# Asymptotic Inverse Spectral Problem for Anharmonic Oscillators 

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Dedicated to V. P. Gurarii on his $50^{\text {th }}$ birthday


#### Abstract

We study perturbations $L=A+B$ of the harmonic oscillator $A=\frac{1}{2}\left(-\partial^{2}+x^{2}-1\right)$ on $\mathbb{R}$, when potential $B(x)$ has a prescribed asymptotics at $\infty, B(x) \sim|x|^{-\alpha} V(x)$ with a trigonometric even function $V(x)=\sum a_{m} \cos \omega_{m} x$. The eigenvalues of $L$ are shown to be $\lambda_{k}=k+\mu_{k}$ with small $\mu_{k}=O\left(k^{-\gamma}\right)$, $\gamma=1 / 2+1 / 4$.

The main result of the paper is an asymptotic formula for spectral fluctuations $\left\{\mu_{k}\right\}$, $$
\mu_{k} \sim k^{-\gamma} \widetilde{V}(\sqrt{2 k})+c / \sqrt{2 k} \quad \text { as } \quad k \rightarrow \infty
$$


whose leading term $\widetilde{V}$ represents the so-called "Radon transform" of $V$,

$$
\tilde{V}(x)=\operatorname{const} \sum \frac{a_{m}}{\sqrt{\omega_{m}}} \cos \left(\omega_{m} x-\pi / 4\right)
$$

as a consequence we are able to solve explicitly the inverse spectral problem, i.e., recover asymptotic part $|x|^{-\alpha} V(x)$ of $B$ from asymptotics of $\left\{\mu_{k}\right\}_{1}^{\infty}$.

The standard spectral problem for a perturbation $L=A+B$ of a differential operator $A$ with the given spectrum $\left\{\lambda_{k}(A)\right\}_{1}^{\infty}$ asks to (approximately) calculate the eigenvalues of $L$ in terms of $\left\{\lambda_{k}(A)\right\}$ and the perturbation. For a "relatively small" perturbation $B$, the $k^{\text {th }}$ eigenvalue of $L$ is

$$
\lambda_{k}(L)=\lambda_{k}(A)+\mu_{k}
$$

so one is asked to calculate spectral fluctuations $\left\{\mu_{k}\right\}_{1}^{\infty}$. The corresponding inverse problem is then to recover $B(x)$ from the given (admissible) sequence of eigenvalues $\left\{\lambda_{k}(\mathrm{~L})\right\}_{1}^{\infty}$ or fluctuations $\left\{\mu_{k}\right\}_{1}^{\infty}$.

Spectral problems were extensively studied in various contexts for both ordinary and partial differential operators. The best known example is the regular Sturm-Liouville problem: $L=\frac{d^{2}}{d x^{2}}+V(x)$ on [0, 1]. The old result of Borg [Bo]
gives the following asymptotics of $\lambda_{k}$,

$$
\begin{equation*}
\lambda_{k} \sim(\pi k)^{2}+\int_{0}^{1} V d x+V_{2 k}+O\left(\frac{1}{k}\right) \tag{1}
\end{equation*}
$$

where $V_{2 k}$ is the $2 k^{\text {th }}$ Fourier coefficient of $V$.
Of course (1) by itself does not provide sufficient data for the inverse problem. The latter was resolved in [Bo] in the "two-spectra" setting: $\left\{\lambda_{k}\right\}_{1}^{\infty} ;\left\{\lambda_{k}^{\prime}\right\}_{1}^{\infty}$, which corresponds to two different values of the boundary-value parameter $h$ : $u+\left.h u^{\prime}\right|_{x=0}=0$.

It turned out that typically the inverse Sturm-Liouville problem does not have a unique solution. Large isospectral families of potentials $V$ exist both in the periodic (Floquet) case [La, MM], where they are given by nonlinear KdV-type evolutions; and also for two-point boundary-value problems [IMT], where the isospectral families are characterized in terms of certain norming constants.

Turning to multivariable problems, i.e., Schrödinger operators $-\Delta+V(x)$, we shall mention two known examples: the " $n$-torus" [ERT] and " $n$-sphere" theory [We, Gu, Ur, Wi].

In both cases there exist natural "isospectral deformations" of $V$ arising from symmetries of the problem: rotations on $S_{n}$, translations and reflections of $T^{n}$. Those were conjectured by Guillemin to be the only isospectral families, so-called "Spectral rigidity problem." This conjecture was proven for "generic potentials" on $T^{n}$ [ERT] and for some classes of spherical harmonics on $S_{n}[\mathrm{Gu}]$.

The method of [Gu] was based on two sets of spectral invariants: the classical "heat-invariants" of Munakshisundaram-Pleijel $\left\{b_{m}(V)\right\}_{m=0}^{\infty}$, obtained by expanding the "heat-kernel" $\operatorname{tr}\left(e^{-t L}\right)$ in powers of $t$, and a new class of spectral invariants, so called band-invariant, introduced Kac-Spencer and Weinstein [We].

Let us briefly outline some basic concepts and results of the $S_{n}$-theory.
The spectrum of the Laplacian $-\Delta$ on $S_{n}$ is well known: the $k^{\text {th }}$ eigenvalue $\lambda_{k}=k(k+n-1)$ has multiplicity $d_{k}=O\left(k^{n-1}\right) ; k=0,1, \ldots$ Introducing perturbation $V(x)$ destroys the underlying rotational symmetry. So each multiple eigenvalue $\lambda_{k}$ splits into the cluster $\left\{\lambda_{k j}=\lambda_{k}+\mu_{k j}\right\}_{j=1}^{d_{k}}$ of simple eigenvalues of $L=-\Delta+V(x)$, whose distribution is described by the probability measure

$$
\begin{equation*}
d \varrho_{k}=\frac{1}{d_{k}} \sum_{j} \delta\left(t-\mu_{k j}\right) \tag{2}
\end{equation*}
$$

As the size of the $k^{\text {th }}$ cluster increases, one is interested in the asymptotic behavior of $\left\{d \varrho_{k}\right\}$ as $k \rightarrow \infty$.

It turned out [We], that the sequence $\left\{d \varrho_{k}\right\}_{1}^{\infty}$ converges to a continuous measure $\beta_{0} d t$ on $\mathbb{R}$ expressed in terms of the Radon transform $\widetilde{V}$ of $V$,

$$
\left\langle f ; \beta_{0}\right\rangle=\iint_{S^{*}\left(S_{n}\right)} f \circ \tilde{V},
$$

integration over the cosphere bundle of $S_{n}$.
Moreover, the following asymptotic expansion analogous to Borg's (1) was derived by Weinstein,

$$
\begin{equation*}
d \varrho_{k} \sim \beta_{0}+\frac{1}{k} \beta_{1}+\frac{1}{k^{2}} \beta_{2}+\ldots \tag{3}
\end{equation*}
$$

Distributions $\left\{\beta_{0} ; \beta_{1}, \ldots\right\}$ form a new class of spectral invariants, called bandinvariants in [We]. The 0 -th invariant $\beta_{0}$ has an immediate implication to the inverse problem. Namely, the Radon transforms of two isospectral potentials $V_{1}$; $V_{2}$ are equally distributed on $S^{*}\left(S_{n}\right)$. If we know the Radon transform $\widetilde{V}$ itself, (rather than its distribution) we could easily solve the inverse problem. But a multivariable function cannot be recovered from its distribution, so $\beta_{0}$ by itself is not enough (cf. [Gu]).

The situation becomes different, however, in the context of the present work, namely perturbations $L=A+B(x)$ of the quantum mechanical harmonic oscillator $A=\frac{1}{2}\left(-\delta^{2}+x^{2}-1\right)$ on $\mathbb{R}$. We loosely call such operators anharmonic oscillators.

The harmonic oscillator is one of few examples (along with Laplacians on $T^{n}$ or $S_{n}$ ) whose spectrum can be explicitly calculated: the $k^{\text {th }}$ eigenvalue $\lambda_{k}=k$ and the eigenfunction $\varphi_{k}=k^{\text {th }}$ Hermite function $e^{x^{2} / 2} \partial^{k}\left(e^{-x^{2}}\right)$.

Our purpose in the present work is to establish the analogue of Borg's formula (1) for operators $L$ and relate its first term to the so-called "Radon transform" of the perturbation. This relation (Theorem 1) enables us to efficiently solve the inverse problems for $A+B$ in the asymptotic context, namely, to link directly "large $x$ asymptotics of $B(x)$ " on one hand and "large $k$-asymptotics of spectral fluctuations $\mu_{k}$ " on the other.

According to our basic "asymptotic" philosophy we shall consider a class of perturbations described by their behavior at $\infty$,

$$
\begin{equation*}
B(x) \sim|x|^{-\alpha} V(x) \quad \text { as } \quad x \rightarrow \infty, \tag{4}
\end{equation*}
$$

with a trigonometric even function

$$
\begin{equation*}
V(x)=\sum a_{m} \cos \omega_{m} x \tag{5}
\end{equation*}
$$

So the input data for the direct problem consists of an exponential $\alpha$ in the algebraic factor as well as frequencies $\left\{\omega_{m}\right\}$ and Fourier coefficients $\left\{a_{m}\right\}$ of the trigonometric part. The inverse problem is then to recover this data (or a portion of it) from asymptotics of spectral fluctuations $\left\{\mu_{k}\right\}_{1}^{\infty}$. One can show that perturbation $B(4)$ is small (compact) relative to $A([\mathrm{RS}])$, so the $k^{\text {th }}$ eigenvalue of $L=A+B$, is

$$
\lambda_{k}=k+\mu_{k} \quad \text { with "small" } \mu_{k} \text {. }
$$

Our main result is the following
Theorem 1. Let $L=A+B$ be the anharmonic oscillator with an even potential $B(x) \sim|x|^{-\alpha} V(x)$ whose trigonometric part $V(x)=\sum a_{m} \cos \omega_{m} x$. The $k^{\text {th }}$ spectral fluctuation of $L$ is asymptotic to

$$
\begin{equation*}
\mu_{k} \sim k^{-\gamma} \tilde{V}(\sqrt{2 k})+\frac{c_{\alpha}}{\sqrt{2 k}} \quad \text { as } \quad k \rightarrow \infty \tag{6}
\end{equation*}
$$

where constants $\gamma=\alpha / 2+1 / 4$,

$$
c_{\alpha}=\cos \left[\frac{\pi}{2}(1-\alpha)\right] \Gamma(1-\alpha) \sum a_{m} \omega_{m}^{\alpha-1}
$$

and $\widetilde{V}$ denotes the "Radon transform" of the trigonometric part:

$$
\begin{equation*}
V \rightarrow \tilde{V}(x)=\sqrt{2 / \pi} \sum \frac{a_{m}}{\sqrt{\omega_{m}}} \cos \left(\omega_{m} \mathrm{x}-\pi / 4\right) \tag{7}
\end{equation*}
$$

Notice that (7) corresponds to a formal fractional derivative operation applied to $V$

$$
\tilde{V}=\frac{1}{\sqrt{\pi}}|\partial|^{-1 / 2}(I-\partial /|\partial|)[V]
$$

We call it "Radon transform" for no other reason than it formally resembles the Radon transform on $S_{n}[\mathrm{We}]$ and plays a similar role in our discussion.

Formula (8) gives the leading coefficient $b_{0}(k)=\tilde{V}(\sqrt{2 k})$ is the asymptotic expansion of $\mu_{k}\left(\right.$ or $\left.\mu_{k}-\frac{c}{\sqrt{k}}\right)$, namely

$$
\begin{equation*}
\mu_{k} \sim k^{-\gamma}\left(b_{0}+\frac{1}{k} b_{1}+\ldots\right) \tag{8}
\end{equation*}
$$

Other coefficients $b_{1}, b_{2} ; \ldots$ could also be calculated, (cf. [Ur]) but we shall not pursue it here.

Expansion (8) is analogous to Borg's (1) or Weinstein (3). One notable difference however is the $k$-dependence (oscillatory behavior) of coefficients $b_{0} ; \ldots$ determined by function $V(x)$.

The latter provides a crucial link to the inverse problem. Precisely, given an admissible sequence of fluctuations $\mu_{k} \sim k^{-\gamma} F(\sqrt{k})$ with a trigonometric function $F(x)$, we proceed in three steps.

1) Exponential $\gamma=\frac{\alpha}{2}+\frac{1}{4}$ (consequently $\alpha=2 \gamma-1 / 2$ ) can be found as an upper bound

$$
\begin{equation*}
\gamma=\sup \left\{p>0: \lim _{k \rightarrow \infty} k^{p} \mu_{k}=0\right\} \tag{9}
\end{equation*}
$$

2) Assuming periodicity (or quasiperiodicity) of $F(x)$ we can reconstruct the period $T$ (or quasiperiods $T_{1} \ldots T_{n}$ ). Notice that the sequence $\{\sqrt{k}\}$ is dense (and uniformly distributed [KN]) modulo any $T>0$ or a tuple $T_{1} ; \ldots T_{n}$.

So any continuous function $F(x)$ on the torus $T=[0, T]$ or $T^{n}=\left[0, T_{1}\right] \times \ldots \times\left[0, T_{n}\right]$ is uniquely determined by its values at $\{\sqrt{k}\}$. Moreover, the oscillation of $F$ on any subsequence

$$
N(x, \varepsilon)=\left\{j:|x-\sqrt{j}|<\varepsilon \bmod \left(T_{1} ; \ldots T_{n}\right)\right\}
$$

has to diminish as $\varepsilon \rightarrow 0$, i.e.,

$$
\begin{gathered}
O(x, \varepsilon)=\sup \left\{\left|i^{\gamma} \mu_{i}-j^{\gamma} \mu_{j}\right|: \text { all pairs } i, j \in N(x, \varepsilon)\right\} \rightarrow 0 \\
\text { as } \varepsilon \rightarrow 0
\end{gathered}
$$

Thus the effective reconstruction procedure will screen all values of $T$ (or a tuple $T_{1} \ldots T_{n}$ ) at which the oscillation $O(x ; \varepsilon)$, as a function of parameters $T\left(T_{1} \ldots T_{n}\right)$, drops down as illustrated in Fig. 1.


Fig. 1. Oscillation as a function of $T$ in the periodic case
3) Once the periods $\left\{T_{1} ; T_{2} ; \ldots\right\}$ are found (or prescribed) we can reconstruct a trigonometric function $F(x)$, i.e., its frequencies and Fourier coefficients.

Namely, we write

$$
F(x)=\sum b_{m} \cos \left(\omega_{m} x-\pi / 4\right)=\sum\left(\frac{b_{m}}{\sqrt{2}} \cos \omega_{m} x+\frac{b_{m}}{\sqrt{2}} \sin \omega_{m} x\right)
$$

$m$ being a tuple of integers $\left(m_{1}, \ldots, m_{n} \ldots\right)$ and $\omega_{m}$ meaning the corresponding frequency

$$
\begin{equation*}
\omega_{m}=2 \pi \sum_{i} \frac{m_{i}}{T_{i}} \tag{10}
\end{equation*}
$$

From the uniform distribution property of $\left\{\sqrt{k} ; \bmod \left(T_{1} ; \ldots T_{n} \ldots\right)\right\}$ we recover the $m^{\text {th }}$ coefficient as

$$
\begin{equation*}
b_{m}=\sqrt{2} \lim _{k \rightarrow \infty} \frac{1}{k} \sum_{1}^{k} j^{\gamma} \mu_{j} \cos \left(\omega_{m} / \sqrt{j}\right) \tag{11}
\end{equation*}
$$

Let us remark that $\cos \omega_{m} x$ could be replaced with $\sin \omega_{m} x$ in (11), as sin and $\cos$ Fourier coefficients of $F$ must be equal by the definition of the Radon transform (Theorem 1) for any even perturbation $B$. This provides an additional compatibility condition on the admissible spectral data $\left\{\mu_{k}\right\}_{1}^{\infty}$.

Finally, the trigonometric part $V$ of perturbation $B$ is obtained by inverting the Radon transform

$$
F \rightarrow V(x)=\sqrt{\frac{\pi}{2}} \sum b_{m} \sqrt{\omega_{m}} \cos \omega_{m} x
$$

where frequencies $\left\{\omega_{m}\right\}$ and coefficients $\left\{b_{m}\right\}$ are given by (10), (11).
We shall summarize the inversion procedure in the following.
Corollary. (i) The admissible spectral data for the inverse problem consists of sequences

$$
\mu_{k} \sim k^{-\gamma} F(\sqrt{k})
$$

with a trigonometric function $F(x)$ whose $\sin$ and $\cos$ Fourier coefficients (11) are equal.
(ii) The inverse problem has a unique solution $B(x) \sim|x|^{-\alpha} \sum a_{m} \cos \omega_{m} x$, whose parameters: $\alpha=2 \gamma-1 / 2$; frequencies $\left\{\omega_{m}\right\}$ and coefficients $\left\{a_{m}=\sqrt{\omega_{m}} b_{m}\right\}$ are given by (9)-(11).

Remark 1. The above uniqueness result can be compared to [IMS] and [MT]. The first paper showed that isospectral classes of the Sturm-Liouville problem have typically a unique even representative, whereas the second derived the same result for the isospectral class of the harmonic oscillator $\frac{d^{2}}{d x^{2}}+x^{2}$.

Our corollary extends these results to "asymptotic" isospectral classes of the harmonic oscillator perturbed by even potentials $B(x)$.

In the rest of the paper we shall outline the proof of Theorem 1. Our argument is based on the "averaging method" of Weinstein [We], whose origins go back to the classical work on celestial mechanics. Precisely, we observe that the spectrum of the harmonic oscillator $A=\frac{1}{2}\left(-\partial^{2}+x^{2}-1\right)$ consists of integers $\{0 ; 1 ; \ldots\}$ and replace perturbation $B$ by the average of its conjugates

$$
\begin{equation*}
\bar{B}=\frac{1}{2 \pi} \int_{0}^{2 \pi} e^{i t A} B e^{-i t A} d t \tag{12}
\end{equation*}
$$

So instead of $L=A+B$ we study operator $\bar{L}=A+\bar{B}$. The main advantage of averaging is that both terms $A$ and $\bar{B}$ now commute. So spectral fluctuations $\left\{\bar{\mu}_{k}\right\}$ become nothing but eigenvalues of $\bar{B}$, which greatly facilitates calculations.

In order to pass from $\left\{\mu_{k}\right\}$ to $\left\{\bar{\mu}_{k}\right\}$ we need to show that spectra of $L$ and $\bar{L}$ become approximately equal as $k \rightarrow \infty$. The reason for the asymptotic proximity of spectra is "almost unitary" equivalence of $L$ and $\bar{L}$. Namely,

Lemma (cf. [We]). (i) There exists a skew symmetric operator $Q$ so that

$$
\begin{equation*}
e^{Q}(A+B) e^{-Q}=A+\bar{B}+" \text { small remainder } R . " \tag{13}
\end{equation*}
$$

(ii) The remainder $R$ satisfies the following operator inequality (in the sense of comparison of selfadjoint operators)

$$
\begin{equation*}
|R|=\left(R^{*} R\right)^{1 / 2} \leqq c A^{-(1 / 4+\alpha)} \quad \text { with constant } \quad c>0 \tag{14}
\end{equation*}
$$

From the lemma we immediately get an estimate of proximity of eigenvalues $\lambda_{k}=k+\mu_{k}$ and $\bar{\lambda}_{k}=k+\bar{\mu}_{k}$,

$$
\begin{equation*}
\left|k^{\alpha / 2+1 / 4}\left(\mu_{k}-\bar{\mu}_{k}\right)\right| \leqq c k^{-\alpha / 2} \rightarrow 0 \quad \text { as } \quad k \rightarrow \infty \tag{15}
\end{equation*}
$$

So large $k$-asymptotics of $\mu_{k}$ and $\bar{\mu}_{k}$ are equal modulo small (higher order) error

$$
\mu_{k}-\bar{\mu}_{k}=O\left(k^{-\gamma-\alpha / 2}\right)
$$

To prove the lemma and to calculate the leading asymptotics of $\left\{\bar{\mu}_{k}\right\}_{1}^{\infty}$ we shall use a form of pseudodofferential calculus to be introduced now.

Symbol classes $S^{m}(-\infty<m<\infty)$ consist of smooth functions $\sigma(x, \xi)$ on the phase-plane $\{(x, \xi)\}=\mathbb{R}^{2}$, which at large radius $r=\sqrt{x^{2}+\xi^{2}}$ admit an asymptotic expansion

$$
\sigma \sim \sum_{j=0}^{\infty} a_{j} r^{m-j} e^{i \omega_{j} r}
$$

Coefficients $\left\{a_{j}\right\}$ and phase-factors $\left\{\omega_{j}\right\}$ are assumed to depend smoothly on the polar angle $\theta=\arccos x / r$. One example is the potential

$$
B(x) \sim|x|^{-\alpha} V(x)=r^{-\alpha}|\cos \theta|^{-\alpha} \sum a_{j} e^{i r \omega_{j} \cos \theta}
$$

which belongs to $S^{-\alpha}$.
The harmonic oscillator $A$ has an elliptic symbol $a(r, \theta)=r^{2} / 2$, whose fractional powers can serve to "gauge" operators of classes $S^{m}$. With each symbol $\sigma(x, \xi)$ we associate a pseudodifferential (or Fourier integral) operator $K=\sigma(x, D) ; D=-i \partial_{x}$, defined by the Weyl convention

$$
(K u)(x)=\frac{1}{2 \pi} \iint e^{i \xi \cdot(x-y)} \sigma\left(\frac{x+y}{2} ; \xi\right) u(y) d \xi d y .
$$

Two basic results of pseudodifferential calculus will be used here.
Proposition. (i) Operators $B=b(x, D)$ of order zero, $\left(b \in S^{0}\right)$ are $L^{2}$-bounded. Consequently $b \in S^{-m}$ implies

$$
\begin{equation*}
\left\|B A^{m / 2} u\right\| \leqq \text { const }\|u\|, \quad \text { equivalently } \quad|B|=\left(B^{*} B\right)^{1 / 2} \leqq C A^{-m / 2} \tag{16}
\end{equation*}
$$

(ii) Product (composition) of two operators $B_{1} \in S^{m_{1}}$ and $B_{2} \in S^{m_{2}}$ has " $A$-order" $=m_{1}+m_{2}$, in the sense of (16), namely

$$
\left|B_{1} B_{2}\right| \leqq \operatorname{const} A^{\left(m_{1}+m_{2}\right) / 2}
$$

The first statement follows from the Calderon-Vaillancourt Theorem and the fact that our class $S^{0}$ is included in the standard class $S_{0,0}=\left\{b(x, \xi):\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} b\right| \leqq C_{\alpha \beta}\right.$ all $\alpha, \beta\}$ (see [Ta, Chap. 13]).

The proof of the second statement is somewhat longer and will be outlined in the Appendix.

Now we proceed to the lemma.
The intertwining operator $Q$ is constructed following [We] as

$$
\begin{equation*}
Q=\frac{1}{2 \pi i} \int_{0}^{2 \pi}(2 \pi-t) B(t) d t, \tag{17}
\end{equation*}
$$

where $B(t)$ denotes the conjugate $e^{i t A} B e^{-i t A}$. Both operators $\bar{B}(12)$ and $Q(17)$ are $\psi D O$ 's of classes $S^{m}$, whose principal symbols can be calculated from symbols of conjugates $\{B(t)\}$. The latter according to the so-called Egorov Theorem is obtained by composing "symbol $B$ " with the Hamiltonian flow $\left\{\exp t H_{a}\right\}$ of symbol $A, a(x, \xi)=\frac{1}{2}\left(x^{2}+\xi^{2}\right)$, i.e.,

$$
\sigma_{B(t)}=\sigma_{B} \circ \exp t H_{a},
$$

where $H_{a}=x \partial_{\xi}-\xi \partial_{x}$. Then

$$
\sigma_{B}=\frac{1}{2 \pi} \int_{0}^{2 \pi} \sigma_{B} \circ \exp t H_{a} d t=\frac{1}{2 \pi} \int_{0}^{2 \pi} B(x \cos t+\xi \sin t) d t
$$

which yields in polar coordinates $(r, \theta)$

$$
\sigma_{B}=\sigma_{B}(r)=\frac{1}{2 \pi} \int_{0}^{2 \pi} B(r \cos t) d t .
$$

Remembering that $B \sim|x|^{-\alpha} V(x)$ we can write $\sigma_{B}$ as

$$
\begin{equation*}
\sigma_{B}=r^{-\alpha} \frac{2}{\pi} \int_{0}^{\pi / 2} \cos ^{-\alpha} t V(r \cos t) d t \tag{18}
\end{equation*}
$$

Given a trigonometric function $V(x)=\sum a_{m} \cos \omega_{m} x$ asymptotics of integral (18) at large $r$ is computed by the stationary phase method. Both end-points of integral (18) contribute to its asymptotics: critical point, $t=0$, and singular point, $t=\pi / 2$.

The contribution of the first $(t=0)$ is

$$
r^{-\alpha-1 / 2} \sqrt{\frac{2}{\pi}} \sum \frac{a_{m}}{\sqrt{\omega_{m}}} \cos \left(\omega_{m} r-\pi / 4\right)=r^{-\alpha-1 / 2} \tilde{V}(r)
$$

the trigonometric series representing the Radon transform $\widetilde{V}$ of $V$ as defined in Theorem 1.

The second point ( $t=\pi / 2$ ) contributes

$$
\frac{c_{\alpha}^{\prime}}{r} \sum a_{m} \omega_{m}^{\alpha-1}=\frac{c_{\alpha}}{r}, \quad \text { with constant } \quad c_{\alpha}=\int_{0}^{\infty} \cos z \frac{d z}{z^{\alpha}}=\cos \left[\frac{\pi}{2}(1-\alpha)\right] \Gamma(1-\alpha)
$$

Combining two contributions we get

$$
\begin{equation*}
\sigma_{B}=r^{-\alpha-1 / 2} \tilde{V}(r)+\frac{c_{\alpha}}{r} \tag{19}
\end{equation*}
$$

Formula (19) will be essential in calculating the eigenvalues of $\bar{B}$. It also shows that the symbol of $\bar{B}$ belongs to our class $S^{-\alpha-1 / 2}$.

Similarly one calculates the principal symbol of the $\psi D O Q,(17)$,

$$
\begin{aligned}
\sigma_{Q}(r, \theta) & =\frac{1}{2 \pi i} \int_{0}^{2 \pi} B(r \cos (t-\theta))(2 \pi-t) d t \\
& =\frac{2 \pi}{i} \sigma_{B}-\frac{1}{2 \pi i} \int_{0}^{2 \pi} B(r \cos (t-\theta)) t d t
\end{aligned}
$$

whose order is also $-(\alpha+1 / 2)$.
To establish the intertwining relation (13) for $e^{Q}$ we start with an easily verified commutation formula

$$
[A ; Q]=\frac{1}{2 \pi} \int_{0}^{2 \pi}(2 \pi-t) i B^{\prime}(t) d t=i(\bar{B}-B)
$$

Application of standard algebra yields

$$
\begin{equation*}
A e^{Q}-e^{Q} A=i(\bar{B}-B)+\frac{i}{2}\{(\bar{B}-B) Q+Q(\bar{B}-B)\}+R_{1} \tag{20}
\end{equation*}
$$

The remainder $R_{1}=\sum_{3}^{\infty} \frac{1}{n!} \sum_{j=0}^{n-1} Q^{n-1-j}(\bar{B}-B) Q^{j}$ in (20) consists of products of $\psi D O$ 's, $Q$, and $\bar{B}-B$, which belong to our classes $S^{-m}$. We can apply the Proposition to calculate the " $A$-order" in the sense of (16) of each product $Q^{n-1-j}(\bar{B}-B) Q^{j}$.

Notice that the second statement of the proposition extends from two to any number of factors $B_{1} B_{2} \ldots B_{n}$, so

$$
\left|Q^{n-1-j}(\bar{B}-B) Q^{j}\right| \leqq C A^{-s} ; \quad \text { with } \quad s=(n-1)(\alpha+1 / 2)+\alpha
$$

We continue algebraic transformations and rewrite (20) as

$$
\begin{equation*}
A e^{Q}-e^{Q} A=e^{Q} \bar{B}-B e^{Q}+R \tag{31}
\end{equation*}
$$

The new remainder $R=\frac{1}{2}[\bar{B}-B ; Q]+\ldots$ has " $A$-order" $=2 \alpha+1 / 2$.
Now the first statement of the proposition (Calderon-Vaillancourt) applies to show that
which proves the lemma.

$$
|R| \leqq c A^{-(\alpha+1 / 4)}
$$

The main lemma reduces the problem of asymptotics of $\left\{\mu_{k}\right\}$ to asymptotics of eigenvalues $\left\{\bar{\mu}_{k}\right\}$ of the average operator $\bar{B}$. To complete the proof of Theorem 1 it remains to observe that the principle symbol (19) of operator $\bar{B}$ is

$$
\sigma_{B}=r^{-\alpha-1 / 2} \widetilde{V}(r)+\frac{c_{\alpha}}{r}
$$

In other words operator $\bar{B}$ represents a function of the operator $\sqrt{2 A}$, which is approximately

$$
\bar{B} \approx(2 A)^{-\gamma} \widetilde{V}(\sqrt{2 A})+c_{\alpha}(2 A)^{-1 / 2}=\bar{B}_{0} .
$$

The error term $\bar{B}-\bar{B}_{0}$ has a lower " $A$-order" according to the proposition. Therefore, the $k^{\text {th }}$ eigenvalue $\bar{\mu}_{k}$ of $\bar{B}$, consequently the $k^{\text {th }}$ fluctuation $\mu_{k}$ of $L$, is approximated by

$$
(2 k)^{-\gamma} \tilde{V}(\sqrt{2 k})+\frac{c_{\alpha}}{\sqrt{2 k}}
$$

as was claimed in Theorem 1.
Remark 2. Asymptotic formula (6) for $\mu_{k}$ also yields the limiting distribution of "average" spectral fluctuations. This result was obtained in the earlier version of our work [Gur]. Namely, by analogy with (2) we introduced measures

$$
\begin{equation*}
d \varrho_{k}(t)=\frac{1}{k} \sum_{1}^{k} \delta\left(t-j^{\gamma} \mu_{j}\right), \quad \text { or } \quad \frac{1}{k} \sum_{1}^{k} \delta\left(t-j^{\gamma}\left(\mu_{j}-\frac{c_{\alpha}}{\sqrt{j}}\right)\right) \quad \text { if } \quad \alpha>1 / 2 \tag{22}
\end{equation*}
$$

The weight factors $\left\{j^{\gamma}\right\}$ in (22) take into account the algebraic rate of decay of $\mu_{k}=O\left(k^{-\gamma}\right)$.

Then we have proved the following
Theorem 2. Sequence d $\varrho_{k}$ converges to a continuous measure $\beta_{0}(t) d t$ on $\mathbb{R}$ ( $0^{\text {th }}$ band invariant) whose density $\beta_{0}(t)$ is equal to the distribution function of the "Radon transform" $\widetilde{V}$, considered on the torus $T$ or $T^{n}$.

Here the trigonometric part $V(x)$ is assumed to be periodic or quasiperiodic. This result follows immediately from Theorem 1 and the equidistribution property [KN] of sequence $\{\sqrt{k}\}$ modulo any period $T$ or a tuple $\left(T_{1} ; \ldots T_{n} \ldots\right)$. Indeed, as $k \rightarrow \infty$,

$$
\left\langle f ; d \varrho_{k}\right\rangle \approx \frac{1}{k} \sum_{1}^{k} f \circ \tilde{V}(\sqrt{2 j}) \rightarrow \frac{1}{T} \int_{0}^{T} f \circ \tilde{V} d x \quad \text { or } \quad \frac{1}{T_{1} \ldots T_{n}} \int_{0}^{T_{1}} \cdots \int_{0}^{T_{n}} f \circ \tilde{V},
$$

But Theorem 2 can also be established directly in the following steps (see [Gur]):
(i) Replacing fluctuations $\left\{\mu_{j}\right\}$ of $A+B$ in (22) by eigenvalues $\left\{\bar{\mu}_{j}\right\}$ of $\bar{B}$ via the lemma, i.e., approximating $d \varrho_{k} \approx d \bar{\varrho}_{k}=\frac{1}{k} \sum_{1}^{k} \delta\left(t-j^{\gamma} \bar{\mu}_{j}\right)$.
(ii) Interpreting $\left\langle f ; d \bar{\varrho}_{k}\right\rangle$ as $\frac{1}{k} \operatorname{trace}\left[f\left(A^{\gamma} \bar{B}\right) \mid E_{k}\right]\left[\right.$ or $A^{\gamma}\left(\bar{B}-c_{\alpha} A^{-1 / 2}\right)$ in case $\alpha \geqq 1 / 2]$. Here operator $f\left(A^{\gamma} B\right)$ is restricted on the linear span $E_{k}$ of the first $k$ eigenfunctions of $A$ (Hermite functions).
(iii) Applying the Szego limit theorem to approximate trace $\left[f\left(A^{y} \bar{B}\right) \mid E_{k}\right]$ by the phase-space integral $\iint_{\frac{1}{2}\left(x^{2}+\xi^{2}\right) \leqq k}$ symbol $f\left(A^{\gamma} \bar{B}\right)$, and finally
(iv) Explicitly calculating the principal symbol of $\bar{B}$ as in (19).

## Appendix. Proof of the Proposition

We want to show that the product of two $\psi D O^{\prime}$ 's, $B$, and $B^{\prime}$, in classes $S^{m}$ and $S^{m^{\prime}}$ has " $A$-order" $=m+m^{\prime}$, i.e.,

$$
\begin{equation*}
\left|B B^{\prime}\right| \leqq c A^{1 / 2\left(m+m^{\prime}\right)} \tag{A1}
\end{equation*}
$$

The standard way to establish (A1) would be a composition (product) formula for Weyl symbols

$$
\begin{equation*}
\operatorname{symbol}\left(B B^{\prime}\right)=b \# b^{\prime} \sim \sum_{k=0}^{\infty} \frac{1}{k!}\left\langle b\left(\frac{\overleftarrow{\partial}_{x} \vec{\partial}_{\xi}-\vec{\partial}_{\xi} \overleftarrow{\partial}_{x}}{2 i}\right)^{k} b^{\prime}\right\rangle . \tag{A2}
\end{equation*}
$$

Here we adopt the convention of equipping the derivative operations $\partial_{x}, \partial_{\xi}$ with arrows that indicate which of two functions, $b$ or $b^{\prime}$, is subjected to it. Unfortunately expansion (A2) does not apply to oscillatory symbols $b=r^{m} e^{i r \omega}$ of classes $S^{m}$, since differentiations $\partial_{x}, \partial_{\xi}$ do not reduce their $r$-order.

The proper "composition rule" of such oscillatory symbols should involve the whole machinery of "Fourier integral operators." However, we are able to circumvent the ensuing technical difficulties by introducing a form of "abstract symbolic calculus" based on the operator $A$. Let us note that (A2) still makes sense (consists of decreasing-order terms) if one of two symbols is of "classical type,"

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} b\right| \leqq c_{\alpha \beta} r^{m-\alpha-\beta} ; \quad r=\sqrt{x^{2}+\xi^{2}},
$$

in particular, $b=r^{s}=\operatorname{symbol} A^{s / 2}$.
Moreover, the commutator of two operators $B \in S^{m}$ (of classical type) and $B^{\prime} \in S^{m^{\prime}}$ belongs to $S^{m+m^{\prime}-1}$. Indeed, the principal symbol of $\left[B ; B^{\prime}\right]$ is equal by (A2) to the Poisson bracket

$$
\left\{b ; b^{\prime}\right\}=\left\langle b\left(\frac{\overleftarrow{\partial}_{x} \vec{\partial}_{\xi}-\vec{\partial}_{x} \overleftarrow{\partial}_{\xi}}{i}\right) b^{\prime}\right\rangle
$$

which in polar coordinates becomes

$$
\frac{1}{r}\left(b_{r} b_{\theta}^{\prime}-b_{\theta} b_{r}^{\prime}\right) \in S^{m+m^{\prime}-1}
$$

It follows, in particular, that

$$
\begin{equation*}
\left\|A^{s-m / 2+1 / 2}\left[A^{-s} ; B\right]\right\| \leqq C, \text { for all } B \text { in } S^{m} \tag{A3}
\end{equation*}
$$

Guided by (16) and (A3) we shall introduce classes $\mathscr{S}_{m}$ of (compact) operators on the Hilbert space $L^{2}(\mathbb{R})$ that satisfy
(a) $\left\|B A^{m / 2}\right\| \leqq C$, consequently $|B|=\left(B^{*} B\right)^{1 / 2} \leqq C A^{-m / 2}$,
(b) iterated commutators $B_{k}=\left[A^{-s_{1}}\left[A^{-s_{2}} \ldots\left[A^{-s_{k}} ; B\right] \ldots\right]\right.$ satisfy $\left\|B_{k} A^{-m / 2-s-k / 2}\right\| \leqq C_{k}$ for all $k=1,2, \ldots$, where $s=s_{1}+s_{1}+\ldots+s_{k}$. Condition (b) essentially means that each commutation operation $B \rightarrow\left[A^{-s} ; B\right]$ sends class $\mathscr{S}_{m}$ into $\mathscr{S}_{m+s}$. We have already shown that class $\mathscr{S}_{m}$ contains all operators in $S^{-m}, m$ playing the role of the "formal $A$-order" of such $B$.

The advantage of embedding $S^{-m}$ into larger classes $\mathscr{S}_{m}$ is that the latter are easy to multiply, maintaining the right order. Namely,

In particular,

$$
\begin{equation*}
\mathscr{S}_{m} \cdot \mathscr{S}_{m^{\prime}} \subseteq \mathscr{S}_{m+m^{\prime}} . \tag{A4}
\end{equation*}
$$

$$
\left|B \cdot B^{\prime}\right| \leqq \text { const } A^{-1 / 2\left(m+m^{\prime}\right)}
$$

as was claimed in part (ii) of the proposition.
To demonstrate (A4) we observe that property (a) in the definition of $\mathscr{S}_{m}$ is equivalent to

$$
\|B u\| \leqq C\left\|A^{-m / 2} u\right\| \quad \text { for all } \quad u \in L^{2}(\mathbb{R})
$$

Choosing any pair $B \in \mathscr{S}_{m} ; B^{\prime} \in \mathscr{S}_{m^{\prime}}$ we need to show
( $\left.\mathrm{a}^{\prime}\right)\left\|B B^{\prime} u\right\| \leqq \mathrm{const}\left\|A^{-\frac{m+m^{\prime}}{2}} u\right\|$,
( $\left.\mathrm{b}^{\prime}\right)\left\|\left[A^{-s_{1}} \ldots\left[A^{-s_{k}} ; B B^{\prime}\right] \ldots\right] u\right\| \leqq C_{k}\left\|A^{-\frac{m+m^{\prime}}{2}-s-k / 2} u\right\|$, with $s=s_{1}+\ldots s_{k}$. To demonstrate ( $\mathrm{a}^{\prime}$ ) we write

$$
\begin{equation*}
\left\|B B^{\prime} u\right\| \leqq C\left\|A^{-m / 2} B^{\prime} u\right\| \leqq C\left(\left\|B^{\prime} A^{-m / 2} u\right\|+\left\|\left[A^{-m / 2} ; B^{\prime}\right] u\right\|\right), \tag{A5}
\end{equation*}
$$

and then apply (a) and (b) with $k=1$ to both norms in (A5).
Similarly ( $\mathrm{b}^{\prime}$ ) with $k=1$ is established by writing $\left[A^{-s} ; B B^{\prime}\right]=\left[A^{-s} ; B\right] B^{\prime}$ $+B\left[A^{-s} ; B^{\prime}\right]$, then estimating each of two terms as in (A5). For instance

$$
\begin{aligned}
\left\|B\left[A^{-s} ; B^{\prime}\right] u\right\| & \leqq C\left\|A^{-m / 2}\left[A^{-s} ; B^{\prime}\right] u\right\| \leqq C\left(\left\|\left[A^{-s} ; B^{\prime}\right] A^{-m / 2} u\right\|\right. \\
& +\left\|\left[A^{-m / 2} ;\left[A^{-s} ; B^{\prime}\right]\right] u\right\| .
\end{aligned}
$$

In a similar fashion one estimaties higher commutators $\left[A^{-s_{1}} \ldots\left[A^{s_{k}} ; B B^{\prime}\right] \ldots\right]$. Thus $B B^{\prime}$ belong to $\mathscr{S}_{m+m^{\prime}}$, which proves the proposition.

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## References

[Bo] Borg, G.: Acta Math. 78, 1-96 (1946)
[Gu] Guillemin, V.: Adv. Math. 42, 283-290, (1981); 27, 273-286 (1978)
[Gur] Gurarie, D.: Band invariants and inverse spectral problems for anharmonic oscillators. Preprint 1986
[ERT] Eskin, G., Ralston, A., Trubowitz, E.: Contemp. Math. 27, 45-56 (1984)
[IMT] Isaacson, E., McKean, H., Trubowitz, E.: Commun. Pure Appl. Math. 36, 767-783 (1983); 37, 1-11 (1984)
[KN] Kuipers, L., Niederreiter, H.: Uniform distribution of sequences. New York: Wiley 1974
[La] Lax, P.: Commun. Pure Appl. Math. 28, 147-199 (1975)
[MM] McKean, H.P., van Moerbecke, P.: Invent. Math. 30, 217-274 (1975)
[RS] Reed, M., Simon, B.: Methods of modern mathematical physics, Vol. II. New York: Academic Press 1975
[Ta] Taylor, M.E.: Pseudodifferential operators. Princeton, NJ: Princeton University Press 1981
[Ur] Uribe, A.: J. Funct. Anal. 59, 3 535-556 (1984)
[We] Weinstein, A.: Duke Math. J. 44, 83-92 (1977)
[Wi] Widom, H.: J. Funct. Anal. 32, 139-147 (1979)

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