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1. INTRODUCTION

In this paper we consider one-dimensional Schrödinger operators, both on the real line \( \mathbb{R} \) and on the lattice \( \mathbb{Z} \). In the first case, they act on \( L^2(\mathbb{R}) \) and have the form

\[
(Hu)(t) = -\frac{d^2}{dt^2} u(t) + V(t)u(t),
\]

while in the second case they act on \( \ell^2(\mathbb{Z}) \) and have the form

\[
(Hu)_n = u_{n+1} + u_{n-1} + V(n)u_n.
\]

We are interested in the so-called ergodic case, where one considers a measured family of potentials defined by the evaluation of a sampling function along the orbits of a dynamical system. Thus, in the first (continuum) case, we have \( V(t) = v(F_t(x)) \), where \( F_t \) is an ergodic flow and in the second (discrete) case \( V(n) = v(f^n(x)) \), where \( f \) is an ergodic invertible map. We denote the implied fixed probability measure by \( \sigma \). We will also assume below that flows, maps, and sampling functions are continuous in some compact phase space \( X \) and that \( \text{supp} \sigma = X \).

By general reasoning, the spectrum of ergodic operators is almost surely constant. In general, the spectral measure is not almost surely independent of \( x \in X \), but the ac part of the spectral measure is. There is much work dedicated to the understanding of the ac part of the spectral measure, with most results so far pointing to very rigid behavior [K1], [DeS], [CJ] (see also [R] for recent developments regarding non-ergodic potentials). Two natural problems in this direction are

**Problem 1.1.** Does the existence of an absolutely continuous component of the spectrum (for almost every \( x \in X \)) imply that the potential is almost periodic?

We recall that an almost periodic potential is one that can be obtained by evaluating a continuous sampling function along an ergodic translation of a compact Abelian group. Another way of formulating this is that the dynamics has a topological almost periodic factor, through which the sampling function factorizes.

**Problem 1.2.** Are all eigenfunctions bounded, for almost every energy, with respect to the ac part of the spectral measure (for almost every \( x \in X \))?

Here, by an eigenfunction associated to energy \( E \) we mean a generalized solution of \( Hu = Eu \) (i.e., without the requirement of belonging to \( L^2(\mathbb{R}) \) or \( \ell^2(\mathbb{Z}) \)).
For the first problem, an affirmative answer has been explicitly conjectured in the discrete case, in what is now known as the Kotani-Last Conjecture (recently popularized in [D], [J], and [S], see also the earlier [KK]).

For the second problem, and also in the discrete case, the affirmative answer would be a particular case of the so-called Schrödinger Conjecture (see [V], §1.7), according to which eigenfunctions should be bounded for almost every energy with respect to the ac part of the spectral measure, for any (possibly non-ergodic) potential. We note that in the continuum case, the corresponding general statement was known to be false if one does not assume the potential to be bounded: The famous counterexample in [MMG] is indeed unbounded both from above and from below (it is also sparse, hence non-ergodic). As it turns out, in such a setup the absolutely continuous spectrum is not constrained by any strict “Parseval-like” bound on the average size of eigenfunctions\(^1\). Such a bound (10) is a key difficulty in our setup, and it is what makes it more similar to another situation of interest (square-integrable potentials), as we will discuss.

**Remark 1.3.** Here is one example of how those conjectures could be used to deduce further regularity properties. It is known that, almost surely in the essential support of ac spectrum, there is a pair of linearly independent complex-conjugate eigenfunctions \(\varphi, \overline{\varphi}\), satisfying \(|\varphi(t)| = U(F_t(x))\) or \(|\varphi_n| = U(f^n(x))\), according to the setting, with \(U : X \to (0, \infty)\) some \(L^2\)-function (depending on \(E\) but independent of \(x\)), see [DeS]. If the dynamics is minimal then the boundedness of the eigenfunctions implies that \(U\) is in fact continuous \(\mathcal{Y}\), so if the dynamics is almost periodic then the absolute value of these eigenfunctions is itself almost periodic.

**Remark 1.4.** There are of course many examples of almost periodic potentials with the ac spectrum, dating from the KAM based work of Dinaburg-Sinai [DS]. KAM approaches do tend to produce bounded eigenfunctions. In [AFK], it has been proved that if \(f\) is an irrational rotation of the circle, and \(v\) is analytic, then (up to taking some sufficiently deep renormalization) a non-standard KAM scheme converges almost everywhere in the essential support of the ac part of the spectral measure, so that eigenfunctions are indeed bounded as predicted by the Schrödinger Conjecture.

In this paper we give negative answers to both problems, in both the discrete and the continuous setting.

**Theorem 1.** There exists a uniquely ergodic map, a sampling function, and a positive measure set \(\Lambda \subset \mathbb{R}\) such that for almost every \(x\), \(\Lambda\) is contained in the essential support of the absolutely continuous spectrum, and for every \(E \in \Lambda\) and almost every \(x\), any non-trivial eigenfunction is unbounded\(^2\).

**Theorem 2.** There exists a weak mixing uniquely ergodic map and a non-constant sampling function such that the spectrum has an absolutely continuous component for every \(x\).

\(^1\)Particularly, moments of growth do not have to be spread out according to the energy, and in fact in [MMG] many eigenfunctions become simultaneously large (in short bursts).

\(^2\)Notice that, by general reasoning, for any \(E\) in the spectrum there exists always some \(x\) with a one-dimensional subspace of bounded eigenfunctions. On the other hand, if the dynamics is minimal, then the existence of an unbounded eigenfunction for some \(x\) implies that there are unbounded eigenfunctions for every \(x\) (with the same \(E\)).
Theorem 3. There exists a uniquely ergodic flow and a sampling function, such that the spectrum is purely absolutely continuous for almost every \( x \), and for almost every energy in the spectrum, and almost every \( x \), any non-trivial eigenfunction is unbounded.

Theorem 4. There exists a weak mixing uniquely ergodic flow and a non-constant sampling function, such that the spectrum is purely absolutely continuous for every \( x \).

We recall that weak mixing means the absence of a measurable almost periodic factor. In particular, potentials associated to non-constant sampling functions are never almost periodic.

Remark 1.5. One may wonder whether there is some natural condition (stronger than lack of almost periodicity) on the dynamics that would prevent the existence of the ac spectrum. After we announced in 2009 the earliest result of this work (a less precise version of Theorem 2), Svetlana Jitomirskaya asked us whether weak mixing would be such a condition. Though our original (unpublished) construction did not yield weak mixing, the underlying mechanism could indeed be used to answer her question, as shown in the argument we present here.

Remark 1.6. Unbounded eigenfunctions can appear with or without almost periodicity: the example provided in the proof of Theorem 3 can be shown to be weak mixing (though it is not done here), while the example provided in the proof of Theorem 4 is almost periodic. In the other direction, the proofs of Theorems 2 and 4 (see Remarks 6.5 and 4.8) show that bounded eigenfunctions are also compatible with weak mixing.

Our methods do give considerably more control on the continuum case (in that we get control on the entire spectrum). The arguments are also much simpler. For this reason, we first develop all arguments in full detail for the continuum case. We then describe more leisurely the additional complications involved in the discrete case.

Remark 1.7. Discrete one-dimensional Schrödinger operators are often considered within the general context of Jacobi matrices (self-adjoint three-diagonal operators). In this more general setting, inverse spectral theory is a powerful technique which allows one to build operators from spectral data. Very recently, the inverse spectral theory approach has been used by Volberg and Yuditskii [VY] to produce a class of ergodic Jacobi matrices with purely absolutely continuous spectrum which are not almost periodic (indeed, within a certain class of spectral data they describe explicit necessary and sufficient conditions for almost periodicity). In particular, they obtain a new class of counterexamples to a natural generalization of the Kotani-Last Conjecture. Their method is completely independent from ours (which is direct), and it would be very interesting to understand the ergodic properties of such examples. Of course, as it is usually very difficult to know when the inverse spectral theory produces Schrödinger operators, it is unclear whether such an approach can be used to obtain counter-examples to the actual Kotani-Last Conjecture.

1.1. Further perspective. Besides its natural interest in the theory of orthogonal polynomials and one-dimensional Schrödinger operators, much of the motivation behind the Schrödinger Conjecture lies in its interpretation as a generalization of the
sought after non-linear version of Carleson’s Theorem about pointwise convergence of the Fourier series of an \(L^2(\mathbb{R}/\mathbb{Z})\)-function. Recall that this theorem (which solved Luzin’s Conjecture) is equivalent to the statement that for any sequence of complex numbers \(\{\lambda_n\}_{n \in \mathbb{N}}\) with \(\sum |\lambda_n|^2 < \infty\), and for almost every \(\theta \in \mathbb{R}\), the series \(\sum \lambda_n e^{2\pi i n \theta}\) is bounded.

One simple formulation of a (conjectural) non-linear version of Carleson’s Theorem goes as follows: for any sequence of \(\text{SL}(2, \mathbb{R})\) matrices \(\{A_j\}_{j \in \mathbb{N}}\) such that

\[
\sum (\ln \|A_j\|)^2 < \infty
\]

and for almost every \(\theta \in \mathbb{R}\), the sequence \(A^{(n)}(\theta) = A_n R_\theta \cdots A_1 R_\theta\) is bounded (here \(R_\theta\) is the rotation of angle \(2\pi \theta\)). To see the connection, notice that it is enough (by polar decomposition) to consider the case where \(A_j = D e^{\lambda_j} R_{\beta_j}\), where \(\beta_j \in \mathbb{R}\) is arbitrary and \(\lambda_j \geq 0\) satisfy \(\sum |\lambda_j|^2 < \infty\). Expand

\[
\begin{pmatrix}
e^{\lambda_j} & 0 \\
0 & e^{-\lambda_j}
\end{pmatrix} = \sum_{k \geq 0} \frac{\lambda_k}{k!} \begin{pmatrix}1 & 0 \\
0 & (-1)^k
\end{pmatrix},
\]

and then expand the product to get

\[
A^{(n)}(\theta) = \sum_{m \geq 0} B_{m,n}(\theta),
\]

where the coefficients of \(B_{m,n}\) are homogeneous polynomials of degree \(m\) on the \(\lambda_j\). Then a direct computation gives, with \(\alpha_j = \sum_{j' \leq j} \beta_{j'}\), \(B_{0,n} = R_{n\theta + \alpha_n}\) (which thus has unit norm), while

\[
B_{1,n} = R_{n\theta + \alpha_n} \sum_{j=1}^{n} \lambda_j \begin{pmatrix} \cos 4\pi(\alpha_j + j\theta) & -\sin 4\pi(\alpha_j + j\theta) \\
-\sin 4\pi(\alpha_j + j\theta) & -\cos 4\pi(\alpha_j + j\theta) \end{pmatrix},
\]

so that Carleson’s Theorem is equivalent to the boundedness of the \(B_{1,n}\).

One reason to hope for the almost sure boundedness of the sequence \(A^{(n)}(\theta)\) is the validity of an analogue of Parseval’s Theorem: Taking \(N(A) = \ln \frac{\|A\| + \|A\|^{-1}}{2}\) (which is asymptotic to \((\ln \|A\|)^2\) when \(\|A\|\) is close to 1), we get

\[
\int N(A^{(n)}(\theta))\,d\theta = \sum_{j=1}^{n} N(A_j).
\]

This presents a quite strict constraint to the construction of any counterexample. Indeed, (7) implies that \(\|A^{(n)}(\theta)\|\) is often bounded: Any growth one sees in a certain moment must be compensated later. This oscillation is rather hard to achieve in the non-linear setting: The product of two large \(\text{SL}(2, \mathbb{R})\) matrices is usually even larger than each factor, unless there is some rather precise alignment between their polar decompositions. However, any such alignment would appear likely to be destroyed when \(\theta\) changes. (Another way to see the difficulty is to recall that the Brownian motion in the hyperbolic plane \(\text{SL}(2, \mathbb{R})/\text{SO}(2, \mathbb{R})\) diverges linearly, while in the real line it is recurrent.)

There is nothing sacred about the above setup (which we chose to start with only for the transparency of the various formulas), and there are many alternative settings where a non-linear analogue of Carleson’s Theorem is expected to hold. The basic theme to keep in mind is the goal of showing almost sure boundedness of square-summable perturbations of a parametrized infinite product of elliptic
matrices, in some setting where some analogue of Parseval’s Theorem holds. For a more detailed discussion (with slightly different setup), see the work of Muscalu-Tao-Thiele [MTT1].

1.1.1. Schrödinger setting. The eigenfunctions of discrete Schrödinger operators [2] with eigenvalue $E$ satisfy a second-order difference equation which can be expressed in matrix form

$$A(E, n, n + 1) \cdot \begin{pmatrix} u_n \\ u_{n-1} \end{pmatrix} = \begin{pmatrix} u_{n+1} \\ u_n \end{pmatrix}$$

where

$$A(E, n, n + 1) = \begin{pmatrix} E - V(n) & -1 \\ 1 & 0 \end{pmatrix}.$$  

Thus writing $A(E, m, n) = A(E, n - 1, n) \cdots A(E, m, m + 1), n > m$, we see that the boundedness of eigenfunctions is equivalent to the boundedness of the sequences $A(E, 0, n)$ and $A(E, -n, 0)$.

It turns out that if $V \in \ell^2(\mathbb{Z})$ then the essential support of the ac spectrum is $(-2, 2)$ (a result of Deift-Killip [DeK]). This is also the set of $E$ such that the unperturbed matrices $\begin{pmatrix} E & -1 \\ 1 & 0 \end{pmatrix}$ are elliptic. Thus the expected “Carleson’s Theorem for Schrödinger operators” is just the Schrödinger Conjecture restricted to potentials in $\ell^2(\mathbb{Z})$. A partial result in this direction, under a stronger decay condition, is obtained in [CKR] (see also [MTT2] for a discussion of the limitations of this approach in the consideration of general potentials in $\ell^2(\mathbb{Z})$).

Why would one want to believe in the Schrödinger Conjecture for ergodic potentials? In our view, it is due to the validity of an inequality, reminiscent of that implied by Parseval’s Theorem in the case of decaying potentials. We state it in terms of eigenfunctions: For almost every $x \in X$, and for almost every $E$ in the essential support $\Lambda$ of the ac spectrum, there is a pair of linearly independent complex-conjugate eigenfunctions $u(E, x, n)$ and $\overline{u(E, x, n)}$, normalized so that the Wronskian $u(E, x, n)\overline{u(E, x, n - 1)} - \overline{u(E, x, n)}u(E, x, n - 1)$ is $i$, such that

$$\frac{1}{2\pi} \int_{\Lambda} |u(E, x, n - 1)|^2 + |u(E, x, n)|^2 dE \leq 1,$$

with equality in the case of the pure ac spectrum. Such an inequality is of course all that is needed to deduce a bound on the average size of the transfer matrices:

$$\frac{1}{4\pi} \int_{\Lambda} \|A(E, m, n)\| + \|A(E, m, n)\|^{-1} dE \leq 1.$$

The ergodic setup has one advantage and one disadvantage with respect to the decaying setup (as far as constructing counterexamples is concerned):

1. There is no need to “spare potential” in trying to promote eigenfunction growth,
2. But potential we introduce must reappear (infinitely often), hence (by the trend of products of large matrices to get larger) one risks promoting “too much growth”, destroying the ac spectrum due to the need to obey (11). (In other words, we have to “spare ac spectrum”.)

---

3In general, $\frac{1}{2\pi} \int_{\Lambda} |u(E, x, n - 1)|^2 + |u(E, x, n)|^2 dE$ is half the sum of the weights of the ac part of the spectral measures associated to $\delta_n$ and $\delta_{n-1}$.  

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The effects of recurrence on transfer matrices growth is hard to neglect: In particular, in the ergodic case, eigenfunctions are known [CJ] to be everywhere bounded in any open interval in the essential support of the ac spectrum (in the $\ell^2$ case, the essential support is an interval, and one certainly can not hope for boundedness everywhere).

In fact, it is quite difficult to achieve the ac spectrum in the ergodic context, which is of course what is behind the formulation of the Kotani-Last Conjecture. One of the known obstructions is Kotani’s Determinism Theorem [K1], which can be stated as follows: If there is some ac spectrum, then the stochastic process $\{V(n)\}_{n \in \mathbb{Z}}$ is deterministic in the sense that perfect knowledge of the past implies perfect knowledge of the future. The Kotani-Last Conjecture can then be seen as an optimistic quantitative generalization of this result: Almost periodicity means that approximate knowledge of $\{V(n)\}_{n \in \mathbb{Z}}$ (i.e., up to $\ell^\infty$-small error) can be obtained by sufficiently precise knowledge of the sequence in a sufficiently long (but finite) interval.

Remark 1.8. Regarding the essential support of the singular part of the spectral measure, it is well known that bounded eigenfunctions can form at most a one-dimensional subspace [Ka]. The existence of bounded potentials having no non-trivial bounded eigenfunctions was first established by Jitomirskaya in [Z] (for certain explicit ergodic potentials with singular continuous spectrum).

1.2. Principles of construction. As discussed above, a key obstacle to the construction of unbounded eigenfunctions in the absolutely continuous spectrum is the need to obey (11). In fact there are other similar constraints that must be satisfied, for instance, $\frac{1}{N} \sum_{k=0}^{N-1} \|A(E,0,k)\|$ is bounded for almost every $E$.

The need for unbounded eigenfunctions to oscillate shows that if we introduce growth, we must cancel it at some later scale. This demands very careful matching of the transfer matrices: If $x,y \in \text{SL}(2,\mathbb{R})$ then we have $\|yx\| \geq \|y\| \|x\| |\sin \omega|$ where $\omega$ is the angle between the contracting eigendirection of $y^*y$ and the expanding eigendirection of $xx^*$. So unless the eigendirections are closely aligned, if $x$ and $y$ are large then $yx$ is even larger. But as energy changes, the eigendirections move, which can easily destroy a precise match, resulting in growth for nearby eigenfunctions, which will result in losses of the ac spectrum. Notice that (11) shows also the need to spread the moments where growth occurs according to the energy, but this creates further complications regarding the interaction of the transfer matrices in those different scales.

Our approach to avoid uncontrolled growth is based on slow deformation of periodic potentials, the spectrum of which consists of bands. In order to create growth in the first place, we must introduce disorder which eats up part of the ac spectrum. In order to lose only an $\epsilon$-proportion of the ac spectrum, we must introduce (in our approach) so little disorder that the corresponding growth is of order $\epsilon$ in the bulk of the bands (this is clearly not enough to win, due to the need to spare the ac spectrum). However, we can produce slightly more growth near the edges. The disorder is introduced by slow deformation, and then we unwind it. The importance of slowness in the deformation procedure is that any introduced eigenfunction growth is also unwound. We get back to a bounded setting which

\footnote{Indeed if $u(E,x,n)$ and $\bar{u}(E,x,n)$ is a pair of complex-conjugate eigenfunctions with Wronskian $i$ then $\lim_{E \to \infty} \frac{1}{\pi N} \sum_{0 \leq k \leq N-1} |u(x,E,k)|^2$ is the derivative of the integrated density of states.}
allows us to iterate the estimates. What we see in the end is that an eigenfunction will tend to pick up oscillation at some time scales. While those oscillations are not absolutely summable, the process is so slow that their sum would still remain bounded unless there is some coherence of the phases. However, at rare random (and it is here one sees the spreading in the energy) time scales the oscillations do become coherent, so the eigenfunction does become unbounded.

We must of course be very careful in our introducing of disorder at each step. Our chosen mechanism is dictated by the setting. In the continuum, it is possible to introduce tiny amounts of rotation (for the transfer matrices), and we proceed by a large variation on the axis of rotation. This does not work in the discrete case, so we must instead create a tiny disturbance on the axis of rotation. In order to do so, we work all the time with one-parameter families of periodic potentials that remain coherent in a large part of the spectrum.

In order to construct non-almost-periodic potentials which have the ac spectrum, we consider again perturbations of periodic potentials. We construct two distinct deformations which are largely coherent, but which have slightly distinct periods (in the continuum case). Iterating each independently for a long time, they will slowly lose the coherence, until it has a definite magnitude. Later on they will become coherent again, and we can match both to construct a new periodic potential with large ac spectrum. Geometrically, the dynamics has slightly different speeds at nearby orbits of the phase space, creating macroscopic sliding in long time scales (think of the horocycle flow), though at some later time scale everything becomes periodic. Sliding is naturally incompatible with almost periodicity. Technically, the discrete case is much more delicate, since we can not produce a tiny difference of periods (as it must be an integer), so we use slow deformation along a coherent family of periodic potentials to construct coherent periodic potentials with discrepant periods.

1.3. **Structure of the paper.** Sections 2, 3 and 4 are dedicated to the continuum case.

In section 2 we provide some supporting material for the main construction, including basic concepts in spectral theory and dynamics and (variations of) known results. As the main construction involves repeated deformation of periodic potentials, we need to ensure that the crucial quantities we must control (measure of the spectrum, Lyapunov exponents, etc.) are stable under this procedure. The necessary continuity statements are straightforward and provided in section 2.5.

In section 3 we prove Theorem 3. The key part (Lemma 3.1) is to construct a suitable deformation of a periodic potential which produces eigenfunction growth without losing much measure of the spectrum. From this one concludes easily just by iterating the procedure and taking an appropriate limit. The deformation produced in Lemma 3.1 is itself obtained by applying a finite but large number of more elementary perturbations. Those more elementary perturbations provide relatively small growth of the eigenfunctions each time, which interact independently to provide large eigenfunction growth (for many energies) by a probabilistic argument. Crucially, it is necessary to show that the elementary perturbations produce much more expected eigenfunction growth than what they cost in terms of the measure of the spectrum. This is carried out in Lemma 3.2, which is the most subtle part of the argument.

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In section 4, we prove Theorem 4. The argument here is much simpler, since we only need to make use of qualitative continuity properties at the spectral side, and the dynamical analysis of the weak mixing properties for the particular sequence of deformations we consider is quite direct.

In sections 5 and 6 we discuss the discrete case. The background material here is mostly analogous to the continuum case, and introduced along the construction.

In section 5, we prove Theorem 1, trying to mimic the strategy used in section 3. While there is a conceptual difference in the implementation, the need to work with coherent parametrized families of periodic potentials (instead of a single one), the central additional difficulty with respect to the continuum case is concentrated in Lemma 5.3 (which plays the role of Lemma 3.2). Essentially we need a more involved construction of the elementary perturbation which involves a new ingredient (sliding) to produce the smooth variation of the axis of rotation of the elliptic matrices which appear.

In section 6, we prove Theorem 2. Here the ideas developed in section 4 can be adapted without significant difficulties.

2. CONTINUUM CASE: PRELIMINARIES

We will make use of the usual SL(2, \mathbb{R}) action on \mathbb{C}: \begin{pmatrix} a & b \\ c & d \end{pmatrix} \cdot z = \frac{az + b}{cz + d}.

Let \( d \) be the hyperbolic distance in the upper half-plane \( \mathbb{H} \), normalized so that \( d(i, t i) = |\ln t| \) for \( t > 0 \). The hyperbolic distance is preserved by the SL(2, \mathbb{R}) action.

Let \( R_{\theta} = \begin{pmatrix} \cos 2\pi \theta & -\sin 2\pi \theta \\ \sin 2\pi \theta & \cos 2\pi \theta \end{pmatrix} \). Note that the subgroup of rotations \( R_{\theta} \) is the stabilizer of \( i \in \mathbb{H} \) for the SL(2, \mathbb{R}) action.

If \( B \in SL(2, \mathbb{R}) \), then \( \|B\| = e^{d(i, B^{-1}i)/2} \) (up to composition with rotations one can assume that \( B \) is diagonal, for which the claim is trivial).

If \( A \in SL(2, \mathbb{R}) \) satisfies \(|\text{tr}A| < 2\), there exists a unique fixed point \( u(A) \) of \( A \) in \( \mathbb{H} \). Moreover, \( A \) is conjugated in SL(2, \mathbb{R}) to a well defined rotation \( R_{\Theta(A)} \), where \( \Theta(A) \in (0, \frac{1}{2}) \cup (\frac{1}{2}, 1) \). The conjugacy matrix \( B \), satisfying \( BAB^{-1} = R_{\Theta(A)} \) is not canonical (one may postcompose \( B \) with rotations), but can be chosen to have the form

\[
B(A) = \frac{1}{(\text{Im } u(A))^{1/2}} \begin{pmatrix} 1 & -\text{Re } u(A) \\ 0 & \text{Im } u(A) \end{pmatrix}.
\]

Notice that \( u \) and \( B \) are analytic functions of \( A \).

Given a continuous function \( V : \mathbb{R} \to \mathbb{R} \), we define the transfer matrices \( A[V](E, t, s) \in SL(2, \mathbb{R}) \) so that \( A[V](E, t, t) = \text{id} \) and

\[
\frac{d}{ds} A[V](E, t, s) = \begin{pmatrix} 0 & -E - V(s) \\ 1 & 0 \end{pmatrix} A[V](E, t, s).
\]

An eigenfunction of the Schrödinger operator with potential \( V \) and energy \( E \) is a solution \( u : \mathbb{R} \to \mathbb{R}^2 \) of \( u(s) = A[V](E,t,s) \cdot u(t) \).

We have the following basic monotonicity property:

**Lemma 2.1.** If \( s > t \) and \(|\text{tr}A[V](E, t, s)| < 2\), then

\[
\frac{d}{dE} \Theta(A[V](E, t, s)) > 0.
\]
Proof. Let us write $A(\cdot)$ for $A[V](\cdot)$. We directly compute
\begin{equation}
\partial_s (A(E, t, s)^{-1} \partial_E A(E, t, s)) = A(E, t, s)^{-1} \begin{pmatrix} 0 & -1 \\ 0 & 0 \end{pmatrix} A(E, t, s).
\end{equation}

For $w \in \mathbb{R}^2$ and $a \in \text{sl}(2, \mathbb{R})$, let $\alpha(a, w)$ be the determinant of the matrix with columns $a$ and $a \cdot w$. Clearly for $A \in \text{SL}(2, \mathbb{R})$ we have $\alpha(A^{-1} a A, w) = \alpha(a, A w)$. We say that $a$ is weakly positive if $\alpha(a, w) \geq 0$ for every $w \in \mathbb{R}^2$ and $\alpha(a, w) > 0$ for some $w \in \mathbb{R}^2$. In particular, if $a$ is weakly positive then $A^{-1} a A$ is weakly positive for every $A \in \text{SL}(2, \mathbb{R})$.

It is easy to see that if $E \mapsto A(E)$ is a $C^1$ family of $\text{SL}(2, \mathbb{R})$ matrices such that $A(E)^{-1} \partial_E A(E)$ is weakly positive, then $\partial_E \Theta(A(E)) > 0$ whenever $|\text{tr} A(E)| < 2$. Thus it is enough to show that $A(E, t, s)^{-1} \partial_E A(E, t, s)$ is weakly positive for $s > t$.

Indeed, for $s > t$ we have
\begin{equation}
A(E, t, s)^{-1} \partial_E A(E, t, s) = \int_t^s A(E, t, q)^{-1} \begin{pmatrix} 0 & -1 \\ 0 & 0 \end{pmatrix} A(E, t, q) dq.
\end{equation}
Note that $\begin{pmatrix} 0 & -1 \\ 0 & 0 \end{pmatrix}$ is weakly positive, and so $A(E, t, q)^{-1} \begin{pmatrix} 0 & -1 \\ 0 & 0 \end{pmatrix} A(E, t, q)$ is weakly positive. Using that the set of weakly positive matrices in $\text{sl}(2, \mathbb{R})$ is clearly an open convex cone, we conclude that $A(E, t, s)^{-1} \partial_E A(E, t, s)$ is weakly positive as well. \hfill \Box

2.1. Periodic case. Assume now that $V$ is periodic of period $T$. In this case we write $A[V](E, t) = A[V](E, t + T)$ and $A[V](E) = A[V](E, 0)$. Note that $\text{tr} A[V](E, t) = \text{tr} A[V](E)$ for all $t \in \mathbb{R}$.

The spectrum $\Sigma = \Sigma(V)$ of the Schrödinger operator with potential $V$ is the set of all $E$ with $|\text{tr} A[V](E)| \leq 2$. Let also $\Omega = \Omega(V)$ be the set of all $E$ with $|\text{tr} A[V](E)| < 2$. We note that $\Sigma \setminus \Omega = \partial \Omega$ consists of isolated points.

For $E \in \Omega$, let $u[V](E, t) = u(A[V](E, t))$ and $u[V](E) = u[V](E, 0)$. We note that it is more common in spectral theory to work with the so-called $m$-functions, defined initially for $E \in \mathbb{H}$ and extended to $\mathbb{R}$ by taking non-tangential limits. The function $u$ considered here (as well as related ones defined below) provide directly accessible dynamical definitions of these boundary values (for one of the $m$-functions) in the regions of the spectrum which play a role in our analysis.

The integrated density of states (i.d.s.) $N$ is absolutely continuous in this case. It satisfies
\begin{equation}
\frac{d}{dE} N(E) = \frac{1}{2\pi T} \int_0^T \frac{1}{\text{Im} u[V](E, t)} dt,
\end{equation}
for each $E \in \Omega$. We note that this formula is obtained by averaging the formula for the spectral measures of the shifts of $V$; see (11).

Note that since $A[V](E, t, s)$ conjugates $A[V](E, t)$ and $A[V](E, s)$, we have $A[V](E, t, s) \cdot u[V](E, t) = u[V](E, s)$.

In the following results, we assume $V$ to be fixed and write $A(\cdot)$ for $A[V](\cdot)$ and $u(\cdot)$ for $u[V](\cdot)$.

**Lemma 2.2.** For almost every $E \in \Sigma$, for every $\epsilon > 0$, there exists $N \in \mathbb{N}$ with the following property. Let $\tilde{u}(E, t) \in \mathbb{H}$ be some (not necessarily periodic) solution...
of $A(E, t, s) \cdot \tilde{u}(E, t) = \tilde{u}(E, s)$. Then
\begin{equation}
\frac{1}{NT} \int_0^{NT} \frac{1}{\text{Im} \, \tilde{u}(E, t)} dt > \frac{1}{T} \int_0^T \frac{1}{\text{Im} \, u(E, t)} dt - \epsilon.
\end{equation}

Proof. Let $E \in \Omega$. If $z \in \mathbb{H}$ is such that $d(z, u(E, t))$ is large then either $z$ is close to $\infty$ or $\text{Im} \, z$ is small. Moreover, if $z$ is close to $\infty$, then $A(E, t) \cdot z$ is not close to $\infty$ (as $|\text{tr} A(E, t)| < 2$, $A(E, t)$ does not have fixed points in $\overline{\mathbb{H}}$), while we obviously have $d(A(E, t) \cdot z, u(E, t)) = d(z, u(E, t))$. Thus if $d(\tilde{u}(E, t), u(E, t))$ is large, then at least one of $\frac{1}{\text{Im} \, u(E, t)}$ and $\frac{1}{\text{Im} \, u(E, t + T)}$ has to be large, so for $N \geq 2$ we have
\begin{equation}
\frac{1}{N} \sum_{j=0}^{N-1} \frac{1}{\text{Im} \, \tilde{u}(E, t + jT)} \geq \frac{1}{\text{Im} \, u(E, t)}.
\end{equation}

We now consider the case of bounded $d(\tilde{u}(E, t), u(E, t))$. Assume further that $\Theta(A(E))$ is irrational. Then for any $0 \leq t \leq T$, as $N$ grows the sequence $\tilde{u}(E, t + jT)$, $0 \leq j \leq N - 1$, is getting equidistributed in the circle of (hyperbolic radius) $d(\tilde{u}(E, t), u(E, t))$ around $u(E, t)$. One directly computes that if $\tilde{z}, \tilde{z}' \in \mathbb{H}$ are symmetric points with respect to some $z \in \mathbb{H}$, then $\frac{1}{2} (\frac{1}{\text{Im} \, \tilde{z}} + \frac{1}{\text{Im} \, \tilde{z}'}) \geq \frac{1}{\text{Im} \, z}$. It follows that for $N$ large
\begin{equation}
\frac{1}{N} \sum_{j=0}^{N-1} \frac{1}{\text{Im} \, \tilde{u}(E, t + jT)} \geq \frac{1}{\text{Im} \, u(E, t)} - \epsilon.
\end{equation}

Since this estimate is uniform on the solution $\tilde{u}(E, t)$ provided $d(\tilde{u}(E, t), u(E, t))$ is bounded, the result follows.

**Lemma 2.3.** For almost every $E \in \Sigma$, for every $t_0 \in \mathbb{R}$ we have
\begin{equation}
\inf_{w \in \mathbb{R}^2, \|w\| = 1} \sup_{t \geq t_0} \|A(E, t_0, t) \cdot w\| = \sup_{t} e^{(d(\tilde{u}(E, t, i) - d(u(E, t, i))/2}.
\end{equation}

Proof. Let us show that (21) holds whenever $E \in \Omega$ is such that $\Theta(A(E))$ is irrational.

Note that $A(E, t_0, t) = B(t)^{-1} R(t_0, t) B(t_0)$ where $B(s) = B(A(E, s))$ and $R(t_0, t) \in \text{SO}(2, \mathbb{R})$. So if $w_- \in \mathbb{R}^2$ is a unit vector such that $\|B(t_0) \cdot w_-\| = \|B(t_0)\|^{-1}$, we must have $\|A(E, t_0, t) \cdot w_-\| \leq \|B(t)\| \|B(t_0)\|^{-1}$. Since $\|B(s)\| = e^{d(u(E, s), i)/2}$, we conclude that the left hand side is not larger than the right hand side in (21).

For the other inequality, let us fix a unit vector $w \in \mathbb{R}^2$ and let $t > t_0$ be such that $\|B(t)\| = e^{d(u(E, t, i))/2}$ is maximal. Let us show that
\begin{equation}
\sup_{n \in \mathbb{N}} \|A(E, t_0, t + nT) \cdot w\| \geq \|B(t)\| \|B(t_0) \cdot w\|.
\end{equation}

Note first that $B(s)$ is $T$-periodic, and $R$ satisfies $R(s, s') = R(s, s'') R(s'', s')$ and $R(s, s + T) = R_\theta$ where $\theta = \Theta(A(E))$. Then we can write
\begin{equation}
A(E, t_0, t + nT) = B(t + nT)^{-1} R_{t_0, t + nT} B(t_0) = B(t)^{-1} R_{n\theta} R_{t_0, t} B(t_0).
\end{equation}

Since $\theta$ is irrational, we may choose $n$ such that $R_{n\theta} R_{t_0, t} B(t_0)$ is almost colinear with a unit vector $w_+$ such that $\|B(t)^{-1} \cdot w_+\| = \|B(t)\|$. Thus $\|A(E, t_0, t + nT) \cdot w\|$ is close to $\|B(t)\| \|B(t_0) \cdot w\|$, thus implying the desired estimate. \qed
2.2. Uniquely ergodic case. Let \( \{F_t : X \to X\}_{t \in \mathbb{R}} \) be a continuous flow which is minimal and uniquely ergodic with invariant probability measure \( \sigma \), and let \( v : X \to \mathbb{R} \) be a continuous function. Let \( A[F,v](E,x,t,s) = A[V](E,t,s) \) with \( V(t) = v(F_t(x)) \). Let \( \Sigma = \Sigma(F,v) \) be the corresponding spectrum (which is \( x \)-independent by minimality). Notice that if \( v \) is non-negative and non-identically vanishing, then \( \Sigma \subset (0, \infty) \).

Let \( N(E) \) be the i.d.s., and \( L(E) \) be the Lyapunov exponent, defined by
\[
L(E) = \lim_{T \to \infty} \frac{1}{T} \int \ln \|A[F,v](E,x,0,T)\|d\sigma(x).
\]

For a uniquely ergodic flow, the i.d.s. is not necessarily absolutely continuous. However, we have the following result (due to Kotani [K1]). Let \( \Sigma_0 \subset \Sigma \) be the set of \( E \) with \( L(E) = 0 \).

**Lemma 2.4.** We have \( \frac{d}{dE} N(E) > 0 \) for almost every \( E \in \Sigma_0 \). Moreover, there exists a measurable function \( u[F,v] : \Sigma_0 \times X \to \mathbb{H} \), unique up to \( \text{Leb} \times \sigma \)-zero measure sets, such that \( A[F,v](E,x,0,t) \cdot u[F,v](E,x) = u[F,v](E,F_t(x)) \). This function satisfies
\[
\frac{d}{dE} N(E) = \frac{1}{2\pi} \int \frac{1}{\text{Im} u[F,v](E,x)}d\sigma(x).
\]

Notice that when \( F \) is \( T \)-periodic, we have \( u[F,v](E,x) = u[V](E) \) for \( V = v(F_t(x)) \).

The following is due to Kotani [K2].

**Theorem 5.** Let \( F_t : X \to X \) be a uniquely ergodic flow, and let \( v : X \to \mathbb{R} \) be continuous. If the Lyapunov exponent vanishes in the spectrum and the i.d.s. is absolutely continuous, then the spectral measures are absolutely continuous for almost every \( x \in X \).

2.3. Solenoidal flows. If \( K \) and \( K' \) are compact Abelian groups, a projection \( K \to K' \) is a continuous surjective homomorphism.

Let \( K \) be a totally disconnected compact Abelian group, and let \( i : \mathbb{Z} \to K \) be a homomorphism with dense image. The solenoid associated to \( (K,i) \) is the compact Abelian group obtained as the quotient of \( K \times \mathbb{R} \) by the subgroup \( \{(i(-j),j)\}_{j \in \mathbb{Z}} \).

It comes with a canonical projection \( \pi : S \to \mathbb{R}/\mathbb{Z} \), \( \pi(x,t) = t \).

Given \( S \) as above, the solenoidal flow on \( S \) is \( F^S_t : S \to S \), \( F^S_t(x,s) = (x,t+s) \).

A time-change of \( F^S \) is a flow \( F^S_t : S \to S \) of the form \( F^S_t(x,s) = (x,h(x,s,t)) \) with \( t \mapsto h(x,s,t) \in C^1 \) for each \( x \) and \( s \), and such that \( w_F(x,s) = \frac{d}{dt} h(x,s,t) \big|_{t=0} \) is a continuous positive function of \( (x,s) \). Notice that any continuous positive function in \( S \) generates a time-change.

Notice that a time-change of a solenoidal flow is uniquely ergodic, and its invariant probability measure is absolutely continuous with respect to the Haar measure, with continuous positive density proportional to \( \frac{1}{w_F} \).

If \( (K,i) \) and \( (K',i') \) are as above, and there is a (necessarily unique) projection \( p_{K,K'} : K \to K' \) such that \( p_{K,K'} \circ i = i' \), then we can define a projection \( p_{S,S'} : S \to S' \) in the natural way.

If \( F \) and \( F' \) are time-changes of \( F^S \) and \( F^{S'} \), we say that \( F \) is \( \epsilon \)-close to the lift of \( F' \) if \( | \ln w_{F'} \circ p_{S,S'} - \ln w | \leq \epsilon \). We say that \( v : S \to \mathbb{R} \) is \( \epsilon \)-close to the lift of \( v' : S' \to \mathbb{R} \) if \( |v' \circ p_{S,S'} - v | \leq \epsilon \).
In all cases we will deal with (e.g., \( K = \mathbb{Z}/n\mathbb{Z} \)), there is a natural choice for the embedding \( i : \mathbb{Z} \to K \). Thus we will often omit the embedding \( i \) from the notation below.

2.4. **Projective limits.** An increasing sequence of compact Abelian groups is the data given by a sequence \( K_j \) of compact Abelian groups, and of projections \( p_{j',j} : K_j' \to K_j \), \( j' > j \) such that \( p_{j_2,j_1} \circ p_{j_3,j_2} = p_{j_3,j_1} \).

Given such an increasing family of compact Abelian groups, there exists a compact Abelian group \( K \) and a sequence of projections \( p_j : K \to K_j \) such that \( p_{j',j} \circ p_j = p_j \) for every \( j' > j \), and the \( p_j \) separate points in \( K \): one takes \( K \) as the set of infinite sequences \( x_j \in K_j \) with \( p_{j',j}(x_j') = x_j \), endowed with the product topology and obvious group structure. We call \( K \) the projective limit of the \( K_j \).

When considering pairs \((K_j,i_j)\) as before, we will assume moreover that the projections are compatible with the embeddings, so that \( p_{j',j} = p_{K_j,K_j'} \).

Notice that if the \( K_j \) are totally disconnected then the projective limit is totally disconnected. Moreover, if \( S_j \) is the solenoid over \( K_j \), then the projective limit \( S \) of the \( S_j \) is the solenoid over \( K \).

An immediate application of projective limits is the following:

**Lemma 2.5.** Let \( S_j \) be an increasing sequence of solenoids, and let \( S \) be the projective limit. Let \( F_j \) be time-changes of the solenoidal flows \( F^{S_j} \). Let \( v_j : S_j \to \mathbb{R} \) be continuous functions. Assume that for \( j' > j \), \((F_j',v_j')\) is \( \epsilon_j \)-close to the lift of \((F_j,v_j)\), where \( \epsilon_j \to 0 \). Then there exists a time-change \( F \) of the solenoidal flow \( F^S \), and a continuous function \( v : S \to \mathbb{R} \) such that \((F,v)\) is \( \epsilon_j \)-close to the lift of \((F_j',v_j)\) for every \( j \).

**Proof.** Define \( S \) as a projective limit of the \( S_j \) and take \( v = \lim v_j \circ p_j \), \( w_F = \lim w_{F_j} \circ p_j \). \( \square \)

2.5. **Lifting properties.** The following is a standard “semi-continuity of the spectrum” property:

**Lemma 2.6.** Let \( F' \) be a time-change of a solenoidal flow \( F^{S'} \), and let \( v' : S' \to \mathbb{R} \) be a continuous function. Then for every \( M > 0 \), \( \epsilon > 0 \), there exists \( \kappa > 0 \) such that if either \((F,v)\) is \( \kappa \)-close to the lift of \((F',v')\), or \((F',v')\) is \( \kappa \)-close to the lift of \((F,v)\), then \( \Sigma(F,v) \cap (-\infty,M] \) is contained in the \( \epsilon \)-neighborhood of \( \Sigma(F',v') \), and \( \Sigma(F',v') \cap (-\infty,M] \) is contained in the \( \epsilon \)-neighborhood of \( \Sigma(F,v) \).

It easily implies

**Lemma 2.7.** Let \( F' \) be a time-change of a solenoidal flow \( F^{S'} \), and let \( v' : S' \to \mathbb{R} \) be a continuous function. Then for every \( M > 0 \), \( \epsilon > 0 \), there exists \( \kappa > 0 \) such that if \((F,v)\) is \( \kappa \)-close to the lift of \((F',v')\), then we have \( |\Sigma(F,v) \cap (-\infty,M] \setminus \Sigma(F',v')| < \epsilon \).

We say that \((F,v)\) is \((\epsilon_1,C_1,M)\)-crooked if there is a set \( \Gamma \subset \Sigma \cap (-\infty,M] \) such that \( |\Sigma \setminus \Gamma \cap (-\infty,M]| < \epsilon_1 \), and for every \( E \in \Gamma \), the set of \( x \in X \) such that

\[
\inf_{w \in \mathbb{R}^2, \|w\|=1} \sup_{t \geq 0} \|A[F,v](E,x,0,t) \cdot w\| > C_1
\]

has \( \sigma \)-measure (strictly) larger than \( 1 - \epsilon_1 \).

This notion provides a quantification of how large the eigenfunctions are for most of the parameters. Largeness can be checked in many cases by bounding the \( u(E,x) \); see Lemma 2.3.
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The following is obvious, but convenient:

**Lemma 2.8.** Let $F'$ be a time-change of a solenoidal flow $F^S'$, and let $v' : S' \to \mathbb{R}$ be a continuous function. Assume that $(F', v')$ is $(\epsilon_1, C_1, M)$-crooked. Then there exists $\kappa > 0$ such that if $(F, v)$ is $\kappa$-close to the lift of $(F', v')$, then $(F, v)$ is $(\epsilon_1, C_1, M)$-crooked.

**Proof.** Just use that $v \circ F_t$ is close to $v' \circ F'_t \circ p_{S,S'}$ for $t$ (arbitrarily) bounded. □

We say that $(F, v)$ is $(\epsilon, M)$-good if
\begin{equation}
\sup_{\Sigma \cap (-\infty, M]} L(E) < \epsilon.
\end{equation}

It also trivially behaves well under lifts:

**Lemma 2.9.** Let $F'$ be a time-change of a solenoidal flow $F^S'$, and let $v' : S' \to \mathbb{R}$ be a continuous function. Assume that $(F', v')$ is $(\epsilon, M)$-good. Then there exists $\kappa > 0$ such that if $(F, v)$ is $\kappa$-close to the lift of $(F', v')$, then $(F, v)$ is $(\epsilon, M)$-good.

We say that $(F, v)$ is $(\epsilon, M)$-nice if
\begin{equation}
N(M) - \int_{-\infty}^{M} \frac{dN}{dE} dE < \epsilon.
\end{equation}

Niceness provides a measure of how absolutely continuous the i.d.s. is.

The following deserves an argument.

**Lemma 2.10.** Let $F'$ be a time-change of a periodic solenoidal flow $F^S'$ and let $v' : S' \to \mathbb{R}$ be a continuous function. Assume that $(F', v')$ is $(\epsilon, M)$-nice. Then there exists $\kappa > 0$ with the following property. Assume that $(F, v)$ is $\kappa$-close to the lift of $(F', v')$, the Lyapunov exponent for $(F, v)$ vanishes in the spectrum, and $|(\Sigma(F', v') \setminus \Sigma(F, v)) \cap (-\infty, M]| < \kappa$. Then $(F, v)$ is $(\epsilon, M)$-nice.

**Proof.** Let $N$, $N'$ be the integrated density of states for $(F, v)$, $(F', v')$. It is clear that $N(M)$ is close to $N'(M)$. Using Lemma 2.2 and Lemma 2.4 we see that for almost every $E' \in \Sigma(F', v')$, for every $\epsilon' > 0$, there exists $\delta > 0$ such that for almost every $E \in \Sigma$ which is $\delta$-close to $E'$, if $\kappa > 0$ is sufficiently small, we have
\begin{equation}
\frac{d}{dE} N(E) > \frac{d}{dE} N'(E') - \epsilon'.
\end{equation}

Integrating over $\Sigma(F, v) \cap (-\infty, M]$, and using that the Lebesgue measure of the spectrum is close, we get $\int_{-\infty}^{M} \frac{d}{dE} N dE$ close to $\int_{-\infty}^{M} \frac{d}{dE} N' dE'$. The result follows. □

2.6. Slow deformation. The following is (a parametrized version of) the basic estimate of Fayad-Krikorian (Proposition 2 in [FK]).

**Lemma 2.11.** Let $J \subset \mathbb{R}$ be a closed interval, and let $A : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R})$ be a smooth function such that $|\text{tr} A(E, t)| < 2$ for $(E, t) \in J \times \mathbb{R}/\mathbb{Z}$. Let $B(E, t) = B(A(E, t))$, $\theta(E, t) = \Theta(A(E, t))$. Then for every $m, k \in \mathbb{N}$, there exists $n(m) \in \mathbb{N}$

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b) This result still holds without assuming periodicity.
and $C_{k,m} > 0$ such that for every $n \geq n(m)$, there exist smooth functions $B_{(m,n)} : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R})$, $\theta_{(m,n)} : J \times \mathbb{R}/\mathbb{Z} \to \mathbb{R}$ such that

1. $\|A_{(m,n)} - R\theta_{(m,n)}\|_{C^k} \leq \frac{C_{k,m}}{n^m}$, where

   $A_{(m,n)}(E, t) = B_{(m,n)} \left( E, t + \frac{n-1}{n} \right) A(E, t) B_{(m,n)}(E, t)^{-1},$

2. $\|B_{(m,n)} - B\|_{C^k} \leq \frac{C_{k,m}}{n^m},$

3. $\|\theta_{(m,n)} - \theta\|_{C^k} \leq \frac{C_{k,m}}{n^m}.$

Proof. Consider first the case $m = 1$. In this case, set $B_{(1,n)} = B$, $\theta_{(1,n)} = \theta$, and the estimate is obvious.

Assume we have proved the result for $m \geq 1$. Let

   $B_{(m+1,n)}(E, t) = B(A_{(m,n)}(E, t)) B_{(m,n)}(E, t),$

$\theta_{(m+1,n)}(E, t) = \Theta(A_{(m,n)}(E, t)).$

The estimates follow from the induction hypothesis. \qed

Under a monotonicity assumption, it yields a parameter estimate:

Lemma 2.12. Under the hypothesis of the previous lemma, assume moreover that $\tilde{\theta}(E) = \int_{\mathbb{R}/\mathbb{Z}} \theta(E, t) dt$ satisfies $\frac{d}{dE} \tilde{\theta}(E) \neq 0$ for every $E \in J$. For $n \in \mathbb{N}$, let $A^{(n)} : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R})$ be given by

   $A^{(n)}(E, t) = A \left( E, t + \frac{n-1}{n} \right) A \left( E, t + \frac{n-2}{n} \right) \cdots A \left( E, t + \frac{1}{n} \right) A(E, t).$

Then there exist functions $\tilde{\theta}^{(n)} : J \to \mathbb{R}/\mathbb{Z}$ such that for every measurable subset $Z \subset \mathbb{R}/\mathbb{Z},$

   $\lim_{n \to \infty} \| \{ E \in J, \tilde{\theta}^{(n)}(E) \in Z \} \| = |Z||J|,$

with the following property. For every $\delta > 0$,

   $\lim_{n \to \infty} \| \text{tr} A^{(n)}(E, t) - 2 \cos 2\pi \tilde{\theta}^{(n)}(E) \|_{C^0(J \times \mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0,$

$\lim_{n \to \infty} \sup_{|\sin 2\pi \tilde{\theta}^{(n)}(E)| > \delta} \| \Theta(A^{(n)}(E, \cdot)) - \tilde{\theta}^{(n)}(E) \|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0,$

$\lim_{n \to \infty} \sup_{|\sin 2\pi \tilde{\theta}^{(n)}(E)| > \delta} \| u(A^{(n)}(E, \cdot)) - u(A(E, \cdot)) \|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{C})} = 0.$

Proof. Let $B_{(m,n)}$, $A_{(m,n)}$ and $\theta_{(m,n)}$ be as in Lemma 2.11 and let

   $A^{(m,n)}(E, t) = B_{(m,n)}(E, t) A_{(m,n)}(E, t) B_{(m,n)}(E, t)^{-1}.$

We have

   $A^{(m,n)}(E, t) = A_{(m,n)} \left( E, t + \frac{n-1}{n} \right) \cdots A_{(m,n)}(E, t).$

Let $\theta_{(m,n)}(E, t) = \sum_{j=0}^{n-1} \theta_{(m,n)}(E, t + \frac{j}{n})$. Then

   $A^{(m,n)}(E, t) - R \theta_{(m,n)} = \sum_{j=1}^{n} H_{(m,n,j)},$
with

$$H_{(m,n,i)} = \sum_{l \neq i} H_{(m,n,l)}$$

where the sum runs over all non-empty sequences $i = (i_1, \ldots, i_j)$ satisfying $0 \leq i_1 < \cdots < i_j < n$, and $H_{(m,n,l)}$ is a product $T_{n-1} \cdots T_0$ with $T_i(E,t) = R_{\theta_{(m,n)}}(E,t+i)$ if $l = i_r$ for every $1 \leq r \leq j$, and $T_j(E,t) = A_{(m,n)}(E,t+i_1) - R_{\theta_{(m,n)}}(E,t+i)$ if $l = i_r$ for some $1 \leq r \leq j$. Then

$$\|H_{(m,n,l)}\|_{C^0} \leq \|A_{(m,n)} - R_{\theta_{(m,n)}}\|_{C^0},$$

$$\|DH_{(m,n,l)}(E,t)\|_{C^0} \leq \|D(A_{(m,n)} - R_{\theta_{(m,n)}})\|_{C^0} \|A_{(m,n)} - R_{\theta_{(m,n)}}\|_{C^0}^{j-1} + (n - j)\|DR_{\theta_{(m,n)}}\|_{C^0} \|A_{(m,n)} - R_{\theta_{(m,n)}}\|_{C^0},$$

where we write $D$ for the total derivative. Thus

$$\|H_{(m,n,l)}\|_{C^1} \leq \frac{C_m^n}{n^{m_j-1}},$$

$$\|A_{(m,n)} - R_{\theta_{(m,n)}}\|_{C^1} \leq \sum_{j=1}^n \frac{C_m^n}{n^{(m-1)j-1}},$$

so that for $m \geq 3$ we have

$$\lim_{n \to \infty} \|A_{(m,n)} - R_{\theta_{(m,n)}}\|_{C^1} = 0.$$ 

Note that

$$\theta_{(m,n)}(E,t) = n \sum_{k \in \mathbb{Z}} e^{2\pi i k t} \int_{\mathbb{R}/\mathbb{Z}} \theta_{(m,n)}(E,t) e^{-2\pi i k t} dt.$$ 

Let $\tilde{\theta}_{(m,n)}(E) = \int_{\mathbb{R}/\mathbb{Z}} \theta_{(m,n)}(E,t) dt$. Then for $m \geq 3$, using that $\theta_{(m,n)}$ is uniformly $C^3$ for fixed $m$ and $n \to \infty$,

$$\lim_{n \to \infty} \sup_{E \in J} \|\theta_{(m,n)}(E, \cdot) - n \tilde{\theta}_{(m,n)}(E)\|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0.$$ 

Since

$$\lim_{n \to \infty} \|\tilde{\theta}_{(m,n)}(E) - \tilde{\theta}(E)\|_{C^1} = 0,$$ 

and the derivative of $\tilde{\theta}(E)$ is non-vanishing, it follows that for $m \geq 3$ and each measurable set $Z \subset \mathbb{R}/\mathbb{Z}$,

$$\lim_{n \to \infty} |\{(E \in J, \tilde{\theta}_{(m,n)}(E) \in Z)\}| = |Z||J|.$$ 

Fix $m \geq 3$ and let $\tilde{\theta}^{(n)} = \tilde{\theta}_{(m,n)}$. Then for $n$ large we get

$$\lim_{n \to \infty} \|\text{tr}A_{(m,n)}(E,t) - 2 \cos 2\pi \tilde{\theta}^{(n)}(E)\|_{C^0(J \times \mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0,$$

$$\lim_{n \to \infty} \sup_{|\sin 2\pi \tilde{\theta}^{(n)}(E)| > \delta} \|u(A_{(m,n)}(E, \cdot)) - i\|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{C})} = 0,$$

$$\lim_{n \to \infty} \sup_{|\sin 2\pi \tilde{\theta}^{(n)}(E)| > \delta} \|\Theta(A_{(m,n)}(E, \cdot)) - \tilde{\theta}^{(n)}\|_{C^1(\mathbb{R}/\mathbb{Z}, SL(2, \mathbb{R}))} = 0.$$ 

Notice that $\text{tr}A_{(m,n)} = \text{tr}A^{(n)}$ and $\Theta(A_{(m,n)}) = \Theta(A^{(n)})$. Moreover, $u(A^{(n)}(E,t)) = B_{(m,n)}(E,t)^{-1}u(A_{(m,n)}(E,t))$. Since

$$\lim_{n \to \infty} \|B_{(m,n)} - B\|_{C^1} = 0,$$
and $B(E, t) \cdot u(A(E, t)) = i$, we conclude that
\begin{equation}
\lim_{n \to \infty} \sup_{|\sin 2\pi \tilde{\theta}^n(E)| > \delta} \|u(A^n(E, \cdot)) - u(A(E, \cdot))\|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{C})} = 0.
\end{equation}

\[\square\]

3. Continuum case: unbounded eigenfunctions

The potential we will produce will be a suitable projective limit of periodic potentials. The actual work we need to do is to inductively construct suitable periodic potentials.

3.1. Construction of periodic potentials. Let $V : \mathbb{R}/T\mathbb{Z} \to \mathbb{R}$ be a continuous function with $V(0) = 0$. For $\delta > 0$, $N \in \mathbb{N}$, $n \in \mathbb{N}$, the $(\delta, N, n)$-padding of $V$ is the continuous function $V' : \mathbb{R}/T'\mathbb{Z} \to \mathbb{R}$, $T' = 2Nt + \delta \sum_{j=0}^{2n-1} \sin 2\pi \frac{j}{2n}$, given by the following conditions:

1. $V'(t) = V(t - a_j)$, $a_j \leq t \leq a_j + NT$, $0 \leq j \leq 2n - 1$,
2. $V'(t) = 0$, $a_j + NT \leq t \leq a_{j+1}$, $0 \leq j \leq 2n - 1$,
3. $a_0 = 0$, $a_{j+1} = a_j + NT + \delta \sin 2\pi \frac{j}{2n}$.

In other words, we repeat the periodic potential $Nn$ times, but with an extra “padding” with a small interval of zeroes every $N$ repetitions. The length of those intervals is slowly modulated, but it is always small (at most $\delta$).

The goal of this section is to establish the following estimate:

Lemma 3.1. Let $V^{(0)} : \mathbb{R}/T^0\mathbb{Z} \to \mathbb{R}$ be a smooth non-constant, non-negative function with $V^{(0)} = 0$ near 0. Then for every $M$, $\xi > 0$, there exists $C > 0$ such that for every $C_0 > 0$, for every $\delta > 0$ sufficiently small, there exist $0 < P < \xi \delta^{-1}$, and sequences $N^{(j)}$, $n^{(j)}$, $1 \leq j \leq P$, such that if we define $V^{(j)} : \mathbb{R}/T^{(j)}\mathbb{Z} \to \mathbb{R}$, $1 \leq j \leq P$ so that $V^{(j)}$ is obtained by $(\delta, N^{(j)}, n^{(j)})$-padding of $V^{(j-1)}$, then there exists a compact subset $\Gamma \subset (0, M] \cap \Omega(V^{(P)})$ such that $|\Sigma(V^{(0)}) \cap (-\infty, M] \setminus \Gamma| < \xi$ and for every $E \in \Gamma$ we have

\begin{equation}
\sup_{t} d(u[V^{(P)}](E, t), i) \geq C_0,
\end{equation}

\begin{equation}
\frac{1}{T^{(P)}} \int_{0}^{T^{(P)}} d(u[V^{(P)}](E, t), i) dt \leq C.
\end{equation}

We will need a few preliminary results.

Lemma 3.2. For every $C > 0$, $M > 0$, there exist $C' > 0$ and $\delta_0 > 0$ with the following property. Let $V : \mathbb{R}/T\mathbb{Z} \to \mathbb{R}$ be a smooth function with $V(t) = 0$ near 0 and let $A(\cdot) = A[V](\cdot)$ and $u(\cdot) = u[V](\cdot)$. Let $E_0 \in \Omega(V) \cap [M^{-1}, M]$. Assume that $C^{-1} < d(u(E_0), E_0^{1/2}i) < C$. Then there exists $\epsilon_0 > 0$ such that for every $0 < \epsilon < \epsilon_0$, for every $\kappa > 0$, for every $0 < \delta < \delta_0$, for every $N$ sufficiently large, for every $n$ sufficiently large, let $V'$ be the $(\delta, N, n)$-padding of $V$, $A'(\cdot) = A[V'](\cdot)$, $u'(\cdot) = u[V'](\cdot)$, we have the following. There exists a compact set $\Lambda \subset \Omega(V') \cap [E_0 - \epsilon, E_0 + \epsilon]$ such that

1. $|\Lambda| > 2(1 - C'\delta)\epsilon$,
2. For $E \in \Lambda$, $d(u'(E), u(E)) < \kappa$ and $C^{-1} < d(u'(E), E_0^{1/2}i) < C$.

\[\text{This neighborhood can be arbitrarily small, but this will influence the constants below that depend on } \epsilon.\]
3. For $E \in \Lambda$,
\[
\sup_{t \in [0,T']} d(u'(E, t), i) \geq \sup_{t \in [0,T]} d(u(E, t), i),
\]
where $B$ so that $1$
\[
\lambda
\]
\[
A \quad \text{and} \quad B
\]
\[
\text{such that for } E \in \Lambda',
\]
\[
\sup_{t \in [0,T']} d(u'(E, t), i) \geq \sup_{t \in [0,T]} d(u(E, t), i) + C'\delta \gamma.
\]
\[
\text{Let}
\]
\[
G(E, t) = D(E) \frac{E^{1/2}}{2\pi} \sin^{2N} \pi t D(E)^{-1} A(E)^N. \quad \text{We have}
\]
\[
A'(E) = G \left( E, \frac{2n - 1}{2n} \right) G \left( E, \frac{2n - 2}{2n} \right) \cdots G \left( E, \frac{1}{2n} \right) G(E, 0).
\]

We can write for $E$ near $E_0$,
\[
B(E) A(E) B(E)^{-1} = R_{\theta(E)},
\]
where $B(E) = B(A(E))$ and $\theta(E) = \Theta(A(E))$. By Lemma 24.1, $\theta$ has non-zero derivative.

Thus we can write
\[
G(E, t) = D(E) \frac{E^{1/2}}{2\pi} \sin^{2N} \pi t D(E)^{-1} R_{N\theta(E)} B(E).
\]

Letting $Q(E) = B(E) D(E)$, we get
\[
\text{tr} G(E, t) = \text{tr} Q(E) R_{E^{1/2}} \sin^{2N} \pi t Q(E)^{-1} R_{N\theta(E)}.
\]

Notice that $Q(E) \notin SO(2, \mathbb{R})$, since $Q(E) \cdot i \neq i$ (here we use that $B(E)^{-1} \cdot i = u(E) \neq E^{1/2}i = D(E) \cdot i$ for $E$ near $E_0$). Thus we can write $Q = R(1) D(0) R(2)$, a product of rotation, diagonal and rotation matrices, depending analytically on $E$.

Then
\[
\text{tr} G(E, t) = \text{tr} D^{(0)}(E) R_{E^{1/2}} \sin^{2N} \pi t D^{(0)}(E)^{-1} R_{N\theta(E)}.
\]

Write $D^{(0)}(E) = \begin{pmatrix} \lambda(E) & 0 \\ 0 & \lambda(E)^{-1} \end{pmatrix}$. We may assume that $\lambda(E) > 1$. Then
\[
\lambda(E) = e^{d(u(E), E^{1/2})/2},
\]
so that $\frac{1}{2\pi} < \ln \lambda(E) < \frac{C}{2}$. Then
\[
\text{tr} G(E, t) = 2 \cos((\delta E^{1/2} \sin^{2N} \pi t) + 2\pi N\theta(E))
\]
\[
- (\lambda(E) - \lambda(E)^{-1})^2 \sin(\delta E^{1/2} \sin^{2N} \pi t) \sin 2\pi N\theta(E).
\]
Thus
\[
|\text{tr} G(E, t) - 2 \cos((\delta E^{1/2} \sin^{2N} \pi t) + 2\pi N\theta(E))| \leq C_1 \delta \sin 2\pi N\theta(E).
\]
We conclude that if $2N\theta(E)$ is at distance at least $C_2\delta$ from $\mathbb{Z}$, then $|\text{tr} G(E, t)| < 2$.

We conclude that for $\epsilon$ sufficiently small, for $N$ sufficiently large, the set of
\[
E \in [E_0 - \epsilon, E_0 + \epsilon]
\]
such that $|\text{tr} G(E, t)| \geq 2$ for some $t$ has a Lebesgue measure of
at most $C_3\delta$. By Lemma 23.12 for $n$ large we will have $|\text{tr}A'(E)| < 2$ for a compact set $\Lambda(\epsilon, \delta, N, n) \subset [E_0 - \epsilon, E_0 + \epsilon]$ of a Lebesgue measure of at least $2(1 - C_4\delta)\epsilon$. We may further assume that for $E \in \Lambda(\epsilon, \delta, N, n)$, the sequence $\{j\theta(E)\}_{0 \leq j \leq N - 1}$ is $\frac{1}{100}$ dense mod 1. Thus for such $E$, and any $w \in \mathbb{H}$, and any $0 \leq t \leq T$, there exists $0 \leq k \leq N - 1$ such that $d(A(E, 0, t)A(E)^k \cdot w, i) \geq d(u(E, t), i) + \frac{1}{2}d(w, u(E))$. Taking $w = u'(E, a_m)$ for some $0 \leq m \leq 2n - 1$ (where $a_j$ is as in the definition of a $(\delta, N, n)$-padding), we get

\begin{equation}
(69) \quad d(u'(E, t + kT + a_m, i) \geq d(u(E, t), i) + \frac{1}{2}d(u'(E, a_m), u(E)).
\end{equation}

In particular,

\begin{equation}
(70) \quad \sup_{t} d(u'(E, t), i) \geq \sup_{t} d(u(E, t), i) + \frac{1}{2} \max_{0 \leq m \leq 2n - 1} d(u'(E, a_m), u(E)).
\end{equation}

Lemma 24.12 shows that $u'(E, a_m)$ is near $u(G(E, \frac{m}{2n}))$ for $n$ large. In particular, $u'(E)$ is near $u(E)$, since $G(E, 0) = A(E)$ and $u'(E, a_n)$ is near $w(E) = u(G(E, 1/2))$.

We want to estimate the hyperbolic distance between $w(E)$ and $u(E)$ in $\mathbb{H}$. Let $w'(E)$ be the fixed point of $D^{(0)}(E)R_{\frac{2}{2\pi}} D^{(0)}(E)^{-1} R_{N\theta(E)}$ in $\mathbb{H}$. Then $w(E) = B^{-1}R^{(i)} \cdot w'(E)$. Since $u(E) = B^{-1}R^{(i)} \cdot i$, it follows that

\begin{equation}
(71) \quad d(w(E), u(E)) = d(w'(E), i).
\end{equation}

But $w'(E)$ is the solution $z \in \mathbb{H}$ of the equation $az^2 + bz + c = 0$, where

\begin{align*}
a &= \cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^{-2} \sin \delta E^{1/2} \cos 2\pi N\theta(E), \\
b &= (\lambda(E)^2 - \lambda(E)^{-2}) \sin \delta E^{1/2} \sin 2\pi N\theta(E), \\
c &= \cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^2 \sin \delta E^{1/2} \cos 2\pi N\theta(E).
\end{align*}

Then

\begin{equation}
(75) \quad \Im w'(E) = \left(1 + \frac{(\lambda(E)^2 - \lambda(E)^{-2}) \sin \delta E^{1/2} \cos 2\pi N\theta(E)}{\cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^{-2} \sin \delta E^{1/2} \cos 2\pi N\theta(E)}ight)^{1/2}.
\end{equation}

Under the condition that $2N\theta(E)$ is at distance at least $C_2\delta$ from $\mathbb{Z}$, we have

\begin{align*}
(76) \quad \left| \frac{(\lambda(E)^2 - \lambda(E)^{-2}) \sin \delta E^{1/2} \cos 2\pi N\theta(E)}{4(\cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^{-2} \sin \delta E^{1/2} \cos 2\pi N\theta(E))^2} \right| &\leq C_5 \delta^2, \\
(77) \quad \left| \frac{(\lambda(E)^2 - \lambda(E)^{-2}) \sin \delta E^{1/2} \cos 2\pi N\theta(E)}{\cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^{-2} \sin \delta E^{1/2} \cos 2\pi N\theta(E)} \right| &\geq C_6^{-1} \delta \cot 2\pi N\theta(E).
\end{align*}

If $2N\theta(E)$ is at distance $C_2\delta < \gamma < 1/4$ from $\mathbb{Z}$ then

\begin{equation}
(78) \quad \left| \frac{(\lambda(E)^2 - \lambda(E)^{-2}) \sin \delta E^{1/2} \cos 2\pi N\theta(E)}{\cos \delta E^{1/2} \sin 2\pi N\theta(E) + \lambda(E)^{-2} \sin \delta E^{1/2} \cos 2\pi N\theta(E)} \right| \geq C_7^{-1} \delta \frac{1}{\gamma},
\end{equation}

so that

\begin{equation}
(79) \quad d(w'(E), i) \geq C_8^{-1} \delta \frac{1}{\gamma}.
\end{equation}
It follows that in this case
\begin{equation}
\sup_{t} u'(E, t) \geq \sup_{t} u(E, t) + C_{0}^{-1} \frac{\delta}{\gamma}.
\end{equation}

For \(\gamma < 1/4\), let \(\Lambda'(\epsilon, \delta, N, n, \gamma)\) be the set of \(E \in \Lambda(\epsilon, \delta, N, n)\) such that \(2N\theta(E)\) is at distance at most \(\gamma\) from \(\mathbb{Z}\). Since \(\theta\) has non-zero derivative, we have \(|\Lambda'(\epsilon, \delta, N, n, \gamma)| \geq \frac{3}{2} \gamma \epsilon\), for \(\epsilon\) small, \(N\) sufficiently large and \(n\) sufficiently large.

To conclude, let us show that if \(E \in \Lambda(\epsilon, \delta, N, n)\) and \(2N\theta(E)\) is \(C_{2}\)-away from \(\mathbb{Z}\), then
\begin{equation}
\frac{1}{T} \int_{0}^{T'} d(u'(E, t), i) dt - \frac{1}{T} \int_{0}^{T} d(u(E, t), i) dt
\end{equation}
is small. The formulas for \(u'\) imply that \(u'(E, t)\) is at bounded hyperbolic distance from some \(u(E, t')\). In fact, if \(a_{j} \leq t \leq a_{j+1}\) then \(u'(E, t)\) is at bounded hyperbolic distance from \(u(E, t - a_{j})\). Moreover, if \(a_{j} \leq t \leq a_{j} + T\), then \(A(0, t - a_{j})^{-1}u'(E, t)\) is near \(u(G(E, \frac{t}{2n}))\). If \(\frac{t}{2n}\) is not close to \(\frac{1}{3}\), the estimates give that the fixed point of \(G(E, \frac{t}{2n})\) is close to \(u(E)\), provided \(N\) is large. It follows that \(u'(E, t)\) is near \(u(E, t - a_{j})\). The result follows.

**Lemma 3.3.** For every \(C > 0\), \(M > 0\), there exist \(C' > 0\) and \(\delta_{0} > 0\) with the following property. Let \(V : \mathbb{R}/T\mathbb{Z} \to \mathbb{R}\) be a smooth non-negative function with \(V(t) = 0\) near \(0\). Let \(\Xi \subset \Omega(V) \cap [M^{-1}, M]\) be a compact subset such that \(C^{-1} < d(u[V](E, E^{1/2}) < C\) for every \(E \in \Lambda\). Then for every \(\kappa > 0\), \(R \in \mathbb{N}\), for every \(0 < \delta < \delta_{0}\), for every \(N\) sufficiently large, for every \(n\) sufficiently large, if \(V' : \mathbb{R}/T'\mathbb{Z} \to \mathbb{R}\) is the \((\delta, N, n)\)-padding of \(v\), then there exists a compact subset \(\Xi' \subset \Xi \cap \Omega(V')\) such that

1. For \(j \geq 0\), the conditional probability that \(E \in \Xi\) belongs to \(\Xi'\), given that \(\frac{T}{R} \leq \sup_{t} d(u[V](E, t), i) < \frac{j+1}{R}\) is at least \(1 - 2C'\).
2. For every \(E \in \Xi'\), \(d(u[V'](E), u[V](E)) < \kappa\) and \(C^{-1} < d(u[V'](E), E^{1/2}) < C\).
3. For every \(E \in \Xi'\),
\begin{equation}
\sup_{t} d(u[V'](E, t), i) \geq \sup_{t} d(u[V'](E, t), i),
\end{equation}

4. For \(j \geq 0\), and for every \(C'\delta < \gamma < 1/4\), the conditional probability that \(E \in \Xi\) belongs to \(\Xi'\) and
\begin{equation}
\sup_{t} d(u[V'](E, t), i) > \sup_{t} d(u[V](E, t), i) + C't^{-1} \frac{\delta}{\gamma},
\end{equation}
given that \(\frac{T}{R} \leq \sup_{t} d(u[V](E, t), i) < \frac{j+1}{R}\) is at least \(\frac{2}{3}\).
5. For every \(E \in \Xi'\),
\begin{equation}
\left| \frac{1}{T'} \int_{0}^{T'} d(u[V'](E, t), i) dt - \frac{1}{T} \int_{0}^{T} d(u[V](E, t), i) dt \right| < \kappa.
\end{equation}

**Proof.** Follows from the previous lemma by a covering argument. (Notice that the statements about conditional probabilities are automatic for large \(j\), since \(\sup_{t} d(u[V](E, t), E^{1/2})\) is bounded by compactness of \(\Xi\).)

**Proof of Lemma 3.3.** Notice that by non-constancy of \(V(0)\), \(u^{(0)}(E) \neq E^{1/2}\) for almost every \(E \in \Sigma(V(0))\). Up to increasing \(M\), we can assume that \(\inf \Sigma(V(0)) >
Then for sufficiently large $C > 0$, there exists a compact subset $\Xi(0) \subset \Sigma(V(0)) \cap (-\infty, M] \cap \Omega(V(0))$ such that $|\Sigma(V(0)) \setminus \Xi(0)| < \xi$, and for every $E \in \Xi(0)$ we have
\begin{equation}
C^{-1} < d(u[V(0)](E), E^{1/2}i) < C,
\end{equation}
and
\begin{equation}
\sup_t d(u[V(0)](E, t), i) < \frac{C}{2}.
\end{equation}
Let $\delta_0 = \delta_0(C, M)$ and $C' = C'(C, M)$ be as in Lemma 3.3.

Let $P$ be maximal so that $(1 - 2C')^P > 1 - \frac{\xi}{2M}$. Choose very small $0 < \delta < \delta_0$, choose $R \in \mathbb{N}$ very large (in particular, much larger than $\delta^{-1}$), and take $\kappa > 0$ very small. Define sequences $V(j)$, $\Xi(j)$, $1 \leq j \leq P$, so that $V(j)$, $\Xi(j)$ is obtained by applying Lemma 3.3 to $V(j-1)$, $\Xi(j-1)$. It follows that
\begin{equation}
|\Sigma(V(0)) \cap (-\infty, M) \setminus \Xi(0)| \leq \frac{\xi}{2} + M(1 - (1 - 2C')^P) < \xi.
\end{equation}
It also follows that
\begin{equation}
\frac{1}{T(P)} \int_0^{T(P)} d(u[V(P)](E, t), i) dt \leq \frac{C}{2} + \kappa P < C.
\end{equation}
Let $Z_j$, $0 \leq j \leq P$, be random variables on $\Xi(0)$ given as follows. If $E \notin \Xi(j)$, let $Z_j = Z_{j-1} + 1$. Otherwise, let $Z_j = \frac{l}{R}$, where $j$ is maximal with $\sup_t d(u[V(j)](E, t), i) \geq \frac{l}{R}$. Let $L \subset \mathbb{N}$ be the set of all $l$ with $4\delta C'\epsilon_1 R < l < C'\epsilon_1^2 R$. We have $Z_0 \geq 0$, and the conditional probability that $Z_j - Z_{j-1} \geq \frac{l}{R} = C'\epsilon_1 \frac{l}{R}$, given $Z_{j-1}$ is at least $\frac{\gamma}{\gamma} = \frac{\delta R}{3MC}$, provided $C'\epsilon_1 < \gamma < 1/4$, i.e. $l \in L$. Consider independent and identically distributed random variables $W_j$, $1 \leq j \leq P$, taking only values of the form $\frac{l}{R}$ with $l = 0$ or $l \in L$, and such that
\begin{equation}
p\left(W_j \geq \frac{l}{R}\right) = \frac{\delta R}{3MC'},
\end{equation}
whenever $l \in L$. Since $Z_0 \geq 0$ and $p(Z_j - Z_{j-1} \geq \frac{l}{R} | Z_{j-1}) \geq p(W_j \geq \frac{l}{R})$ for every $l \in \mathbb{Z}$, we get
\begin{equation}
p\left(Z(P) \geq \frac{m}{R}\right) \geq p\left(\sum_{j=1}^P W_j \geq \frac{m}{R}\right)
\end{equation}
for every $m \in \mathbb{Z}$.

To conclude, it is enough to show that $p(\sum_{j=1}^P W_j < C_0) < \xi/2$, provided $\delta$ is sufficiently small.

By (89), for $C'' < k < -C'' - \ln \delta$, we have $p(2^k \delta < W_j < 2^k \delta + 1) > C''^{-1}2^{-k}$ (here $C''$ is an appropriately large constant depending on $C'$, and we are using that $R$ is much larger than $\delta^{-1}$). We also have $P \geq C''^{-1}2^{-1}$ for some constant $C''$ depending on $M, \xi$ and $C'$. By the Law of Large Numbers, for each $D \in \mathbb{N}$, and each $C'' < k \leq D$, if $\delta$ is sufficiently small, then with probability at least $1 - \frac{\xi}{4D}$, we will have $2^k \delta < W_j < 2^{k+1} \delta$ for a set of 1 \leq j \leq P of cardinality at least $C''^{-1}2^{-k-1}P \geq C''^{-1}2^{-k-1}2^{-k}2^{-1}. This implies that, with probability at least $1 - \frac{\xi}{4}$, $\sum W_j \geq \frac{D - C''}{2C''2^{-k}2^{-1}}. The result follows by taking $D \geq C'' + 2C''C''C_0$. \hfill \Box
3.2. Passing to the limit.

Lemma 3.4. Let $F_t : S \to S$ be a $T$-periodic time-change of a solenoidal flow, and let $V : \mathbb{R}/T\mathbb{Z} \to \mathbb{R}$, $v : S \to \mathbb{R}$ be continuous functions such that $V(t) = v(F_t(0))$ is smooth. Assume that $V(t) = 0$ for $T - \epsilon_0 \leq t \leq T$ for some $0 < \epsilon_0 < T$. Then for $0 < \delta < \epsilon_0$, and for every $N, \epsilon \in \mathbb{N}$, the $(\delta, N, \epsilon)$-padding $V'$ of $V$ has the form $V'(t) = v'(F'_t(0))$, where $F'_t : S' \to S'$ is a $T'$-periodic time-change of a solenoidal flow, $v' : S' \to \mathbb{R}$ is continuous, and $(F', v')$ is $\frac{\delta}{\epsilon_0}$-close to a lift of $(F, v)$.

Proof. Let $U_0 = \{F_t(0), T - \epsilon_0 < t < T\}$. We take $S'$ as the $Nn$-cyclic cover of $S$. Let $U_0' \to S$ be the corresponding projection. Let $U_{0,j}$, $0 \leq j \leq 2nN - 1$, be the connected components of $U_0' = p^{-1}(U_0)$, labeled so that they are positively cyclically ordered and such that the right boundary of $U_0'_{2nN-1}$ is 0. Then there exists a continuous non-positive function $\rho : S' \to \mathbb{R}$ such that $\rho = 0$ outside $U_0'$ and any $U_{0,j}$ with $j$ not divisible by $N$, and such that

$$\int_{U_{0,j}} \frac{1}{e^{\rho(x)}w_F(p(x))} dx = \epsilon_0 + \delta \sin^{2N} \pi \frac{j}{N}$$

is equal to if $j$ is divisible by $N$. Indeed, we can take $\|\rho\|_{C^0}$ arbitrarily close to $\log \frac{\epsilon_0}{\delta}$, and hence less than $\frac{\delta}{\epsilon_0}$.

The result then follows with $v' = v \circ p$, and $w_{F'} = e^{\rho}w_F \circ p$. \qed

Lemma 3.5. Let $F_t : S \to S$ be a periodic time-change of a solenoidal flow, and let $v : S \to \mathbb{R}$ be continuous non-constant non-negative function such that $t \mapsto v(F_t(0))$ is smooth and $v(F_t(0)) = 0$ for $t$ near 0. Then for every $\epsilon_1, C_1, M, \epsilon > 0$, there exists a periodic time-change of a solenoidal flow $F' : S' \to S'$, and a continuous non-constant non-negative function such that $t \mapsto v'(F'_t(0))$ is smooth, $v'(F'_t(0)) = 0$ for $t$ near 0, $(F', v')$ is $\epsilon$-close to a lift of $(F, v)$, and $(F', v')$ is $(\epsilon_1, C_1, M)$-crooked.

Proof. Let $\epsilon_0 > 0$ be such that $v(F_t(0)) = 0$ for $T - \epsilon_0 \leq t \leq T$. Apply Lemma 3.1 to $V(0)(t) = v'(F'_t(0))$, with parameters $C_0$ and $\xi < \epsilon_0 \epsilon$ to be specified below, to get $\Gamma$ and $P$, and then Lemma 3.4 $P$ times, to get $(F', v')$ with $V'(t) = v'(F'_t(0))$ such that $(F', v')$ is $\frac{\delta}{\epsilon_0}$-close to a lift of $(F, v)$. By Lemma 2.7 if $\xi$ is small then $|\Sigma(V'(t))| \cap (-\infty, M]\Sigma(V(t))) < \epsilon_1/2$, and if additionally $0 < \xi < \epsilon_1/2$ we conclude that $|\Sigma(V'(t))| \cap (-\infty, M]\Gamma| < \epsilon_1$. Fix such $\xi$ and let $C$ be as in Lemma 3.1.

By (57), for every $E \in \Gamma$, the set $Z_E$ of all $t \in \mathbb{R}/T(P)\mathbb{Z}$ with

$$d(vV(P)(E, t), i) \leq \frac{C}{\epsilon_1}$$

has a measure of at least $(1 - \epsilon_1)T(P)$. By Lemma 28 for almost every $E \in \Gamma$, $t_0 \in Z_E$ implies that

$$\inf_{w \in \mathbb{R}^2, \|w\| = 1} \sup_{t > t_0} \|A[vV(P)](E, t_0, t)\| \geq \left(\frac{C_0 \epsilon_1}{C_1}\right)^{1/2}.$$ 

So by taking $C_0 = C_2^2C/\epsilon_1$, we get that $(F', v')$ is $(\epsilon_1, C_1, M)$-crooked. \qed

Proof of Theorem 3. Let $V(0) : \mathbb{R}/\mathbb{Z} \to \mathbb{R}$ be a smooth non-constant non-negative periodic function with $V(0)(t) = 0$ near 0. We can see $\mathbb{R}/\mathbb{Z}$ as the solenoid $S^{(0)}$ corresponding to the trivial group. Let $F^{(0)}$ be the solenoidal flow on $S^{(0)}$. 

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Let $\kappa_0 = 1$. Apply Lemma 3.5 inductively to obtain a sequence $(F^{(j)}, v^{(j)})$, $j \geq 1$ such that

1. $(F^{(j)}, v^{(j)})$ is $(2^{-j}, 2^j)$-crooked,
2. $|\Sigma(F^{(j)}, v^{(j)})| - \Sigma(F^{(j-1)}, v^{(j-1)}) \cap (-\infty, 2^j) < \kappa_j$,
3. $(F^{(j)}, v^{(j)})$ is $\kappa_j$-close to the lift of $(F^{(j-1)}, v^{(j-1)})$,
4. $\kappa_j < \kappa_{j-1}/2$ is chosen so small that if $(F, v)$ is $2\kappa_j$-close to the lift of $(F^{(j-1)}, v^{(j-1)})$, and $|\Sigma(F, v) - \Sigma(F^{(j-1)}, v^{(j-1)}) \cap (-\infty, 2^j) < 2\kappa_j$, then $(F, v)$ is $(2^{-j}, 2^j)$-nice (use Lemma 2.10), $(2^{-j}, 2^j)$-good (use Lemma 2.9), and if $j \geq 2$, $(2^{-j-1}, 2^{-1})$-crooked (use Lemma 2.8).

Let $(F, v)$ be the limit of the $(F^{(j)}, v^{(j)})$ obtained by Lemma 2.5. Then $(F, v)$ is $(2^{-j}, 2^j)$-good for all $j$, so the Lyapunov exponent vanishes in the spectrum. Moreover, $(F, v)$ is $(2^{-j}, 2^j)$-nice for all $j$, so the i.d.s. is absolutely continuous. By Theorem 3 for almost every $x \in S$ the spectral measure is purely absolutely continuous. By construction, $(F, v)$ is $(2^{-j}, 2^j)$-crooked for all $j$. Thus for almost every $E$ in the spectrum, for almost every $x \in S$, all non-trivial eigenfunctions are unbounded. \hfill \Box

4. Continuum case: breaking almost periodicity

The example discussed in the previous section can be verified to be not almost periodic. Here we will discuss a simpler example that will be easier to analyze.

4.1. Spectral measure. Given a bounded continuous function $V : \mathbb{R} \to \mathbb{R}$, we denote by $\mu_V$ the spectral measure for the line Schrödinger operator. It has some basic continuity property (see [RS], Theorem VIII.20):

**Lemma 4.1.** Let $V : \mathbb{R} \to \mathbb{R}$ be a bounded continuous function, and let $V^{(n)} : \mathbb{R} \to \mathbb{R}$, $n \in \mathbb{N}$ be a sequence of uniformly bounded continuous functions such that $V^{(n)} \to V$ uniformly on compact subsets of $\mathbb{R}$. Then $\int \phi d\mu_{V^{(n)}} \to \int \phi d\mu_V$ for every compactly supported continuous function $\phi : \mathbb{R} \to \mathbb{R}$.

We will only need explicit formulas for the spectral measure in the case of periodic potentials. Let $V : \mathbb{R}/TZ \to \mathbb{R}$ be continuous, and denote its shifts by $V_s : t \mapsto V(s + t)$. Then $\mu_{V_s}$ is absolutely continuous and

$$
\frac{d}{dE} \mu_{V_s} = \frac{1}{\text{Im } u[V](E, s)}.
$$

This formula is an immediate consequence of the expression of the spectral measure in terms of the $m$-functions (see for instance [KI]).

For $C > 0$, let $\mu_{V_s} \in$ be the restriction of $\mu_{V_s}$ to the set of $E$ with $|\frac{d}{dE} \mu_{V_s}| < C$. We say that a periodic $v$ is $(\epsilon, C, M)$-uniform if $\mu_{V_s}(-\infty, M] - \mu_{V_s}(-\infty, M] < \epsilon$ for every $s$.

We clearly have

**Lemma 4.2.** For every periodic $V$, $\epsilon > 0$, $M > 0$, there exists $C > 0$ such that $V$ is $(\epsilon, C, M)$-uniform.

\footnote{Notice that since $\Sigma(F, v)$ contains $\bigcap_{j' \geq j} \Sigma(F^{(j')}, v^{(j')})$ (see Lemma 2.10), we must have $|\Sigma(F, v) - \Sigma(F^{(j-1)}, v^{(j-1)}) \cap (-\infty, 2^j) | < 2\kappa_j$.}
4.2. Weak mixing. Let $F'_t : S' \to S'$ be a time-change of a periodic solenoidal flow. We say that a time-change $F_t : S \to S$ of a solenoidal flow is $(N,F')$-mixed if $S$ projects onto $S'$ (through $p$), and for every $1 \leq j \leq N$, there exists $t_j > 0$ and compact subsets $U_j, V_j \subset S$ with Haar measure strictly larger than $1/3$, such that for each $x \in U_j$ there exists $|t| < \frac{1}{N}$ such that $p(F_{t_j}(x)) = F'_{t_j}(p(x))$, and for each $x \in V_j$, there exists $|t - \frac{1}{N}| < \frac{1}{N}$ such that $p(F_{t_j}(x)) = F'_{t_j}(p(x))$.

**Lemma 4.3.** Let $F$ be $(N,F')$-mixed. Then there exists $\kappa > 0$ such that if $\tilde{F}$ is $\kappa$-close to the lift of $F$, then $\tilde{F}$, if $(N,F')$-mixed.

**Lemma 4.4.** Let $F$ be the projective limit of $F^{(n)}$, and assume that for every $N \in \mathbb{N}$, for every $n$ sufficiently large, $F$ is $(N,F^{(n)})$-mixed. Then $F$ is weak mixing.

**Proof.** If $F$ is not weakly mixing, there exists a non-trivial eigenfunction, i.e., a measurable function $\psi : S \to \mathbb{S}^1$ such that $\psi \circ F = e^{2\pi i \theta} \psi$ for some $\theta \in \mathbb{R} \setminus \{0\}$. Let $\psi^{(n)} : S^{(n)} \to \mathbb{R}$ be the expected value of $\psi$ on $p_{S,S^{(n)}}^{-1}(x)$ (with respect to the Haar measure on $S$). Then $\psi^{(n)} \circ p_{S,S^{(n)}}$ converges to $\psi$ almost everywhere (Martingale Convergence Theorem).

Since $\psi$ is an eigenfunction, $t \mapsto \psi^{(n)}(F_t^{(n)}(0))$ is continuous, uniformly on $t$ and $n$.

By the definition of projective limit, 

$$\lim_{n \to \infty} \sup_{x \in S^{(n)}} \sup_{0 \leq t \leq 1} |\psi^{(n)}(F_t^{(n)}(x)) - e^{2\pi i \theta} \psi^{(n)}(x)| = 0. \tag{95}$$

Thus, for $x \in U_j$, $\psi^{(n)}(p_{S,S^{(n)}}(F_t(x)))$ is close to $\psi^{(n)}(p_{S,S^{(n)}}(x))$. Since the first is close to $\psi(F_t(x))$ and the second is close to $\psi(x)$ for most $x \in U_j$, this shows that $\theta t_j$ is close to an integer.

A similar argument using $V_j$ shows that $\theta(t_j - \frac{1}{N})$ is close to an integer, so that $\theta \frac{1}{N}$ is close to an integer. Since $1 \leq j \leq N$ is arbitrary, we conclude that $\theta = 0$. \(\square\)

4.3. The construction. Let $V : \mathbb{R}/T \mathbb{Z} \to \mathbb{R}$ be a continuous function with $V(0) = 0$. For $\delta > 0$, $n \in \mathbb{N}$, the $(\delta, n)$-padding (a simplified version of a $(\delta, N, n)$-padding) of $V$ is the continuous function $V' : \mathbb{R}/T' \mathbb{Z} \to \mathbb{R}$, $T' = 2nT + \delta n$, given by the following conditions:

1. $V'(t) = V(t - a_1), a_1 \leq t \leq a_1 + T, 0 \leq j \leq 2n - 1,$
2. $V'(t) = 0, a_2 + T \leq t \leq a_2 + 1, 0 \leq j \leq 2n - 1,$
3. $a_j = jT, 0 \leq j \leq n, a_j = jT + (j - n)\delta, n + 1 \leq j \leq 2n$.

**Lemma 4.5.** Let $F, v, V, \epsilon_0$ be as in Lemma 3.4. Then for every $\delta > 0$ sufficiently small, for every $N \in \mathbb{N}$, for every $n$ sufficiently large, the $(\delta, n)$-padding $V'$ of $V$ has the form $V'(t) = v'(F_t^S(0))$, where $F'_t : S' \to S'$ is a $T'$-periodic time-change of a solenoidal flow, $v' : S' \to \mathbb{R}$ is continuous, $(F', v')$ is $(\delta, n)$-close to a lift of $(F, v)$, and $(F', v')$ is $(N,F)$-mixed.

**Proof.** Let $N_S \in \mathbb{N}$ be the period of the solenoidal flow $F_t^S$. Define $S'$ as the $2n$-cover of $S$. Define a continuous function $\rho : S \to \mathbb{R}$ supported on $\{F_t(0), T - \epsilon_0 < t < T\}$ such that $\int_0^T e^{-\rho(F_t(0))} dt = T + \delta$. As in Lemma 3.4, we can choose $\rho$ with $\|\rho\|_{C^0} < \frac{\delta}{\epsilon_0}$. Let $\rho' : S' \to \mathbb{R}$ be defined so that $\rho' = 0$ on $[0, nN_S]$ and $\rho' = \rho \circ p_{S', S}$ on $[nN_S, 2nN_S]$.

Let $F'$ be the solenoidal flow with $w_{F'} = e^{\theta w} \circ p_{S', S}$, and let $v' = v \circ p_{S', S}$. All properties, but the last one, follow as in Lemma 3.4. For the last property,
notice that if $t_j = \lceil \frac{j}{N} \rceil (T + \delta)$, then for $x \in \{F'_s(0), 0 \leq s \leq nT - t_j\}$ we have $p_{s'} \circ F'_s(x) = F'_s(\delta(Ns + \delta))$, which belongs to $\{F_s(p_{s'}(s)), \frac{1}{N}, \delta \leq s \leq \frac{1}{N}\}$, and for $x \in \{F'_s(0), nT \leq s \leq T' - t_j\}$ we have $p_{s'} \circ F'_s(x) = p_{s'}(s)$. \hfill $\square$

**Lemma 4.6.** Let $V : \mathbb{R}/\mathbb{T}$ be a continuous function with $V(0) = 0$. If $V$ is $(\epsilon, C, M)$-uniform, then for $\delta > 0$ sufficiently small, for every $n \in \mathbb{N}$, if $V'$ is the $(\delta, n)$-padding of $V$, then $V'$ is $(\epsilon, C, M)$-uniform.

**Proof.** Let $A(\cdot) = A[V](\cdot)$. Let $J \subset \Omega(V) \cap (-\infty, M]$ be a finite union of closed intervals such that

$$\sup V_s(-\infty, M] - \mu_{V_s} J < \epsilon_0 < \epsilon,$$

where $V_s$ is the shift of $V$ and $\mu_{V_s} C$ is the truncation of the spectral measure.

If $E \in \Omega(V)$, then $B(E) A(E) B(E)^{-1} = R_{\delta(E)}$, where $B = B(A(E))$ and $\theta(E) = \Theta(A(E))$ are analytic functions and $\frac{dE}{dE} \theta(E) > 0$. Let

$$A_\delta(E) = D(E) R_{\delta(E)} D(E)^{-1} A(E),$$

where $D(E) = \begin{pmatrix} E^{1/4} & 0 \\ 0 & E^{-1/4} \end{pmatrix}$. Then

$$A'(E) = A_\delta(E)^n A(E)^n,$$

where $A'(\cdot) = A[V'](\cdot)$.

For every $\kappa > 0$, for $\delta > 0$ sufficiently small, it is clear that for every $0 \leq t \leq \delta$, for every $E \in J$, we have $d(H(t) \cdot u(E), u(E)) < \kappa$, where $H(t) = D(E) R_{\delta(E)} D(E)^{-1}$ is the exponential of $t \begin{pmatrix} 0 & -E \\ 1 & 0 \end{pmatrix}$ and $u(\cdot) = u[V](\cdot)$.

Notice that for $\delta > 0$ sufficiently small, we have $|\text{tr} A_\delta(E)| < 2$ for $E \in J$. Moreover, $B_\delta(E) = B(A_\delta(E))$ and $\theta_\delta(E) = \Theta(A_\delta(E))$ converge to $B(E)$ and $\theta(E)$, when $\delta \to 0$, as analytic functions of $E \in J$. In particular,

$$\lim_\delta \sup \max_\delta \sup_n \sup_{E \in J} \|B(E) A'(E) B(E)^{-1} - R_{n(\theta(E) + \theta_\delta(E))}\| = 0.$$

For $0 < \eta < 1/2$, let $J_{\delta, n, \eta} \subset J$ be the set of all $E$ such that $2n(\theta(E) + \theta_\delta(E))$ is at distance at least $\eta$ from $\mathbb{Z}$. Since $\frac{d}{dE} \theta(E) > 0$ and $\frac{d}{dE} \theta_\delta(E) > 0$, we get, for every $\delta > 0$ small,

$$\lim_{n \to \infty} |J_{\delta, n, \eta}| = (1 - 2\eta) |J|.$$

For every $0 < \eta < 1/2$ and $\kappa > 0$, if $\delta$ is sufficiently small, then for every $n$ and for every $E \in J_{\delta, n, \eta}$, we have $E \in \Omega(V')$ and $d(u'(E), u_\delta(E)) + d(u_\delta(E), u(E)) < \kappa$, where $u' \cdot = u[V'](\cdot)$ and $u_\delta(E) = u[A_\delta(E)]$. Notice that $A(E)^{j/2} u'(E) = u'(E, a_j)$ for $0 \leq j \leq n$, and $A_\delta(E)^{2n-j} u'(E) = u'(E)$ for $n \leq j \leq 2n$. In particular,

$$d(u'(E, a_j), u(E)) < \kappa.$$

Thus for $0 \leq j \leq 2n - 1$ and $a_j \leq t \leq a_j + T$, we get

$$d(u'(E, t), u(E, t - a_j)) < \kappa.$$
For $n \leq j \leq 2n - 1$ and $a_j + T \leq t \leq a_j + T + \delta$, we have $H(a_j + T + \delta - t)u'(E, t) = u'(E, a_j + 1)$, so that
\begin{equation}
(103) \quad d(u'(E, t), u(E)) = d(u'(E, a_j + 1), H(a_j + T + \delta - t) \cdot u(E)) \\
\leq d(u'(E, a_j + 1), u(E)) + d(H(a_j + T + \delta - t) \cdot u(E)), u(E)) < 2\kappa.
\end{equation}
It follows that for each $0 \leq t' \leq T'$ we can find some $0 \leq t(t') \leq T$, defined by $t(t') = t' - a_j$ if $a_j \leq t' \leq a_j + T$ for some $j$, and $t(t') = 0$, if $a_j + T \leq t \leq a_j + 1$ for some $j$ such that for every $E \in J_{\delta,n,\eta}$, we have $d(u'(E, t'), u(E), t) < 2\kappa$.

It follows that for every $0 < \eta < 1/2$ and $\kappa > 0$, for $\delta > 0$ sufficiently small, for every $n$ sufficiently large,
\begin{equation}
(104) \quad \mu_{V_t', C}(J_{\delta,n,\eta}) \geq e^{-2\kappa} \mu_{V_t, C}(J_{\delta,n,\eta}) \\
\geq \mu_{V_t, C}(J) - (1 - e^{-2\kappa})\mu|J| - e^{-2\kappa}\mu|J \setminus J_{\delta,n,\eta}| \\
\geq \mu_{V_t, C}(J) - 2(\kappa + \eta)\mu|J|
\end{equation}
(where $\mu_{V_t', C}$ denotes the truncation of the spectral measure $\mu_{V_t'}$ for the shift of $V'$). Thus, if $\eta + \kappa$ is sufficiently small, we get
\begin{equation}
(105) \quad \mu_{V_t', C}(-\infty, M) \geq \mu_{V_t, C}(J) - \frac{\epsilon - \epsilon_0}{2}.
\end{equation}
Notice that $t(t')$ is such that for every $C_1 > 0$, $\epsilon_1 > 0$, we have, for every $\delta > 0$ sufficiently small, for every $n \in \mathbb{N}$,
\begin{equation}
(106) \quad \sup_{|s| \leq C_1} |V'(t' + s) - V(t + s)| < \epsilon_1.
\end{equation}
By Lemma 4.1 if $\delta > 0$ is sufficiently small we have
\begin{equation}
(107) \quad \mu_{V_t'}(-\infty, M) < \mu_{V_t}(-\infty, M) + \frac{\epsilon - \epsilon_0}{2}.
\end{equation}
Together with (106), it follows that
\begin{equation}
(108) \quad \mu_{V_t', C}(-\infty, M) \leq \mu_{V_t, C}(-\infty, M) + \epsilon,
\end{equation}
as desired. \hfill \Box

Remark 4.7. The construction also gives that for every $\kappa > 0$, for a subset of $\Omega(V)$ whose complement has arbitrarily small measure, we have
\begin{equation}
(109) \quad \sup_{t'} d(u[V'](E, t'), i) \leq \sup_t d(u[V](E, t), i) + 2\kappa.
\end{equation}

Proof of Theorem 1. \textbf{Define a sequence of $T^{(n)}$-periodic time-changes of solenoidal flows $F_t^{(n)} : S^{(n)} \to S^{(n)}$ and a sequence of continuous functions $v^{(n)} : S^{(n)} \to \mathbb{R}$ in the following way.}

First take $T^{(0)} = 1$, $S^{(0)} = \mathbb{R}/\mathbb{Z}$, $F_t^{(0)} = F_t^{S^{(0)}}$, and $v^{(0)} : \mathbb{R}/\mathbb{Z} \to \mathbb{R}$ a non-constant smooth function with $v^{(0)}(0) = 0$ near 0. Let $\kappa_0 = 1$. Then for $j \geq 1$,
1. Choose $C_{j-1} > 0$ so that $t \mapsto v^{(j-1)}(F_t^{(j-1)}(0))$ is $(2^{-j}, C_{j-1}, 2^{-j})$-uniform.
2. Choose $(F^{(j)}, v^{(j)})$ so that it is $(2^{-j'}, C_{j'}, 2^{j'})$-uniform for all $0 \leq j' \leq j - 1$, $F^{(j)}$ is $(2^{-j-1}, F^{(j-1)})$-mixed, and $(F^{(j)}, v^{(j)})$ is $\kappa_{j-1}$-close to a lift of $(F^{(j-1)}, v^{(j-1)})$.
3. Let $0 < \kappa_j < \kappa_{j-1}/2$ be such that if $(F, v)$ is $2\kappa_j$-close to the lift of $(F^{(j)}, v^{(j)})$ then $(F, v)$ is $(2^{-j-1}, F^{(j-1)})$-mixed.
The first step is an application of Lemma 4.2, the second is an application of Lemmas 4.6 (notice that by the previous choices, \((F^{j-1}, v^{j-1})\) is \((2-j', C'_j, 2j')\)-uniform for all \(0 \leq j' \leq j - 1\) and 4.5 and the third is an application of Lemma 4.3.

Let \(S\) be the projective limit of the \(S^{(j)}\) and let \((F, v)\) be the projective limit of the \((F^{(j)}, v^{(j)})\). Then \(F\) is \((2^j, F^{(j)})\)-mixed for all \(j \geq 1\), so it is weak mixing by Lemma 4.4. We also have that for every \(x \in S\), \(V_x^{(j)} : t \mapsto v^{(j)}(F_t^{(j)}(p_{S,S^{(j)}}(x)))\) converges to \(V_x : t \mapsto v(F_t(x))\) uniformly on compacts. It follows that the spectral measure \(\mu = \mu_{V_j}\) is the limit of the spectral measures \(\mu_{V^{(j)}}\). For every \(C > 0\), and up to taking a subsequence, the truncations \(\mu_{V^{(j)}}_{C}\) converge to a measure \(\mu_C \leq \mu\) which is absolutely continuous with density bounded by \(C\).

Then we have

\[
\mu(-\infty, 2^j) - \int_{-\infty}^{2^j} \frac{d\mu}{dE} dE \leq \lim_{k \to \infty} \mu_{V^{(k)}}(-\infty, 2^j) - \mu_{V^{(k)}_{C_j}}(-\infty, 2^j) \leq 2^{-j}.
\]

The result follows.

**Remark 4.8.** Notice that the construction allows us to obtain (by Remark 4.7), that for \(k \in \mathbb{N}\) there exists a subset \(\Gamma^{(k)} \subset \Omega(V^{(k)}) \cap \Omega(V^{(k+1)})\) such that \(|\Omega^{(k)} \setminus \Gamma^{(k)}| \leq 2^{-k}\) and

\[
\sup_t d(u[V^{(k+1)}](E, t), i) \leq \sup_t d(u[V^{(k)}](E, t), i) + 2^{-k}.
\]

Moreover, by Lemma 4.7 we may also assume that \(|\Sigma(F, v) \setminus \Omega(V^{(k)})| \leq 2^{-k}\). It follows that for almost every \(E \in \Sigma(F, v)\), there exists \(C(E) > 0\) such that \(E \in \Omega(V^{(k)})\) for every \(k\) sufficiently large and \(\sup_t d(u[V^{(k)}](E, t), i) \leq C(E)\). This implies that \(\sup_t \sup_s \|A[V^{(k)}](E, t, s)\| \leq e^{C(E)}\) and hence

\[
\sup_t \sup_s \|A[F, v](E, x, t, s)\| \leq e^{C(E)},
\]

so that every eigenfunction with such an energy must be bounded.

5. DISCRETE CASE: UNBOUNDED EIGENFUNCTIONS

5.1. SCHROEDINGER COCYCLES. Given a function \(V : \mathbb{Z} \to \mathbb{R}\), we define the transfer matrices \(A[V](E, m, n)\) so that \(A[V](E, m, m) = \text{id}\),

\[
A[V](E, m, n + 1) = \begin{pmatrix} E - V(n) & -1 \\ 1 & 0 \end{pmatrix} A[V](E, m, n).
\]

An eigenfunction of the Schrödinger operator with potential \(V\) and energy \(E\) is a solution of \(u_n = A[V](E, m, n) \cdot u_{n-1}\).

**Lemma 5.1.** If \(n > m\) and \(|\text{tr} A[V](E, m, n)| < 2\) then

\[
\frac{d}{dE} \Theta(A[V](E, m, n)) < 0.
\]

**Proof.** As in the proof of Lemma 2.4, it is enough to show that \(-A(E, m, n)^{-1}\partial_E A(E, m, n)\) is weakly positive, where we write \(A(\cdot)\) for \(A[V](\cdot)\).
Note that
\begin{equation}
A(E, m, k + 1)A(E, m, k + 1) = A(E, m, k)^{-1} \partial_E A(E, m, k) + A(E, m, k)^{-1} \begin{pmatrix}
0 & 0 \\
1 & 0
\end{pmatrix} A(E, m, k),
\end{equation}
so that
\begin{equation}
-A(E, m, n)\partial_E A(E, m, n) = \sum_{k=m}^{n-1} A(E, m, k)^{-1} \begin{pmatrix}
0 & 0 \\
1 & 0
\end{pmatrix} A(E, m, k)
\end{equation}
is a sum of weakly positive matrices in \( \text{sl}(2, \mathbb{R}) \), and hence is weakly positive as well. \( \square \)

Assume now that \( V \) is periodic of period \( N \). In this case we write \( A(V)(E, n) = A(V)(E, n + N) \) and \( A(V)(E) = A(V)(E, 0) \). Note that \( \text{tr}A(V)(E, n) = \text{tr}A(V)(E) \) for all \( n \in \mathbb{N} \). Then the spectrum \( \Sigma = \Sigma(V) \) of the Schrödinger operator with potential \( V \) is the set of all \( E \) with \( |\text{tr}A(V)(E)| \leq 2 \). Let also \( \Omega = \Omega(V) \) be the set of all \( E \) with \( |\text{tr}A(V)(E)| < 2 \). We note that \( \Sigma \setminus \Omega = \partial \Omega \) consists of finitely many points. For \( E \in \Omega(V) \), put \( u[V](E, n) = u[A(V)(E, n)] \) and \( u[V](E) = u[V](E, 0) \).

Let \( f : X \to X \) be a minimal uniquely ergodic map with invariant probability measure \( \sigma \). Given \( v : X \to \mathbb{R} \) continuous, we let \( A[f, v](E, x, m, n) = A(V)(E, m, n) \) where \( V(n) = v(f^n(x)) \). We define the Lyapunov exponent
\begin{equation}
L(E) = \lim_{n \to \infty} \frac{1}{n} \int \ln \|A[f, v](E, x, 0, n)\| d\sigma(x).
\end{equation}

We will use the following criterion for the existence of the ac spectrum (Ishii-Pastur, Kotani, Last-Simon).

**Theorem 6** (see [D]). The ac part of the spectral measure of the discrete Schrödinger operator with potential \( V(n) = v(f^n(x)) \) is equivalent to the restriction of Lebesgue measure to \( \{L(E) = 0\} \).

**Remark 5.2.** The fact that the essential support of the ac spectrum is contained in \( \{L(E) = 0\} \), for almost every \( x \), is the Ishii-Pastur Theorem. Kotani’s Theorem gives the reverse inclusion, still for almost every \( x \). Those results apply for general ergodic dynamics. Last-Simon and Kotani proved that for minimal dynamics the essential support of the ac spectrum is constant everywhere (and not only almost everywhere).

### 5.2. Construction of families of periodic potentials.

Let \( \mathcal{V} : \mathbb{R}/N_0 \mathbb{Z} \times \mathbb{Z}/N_1 \mathbb{Z} \to \mathbb{R} \) be a continuous function. We think of \( \mathcal{V} \) as a one-parameter family (parametrized by \( \mathbb{R}/N_0 \mathbb{Z} \)) of periodic potentials \( \mathcal{V}_t(\cdot) = \mathcal{V}(t, \cdot) \) (of period \( N_1 \)).

We define some basic operations on such a \( \mathcal{V} \). First, for \( n \in \mathbb{N} \), the \( n \)-repetition \( \mathcal{V}' : \mathbb{R}/N_0 \mathbb{Z} \times \mathbb{Z}/nN_1 \mathbb{Z} \to \mathbb{R} \) of \( \mathcal{V} \) is given by \( \mathcal{V}'(t, j) = \mathcal{V}(t, j) \). We obviously have
\begin{equation}
A[\mathcal{V}'_t](E, m) = A[\mathcal{V}_t](E, m)^n.
\end{equation}

Secondly, given some \( n \in \mathbb{N} \), we define the \( n \)-twist \( \mathcal{V}'' : \mathbb{R}/N_0 \mathbb{Z} \times \mathbb{Z}/nN_1 \mathbb{Z} \to \mathbb{R} \) of \( \mathcal{V} \) by \( \mathcal{V}''(t, j) = \mathcal{V}(t + N_0 \frac{k}{n}, l) \), whenever \( j = kN_1 + l \) with \( 0 \leq j \leq n - 1 \) and \( 0 \leq l \leq N_1 - 1 \). We notice that
\begin{equation}
A[\mathcal{V}'_t](E) = A[\mathcal{V}_{t + N_0 \frac{n-1}{n}}](E) \cdots A[\mathcal{V}_t](E).
\end{equation}
For the third operation, we will make use of some fixed smooth function \( \Psi : [-1, 2] \to [0, 1] \), with \( \Psi = 0 \) in a neighborhood of \(-1\) and \(2\), and \( \Psi = 1 \) in a neighborhood of \([0, 1]\). We also assume that \( N_1 \geq 3 \). Then for \( \delta > 0 \) and \( n \in \mathbb{N} \), we define the \((\delta, n)\)-slide \( \mathcal{V}^\prime : \mathbb{R}/2nN_0 \mathbb{Z} \times \mathbb{Z}/3N_1 \mathbb{Z} \to \mathbb{R} \) of \( v \) by

\[
\mathcal{V}^\prime(t, j) = \mathcal{V}(t, j), \quad 0 \leq j \leq 2N_1 - 1,
\]

\[
\mathcal{V}^\prime(t, j) = \mathcal{V}(t, j), \quad 2N_1 \leq j \leq 3N_1 - 1 \quad t \in [0, nN_0 - 1] \cup [nN_0 + 2, 2nN_0],
\]

and

\[
\mathcal{V}^\prime(t, j) = \mathcal{V}(t+\delta \Psi(t-nN_0), j), \quad 2N_1 \leq j \leq 3N_1 - 1 \quad t \in [nN_0 - 1, nN_0 + 2].
\]

Notice that we have

\[
A[\mathcal{V}_t^\prime](E) = A[\mathcal{V}_t](E)^3, \quad t \in [0, nN_0 - 1] \cup [nN_0 + 2, 2nN_0],
\]

\[
A[\mathcal{V}_t^\prime](E) = A[\mathcal{V}_{t+\delta \Psi(t-nN_0)}](E) \cdot A[\mathcal{V}_t](E)^2, \quad t \in [nN_0 - 1, nN_0 + 2].
\]

**Lemma 5.3.** Fix some closed interval \( J \subset \mathbb{R} \) and let \( u_0 : J \times [-1, 2] \to \mathbb{H} \) be a smooth function with

\[
\sup_{t \in [0, 1]} \left| \frac{d}{dt} u_0(E, t) \right| > 0
\]

for every \( E \in J \). There exists \( \epsilon_1 > 0 \), \( C' > 0 \) and \( \delta_0 > 0 \) with the following property. Let \( \mathcal{V} : \mathbb{R}/N_0 \mathbb{Z} \times \mathbb{Z}/N_1 \mathbb{Z} \to \mathbb{R} \) be a smooth function. Let \( E_0 \in \text{int} J \cap \bigcap \Omega(\mathcal{V}_t) \). Assume that \([-1, 2] \ni t \mapsto u[\mathcal{V}_t](E_0) \) is (strictly) \( \epsilon_1 \)-close in the \( C^1 \)-topology to \([-1, 2] \ni t \mapsto u_0(E_0, t) \). Then there exists \( \epsilon_0 > 0 \) such that for every \( 0 < \epsilon < \epsilon_0 \), for every \( \kappa > 0 \), for every \( 0 < \delta < \delta_0 \), for every \( N_2 \) sufficiently large, for every \( N_3 \) sufficiently large, for every \( N_4 \) sufficiently large, for every \( N_5 \) sufficiently large, if \( \mathcal{V}^\prime : \mathbb{R}/2N_4 \mathbb{Z} \times \mathbb{Z}/3N_5 \mathbb{Z} \to \mathbb{Z}/N_2 \mathbb{Z} \mathbb{N}_2 \mathbb{N}_1 \mathbb{Z} \) is the \( N_5 \)-twist of the \((\delta, N_4)\)-slide of the \( N_2 \)-repetition of the \( N_3 \)-twist of \( \mathcal{V} \), then there exists a compact set

\[
\Lambda \subset [E_0 - \epsilon, E_0 + \epsilon] \cap \bigcap \Omega(\mathcal{V}_t^\prime)
\]

such that

1. \( |\Lambda| > 2(1 - C'\delta)\epsilon \)
2. For \( E \in \Lambda \), \([-1, 2] \ni t \mapsto u[\mathcal{V}_t^\prime](E) \) is (strictly) \( \epsilon_1 \)-close in the \( C^1 \)-topology to \([-1, 2] \ni t \mapsto u_0(E, t) \).
3. For \( E \in \Lambda \),

\[
\frac{1}{6N_5 N_4 N_3 N_2 N_1 N_0} \sum_{j \in \mathbb{Z}/3N_5 N_3 N_2 N_1 \mathbb{Z}} \int_0^{2N_4 N_0} d(u[\mathcal{V}_t^\prime](E, j), i) dt - \frac{1}{N_1 N_0} \sum_{j \in \mathbb{Z}/N_1 \mathbb{Z}} \int_0^{N_0} d(u[\mathcal{V}_t](E, j), i) dt < \kappa,
\]

4. For \( E \in \Lambda \),

\[
\inf_t \sup_j d(u[\mathcal{V}_t^\prime](E, j), i) \geq \sup_j \sup_t d(u[\mathcal{V}_t](E, j), i) - \kappa,
\]
5. For any $C'\delta < \gamma < C'^{-1}$, there exists a compact set $\Lambda' \subset \Lambda$ with $|\Lambda'| > \gamma \epsilon$ such that for $E \in \Lambda'$,

\begin{equation}
\inf \sup \limits_{t,j} d(u[V_t'](E,j),i) \geq \sup \sup \limits_{t,j} d(u[V_t](E,j),i) + C'^{-1} \frac{\delta}{\gamma} - \kappa.
\end{equation}

Proof. Write $V''''$ for the $N_3$-twist of $V$, $V'''$ for the $N_2$-repetition of $V''''$, $V''$ for the $(\delta, N_4)$-slide of $V''''$.

In the first step, going from $V$ to $V''''$, we obtain, using Lemma 2.12 a set of good energies $E \in [E_0 - \epsilon, E_0 + \epsilon] \cap \bigcap_{i} \Omega_i$ of measure at least $2\epsilon(1 - \delta)$, such that $t \rightarrow u[V'''](E)$ is $C^1$ close to $t \rightarrow u[V_t](E)$, and letting $B(E,t) = B(A[V'''](E), \theta(E,t) = \Theta(A[V'''](E))$, we have $N_2(\sup, \theta(E,t) - d, \theta(E,t))$ arbitrarily small. Moreover, the random variables $\theta(E,0)$ near $E_0$ are becoming equidistributed in $\mathbb{R}/\mathbb{Z}$ as $N_3$ grows. We also have that

\begin{equation}
\inf \sup \limits_{t,j} d(u[V'''](E,j)) \geq \sup \sup \limits_{t,j} d(u[V_t](E,j),i) - \frac{\kappa}{3}.
\end{equation}

The second step does not change much, since $E \in \Omega_i$ provided $N_2 \theta(E,t)$ is not an integer, and in this case $u[V'''](E,j) = u[V'''](E,j)$. On the other hand, if $\{\theta(E,0)\}_{0 \leq j \leq N_3 - 1}$ is $1/100$ dense mod 1, then for any $w \in \mathbb{H}$, and for every $t$,

\begin{equation}
\sup \limits_{0 \leq j \leq N_3 N_2 N_1 - 1} d(A[V'''](E,j), w, i) \geq d(u[V'''](E,l), i) + \frac{1}{2} d(w, u[V'''](E)).
\end{equation}

The condition on $\theta(E,0)$ demands the exclusion of some energies, but of arbitrarily small measure (which we can take less than $2\delta \epsilon$), provided $N_2$ is large.

Going from $V''''$ to $V''$, for such good energies $E$ we obtain

\begin{equation}
\text{tr} A[V'''](E) = 2 \cos 2N_2 \pi (2\theta(E,t) + \theta(E,t + \delta \tilde{\Psi}(t))) - (\lambda(E,t) - \lambda(E,t)^{-1})^2 \sin 4N_2 \pi \theta(E,t) \sin 2N_2 \pi \theta(E,t + \delta \tilde{\Psi}(t)),
\end{equation}

with

\begin{equation}
\lambda(E,t) = e^{d(u[V'''](E), u[V'''](E,t + \delta \tilde{\Psi}(t)))}/2,
\end{equation}

and $\tilde{\Psi} : \mathbb{R}/2N_4 N_0 \mathbb{Z} \rightarrow [0,1]$ is given by $\tilde{\Psi}(t) = 0$ if $t \in [0, N_4 N_0 - 1] \cup [N_4 N_0 + 2, 2N_4 N_0]$ and $\tilde{\Psi}(t) = \Psi(t - N_4 N_0)$ if $t \in [N_4 N_0 - 1, N_4 N_0 + 2]$. Notice that $\lambda(E,t)^{-1}$ vanishes if $t \in [0, N_4 N_0 - 1] \cup [N_4 N_0 + 2, 2N_4 N_0]$, is at most of order $\delta$ everywhere, and gets to be of precisely order $\delta$ for some $t \in [N_0, N_0 + 1]$ (here we use (I.15)).

We now exclude $E$ with $\sin 6N_2 \pi \theta(E,0)$ of order $\delta$. The excluded set of energies has measure of order $2\delta \epsilon$. For the remaining energies, $|\text{tr} A[V'''](E)| < 2 - \delta^2$ for all $t$.

By (131), for every $t$, using that $u[V'''](E,j) = A[V'''](E,j) \cdot u[V'''](E)$ and also $A[V'''](E,j) = A[V'''](E) \cdot j$ for $0 \leq j \leq 2N_3 N_2 N_1$, we have

\begin{equation}
\sup \limits_{j} d(u[V'''](E,j), i) \geq \sup \sup \limits_{j} d(u[V'''](E,j), i) + \frac{1}{2} d(u[V'''](E), u[V'''](E)),
\end{equation}
and together with (130) we get

\begin{equation}
\sup_j d(u[\mathcal{V}_t^\dagger](E),i) \geq \sup_t \sup_j d(u[\mathcal{V}_t](E),i) + \frac{1}{2} d(u[\mathcal{V}_t^\dagger](E),u[\mathcal{V}_t^\dagger](E)) - \frac{\kappa}{3}.
\end{equation}

In particular, we always have

\begin{equation}
\sup_j d(u[\mathcal{V}_t^\dagger](E),i) \geq \sup_t \sup_j d(u[\mathcal{V}_t](E),i) - \frac{\kappa}{3}.
\end{equation}

We compute the distance from $u[\mathcal{V}_t^\dagger](E)$ to $u[\mathcal{V}_t^\dagger](E)$. It is equal to the distance from $w'(E,t)$ to $i$ where $w'(E,t)$ is the solution $z \in \mathbb{H}$ of the equation $az^2 + bz + c = 0$, where

\begin{align}
a &= \cos 2N_2\pi\theta(E,t + \delta \tilde{\Psi}(t)) \sin 4N_2\pi\theta(E,t) \\
&\quad + \lambda(E,t)^{-2} \sin 2N_2\pi\theta(E,t + \delta \tilde{\Psi}(t)) \cos 4N_2\pi\theta(E,t),
\end{align}

\begin{align}
b &= (\lambda(E,t)^2 - \lambda(E,t)^{-2}) \sin 2N_2\pi\theta(E,t + \delta \tilde{\Psi}(t)) \sin 4N_2\pi\theta(E,t),
\end{align}

\begin{align}
c &= \cos 2N_2\pi\theta(E,t + \delta \tilde{\Psi}(t)) \sin 4N_2\pi\theta(E,t) \\
&\quad + \lambda(E,t)^2 \sin 2N_2\pi\theta(E,t + \delta \tilde{\Psi}(t)) \cos 4N_2\pi\theta(E,t) = 0.
\end{align}

If the distance from $N_2\theta(E,0)$ to $\frac{1}{3} + \mathbb{Z}$ is exactly $\gamma$, with $C_2\delta < \gamma < C_2^{-1}$, then

\begin{equation}
C_3^{-1} \frac{\lambda(E,t)^2 - \lambda(E,t)^{-2}}{\gamma} \leq d(w'(E,t),i) \leq C_3 \frac{\lambda(E,t)^2 - \lambda(E,t)^{-2}}{\gamma}.
\end{equation}

Using that $\lambda(E,t) - 1$ does become of order $\delta$ for some $t$, we get, for such $E$,

\begin{equation}
\sup_t \sup_j d(u[\mathcal{V}_t^\dagger](E),j),i) \geq \sup_t \sup_j d(u[\mathcal{V}_t](E),j),i) + C_4^{-1} \frac{\delta}{\gamma} - \frac{\kappa}{3}.
\end{equation}

On the other hand, if one only excludes the energies with $\sin 6N_2\pi\theta(E,0)$ of order $\delta$, we still get that $d(w'(E,t),i)$ is uniformly bounded as $N_4$ grows, which implies that $\sup_t \sup_j d(u[\mathcal{V}_t^\dagger](E),j),i)$ is uniformly bounded as $N_4$ grows.

Proceeding with the last step, we get

\begin{equation}
\inf \sup_t d(u[\mathcal{V}_t^\dagger](E),j),i) \geq \sup_t \sup_j d(u[\mathcal{V}_t](E),j),i) - \frac{2\kappa}{3},
\end{equation}

while for $2C_2\delta < \gamma < C_2^{-1}$ and a set of $E$ of probability of order $\gamma$ we get

\begin{equation}
\inf \sup_t d(u[\mathcal{V}_t^\dagger](E),j),i) \geq \sup_t \sup_j d(u[\mathcal{V}_t](E),j),i) + C_4^{-1} \frac{\delta}{\gamma} - \frac{2\kappa}{3}.
\end{equation}

It remains to check that the average of $d(u[\mathcal{V}_t^\dagger](E),j),i)$ is close to the average of $d(u[\mathcal{V}_t](E),j),i)$. The first, second, and fourth steps clearly do not increase the average significantly. For the third step, we have $u[\mathcal{V}_t^\dagger](E,j) = u[\mathcal{V}_t^\dagger](E,j)$ except when $t \in [N_4N_0 - 1, N_4N_0 + 2]$. Since $d(u[\mathcal{V}_t^\dagger](E),j),i)$ remains bounded as $N_4$ grows, we conclude that the average can not be increased significantly in the third step as well.
With this result in hands, analogues of Lemmas 3.3 and 3.1 can be easily obtained. We state the conclusion:

**Lemma 5.4.** Fix some closed interval $J \subset \mathbb{R}$ and let $u_0 : J \times [-1, 2] \to \mathbb{H}$ be a smooth function with

$$
(144) \quad \sup_{t \in [0, 1]} \left| \frac{d}{dt} u_0(E, t) \right| > 0
$$

for every $E \in J$. There exists $\epsilon_1 > 0$ with the following property. Let $V^{(0)} : \mathbb{R}/N_0^{(0)} \mathbb{Z} \times \mathbb{Z}/N_1^{(0)} \mathbb{Z} \to \mathbb{R}$ be a smooth function. Let $\Gamma_0 \subset J \cap \bigcap_t \Omega(V_t^{(0)})$ be a compact set of $E$ such that $[-1, 2] \ni t \mapsto u[V_t^{(0)}](E)$ is (strictly) $\epsilon_1$-close to $[-1, 2] \ni t \mapsto u_0(E, t)$ in the $C^1$-topology.

Let $C > 0$ be such that

$$
(145) \quad \sup_{E \in \Gamma_0} \frac{1}{N_0 N_1} \sum_{j=0}^{N_1-1} \int_0^{N_0} d(u[V^{(0)}](E, j), i) dt < C.
$$

Then for every $\xi > 0$, $C_0 > 0$, for every $\delta > 0$ sufficiently small, there exist $0 < P < \xi \delta^{-1}$, and sequences $N^{(j)}_1$, $1 \leq j \leq P$, $2 \leq l \leq 5$, such that if we define $V^{(j)}$, $1 \leq j \leq P$ so that $V^{(j)}$ is obtained by $N_1^{(j)}$-twist of the $(\delta, N_1^{(j)})$-slide of the $N_2^{(j)}$-repetition of the $N_3^{(j)}$-twist of $V^{(j-1)}$, then there exists a compact subset $\Gamma \subset \Gamma_0 \cap \bigcap_t \Omega(V_t^{(P)})$ such that $|\Gamma_0 \setminus \Gamma| < \xi$, and for every $E \in \Gamma$, letting $N_0' = N_0 2^P \prod_{j=1}^P N_4^{(j)}$ and $N_1' = N_1 3^P \prod_{j=1}^P N_5^{(j)} N_3^{(j)} N_2^{(j)}$, we have

$$
(146) \quad \inf_j \sup_t d(u[V_t^{(P)}](E, j), i) \geq C_0,
$$

$$
(147) \quad \frac{1}{N_0 N_1'} \sum_{j=0}^{N_1'-1} \int_0^{N_0'} d(u[V_t^{(P)}](E, j), i) < C.
$$

Moreover, for $E \in \Gamma$, $[-1, 2] \ni t \mapsto u[V_t^{(P)}](E)$ is (strictly) $\epsilon_1$-close to $[-1, 2] \ni t \mapsto u_0(E, t)$ in the $C^1$-topology.

**Remark 5.5.** In the setting of the previous lemma, we have the following extra information on $V^{(P)}$. There exists $n \in \mathbb{N}$ such that for every $E \in \Gamma$ we have

$$
(148) \quad \inf_{w \in \mathbb{R}^2, \|w\| = 1} \sup_{0 \leq t \leq n} A[V_t^{(P)}](E, j, j + l) > e^{(C_0 - 2C)/4}/2,
$$

except for a set of $(t, j)$ of measure less than $C_0^{-1/2}$. The argument is similar to the proof of Lemma 2.3. Indeed, if $(t, j)$ is such that $d(u[V_t^{(P)}](E, j), i) \leq C_0^{1/2} C$ and $l \in \mathbb{N}$ is such that $d(u[V_t^{(P)}](E, j + l), i) \geq C_0$, then $A[V_t^{(P)}](E, j, j + l + k N_1)$ decomposes as a product $B(t, j + l)^{-1} R_{\theta + k \theta} B(t, j)$, where $B(t, m) = B(A[V_t^{(P)}](E, m))$, and $\theta = \Theta(A[V_t^{(P)}](E, j))$. Since $20 \notin \mathbb{Z}$, this implies that for any $w$ we can find $k$ such that $R_{\theta + k \theta} B(t, j) \cdot w$ has angle at most $\pi/4$ with the direction most expanded by $B(t, j + l)^{-1}$, which gives the estimate $\|B(t, j + l)^{-1} R_{\theta + k \theta} B(t, j) \cdot w\| \geq \|B(t, j + l)\|^{-1}/\sqrt{2}$. 

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5.3. Construction of almost periodic dynamics. Let $N_0, N_1 \in \mathbb{N}$, and let $a \in \mathbb{Q}$ be an integer multiple of $\frac{N_0}{N_1}$. Consider a smooth family of periodic potentials $\mathcal{V} : \mathbb{R}/N_0\mathbb{Z} \times \mathbb{Z}/N_1\mathbb{Z} \to \mathbb{R}$. It is natural to consider this periodic family as arising from the non-ergodic dynamics $(t, j) \mapsto (t, j + 1)$ on $\mathbb{R}/N_0\mathbb{Z} \times \mathbb{Z}/N_1\mathbb{Z}$, in the obvious way. But we can also think of it as arising from the dynamics $(t, j) \mapsto (t + a, j + 1)$, by considering the sampling function $v : \mathbb{R}/N_0\mathbb{Z} \times \mathbb{Z}/N_1\mathbb{Z} \to \mathbb{R}$ defined by $V(t, j) = v(t + ja, j)$.

Such a point of view is advantageous in that it allows us to consider our three operations on potentials as “small perturbations”.

Take $V'$ to be the $n$-repetition of $V$. Defining $v' : \mathbb{R}/N_0\mathbb{Z} \times \mathbb{Z}/nN_1\mathbb{Z} \to \mathbb{R}$ by $V'(t, j) = v'(t + ja, j)$, we obviously still get $v'(t, j) = v(t, j)$.

Take $V'$ to be the $n$-twist of $V$. Set $a' = a + \frac{N_0}{nN_1}$. Defining $v' : \mathbb{R}/N_0\mathbb{Z} \times \mathbb{Z}/nN_1\mathbb{Z} \to \mathbb{R}$ by $V'(t, j) = v'(t + ja', j)$, we see that $\sup_j v'(t, j) = v(t, j)$ becomes small for large $n$.

Take $V'$ to be the $(\delta, n)$-slide of $V$. Defining $v' : \mathbb{R}/2nN_0\mathbb{Z} \times \mathbb{Z}/3N_1\mathbb{Z} \to \mathbb{R}$ by $V'(t, j) = v'(t + ja, j)$, we see that $\sup_j v'(t, j) = v(t, j)$ becomes small for large $n$.

Choose $0 < \lambda_0 < 1/2$, and let $\Gamma = [-2 + 4\lambda_0, 2 - 4\lambda_0]$. Let $u_0(E, t)$ be the fixed point of $\left( E - 2\lambda_0 \cos 2\pi t \over 1 - 0 \right)$. Let $\epsilon_1 > 0$ be as in Lemma 5.4.

Let $C_1 > \int_0^1 d(u_0(E, t), i) dt$.

We now produce sequences $V^j, v^j : \mathbb{R}/N_0,j\mathbb{Z} \times \mathbb{Z}/N_1,j\mathbb{Z} \to \mathbb{R}$, $a_j \in \mathbb{Q}$, and compact sets $\Gamma_j$ as follows.

First set $N_0 = N_1 = 1$, $V^0(t, j) = 2\lambda_0 \cos 2\pi t, v^0 = V^0, a_0 = 0, \Gamma_0 = J$.

We now apply Lemma 5.4 to obtain $\Gamma_1 \subset \Gamma_0$ with $\Gamma_0 \setminus \Gamma_1$ arbitrarily small, and some $V^1 : \mathbb{R}/N_0,j\mathbb{Z} \times \mathbb{Z}/N_1,j\mathbb{Z} \to \mathbb{R}$ such that for $E \in \Gamma_1$

\begin{equation}
\inf_j \sup_i d(V_1^j(E, j), i) \geq 2C_1 + 4,
\end{equation}

\begin{equation}
\frac{1}{N_0,jN_1,j} \sum_{j=0}^{N_1,j-1} \int_0^{N_0,j} d(V_1^j(E, j), i) < C_1,
\end{equation}

and moreover, $[−1, 2] \ni t \mapsto u[V_1^j](E, j)$ is (strictly) $\epsilon_1$-close to $[−1, 2] \ni t \mapsto u_0(E, t)$ in the $C^1$-topology. Using Remark 5.5, we see that there exists $q_1 \in \mathbb{N}$ such that for every $E \in \Gamma_1$,

\begin{equation}
\inf_{w \in \mathbb{R}^2, \|w\| = 1} \sup_{0 \leq l \leq q_1} \|A[V_1^j](E, j, j + l) \cdot w\| > \frac{e}{2},
\end{equation}

except for a set of $(t, j)$ of measure less than $(2C_1 + 4)^{-1/2}$.

Moreover, we can alternatively realize $V^1(t, j) = v^1(t + a_1,j, j)$ for appropriate $a_1$, so that $\sup_j |v^1(t, j) - v^0(t, j)|$ is arbitrarily small. Notice that $|a_1 - a_0|$ can be taken arbitrarily small but non-zero.

We continue by induction, obtaining a decreasing sequence $\Gamma_k$, and $V^k, v^k, a_k$ such that for $k \geq 2$ and $E \in \Gamma_k$ we have

1. $\inf_j \sup_i d(u[V_k^j](E, j), i) \geq 2C_1 + 4k$,
2. $\frac{1}{N_0,kN_1,k} \sum_{j=0}^{N_1,k-1} \int_0^{N_0,k} d(u[V_k^j](E, j), i) < C_1$,
3. $[-1, 2] \ni t \mapsto u(V_t^k)(E, j)$ is (strictly) $\epsilon_1$-close to $[-1, 2] \ni t \mapsto u_0(E, t)$ in the $C^1$-topology,
4. $\sup_{t} \sup_{j} |v^k(t, j) - v^{k-1}(t, j)| < 2^{-k}$,
5. There exists $q_k \in \mathbb{N}$ such that for every $E \in \Gamma_k$,

$$\inf_{w \in \mathbb{R}^2, \|w\|=1} \sup_{0 \leq l \leq q_k} \|A[V_t^k](E, j, j + l) \cdot w\| > \frac{e^k}{2},$$

except for a set of $(t, j)$ of measure (strictly) less than $(2C_1 + 4k)^{-1/2}$.
6. $|a_k - a_{k-1}|$ is non-zero but smaller than $\frac{1}{2^{k-1}N_{1, k-1}}$.

We now construct the sampling function and the dynamics.

Let $K$ be the projective limit of $\mathbb{Z}/N_{1, j}\mathbb{Z}$ (a Cantor group), and let $S$ be the projective limit of $\mathbb{R}/N_{0, j}\mathbb{Z}$ (a solenoid). Then $v(t, j) = \lim v^k(t, j)$ defines a continuous function on $S \times K$ (for simplicity, we omit the projections $S \times K \to \mathbb{R}/N_{0, k}\mathbb{Z} \times \mathbb{Z}/N_{1, k}\mathbb{Z}$ from the notation). This is the sampling function.

Notice that $a = \lim a_k$ is irrational, since $a_k$ are rational with denominators at most $N_{1, k}$ and $0 < |a - a_k| \leq \frac{1}{2^{k-1}N_{1, k}}$. Thus $f(t, j) = (t + a, j + 1)$ is a uniquely ergodic translation in the compact Abelian group $S \times K$. This is the base dynamics.

Let $\Gamma = \bigcap \Gamma_k$. By construction, $\Gamma$ is a compact set of positive Lebesgue measure.

Notice that

$$\sup_{x \in \Gamma} \sup_{0 \leq l \leq n} \|A[f, v](E, x, 0, l) - A[V_t^k](E, j, j + l)\|$$

(with $(t, j)$ the projection of $x$) can be made arbitrarily small, for any $n$ chosen after $v^k$ is constructed, but before $v^{k+1}$ is constructed. Choosing parameters growing sufficiently fast we get

$$\lim_{n \to \infty} \sup_{x} \frac{1}{n} \ln \|A[f, v](E, x, 0, n)\| = 0,$$

i.e., the Lyapunov exponent vanishes for $E \in \Gamma$, so that $\Gamma$ is contained in the essential support of the ac spectrum for every $x$, and moreover, for every $k \geq 1$ and $E \in \Gamma_k$,

$$\inf_{w \in \mathbb{R}^2, \|w\|=1} \sup_{0 \leq l \leq q_k} \|A[f, v](E, x, 0, l) \cdot w\| \geq \frac{e^k}{4},$$

except for a set of $x$ of measure less than $(2C_1 + 4k)^{-1/2}$. Thus, for every $E \in \Gamma$, for almost every $x$ we have

$$\inf_{w \in \mathbb{R}^2, \|w\|=1} \lim_{l \to \infty} \sup_{t} \|A[f, v](E, x, 0, l) \cdot w\| = \infty,$$

which is the desired eigenfunction growth.

6. Discrete case: breaking almost periodicity

6.1. Slow deformation. The following are variations of Lemmas 2.11 and 2.12 and we leave the proofs for the reader.

**Lemma 6.1.** Let $J \subset \mathbb{R}$ be a closed interval, let $N \in \mathbb{N}$, and let $A : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R})$ be a smooth function such that $|\text{tr}A^{(N)}(E, t)| < 2$ for $(E, t) \in J \times \mathbb{R}/\mathbb{Z}$, where

$$A^{(N)}(E, t) = A \left( E, t + \frac{N - 1}{N} \right) \cdots A(E, t).$$
Let \( B(E, t) = B(A^{(N)}(E, t)) \) and let \( \theta(E, t) \) be a smooth function satisfying

\[
B \left( E, t + \frac{1}{N} \right) A(E, t) B(E, t)^{-1} = R_{\theta(E, t)}.
\]

Then for every \( m, k \in \mathbb{N} \), there exists \( n(m) \in \mathbb{N} \) and \( C_{k,m} > 0 \) such that for every \( n \geq n(m) \), there exist smooth functions \( B_{(m,n)} : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R}) \), \( \theta_{(m,n)} : J \times \mathbb{R}/\mathbb{Z} \to \mathbb{R} \) such that

\[
1. \| A_{(m,n)} - R_{\theta_{(m,n)}} \|_{C^k} \leq \frac{C_{k,m}}{n^m}, \quad \text{where}
\]

\[
2. \| B_{(m,n)} - B \|_{C^k} \leq \frac{C_{k,m}}{n^m},
\]

\[
3. \| \theta_{(m,n)} - \theta \|_{C^k} \leq \frac{C_{k,m}}{n}.
\]

**Lemma 6.2.** Under the hypothesis of the previous lemma, assume moreover that \( \tilde{\theta}(E) = \int_{\mathbb{R}/\mathbb{Z}} \theta(E, t) dt \) satisfies \( \frac{d}{dE} \tilde{\theta}(E) \neq 0 \) for every \( E \in J \). For \( n \in \mathbb{N} \), let \( A^{(N+n)} : J \times \mathbb{R}/\mathbb{Z} \to \text{SL}(2, \mathbb{R}) \) be given by

\[
A^{(N+n)}(E, t) = A \left( E, t + (nN - 1) \frac{n+1}{nN} \right) A \left( E, t + (nN - 2) \frac{n+1}{nN} \right) \cdots A(E, t).
\]

Then there exist functions \( \hat{\theta}^{(n)} : J \to \mathbb{R}/\mathbb{Z} \) such that for every measurable subset \( Z \subset \mathbb{R}/\mathbb{Z} \),

\[
\lim_{n \to \infty} |\{ E \in J, \hat{\theta}^{(n)}(E) \in Z \}| = |Z||J|,
\]

\[
\lim_{n \to \infty} |\{ E \in J, \hat{\theta}^{(n)}(E) + \hat{\theta}^{(2n)}(E) \in Z \}| = |Z||J|,
\]

with the following property. For every \( \delta > 0 \),

\[
\lim_{n \to \infty} \left\| \text{tr} A^{(N+n)}(E, t) - 2 \cos 2\pi \hat{\theta}_{(m,n)}(E) \right\|_{C^0(J \times \mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0,
\]

\[
\lim_{n \to \infty} \sup_{|\sin 2\pi \hat{\theta}^{(n)}(E)| > \delta} \| \Theta(A^{(N+n)}(E, \cdot)) - \hat{\theta}^{(n)}(E) \|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{R})} = 0,
\]

\[
\lim_{n \to \infty} \sup_{|\sin 2\pi \hat{\theta}^{(n)}(E)| > \delta} \| \text{u}(A^{(N+n)}(E, \cdot)) - \text{u}(A^{(N)}(E, \cdot)) \|_{C^1(\mathbb{R}/\mathbb{Z}, \mathbb{C})} = 0.
\]

### 6.2. The construction

In this section, we will interpret a continuous function \( \mathcal{V} : \mathbb{R}/N\mathbb{Z} \to \mathbb{R} \) as a one-parameter family of \( N \)-periodic potentials \( \mathcal{V}_t(j) = \mathcal{V}(t+j) \).

For \( n \in \mathbb{N} \), we define the \( n \)-crumbling \( \mathcal{V}' : \mathbb{R}/3nN\mathbb{Z} \to \mathbb{R} \) of \( \nu \) by

1. \( \mathcal{V}'(t) = \mathcal{V}(\frac{n+1}{n} t), \quad t \in [0, nN] \),
2. \( \mathcal{V}'(t) = \mathcal{V}(\frac{2n+1}{2n} (t-nN)), \quad t \in [nN, 3nN] \).

**Lemma 6.3.** Let \( \mathcal{V} : \mathbb{R}/N\mathbb{Z} \) be a smooth function which is constant near 0. Then for every \( \delta > 0 \), for every \( n \) sufficiently large, letting \( \mathcal{V}' \) be the \( n \)-crumbling of \( \mathcal{V} \), we have

\[
|\lim_{t \to \infty} \Omega(\mathcal{V}_t) \setminus \lim_{t \to \infty} \Omega(\mathcal{V}'_t)| < \delta.
\]

**Proof.** Fix a compact interval \( J \subset \bigcap_t \Omega(\mathcal{V}_t) \). Apply Lemma 6.2 to \( A(E, t) = \begin{pmatrix} E - e^{(Nt)} & -1 \\ 1 & 0 \end{pmatrix} \). It yields a sequence \( \hat{\theta}^{(n)}(E) \).
If $\mathcal{V}'$ is the $n$-crumbling of $v$, then for $t \in [0, 1]$, we get $A[\mathcal{V}_t'][E, 0, nN] = A^{(N+n)}(E, \frac{n+1}{nN} t)$ and $A[\mathcal{V}_t'](E, nN, 3nN) = A^{(N+2n)}(E, \frac{2n+1}{2nN} t)$. Thus for $t \in [0, 1]$ we have

(166) \[ A[\mathcal{V}_t'](E) = A^{(N+2n)}(E, \frac{2n+1}{2nN} t) A^{(N+n)}(E, \frac{n+1}{nN} t). \]

As long as $|\sin 2\pi \tilde{\theta}^{(n)}|$ and $|\sin 2\pi \tilde{\theta}^{(2n)}|$ are not too small, we can write, for $t \in [0, 1]$,

(167) \[ A[\mathcal{V}_t''](E) = B^{(2n)}(E, t)^{-1} R_{\theta^{(2n)}(E, t)} B^{(2n)}(E, t) B^{(n)}(E, t)^{-1} R_{\theta^{(n)}(E, t)} B^{(n)}(E, t), \]

where $B^{(m)} = B(A^{(N+m)}(E, t))$ and $\theta^{(m)}(E, t) = \Theta(A^{(N+m)}(E, t))$. Notice that $B^{(n)}$ and $B^{(2n)}$ are both $C^1$-close to $B(A[\mathcal{V}_t'](E))$ as functions of $t \in [0, 1]$. Moreover, $\theta^{(n)}(E, t)$ is close to $\tilde{\theta}^{(n)}(E)$ and $\theta^{(2n)}(E, t)$ is close to $\tilde{\theta}^{(2n)}(E)$.

It follows that for $t \in [0, 1]$, $\text{tr} A[\mathcal{V}_t''](E)$ is close to $2 \cos 2\pi(\tilde{\theta}^{(2n)}(E) + \tilde{\theta}^{(n)}(E))$. Thus, as long as $|\sin 2\pi(\tilde{\theta}^{(n)} + \tilde{\theta}^{(2n)})|$ is not small, we have $|\text{tr} A[\mathcal{V}_t''](E)| < 2$ for every $t \in [0, 1]$. Since $\text{tr} A[\mathcal{V}_t''](E)$ is 1-periodic, this implies that $|\text{tr} A[\mathcal{V}_t''](E)| < 2$ for all $t$.

**Remark 6.4.** One also easily gets from this construction,

(168) \[ \sup \limits_t d(u[\mathcal{V}_t'', E, i]) < \sup \limits_t d(u[\mathcal{V}_t', E, i]) + \delta \]

except for a set of $E \in \cap, \Omega(\mathcal{V}_t) \cap \cap, \Omega(\mathcal{V}_t')$ of arbitrarily small measure.

**Proof of Theorem 1.** Starting with a smooth non-constant function $\mathcal{V}^{(0)} : \mathbb{R}/\mathbb{Z} \to \mathbb{R}$, apply Lemma [3] successively to obtain a sequence $\mathcal{V}^{(k)} : \mathbb{R}/N^{(k)} \to \mathbb{R}$ such that $\mathcal{V}^{(k)}$ is the $n_k$-crumbling of $\mathcal{V}^{(k-1)}$, and compact sets $\Gamma^{(k)} \subset \cap, \Omega(\mathcal{V}^{(k)})$ with $\Gamma^{(k)} \subset \Gamma^{(k-1)}$ and $\lim_{k \to \infty} |\Gamma^{(k)}| > 0$. By taking parameters $n_k$ growing sufficiently fast, we ensure that for $E \in \Gamma^{(k+1)}$ we have

(169) \[ \sup \limits_{1 \leq j \leq N^{(k)}} \sup \limits_t \| A[\mathcal{V}^{(k+1)}_t](E, 0, j) \| - \sup \limits_t \| A[\mathcal{V}^{(k)}_t](E, 0, j) \| \leq \frac{1}{2^k}, \]

(170) \[ \sup \limits_t \frac{1}{N^{(k+1)}} \ln \| A[\mathcal{V}^{(k+1)}_t](E, 0, N^{(k+1)}) \| \leq \frac{1}{2^k}. \]

We now turn to the dynamical realization. Let $\tilde{N}^{(k)}$ be defined by $\tilde{N}^{(0)} = 1$, $\tilde{N}^{(k)} = (3n_k + 2)\tilde{N}^{(k-1)}$. We first construct $N^{(k)}$-periodic time-changes $F^{(k)}$ of the solenoidal flow on $S^{(k)} = \mathbb{R}/\tilde{N}^{(k)} \mathbb{Z}$ such that $\mathcal{V}^{(0)}(p_{S^{(k)}}, S^{(0)}(F^{(k)}(0))) = \mathcal{V}^{(k)}(t)$. We first take $F^{(0)}$ to be just the solenoidal flow on $S^{(0)}$. Now define inductively

(171) \[ w_{F^{(k+1)}}(t) = w_{F^{(k)}}(t) e^{|\rho^{(k+1)}(t)|} \]

for a suitable function $\rho^{(k+1)}$. Here it is enough to take $\rho^{(k+1)} = \ln \frac{n_{k+1} + 1}{n_k}$ on $[0, (n_{k+1} + 1)\tilde{N}^{(k)}]$, $\rho^{(k+1)} = \ln \frac{2n_{k+1} + 1}{2n_{k+1}}$ on $[(n_{k+1} + 1)\tilde{N}^{(k)} + \epsilon, (3n_{k+1} + 2)\tilde{N}^{(k)} - \epsilon]$. 

for suitably small $\epsilon$, and such that
\begin{align}
(172) \quad \int_{(n_{k+1}+1)\hat{N}(k)}^{(n_{k+1}+1)\hat{N}(k)+\epsilon} \frac{1}{w_{F(k)}(t)e^{\mu_{k+1}(t)}} dt &= \frac{2n_{k+1}}{2n_{k+1}+1} \int_0^\epsilon \frac{1}{w_{F(k)}(t)} dt, \\
(173) \quad \int_{(3n_{k+1}+2)\hat{N}(k)}^{(3n_{k+1}+2)\hat{N}(k)} \frac{1}{w_{F(k)}(t)e^{\mu_{k+1}(t)}} dt &= \frac{2n_{k+1}}{2n_{k+1}+1} \int_{S_k}^{\hat{N}(k)} \frac{1}{w_{F(k)}(t)} dt.
\end{align}
Notice that by taking parameters growing sufficiently fast, we can take $F^{(k+1)}$ close to the lift of $F^{(k)}$.

Let $S$ be the projective limit of $\mathbb{R}/\hat{N}(k)\mathbb{Z}$, and let $v : S \to \mathbb{R}$ be given by $v(x) = \mathcal{V}^{(0)}(p_{S,S(0)}(x))$. Let $F_t : S \to S$ be the projective limit of the $F^{(k)}_t$. The base dynamics will be the time-one map $F_1$ and the sampling function will be $v$.

By (169), for every $k$, if $E \in \Gamma = \bigcap \Gamma^{(k)}$,
\begin{equation}
(174) \quad \sup_{1 \leq j \leq N(k)} \left| \sup_x \ln \|A[F_1,v](E,x,0,j)\| - \sup_x \ln \|A[\mathcal{V}^{(k)}](E,0,j)\| \right| \leq \frac{1}{2^k-1},
\end{equation}
and together with (170) we get, for $E \in \Gamma$
\begin{equation}
(175) \quad \sup_x \frac{1}{N(k+1)} \ln \|A[F_1,v](E,x,0,N^{(k+1)})\| \leq \frac{1}{2^k} + \frac{1}{N(k+1)2^k} \leq \frac{1}{2^k-1},
\end{equation}
so that the Lyapunov exponent (with respect to any $F_1$-invariant measure) must vanish over $\Gamma$.

To conclude, let us show that the flow $F$ is weak mixing: This implies that the discrete dynamics $F_1$ is weak mixing as well, and since $F$ is minimal and uniquely ergodic, it also implies that $F_1$ is minimal and uniquely ergodic, so that $\Gamma$ is contained in the essential support of the absolutely continuous spectrum for every $x$. In order to do this, we notice that for $0 \leq j \leq n_{k+1} - 1$
\begin{equation}
(176) \quad p_{S^{(k+1)},S^{(k)}}(F^{(k+1)}_{jN(k)}(t)) = F^{(k)}_{j/n_{k+1}}(p_{S^{(k+1)},S^{(k)}}(t)),
\end{equation}
as long as $t \in [0,(n_{k+1}+1)(1-\frac{j}{n_{k+1}})\hat{N}(k)]$. On the other hand, for $0 \leq j \leq 2n_{k+1} - 1$
\begin{equation}
(177) \quad p_{S^{(k+1)},S^{(k)}}(F^{(k+1)}_{j\hat{N}(k)}(t)) = F^{(k)}_{2j/n_{k+1}}(p_{S^{(k+1)},S^{(k)}}(x)),
\end{equation}
as long as $t \in [(n_{k+1} + 1)\hat{N}(k) + 1, \hat{N}(k) + \frac{2n_{k+1} + 1}{2n_{k+1}} \hat{N}(k) - 1]$.

The conclusion proceeds along the same line as in Lemma 4,4. Take a measurable eigenfunction $\psi$ taking values on the unit circle, associated to an eigenvalue $\theta \neq 0$, so that $\psi \circ F_t = e^{2\pi i \theta t} \psi$. Taking conditional expectations, we obtain $\psi^{(j)}$ on $S^{(j)}$, taking values on the closed unit disk, with $\lim \psi^{(j)}(p_{S,S(j)}(x)) = \psi(x)$ for almost every $x$. We then conclude from (176) and (177) that $\frac{\theta_j}{2n_{k+1}}$ is close to an integer for $1 \leq j \leq [n_{k+1}/2]$. This contradicts $\theta \neq 0$.

Remark 6.5. Using Remark 6.4 we can ensure in the construction that
\begin{equation}
(178) \quad C = \sup_k \sup_{E \in \Gamma_k} \sup_t \sup \left| \int d(\mathcal{V}^{(k)}(E),i) < \infty.
\end{equation}
This implies that
\begin{equation}
(179) \quad \sup_k \sup_{E \in \Gamma_k} \sup_t \sup_j \|A[\mathcal{V}^{(k)}(E,0,j)]\| \leq e^C,
\end{equation}
and by (174),

$$\sup_{E \in \Gamma} \sup_{x} \sup_{j} \|A[F_1, v](E, x, 0, j)\| \leq e^C,$$

so all eigenfunctions with energies in $\Gamma$ are bounded.

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