right|}
$$

On the other hand, for $j>l$, if $\zeta \in U_{l}^{+}$, then

$$
\operatorname{Re}\left(\tau_{j}(\zeta)\right)^{2} \leq \operatorname{Re}\left(\tau_{l+1}(\zeta)\right)^{2}=\left(\operatorname{Re} \tau_{l}(\zeta)\right)^{2}-2 l-\left(\operatorname{Im} \tau_{l}(\zeta)\right)^{2}-1 \leq-2 l\left(1-\gamma^{2}\right)
$$

Since $\operatorname{Im} \tau_{j}(\zeta)>0$ for $j>l$ and $\zeta \in B_{l}(\sqrt{2 l-1})$, this is enough to show that

$$
\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{j}(\zeta)\right)_{-}}}{\left|\tau_{j}(\zeta)\right|}<\frac{1}{2}
$$

for $\zeta \in U_{l}^{+}$and $l$ sufficiently large. Then $\left\|V R_{0}(\zeta) \chi\right\|<\frac{1}{2}$, and $I+V R_{0}(\zeta) \chi$ is invertible.
For $U_{l}^{-}$, choose $M_{-}>0$ so that $16\|V\|_{L^{\infty}}<M_{-}^{2}$. Then using (3-2) and $\left\|R_{0,0}(\lambda)\right\| \leq 1 /\left(\operatorname{dist}\left(\lambda^{2},[0, \infty)\right)\right)$ for $\operatorname{Im} \lambda>0$, for $\zeta \in U_{l}^{-}$,

$$
\left\|V R_{0}(\zeta) \chi \sum_{j \geq l} \mathcal{P}_{j}\right\| \leq\|V\|_{L^{\infty}} \sup _{j \geq l} \frac{1}{\left(\operatorname{dist} \tau_{j}^{2},[0, \infty)\right)} \leq \frac{8\|V\|_{L^{\infty}}}{M_{-}^{2}} \leq \frac{1}{2}
$$

Next we show that

$$
\left\|V R_{0}(\zeta) \sum_{0 \leq j<l} \mathcal{P}_{j} \chi\right\| \leq \frac{1}{2}
$$

in $U_{l}^{-}$for sufficiently large $l$. Using the orthogonality of the projections $\sum_{j \geq l} \mathcal{P}_{j}$ and $\sum_{0 \leq j<l} \mathcal{P}_{j}$ this will complete our proof that $I+V R_{0} \chi$ is invertible. Note that

$$
\tau_{l-1}^{2}=2 l-\left(\operatorname{Im} \tau_{l}\right)^{2}+\left(\operatorname{Re} \tau_{l}\right)^{2}-1+2 i \operatorname{Re}\left(\tau_{l}\right) \operatorname{Im}\left(\tau_{l}\right)
$$

Thus $\left|\tau_{l-1}\right| \geq \sqrt{\left(1-\gamma^{2}\right) 2 l}+O(1)$ and $-\operatorname{Im}\left(\tau_{l-1}\right) \leq 2 \alpha / \sqrt{1-\gamma^{2}}+O\left(l^{-1 / 2}\right)$, so for $l$ sufficiently large,

$$
\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{l-1}(\zeta)\right)_{-}}}{\left|\tau_{l-1}(\zeta)\right|}<\frac{1}{2}
$$

for $\zeta \in U_{l}^{-}$. But if $0 \leq j<l-1$ and $\zeta \in U_{l}^{-}$,

$$
\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{j}(\zeta)\right)_{-}}}{\left|\tau_{j}(\zeta)\right|}<\frac{C_{1} e^{C_{2}\left(\operatorname{Im} \tau_{l-1}(\zeta)\right)_{-}}}{\left|\tau_{l-1}(\zeta)\right|}
$$

This ensures that

$$
\left\|V R_{0}(\zeta) \sum_{0 \leq j<l} \mathcal{P}_{l} \chi\right\|<\frac{1}{2}
$$

so that $I+V R_{0}(\zeta) \chi$ is invertible on $U_{l}^{-}$for $l$ sufficiently large.

We remark that we have not made an effort to optimize the results of Lemma 6.1, as in this paper we are concentrating instead on regions near the thresholds, where, as we have seen, resonances can occur.

Before proving Corollary 1.4, we discuss $\hat{Z}$ and the boundary of the physical space a bit more. To motivate the discussion, consider the simpler case of the Schrödinger operator $-\Delta_{0}+V_{0}$ on $\mathbb{R}^{d}$, where we use $\lambda^{2}$ as the spectral parameter in defining the (scattering) resolvent. Thus, given a value $E>0$, there are two points, $\pm \sqrt{E}$ corresponding to the spectral parameter $E$ on the boundary of the physical space, with $R_{V_{0}, 0}( \pm \sqrt{E})=\left(-\Delta_{0}+V_{0}-(\sqrt{E} \pm i 0)\right)^{-1}$.

There is a similar phenomena in the case of $-\Delta+V$ on $\mathbb{R}^{d} \times \mathbb{S}^{1}$, but it is notationally harder to describe. Given $E>0$, let $\sqrt{E} \pm i 0 \in \hat{Z}$ be the points on $\hat{Z}$ with $R_{V}(\sqrt{E} \pm i 0)=(-\Delta+V-E \mp i 0)^{-1}$. Equivalently, we could define $\sqrt{E} \pm i 0$ to be the point in $\hat{Z}$ with $\tau_{j}(\sqrt{E} \pm i 0)= \pm \sqrt{E-j^{2}}$ if $j^{2} \leq E$, and $\tau_{j}(\sqrt{E} \pm i 0)=i \sqrt{j^{2}-E}$ if $j^{2}>E$. By our definition of $B_{l}(\rho)$, if $l_{E}=\lfloor\sqrt{E}\rfloor$ and $l_{E}>0$, then $\sqrt{E}+i 0 \in B_{l_{E}}\left(\sqrt{2 l_{E}-1}\right)$, but $\sqrt{E}-i 0 \notin B_{l_{E}}\left(\sqrt{2 l_{E}-1}\right)$. Thus there is some sense in which we have been studying only "half" of the boundary of the physical space. However, we shall see in Lemma 6.2 that this suffices for understanding the behavior of the resolvent, if we consider both the resolvent of $-\Delta+V$ and that of $-\Delta+\bar{V}$.

Thus, to fully cover points on the boundary of the physical space, we need to define another type of open set in $\hat{Z}$, analogous to $B_{l}(\rho)$. For $l \in \mathbb{N}$ and $\rho>0$, denote by $B_{l}^{ \pm}(\rho)$ the connected component of $\left\{\zeta \in \hat{Z}:\left|\tau_{l}(\zeta)\right|<\rho\right\}$ which intersects the physical space and includes a region with $\pm \operatorname{Re} \tau_{0}(\zeta)>0$. With the + sign, we get the set $B_{l}(\rho)$ defined in the introduction: $B_{l}^{+}(\rho)=B_{l}(\rho)$. If $l_{E}=\lfloor\sqrt{E}\rfloor$ and $\sqrt{E}-l_{E}<\rho$, then the point $\sqrt{E}-i 0$ corresponding to $E$ on the boundary of the physical space as defined above has $\sqrt{E}-i 0 \in B_{l_{E}}^{-}(\rho)$. Hence any point on the boundary of the physical space lies in

$$
B_{0}^{+}(1) \cup\left(\bigcup_{l=1}^{\infty} B_{l}^{+}(\sqrt{2 l-1})\right) \cup\left(\bigcup_{l=1}^{\infty} B_{l}^{-}(\sqrt{2 l-1})\right)
$$

As before, we make the choice of $\sqrt{2 l-1}$ for $\rho$ as that is the largest value of $\rho$ for which $B_{l}^{ \pm}(\rho)$ contains only a single point corresponding to a threshold. For certain combinations of $l$ and $\rho$ it can happen that $B_{l}^{+}(\rho)=B_{l}^{-}(\rho)$.

Consider a Schrödinger operator on $d$-dimensional Euclidean space with potential $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and scattering resolvent $R_{V_{0}, 0}(\lambda)$. When $\operatorname{Im} \lambda>0$, that is $\lambda$ is in the physical space,

$$
R_{V_{0}, 0}(\lambda)=\left(-\Delta_{0}+V_{0}-\lambda^{2}\right)^{-1}=\left(\left(-\Delta_{0}+\bar{V}_{0}-\bar{\lambda}^{2}\right)^{-1}\right)^{*}=\left(R_{\bar{V}_{0} 0}(-\bar{\lambda})\right)^{*}
$$

Here $\bar{V}_{0}$ and $\bar{\lambda}$ denote the usual complex conjugates. For odd $d$ the identity $R_{V_{0}, 0}(\lambda)=\left(R_{\bar{V}_{0} 0}(-\bar{\lambda})\right)^{*}$ then holds by meromorphic continuation for all $\lambda \in \mathbb{C}$. In particular, this implies $\lambda_{0}$ is a pole of $R_{V_{0}, 0}(\lambda)$ if and only if $-\bar{\lambda}_{0}$ is a pole of $R_{\bar{V}_{0} 0}(\lambda)$. For real-valued $V$, this is the well-known symmetry of resonances for symmetric Schrödinger operators in odd dimensions.

We turn to the analog of this result for $R_{V}$, which is shown in a similar way. Suppose $\zeta$ is in the physical space, here identified with the upper half-plane, so that $R_{V}(\zeta)=\left(-\Delta+V-\zeta^{2}\right)^{-1}$. Thus $\left(R_{\bar{V}}(-\bar{\zeta})\right)^{*}=R_{V}(\zeta)$. For general $\zeta \in \hat{Z}$, we define $-\zeta^{\dagger} \in \hat{Z}$ to be the point in $\hat{Z}$ which satisfies $\tau_{j}\left(-\zeta^{\dagger}\right)=-\bar{\tau}_{j}(\zeta)$ for all $j$. This is an antiholomorphic mapping, and if $\zeta$ is in the physical space,
identified with the upper half-plane, the mapping $\zeta \mapsto-\zeta^{\dagger}$ agrees with the mapping $\zeta \mapsto-\bar{\zeta}$. Then the identity

$$
\begin{equation*}
\left(R_{\bar{V}}\left(-\zeta^{\dagger}\right)\right)^{*}=R_{V}(\zeta), \quad \text { where } \tau_{j}\left(-\zeta^{\dagger}\right)=-\bar{\tau}_{j}(\zeta), \text { for all } j \in \mathbb{N}_{0} \tag{6-1}
\end{equation*}
$$

holds for all $\zeta \in \hat{Z}$ by meromorphic continuation. In particular, this means that $\zeta_{0} \in \hat{Z}$ is a pole of $R_{V}(\zeta)$ if and only if $-\zeta_{0}^{\dagger}$ is a pole of $R_{\bar{V}}(\zeta)$. Note that if $\zeta \in B_{l}^{+}(\rho)=B_{l}(\rho)$, then $-\zeta^{\dagger} \in B_{l}^{-}(\rho)$. Thus to study the poles of $R_{V}(\zeta)$ in $B_{l}^{-}(\rho)$ it suffices to study the poles of $R_{\bar{V}}(\zeta)$ in $B_{l}^{+}(\rho)=B_{l}(\rho)$. Likewise, an estimate on $R_{\bar{V}}$ in $B_{l}^{+}(\sqrt{2 l-1})$ implies an estimate on $R_{V}$ in $B_{l}^{-}(\sqrt{2 l-1})$.

We summarize these results in the following lemma.
Lemma 6.2. If $V_{0} \in L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, then $\lambda_{0}$ is a pole of $R_{V_{0}, 0}(\lambda)$ if and only if $-\bar{\lambda}_{0}$ is a pole of $R_{\bar{V}_{0} 0}(\lambda)$. Let $V \in L_{c}^{\infty}(X)$. Then $\zeta_{0} \in \hat{Z}$ is a pole of $R_{V}(\zeta)$ if and only if $-\zeta_{0}^{\dagger}$ is a pole of $R_{\bar{V}}(\zeta)$. Here $\bar{\lambda}_{0}, \bar{V}$, and $\bar{V}_{0}$ are the complex conjugates of $\lambda_{0}, V$, and $V_{0}$, respectively, and $-\zeta^{\dagger}$ is as defined in (6-1).

We define a distance on $\hat{Z}$ as follows: for $\zeta, \zeta^{\prime} \in \hat{Z}$,

$$
\begin{equation*}
d_{\hat{Z}}\left(\zeta, \zeta^{\prime}\right) \stackrel{\text { def }}{=} \sup _{j}\left|\tau_{j}(\zeta)-\tau_{j}\left(\zeta^{\prime}\right)\right| \tag{6-2}
\end{equation*}
$$

That this is well defined and a metric is shown in [Christiansen and Datchev 2021, Section 5.1]. Note that if $\zeta, \zeta^{\prime} \in \hat{Z}$ satisfy $\tau_{j}(\zeta) \neq-\tau_{j}\left(\zeta^{\prime}\right)$, then since $\tau_{j}(\zeta)^{2}-\tau_{j}\left(\zeta^{\prime}\right)^{2}=\tau_{l}(\zeta)^{2}-\tau_{l}\left(\zeta^{\prime}\right)^{2}$,

$$
\left|\tau_{j}(\zeta)-\tau_{j}\left(\zeta^{\prime}\right)\right|=\left|\tau_{l}(\zeta)-\tau_{l}\left(\zeta^{\prime}\right)\right|\left|\frac{\tau_{l}(\zeta)+\tau_{l}\left(\zeta^{\prime}\right)}{\tau_{j}(\zeta)+\tau_{j}\left(\zeta^{\prime}\right)}\right|
$$

In particular, this implies that for any $\rho>0$ there is an $L=L(\rho)$ so that if $l \geq L$ and $\zeta, \zeta^{\prime} \in B_{l}(\rho)$ then

$$
d_{\hat{Z}}\left(\zeta, \zeta^{\prime}\right)=\left|\tau_{l}(\zeta)-\tau_{l}\left(\zeta^{\prime}\right)\right| .
$$

Proof of Corollary 1.4. Recall our hypotheses include that $V$ is real-valued, ensuring that $V_{0}$ is real-valued as well.

The operator-valued function $R_{V}(\zeta)$ has a sequence $\left\{\zeta_{j}^{b}\right\}$ of poles satisfying $\left|\tau_{0}\left(\zeta_{j}^{b}\right)\right| \rightarrow \infty$ as $j \rightarrow \infty$ and $d_{\hat{Z}}\left(\zeta_{j}^{\mathrm{b}}\right.$, physical space $) \rightarrow 0$ only if $R_{V}(\zeta)$ has infinitely many poles in $\bigcup_{l=1}^{\infty} B_{l}(\sqrt{2 l-1})$ or infinitely many poles in $\bigcup_{l=1}^{\infty} B_{l}^{-}(\sqrt{2 l-1})$ (or both). If $R_{V}(\zeta)$ has infinitely many poles in $\bigcup_{l=1}^{\infty} B_{l}^{-}(\sqrt{2 l-1})$, then by Lemma $6.2, R_{\bar{V}}(\zeta)=R_{V}(\zeta)$ has infinitely many poles in $\bigcup_{l=1}^{\infty} B_{l}(\sqrt{2 l-1})$. Thus it suffices to study sequences of poles in $\bigcup_{l=1}^{\infty} B_{l}(\sqrt{2 l-1})$.

Note that while $B_{l}(\sqrt{2 l-1})$ contains only a single threshold, $B_{l}(\sqrt{2 l-1})$ and $B_{l+1}(\sqrt{2 l+1})$ are not disjoint and in fact have substantial overlap which contains an interval of the continuous spectrum. Moreover, for $l$ sufficiently large the sets $U_{l}^{+}$and $U_{l+1}^{-}$of Lemma 6.1 have nontrivial intersection. Applying Lemma 6.1 we see that in order to have a sequence of resonances contained in $\bigcup_{l=1}^{\infty} B_{l}(\sqrt{2 l-1})$ and approaching the continuous spectrum (and with $\left|\tau_{0}\right| \rightarrow \infty$ ), the resonances must lie in $\bigcup_{l=1}^{\infty} B_{l}(M)$ for some $M$. But then the corollary follows from an application of Theorems 1.1-1.3.

We now have the ingredients we need to prove Theorem 1.9.

Proof of Theorem 1.9. The hypotheses on $-d^{2} / d x^{2}+V_{0}$ and the expression (3-2) mean that the resolvent $R_{V_{0}}(\zeta)$ has no poles on the boundary of the physical space. Moreover, since for any $\tilde{\chi} \in C_{c}^{\infty}(\mathbb{R})$ there is a constant $C$ so that $\left\|\tilde{\chi} R_{V_{0}, 0}(\lambda) \tilde{\chi}\right\| \leq C$ for all $\lambda \in \mathbb{R} \cup i[0, \infty)$, for any $\chi \in C_{c}^{\infty}(X)$ there is a $C_{1}>0$ so that $\left\|\chi R_{V_{0}}(\zeta) \chi\right\| \leq C_{1}$ for all $\zeta$ in the boundary of the physical space.

Corollary 1.4 shows that there are no poles of the resolvent $R_{V}$ in the continuous spectrum at high energy. Proposition 5.4 and Lemma 5.5 show that when $\zeta$ is in the boundary of the physical space and $\zeta \in B_{l}(\sqrt{2 l-1})$, the cut-off resolvent of $-\Delta+V$ satisfies $\left\|\chi R_{V}(\zeta) \chi-\chi R_{V_{0}}(\zeta) \chi\right\|=O\left(l^{-1 / 2}\right)$. Thus $\left\|\chi R_{V_{0}}(\zeta) \chi\right\|$ is uniformly bounded on the boundary of the physical space when $\left|\tau_{0}(\zeta)\right|$ is sufficiently large. Hence by [Christiansen and Datchev 2021, Theorem 5.6] the hypotheses of [Christiansen and Datchev 2022, Theorem 4.1] hold. Theorem 1.9 then follows directly.

## 7. Larger neighborhoods of the threshold $\boldsymbol{l}^{\mathbf{2}}$

In this section we consider poles of $R_{V}(\zeta)$ in neighborhoods $B_{l}(\alpha \log l)$ and $B_{l}\left(\alpha(\log l)^{1-\epsilon}\right)$ of the $l$-th threshold. We prove Theorem 1.5 for potentials with $V_{0} \equiv 0$ and the related, but weaker, Theorem 7.1 which holds for a general potential $V \in L_{c}^{\infty}(X)$.

The proof of Theorem 1.5 is similar to that of the proof of Theorem 1.2 for $L^{\infty}$ potentials.
Proof of Theorem 1.5. Choose $\chi \in L_{c}^{\infty}(X), \chi V=V$, and $\chi$ independent of $\theta$. We write
$\chi R_{0} V R_{0} \chi$
$=\chi R_{0} \mathcal{P}_{l} V R_{0} \mathcal{P}_{l} \chi+\chi R_{0}\left(1-\mathcal{P}_{l}\right) V R_{0} \mathcal{P}_{l} \chi+\chi \mathcal{P}_{l} R_{0} V R_{0}\left(1-\mathcal{P}_{l}\right) \chi+\chi R_{0}\left(1-\mathcal{P}_{l}\right) V R_{0}\left(1-\mathcal{P}_{l}\right) \chi$.
Let $\alpha^{\prime}>0$, and let $\zeta \in B_{l}\left(\alpha^{\prime}|\log l|\right)$, where $l$ is large enough that $B_{l}\left(\alpha^{\prime}|\log l|\right)$ contains only a single point of $\hat{Z}$ which corresponds to a threshold. Let $\zeta \in B_{l}\left(\alpha^{\prime} \log l\right)$ satisfy $\left|\tau_{l}(\zeta)\right| \geq 1$. Then by Lemma 3.2,

$$
\left\|\chi R_{0}(\zeta)\left(1-\mathcal{P}_{l}\right) \chi\right\|=O\left(l^{-1 / 2}\right)
$$

and by (3-1) and [Dyatlov and Zworski 2019, Theorem 3.1],

$$
\left\|\chi R_{0}(\zeta) \mathcal{P}_{l} \chi\right\|=O\left(\frac{e^{C\left(\operatorname{Im} \tau_{l}(\zeta)\right)_{-}}}{\left|\tau_{l}(\zeta)\right|}\right)
$$

for some $C>0$. Using this estimate and $\mathcal{P}_{l} V \mathcal{P}_{l}=O\left(l^{-\delta}\right)$ in (7-1) shows

$$
\left\|\chi R_{0}(\zeta) V R_{0}(\zeta) \chi\right\|=O\left(l^{-\delta} e^{2 C\left(\operatorname{Im} \tau_{l}(\zeta)\right)_{-}}\right)
$$

Thus from (7-1) there is a $C_{1}>0$ so that $I+V R_{0}(\zeta) \chi$ is invertible if $l$ is sufficiently large, $\zeta \in B_{l}\left(\alpha^{\prime} \log l\right)$, $\left|\tau_{l}(\zeta)\right| \geq 1$, and $e^{2 C\left(\operatorname{Im} \tau_{l}(\zeta)\right)-} \leq C_{1} l^{\delta}$. This last item may be ensured by requiring $\left|\tau_{l}\right| \leq \alpha \log l$, for suitably chosen $\alpha>0, \alpha \leq \alpha^{\prime}$, and taking $l$ sufficiently large. Recall that $-\Delta+V$ has no resonances in regions where $I+V R_{0} \chi$ is invertible, see Proposition 4.3.

Applying Theorems 1.2 and 1.3 shows that if $d=1$ there is a $c_{0}>0$ so that when $l$ is sufficiently large the region $\left\{\zeta \in B_{l}(\alpha \log l): 1 \geq\left|\tau_{l}(\zeta)\right|>c_{0} l^{-\delta}\right\}$ contains no resonances, and if $d>1$ there are no resonances in $B_{l}(1)$ for $l$ sufficiently large.

A similar proof gives the next theorem.

Theorem 7.1. Let $V \in L_{c}^{\infty}(X)$ satisfy (1-1), and let $\epsilon>0$. Then there is a $c_{0}=c_{0}(\epsilon, V)>0$ so that for $l$ sufficiently large, the region

$$
\left\{\zeta \in B_{l}\left(c_{0}(\log l)^{1 /(d+\epsilon)}\right):\left|\tau_{l}(\zeta)-\lambda^{\prime}\right| \geq\left(1+\left|\lambda^{\prime}\right|^{2}\right)^{-(d+\epsilon) / 2} \text { for every } \lambda^{\prime} \in \mathbb{C}: m_{V_{0}, 0}\left(\lambda^{\prime}\right)>0\right\}
$$

contains no poles of $R_{V}(\zeta)$.
Proof. We assume $V^{\#}=V-V_{0} \not \equiv 0$, since otherwise there is nothing to prove.
Choose $\chi \in L_{c}^{\infty}(X)$ so that $\chi V=V$ and $\chi$ is independent of $\theta$. We may think of $\chi$ as an element of $L_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ as well.

Set

$$
A_{\epsilon} \stackrel{\text { def }}{=}\left\{\lambda \in \mathbb{C}:\left|\lambda-\lambda^{\prime}\right| \geq\left(1+\left|\lambda^{\prime}\right|^{2}\right)^{-(d+\epsilon) / 2} \text { for every } \lambda^{\prime} \in \mathbb{C}: m_{V_{0}, 0}\left(\lambda^{\prime}\right)>0\right\}
$$

We shall use, from the proof of [Dyatlov and Zworski 2019, Theorem 3.54], that there is a $C>0$ so that

$$
\begin{equation*}
\left\|\left(I+V_{0} R_{0,0}(\lambda)\right)^{-1}\right\| \leq C \exp \left(C|\lambda|^{d+\epsilon}\right) \quad \text { if } \lambda \in A_{\epsilon} \tag{7-2}
\end{equation*}
$$

Choose $\alpha^{\prime}>0$. If $\zeta \in B_{l}\left(\alpha^{\prime} \log l\right)$,

$$
\begin{aligned}
\chi R_{V_{0}}(\zeta) \mathcal{P}_{l} \chi & =\chi R_{V_{0}, 0}\left(\tau_{l}(\zeta)\right) \mathcal{P}_{l} \chi \\
& =\chi R_{0,0}\left(\tau_{l}(\zeta)\right) \chi\left(I+V_{0} R_{0,0}\left(\tau_{l}(\zeta)\right) \chi\right)^{-1} \mathcal{P}_{l}
\end{aligned}
$$

Thus, if $\zeta \in B_{l}\left(\alpha^{\prime} \log l\right)$ with $\tau_{l} \in A_{\epsilon}$ and $\left|\tau_{l}(\zeta)\right| \geq 1$, then

$$
\begin{align*}
\left\|\chi R_{V_{0}}(\zeta) \mathcal{P}_{l} \chi\right\| & \leq C \exp \left(C\left(\operatorname{Im} \tau_{l}(\zeta)_{-}\right)\right) \exp \left(C\left|\tau_{l}(\zeta)\right|^{d+\epsilon}\right) \\
& \leq C \exp \left(C\left|\tau_{l}(\zeta)\right|^{d+\epsilon}\right) \tag{7-3}
\end{align*}
$$

Here and below we allow the constant $C$ to change from line to line, and note that it depends on $V, \epsilon$, and $\chi$, but not $l$.

Let $\zeta \in B_{l}\left(\alpha^{\prime} \log l\right)$ with $\tau_{l} \in A_{\epsilon}$ and $\left|\tau_{l}(\zeta)\right| \geq 1$. Writing $\chi R_{V_{0}} \chi$ as in (7-1) and applying Lemma 3.2 and (7-3), we find that for these $\zeta$, if $l$ is sufficiently large,

$$
\begin{equation*}
\left\|\chi R_{V_{0}}(\zeta) V^{\#} R_{V_{0}}(\zeta) \chi\right\| \leq C_{1} l^{-\delta} \exp \left(C_{1}\left|\tau_{l}(\zeta)\right|^{d+\epsilon}\right) \tag{7-4}
\end{equation*}
$$

for some $C_{1}$. Now we can choose $c_{0}>0$ sufficiently small and $L>0$ sufficiently large so that

$$
\text { if }\left|\tau_{l}(\zeta)\right| \leq c_{0}(\log l)^{1 /(d+\epsilon)} \text { and } l>L \quad \text { then } C_{1} l^{-\delta} \exp \left(C_{1}\left|\tau_{l}(\zeta)\right|^{d+\epsilon}\right) \leq \frac{1}{2}
$$

ensuring that $I+V^{\#} R_{V_{0}}(\zeta) \chi$ is invertible.
Recalling that with $V^{\#}$ nontrivial if $I+V^{\#} R_{V_{0}}(\zeta) \chi$ is invertible then $\zeta$ is not a resonance of $-\Delta+V$ proves the theorem.

## 8. Expansion of $\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}} \chi\right)^{-1} V^{\#} \mathcal{P}_{l}$ for smooth $V$

This section contains preliminary computations which allow us to refine some of our results when $V$ is smooth. We begin with a straightforward lemma about Schrödinger operators on $\mathbb{R}^{d}$.

Lemma 8.1. Let $V_{0}, \chi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and $J \in \mathbb{N}$. Then as an operator from $H^{s}\left(\mathbb{R}^{d}\right)$ to $H^{s-2 J}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
\chi R_{V_{0}, 0}(\lambda) \chi=-\sum_{j=1}^{J} \frac{1}{\lambda^{2 j}} \chi\left(-\Delta_{0}+V_{0}\right)^{j-1} \chi+\frac{1}{\lambda^{2 J}} \chi R_{V_{0}, 0}(\lambda)\left(-\Delta_{0}+V_{0}\right)^{J} \chi \tag{8-1}
\end{equation*}
$$

Proof. First assume $\lambda$ is in the physical region, that is, $\operatorname{Im} \lambda>0$. Then the $J=1$ case follows from rearranging the equality

$$
\left(-\Delta_{0}+V_{0}-\lambda^{2}\right) R_{V_{0}, 0}(\lambda)=R_{V_{0}, 0}(\lambda)\left(-\Delta_{0}+V_{0}-\lambda^{2}\right)=I
$$

to get

$$
R_{V_{0}, 0}(\lambda)=\frac{1}{\lambda^{2}}\left(-I+R_{V_{0}, 0}(\lambda)\left(-\Delta_{0}+V_{0}\right)\right)
$$

The general case follows by induction.
Since both sides of (8-1) have meromorphic continuations to the complex plane, the equality holds for all $\lambda$.

We shall use the following Hilbert spaces: for $n \in \mathbb{N}_{0}$,

$$
H_{(0, n)}(X) \stackrel{\text { def }}{=}\left\{u \in L^{2}(X): \frac{\partial^{\alpha}}{\partial x^{\alpha}} u \in L^{2}(X) \text { if }|\alpha| \leq n\right\} \quad \text { with }\|u\|_{H_{(0, n)}}^{2}=\sum_{|\alpha| \leq n}\left\|\frac{\partial^{\alpha}}{\partial x^{\alpha}} u\right\|_{L^{2}(X)}^{2}
$$

Here we use the usual multi-index notation for $\alpha=\left(\alpha_{1}, \ldots, \alpha_{d}\right)$. This allows us to indicate mapping properties of operators which act differently in the $x$ and $\theta$ variables.

One of the main results of this section is the following proposition. Recall that $R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$ is defined in (5-3).

Proposition 8.2. Let $V, \chi \in C_{c}^{\infty}(X)$ satisfy $\chi V=V$. In addition, suppose $\chi$ is independent of $\theta$. Let $\lambda_{0} \in \mathbb{C}$, and suppose $R_{V_{0}, 0}(\lambda)$ is analytic on $0<\left|\lambda-\lambda_{0}\right| \leq \epsilon$. Then, for $R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$ and $\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right)$,

$$
\left\|\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \mathcal{P}_{l}+\frac{1}{l^{2}} \sum_{\substack{k \in \mathbb{Z} \\ k \neq 0}}\left(\frac{\tau_{l}^{2}-k^{2}}{4 k^{2}} V_{-k} V_{k}-\frac{V_{-k}}{4 k^{2}}\left(-\Delta_{0}+V_{0}\right) V_{k}\right) \mathcal{P}_{l}\right\|_{H_{(0,8)}(X) \rightarrow L^{2}(X)}=O\left(l^{-3}\right)
$$

where the error is uniform on $D_{l}\left(\lambda_{0}, \epsilon\right)$ for $l$ sufficiently large.
To prove this proposition we use Lemmas 8.3-8.6. In each of these, $V, \lambda_{0}, R_{V_{0}}^{\text {reg }}(\zeta)$, and $\epsilon$ are as in Proposition 8.2. Some of these computations rely on the identity $e^{ \pm i k \theta} e^{ \pm i l \theta}=e^{ \pm i(k+l) \theta}$ and hence use the structure of the eigenfunctions of the Laplacian on $\mathbb{S}^{1}$ in an essential way.

For $l \in \mathbb{N}$, let $\mathcal{P}_{l \pm}: L^{2}(X) \rightarrow L^{2}(X)$ denote orthogonal projection onto $L^{2}\left(\mathbb{R}_{x}^{d}\right) e^{ \pm i l \theta}$, so that

$$
\left(\mathcal{P}_{l \pm} f\right)(x, \theta)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f\left(x, \theta^{\prime}\right) e^{ \pm i l\left(\theta-\theta^{\prime}\right)} d \theta^{\prime}
$$

for $l>0$ and $\mathcal{P}_{l}=\mathcal{P}_{l+}+\mathcal{P}_{l-}$.

Lemma 8.3. Under the hypotheses of Proposition 8.2,

$$
\left\|\mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) V^{\#} \mathcal{P}_{l}-\frac{1}{l^{2}} \sum_{k \in \mathbb{Z}, k \neq 0}\left(\frac{\tau_{l}^{2}-k^{2}}{4 k^{2}} V_{-k} V_{k}-\frac{V_{-k}}{4 k^{2}}\left(-\Delta_{0}+V_{0}\right) V_{k}\right) \mathcal{P}_{l}\right\|_{H_{(0, n+6)} \rightarrow H_{(0, n)}}=O\left(l^{-3}\right)
$$

uniformly for $\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right)$ when $l$ is sufficiently large.
Proof. Since $V \in C_{c}^{\infty}(X)$, we have $\left\|V_{m}\right\|_{L^{\infty}}=O\left(|m|^{-N}\right)$ for any $N$, so $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-N}\right)$. Thus, choosing $l$ sufficiently large so that (5-4) holds, it suffices to consider

$$
\mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}=\mathcal{P}_{l} V^{\#} R_{V_{0}}(\zeta)\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}
$$

Then

$$
\begin{aligned}
\mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l} & =\sum_{ \pm} \sum_{\substack{k \in \mathbb{Z} \\
0<|k|, k \neq-l}} V_{\mp k} R_{V_{0}, 0}\left(\tau_{l+k}\right) V_{ \pm k} \mathcal{P}_{l \pm} \\
& =\sum_{ \pm} \sum_{\substack{k \in \mathbb{Z} \\
0<|k|<l^{1 / 2}}} V_{\mp k} R_{V_{0}, 0}\left(\tau_{l+k}\right) V_{ \pm k} \mathcal{P}_{l \pm}+O_{L^{2} \rightarrow L^{2}}\left(l^{-N}\right)
\end{aligned}
$$

Here we use the rapid decay of $\left\|V_{m}\right\|$ to bound the error obtained when we restrict the values of $k$ in the sum. Using Lemma 8.1 with $J=3$ gives

$$
\begin{equation*}
\mathcal{P}_{l} V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}=\sum_{ \pm} \sum_{\substack{k \in \mathbb{Z} \\ 0<|k|<l^{1 / 2}}} V_{\mp k}\left(\frac{-1}{\tau_{l+k}^{2}}-\frac{1}{\tau_{l+k}^{4}}\left(-\Delta_{0}+V_{0}\right)\right) V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-3}\right) \tag{8-2}
\end{equation*}
$$

where the error is as an operator from $H_{(0, n+6)}(X)$ to $H_{(0, n)}(X)$ and is uniform in $D_{l}\left(\lambda_{0}, \epsilon\right)$. Since we have restricted $|k|$ to be relatively small compared with $l$, we can expand $\tau_{l \pm k}$ asymptotically in $l$. Thus, with each sum over $k \in \mathbb{Z}$ with $0<|k|<l^{1 / 2}$, using $\tau_{l \pm k}^{2}=\tau_{l}^{2} \mp 2 l k-k^{2}$ gives

$$
\begin{align*}
\sum_{0<|k|<l^{1 / 2}} \frac{1}{\tau_{l+k}^{2}} V_{-k} V_{k} & =\frac{1}{2} \sum_{0<|k|<l^{1 / 2}}\left(\frac{1}{\tau_{l+k}^{2}}+\frac{1}{\tau_{l-k}^{2}}\right) V_{-k} V_{k}=\sum_{0<|k|<l^{1 / 2}} \frac{\tau_{l}^{2}-k^{2}}{\left(\tau_{l}^{2}-k^{2}\right)^{2}-4 k^{2} l^{2}} V_{-k} V_{k} \\
& =\frac{-1}{4 l^{2}} \sum_{0<|k|<l^{1 / 2}}\left(\frac{\tau_{l}^{2}-k^{2}}{k^{2}}\right) V_{-k} V_{k}+O\left(l^{-4}\right) \tag{8-3}
\end{align*}
$$

Here and below the error is uniform in $D_{l}\left(\lambda_{0}, \epsilon\right)$ when $l$ is sufficiently large.
For the second term in (8-2), we write

$$
\begin{aligned}
\sum_{0<|k|<l^{1 / 2}} \frac{1}{\tau_{l+k}^{4}} V_{\mp k}\left(-\Delta_{0}+V_{0}\right) V_{ \pm k} & =\sum_{0<|k|<l^{1 / 2}} \frac{1}{\left(\tau_{l}^{2}-2 l k-k^{2}\right)^{2}} V_{\mp k}\left(-\Delta_{0}+V_{0}\right) V_{ \pm k} \\
& =\frac{1}{4 l^{2}} \sum_{0<|k|<l^{1 / 2}} \frac{1}{k^{2}} V_{\mp k}\left(-\Delta_{0}+V_{0}\right) V_{ \pm k}+O\left(l^{-3}\right)
\end{aligned}
$$

Note that

$$
\begin{equation*}
\sum_{0<|k|<l^{1 / 2}} \frac{1}{k^{2}} V_{\mp k}\left(-\Delta_{0}+V_{0}\right) V_{ \pm k}=\sum_{0<|k|<l^{1 / 2}} \frac{1}{k^{2}} V_{-k}\left(-\Delta_{0}+V_{0}\right) V_{k} \tag{8-4}
\end{equation*}
$$

since the sum is over $k \in \mathbb{Z}$, with $0<|k|<l^{1 / 2}$. The rapid decay in $m$ of $\left\|V_{m}\right\|_{C^{p}}$ means we can replace the sums in (8-3) and (8-4) over $0<|k|<l^{1 / 2}$ by sums over all nonzero $k \in \mathbb{Z}$, with an error which is $O\left(l^{-N}\right)$.

The next lemma is an algebraic identity.
Lemma 8.4. For any $V \in C_{c}^{\infty}(X)$

$$
\sum_{\substack{m, j \in \mathbb{Z} \\ m, j \neq 0, m \neq-j}} \frac{1}{j(j+m)} V_{m} V_{j} V_{-m-j}=0
$$

We give two different proofs.
Proof. For this proof, we show that for each $j_{0} \neq 0, m_{0} \neq 0$ the coefficient of $V_{j_{0}} V_{m_{0}} V_{-j_{0}-m_{0}}$ in the sum is zero. This proof is purely algebraic in nature.

If $m_{0} \neq \pm j_{0}$, then there are six possibilities for the pair $(j, m)$ which will give a term containing $V_{m_{0}} V_{j_{0}} V_{-m_{0}-j_{0}}:\left(j_{0}, m_{0}\right),\left(m_{0}, j_{0}\right),\left(-m_{0}-j_{0}, m_{0}\right),\left(m_{0},-j_{0}-m_{0}\right),\left(j_{0},-m_{0}-j_{0}\right),\left(-m_{0}-j_{0}, j_{0}\right)$. Thus the sum of the coefficients of $V_{m_{0}} V_{j_{0}} V_{-m_{0}-j_{0}}$ is

$$
\frac{1}{j_{0}\left(j_{0}+m_{0}\right)}+\frac{1}{m_{0}\left(j_{0}+m_{0}\right)}+\frac{1}{j_{0}\left(j_{0}+m_{0}\right)}-\frac{1}{j_{0} m_{0}}-\frac{1}{j_{0} m_{0}}+\frac{1}{m_{0}\left(j_{0}+m_{0}\right)}=0
$$

A similar argument when $j_{0}=m_{0}$ shows the coefficient of $V_{j_{0}}^{2} V_{-2 j_{0}}$ is zero as well.
Alternate proof of Lemma 8.4. For this proof, we use that $V_{j}$ is the $j$-th Fourier coefficient of $V$. Though in our applications $V_{j}$ depends on $x$, that dependence is not important here so we will suppress it.

Set

$$
W(\theta)=\sum_{j \neq 0} \frac{1}{j} V_{j} e^{i j \theta}
$$

and note $d / d \theta W(\theta)=V(\theta)-V_{0}$. Then

$$
\begin{equation*}
\int_{0}^{2 \pi}\left(V(\theta)-V_{0}\right)(W(\theta))^{2} d \theta=\left.\frac{1}{3}(W(\theta))^{3}\right|_{0} ^{2 \pi}=0 \tag{8-5}
\end{equation*}
$$

by the fundamental theorem of calculus. But

$$
\begin{align*}
\sum_{\substack{m, j \in \mathbb{Z} \\
m, j \neq 0 \\
m \neq-j}} \frac{1}{j(j+m)} V_{m} V_{j} V_{-m-j} & =\sum_{\substack{m, j \in \mathbb{Z} \\
m, j \neq 0 \\
m \neq-j}} \frac{-1}{j m} V_{m} V_{j} V_{-m-j} \\
& =-\int_{0}^{2 \pi}\left(V(\theta)-V_{0}\right)(W(\theta))^{2} d \theta \tag{8-6}
\end{align*}
$$

where the last equality uses $e^{i j \theta} e^{i m \theta}=e^{i(j+m) \theta}$ and the fact that the integral of a function over a circle is its zeroth Fourier coefficient. Combining (8-5) and (8-6) proves the lemma.

Lemma 8.5. Under the hypotheses of Proposition 8.2, if $l$ is sufficiently large

$$
\left\|\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}\right\|_{H_{(0, n+6)(X)} \rightarrow H_{(0, n)(X)}}=O\left(l^{-3}\right) \quad \text { uniformly for } \zeta \in D_{l}\left(\lambda_{0}, \epsilon\right)
$$

Proof. Again we use that $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-N}\right)$ for any $N$. This implies

$$
\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}=\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(I-\mathcal{P}_{l}\right)\right)^{2} V^{\#} \mathcal{P}_{l}+O_{L^{2} \rightarrow L^{2}}\left(l^{-N}\right)
$$

Note that for $\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right)$ and $l$ sufficiently large, $R_{V_{0}}^{\mathrm{reg}}(\zeta)\left(I-\mathcal{P}_{l}\right)=R_{V_{0}}(\zeta)\left(I-\mathcal{P}_{l}\right)$. Then $\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{2} V^{\#} \mathcal{P}_{l}$

$$
\begin{align*}
& =\mathcal{P}_{l} \sum_{ \pm} e^{ \pm i(j+k+m) \theta} \sum_{\substack{k, m, j \in \mathbb{Z} \\
k, m+k \neq 0,-2 l \\
m, j \neq 0}} V_{ \pm j} R_{V_{0}, 0}\left(\tau_{l+k+m}\right) V_{ \pm m} R_{V_{0}, 0}\left(\tau_{l+k}\right) V_{ \pm k} \mathcal{P}_{l \pm} \\
& =\sum_{ \pm \begin{array}{c}
k, m+k \neq 0,-2 l \\
m \neq 0, k, m \in \mathbb{Z}
\end{array}} V_{\mp(k+m)} R_{V_{0}, 0}\left(\tau_{l+k+m}\right) V_{ \pm m} R_{V_{0}, 0}\left(\tau_{l+k}\right) V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-N}\right) \tag{8-7}
\end{align*}
$$

By Lemma 8.1, for $k, m+k \neq 0,-2 l$,

$$
\begin{aligned}
&\left\|\chi R_{V_{0}, 0}\left(\tau_{l+k+m}\right) V_{ \pm m} R_{V_{0}, 0}\left(\tau_{l+k}\right) V_{ \pm k}-\frac{1}{\tau_{l+k+m}^{2} \tau_{l+k}^{2}} \chi V_{ \pm m} V_{ \pm k}\right\|_{H^{n+6}\left(\mathbb{R}^{d}\right) \rightarrow H^{n}\left(\mathbb{R}^{d}\right)} \\
&=O\left(l^{-3}\left\|V_{ \pm k}\right\|_{C^{6+n}}\left\|V_{ \pm m}\right\|_{C^{6+n}}\right)
\end{aligned}
$$

This implies (with sums still over $\mathbb{Z}$ ), using $\left\|V_{m}\right\|_{C^{p}}=O\left(|m|^{-N}\right)$, that

$$
\begin{align*}
\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{2} V^{\#} \mathcal{P}_{l} & =\sum_{ \pm} \sum_{k, m, k+m \neq 0,-2 l} \frac{1}{\tau_{l+k+m}^{2} \tau_{l+k}^{2}} V_{\mp(k+m)} V_{ \pm m} V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-3}\right) \\
& =\sum_{ \pm} \sum_{0<|k|,|k+m|<l^{1 / 2}, m \neq 0} \frac{1}{\tau_{l+k+m}^{2} \tau_{l+k}^{2}} V_{\mp(k+m)} V_{ \pm m} V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-3}\right) \\
& =\sum_{ \pm} \sum_{0<|k|,|k+m|<l^{1 / 2}, m \neq 0} \frac{1}{4 l^{2} k(k+m)} V_{\mp(k+m)} V_{ \pm m} V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-3}\right) \\
& =\sum_{ \pm} \sum_{0 \neq k, k+m, m} \frac{1}{4 l^{2} k(k+m)} V_{\mp(k+m)} V_{ \pm m} V_{ \pm k} \mathcal{P}_{l \pm}+O\left(l^{-3}\right) \tag{8-8}
\end{align*}
$$

Here errors are as operators from $H_{(0, n+6)}(X)$ to $H_{(0, n)}(X)$, and are uniform in $D_{l}\left(\lambda_{0}, \epsilon\right)$ when $l$ is sufficiently large. But the final sum in (8-8) is zero by Lemma 8.4.

Lemma 8.6. Under the hypotheses of Proposition 8.2 for $j \geq 3, j \in \mathbb{N}$, and $l$ sufficiently large,

$$
\left\|\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta)\right)^{j} V^{\#} \mathcal{P}_{l}\right\|_{H_{(0,8)}(X) \rightarrow L^{2}(X)}=O\left(l^{-3}\right)
$$

uniformly for $\zeta \in D_{l}\left(\lambda_{0}, \epsilon\right)$.
Proof. By Lemma 8.5,

$$
\left\|\mathcal{P}_{l}\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}\right\|_{H_{(0, n+6)} \rightarrow H_{(0, n)}}=O\left(l^{-3}\right)
$$

This gives

$$
\begin{align*}
\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{3} V^{\#} \mathcal{P}_{l} & =V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(I-\mathcal{P}_{l}\right)\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}+V^{\#} R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l}\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l} \\
& =V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(I-\mathcal{P}_{l}\right)\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}+O\left(l^{-3}\right) \tag{8-9}
\end{align*}
$$

as an operator from $H_{(0, n+6)}(X)$ to $H_{(0, n)}(X)$. Using that $\mathcal{P}_{l}$ commutes with $R_{V_{0}}$ and $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-N}\right)$ for any $N$ gives

$$
\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{2} V^{\#} \mathcal{P}_{l}=\left(V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{2} V^{\#} \mathcal{P}_{l}+V^{\#} R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l} V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}+O_{L^{2} \rightarrow L^{2}}\left(l^{-N}\right)
$$

Using this in (8-9) yields

$$
\begin{align*}
& \left(V^{\#} R_{V_{0}}^{\mathrm{reg}}\right)^{3} V^{\#} \mathcal{P}_{l}=\left(V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{3} V^{\#} \mathcal{P}_{l} \\
& \quad+V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l} V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}+O_{H_{(0, n+6)} \rightarrow H_{(0, n)}}\left(l^{-3}\right) \tag{8-10}
\end{align*}
$$

For large $l$, Lemma 8.1 applied with $J=1$ shows that

$$
\left\|\left(V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{3} V^{\#} \mathcal{P}_{l}\right\|_{H_{(0,6)}(X) \rightarrow L^{2}(X)}=O\left(l^{-3}\right)
$$

Choose $\chi \in C_{c}^{\infty}(X)$ independent of $\theta$ so that $V \chi=V$. We write the second term on the right in (8-10) as the composition of three operators, with the grouping indicated below by the brackets:

$$
\begin{align*}
& V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l} V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l} \\
&=\left[V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#}\right]\left[\chi R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l} \chi\right]\left[\mathcal{P}_{l} V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#} \mathcal{P}_{l}\right] \tag{8-11}
\end{align*}
$$

By Lemma 8.1,

$$
\left\|V^{\#} R_{V_{0}}\left(I-\mathcal{P}_{l}\right) V^{\#}\right\|_{H_{(0, n+2)} \rightarrow H_{(0, n)}}=O\left(l^{-1}\right) .
$$

The second operator, $\chi R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l} \chi$, is bounded. By Lemma 8.3, the third is $O\left(l^{-2}\right)$ as an operator from $H_{(0, n+6)}$ to $H_{(0, n)}$. Thus we have proved the lemma when $j=3$.

The case of $j>3$ follows from the $j=3$ case.
We now can prove Proposition 8.2.
Proof of Proposition 8.2. For $l$ sufficiently large, on $D_{l}\left(\lambda_{0}, \epsilon\right)$,

$$
\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \mathcal{P}_{l}=\mathcal{P}_{l} \sum_{j=0}^{\infty}\left(-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{j} V^{\#} \mathcal{P}_{l}
$$

The proposition then follows from an application of Lemmas 8.3, 8.5, and 8.6, and recalling that $\left\|\mathcal{P}_{l} V^{\#} \mathcal{P}_{l}\right\|=O\left(l^{-N}\right)$.

The proof of Theorem 1.6 uses the next lemma, which computes an expression related to the leading term of

$$
\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(\zeta_{l}(z)\right)\right)^{-1} V^{\#} \mathcal{P}_{l}
$$

Lemma 8.7. Suppose $V \in C_{c}^{\infty}(X)$ and $u \in H^{2}\left(\mathbb{R}^{d}\right)$ satisfies $\left(-\Delta_{0}+V_{0}-\lambda_{0}^{2}\right) u=0$. Then $-\int_{\mathbb{R}^{d}} u\left(\left(z^{2}-k^{2}\right) V_{-k} V_{k} u-V_{-k}\left(-\Delta_{0}+V_{0}\right)\left(V_{k} u\right)\right) d x=\int_{\mathbb{R}^{d}}\left(\left(k^{2}+\lambda_{0}^{2}-z^{2}\right) u^{2} V_{-k} V_{k}+u^{2} \nabla_{0} V_{-k} \cdot \nabla_{0} V_{k}\right) d x$.

Proof. We first compute $\int_{\mathbb{R}} u V_{-k}\left(-\Delta_{0}+V_{0}\right)\left(V_{k} u\right) d x$. Expanding and then integrating by parts yields

$$
\begin{align*}
\int_{\mathbb{R}^{d}} u V_{-k}\left(-\Delta_{0}\right. & \left.+V_{0}\right)\left(V_{k} u\right) d x \\
& =-\int_{\mathbb{R}^{d}}\left(u^{2} V_{-k} \Delta_{0} V_{k}+2 V_{-k} u \nabla_{0} V_{k} \cdot \nabla_{0} u\right) d x+\int_{\mathbb{R}^{d}} u V_{-k} V_{k}\left(-\Delta_{0}+V_{0}\right) u d x \\
& =-\int_{\mathbb{R}^{d}} u^{2} V_{-k} \Delta_{0} V_{k} d x+\int_{\mathbb{R}^{d}} u^{2} \sum_{j=1}^{d} \frac{\partial}{\partial x_{j}}\left(V_{-k} \frac{\partial}{\partial x_{j}} V_{k}\right) d x+\lambda_{0}^{2} \int_{\mathbb{R}^{d}} u^{2} V_{-k} V_{k} d x \\
& =\int_{\mathbb{R}^{d}} u^{2} \nabla_{0} V_{-k} \cdot \nabla_{0} V_{k} d x+\lambda_{0}^{2} \int_{\mathbb{R}^{d}} u^{2} V_{-k} V_{k} d x . \tag{8-12}
\end{align*}
$$

Using this, we find
$\int_{\mathbb{R}^{d}}\left(\left(z^{2}-k^{2}\right) V_{-k} V_{k} u^{2}-u V_{-k}\left(-\Delta_{0}+V_{0}\right)\left(V_{k} u\right)\right) d x=-\int_{\mathbb{R}^{d}}\left(\left(\left(k^{2}+\lambda_{0}^{2}-z^{2}\right) V_{-k} V_{k}+\nabla_{0} V_{-k} \cdot \nabla_{0} V_{k}\right) u^{2}\right) d x$, completing the proof.

The proof of the next lemma uses some of the same ideas as that of Proposition 8.2. This result will be used in the proof of Theorem 1.7.

Lemma 8.8. Suppose $V \in C_{c}^{\infty}(X ; \mathbb{R})$. Let $\lambda_{0} \in i \mathbb{R}$ be a simple pole of $R_{V_{0}, 0}(\lambda)$ with residue of rank 1 . Let $M>\left|\lambda_{0}\right|$ and $N \in \mathbb{N}$, and suppose $R_{V_{0}, 0}(\lambda)-\Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right)$ is analytic for $\left|\lambda-\lambda_{0}\right| \leq \epsilon$. Then if $\chi \in C_{c}^{\infty}(X ; \mathbb{R})$ is independent of $\theta$ and satisfies $V \chi=V$, there exists an $s=s(N) \in \mathbb{N}$ and an $A_{N}=A_{N}\left(\tau_{l}, l\right): H_{(0, s)}(X) \rightarrow L^{2}(X)$ so that forl sufficiently large,

$$
\begin{equation*}
\left\|\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}, 0}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \mathcal{P}_{l}-A_{N}\left(\tau_{l}(\zeta), l\right)\right\|_{H_{(0, s)}(X) \rightarrow L^{2}(X)}=O\left(l^{-N}\right) \tag{8-13}
\end{equation*}
$$

uniformly for $\zeta \in \bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$. Moreover, $A_{N}(z, l)$ depends analytically on $z$ in the set $\left\{z \in \mathbb{C}:\left|z-\lambda_{0}\right| \leq \epsilon\right\}$ and if $z \in i \mathbb{R}$, then $A_{N}(z, l)$ is symmetric on $C_{c}^{\infty}(X) \subset L^{2}(X)$. Furthermore,

$$
\left\|\mathcal{P}_{l \pm} A_{N} \mathcal{P}_{l \mp}\right\|_{H_{(0, s)}(X) \rightarrow L^{2}(X)}=O\left(l^{-N}\right)
$$

for any $N$.
Proof. By Lemma 5.2, if $j>2 N$, then on $\bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$ we have

$$
\left\|\left(V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{j}\right\|_{L^{2}(X) \rightarrow L^{2}(X)}=O\left(l^{-N}\right)
$$

Thus

$$
\begin{equation*}
\left\|\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1}-\sum_{j=0}^{2 N}\left(-V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{j}\right\|_{L^{2}(X) \rightarrow L^{2}(X)}=O\left(l^{-N}\right) \tag{8-14}
\end{equation*}
$$

Now we write, for $l$ sufficiently large,

$$
\begin{equation*}
R_{V_{0}}^{\mathrm{reg}}=R_{V_{0}}^{\mathrm{reg}} \mathcal{P}_{l}+R_{V_{0}}\left(I-\mathcal{P}_{l}\right) \tag{8-15}
\end{equation*}
$$

From our assumptions on $V_{0}$ and the pole of $R_{V_{0}, 0}$ at $\lambda_{0}$, there is a $u \in C^{\infty}\left(\mathbb{R}^{d} ; \mathbb{R}\right)$ so that for $\left|\lambda-\lambda_{0}\right| \leq \epsilon$, $R_{V_{0}, 0}(\lambda)-i /\left(\lambda-\lambda_{0}\right) u \otimes u$ is analytic. Then for $l$ sufficiently large

$$
R_{V_{0}}^{\mathrm{reg}}(\zeta) \mathcal{P}_{l}=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right) \mathcal{P}_{l}=R_{V_{0}, 0}\left(\tau_{l}(\zeta)\right) \mathcal{P}_{l}-\frac{i}{\tau_{l}(\zeta)-\lambda_{0}}(u \otimes u) \mathcal{P}_{l}
$$

If $\tau_{l}=\tau_{l}(\zeta) \in i \mathbb{R}$ and $\zeta \in \bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$, the operator $\chi R_{V_{0}}^{\mathrm{reg}}(\zeta) \mathcal{P}_{l} \chi$ is symmetric on $C_{c}^{\infty}(X)$. On the other hand, for $k \neq l$, writing $\tau_{k}$ for $\tau_{k}(\zeta)$ and using Lemma 8.1 yields

$$
\begin{align*}
& \chi R_{V_{0}} \mathcal{P}_{k} \chi=\chi R_{V_{0}, 0}\left(\tau_{k}\right) \mathcal{P}_{k} \chi \\
& \begin{aligned}
= & -\chi \sum_{j=1}^{N} \frac{1}{\left(\tau_{l}^{2}+l^{2}-k^{2}\right)^{j}}\left(-\Delta_{0}+V_{0}\right)^{j-1} \mathcal{P}_{k} \chi \\
& +\chi \frac{1}{\left(\tau_{l}^{2}+l^{2}-k^{2}\right)^{N}} R_{V_{0}}\left(\tau_{k}\right)\left(-\Delta_{0}+V_{0}\right)^{N} \mathcal{P}_{k} \chi .
\end{aligned}
\end{align*}
$$

If $\tau_{l}^{2} \in \mathbb{R}$, then

$$
\chi \frac{1}{\left(\tau_{l}^{2}+l^{2}-k^{2}\right)^{j}}\left(-\Delta_{0}+V_{0}\right)^{j-1} \mathcal{P}_{k} \chi
$$

is symmetric on $C_{c}^{\infty}(X)$. Set

$$
\begin{equation*}
T_{N}=T_{N}\left(\tau_{l}, l\right)=R_{V_{0}, 0}^{\mathrm{reg}}\left(\tau_{l}\right) \mathcal{P}_{l}-\sum_{k \neq l} \sum_{j=1}^{N} \frac{1}{\left(\tau_{l}^{2}+l^{2}-k^{2}\right)^{j}}\left(-\Delta_{0}+V_{0}\right)^{j-1} \mathcal{P}_{k} \tag{8-17}
\end{equation*}
$$

Note that $T_{N}$ is an analytic operator-valued function of $\tau_{l}$ for $\zeta \in \bar{D}_{l}\left(\lambda_{0}, \epsilon\right)$, where $\left|\tau_{l}-\lambda_{0}\right| \leq \epsilon$. Using (8-16),

$$
\left\|\chi\left(R_{V_{0}}^{\mathrm{reg}}-T_{N}\right) \chi\right\|_{H_{(0,2 N+t)}(X) \rightarrow H_{(0, t)}(X)}=O\left(l^{-N}\right)
$$

if $\left|\tau_{l}-\lambda_{0}\right| \leq \epsilon$, and $\chi T_{N}\left(\tau_{l}, l\right) \chi$ is symmetric on $C_{c}^{\infty}(X)$ if $\tau_{l} \in i \mathbb{R}$. Moreover, by (8-14),

$$
\left\|\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}\left(\zeta_{l}\left(\tau_{l}\right)\right) \chi\right)^{-1}-\sum_{j=0}^{2 N}\left(-V^{\#} T_{N}\left(\tau_{l}, l\right)\right)^{j} \chi\right\|_{H_{(0, s(N))} \rightarrow L^{2}}=O\left(l^{-N}\right)
$$

if $s(N) \geq 4 N^{2}$. Thus if we define

$$
\begin{equation*}
A_{N}=A_{N}\left(\tau_{l}, l\right)=\mathcal{P}_{l} \sum_{j=0}^{2 N}\left(-V^{\#} T_{N}\right)^{j} V^{\#} \mathcal{P}_{l} \tag{8-18}
\end{equation*}
$$

then $A_{N}$ satisfies (8-13), $A_{N}$ is an analytic function of $\tau_{l}$ if $\left|\tau_{l}-\lambda_{0}\right| \leq \epsilon$, and $A_{N}\left(\tau_{l}, l\right)$ is symmetric on $C_{c}^{\infty}(X)$ if $\tau_{l} \in i \mathbb{R}$.

To show that $\left\|\mathcal{P}_{l \pm} A_{N} \mathcal{P}_{l \mp}\right\|_{H_{(0, s)} \rightarrow L^{2}}=O\left(l^{-N}\right)$, consider a term $\mathcal{P}_{l+}\left(V^{\#} T_{N}\right)^{j} V^{\#} \mathcal{P}_{l-}$. We write

$$
\mathcal{P}_{l+}\left(V^{\#} T_{N}\right)^{j} V^{\#} \mathcal{P}_{l-}=\sum_{\substack{m_{1}+m_{2}+\cdots+m_{j+1}=2 l \\ m_{k} \neq 0}} V_{m_{1}} e^{i m_{1} \theta} T_{N} V_{m_{2}} e^{i m_{2} \theta} T_{N} \cdots V_{m_{j}} e^{i m_{j} \theta} T_{N} V_{m_{j+1}} e^{i m_{j+1} \theta} \mathcal{P}_{l-}
$$

Thus we see that at least one $m_{n}$ must have absolute value at least $2 l /(j+1)$. Since $\left\|V_{m}\right\|_{C^{r}}=O\left(|m|^{-p}\right)$ for any fixed $r$, any $p$, we obtain

$$
\left\|\mathcal{P}_{l+}\left(V^{\#} T_{N}\right)^{j} V^{\#} \mathcal{P}_{l-}\right\|_{H_{(0, s)} \rightarrow L^{2}}=O\left(l^{-N}\right)
$$

for some sufficiently large $s$. Thus the result for $\mathcal{P}_{l+} A_{N} \mathcal{P}_{l-}$ follows from our expression (8-18) for $A_{N}$. The result for $\mathcal{P}_{l-} A_{N} \mathcal{P}_{l+}$ follows similarly.

## 9. Proofs of the smooth case of Theorem 1.2 and Theorem 1.3

The first application of our results in the previous section is to improve the localization of the resonances when $V \in C_{c}^{\infty}(X)$.

Proof of Theorem 1.2 for $V \in C_{c}^{\infty}(X)$. Let $\lambda_{j} \in \Lambda_{\rho}$ and choose $\epsilon>0$ so that there are no poles of $R_{V_{0}, 0}(\lambda)$ in $0<\left|\lambda-\lambda_{j}\right| \leq \epsilon$. We will show that there is a $C_{j}>0$ so that there are no poles of $R_{V}(\zeta)$ in $\zeta \in D_{l}\left(\lambda_{j}, \epsilon\right)$ with $\left|\tau_{l}(\zeta)-\lambda_{j}\right|>C_{j} l^{-2 /\left(m_{V_{0}, 0}\left(\lambda_{j}\right)\right)}$ when $l$ is sufficiently large.

Choose $\chi \in C_{c}^{\infty}(X)$ so that $\chi V=V$ and $\chi$ is independent of $\theta$. As previously, if $l$ is sufficiently large,

$$
R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{j}, l\right)=R_{V_{0}}(\zeta)-\left.\Xi\left(R_{V_{0}, 0}, \lambda_{j}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}
$$

and note that $R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{j}, l\right)$ is analytic on $\bar{D}_{l}\left(\lambda_{j}, \epsilon\right)$. By (3-7), any poles of $R_{V}(\zeta)$ in $D_{l}\left(\lambda_{j}, \epsilon\right)$ are points at which $I+\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{j}\right) \chi \mathcal{P}_{l}$ has nontrivial null space.

Using the smoothness of $V$, for any fixed $s \in \mathbb{N}$ there is a constant $C>0$ (depending on $s, V_{0}, \lambda_{j}$ ) with

$$
\begin{equation*}
\left\|\left.V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{j}\right)\right|_{\lambda=\tau_{l}(\zeta)} \chi \mathcal{P}_{l}\right\|_{L^{2}(X) \rightarrow H_{(0, s)}(X)} \leq \frac{C}{\left|\tau_{l}(\zeta)-\lambda_{j}\right|^{m_{V_{0}, 0}\left(\lambda_{j}\right)}} \tag{9-1}
\end{equation*}
$$

[Dyatlov and Zworski 2019, Theorems 2.5, 2.7, 3.9, and 3.17]. Thus on $D_{l}\left(\lambda_{j}, \epsilon\right)$, for $l$ sufficiently large by Proposition 8.2,

$$
\left\|\left.\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{j}\right)\right|_{\lambda=\tau_{l}(\zeta)} \chi \mathcal{P}_{l}\right\|_{L^{2}(X) \rightarrow L^{2}(X)} \leq \frac{C}{l^{2}\left|\tau_{l}(\zeta)-\lambda_{j}\right|^{m_{V_{0}, 0}\left(\lambda_{j}\right)}}
$$

for some $C$. Thus there is a $C_{j}>0$ so that if $\zeta \in D_{l}\left(\lambda_{j}, \epsilon\right)$ and $\left|\tau_{l}(\zeta)-\lambda_{j}\right|>C_{j} l^{-2 / m_{V_{0}, 0}\left(\lambda_{j}\right)}$, then $I+\mathcal{P}_{l}\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{j}\right) \mathcal{P}_{l}$ is invertible, and $\zeta$ is not a resonance.

Since $\lambda_{j} \in \Lambda_{\rho}$ is arbitrary, $\Lambda_{\rho}$ contains only finitely many elements and we have already proved the theorem for the case of an $L^{\infty}$ potential $V$, this suffices to prove the smooth version of the theorem.

The proof of the smooth case of Theorem 1.3 is almost identical, given our earlier results.
Proof of Theorem 1.3 for $V \in C_{c}^{\infty}(X)$. Recall that we have already proved the $L^{\infty}$ case of this theorem. Thus, the proof follows just as in the proof of the smooth case of Theorem 1.2, except that estimate (9-1) is replaced by

$$
\left\|\left.V^{\#} \Xi\left(R_{V_{0}, 0}, 0\right)\right|_{\lambda=\tau_{l}(\zeta)} \chi \mathcal{P}_{l}\right\|_{L^{2}(X) \rightarrow H_{(0, s)}(X)} \leq \frac{C}{\left|\tau_{l}(\zeta)\right|^{r}}
$$

## 10. Proofs of Theorems 1.6 and 1.7

We prove Theorems 1.6 and 1.7 in this section, using results of Section 8.
Before turning to the proofs of the theorems, we say something more about the function $u$ of (1-3). The mapping properties of the resolvent mean that for any $\epsilon>0$ away from its poles, we have the map

$$
R_{V_{0}, 0}(\lambda): e^{-(\epsilon+\max (0,-\operatorname{Im} \lambda))|x|} L^{2}\left(\mathbb{R}^{d}\right) \rightarrow e^{(\epsilon+\max (0,-\operatorname{Im} \lambda))|x|} L^{2}\left(\mathbb{R}^{d}\right)
$$

With $R_{V_{0}, 0}(\lambda)^{t}$ denoting the transpose, we have the symmetry $R_{V_{0}, 0}(\lambda)^{t}=R_{V_{0}, 0}(\lambda)$, checked first for $\operatorname{Im} \lambda>0$ and then holding by analytic continuation for all $\lambda$. This implies that if $R_{V_{0}, 0}(\lambda)$ has a simple pole of rank 1 at $\lambda_{0}$, then there is a $u \in e^{(\epsilon+\max (0,-\operatorname{Im} \lambda))|x|} L^{2}\left(\mathbb{R}^{d}\right)$ so that (1-3) holds, where the operator $u \otimes u$ is understood as an operator between weighted $L^{2}$ spaces.

Now we turn more directly to the proofs, beginning with a preliminary lemma.
Lemma 10.1. Let $\lambda_{0}$ be a pole of $R_{V_{0}, 0}$ and set $R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$. Let $\chi \in C_{c}^{\infty}(X)$ be independent of $\theta$ and satisfy $\chi V=V$, with $\chi$ nontrivial. Suppose $R_{V_{0}, 0}(\lambda)$ is analytic for $0<\left|\lambda-\lambda_{0}\right| \leq \epsilon$. Then there is an $L>0$ so that for $l>L$, if $\zeta_{0} \in D_{l}\left(\lambda_{0}, \epsilon\right)$, then

$$
M\left(I+V^{\#} R_{V_{0}}(\zeta) \chi, \zeta_{0}\right)=M\left(I+\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right) \upharpoonright_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}, \zeta_{0}\right)
$$

Proof. By Lemma 5.2, there is an $L>0$ so that $I+V^{\#} R_{V_{0}}^{\text {reg }}(\zeta) \chi$ is invertible on $D_{l}\left(\lambda_{0}, \epsilon\right)$ for $l>L$. Then if $l>L$ and $\zeta_{0} \in D_{l}\left(\lambda_{0}, \epsilon\right)$,

$$
\begin{aligned}
M\left(I+V^{\#} R_{V_{0}} \chi, \zeta_{0}\right) & =M\left(\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)\left(I+\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right) \upharpoonright_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}\right), \zeta_{0}\right) \\
& =M\left(I+\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}(\lambda), \lambda_{0}\right) \upharpoonright_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}, \zeta_{0}\right)
\end{aligned}
$$

where the second equality uses Lemma 4.1.
Given $f \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, define $h_{ \pm l} \in C_{c}^{\infty}(X)$ by $h_{ \pm l}(x, \theta)=f(x) e^{ \pm i l \theta} / \sqrt{2 \pi}$. For $z_{0} \in \mathbb{C}$ and an operator $A: H_{(0, s)}(X) \rightarrow L^{2}(X)$ set

$$
\begin{equation*}
\mathcal{D}_{A}(z)=\operatorname{det}\left(I+\frac{i}{z-z_{0}}\left(A h_{l} \otimes h_{-l}+A h_{-l} \otimes h_{l}\right)\right) \tag{10-1}
\end{equation*}
$$

Here "det" is the Fredholm determinant. In this special case it is easily calculated to be

$$
\begin{equation*}
\mathcal{D}_{A}(z)=\frac{1}{\left(z-z_{0}\right)^{2}}\left\{\left(z-z_{0}+i \int_{X} h_{-l}\left(A h_{l}\right)\right)\left(z-z_{0}+i \int_{X} h_{l}\left(A h_{-l}\right)\right)+\int_{X} h_{-l}\left(A h_{-l}\right) \int_{X} h_{l}\left(A h_{l}\right)\right\} \tag{10-2}
\end{equation*}
$$

Proposition 10.2. Let $z_{0} \in \mathbb{C}, \epsilon>0$, and set $U_{\epsilon}=\left\{z \in \mathbb{C}:\left|z-z_{0}\right|<\epsilon\right\}$. Suppose there are $L_{1}, m_{0} \geq \frac{1}{2}$ and $s \in \mathbb{N}$ so that for $l>L_{1}, l \in \mathbb{N}$ and $z \in U_{\epsilon}$ there are linear operators $S_{l}=S_{l}(z)$ and $T_{l}=T_{l}(z)$ mapping $H_{(0, s)}(X)$ to $L^{2}(X)$ which are operator-valued functions analytic on $U_{\epsilon}$ satisfying:

- $\sup _{z \in U_{\epsilon}}\left\|\mathcal{P}_{l} S_{l}(z) \mathcal{P}_{l}-T_{l}(z) \mathcal{P}_{l}\right\|_{H_{(0, s)}(X) \rightarrow L^{2}(X)}=O\left(l^{-m_{0}}\right)$,
- $T_{l}(z) \mathcal{P}_{l}=\mathcal{P}_{l+} T_{l}(z) \mathcal{P}_{l+}+\mathcal{P}_{l-} T_{l}(z) \mathcal{P}_{l-}$ and $\sup _{z \in U_{\epsilon}}\left\|T_{l}(z)\right\|_{H_{(0, s)}(X) \rightarrow L^{2}(X)}=O\left(l^{-1 / 2}\right)$.

Then given $f \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, for $l$ sufficiently large the functions $\left(z-z_{0}\right)^{2} \mathcal{D}_{S_{l}}(z)$ and $\left(z-z_{0}\right)^{2} \mathcal{D}_{T_{l}}(z)$ have exactly two zeros, counted with multiplicity, in $U_{\epsilon}$, and they lie in $U_{\epsilon / 2}$. Moreover, there is a labeling of these two sets of zeros as $z_{S_{l} \pm}$ and $z_{T_{l} \pm}$, so that $\left|z_{S_{l} \pm}-z_{T_{l} \pm}\right|=O\left(l^{-m_{0}}\right)$.

Proof. By translating if necessary, we may assume $z_{0}=0$.
Our assumptions on $T_{l}$ imply that $F_{ \pm}(z)=F_{ \pm}(z ; l) \stackrel{\text { def }}{=} z+i \int_{X} h_{\mp l}\left(T_{l}(z) h_{ \pm l}\right)$ is analytic on $U_{\epsilon}$ and satisfies $F_{ \pm}(z)=z+O\left(l^{-1 / 2}\right)$ uniformly on $U_{\epsilon}$. Applying Rouché's theorem to the pair $F_{ \pm}(z)$ and the function $z$, we see that $F_{ \pm}$has, for $l$ sufficiently large, exactly one zero in the set $U_{\epsilon / 4}$ and no zeros in $U_{\epsilon} \backslash U_{\epsilon / 4}$. We label this zero as $z_{T_{l} \pm}$. Since $\int_{X} h_{ \pm l}\left(T_{l} h_{ \pm l}\right)=0$, we have that $z^{2} \mathcal{D}_{T_{l}}(z)=F_{+}(z) F_{-}(z)$ and $z_{T_{l} \pm}$ are the zeros of $z^{2} \mathcal{D}_{T_{l}}$.

We write

$$
\begin{equation*}
F_{ \pm}(z ; l)=z+i \int_{X} h_{\mp l}\left(T_{l}(z) h_{ \pm l}\right)=\left(z-z_{T_{l} \pm}\right) \varphi_{ \pm}(z ; l) \tag{10-3}
\end{equation*}
$$

with $\varphi_{ \pm}$analytic on $U_{\epsilon}$ for $l$ sufficiently large. An application of the maximum principle shows that there is a $C>0$ independent of $l$ so that for $l$ sufficiently large,

$$
\begin{equation*}
\frac{1}{C} \leq\left|\varphi_{ \pm}(z ; l)\right| \leq C \quad \text { for all } z \in U_{3 \epsilon / 4} \tag{10-4}
\end{equation*}
$$

Next consider the intermediary

$$
G_{ \pm}(z)=G_{ \pm}(z ; l) \stackrel{\text { def }}{=} z+i \int_{X} h_{\mp l}\left(S_{l}(z) h_{ \pm l}\right)=z+i \int_{X} h_{\mp l}\left(T_{l}(z) h_{ \pm l}\right)+O\left(l^{-m_{0}}\right)
$$

Our estimate $G_{ \pm}-F_{ \pm}=O\left(l^{-m_{0}}\right)$, (10-3), and (10-4) allow an application of Rouché's theorem to the pair $F_{ \pm}, G_{ \pm}$on a disk with center $z_{T_{l} \pm}$ and radius $c_{0} l^{-m_{0}}$ for an appropriate choice of $c_{0}>0$ and for $l$ sufficiently large. This shows that for $l$ sufficiently large, $G_{ \pm}$has exactly one zero (counting multiplicity) in $U_{\epsilon / 3}$. We label this zero $z_{I, l, \pm}$ (the " $I$ " here stands for intermediate, as this is an intermediate step). We have shown $\left|z_{I, l, \pm}-z_{T_{l}, \pm}\right|=O\left(l^{-m_{0}}\right)$. As before, by the maximum principle we may write

$$
\begin{equation*}
G_{ \pm}(z ; l)=\left(z-z_{I, l, \pm}\right) \varphi_{I \pm}(z ; l), \quad \text { with } \frac{1}{C} \leq\left|\varphi_{I \pm}(z ; l)\right| \leq C, \quad \text { for all } z \in U_{3 \epsilon / 4} \tag{10-5}
\end{equation*}
$$

for some constant $C$ independent of $l$, and for $l$ sufficiently large.
Now consider $z^{2} \mathcal{D}_{S_{l}}(z)$. By our assumptions on $S_{l}$ and $T_{l}$,

$$
z^{2} \mathcal{D}_{S_{l}}(z)=G_{+}(z) G_{-}(z)+O\left(l^{-2 m_{0}}\right)=\left(z-z_{I, l,+}\right)\left(z-z_{I, l,-}\right) \varphi_{I+}(z) \varphi_{I-}(z)+O\left(l^{-2 m_{0}}\right)
$$

Thus we can apply Rouché's theorem again, this time to the pair $z^{2} \mathcal{D}_{S_{l}}(z)$ and $G_{+}(z ; l) G_{-}(z ; l)$ at a distance proportional to $l^{-m_{0}}$ of $z_{I, l, \pm}$, proving the proposition.

We apply this proposition in the proof of Theorem 1.6.
Proof of Theorem 1.6. We assume that $V^{\#} \not \equiv 0$, since otherwise there is nothing to prove. Choose $\chi \in C_{c}^{\infty}(X)$ with $\chi V=V$, and $\chi$ independent of $\theta$.

Let $R_{V_{0}}^{\mathrm{reg}}(\zeta)=R_{V_{0}}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$, and let $\epsilon, L>0$ be as in Lemma 10.1. For $l>L$ the function

$$
F_{l}(\zeta) \stackrel{\text { def }}{=}\left(\tau_{l}(\zeta)-\lambda_{0}\right)^{2} \operatorname{det}\left(I+\left.\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{0}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}\right)
$$

is analytic on $D_{l}\left(\lambda_{0}, \epsilon\right)$. Moreover, the order of vanishing of $F_{l}$ at $\zeta_{0} \in D_{l}\left(\lambda_{0}, \epsilon\right)$ is given by

$$
M\left(I+\left.\left(I+V^{\#} R_{V_{0}}^{\mathrm{reg}}(\zeta) \chi\right)^{-1} V^{\#} \Xi\left(R_{V_{0}, 0}, \lambda_{0}\right)\right|_{\lambda=\tau_{l}(\zeta)} \mathcal{P}_{l}, \zeta_{0}\right)+m_{V_{0}}\left(\zeta_{0}\right)
$$

see [Gohberg and Sigal 1971, Theorem 5.1]. Note that for $\zeta_{0} \in D_{l}\left(\lambda_{0}, \epsilon\right)$ and $l$ sufficiently large, $m_{V_{0}}\left(\zeta_{0}\right) \neq 0$ if and only if $\tau_{l}\left(\zeta_{0}\right)=\lambda_{0}$. For $\lambda_{0} \neq 0$, combining this with Lemmas 10.1 and 4.4 , we see that the poles of $R_{V}$ in $D_{l}\left(\lambda_{0}, \epsilon\right)$ are, for $l>L$, given by the zeros of $F_{l}$, and the multiplicities agree. If $\lambda_{0}=0$, the same is true, but as in the proof of Theorem 1.3 we use Lemmas 5.6, 5.7, and 5.9.

To prove the theorem, we will apply Proposition 10.2 with the following choices: $z=\tau_{l}(\zeta), z_{0}=\lambda_{0}$, $f(x)=\chi(x) u(x)$ so that $h_{ \pm l}(x, \theta)=\chi(x) u(x) e^{ \pm i l \theta} / \sqrt{2 \pi}$,

$$
\begin{aligned}
& S_{l}=S_{l}(z) \\
&=\left(I+V^{\#} R_{V}^{\mathrm{reg}}\left(\zeta_{l}(z)\right)\right)^{-1} V^{\#} \mathcal{P}_{l} \\
& T_{l}=T_{l}(z)=\frac{-1}{l^{2}} \sum_{k \neq 0}\left(\frac{z^{2}-k^{2}}{4 k^{2}} V_{-k} V_{k}-\frac{1}{4 k^{2}} V_{-k}\left(-\Delta_{0}+V_{0}\right) V_{k}\right) \mathcal{P}_{l}
\end{aligned}
$$

and $s=8$. By Proposition 8.2 we have, in the notation of Proposition 10.2, $m_{0}=3$. Note that using the coordinate $z=\tau_{l}(\zeta)$, we have $F_{l}\left(\zeta_{l}(z)\right)=\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{S_{l}}(z)$, where $\mathcal{D}_{S_{l}}$ is as defined via (10-1).

The function $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{T_{l}}(z)$ has a single zero of multiplicity 2 in $U_{\epsilon}$, and by Lemma 8.7 this is the zero of

$$
z-\lambda_{0}+\frac{i}{4 l^{2}} \sum_{k \neq 0} \int_{\mathbb{R}}\left(\frac{k^{2}+\lambda_{0}^{2}-z^{2}}{k^{2}} u^{2} V_{-k} V_{k}+\frac{u^{2} \nabla_{0} V_{-k} \cdot \nabla_{0} V_{k}}{k^{2}}\right)
$$

near $z=\lambda_{0}$. This zero is given by

$$
z_{T_{l} \pm}=\lambda_{0}-\frac{i}{4 l^{2}} \sum_{k \neq 0} \int_{\mathbb{R}}\left(u^{2} V_{-k} V_{k}+\frac{u^{2} \nabla_{0} V_{-k} \cdot \nabla_{0} V_{k}}{k^{2}}\right)+O\left(l^{-4}\right)
$$

By Proposition 10.2, the zeros of $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{S_{l}}(z)$ in $U_{\epsilon}$ are within $O\left(l^{-m_{0}}\right)=O\left(l^{-3}\right)$ of the zero (of multiplicity 2) of $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{T_{l}}(z)$ in $U_{\epsilon}$, thus completing the proof.

The proof of Theorem 1.7 is similar.
Proof of Theorem 1.7. We prove the theorem by showing that for any $N \in \mathbb{N}$ there is an $\epsilon>0$ so that for $l \in \mathbb{N}$ sufficiently large if $\zeta_{l}^{b} \in D_{l}\left(\lambda_{0}, \epsilon\right)$ and $\zeta_{l}^{b}$ is a pole of $R_{V}(\zeta)$, then $\operatorname{Re} \tau_{l}\left(\zeta_{l}^{b}\right)=O\left(l^{-N}\right)$.

Choose $\chi \in C_{c}^{\infty}(X ; \mathbb{R})$ so that $\chi V=V$ and $\chi$ is independent of $\theta$. Choose $\epsilon, L>0$, as in Lemma 10.1. Let $u \in C^{\infty}\left(\mathbb{R}^{d}\right)$ be such that $R_{V_{0}, 0}(\lambda)-i /\left(\lambda-\lambda_{0}\right) u \otimes u$ is analytic for $\lambda$ near $\lambda_{0}$. Our assumptions on $V$ and $\lambda_{0}$ imply that $u$ is real-valued. We apply Proposition 10.2 in a way very similar to the proof of Theorem 1.6. We make the following choices: $z=\tau_{l}(\zeta), z_{0}=\lambda_{0}, h_{ \pm l}(x, \theta)=\chi(x) u(x) e^{ \pm i l \theta} / \sqrt{2 \pi}$, and $S_{l}=S_{l}(z)=\left(I+V^{\#} R_{V}^{\mathrm{reg}}\left(\zeta_{l}(z)\right) \chi\right)^{-1} V^{\#} \mathcal{P}_{l}$, where $R_{V}^{\mathrm{reg}}(\zeta)=R_{V}^{\mathrm{reg}}\left(\zeta ; \lambda_{0}, l\right)$. For $l$ sufficiently large, $S_{l}$ is analytic on $U_{\epsilon}$. Let $A_{N}=A_{N}(z, l)$ be the operator from Lemma 8.8, and set

$$
T_{l}=T_{l}(z ; N)=\mathcal{P}_{l+} A_{N} \mathcal{P}_{l+}+\mathcal{P}_{l-} A_{N} \mathcal{P}_{l-}
$$

By Lemma 8.8 , there is an $s \in \mathbb{N}$ so that

$$
\left\|\mathcal{P}_{l} S_{l}(z) \mathcal{P}_{l}-T_{l}(z)\right\|_{H_{(0, s)}(X) \rightarrow L^{2}(X)}=O\left(l^{-N}\right)
$$

uniformly for $z \in U_{\epsilon}$. Thus for our application of Proposition 10.2 we have $m_{0}=N$.
Following the proof of Theorem 1.6, the poles of $R_{V}$ in $D_{l}\left(\lambda_{0}, \epsilon\right)$ are determined by the zeros of $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{S_{l}}(z)$ in $U_{\epsilon}$, using $U_{\epsilon} \ni z=\tau_{l}(\zeta)$. By Proposition 10.2, these zeros are approximated by those
of $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{T_{l}}(z)$ in $U_{\epsilon}$, with an error which is $O\left(l^{-N}\right)$. We compete the proof by showing that for $l$ sufficiently large the zeros of $\mathcal{D}_{T_{l}}(z)$ in $U_{\epsilon}$ lie on the imaginary axis.

Set $a_{ \pm}(z ; l) \stackrel{\text { def }}{=} \int_{X} h_{\mp l}\left(T_{l}(z) h_{ \pm l}\right)=\int_{X} \overline{h_{ \pm l}}\left(T_{l}(z) h_{ \pm l}\right)$. From Lemma 8.8 and the definition of $T_{l}$, if $z \in U_{\epsilon} \cap i \mathbb{R}$, then $T_{l}(i z)$ is symmetric on $C_{c}^{\infty}(X) \subset L^{2}(X)$. In particular, this implies that if $z \in i \mathbb{R} \cap U_{\epsilon}$ then $a_{ \pm}(z ; l) \in \mathbb{R}$. Since $a_{ \pm}(z ; l)$ is analytic for $z \in U_{\epsilon}$ and is real-valued for $z \in i \mathbb{R} \cap U_{\epsilon}$, we must have

$$
\begin{equation*}
a_{ \pm}(z ; l)=\bar{a}_{ \pm}(-\bar{z} ; l) \quad \text { for } z \in U_{\epsilon} \tag{10-6}
\end{equation*}
$$

We remark that since $\lambda_{0} \in i \mathbb{R}$, we have $z \in U_{\epsilon}$ if and only if $-\bar{z} \in U_{\epsilon}$.
From the proof of Proposition 10.2, the zeros of $\left(z-\lambda_{0}\right)^{2} \mathcal{D}_{T_{l}}(z)$ in $U_{\epsilon}$ are given by the zeros of $z-\lambda_{0}+i a_{ \pm}(z, l)$ in $U_{\epsilon}$, and there is, for $l$ sufficiently large, exactly one such zero for each choice of $\pm$. We denote these zeros by $z_{T_{l} \pm}$ and focus on the zero for the " + " sign, $z_{T_{l}+}$. Using $\lambda_{0} \in i \mathbb{R}$,

$$
\begin{aligned}
z_{T_{l}+}-\lambda_{0}+i a_{+}\left(z_{T_{l}+} ; l\right)=0=\overline{z_{T_{l}+}-\lambda_{0}+i a_{+}\left(z_{T_{l}+} ; l\right)} & =-\left(-\overline{z_{T_{l}+}}-\lambda_{0}+i \bar{a}_{+}\left(z_{T_{l}+} ; l\right)\right) \\
& =-\left(-\overline{z_{T_{l}+}}-\lambda_{0}+i a_{+}\left(-\overline{z_{T_{l}+}} ; l\right)\right)
\end{aligned}
$$

where the last equality uses (10-6). Hence $-\overline{z_{T_{l}+}}$ is also a zero of $z-\lambda_{0}+i a_{+}(z ; l)$ in $U_{\epsilon}$, and since there is exactly one such zero, it must be that $-\overline{z_{T_{l}+}}=z_{T_{l}+}$, and thus $z_{T_{l}+} \in i \mathbb{R}$. The same argument shows $z_{T_{l}-} \in i \mathbb{R}$.

## 11. Proof of Theorem 1.8 , the resonant uniqueness of $V \equiv 0$ when $d=1$

Theorem 1.8 , a result on the resonant rigidity of the zero potential on $\mathbb{R} \times \mathbb{S}^{1}$, follows rather directly from Theorems 1.1, 1.3, and 1.6.
Proof of Theorem 1.8. Suppose $X=\mathbb{R} \times \mathbb{S}^{1}$ and $V$ is as in Theorem 1.8. Then by Theorems 1.1 and 1.3, the one-dimensional operator $-d^{2} / d x^{2}+V_{0}$ on $\mathbb{R}$ must have a resonance at the origin and nowhere else, and this resonance must have multiplicity 1 . But since $V_{0} \in L_{c}^{\infty}(\mathbb{R})$, by well-known results for one-dimensional Schrödinger operators, $V_{0} \equiv 0$; see for example [Zworski 1987].

The operator $R_{0,0}(\lambda)-i /(2 \lambda) 1 \otimes 1$ is analytic at the origin. Using this in Theorem 1.6 along with the fact that $R_{V}$ has poles at a sequence of thresholds tending to infinity, we find

$$
\sum_{k \neq 0} \frac{1}{k^{2}} \int_{\mathbb{R}}\left(k^{2} V_{k} V_{-k}+V_{k}^{\prime} V_{-k}^{\prime}\right)(x) d x=0
$$

But since $V_{-k}(x)=\bar{V}_{k}(x)$ for a real-valued potential $V$, this implies $V_{k} \equiv 0$ for all $k$, and hence $V \equiv 0$.

## 12. The potential $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$ on $\mathbb{R} \times \mathbb{S}^{1}$

In this section we investigate the resonances near the $l$-th threshold of the Schrödinger operator with potential $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$ on $X=\mathbb{R} \times \mathbb{S}^{1}$. Here $\chi_{I_{0}}(x)$ is the characteristic function of the interval $I_{0}=[-1,1]$, so $\chi_{I_{0}}(x)=1$ if $|x| \leq 1$ and $\chi_{I_{0}}(x)=0$ if $|x|>1$. This potential has $V_{0} \equiv 0$ so that $V^{\#}=V$. Proposition 12.3 shows that the resonances nearest the threshold, which correspond to perturbations of the pole at the origin for $R_{0,0}(\lambda)$, are, for this potential, localized in a different way than for smooth potentials;
compare Theorem 1.6. By Proposition 12.6, there is a sense in which Theorem 1.5 is sharp. We remark that some of the computations of this section are reminiscent of those found in [Drouot 2018, Section 2].

In all of this section,

$$
V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta \quad \text { and } \quad X=\mathbb{R} \times \mathbb{S}^{1}
$$

We will use this preliminary lemma.
Lemma 12.1. For $\lambda, \lambda^{\prime} \in \mathbb{C}, \lambda \neq \pm \lambda^{\prime}$,

$$
\begin{align*}
& \chi_{I_{0}} R_{0,0}(\lambda) \chi_{I_{0}} R_{0,0}\left(\lambda^{\prime}\right) \chi_{I_{0}} \\
& \quad=\frac{1}{\left(\lambda^{\prime}\right)^{2}-\lambda^{2}} \chi_{I_{0}}\left(R_{0,0}\left(\lambda^{\prime}\right)-R_{0,0}(\lambda)\right) \chi_{I_{0}}+\frac{i}{4 \lambda \lambda^{\prime}\left(\lambda+\lambda^{\prime}\right)} e^{i\left(\lambda+\lambda^{\prime}\right)}\left(\phi_{\lambda} \otimes \phi_{\lambda^{\prime}}+\phi_{-\lambda} \otimes \phi_{-\lambda^{\prime}}\right) \tag{12-1}
\end{align*}
$$

where

$$
\phi_{ \pm \lambda}(x)=e^{ \pm i \lambda x} \chi_{I_{0}}(x)
$$

Moreover, if $\tau \in \mathbb{C}, \tau \neq \pm \lambda$, applying the operator $\chi_{I_{0}} R_{0,0}(\tau)$ to the function $\chi_{I_{0}}(x) e^{i \lambda x}$ yields

$$
\begin{equation*}
\left(\chi_{I_{0}} R_{0,0}(\tau) \chi_{I_{0}} e^{i \lambda \star}\right)(x)=\chi_{I_{0}}(x)\left(\frac{1}{\lambda^{2}-\tau^{2}} e^{i \lambda x}+\frac{1}{2 \tau(\lambda-\tau)} e^{-i \lambda} e^{i \tau(1+x)}+\frac{1}{2 \tau(\tau+\lambda)} e^{i \lambda} e^{i \tau(1-x)}\right) . \tag{12-2}
\end{equation*}
$$

Proof. The first can be seen, for example, by using (3-1), the explicit expression for the Schwartz kernel of $R_{0,0}$, and evaluating

$$
\int_{-1}^{1} e^{i \lambda\left|x-x^{\prime \prime}\right|+i \lambda^{\prime}\left|x^{\prime \prime}-x^{\prime}\right|} d x^{\prime \prime}
$$

for $|x|,\left|x^{\prime}\right| \leq 1$. Likewise, (12-2) follows from an explicit computation using (3-1).
12A. Resonances near the threshold $\tau_{l}=0$ for $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$. Since in this section we concentrate on the resonance near the threshold, we work on $B_{l}(1)$. A preliminary step is the following.
Lemma 12.2. Let $R_{0}^{\mathrm{reg}}(\zeta)=R_{0}^{\mathrm{reg}}(\zeta ; 0, l)$. Then for $l$ sufficiently large, uniformly on $B_{l}(1)$,

$$
\left\|\mathcal{P}_{l}\left(\left(I+V R_{0}^{\mathrm{reg}}(\zeta) \chi_{I_{0}}\right)^{-1} V+V R_{0}^{\mathrm{reg}}(\zeta) V+\left(V R_{0}^{\mathrm{reg}}(\zeta)\right)^{3} V\right) \mathcal{P}_{l}\right\|=O\left(l^{-2}\right)
$$

Proof. Using the Neumann series,

$$
\left(I+V R_{0}^{\mathrm{reg}}(\zeta) \chi_{I_{0}}\right)^{-1} V=\sum_{j=0}^{\infty}\left(-V R_{0}^{\mathrm{reg}}(\zeta)\right)^{j} V
$$

By Lemma 5.2, $\left\|\left(-V R_{0}^{\mathrm{reg}}(\zeta)\right)^{j}\right\|=O\left(l^{-2}\right)$ on $B_{l}(1)$ if $j \geq 4$ and $l$ is sufficiently large. This ensures the Neumann series for $\left(I+V R_{0}^{\text {reg }}(\zeta) \chi_{I_{0}}\right)^{-1}$ converges, and

$$
\left\|\left(I+V R_{0}^{\mathrm{reg}}(\zeta) \chi_{I_{0}}\right)^{-1} V-\sum_{j=0}^{3}\left(-V R_{0}^{\mathrm{reg}}(\zeta)\right)^{j} V\right\|=O\left(l^{-2}\right)
$$

on $B_{l}(1)$.
Now we note that our explicit expression for $V$ means that $\mathcal{P}_{l} V \mathcal{P}_{l}=0$. Likewise, it implies that $\mathcal{P}_{l}\left(V R_{0}^{\mathrm{reg}}(\zeta)\right)^{2} V \mathcal{P}_{l}=0$, completing the proof.

Proposition 12.3. Forl sufficiently large, the poles of $R_{V}(\zeta)$ in $B_{l}(1)$ satisfy

$$
\tau_{l}(\zeta)=\frac{1}{4 l \sqrt{2 l}}\left(-1-i+e^{i 2 \sqrt{2 l}}\right)+O\left(l^{-2}\right)
$$

Proof. We give a proof similar to that of Theorem 1.6 using Proposition 10.2.
Let $R_{0}^{\text {reg }}$ be as in Lemma 12.2, and restrict $\zeta$ to $\zeta \in B_{l}(1)$. Note

$$
R_{0,0}(\lambda)-\frac{i}{2 \lambda} 1 \otimes 1
$$

is regular at $\lambda=0$. Set $z=\tau_{l}(\zeta)$,

$$
S_{l}(z)=\left(I+V R_{0}^{\mathrm{reg}}\left(\zeta_{l}(z)\right) \chi_{I_{0}}\right)^{-1} V \mathcal{P}_{l}, \quad \text { and } \quad h_{ \pm l}(x, \theta)=\frac{1}{\sqrt{2 \pi}} \chi_{I_{0}}(x) e^{ \pm i l \theta}
$$

We use $\mathcal{D}_{S_{l}}$ as is defined by (10-1) and $U_{\epsilon}$ as in Proposition 10.2. Then just as in the proof of Theorem 1.6, the poles of $R_{V}$ in $B_{l}(1)$ are identified via $z=\tau_{l}(\zeta)$ with the zeros of $z^{2} \mathcal{D}_{S_{l}}(z)$ in $U_{1}$. Set $z_{0}=0$ and $T_{l}=\mathcal{P}_{l}\left(-V R_{0}^{\mathrm{reg}}(\zeta) V-\left(V R_{0}^{\mathrm{reg}}(\zeta)\right)^{3} V\right) \mathcal{P}_{l}$. Then by Lemma 12.2, in our application of Proposition 10.2 we can take $s=0$ and $m_{0}=2$. We claim that uniformly for $z \in U_{1}$,

$$
\begin{equation*}
z^{2} \mathcal{D}_{T_{l}}(z)=\left(z+\frac{1}{2(2 l)^{3 / 2}}\left(1-e^{2 i \sqrt{2 l}}+i\right)+O\left(l^{-2}\right)\right)^{2} \tag{12-3}
\end{equation*}
$$

Assuming for the moment that (12-3) holds, this shows that the two zeros (when counted with multiplicity) of $z^{2} \mathcal{D}_{T_{l}}(z)$ in $U_{1}$ satisfy

$$
z=\frac{-1-i+e^{2 i \sqrt{2 l}}}{2(2 l)^{3 / 2}}+O\left(l^{-2}\right)
$$

An application of Proposition 10.2 and Lemma 12.2 then proves the proposition.
We now turn to showing (12-3). We use

$$
\begin{equation*}
R_{0}^{\mathrm{reg}}\left(\zeta_{l}(z)\right) V \mathcal{P}_{l}=\sum_{ \pm}\left(e^{ \pm i \theta} R_{0,0}\left(\tau_{l+1}\right)+e^{\mp i \theta} R_{0,0}\left(\tau_{l-1}\right)\right) \chi_{I_{0}} \mathcal{P}_{l \pm} \tag{12-4}
\end{equation*}
$$

where $\tau_{l \pm 1}=\tau_{l \pm 1}\left(\zeta_{l}(z)\right)$, so that

$$
\begin{equation*}
\mathcal{P}_{l} V R_{0}^{\mathrm{reg}}\left(\zeta_{l}(z)\right) V \mathcal{P}_{l}=\chi_{I_{0}}\left(R_{0,0}\left(\tau_{l-1}\right)+R_{0,0}\left(\tau_{l+1}\right)\right) \chi_{I_{0}} \mathcal{P}_{l} \tag{12-5}
\end{equation*}
$$

Then using (12-2) gives

$$
\begin{equation*}
\int_{X} h_{\mp l} V R_{0}^{\mathrm{reg}}\left(\zeta_{l}(z)\right) V h_{ \pm l}=\frac{-i}{2(2 l)^{3 / 2}}\left(1-e^{2 i \sqrt{2 l}}\right)+\frac{1}{2(2 l)^{3 / 2}}+O\left(l^{-2}\right) \tag{12-6}
\end{equation*}
$$

uniformly on $U_{1}$. Now note

$$
\begin{equation*}
\int_{X} h_{\mp l}\left(V R_{0}^{\mathrm{reg}}\right)^{3} V h_{ \pm l}=\int_{X}\left(V R_{0}^{\mathrm{reg}} V h_{\mp l}\right)\left(\chi_{I_{0}}\left(R_{0}^{\mathrm{reg}} V\right)^{2} h_{ \pm l}\right) \tag{12-7}
\end{equation*}
$$

By (12-2),

$$
\left\|\left(V R_{0}^{\mathrm{reg}} V h_{\mp l}\right)\right\|=O\left(l^{-1}\right) \quad \text { and } \quad\left\|\chi_{I_{0}}\left(R_{0}^{\mathrm{reg}} V\right)^{2} h_{ \pm l}\right\|=O\left(l^{-1}\right)
$$

Using the expression for $\mathcal{D}_{T_{l}}$ as in (10-2) and equations (12-5)-(12-7) completes the proof of (12-3).

12B. Existence of poles of $R_{V}$ within $\approx \log l$ of the $l$-th threshold, for $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$. As a point of comparison with Theorem 1.5 , for the special case $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$ on $X=\mathbb{R} \times \mathbb{S}^{1}$ we consider the existence of poles of $R_{V}(\zeta)$ in $D_{l}(\alpha \log l)$ with $\left|\tau_{l}(\zeta)\right|>1$.

Again, we use the coordinate $z=\tau_{l}(\zeta)$ on $B_{l}(\alpha \log l)$, and the functions $\phi_{\lambda}$ are as defined in Lemma 12.1.
Lemma 12.4. Let $\alpha>0$ be fixed, and set $z=\tau_{l}(\zeta)$. For $l$ sufficiently large, uniformly on $B_{l}(\alpha \log l) \backslash B_{l}(1)$ we have

$$
\left\|\mathcal{P}_{l}\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1} V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}+\left(f_{+} \otimes \phi_{z}+f_{-} \otimes \phi_{-z}\right) \mathcal{P}_{l}-\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}\right\| x .
$$

where

$$
f_{ \pm}(x)=f_{ \pm}(x, z, l)=\frac{i e^{i z}}{4 z} \chi_{I_{0}}(x)\left(\frac{e^{i \tau_{l+1}}}{\tau_{l+1}\left(z+\tau_{l+1}\right)} \phi_{ \pm \tau_{l+1}}+\frac{e^{i \tau_{l-1}}}{\tau_{l-1}\left(z+\tau_{l-1}\right)} \phi_{ \pm \tau_{l-1}}\right)
$$

For notational simplicity, we have written $\tau_{l \pm 1}$ for $\tau_{l \pm 1}\left(\zeta_{l}(z)\right)$.
Proof. We use

$$
\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1}=\sum_{j=0}^{\infty}\left(-V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{j}
$$

since $\left\|V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right\|=O\left(l^{-1 / 2}\right)$. This estimate, along with others in this proof, are uniform for $\zeta \in B_{l}(\alpha \log l) \backslash B_{l}(1)$. By Lemma 12.1, (3-1), and the explicit expression for $V$, we see that

$$
\left\|\chi_{I_{0}} R_{0}(\zeta)\left(I-\mathcal{P}_{l}\right) V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}\right\|=O\left(e^{2(\operatorname{Im} z)_{-}-} /(l|z|)\right) \quad \text { for } \zeta \in B_{l}(\alpha \log l)
$$

for $l$ sufficiently large. Moreover, this same lemma implies that if $|j-l| \leq 2$, then

$$
\left\|\chi_{I_{0}}\left(V R_{0}(\zeta)\left(I-\mathcal{P}_{l}\right)\right)^{2} \chi_{I_{0}} \mathcal{P}_{j}\right\|=O\left(l^{-3 / 2}\right)
$$

uniformly on $B_{l}(\alpha \log l)$. This ensures that

$$
\begin{align*}
&\left\|\left(\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1}-\sum_{j=0}^{2}\left(-V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{j}\right) V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}\right\| \\
&=O\left(\frac{1}{l^{5 / 2}|z|} e^{2(\operatorname{Im} z)_{-}}\right) \tag{12-9}
\end{align*}
$$

Since, as in the proof of Proposition 12.3, $\mathcal{P}_{l} V \mathcal{P}_{l}=0$ and $\mathcal{P}_{l}\left(V R_{0}\left(I-\mathcal{P}_{l}\right)\right)^{2} V \mathcal{P}_{l}=0$, it suffices to use $-\mathcal{P}_{l} V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right) V R_{0}(\zeta) \mathcal{P}_{l}$ to approximate $\mathcal{P}_{l}\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1} V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}$ with the desired accuracy.

Using Lemma 12.1 and its notation,

$$
\begin{aligned}
& \mathcal{P}_{l} V R_{0}\left(\zeta_{l}(z)\right)\left(I-\mathcal{P}_{l}\right) V R_{0}\left(\zeta_{l}(z)\right) \chi_{I_{0}} \mathcal{P}_{l} \\
& =\chi_{I_{0}}\left(R_{0,0}\left(\tau_{l+1}\right) \chi_{I_{0}} R_{0,0}(z)+R_{0,0}\left(\tau_{l-1}\right) \chi_{I_{0}} R_{0,0}(z)\right) \chi_{I_{0}} \mathcal{P}_{l} \\
& = \\
& \quad \frac{1}{\tau_{l+1}^{2}-z^{2}} \chi_{I_{0}}\left(R_{0,0}\left(\tau_{l+1}\right)-R_{0,0}(z)\right) \chi_{I_{0}} \mathcal{P}_{l}+\frac{i e^{i\left(z+\tau_{l+1}\right)}}{4 z \tau_{l+1}\left(z+\tau_{l+1}\right)}\left(\phi_{\tau_{l+1}} \otimes \phi_{z}+\phi_{-\tau_{l+1}} \otimes \phi_{-z}\right) \mathcal{P}_{l} \\
& \quad \quad \quad+\frac{1}{\tau_{l-1}^{2}-z^{2}} \chi_{I_{0}}\left(R_{0,0}\left(\tau_{l-1}\right)-R_{0,0}(z)\right) \chi_{I_{0}} \mathcal{P}_{l}+\frac{i e^{i\left(z+\tau_{l-1}\right)}}{4 z \tau_{l-1}\left(z+\tau_{l-1}\right)}\left(\phi_{\tau_{l-1}} \otimes \phi_{z}+\phi_{-\tau_{l-1}} \otimes \phi_{-z}\right) \mathcal{P}_{l} .
\end{aligned}
$$

Note that

$$
\left\|\frac{1}{\tau_{l \pm 1}^{2}-z^{2}} \chi_{I_{0}} R_{0,0}\left(\tau_{l \pm 1}\right) \chi_{I_{0}}\right\|=O\left(l^{-3 / 2}\right)
$$

and

$$
\left\|\left(\frac{1}{\tau_{l+1}^{2}-z^{2}}+\frac{1}{\tau_{l-1}^{2}-z^{2}}\right) \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}}-\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}}\right\|=O\left(l^{-4}|z|^{-1} e^{2(\operatorname{Im} z)_{-}}\right)
$$

This gives

$$
\begin{align*}
& \mathcal{P}_{l}\left(V R_{0}\left(\zeta_{l}(z)\right)\left(I-\mathcal{P}_{l}\right) V R_{0,0}\left(\zeta_{l}(z)\right) \chi_{I_{0}} \mathcal{P}_{l}\right. \\
& \begin{aligned}
=\frac{i e^{i\left(z+\tau_{l+1}\right)}}{4 z \tau_{l+1}\left(z+\tau_{l+1}\right)}\left(\phi_{\tau_{l+1}} \otimes \phi_{z}\right. & \left.+\phi_{-\tau_{l+1}} \otimes \phi_{-z}\right) \mathcal{P}_{l}+\frac{i e^{i\left(z+\tau_{l-1}\right)}}{4 z \tau_{l-1}\left(z+\tau_{l-1}\right)}\left(\phi_{\tau_{l-1}} \otimes \phi_{z}+\phi_{-\tau_{l-1}} \otimes \phi_{-z}\right) \mathcal{P}_{l} \\
& \quad-\frac{1}{2 l^{2}} R_{0,0}(z) \mathcal{P}_{l}+O_{L^{2} \rightarrow L^{2}}\left(\frac{1}{l^{5 / 2}|z|} e^{2(\operatorname{Im} z)_{-}}\right)+O_{L^{2} \rightarrow L^{2}}\left(l^{-3 / 2}\right),
\end{aligned}
\end{align*}
$$

and completes the proof.
Note that the functions $f_{ \pm}$and $\phi_{ \pm}$in Lemma 12.4 depend holomorphically on $z$ in the set

$$
\{z \in \mathbb{C}: 1 \leq z \leq \alpha \log l\}
$$

The function $g_{l}$ of the next lemma appears in the proof of Proposition 12.6, as its zeros approximate the locations of the poles of $R_{V}(\zeta)$ away from the threshold in $B_{l}(\alpha \log l)$, if $\alpha<1$. A discussion of the Lambert $W$ function can be found, for example, in [Corless et al. 1996]. This next lemma is very similar to [Drouot 2018, Lemma 2.4].

Lemma 12.5. The zeros of

$$
g_{l}(z) \stackrel{\operatorname{def}}{=}\left(1-\frac{1}{z 8 l \sqrt{2 l}} e^{2 i(\sqrt{2 l}+z)}\right)^{2}-\left(\frac{1}{8 l z \sqrt{2 l}}\left(i e^{2 i z}+e^{2 i z}\right)\right)^{2}
$$

are given by $z_{v}^{ \pm}=z_{v}^{ \pm}(l)=\frac{i}{2} \mathcal{W}_{v}\left(\left(-i e^{2 i \sqrt{2 l}} \mp i \pm 1\right) /(4 l \sqrt{2 l})\right)$, where $\mathcal{W}_{v}$ is the $v-$ th branch of the Lambert $W$ function. In particular, we have $z_{1}^{+} \sim-\frac{3 i}{4} \log l$. Moreover, for $l$ sufficiently large there is an $r_{0}>0$ independent of $l$ so that if $w \in \mathbb{C}$ and $|w|<r_{0}$, then

$$
\begin{equation*}
\left|g_{l}\left(z_{1}^{+}(l)+w\right)\right| \geq \frac{2}{3}|w| \tag{12-11}
\end{equation*}
$$

Proof. The zeros of $g_{l}$ are solutions of

$$
1-\frac{1}{z 8 l \sqrt{2 l}} e^{2 i(\sqrt{2 l}+z)}= \pm \frac{1}{8 l z \sqrt{2 l}}\left(i e^{2 i z}+e^{2 i z}\right)
$$

and so satisfy

$$
z e^{-2 i z}=\frac{1}{8 l \sqrt{2 l}}\left(e^{2 i \sqrt{2 l}} \pm 1 \pm i\right)
$$

Solutions of this equation are given by

$$
z_{v}^{ \pm}=\frac{i}{2} \mathcal{W}_{v}\left(\frac{1}{4 l \sqrt{2 l}}\left(-i e^{2 i \sqrt{2 l}} \mp i \pm 1\right)\right)
$$

From [Corless et al. 1996, (4.20)], we have $z_{1}^{+} \sim-\frac{3 i}{4} \log l$ as $l \rightarrow \infty$.
To finish the proof, we set $\gamma=1 /(8 l \sqrt{2 l})$ and write

$$
g_{l}(z)=\left(1+\frac{\gamma}{z} e^{2 i z}\left(-e^{2 i \sqrt{2 l}}-1-i\right)\right)\left(1+\frac{\gamma}{z} e^{2 i z}\left(-e^{2 i \sqrt{2 l}}+1+i\right)\right)
$$

Now we evaluate at $z=z_{1}^{+}+w$, with $w \in \mathbb{C},|w|$ small, to find

$$
\begin{aligned}
g_{l}\left(z_{1}^{+}+w\right) & =\left(1+\frac{z_{1}^{+} e^{2 i w}}{z_{1}^{+}+w} \frac{\gamma}{z_{1}^{+} e^{-2 i z_{1}^{+}}}\left(-e^{2 i \sqrt{2 l}}-1-i\right)\right)\left(1+\frac{z_{1}^{+} e^{2 i w}}{z_{1}^{+}+w} \frac{\gamma}{z_{1}^{+} e^{-2 i z_{1}^{+}}}\left(-e^{2 i \sqrt{2 l}}+1+i\right)\right) \\
& =\left(1-\frac{z_{1}^{+} e^{2 i w}}{z_{1}^{+}+w}\right)\left(1+\frac{z_{1}^{+} e^{2 i w}}{z_{1}^{+}+w} \frac{-e^{2 i \sqrt{2 l}}+1+i}{e^{2 i \sqrt{2 l}}+1+i}\right),
\end{aligned}
$$

where for the second equality we have used $z_{1}^{+} e^{-2 i z_{1}^{+}}=\gamma\left(e^{2 i \sqrt{2 l}}+1+i\right)$. This gives, then, recalling $\left|z_{1}^{+}\right| \rightarrow \infty$ as $l \rightarrow \infty$,

$$
g_{l}\left(z_{1}^{+}+w\right)=\left(-2 i w+O\left(|w| /\left|z_{1}^{+}\right|\right)+O\left(|w|^{2}\right)\right)\left(\frac{2(i+1)}{e^{2 i \sqrt{2 l}}+1+i}+O(|w|)\right)
$$

for $|w|$ small. Then there is a $r_{0}>0$ independent of $l$ so that for $l$ sufficiently large and $|w|<r_{0}$, $\left|g_{l}\left(z_{1}^{+}+w\right)\right|>\frac{2}{3}|w|$.
Proposition 12.6. For $V(x, \theta)=2 \chi_{I_{0}}(x) \cos \theta$ and $l$ sufficiently large, $R_{V}(\zeta)$ has a pole at a point $\zeta_{l}^{+} \in B_{l}\left(\frac{7}{8} \log l\right)$ with $\zeta_{l}^{+}$satisfying

$$
\tau_{l}\left(\zeta_{l}^{+}\right)=\frac{i}{2} \mathcal{W}_{1}\left(\frac{1}{4 l \sqrt{2 l}}\left(i e^{2 i \sqrt{2 l}}-i+1\right)\right)+O\left(l^{-1 / 2+\epsilon}\right)
$$

for any $\epsilon>0$.
Proof. We continue to use $z=\tau_{l}(\zeta)$ and work in a region with $1<|z|<\frac{7}{8} \log l$.
Using Lemma 12.4,

$$
\mathcal{P}_{l}\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1} V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}=F \mathcal{P}_{l}+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}+A
$$

where, with notation from Lemma 12.4,

$$
F=F(z, l)=-f_{+} \otimes \phi_{z}-f_{-} \otimes \phi_{-z}
$$

and $\|A\|=O\left(l^{-5 / 2} e^{(2 \operatorname{Im} z)_{-}}\right)+O\left(l^{-3 / 2}\right)$ on $B_{l}\left(\frac{7}{8} \log l\right) \backslash B_{1}(l)$. We recall that the poles of $R_{V}$ in $B_{l}\left(\frac{7}{8} \log l\right) \backslash B_{1}(l)$ are the zeros of $I+\mathcal{P}_{l}\left(I+V R_{0}(\zeta) \chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1} V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l}$ in $B_{l}\left(\frac{7}{8} \log l\right) \backslash B_{1}(l)$. We write

$$
\begin{align*}
I+\mathcal{P}_{l}\left(I+V R_{0}(\zeta)\right. & \left.\chi_{I_{0}}\left(I-\mathcal{P}_{l}\right)\right)^{-1} V R_{0}(\zeta) \chi_{I_{0}} \mathcal{P}_{l} \\
& =\left(I+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}\right)\left(I+\left(I+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}\right)^{-1}\left(F \mathcal{P}_{l}+A\right)\right) \tag{12-12}
\end{align*}
$$

since

$$
I+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}
$$

is invertible here. For notational convenience, set

$$
S=S_{l}=\left(I+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}\right)^{-1}
$$

and note that

$$
S=I-\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}} \mathcal{P}_{l}+O_{L^{2} \rightarrow L^{2}}\left(l^{-4} e^{4(\operatorname{Im} z)_{-}}\right)
$$

We first consider the poles of $I+S F \mathcal{P}_{l}$. These poles are given by the zeros of the function

$$
\widetilde{\mathcal{D}}_{l}(z) \stackrel{\text { def }}{=} \operatorname{det}\left(I+S F \mathcal{P}_{l \pm}\right)=\left(1-\int_{\mathbb{R}}\left(S f_{+}\right) \phi_{z}\right)\left(1-\int_{\mathbb{R}}\left(S f_{-}\right) \phi_{-z}\right)-\left(\int_{\mathbb{R}}\left(S f_{-}\right) \phi_{z}\right)\left(\int_{\mathbb{R}}\left(S f_{+}\right) \phi_{-z}\right)
$$

with twice the multiplicity. A computation and use of the approximations $\tau_{l+1}=i \sqrt{2 l}+O\left(l^{-1 / 2}\right)$ and $\tau_{l-1}=\sqrt{2 l}+O\left(l^{-1 / 2}\right)$ show that

$$
\widetilde{\mathcal{D}}_{l}(z)=g_{l}(z)+O\left(l^{-3 / 2}\right)+O\left(l^{-2} \log l e^{2(\operatorname{Im} z)_{-}}\right)
$$

where $g_{l}$ is the function of Lemma 12.5. We note that both $g_{l}$ and $\widetilde{\mathcal{D}}_{l}$ are analytic in $z$ if $1<|z|<\frac{7}{8} \log l$. We use $z_{1}^{+}(l)$ as in Lemma 12.5. Recalling that $\operatorname{Im} z_{1}^{+} \sim-\frac{3}{4} \log l$, the estimate (12-11) combined with Rouché's theorem shows that $\widetilde{\mathcal{D}}_{l}(z)$ has a zero within $O\left(l^{-1 / 2+\epsilon}\right)$, for any $\epsilon>0$, of $z_{1}^{+}(l)$. This, in turn, means that

$$
\left(I+S F \mathcal{P}_{l}\right)^{-1}=\left(I+\left(I+\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}}\right)^{-1} F \mathcal{P}_{l}\right)^{-1}
$$

has a single pole of multiplicity two at a point satisfying $z=z_{1}^{+}(l)+O\left(l^{-1 / 2+\epsilon}\right)$. Moreover, we can find a $c_{0}=c_{0}(\epsilon)$ so that

$$
\left\|\left(I+\left(I-\frac{1}{2 l^{2}} \chi_{I_{0}} R_{0,0}(z) \chi_{I_{0}}\right) F \mathcal{P}_{l}\right)^{-1}\right\|=O\left(l^{1+\epsilon}\right)
$$

when the distance from $z$ to the pole is given by $c_{0} l^{-1 / 2+\epsilon}$.
Now using our estimate on $\|A\|$ we can apply the operator Rouché theorem to the pair $I+S F \mathcal{P}_{l}$ and $I+S F \mathcal{P}_{l}+S A$, to find that $I+S F \mathcal{P}_{l}+S A$ has two poles (when counted with multiplicity) which are, using the $z$-coordinate, within $O\left(l^{-1 / 2+\epsilon}\right)$ of $z_{1}^{+}(l)$.

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