# Inverse spectral results for the Schrödinger operator in Sobolev spaces 

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#### Abstract

We provide sufficiently large sets of eigenvalues that determine the potential of a Schrödinger operator on the unit interval $[0,1]$ partially known on $[a, 1]$ and belonging to $W^{k, p}$ in a neighbourhood of $a(k \in \mathbb{N} \cup\{0\}, p \in[1,+\infty])$. The number of these given eigenvalues depends on ( $a, k, p$ ).


## 1 Introduction and statement of the results

This paper is concerned with the Schrödinger operator

$$
\begin{equation*}
A_{q, h, H}=-\frac{d^{2}}{d x^{2}}+q \tag{1}
\end{equation*}
$$

defined on the unit interval with real-valued potentials $q$ belonging to $L^{1}((0,1))$. This operator is associated with the boundary conditions

$$
\begin{equation*}
u^{\prime}(0)+h u(0)=0, \quad u^{\prime}(1)+H u(1)=0 \tag{2}
\end{equation*}
$$

where $h, H$ are real numbers and where the notation 'stands for the derivative with respect to the variable $x$. It is well-known that, for each $(q, h, H) \in L^{1}([0,1]) \times \mathbb{R}^{2}$ the operator $A_{q, h, H}$ is a selfadjoint operator in $L^{2}([0,1])$. Its spectrum $\sigma\left(A_{q, h, H}\right)$ is an increasing and non-bounded sequence of non degenerate eigenvalues denoted by $\left(\lambda_{j}(q, h, H)\right)_{j \in \mathbb{N} \cup\{0\}}$.

Our purpose here is to provide sets of eigenvalues sufficiently large in order to determine a potential that is already known on $[a, 1]$ (for some given $\left.a \in\left(0, \frac{1}{2}\right]\right)$ when it belongs to some $W^{k, p}$ space. This problem has been initiated in 1978 by [HL] in the special case $a=\frac{1}{2}$ and for potentials in $L^{1}([0,1])$. In 2000, the problem is studied in [GS] for any $a$ and for potentials in $C^{2 k}$ near $a(k \in \mathbb{N} \cup\{0\})$. Results like one spectrum and half of another one added to the knowledge of the potential on $\left[\frac{3}{4}, 1\right]$ uniquely determine the potential are derived in [DGS1] and [DGS2]. Potentials in $L^{p}$ spaces are considered in [Ho] and [AR].

[^0]Recently, potentials in $W^{k, p}([0, a])$ are considered in [AFR] for any $a(p \in[1,+\infty])$, with however the restriction $k \in\{0,1,2\}$. We have conjectured in [AFR] that the result in [AFR] should be valid for all $k \in \mathbb{N} \cup\{0\}$. This is one of our aim here to get rid of this condition on $k$ and to consider all $k \in \mathbb{N} \cup\{0\}$. Our second goal is to replace regularity hypotheses of $q_{1}, q_{2}$ on $[0, a]$ by regularity hypotheses on $q_{1}, q_{2}$ only on an arbitrary small neighborhood of $a$ (as in [GS]).

The following function is involved in the statement of the main theorem (Theorem 1.1). For any complexvalued sequence $\alpha=\left(\alpha_{j}\right)_{j \in \mathbb{N} \cup\{0\}}$ and for all $t \geq 0$, we set

$$
\begin{equation*}
n_{\alpha}(t)=\sharp\left\{j \in \mathbb{N} \cup\{0\} \quad|\quad| \alpha_{j} \mid \leq t\right\} . \tag{3}
\end{equation*}
$$

The main result of the paper is the following.

## Theorem 1.1.

Set $q_{1}, q_{2} \in L^{1}((0,1))$. Fix $a \in\left(0, \frac{1}{2}\right]$ and suppose that $q_{1}=q_{2}$ on $[a, 1]$.
Let $k \in \mathbb{N} \cup\{0\}$ and $p \in[1,+\infty]$. Assume that $q_{1}, q_{2} \in W^{k, 1}((a-\varepsilon, a))$ with $q_{1}-q_{2} \in W^{k, p}((a-\varepsilon, a))$ for some arbitrary small $\varepsilon \in(0, a)$. If $k \geq 1$ assume in addition that $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+\varepsilon))$ with any arbitrary small $\varepsilon \in(0, a)$.

Fix the real numbers $h_{1}, h_{2}$ and $H$. Assume that a set of common eigenvalues $S \subseteq \sigma\left(A_{q_{1}, h_{1}, H}\right) \cap$ $\sigma\left(A_{q_{2}, h_{2}, H}\right)$ verifies either

$$
\begin{equation*}
n_{S}(t) \geq 2 a n_{\sigma(A)}(t)-\frac{k}{2}+\frac{1}{2 p}-\frac{1}{2}-a, t \in \sigma(A), \text { t large enough } \tag{H}
\end{equation*}
$$

or assume that there exists a real number $C$ such that

$$
2 a n_{\sigma(A)}(t)+C \geq n_{S}(t) \geq 2 a n_{\sigma(A)}(t)-\frac{k}{2}+\frac{1}{2 p}-2 a, t \in S, \text { large enough }
$$

where in $(H)$ and $\left(H^{\prime}\right)$ the operator $A$ denotes either $A_{q_{1}, h_{1}, H}$ or $A_{q_{2}, h_{2}, H}$.
Then $h_{1}=h_{2}$ and $q_{1}=q_{2}$.

Let us emphasize here that the case $p=+\infty$ is considered in Theorem 1.1. In that case, the term $\frac{1}{p}$ in the hypotheses $(H)$ or $\left(H^{\prime}\right)$ is suppressed. Also note that only the difference of the two potentials needs to be in $W^{k, p}$ and $C^{k-1}$ near $a$.

One may replace the assumptions on $q_{1}$ and $q_{2}$ in Theorem 1.1 by the more concise (but stronger) hypotheses: $q_{1}, q_{2} \in L^{1}((0,1)) \cap W^{k, 1}((a-\varepsilon, a))$ with $q_{1}-q_{2} \in W^{k, p}((a-\varepsilon, a+\varepsilon))$ (since $W^{k, p}((a-\varepsilon, a+$ $\varepsilon)) \subset C^{k-1}([a-\varepsilon, a+\varepsilon])$ for $\left.k \geq 1\right)$.

Let us gives two corollaries of Theorem 1.1. The first one concerns the particular case $k=0, p=1$ and $a=\frac{1}{4}$. It is already given in [AR] (where $k=0$ ), it is however recalled here in order to emphasize on the role of $\left(H^{\prime}\right)$ in Theorem 1.1. Namely, this corollary may be proved using the assumption $\left(H^{\prime}\right)$ while it is not be derived assuming $(H)$ (see [AR]). It is written in a short way.

Corollary 1.2. Suppose that $q$ belongs to $L^{1}((0,1))$ and $H \in \mathbb{R}$. Then the even (resp. odd) spectrum $\left(\lambda_{2 j}(q, h, H)\right)_{j \geq 0}\left(\right.$ resp. $\left.\left(\lambda_{2 j+1}(q, h, H)\right)_{j \geq 0}\right),\left.q\right|_{\left[0, \frac{1}{4}\right]}$ and $H$ uniquely determine $h$ and the potential $q$ on all of $[0,1]$.

The second corollary is Theorem 1.1 in the particular case $p=+\infty$ and $a=\frac{1}{2}$ using hypothesis ( $H$ ). It allows us to remove a precise number of eigenvalues when the potentials (and their difference) are sufficiently regular. It slightly improves one of the results established in [GS]. The result in [GS] is the same as Corollary 1.3 but the potentials satisfy $q_{1}, q_{2} \in L^{1}((0,1)) \cap C^{2 k}\left(\left(\frac{1}{2}-\varepsilon, \frac{1}{2}+\varepsilon\right)\right)$ for some small positive $\varepsilon$.

Corollary 1.3. Let $k \in \mathbb{N} \cup\{0\}$. Suppose that $q_{1}$ and $q_{2}$ belong to $L^{1}((0,1)) \cap W^{2 k, \infty}\left(\left(\frac{1}{2}-\varepsilon, \frac{1}{2}\right)\right)$ and if $k \geq 1$ also assume that the difference $q_{1}-q_{2}$ is in $C^{2 k-1}\left(\left(\frac{1}{2}-\varepsilon, \frac{1}{2}+\varepsilon\right)\right)$ for some $\varepsilon>0$. Suppose that $q_{1}=q_{2}$ on $\left[\frac{1}{2}, 1\right]$. Let $h_{1}, h_{2}, H \in \mathbb{R}$.

If $\sigma\left(A_{q_{1}, h_{1}, H}\right)=\sigma\left(A_{q_{2}, h_{2}, H}\right)$ excepted for at most $k+1$ eigenvalues, then $h_{1}=h_{2}$ and $q_{1}=q_{2}$.

Also note that this implies that, if $q \in L^{1}((0,1))$ is $L^{\infty}$ near $x=\frac{1}{2}$ then $q$ on $\left[0, \frac{1}{2}\right], H$ and all the eigenvalues of $\sigma\left(A_{q, h, H}\right)$ excepted one, uniquely determine $h$ and $q$ on $[0,1]$.

The proof of Theorem 1.1 relies on the same strategy as in [AFR] excepted that [AFR, Proposition 3.1] is replaced by Proposition 1.4 below. Let us also mention that our proof is different from the result in [GS] which and is based on Weyl-Titchmarsh functions.

The estimate in Proposition 1.4 is the same as the one in [AFR, Proposition 3.1] but the assumption on $k$ and on the regularity on $q_{1}-q_{2}$ are largely weakened. Firstly, $k \in\{0,1,2\}$ in [AFR] is replaced here with $k \in \mathbb{N} \cup\{0\}$. Secondly, the hypotheses $q_{1}, q_{1} \in W^{k, 1}([0,1])$ and $q_{1}-q_{2} \in W^{k, p}([0, a])$ in [AFR] is now replaced by the hypotheses on $q_{1}, q_{2}$ in Theorem 1.1, namely, $q_{1}, q_{2} \in L^{1}((0,1)) \cap W^{k, 1}((a-\varepsilon, a))$ with $q_{1}-q_{2} \in W^{k, p}((a-\varepsilon, a))$ added when $k \geq 1$ to $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+\varepsilon))$ (for any arbitrary small $\varepsilon \in(0, a))$. Note that Sobolev's imbedding implies that $W^{k, p}((a-\varepsilon, a)) \subset C^{k-1}([a-\varepsilon, a])$ when $k \geq 1$.

We now define the entire function $f$ which is involved in Proposition 1.4. Fix $q \in L^{1}((0,1))$ and fix $H, h \in \mathbb{R}$. For any $z$ in $\mathbb{C}$, let $\psi(\cdot, z, q, h)$ be defined on $[0,1]$ as the solution to $-\frac{d^{2} \psi}{d x^{2}}+q \psi=z^{2} \psi$, $\psi(0)=1, \psi^{\prime}(0)=-h$. It is known that $\psi(x, \cdot, q, h)$ is an entire function ([LG]).

For all $z \in \mathbb{C}$, let us define

$$
\begin{equation*}
f(z)=\int_{0}^{a}\left(\psi\left(x, z, q_{1}, h_{1}\right) \psi\left(x, z, q_{2}, h_{2}\right)-\frac{1}{2}\right)\left(q_{1}(x)-q_{2}(x)\right) d x \tag{4}
\end{equation*}
$$

Proposition 1.4. Set $a \in\left(0, \frac{1}{2}\right]$, fix $k \in \mathbb{N} \cup\{0\}$ and let $p \in[1,+\infty]$. Fix $q_{1}, q_{2} \in L^{1}((0,1)) \cap$ $W^{k, 1}((a-\varepsilon, a))$ such that $q_{1}-q_{2} \in W^{k, p}((a-\varepsilon, a))$ and assume furthermore that $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+$ $\varepsilon)$ ) when $k \geq 1$ for some arbitrary small $\varepsilon \in(0, a)$. Then there is a real positive number $C$ independent of $z \in \mathbb{C}$ and $\varepsilon^{\prime}>0$ such that $|f(z)| \leq C \frac{e^{2|\Im z| a}}{|\Im z|^{k+1-\frac{1}{p}}}\left(e^{-\varepsilon^{\prime}|\Im z|}+o(1)\right)$ as $\varepsilon^{\prime} \rightarrow 0^{+}$uniformly in $z \in \mathbb{C} \backslash \mathbb{R}$.

Proof of Theorem 1.1: It is the same as the one of [AFR, Theorem 1] when replacing [AFR, Proposition 3.1] by Proposition 1.4 above. For the sake of completeness let us recall very briefly here the main steps (see [AFR] for more details). Suppose that $a, k, p, q_{1}, q_{2}, h_{1}, h_{2}$ satisfy the same assumptions as the ones in Theorem 1.1. Define the $s_{j}, j \in \mathbb{N}$ as the strictly increasing sequence being in $S$ and define the set $S^{\frac{1}{2}}=\left\{ \pm \sqrt{s_{j}}, j \in \mathbb{N}\right\}$. We also define for any set of complex numbers $\alpha, N_{\alpha}(R)=\int_{0}^{R} \frac{n_{\alpha}(t)}{t} d t$, for any $R>0$ and where $n_{\alpha}(t)$ is given in (3). On one hand, the hypothesis $(H)$ or $\left(H^{\prime}\right)$ implies that the sequence $\left(N_{S^{\frac{1}{2}}}\left(\sqrt{s_{j}}\right)-\frac{4 a}{\pi} \sqrt{s_{j}}+\left(k+1-\frac{1}{p}\right) \ln \sqrt{s_{j}}\right)_{j \in \mathbb{N}}$ is bounded from below ([AFR, Prop. 4.1 and 4.2]). On the other side, using Proposition 1.4 above, using [AFR, Prop. 4.3] and assuming that $f$ is not entirely vanishing in order to use Jensen's Theorem, we deduce that $\lim _{R \rightarrow+\infty} N_{f^{-1}(0)}(R)-\frac{4 a}{\pi} R+$ $\left(k+1-\frac{1}{p}\right) \ln R=-\infty$. The last two points combined to $N_{f^{-1}(0)} \geq N_{S^{\frac{1}{2}}}$ (see (23) in [AFR]) lead to a contradiction if $f$ is not entirely vanishing. The fact that $f \equiv 0$ implies that $\left(q_{1}, h_{1}\right)=\left(q_{2}, h_{2}\right)$ is already proved in [L].

The rest of this paper is therefore concerned with the proof of Proposition 1.4. The main difference here is that we imply the transformation operators [L] (see also references therein and see [Le],[LS], $[\mathrm{M}], \ldots$ ) instead of using expansions of the fundamental solutions to $A_{q, h, H} y=z y$

Proposition 1.4 is derived in the next section. The case of Dirichlet boundary conditions is considered in Appendix A.

## 2 Proof of Proposition 1.4

The proof of Proposition 1.4 shall follow from Lemmata 2.1-2.7 below.

We first start with the definition of the transformation operators (see [L], [Le], $[\mathrm{LS}],[\mathrm{M}]$ and references therein). We shall use in the following the kernel $\tilde{L}$ (see (11) below) computed in [L]. This kernel is expressed in terms of the kernel $L$ (see (7) below). Its properties are taken from $[\mathrm{M}]$.

We first recall the definition of $L$ given by [M]. To do this we first define the kernel $K$.

Suppose $q \in L_{\text {loc }}^{1}((0,1))$. There exists a kernel $K \equiv K(x, t)$ for $0 \leq x \leq 1$ and $-x \leq t \leq x$ (see [M]) such that, for each $z \in \mathbb{C}$, the solution $\alpha \equiv \alpha(x, z)$ to

$$
\begin{equation*}
-\alpha^{\prime \prime}+q \alpha=z^{2} \alpha \quad[0,1], \quad \alpha(0, z)=1, \alpha^{\prime}(0, z)=i z \tag{5}
\end{equation*}
$$

may be expressed as

$$
\begin{equation*}
\alpha(x, z)=e^{i z x}+\int_{0}^{x} K(x, t) e^{i z t} d t, \quad x \in[0,1] . \tag{6}
\end{equation*}
$$

Let

$$
\begin{equation*}
L(x, t)=-h+K(x, t)+K(x,-t)-h \int_{t}^{x}(K(x, \xi)-K(x,-\xi)) d \xi \tag{7}
\end{equation*}
$$

for $0 \leq t \leq x \leq 1$. One then obtains that, for each $z \in \mathbb{C}$, the solution $\beta \equiv \beta(x, z)$ to

$$
\begin{equation*}
-\beta^{\prime \prime}+q \beta=z^{2} \beta \quad[0,1], \quad \beta(0, z)=1, \beta^{\prime}(0, z)=-h \tag{8}
\end{equation*}
$$

may be expressed as

$$
\begin{equation*}
\beta(x, z)=\cos z x+\int_{0}^{x} L(x, t) \cos z t d t, \quad x \in[0,1] . \tag{9}
\end{equation*}
$$

Let us denote respectively by $I$ and $T_{L}$ the identity operator and Volterra operator with kernel $L(x, t)$. Let $\beta_{0}(x, z)=\cos z x(z \in \mathbb{C})$. With these notations (9) is also written as $\beta(\cdot, z)=\left(I+T_{L}\right) \beta_{0}(\cdot, z)$ for any $z \in \mathbb{C}$. Fix $q \in L^{1}((0,1))$ and $h \in \mathbb{R}$. The main point is that the operator $\left(I+T_{L}\right)$ maps the solution to (8) with $q$ identically vanishing and $h=0$ to solution to (8) with the potential $q$ and the parameter $h$.

Fix $q_{j}$ in $L^{1}((0,1))$ and $h_{j} \in \mathbb{R}$ for $j=1,2$. Set $L_{j}$ the function defined in (7) associated to $q=q_{j}$ and $h=h_{j}, j=1,2$ and extended for $t \in[-x, 0]$ by setting $L(x, t)=L(x,-t)$. With these notations one obtains (see [L, Appendix IV]),

$$
\begin{equation*}
\psi\left(x, z, q_{1}, h_{1}\right) \psi\left(x, z, q_{2}, h_{2}\right)-\frac{1}{2}=\frac{1}{2} \cos 2 z x+\frac{1}{2} \int_{-x}^{x} \tilde{L}(x, \tau) \cos 2 z \tau d \tau \tag{10}
\end{equation*}
$$

with $x \in[0,1]$ and where

$$
\tilde{L}(x, \tau)=\left\{\begin{array}{l}
2\left(L_{1}(x, x-2 \tau)+L_{2}(x, x-2 \tau)\right)+\int_{-x+2 \tau}^{x} L_{1}(x, s) L_{2}(x, s-2 \tau) d s \text { if } \tau>0  \tag{11}\\
\int_{-x}^{x+2 \tau} L_{1}(x, s) L_{2}(x, s-2 \tau) d s \text { if } \tau<0
\end{array}\right.
$$

Throughout the paper we suppose that $a$ and $\varepsilon$ are fixed in $\left(0, \frac{1}{2}\right]$ and $(0, a)$ respectively. Let us first decompose $f$ as

$$
f(z)=f_{a-\varepsilon}(z)+f_{a}(z)
$$

with

$$
f_{a-\varepsilon}(z)=\int_{0}^{a-\varepsilon}\left(\psi\left(x, z, q_{1}, h_{1}\right) \psi\left(x, z, q_{2}, h_{2}\right)-\frac{1}{2}\right)\left(q_{1}(x)-q_{2}(x)\right) d x
$$

and

$$
f_{a}(z)=\int_{a-\varepsilon}^{a}\left(\psi\left(x, z, q_{1}, h_{1}\right) \psi\left(x, z, q_{2}, h_{2}\right)-\frac{1}{2}\right)\left(q_{1}(x)-q_{2}(x)\right) d x
$$

for all $z \in \mathbb{C}$. The function $f_{a-\varepsilon}$ is easily estimated.

Lemma 2.1. For $q_{1}$ and $q_{2}$ in $L^{1}((0,1))$ we have

$$
f_{a-\varepsilon}(z)=O\left(e^{2|\operatorname{Im} z|(a-\varepsilon)}\right)
$$

uniformly in $z \in \mathbb{C}$.

Proof of Lemma 2.1: It follows from the asymptotic expansions of the function $\psi$.
Namely, $\psi(x, z, q, h)=O\left(e^{|I m z| x}\right)$ uniformly for $(z, x) \in \mathbb{C} \times[0,1]$ (see [LG]) and $q_{1}-q_{2} \in L^{1}((0,1))$ directly implies the stated estimate on $f_{a-\varepsilon}(z)$.

In view of (10) the function $f_{a}$ is split as

$$
f_{a}(z)=f_{0}(z)+\tilde{f}(z), \forall z \in \mathbb{C}
$$

where

$$
\begin{equation*}
f_{0}(z)=\frac{1}{2} \int_{a-\varepsilon}^{a} \cos 2 z x\left(q_{1}(x)-q_{2}(x)\right) d x, \forall z \in \mathbb{C} \tag{12}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{f}(z)=\frac{1}{2} \int_{a-\varepsilon}^{a}\left(\int_{-x}^{x} \tilde{L}(x, \tau) \cos 2 z \tau d \tau\right)\left(q_{1}(x)-q_{2}(x)\right) d x, \forall z \in \mathbb{C} \tag{13}
\end{equation*}
$$

We shall first estimate the function $f_{0}$. For any $k \in \mathbb{N}$, let $c_{k}$ be the $k^{t h}$ integral of the cosine function verifying $c_{k}^{(l)}(0)=0, l=0, \ldots, k-1$, that is to say, $c_{k}(x)=\int_{0}^{x} \int_{0}^{t_{k}} \cdots \int_{0}^{t_{2}} \cos \left(t_{1}\right) d t_{1} \ldots d t_{k}$. For any $l \in \mathbb{N}$ and for any sufficiently smooth function $g$ depending only on one variable, $g^{(l)}$ denotes its $l^{\text {th }}$ derivative.

Lemma 2.2. Fix $k \in \mathbb{N} \cup\{0\}$. Let $q_{1}, q_{2} \in W^{k, 1}((a-\varepsilon, a))$ with $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+\varepsilon))$ if $k \geq 1$. There exist $k$ complex numbers $L_{0, l}(z)(l=1, \ldots, k)$ satisfying

$$
\begin{equation*}
L_{0, l}(z)=O\left(e^{2|\Im z|(a-\varepsilon)}\right) \tag{14}
\end{equation*}
$$

uniformly in $z \in \mathbb{C}$ and such that

$$
\begin{equation*}
f_{0}(z)=\sum_{l=1}^{k} \frac{L_{0, l}(z)}{(2 z)^{l}}+\int_{a-\varepsilon}^{a} \frac{c_{k}(2 z x)}{(-2 z)^{k}}\left(q_{1}-q_{2}\right)^{(k)}(x) d x \tag{15}
\end{equation*}
$$

for all $z \in \mathbb{C} \backslash\{0\}$.

Proof of Lemma 2.2: Clearly one may suppose that $k \geq 1$. Then one can integrate by parts $k$ times the r.h.s. of (12) since $q_{1}-q_{2} \in W^{k, 1}((a-\varepsilon, a))$. Since $q_{1}-q_{2} \in C^{k-1}$ near $a$ and using $q_{1}-q_{2} \equiv 0$ on $[a, 1]$ we see that $\left(q_{1}-q_{2}\right)^{(l)}(a)=0, l=0, \ldots k-1$. This shows that the $k$ boundary terms at $x=a$ are vanishing. It remains $k$ boundary terms at $x=a-\varepsilon$. These terms lead to $\sum_{l=1}^{k} \frac{L_{0, l}(z)}{2 z^{l}}$ with the $L_{0, l}(z)=(-1)^{l-1} c_{l}(2 z(a-\varepsilon))\left(q_{1}-q_{2}\right)^{(l-1)}(a-\varepsilon)$. Using $|\sin z| \leq e^{|\Im z|}$ and $|\cos z| \leq e^{|\Im z|}$ for all $z \in \mathbb{C}$ one clearly gets (14) and (15).

Next and in order to deal with $\tilde{f}$ we write using Fubini's theorem that

$$
\tilde{f}=f_{1}+f_{2}+f_{3}
$$

with

$$
\begin{gather*}
f_{1}(z)=\int_{a-\varepsilon}^{a} \int_{\tau}^{a} \tilde{L}(x, \tau)\left(q_{1}(x)-q_{2}(x)\right) \cos 2 z \tau d x d \tau  \tag{16}\\
f_{2}(z)=\int_{-a}^{-a+\varepsilon} \int_{-\tau}^{a} \tilde{L}(x, \tau)\left(q_{1}(x)-q_{2}(x)\right) \cos 2 z \tau d x d \tau \tag{17}
\end{gather*}
$$

and

$$
\begin{equation*}
f_{3}(z)=\int_{a-\varepsilon}^{a} \int_{-a+\varepsilon}^{a-\varepsilon} \tilde{L}(x, \tau)\left(q_{1}(x)-q_{2}(x)\right) \cos 2 z \tau d \tau d x \tag{18}
\end{equation*}
$$

for all $z \in \mathbb{C}$. Consequently, we shall only consider $f_{1}$ and $f_{3}$ in the sequel since the treatment of $f_{2}$ would be similar to $f_{1}$ making the change of variables $\tau \mapsto-\tau$ in $f_{2}$.

Set

$$
\begin{equation*}
w(\tau)=\int_{\tau}^{a} \tilde{L}(x, \tau)\left(q_{1}(x)-q_{2}(x)\right) d x \tag{19}
\end{equation*}
$$

for any $\tau \in(a-\varepsilon, a)$. That is to say,

$$
\begin{equation*}
f_{1}(z)=\int_{a-\varepsilon}^{a} w(\tau) \cos 2 z \tau d \tau \tag{20}
\end{equation*}
$$

for all $z \in \mathbb{C}$. In order to integrate by parts the r.h.s. of (20), we need that $w$ defined in (19) belongs to $W^{k, 1}((a-\varepsilon, a))$. It is actually in $W^{k, \infty}([a-\varepsilon, a])$. This is precisely the purpose of Lemma 2.4 below with the help of Lemma 2.3.

In the sequel, for any sufficiently smooth function $g$ depending on the variables $\left(x_{1}, \ldots, x_{n}\right), \partial_{j_{1}, \ldots, j_{l} g}$ stands for the derivative of order $l$ of $g$ with respect the variables $x_{j_{1}}, \ldots, x_{j_{l}}$ (with $j_{1}, \ldots, j_{l} \in\{1, \ldots, n\}$, $l \in \mathbb{N}$ ) and $\partial_{j}^{m} g$ denotes the derivative of order $m$ of $g$ with respect the variable $x_{j}$ (where $j \in\{1, \ldots, n\}$, $m \in \mathbb{N})$.

Let us recall that the kernel $\tilde{L}$ is written in terms of the two kernels $L_{1}$ and $L_{2}$ and these two kernels $L_{j}$ $(j=1,2)$ are expressed in (7) with the functions $K_{j}$ corresponding to $q=q_{j}$.

Set $T_{a, \varepsilon}$ be the triangle $\{a-\varepsilon \leq t \leq x \leq a\}$ and let $D_{a, \varepsilon}$ be the diagonal $D_{a, \varepsilon}=\{(\tau, \tau) \mid \tau \in[a-\varepsilon, a]\}$. Let us recall here that in this section $\varepsilon$ is fixed in $(0, a)$.

Lemma 2.3. (i) Fix $k \in \mathbb{N} \cup\{0\}$ and $q \in L^{1}((0,1)) \cap W^{k, 1}((a-\varepsilon, a))$. Then, the kernel $K$ associated to $q$ belongs to $C^{k}\left(T_{a, \varepsilon}\right)$.
(ii) Suppose that $q \in L^{1}((0,1)) \cap W^{k, 1}((a-\varepsilon, a))$ for some $k \in \mathbb{N} \cup\{0\}$. Then, the kernel $L$ defined in (7) corresponding to $q$ is in $C^{k}\left(T_{a, \varepsilon}\right)$.
(iii) Assume that $q_{1}, q_{2} \in L^{1}((0,1)) \cap W^{k, 1}((a-2 \varepsilon, a))$ for $k \geq 0$ (with $0<\varepsilon<\frac{a}{2}$ ). If $k=0$ then the kernel $\tilde{L}$ given by (11) is in $C^{0}\left(T_{a, \varepsilon}\right)$. When $k \geq 1$ then $\partial_{2}^{j} \tilde{L} \in C^{0}\left(T_{a, \varepsilon}\right)$ for all $0 \leq j \leq k$ and $\left[\left.\left(\partial_{2}^{l} \tilde{L}\right)\right|_{D_{a, \varepsilon}}\right]^{(\alpha)} \in C^{0}\left(D_{a, \varepsilon}\right)$ for $l+\alpha \leq k$ (with $l \geq 0$ and $\left.\alpha \geq 0\right)$.

## Proof of Lemma 2.3:

(i) It is proved in Theorem 1.2.1 in $[\mathrm{M}]$ (see also Problem 1 in $[\mathrm{M}]$ ) that if $q \in L_{l o c}^{1}((0,1))$ then the kernel $K$ belongs to $C^{0}(T)$ where $T=\{0 \leq t \leq x \leq 1\}$. When $k \geq 1$, if $q \in W^{k, 1}((0,1))$ then $q \in C^{k-1}([0,1])$ and it is derived in Theorem 1.2.2 $([\mathrm{M}])$ that $K \in C^{k}(T)$.

Here $q \in L^{1}((0,1))$ then $K$ exists and is continuous on $T$ and the same arguments as in $([\mathrm{M}])$ show that $K \in C^{k}\left(T_{a, \varepsilon}\right)$ when $q \in W^{k, 1}((a-\varepsilon, a))(k \geq 1)$.
(ii) From the definition of $L$ (see (7)) and (i) we only have to check that $I$ defined by $I(x, t)=\int_{t}^{x} K(x, \xi) d \xi$ verifies $I \in C^{k}\left(T_{a, \varepsilon}\right)$ when $q \in W^{k, 1}((a-\varepsilon, a))(k \geq 0)$.

If $k=0$ then $K \in C^{0}(T)$ and $I \in C^{0}\left(T_{a, \varepsilon}\right)$.
If $k \geq l_{1} \geq 1$ then

$$
\begin{equation*}
\partial_{1}^{l_{1}} I(x, t)=\sum_{i, j \geq 0, i+j=l_{1}-1}\left[\left.\left(\partial_{1}^{j} K\right)\right|_{D_{a, \varepsilon}}\right]^{(i)}(x)+\int_{t}^{x} \partial_{1}^{l_{1}} K(x, \xi) d \xi \tag{21}
\end{equation*}
$$

for all $(x, t) \in T_{a, \varepsilon}$.
If $k \geq l_{2} \geq 1$ then

$$
\begin{equation*}
\partial_{2}^{l_{2}} I(x, t)=-\partial_{2}^{l_{2}-1} K(x, t) \tag{22}
\end{equation*}
$$

for all $(x, t) \in T_{a, \varepsilon}$.
Thus, if $l_{1} \geq 1, l_{2} \geq 1$ with $l_{1}+l_{2} \leq k$ then,

$$
\begin{equation*}
\partial_{1}^{l_{1}} \partial_{2}^{l_{2}} I(x, t)=-\partial_{1}^{l_{1}} \partial_{2}^{l_{2}-1} K(x, t) \tag{23}
\end{equation*}
$$

for any $(x, t)$ in $T_{a, \varepsilon}$. In view of (21)(22) (23) and according to (i) we see that $I \in C^{k}\left(T_{a, \varepsilon}\right)$ when $k \geq 1$.
(iii) From the definition (11) and following the point (ii) above it is sufficient to verify that $J$ satisfies $\partial_{2}^{k} J \in C^{0}\left(T_{a, \varepsilon}\right)$ and $\left[\left.\partial_{2}^{l} J\right|_{D_{a, \varepsilon}}\right]^{(\alpha)} \in C^{0}\left(D_{a, \varepsilon}\right)$ when $l+\alpha \leq k$ where the function $J$ is defined by

$$
\begin{equation*}
J(x, \tau)=\int_{-x+2 \tau}^{x} L_{1}(x, s) L_{2}(x, s-2 \tau) d s \tag{24}
\end{equation*}
$$

for all $(x, \tau) \in T_{a, \varepsilon}$.
If $k=0$ then $L_{1}$ and $L_{2}$ are continuous on $T_{a, 2 \varepsilon}$ and $J \in C^{0}\left(T_{a, \varepsilon}\right)$.

Suppose $k \geq 1$. One may differentiate the r.h.s of (24) $k$ times with respect to the second variable. Indeed, one gets

$$
\begin{align*}
\partial_{2}^{k} J(x, \tau)= & \sum_{i, j \geq 0, i+j=k-1} 2^{i}(-2)^{j+1} \partial_{2}^{i} L_{1}(x,-x+2 \tau) \partial_{2}^{j} L_{2}(x,-x)+ \\
& (-2)^{k} \int_{-x+2 \tau}^{x} L_{1}(x, s) \partial_{2}^{k} L_{2}(x, s-2 \tau) d s \tag{25}
\end{align*}
$$

for all $(x, \tau)$ in $T_{a, \varepsilon}$. According to (ii), this implies that $\partial_{2}^{k} J \in C^{0}\left(T_{a, \varepsilon}\right)$. Moreover, on the diagonal $D_{a, \varepsilon}$ the last integral in (25) vanishes and we obtain after differentiating $\alpha$ times that,

$$
\left[\left.\left(\partial_{2}^{l} J\right)\right|_{D_{a, \varepsilon}}\right]^{(\alpha)}(x)=\sum_{i+j=l-1} \alpha_{1}+\alpha_{2}=\alpha<i c_{i j \alpha_{1} \alpha_{2}} \partial_{2}^{i}\left(\partial_{1}+\partial_{2}\right)^{\alpha_{1}} L_{1}(x, x) \partial_{2}^{j}\left(\partial_{1}+\partial_{2}\right)^{\alpha_{2}} L_{2}(x,-x)
$$

for some numerical real number $c_{i j \alpha_{1} \alpha_{2}}$, for any $x \in[a-\varepsilon, a]$ and for all $l+\alpha \leq k$. Since $i+\alpha_{1} \leq l+\alpha$, $j+\alpha_{2} \leq l+\alpha, l+\alpha \leq k$ and since $L_{1}$ and $L_{2}$ are $C^{k}\left(T_{a, \varepsilon}\right)$ then $\left[\left.\left(\partial_{2}^{k} J\right)\right|_{D_{a, \varepsilon}}\right]^{(\alpha)}$ is continuous on $[a-\varepsilon, a]$.

Lemma 2.4. Set $k \in \mathbb{N} \cup\{0\}$ and let $q_{1}, q_{2} \in L^{1}((0,1)) \cap W^{k, 1}((a-2 \varepsilon, a))$. Then the function $w$ defined in (19) belongs to $W^{k, \infty}((a-\varepsilon, a))$.

Proof of Lemma 2.4: From (19) it is clear that

$$
\begin{equation*}
w^{(j)}(\tau)=\sum_{l, m \geq 0, l+m=j-1} \sum_{\alpha, \beta \geq 0, \alpha+\beta=m} c_{j l m \alpha \beta}\left[\left.\left(\partial_{2}^{l} \tilde{L}\right)\right|_{D_{a, \varepsilon}}\right]^{(\alpha)}(\tau)\left(q_{1}-q_{2}\right)^{(\beta)}(\tau)+\int_{\tau}^{a} \partial_{2}^{j} \tilde{L}(x, \tau)\left(q_{1}-q_{2}\right)(x) d x \tag{26}
\end{equation*}
$$

for all $\tau \in[a-\varepsilon, a]$ and for some numerical coefficients $c_{j l m \alpha \beta}$ provided that the r.h.s. is well-defined. If $j=0$ the first term in the r.h.s. of the equality above is omitted. Let us verify that $w^{(j)} \in L^{\infty}((a-\varepsilon, a))$ for $j \leq k$. Since $l+\alpha \leq l+m \leq k-1$ then $\tau \mapsto\left[\left.\left(\partial_{2}^{l} \tilde{L}\right)\right|_{D_{a, \varepsilon}}\right]^{(\alpha)}(\tau) \in L^{\infty}((a-\varepsilon, a))$ by Lemma 2.3 (iii). Since $\beta \leq m \leq k-1$ then $\left(q_{1}-q_{2}\right)^{(\beta)} \in L^{\infty}((a-\varepsilon, a))$. Thus, the first term in the r.h.s. of the above equality is in $L^{\infty}([a-\varepsilon, a])$ as a function of the variable $\tau$. Furthermore, $\left(q_{1}-q_{2}\right) \in L^{1}((a-\varepsilon, a))$ and $\partial_{2}^{j} \tilde{L} \in L^{\infty}\left(T_{a, \varepsilon}\right)$ by Lemma 2.3 (iii) imply that the second term in the r.h.s. of the above equality for all $j \in \mathbb{N} \cup\{0\}$ with $j \leq k$ is also in $L^{\infty}((a-\varepsilon, a))$ as a function of the variable $\tau$.

With this Lemma, we are now able to integrate by parts $k$ times the function in the r.h.s. of (20). We recall that the functions $c_{k}$ are defined before Lemma 2.2.

Lemma 2.5. Let $k \in \mathbb{N} \cup\{0\}$. Set $q_{1}, q_{2} \in L^{1}((0,1)) \cap W^{k, 1}((a-2 \varepsilon, a))$ and if $k \geq 1$ assume in addition that $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+\varepsilon))$. One has

$$
f_{1}(z)=\sum_{l=1}^{k} \frac{L_{1, l}(z)}{(2 z)^{l}}+\frac{1}{(-2 z)^{k}} \int_{a-\varepsilon}^{a} w^{(k)}(\tau) c_{k}(2 z \tau) d \tau
$$

for all $z \in \mathbb{C}$, for $i=1,2$ where the $L_{1, l}(z)$ are $k$ real numbers satisfying

$$
\begin{equation*}
L_{1, l}(z)=O\left(e^{2|\Im z|(a-\varepsilon)}\right) \tag{27}
\end{equation*}
$$

for all $z \in \mathbb{C} \backslash\{0\}$.

Proof of Lemma 2.5: It suffices to suppose $k \geq 1$. As in Lemma 2.2, the proof follows from $k$ integrations by parts. These are justified by the regularity of $w$ provided by Lemma 2.4. Note also that all the boundary terms at $\tau=a$ are vanishing. Indeed, in view of (26) one sees that $w^{(\beta)}(a)=0, \beta=0, \ldots, k-1$ since $q_{1}-q_{2} \in C^{k-1}((a-\varepsilon, a+\varepsilon))$ and $q_{1}-q_{2}=0$ on $[a, 1]$ and since the last integral vanishes. Therefore $L_{1, l}(z)=(-1)^{l-1} c_{l}(2 z(a-\varepsilon)) w^{(l-1)}(a-\varepsilon)$. From Lemma 2.6 the function $w^{(k)} \in C^{0}([a-\varepsilon, a])$ and using again $\left|c_{k}(2 z \tau)\right| \leq e^{2|\Im z| \tau}$ one gets the estimate (27).

Finally we consider $f_{3}(z)$ defined in (18).
Lemma 2.6. Let $k \in \mathbb{N} \cup\{0\}$. Let $q_{1}, q_{2} \in L^{1}((0,1))$. One has

$$
f_{3}(z)=O\left(e^{2|\Im z|(a-\varepsilon)}\right)
$$

for all $z \in \mathbb{C}$.

Proof of Lemma 2.6: it follows directly from (18) with $\tilde{L} \in C^{0}(T)$ and $|\cos 2 z \tau| \leq e^{2|\Im z| \tau}$ for all $z \in \mathbb{C}$ and all $\tau \in \mathbb{R}$ together with $q_{1}-q_{2} \in L^{1}((0,1))$.

We are now ready to derive Proposition 1.4. Let us first recall the following result (see $[\mathrm{L}]$ and see also Lemma 3.2 in [AFR] for a short proof replacing 0 by $b$ ).

Lemma 2.7. Let $a, b \in(0,1]$ with $b<a$. Suppose that the function $u$ defined on $[0,1] \times \mathbb{C}$ satisfies $|u(x, z)|=O\left(e^{2|\Im z| x}\right)$ and let $v \in L^{p}([0,1])$ with $1 \leq p \leq+\infty$. Set $g(z)=\int_{b}^{a} u(x, z) v(x) d x$. There is a real positive number $C$ depending only on $p$ and $\|v\|_{L^{p}([b, a])}$ such that for any $\varepsilon^{\prime}>0$ there is a real positive number $\delta_{\varepsilon^{\prime}}$ depending only on $\varepsilon^{\prime}, p, a, b$ and $\|v\|_{L^{p}([b, a])}$ verifying

$$
\lim _{\varepsilon^{\prime} \rightarrow 0} \delta_{\varepsilon^{\prime}}=0
$$

and

$$
|g(z)| \leq C \frac{e^{2|\Im z| a}}{|\Im z|^{1-\frac{1}{p}}}\left(e^{-\varepsilon^{\prime}|\Im z|}+\delta_{\varepsilon^{\prime}}\right)
$$

Proof of Proposition 1.4: Without loss of generality we suppose that $q_{1}$ and $q_{2}$ are in $W^{k, 1}((a-2 \varepsilon, a))$ instead of $W^{k, 1}((a-\varepsilon, a))$. Let us denote by $w_{1}$ the preceding function $w$ defined in (19) associated to $f_{1}$ and by $w_{2}$ the similar one corresponding to $f_{2}$. Using the estimates for $f_{a-\varepsilon}, f_{0}, f_{1}$ (and the analogous one for $f_{2}$ ) and $f_{3}$ in Lemma 2.1, 2.2 2.5 and 2.6 respectively, one has

$$
\begin{equation*}
f(z)=\sum_{l=0}^{k} O\left(\frac{e^{2|\Im z|(a-\varepsilon)}}{|z|^{l}}\right)+\frac{1}{(-2 z)^{k}} \int_{a-\varepsilon}^{a} c_{k}(2 z x)\left(q_{1}-q_{2}+w_{1}+w_{2}\right)^{(k)}(x) d x \tag{28}
\end{equation*}
$$

Since $\left(q_{1}-q_{2}\right) \in W^{k, p}((a-\varepsilon, a))$ and $w_{1}, w_{2} \in W^{k, \infty}((a-\varepsilon, a))$ one concludes with Lemma 2.7 that the integral term in the r.h.s. of (28) is bounded by $\frac{e^{2|\Im z| a}}{|\Im z|^{k+1-\frac{1}{p}}}\left(e^{-\varepsilon^{\prime}|\Im z|}+o(1)\right)$ as $\varepsilon^{\prime} \rightarrow 0^{+}$uniformly in $z \in \mathbb{C}$.

In the sum in the r.h.s. of (28), the functions $f_{a-\varepsilon}$ and $f_{3}$ are contributing for $l=0$. Writing

$$
\frac{e^{2|\Im z|(a-\varepsilon)}}{|\Im z|^{l}} \leq \frac{e^{2|\Im z| a}}{|\Im z|^{k+1-\frac{1}{p}}} e^{-|\Im z| \varepsilon^{\prime}} e^{-|\Im z| \varepsilon}|\Im z|^{k+1-\frac{1}{p}-l}
$$

for all $\varepsilon^{\prime} \leq \varepsilon$ and since $e^{-|\Im z| \varepsilon}|\Im z|^{k+1-\frac{1}{p}-l} \leq C_{l}$ for all $z \in \mathbb{C}$ and for some $C_{l}$ depending on $l$ (and $\varepsilon$ ) one sees that the sum in the r.h.s. of (28) is $o\left(\frac{e^{2|\Im z| a}}{|\Im z|^{k+1-\frac{1}{p}}}\right)$ as $\varepsilon^{\prime} \rightarrow 0^{+}$uniformly in $z \in \mathbb{C}$. These two points complete the proof of Proposition 1.4.

## 3 Appendix A

The case of Dirichlet boundary conditions

$$
u(0)=0, \quad u(1)=0
$$

corresponding to $h=H=\infty$ may be considered similarly to the case of finite $h$ and $H$. Let us also mention that the cases $(h=\infty, H<\infty)$ and $(h<\infty, H=\infty)$ may be not treated analogously entirely, the reason being that the leading term in the asymptotic expansions of the square roots of the sequences of eigenvalues is (up to the $\pi$ factor) an half-integer whereas it is an integer for the cases $(h<\infty, H<\infty)$ and $(h=\infty, H=\infty)$ and one cannot follow [AFR, Section 4].

Theorem 3.1. Under the hypotheses of Theorem 1.1 with $H=\infty$ one concludes that, $h_{1}=h_{2}=\infty$ and $q_{1}=q_{2}$ on $[0,1]$.

Proof of Theorem 3.1: it is a direct modification of the proof of Theorem 1.1 when setting

$$
f(z)=z^{2} \int_{0}^{a}\left(\psi\left(x, z, q_{1}\right) \psi\left(x, z, q_{2}\right)\right)\left(q_{1}(x)-q_{2}(x)\right) d x
$$

where $\psi(\cdot, z, q)$ defined on $[0,1]$ is the solution to $-\frac{d^{2} \psi}{d x^{2}}+q \psi=z^{2} \psi$ with the initial conditions $\psi(0, z, q)=0$, $\psi^{\prime}(0, z, q)=1$.

Therefore we emphasize here on the main changes comparing to the case ( $h<\infty, H<\infty$ ). Note that the missing factor $\frac{1}{2}$ in the definition of $f$ comes from the fact that $\int_{0}^{1} q(x) d x$ is a spectral invariant in the Dirichlet case. The sequence of eigenvalues is denoted by $\left(\lambda_{j}(q)\right)_{j \geq 1}$.

The point of adding the $z^{2}$ factor in the definition of $f$ is the following. On one side, 0 is now a supplementary zero of order two for the function $f$ compensating the missing eigenvalue $\lambda_{0}(q)$. Furthermore
the leading term in the asymptotic expansion of $\left(\sqrt{\lambda_{j}(q)}\right)$ is the same. In particular, when setting

$$
S^{\frac{1}{2}}=\left\{0,0, \pm \sqrt{s_{j}}, j \geq 1\right\}
$$

the same analysis as in [AFR, Section 4] holds and one obtains exactly the same estimates on $N_{S^{\frac{1}{2}}}$ as the ones in [AFR, Section 4]. On the other side, the functions $\psi\left(\cdot, z, q_{1}\right)$ and $\psi\left(\cdot, z, q_{2}\right)$ are written using a transformation operator with kernel $L_{1}$ and $L_{2}$ starting from $\frac{\sin z x}{z}$ instead of $\cos z x$. These two factors $\frac{1}{z}$ vanish with the added $z^{2}$ factor.

Moreover, the kernel $L_{j}(j=1,2)$ is simpler since it is essentially (up to a change of sign) the same kernel as in the case of finite $h_{j}$ and $H$ with $h_{j}=0$ (see [M]). Therefore the results concerning the regularity properties of the kernels involved in Section 2 are unchanged and Proposition 1.4 holds for the function $f$ defined above.

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