# DRY TEN MARTINI PROBLEM FOR THE NON-SELF-DUAL EXTENDED HARPER'S MODEL 

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#### Abstract

In this paper we prove the dry version of the Ten Martini problem: Cantor spectrum with all gaps open, for the extended Harper's model in the non-self-dual region for Diophantine frequencies.


## 1. Introduction

The study of independent electrons on a two-dimensional lattice exposed to a perpendicular magnetic field and periodic potentials can be reduced via an appropriate choice of gauge field to the study of discrete one-dimensional quasiperiodic Jacobi matrices. The most extensively studied case is the almost Mathieu operator (AMO) acting on $l^{2}(\mathbb{Z})$ defined by

$$
\left(H_{\lambda, \alpha, \theta} u\right)_{n}=u_{n+1}+u_{n-1}+2 \lambda \cos 2 \pi(\theta+n \alpha) u_{n} .
$$

This is a one-dimensional tight-binding model with anisotropic nearest neighbor couplings in general. A more general model, called the extended Harper's model (EHM), is the operator acting on $l^{2}(\mathbb{Z})$ defined by:

$$
\left(H_{\lambda, \alpha, \theta} u\right)_{n}=c(\theta+n \alpha) u_{n+1}+\tilde{c}(\theta+(n-1) \alpha) u_{n-1}+2 \cos 2 \pi(\theta+n \alpha) u_{n} .
$$

where $c(\theta)=\lambda_{1} e^{-2 \pi i\left(\theta+\frac{\alpha}{2}\right)}+\lambda_{2}+\lambda_{3} e^{2 \pi i\left(\theta+\frac{\alpha}{2}\right)}$ and $\tilde{c}(\theta)=\lambda_{1} e^{2 \pi i\left(\theta+\frac{\alpha}{2}\right)}+\lambda_{2}+$ $\lambda_{3} e^{-2 \pi i\left(\theta+\frac{\alpha}{2}\right)}$. It is obtained when both the nearest neighbor coupling (expressed through $\lambda_{2}$ ) and the next-nearest couplings (expressed through $\lambda_{1}$ and $\lambda_{3}$ ) are included. This model includes AMO as a special case (when $\lambda_{1}=\lambda_{3}=0$ ).

For the AMO, it was proved in [5] that the spectrum is a Cantor set for any $\alpha \in \mathbb{R} \backslash \mathbb{Q}$ and $\lambda \neq 0$. This is the Ten Martini problem dubbed by Barry Simon, after an offer of Mark Kac. A much more difficult problem, known as the dry version of the Ten Martini problem, is to prove that the spectrum is not only a Cantor set, but that all gaps predicted by the Gap-Labelling theorem [10, [15] are open. The first result was obtained for Liouvillean $\alpha$ [12], and later it was proved for a set $(\lambda, \alpha)$ of positive Lebesgue measure [16]. The most recent result is [6], in which they were able to deal with all Diophantine frequencies and $\lambda \neq 1$. A solution for all irrational frequencies and $\lambda \neq 1$ was also recently announced in [9].

Recently, there have been several important advances on the spectral theory of the EHM: purely point spectrum for Diophantine $\alpha$ and a.e. $\theta$ in the positive Lyapunov exponent region [13]; the exact formula for Lyapunov exponent for all

[^0]coupling constants [14]; the spectral decomposition for a.e. $\alpha$ [7]. However the results that study the spectrum as a set have not been obtained for the EHM.

For EHM, depending on the values of the parameters $\lambda_{1}, \lambda_{2}, \lambda_{3}$, we could divide the parameter space into three regions as shown in the picture below:


> region $I: 0<\max \left(\lambda_{1}+\lambda_{3}, \lambda_{2}\right)<1$, region II $: 0<\max \left(\lambda_{1}+\lambda_{3}, 1\right)<\lambda_{2}$ region III $: 0<\max \left(1, \lambda_{2}\right)<\lambda_{1}+\lambda_{3}$.

According to the action of the duality transformation $\sigma: \lambda=\left(\lambda_{1}, \lambda_{2}, \lambda_{3}\right) \rightarrow \hat{\lambda}=$ $\left(\frac{\lambda_{3}}{\lambda_{2}}, \frac{1}{\lambda_{2}}, \frac{\lambda_{1}}{\lambda_{2}}\right)$, region I and region II are dual to each other and region III is a selfdual region. Region I is the positive Lyapunov exponent region, which is a natural extension of the segment $\left\{\lambda_{1}+\lambda_{3}=0,0<\lambda_{2}<1\right\}$ corresponding to the case $\lambda>1$ in the AMO. Region II is the subcritical region, which is an extension of the segment $\left\{\lambda_{1}+\lambda_{3}=0,1<\lambda_{2}\right\}$ corresponding to the case $\lambda<1$ in the AMO.

In this paper we prove the dry version of the Ten Martini problem in region I and region II under the Diophantine condition.

Let $p_{n} / q_{n}$ be the continued fraction appoximants of $\alpha \in \mathbb{R} \backslash \mathbb{Q}$. Let

$$
\beta(\alpha)=\limsup _{n \rightarrow \infty} \frac{\ln q_{n+1}}{q_{n}} .
$$

If $\beta(\alpha)=0$, we say $\alpha$ satisfies the Diophantine condition, denoted by $\alpha \in$ DC. It is easily seen that such $\alpha$ form a full measure subset of $\mathbb{T}$.

It is known that when $E$ is in the closure of a spectral gap, the integrated density of states (IDS) $N(E) \in \alpha \mathbb{Z}+\mathbb{Z}$ (refer to (2.5) for the definition of IDS) [10], [15]. Here we prove the inverse is true.

Theorem 1.1. If $\alpha \in \mathrm{DC}$ and $\lambda$ belongs to region I or region II, all possible spectral gaps are open.

Remark 1.1. We note the Dry Ten Martini problem has not yet been solved for the self-dual AMO. In the self-dual region III, Cantor spectrum is known in the isotropic case (when $\lambda_{1}=\lambda_{3}$ ); see Fact 2.1 in [7]. In fact one could prove the operator has zero Lebesgue measure spectrum for all frequencies.

Remark 1.2. In regions I and II, for Liouvillean $\alpha$ (where $\beta(\alpha)$ is large), it is not clear whether even the Cantor spectrum holds. The proof may require a non-trivial adjustment of the proof for AMO in [12].

We first establish almost localization (see section 3.1) in region I. Then a quantitative version of Aubry duality to obtain almost reducibility (see section 3.2) in region II which enables us to deal with all energies whose rotation numbers are $\alpha$-rational.

Thus the strategy follows that of [6], but we need to extend the almost localization and quantitative duality, as well as the final argument to our Jacobi setting, which is non-trivial on a technical level. At the same time unlike [6], we only deal with a short-range dual operator, leading to a significant streamlining of some arguments of [6].

We organize the paper as follows: in section 2 we present some preliminaries, in section 3 we state our main results about almost localization and almost reducibility, relying on which we provide a proof of Theorem 1.1 In sections 4 and 5 we prove the main results that we present in section 3.

## 2. Preliminaries

2.1. Cocycles. Let $\alpha \in \mathbb{R} \backslash \mathbb{Q}$ and $A \in C^{0}\left(\mathbb{T}, M_{2}(\mathbb{C})\right)$ measurable with $\log \|A(x)\| \in$ $L^{1}(\mathbb{T})$. The quasiperiodic cocycle $(\alpha, A)$ is the dynamical system on $\mathbb{T} \times \mathbb{C}^{2}$ defined by $(\alpha, A)(x, v)=(x+\alpha, A(x) v)$. The Lyapunov exponent is defined by

$$
L(\alpha, A)=\lim _{n \rightarrow \infty} \frac{1}{n} \int_{\mathbb{T}} \log \left\|A_{n}(x)\right\| \mathrm{d} x=\inf _{n} \frac{1}{n} \int_{\mathbb{T}} \log \left\|A_{n}(x)\right\| \mathrm{d} x .
$$

where

$$
\left\{\begin{array}{l}
A_{n}(x)=A(x+(n-1) \alpha) \cdots A(x) \text { for } n \geq 0 \\
A_{n}(x)=A^{-1}(x+n \alpha) \cdots A^{-1}(x-\alpha) \text { for } n<0
\end{array}\right.
$$

Lemma 2.1 (e.g. [6]). Let ( $\alpha, A$ ) be a continuous cocycle; then for any $\delta>0$ there exists $C_{\delta}>0$ such that for any $n \in \mathbb{N}$ and $\theta \in \mathbb{T}$ we have

$$
\left\|A_{n}(\theta)\right\| \leq C_{\delta} e^{(L(\alpha, A)+\delta) n}
$$

We say that $(\alpha, A)$ is uniformly hyperbolic if there exists continuous splitting $\mathbb{C}^{2}=E^{s}(x) \oplus E^{u}(x), x \in \mathbb{T}$ such that for some constant $C, \eta>0$ and all $n \geq 0$, $\left\|A_{n}(x) v\right\| \leq C e^{-\eta n}\|v\|$ for $v \in E^{s}(x)$ and $\left\|A_{-n}(x) v\right\| \leq C e^{-\eta n}\|v\|$ for $v \in E^{u}(x)$.

Given two complex cocycles $\left(\alpha, A^{(1)}\right)$ and $\left(\alpha, A^{(2)}\right)$, we say they are complex conjugate to each other if there is $M \in C^{0}(\mathbb{T}, S L(2, \mathbb{C}))$ such that

$$
M^{-1}(x+\alpha) A^{(1)}(x) M(x)=A^{(2)}(x)
$$

We assume now that $A$ is a real cocycle, $A \in C^{0}(\mathbb{T}, S L(2, \mathbb{R}))$. The notation of real conjugacy (between real cocycles) is the same as before, except that we look for $M \in C^{0}(\mathbb{T}, \operatorname{PSL}(2, \mathbb{R}))$. A reason why we look for $M \in C^{0}(\mathbb{T}, P S L(2, \mathbb{R}))$ instead of $M \in C^{0}(\mathbb{T}, S L(2, \mathbb{R}))$ is given by the following well-known result.

Theorem 2.2. Let $(\alpha, A)$ be uniformly hyperbolic, assume $\alpha \in \mathrm{DC}$ and $A$ analytic. Then there exists $M \in C^{\omega}\left(\mathbb{T}, \operatorname{PSL}(2, \mathbb{R}) \sqrt{1}\right.$ such that $M^{-1}(x+\alpha) A(x) M(x)$ is constant.

We say $(\alpha, A)$ is (analytically) reducible if it is real conjugate to a constant cocycle by an analytic conjugacy.

Let

$$
R_{\theta}=\left(\begin{array}{cc}
\cos 2 \pi \theta & -\sin 2 \pi \theta \\
\sin 2 \pi \theta & \cos 2 \pi \theta
\end{array}\right)
$$

Any $A \in C^{0}(\mathbb{T}, P S L(2, \mathbb{R}))$ is homotopic to $x \rightarrow R_{\frac{k}{2} x}$ for some $k \in \mathbb{Z}$ called the degree of $A$, denoted by $\operatorname{deg} A=k$.

Assume now that $A \in C^{0}(\mathbb{T}, S L(2, \mathbb{R}))$ is homotopic to identity. Then there exists $\phi: \mathbb{R} / \mathbb{Z} \times \mathbb{R} / \mathbb{Z} \rightarrow \mathbb{R}$ and $v: \mathbb{R} / \mathbb{Z} \times \mathbb{R} / \mathbb{Z} \rightarrow \mathbb{R}^{+}$such that

$$
A(x)\binom{\cos 2 \pi y}{\sin 2 \pi y}=v(x, y)\binom{\cos 2 \pi(y+\phi(x, y))}{\sin 2 \pi(y+\phi(x, y))}
$$

The function $\phi$ is called a lift of $A$. Let $\mu$ be any probability on $\mathbb{R} / \mathbb{Z} \times \mathbb{R} / \mathbb{Z}$ which is invariant under the continuous map $T:(x, y) \mapsto(x+\alpha, y+\phi(x, y))$, projecting over Lebesgue measure on the first coordinate. Then the number

$$
\rho(\alpha, A)=\int \phi d \mu \bmod \mathbb{Z}
$$

is independent of the choices of $\phi$ and $\mu$, and is called the fibered rotation number of $(\alpha, A)$.

It can be proved directly by the definition that

$$
\begin{equation*}
|\rho(\alpha, A)-\theta|<C\left\|A-R_{\theta}\right\|_{0} . \tag{2.1}
\end{equation*}
$$

If $\left(\alpha, A^{(1)}\right)$ and $\left(\alpha, A^{(2)}\right)$ are real conjugate, $M^{-1}(x+\alpha) A^{(2)}(x) M(x)=A^{(1)}(x)$, and $M: \mathbb{R} / \mathbb{Z} \rightarrow \operatorname{PSL}(2, \mathbb{R})$ has degree $k$, then

$$
\begin{equation*}
\rho\left(\alpha, A^{(1)}\right)=\rho\left(\alpha, A^{(2)}\right)-k \alpha / 2 \tag{2.2}
\end{equation*}
$$

For uniformly hyperbolic cocycles there is the following well-known result.
Theorem 2.3. Let $(\alpha, A)$ be a uniformly hyperbolic cocycle, with $\alpha \in \mathbb{R} \backslash \mathbb{Q}$. Then $2 \rho(\alpha, A) \in \alpha \mathbb{Z}+\mathbb{Z}$.
2.2. Extended Harper's model. We consider the extended Harper's model $\left\{H_{\lambda, \theta}\right\}_{\theta \in \mathbb{T}}$. The formal solution to $H_{\lambda, \theta} u=E u$ can be reconstructed via the following equation:

$$
\binom{u_{n+1}}{u_{n}}=A_{\lambda, E}(\theta+n \alpha)\binom{u_{n}}{u_{n-1}}
$$

where $A_{\lambda, E}(\theta)=\frac{1}{c(\theta)}\left(\begin{array}{cc}E-2 \cos 2 \pi \theta & -\tilde{c}(\theta-\alpha) \\ c(\theta) & 0\end{array}\right)$. Notice that since $A_{\lambda, E}(\theta) \notin$ $S L(2, \mathbb{R})$, we introduce the following matrix (see Lemma A.2):

$$
\begin{aligned}
\tilde{A}_{\lambda, E}(\theta) & =\frac{1}{\sqrt{|c|(\theta)|c|(\theta-\alpha)}}\left(\begin{array}{cc}
E-2 \cos 2 \pi \theta & -|c|(\theta-\alpha) \\
|c|(\theta) & 0
\end{array}\right) \\
& =Q_{\lambda}(\theta+\alpha) A_{\lambda, E}(\theta) Q_{\lambda}^{-1}(\theta)
\end{aligned}
$$

[^1]where $|c|(\theta)=\sqrt{c(\theta) \tilde{c}(\theta)}$ (which is not the same as $|c(\theta)|=\sqrt{c(\theta) \overline{c(\theta)}}$ when $\theta \notin \mathbb{T}$ ) and $Q_{\lambda}(\theta)$ is analytic on $|\operatorname{Im} \theta| \leq \frac{\epsilon_{1}}{2 \pi}$.

The spectrum of $H_{\lambda, \theta}$ denoted by $\Sigma_{\lambda}$, does not depend on $\theta$ [8] and it is the set of $E$ such that ( $\alpha, \tilde{A}_{\lambda, E}$ ) is not uniformly hyperbolic.

The Lyapunov exponent is defined by $L_{\lambda}(E)=L\left(\alpha, A_{\lambda, E}\right)=L\left(\alpha, \tilde{A}_{\lambda, E}\right)$.
For a matrix-valued function $M(\theta)$, let $M_{\epsilon}(\theta)=M(\theta+i \epsilon)$ be the phasecomplexified matrix.

In [4], Avila divides all the energies in the spectrum into three catagories: supercritical, namely the energy with positive Lyapunov exponent; subcritical, namely the energy whose Lyapunov exponent of the phase-complexified cocycle is identically equal to zero in a neighborhood of $\epsilon=0$; critical, otherwise.

The following theorem is shown in [14] (see also the appendix):
Theorem 2.4. Extended Harper's model is supercritical in region I and subcritical in region II. Indeed

- when $\lambda$ belongs to region II, $L_{\lambda}(E)=L\left(\alpha, A_{\lambda, E, \epsilon}\right)=L\left(\alpha, \tilde{A}_{\lambda, E, \epsilon}\right)=0$ on $|\epsilon| \leq \frac{1}{2 \pi} \epsilon_{1}(\lambda)$,
- when $\lambda$ belongs to region II, we have $\hat{\lambda}=\left(\frac{\lambda_{3}}{\lambda_{2}}, \frac{1}{\lambda_{2}}, \frac{\lambda_{1}}{\lambda_{2}}\right)$ belongs to region $I$ and

$$
\begin{equation*}
L_{\hat{\lambda}}(E)=\epsilon_{1}(\lambda), \tag{2.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\epsilon_{1}(\lambda)=\ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{\max \left(\lambda_{1}+\lambda_{3}, 1\right)+\sqrt{\max \left(\lambda_{1}+\lambda_{3}, 1\right)^{2}-4 \lambda_{1} \lambda_{3}}}>0 . \tag{2.4}
\end{equation*}
$$

Fix a $\theta$ and $f \in l^{2}(\mathbb{Z})$. Let $\mu_{\lambda, \theta}^{f}$ be the spectral measure of $H_{\lambda, \theta}$ corresponding to $f$,

$$
\left\langle\left(H_{\lambda, \theta}-z\right)^{-1} f, f\right\rangle=\int_{\mathbb{R}} \frac{1}{E-z} \mathrm{~d} \mu_{\lambda, \theta}^{f}(E)
$$

for $z$ in the resolvent set $\mathbb{C} \backslash \Sigma_{\lambda}$.
The integrated density of states (IDS) is the function $N_{\lambda}: \mathbb{R} \rightarrow[0,1]$ defined by

$$
\begin{equation*}
N_{\lambda}(E)=\int_{\mathbb{T}} \mu_{\lambda, \theta}^{f}(-\infty, E] \mathrm{d} \theta \tag{2.5}
\end{equation*}
$$

where $f \in l^{2}(\mathbb{Z})$ is such that $\|f\|_{l^{2}(\mathbb{Z})}=1$. It is a continuous non-decreasing surjective funtion.

Notice that $\tilde{A}_{\lambda, E}(\theta) \in S L(2, \mathbb{R})$ is homotopic to identity in $C^{0}(\mathbb{T}, S L(2, \mathbb{R}))$, in fact just consider

$$
H_{t}(\lambda, E, \theta)=\frac{1}{\sqrt{|c|(\theta)|c|(\theta-t \alpha)}}\left(\begin{array}{cc}
t(E-v(\theta)) & -|c|(\theta-t \alpha) \\
|c|(\theta) & 0
\end{array}\right)
$$

which establishes a homotopy of $\tilde{A}_{\lambda, E}(\theta)$ to $R_{\frac{1}{4}}$ and hence to the identity. Therefore we can define the rotation number $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)$. Let $\rho_{\lambda}(E)=\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)$. Notice that $\rho_{\lambda}(E)$ is associated to the operator

$$
\left(\tilde{H}_{\lambda, \theta} u\right)_{n}=|c|(\theta+n \alpha) u_{n+1}+|c|(\theta+(n-1) \alpha) u_{n-1}+2 \cos 2 \pi(\theta+n \alpha) u_{n} .
$$

It is easily seen that for each $\theta, \tilde{H}_{\lambda, \theta}$ and $H_{\lambda, \theta}$ differ by a unitary operator, thus they share the same spectrum and integrated density of states, $\tilde{N}_{\lambda}(E)=N_{\lambda}(E)$. The relation between the integrated density of states and rotation number of $\tilde{H}_{\lambda, \theta}$ yields

$$
\begin{equation*}
N_{\lambda}(E)=\tilde{N}_{\lambda}(E)=1-2 \rho_{\lambda}(E) \tag{2.6}
\end{equation*}
$$

2.3. The dual model. It turns out the spectrum $\Sigma_{\lambda}$ of $H_{\lambda, \theta}$ is related to the spectrum $\Sigma_{\hat{\lambda}}$ of $H_{\hat{\lambda}, \theta}$ in the following way:

$$
\Sigma_{\lambda}=\lambda_{2} \Sigma_{\hat{\lambda}}
$$

by Aubry duality. This map $\sigma: \lambda \rightarrow \hat{\lambda}$ establishes the duality between region I and region II. The IDS $N_{\lambda}(E)$ of $H_{\lambda, \theta}$ coincide with the $\operatorname{IDS} N_{\hat{\lambda}}\left(E / \lambda_{2}\right)$ of $H_{\hat{\lambda}, \theta}$. Since $\Sigma_{\lambda}=\lambda_{2} \Sigma_{\hat{\lambda}}$, we have the following
Theorem 2.5 (11, [17]). For any $\lambda, \theta$, there exists a dense set of $E \in \Sigma_{\lambda}$ such that there exists a non-zero solution of $H_{\hat{\lambda}, \theta} u=\frac{E}{\lambda_{2}} u$ with $\left|u_{k}\right| \leq 1+|k|$.
2.4. Bounded eigenfunction for every energy. The next result from [6] allows us to pass from a statement of every $\theta$ to every $E$.

Theorem 2.6 ( $[6])$. If $E \in \Sigma_{\lambda}$, then there exists $\theta(E) \in \mathbb{T}$ and a bounded solution of $H_{\hat{\lambda}, \alpha, \theta} u=\frac{E}{\lambda_{2}} u$ with $u_{0}=1$ and $\left|u_{k}\right| \leq 1$.

### 2.5. Localization and reducibility.

Theorem 2.7. Given $\alpha$ irrational, $\theta \in \mathbb{R}$ and $\lambda$ in region $I I$, fix $E \in \Sigma_{\lambda}$, and suppose $H_{\hat{\lambda}, \theta} u=\frac{E}{\lambda_{2}} u$ has a non-zero exponentially decaying eigenfunction $u=$ $\left\{u_{k}\right\}_{k \in \mathbb{Z}},\left|u_{k}\right| \leq e^{-c|k|}$ for $k$ large enough. Then the following hold:

- (A) If $2 \theta \notin \alpha \mathbb{Z}+\mathbb{Z}$, then there exists $M: \mathbb{R} / \mathbb{Z} \rightarrow S L(2, \mathbb{R})$ analytic, such that

$$
M^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) M(x)=R_{ \pm \theta}
$$

In this case $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)= \pm \theta+\frac{m}{2} \alpha \bmod \mathbb{Z}$, where $m=\operatorname{deg} M$ (here since $M \in S L(2, \mathbb{R})$, we have that $m$ is an even number) and $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right) \notin$ $\alpha \mathbb{Z}+\mathbb{Z}$.

- (B) If $2 \theta \in \alpha \mathbb{Z}+\mathbb{Z}$ and $\alpha \in \mathrm{DC}$, then there exists $M: \mathbb{R} / \mathbb{Z} \rightarrow \operatorname{PSL}(2, \mathbb{R})$ analytic, such that

$$
M^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) M(x)=\left(\begin{array}{cc} 
\pm 1 & a \\
0 & \pm 1
\end{array}\right)
$$

with $a \neq 0$. In this case $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)=\frac{m}{2} \alpha \bmod \mathbb{Z}$, where $m=\operatorname{deg} M$, i.e. $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right) \in \alpha \mathbb{Z}+\mathbb{Z}$.
Proof. Let $u(x)=\sum_{k \in \mathbb{Z}} \hat{u}_{k} e^{2 \pi i k x}, U(x)=\binom{e^{2 \pi i \theta} u(x)}{u(x-\alpha)}$. Then

$$
\begin{aligned}
& A_{\lambda, E}(x) U(x)=e^{2 \pi i \theta} U(x+\alpha), \\
& \tilde{A}_{\lambda, E}(x) \tilde{U}(x)=e^{2 \pi i \theta} \tilde{U}(x+\alpha)
\end{aligned}
$$

Notice $\tilde{U}(x)=Q_{\lambda}(x) U(x)$ is analytic in $|\operatorname{Im} x|<\frac{\tilde{c}}{2 \pi}$, where $\tilde{c}=\min \left(\epsilon_{1}, c\right), \epsilon_{1}$ as in (2.4) and $Q_{\lambda}$ as in Lemma A.2. Define $\overline{\tilde{U}(x)}$ to be the complex conjugate of
$\tilde{U}(x)$ on $\mathbb{T}$ and its analytic extension to $|\operatorname{Im} x|<\frac{\tilde{c}}{2 \pi}$. Let $M(x)$ be the matrix with columns $\tilde{U}(x)$ and $\tilde{U}(x)$. Then,

$$
\tilde{A}_{\lambda, E}(x) M(x)=M(x+\alpha)\left(\begin{array}{cc}
e^{2 \pi i \theta} & 0 \\
0 & e^{-2 \pi i \theta}
\end{array}\right) \quad \text { on } \mathbb{T} \text {. }
$$

Then since $\operatorname{det} M(x+\alpha)=\operatorname{det} M(x)$, we know $\operatorname{det} M(x)$ is a constant on $\mathbb{T}$.
Case 1. If $\operatorname{det} M(x) \neq 0$, then let $M(x)=\tilde{M}(x)\left(\begin{array}{cc}1 & 1 \\ i & -i\end{array}\right)$. Then

$$
\tilde{M}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) \tilde{M}(x)=R_{\theta}=\left(\begin{array}{cc}
\cos 2 \pi \theta & -\sin 2 \pi \theta \\
\sin 2 \pi \theta & \cos 2 \pi \theta
\end{array}\right)
$$

Case 2. If $\operatorname{det} M(x)=0$, then if we denote $\tilde{U}(x)=\binom{u_{1}(x)}{u_{2}(x)}$, then $\operatorname{det} M(x)=0$ means there exists $\eta(x)$ such that $u_{1}(x)=\eta(x) \overline{u_{1}(x)}$ and $u_{2}(x)=\eta(x) \overline{u_{2}(x)}$. This implies that $\eta(x) \in \mathbb{C}^{\omega}(\mathbb{T}, \mathbb{C})$, and $|\eta(x)|=1$ on $\mathbb{T}$. Therefore there exists $\phi(x) \in$ $\mathbb{C}^{\omega}(\mathbb{R} / 2 \mathbb{Z}, \mathbb{C})$ such that $\phi^{2}(x)=\eta(x)$ and $|\phi(x)|=1$. It is easy to see $\overline{\phi(x)} u_{1}(x)=$ $\phi(x) \overline{u_{1}(x)}$ and $\overline{\phi(x)} u_{2}(x)=\phi(x) \overline{u_{2}(x)}$. Then we define $W(x)=\left(\frac{\phi(x)}{\overline{\phi(x)}} u_{2}(x)\right)$, it is a real vector on $\mathbb{R} / 2 \mathbb{Z}$ with $W(x+1)= \pm W(x)$, and $\tilde{U}(x)=\phi(x) W(x)$. Now let us define $\tilde{M}(x)$ to be the matrix with columns $W(x)$ and $\frac{1}{\|W(x)\|^{-2}} R_{\frac{1}{4}} W(x)$; then $\operatorname{det} \tilde{M}(x)=1$ and $\tilde{M}(x) \in \operatorname{PSL}(2, \mathbb{R})$. Since

$$
\tilde{A}_{\lambda, E}(x) W(x)=\frac{e^{2 \pi i \theta} \phi(x+\alpha)}{\phi(x)} W(x+\alpha)
$$

we have

$$
\tilde{A}_{\lambda, E}(x) \tilde{M}(x)=\tilde{M}(x+\alpha)\left(\begin{array}{cc}
d(x) & \tau(x) \\
0 & d(x)^{-1}
\end{array}\right)
$$

where $d(x)=\frac{e^{2 \pi i \theta} \phi(x+\alpha)}{\phi(x)},|d(x)|=1$ and $d(x)$ being a real number, therefore $d(x)= \pm 1$. Also $\tau(x) \in \mathbb{C}^{\omega}(\mathbb{R} / 2 \mathbb{Z}, \mathbb{C})$. But in fact $\tilde{M}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) \tilde{M}(x)$ is well defined on $\mathbb{T}$. Therefore $\tau(x) \in \mathbb{C}^{\omega}(\mathbb{T}, \mathbb{C})$. Now since we assumed $\alpha \in \mathrm{DC}$, we can further reduce $\tau(x)$ to the constant $\tau=\int_{\mathbb{T}} \tau(x) \mathrm{d} x$. In fact there exists $\psi(x) \in \mathbb{C}^{\omega}(\mathbb{T}, \mathbb{C})$ such that $-\psi(x+\alpha)+\psi(x)+\tau(x)=\int_{\mathbb{T}} \tau(x) \mathrm{d} x$. This implies

$$
\left(\begin{array}{cc}
1 & -\psi(x+\alpha) \\
0 & 1
\end{array}\right) \tilde{M}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) \tilde{M}(x)\left(\begin{array}{cc}
1 & \psi(x) \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc} 
\pm 1 & \tau \\
0 & \pm 1
\end{array}\right)
$$

In fact if $\operatorname{det} M(x)=0$, then $\frac{e^{2 \pi i \theta} \phi(x+\alpha)}{\phi(x)}= \pm 1$, which implies that $2 \theta \in \alpha \mathbb{Z}+\mathbb{Z}$. Therefore if $2 \theta \notin \alpha \mathbb{Z}+\mathbb{Z}$, we must be in case (A). If on the other hand, $2 \theta \in \alpha \mathbb{Z}+\mathbb{Z}$, $2 \theta=k \alpha+n$, suppose $\tilde{M}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) \tilde{M}(x)=R_{\theta}$; then

$$
R_{-\frac{k}{2}(x+\alpha)} \tilde{M}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) \tilde{M}(x) R_{\frac{k}{2} x}=R_{\frac{n}{2}}= \pm I
$$

leading to a contradiction. Therefore if $2 \theta \in \alpha \mathbb{Z}+\mathbb{Z}$, we must be in case (B).
2.6. Continued fractions. Let $\left\{q_{n}\right\}$ be the denominators of the continued fraction approximants of $\alpha$. We recall the following properties:

$$
\begin{aligned}
\left\|q_{n} \alpha\right\|_{\mathbb{R} / \mathbb{Z}} & =\inf _{1 \leq|k| \leq q_{n+1}-1}\|k \alpha\|_{\mathbb{R} / \mathbb{Z}} \\
\frac{1}{2 q_{n+1}} & \leq\left\|q_{n} \alpha\right\|_{\mathbb{R} / \mathbb{Z}} \leq \frac{1}{q_{n+1}}
\end{aligned}
$$

Recall that the Diophantine condition of $\alpha$ is $\beta(\alpha)=\lim \sup _{n \rightarrow \infty} \frac{\ln q_{n+1}}{q_{n}}=0$. Thus for any $\xi>0$, there exists $C_{\xi}>0$ such that

$$
\begin{equation*}
\|k \alpha\|_{\mathbb{R} / \mathbb{Z}} \geq C_{\xi} e^{-\xi|k|} \text { for any } k \neq 0 \tag{2.7}
\end{equation*}
$$

Lemma 2.8 ([5). Let $\alpha \in \mathbb{R} \backslash \mathbb{Q}, x \in \mathbb{R}$ and $0 \leq l_{0} \leq q_{n}-1$ be such that

$$
\left|\sin \pi\left(x+l_{0} \alpha\right)\right|=\inf _{0 \leq l \leq q_{n}-1}|\sin \pi(x+l \alpha)|
$$

then for some absolute constant $C_{1}>0$,

$$
-C_{1} \ln q_{n} \leq \sum_{0 \leq l \leq q_{n}-1, l \neq l_{0}} \ln |\sin \pi(x+l \alpha)|+\left(q_{n}-1\right) \ln 2 \leq C_{1} \ln q_{n}
$$

Lemma 2.9 ( 6 ). Let $1 \leq r \leq\left[q_{n+1} / q_{n}\right]$. If $p(x)$ has essential degree at most $k=r q_{n}-1$ and $x_{0} \in \mathbb{R} / \mathbb{Z}$, then for some absolute constant $C_{2}$,

$$
\|p(x)\|_{0} \leq C_{2} q_{n+1}^{C_{2} r} \sup _{0 \leq j \leq k}\left|p\left(x_{0}+j \alpha\right)\right| .
$$

## 3. Main estimates and proof of Theorem 1.1

### 3.1. Almost localization for every $\theta$.

Definition 3.1. Let $\alpha \in \mathbb{R} \backslash \mathbb{Q}, \theta \in \mathbb{R}, \epsilon_{0}>0$. We say that $k$ is an $\epsilon_{0}$-resonance of $\theta$ if $\|2 \theta-k \alpha\| \leq e^{-\epsilon_{0}|k|}$ and $\|2 \theta-k \alpha\|=\min _{|l| \leq|k|}\|2 \theta-l \alpha\|$.
Definition 3.2. Let $0=\left|n_{0}\right|<\left|n_{1}\right|<\ldots$ be the $\epsilon_{0}$-resonances of $\theta$. If this sequence is infinite, we say $\theta$ is $\epsilon_{0}$-resonant, otherwise we say it is $\epsilon_{0}$-non-resonant.

Definition 3.3. We say the extended Harper's model $\left\{H_{\lambda, \alpha, \theta}\right\}_{\theta}$ exhibits almost localization if there exists $C_{0}, C_{3}, \epsilon_{0}, \tilde{\epsilon}_{0}>0$, such that for every solution $\phi$ to $H_{\lambda, \alpha, \theta} \phi=E \phi$ satisfying $\phi(0)=1$ and $|\phi(m)| \leq 1+|m|$, and for every $C_{0}\left(1+\left|n_{j}\right|\right)<$ $|k|<C_{0}^{-1}\left|n_{j+1}\right|$, we have $|\phi(k)| \leq C_{3} e^{-\tilde{\epsilon}_{0}|k|}$ (where $n_{j}$ are the $\epsilon_{0}$-resonances of $\theta$ ).
Theorem 3.1. If $\lambda$ belongs to region II, $\left\{H_{\hat{\lambda}, \alpha, \theta}\right\}_{\theta}$ is almost localized for every $\alpha \in \mathrm{DC}$.

Remark 3.1. It is clear from Theorem 3.1 that almost localization implies localization for non-resonant $\theta$.

We will actually prove the following explicit lemma:
Lemma 3.2. Let $\lambda$ be in region II. Let $C_{4}$ be the absolute constant in Lemma 4.3, $\epsilon_{1}=\epsilon_{1}(\lambda)$ be as in (2.4); then for any $0<\epsilon_{0}<\frac{\epsilon_{1}}{100 C_{4}}$, there exists constant $C_{3}>0$, which depends on $\lambda, \alpha$ and $\epsilon_{0}$, so that for every solution $u$ of $H_{\hat{\lambda}, \alpha, \theta} u=$ Eu satisfying $u(0)=1$ and $\left|u_{k}\right| \leq 1+|k|$, if $3\left(\left|n_{j}\right|+1\right)<|k|<\frac{1}{3}\left|n_{j+1}\right|$, then $\left|u_{k}\right| \leq C_{3} e^{-\frac{\epsilon_{1}}{5}|k|}$, where $\left\{n_{j}\right\}$ are the $\epsilon_{0}$-resonances of $\theta$.

The proof of Lemma 3.2 (and thus of Theorem 3.1) is given in section 4.
3.2. Almost reducibility. Let $\lambda$ be in region II. For every $E \in \Sigma_{\lambda}$, let $\theta(E) \in \mathbb{T}$ be given in Theorem [2.6. Let $0<\epsilon_{0}<\frac{\epsilon_{1}}{100 C_{4}}$ and $\left\{n_{j}\right\}$ be the set of $\epsilon_{0}$-resonances of $\theta(E)$. Then for some positive constants $N_{0}, C$ and $c$, independent of $E$ and $\theta$, we have the following theorem.

Theorem 3.3. For any fixed $j$, with $N_{0}<n=\left|n_{j}\right|+1<\infty$, let $N=\left|n_{j+1}\right|, L^{-1}=$ $\left\|2 \theta-n_{j} \alpha\right\|$. Then there exists $W: \mathbb{T} \rightarrow S L(2, \mathbb{R})$ analytic such that $|\operatorname{deg} W| \leq C n$, $\|W\|_{0} \leq C L^{C}$ and $\left\|W^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) W(x)-R_{\mp \theta}\right\| \leq C e^{-c N}$.

Remark 3.2. Notice that this theorem requires $n>N_{0}$, which is not always ensured when $\theta(E)$ is non-resonant, however in that case we have localization for $H_{\hat{\lambda}, \alpha, \theta}$ instead of almost localization. We will prove Theorem 3.3 in section 5.
3.3. Spectral consequences of almost reducibility. Let $\epsilon_{1}=\epsilon_{1}(\lambda)$ and $C_{4}$ be as in Lemma 3.2 ,

Theorem 3.4. Assume $\alpha \in$ DC. For $\lambda$ in region II, fix $E \in \Sigma_{\lambda}$. Assume $\theta(E) \in \mathbb{T}$ is such that $H_{\hat{\lambda}, \alpha, \theta} u=\frac{E}{\lambda_{2}} u$ has solution satisfying $u_{0}=1$ and $\left|u_{k}\right| \leq 1$. Let $C$ be the constant in Theorem 3.3. Then $\theta(E)$ and $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)$ have the following relation:

- (A) If $\theta$ is $\epsilon_{0}$-non-resonant for some $\frac{\epsilon_{1}}{100 C_{4}}>\epsilon_{0}>0$, then $2 \theta \in \mathbb{Z} \alpha+\mathbb{Z}$ if and only if $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right) \in \mathbb{Z} \alpha+\mathbb{Z}$.
- (B) If $\theta$ is $\epsilon_{0}$-resonant for some $\frac{\epsilon_{1}}{100 C_{4}}>\epsilon_{0}>0$, then $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)$ is $\frac{\epsilon_{0}}{C+2}-$ resonant.

Proof. (A) When $\theta$ is $\epsilon_{0}$-non-resonant for some $\frac{\epsilon_{1}}{100 C_{4}}>\epsilon_{0}>0$, Theorem 3.1implies $H_{\hat{\lambda}, \alpha, \theta}$ has exponentially decaying eigenfunction. Then applying Theorem [2.7 we get $2 \theta \in \mathbb{Z} \alpha+\mathbb{Z}$ if and only if $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right) \in \mathbb{Z} \alpha+\mathbb{Z}$.
(B) Assume $\theta$ is $\epsilon_{0}$-resonant for some $\frac{\epsilon_{1}}{100 C_{4}}>\epsilon_{0}>0$. Fix any $\xi<\frac{\epsilon_{0}}{2 C+2}$; then there exists $C_{\xi}>0$ such that for any $k \neq 0$ we have $\|k \alpha\| \geq C_{\xi} e^{-\xi|k|}$. Now take an $\epsilon_{0}$-resonance $n_{j}$ of $\theta$ such that $n=\left|n_{j}\right|>\max \left(\frac{-\ln C_{\xi} / 2}{\epsilon_{0}-(2 C+2) \xi}, N_{0}\right)$. Then there exists $|m| \leq C n$ such that $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right)-m \alpha=-2 \theta$. Then

$$
\left\|2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right)-\left(m-n_{j}\right) \alpha\right\|=\left\|2 \theta-n_{j} \alpha\right\|<e^{-\epsilon_{0} n} \leq e^{-\frac{\epsilon_{0}}{C+2}\left|m-n_{j}\right|} .
$$

Take any $|l| \leq\left|m-n_{j}\right|, l \neq m-n_{j}$. Then

$$
\left\|\left(l-\left(m-n_{j}\right)\right) \alpha\right\| \geq C_{\xi} e^{-2 \xi\left|m-n_{j}\right|}>2 e^{-\epsilon_{0} n}>2\left\|2 \rho\left(\alpha, \tilde{A}_{E}\right)-\left(m-l_{0}\right) \alpha\right\| .
$$

Thus $\left\|2 \rho\left(\alpha, \tilde{A}_{E}\right)-l \alpha\right\|>\left\|2 \rho\left(\alpha, \tilde{A}_{E}\right)-\left(m-n_{j}\right) \alpha\right\|$ for any $|l| \leq\left|m-n_{j}\right|, l \neq m-n_{j}$. This by definition means $\rho\left(\alpha, \tilde{A}_{\lambda, E}\right)$ is $\frac{\epsilon_{0}}{C+2}$-resonant.

Now based on Theorem [3.4, we can complete the proof of the dry version of the Ten Martini problem for extended Harper's model in regions I and II.

Proof of Theorem 1.1. It is enough to consider $\lambda$ in region II. Let $E \in \Sigma_{\lambda}$ be such that $N_{\lambda}(E) \in \mathbb{Z} \alpha+\mathbb{Z}$. We are going to show $E$ belongs to the boundary of a component of $\mathbb{R} \backslash \Sigma_{\lambda}$. Now by (2.6) we have $2 \rho\left(\alpha, \tilde{A}_{\lambda, E}\right) \in \alpha \mathbb{Z}+\mathbb{Z}$, thus by Theorem [3.4, $2 \theta(E) \in \alpha \mathbb{Z}+\mathbb{Z}$. By Theorem [2.7, this means there exist $M(x) \in$ $C_{h}^{\omega}(\mathbb{T}, \operatorname{PSL}(2, \mathbb{R}))$ such that $M^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) M(x)=\left(\begin{array}{cc} \pm 1 & a \\ 0 & \pm 1\end{array}\right)$. Without
loss of generality, we assume $M^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) M(x)=\left(\begin{array}{ll}1 & a \\ 0 & 1\end{array}\right)$. Let $\tilde{M}(x)=$ $\frac{M(x)}{\sqrt{|c|(x-\alpha)}}$; then

$$
\tilde{M}^{-1}(x+\alpha)\left(\begin{array}{cc}
\frac{E-v(x)}{|c|(x)} & -\frac{|c|(x-\alpha)}{|c|(x)} \\
1 & 0
\end{array}\right) \tilde{M}(x)=\left(\begin{array}{ll}
1 & a \\
0 & 1
\end{array}\right) .
$$

Now let $\tilde{M}(x)=\left(\begin{array}{ll}M_{11}(x) & M_{12}(x) \\ M_{21}(x) & M_{22}(x)\end{array}\right)$. Then $M_{21}(x)=M_{11}(x-\alpha)$ and $M_{22}(x)=$ $M_{12}(x-\alpha)-a M_{11}(x-\alpha)$ and

$$
\begin{aligned}
& \tilde{M}^{-1}(x+\alpha)\left(\begin{array}{cc}
\frac{E+\epsilon-v(x)}{|c|(x)} & -\frac{|c|(x-\alpha)}{|c|(x)} \\
1
\end{array}\right) \tilde{M}(x) \\
& \quad=\left(\begin{array}{cc}
1 & a \\
0 & 1
\end{array}\right)+\epsilon\left(\begin{array}{cc}
M_{11}(x) M_{12}(x)-a M_{11}^{2}(x) & M_{12}^{2}(x)-a M_{11}(x) M_{12}(x) \\
\triangleq-M_{11}^{2}(x) & -M_{11}(x) M_{12}(x)
\end{array}\right) \\
& \triangleq M_{0}+\epsilon M_{1}(x)
\end{aligned}
$$

Now we look for $Z_{\epsilon}(x)$ of the form $e^{\epsilon Y(x)}$ such that

$$
Z_{\epsilon}^{-1}(x+\alpha)\left(M_{0}+\epsilon M_{1}(x)\right) Z_{\epsilon}(x)=M_{0}+\epsilon\left[M_{1}\right]+O\left(\epsilon^{2}\right) .
$$

We then just need to solve the equation:

$$
\begin{aligned}
& \left(I-\epsilon Y(x+\alpha)+O\left(\epsilon^{2}\right)\right)\left(M_{0}+\epsilon M_{1}(x)\right)\left(I+\epsilon Y(x)+O\left(\epsilon^{2}\right)\right) \\
& \quad=M_{0}+\epsilon\left[M_{1}\right]+O\left(\epsilon^{2}\right) .
\end{aligned}
$$

It is sufficient to solve the coholomogical equation:

$$
Y(x+\alpha) M_{0}-M_{0} Y(x)=M_{1}(x)-\left[M_{1}\right],
$$

which is guaranteed by the Diophantine condition on $\alpha$. Thus

$$
\begin{aligned}
& \left(M(x+\alpha) Z_{\epsilon}(x+\alpha)\right)^{-1} \tilde{A}_{\lambda, E}(x)\left(M(x) Z_{\epsilon}(x)\right) \\
& \quad=\left(\begin{array}{cc}
1+\epsilon\left[M_{11} M_{12}\right]-a \epsilon\left[M_{11}^{2}\right] & a+\epsilon\left[M_{12}^{2}\right]-a \epsilon\left[M_{11} M_{12}\right] \\
-\epsilon\left[M_{11}^{2}\right] & 1-\epsilon\left[M_{11} M_{12}\right]
\end{array}\right)+O\left(\epsilon^{2}\right) \\
& \quad \triangleq M_{\epsilon}+O\left(\epsilon^{2}\right) .
\end{aligned}
$$

Notice that $\tilde{A}_{\lambda, E}$ is uniformly hyperbolic iff $\operatorname{Trace}\left(M_{\epsilon}\right)>2$ which is fulfilled when $-a \epsilon\left[M_{11}^{2}\right]>0$. Thus for $\epsilon$ small, satisfying $-a \epsilon\left[M_{11}^{2}\right]>0, E+\epsilon \notin \Sigma_{\lambda}$, which means this spectral gap is open.

## 4. Almost localization in region I

In this section we will prove Lemma 3.2 For fixed $\lambda$ in region II and $E$, let $D_{\hat{\lambda}, E}(\theta)=c_{\hat{\lambda}}(\theta) A_{\hat{\lambda}, E}(\theta)$, where $c_{\hat{\lambda}}(\theta)=\frac{\lambda_{3}}{\lambda_{2}} e^{-2 \pi i\left(\theta+\frac{\alpha}{2}\right)}+\frac{1}{\lambda_{2}}+\frac{\lambda_{1}}{\lambda_{2}} e^{2 \pi i\left(\theta+\frac{\alpha}{2}\right)}$. Regarding the Lyapunov exponent, we recall the following result in [14]:

$$
L\left(\alpha, A_{\hat{\lambda}, E}\right)=L\left(\alpha, D_{\hat{\lambda}, E}\right)-\int_{\mathbb{T}} \ln \left|c_{\hat{\lambda}}(\theta)\right| \mathrm{d} \theta \triangleq \tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|>0
$$

where $\tilde{L}=\ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{2}}$ and $\int \ln \left|c_{\hat{\lambda}}\right|=\ln \frac{\max \left(\lambda_{1}+\lambda_{3}, 1\right)+\sqrt{\max \left(\lambda_{1}+\lambda_{3}, 1\right)^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{2}}$.

Proof of of Lemma 3.2. Suppose $u$ is a solution satisfying the condition of Lemma 3.2. For an interval $I=\left[x_{1}, x_{2}\right]$, let $\Gamma_{I}$ be the coupling operator between $I$ and $\mathbb{Z} \backslash I$ :

$$
\Gamma_{I}(i, j)=\left\{\begin{array}{lr}
\tilde{c}\left(\theta+\left(x_{1}-1\right) \alpha\right), & (i, j)=\left(x_{1}, x_{1}-1\right) \\
c\left(\theta+\left(x_{1}-1\right) \alpha\right), & (i, j)=\left(x_{1}-1, x_{1}\right), \\
\tilde{c}\left(\theta+x_{2} \alpha\right), & (i, j)=\left(x_{2}+1, x_{2}\right), \\
c\left(\theta+x_{2} \alpha\right), & (i, j)=\left(x_{2}, x_{2}+1\right), \\
0, & \text { otherwise }
\end{array}\right.
$$

Let $H_{I}=R_{I} H_{\hat{\lambda}, \theta} R_{I}^{*}$ be the restricted operator of $H_{\hat{\lambda}, \theta}$ to $I$. Then for $x \in I$, we have $\left(H_{I}+\Gamma_{I}-E\right) u(x)=0$. Thus $u(x)=G_{I} \Gamma_{I} u(x)$, where $G_{I}=\left(E-H_{I}\right)^{-1}$. By matrix multiplication,

$$
\begin{aligned}
u(x) & =\sum_{y \in I,(y, z) \in \Gamma_{I}} G_{I}(x, y) \Gamma_{I}(y, z) u(z) \\
& =\tilde{c}\left(\theta+\left(x_{1}-1\right) \alpha\right) G_{I}\left(x, x_{1}\right) u\left(x_{1}-1\right)+c\left(\theta+x_{2} \alpha\right) G_{I}\left(x, x_{2}\right) u\left(x_{2}+1\right) .
\end{aligned}
$$

Let us denote $P_{k}(\theta)=\operatorname{det}\left(E-H_{[0, k-1]}(\theta)\right)$. Then the $k$-step matrix $D_{\hat{\lambda}, E, k}(\theta)$ satisfies

$$
D_{\hat{\lambda}, E, k}(\theta)=\left(\begin{array}{cc}
P_{k}(\theta) & -\tilde{c}(\theta-\alpha) P_{k-1}(\theta+\alpha) \\
c(\theta+(k-1) \alpha) P_{k-1}(\theta) & -\tilde{c}(\theta-\alpha) c(\theta+(k-1) \alpha) P_{k-2}(\theta+\alpha)
\end{array}\right) .
$$

This relation between $P_{k}(\theta)$ and $D_{\hat{\lambda}, E, k}(\theta)$ gives a general upper bound of $P_{k}(\theta)$ in terms of $\tilde{L}$. Indeed by Lemma 2.1 for any $\epsilon>0$ there exists $C(\epsilon)>0$ so that

$$
\left|P_{n}(\theta)\right| \leq C(\epsilon) e^{(\tilde{L}+\epsilon) n} \text { for any } n \in \mathbb{N}
$$

By Cramer's rule,

$$
\begin{aligned}
\left|G_{I}\left(x_{1}, y\right)\right| & =\prod_{j=x_{1}}^{y-1}|c(\theta+j \alpha)|\left|\frac{\operatorname{det}\left(E-H_{\left[y+1, x_{2}\right]}(\theta)\right)}{\operatorname{det}\left(E-H_{I}(\theta)\right)}\right| \\
& =\prod_{j=x_{1}}^{y-1}|c(\theta+j \alpha)|\left|\frac{P_{x_{2}-y}(\theta+(y+1) \alpha)}{P_{k}\left(\theta+x_{1} \alpha\right)}\right|, \\
\left|G_{I}\left(y, x_{2}\right)\right| & =\prod_{j=y+1}^{x_{2}}|c(\theta+j \alpha)|\left|\frac{\operatorname{det}\left(E-H_{\left[x_{1}, y-1\right]}(\theta)\right)}{\operatorname{det}\left(E-H_{I}(\theta)\right)}\right| \\
& =\prod_{j=y+1}^{x_{2}}|c(\theta+j \alpha)|\left|\frac{P_{y-x_{1}}\left(\theta+x_{1} \alpha\right)}{P_{k}\left(\theta+x_{1} \alpha\right)}\right| .
\end{aligned}
$$

Notice that $P_{k}(\theta)$ is an even function about $\theta+\frac{k-1}{2} \alpha$, it can be written as a polynomial of degree $k$ in $\cos 2 \pi\left(\theta+\frac{k-1}{2} \alpha\right)$. Let $P_{k}(\theta)=Q_{k}\left(\cos 2 \pi\left(\theta+\frac{k-1}{2} \alpha\right)\right)$. Let $M_{k, r}=\left\{\theta \in \mathbb{T},\left|Q_{k}(\cos 2 \pi \theta)\right| \leq e^{(k+1) r}\right\}$.

Definition 4.1. Fix $m>0$. A point $y \in \mathbb{Z}$ is called $(k, m)$-regular if there exists an interval $\left[x_{1}, x_{2}\right]$ containing $y$, where $x_{2}=x_{1}+k-1$ such that

$$
\left|G_{I}\left(y, x_{i}\right)\right| \leq e^{-m\left|y-x_{i}\right|} \text { and } \operatorname{dist}\left(y, x_{i}\right) \geq \frac{1}{3} k \text { for } i=1,2,
$$

otherwise $y$ is called $(k, m)$-singular.

Lemma 4.1. Suppose $y \in \mathbb{Z}$ is $\left(k, \tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|-\rho\right)$-singular. Then for any $\epsilon>0$ and any $x \in \mathbb{Z}$ satisfying $y-\frac{2}{3} k \leq x \leq y-\frac{1}{3} k$, we have $\theta+\left(x+\frac{1}{2}(k-1)\right) \alpha$ belongs to $M_{k, \tilde{L}-\frac{1}{3} \rho+\epsilon}$ for $k>k(\lambda, \epsilon, \rho)$.
Proof. Suppose there exists $\epsilon>0$ and $x_{1}: y-(1-\delta) k \leq x_{1} \leq y-\delta k$, such that $\theta+\left(x_{1}+\frac{1}{2}(k-1)\right) \alpha$ does not belong to $M_{k, \tilde{L}-\frac{1}{3} \rho+\epsilon}$, that is, $\left|P_{k}\left(\theta+x_{1} \alpha\right)\right|>$ $e^{(k+1)(\tilde{L}-\rho \delta+\epsilon)}$,

$$
\begin{aligned}
\left|G_{I}\left(x_{1}, y\right)\right| & \leq \prod_{j=x_{1}}^{y-1}\left|c_{\hat{\lambda}}(\theta+j \alpha)\right| e^{\left(k-\left|x_{1}-y\right|\right)(\tilde{L}+\epsilon)} e^{-(k+1)\left(\tilde{L}-\frac{1}{3} \rho+\epsilon\right)} \\
& <e^{-\left(\tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|-\rho\right)\left|y-x_{1}\right|} \text { for } k>k(\lambda, \epsilon, \rho) .
\end{aligned}
$$

Similarly

$$
\left|G_{I}\left(x_{2}, y\right)\right| \leq e^{-\left(\tilde{L}-\int \ln \left|c_{\grave{\lambda}}\right|-\rho\right)\left|y-x_{2}\right|}
$$



Definition 4.2. We say that the set $\left\{\theta_{1}, \ldots, \theta_{k+1}\right\}$ is $\gamma$-uniform if

$$
\max _{x \in[-1,1]} \max _{i=1, \ldots, k+1} \prod_{j=1, j \neq i}^{k+1} \frac{\left|x-\cos 2 \pi \theta_{j}\right|}{\left|\cos 2 \pi \theta_{i}-\cos 2 \pi \theta_{j}\right|}<e^{k \gamma} .
$$

Lemma 4.2. Let $\gamma_{1}<\gamma$. If $\theta_{1}, \ldots, \theta_{k+1} \in M_{k, \tilde{L}-\gamma}$, then $\left\{\theta_{1}, \ldots, \theta_{k+1}\right\}$ is not $\gamma_{1-}$ uniform for $k>k\left(\gamma, \gamma_{1}\right)$.

Proof. Otherwise, using Lagrange interpolation form we can get $\left|Q_{k}(x)\right|<e^{k \tilde{L}}$ for all $x \in[-1,1]$. This implies $\left|P_{k}(x)\right|<e^{k \tilde{L}}$ for all $x$. But by Herman's subharmonic function argument, $\int_{\mathbb{R} / \mathbb{Z}} \ln \left|P_{k}(x)\right| \mathrm{d} x \geq k \tilde{L}$. This is impossible.

Now take $\xi$ and $\epsilon_{0}$ such that $0<1000 \xi<\epsilon_{0}$. Then for $\left|n_{j+1}\right|>N(\xi)$ we have

$$
\begin{aligned}
2 e^{-4 \xi\left|n_{j+1}\right|} & \leq C_{\xi} e^{-2 \xi\left|n_{j+1}\right|} \leq\left\|\left(n_{j+1}-n_{j}\right) \alpha\right\| \\
& =\left\|n_{j+1} \alpha-2 \theta+2 \theta-n_{j} \alpha\right\| \leq 2\left\|2 \theta-n_{j} \alpha\right\| \leq 2 e^{-\epsilon_{0}\left|n_{j}\right|}
\end{aligned}
$$

which yields that

$$
\begin{equation*}
\left|n_{j+1}\right|>\frac{\epsilon_{0}}{4 \xi}\left|n_{j}\right|>250\left|n_{j}\right| . \tag{4.1}
\end{equation*}
$$

Without loss of generality, assume $3\left(\left|n_{j}\right|+1\right)<y<\frac{\left|n_{j+1}\right|}{3}$ and $y>N(\xi)$. Select $n$ such that $q_{n} \leq \frac{y}{8}<q_{n+1}$ and let $s$ be the largest positive integer satisfying $s q_{n} \leq \frac{y}{8}$. Set $I_{1}, I_{2} \subset \mathbb{Z}$ as follows:

$$
\begin{aligned}
& I_{1}=\left[1-2 s q_{n}, 0\right] \text { and } I_{2}=\left[y-2 s q_{n}+1, y+2 s q_{n}\right], \text { if } n_{j}<0, \\
& I_{1}=\left[0,2 s q_{n}-1\right] \text { and } I_{2}=\left[y-2 s q_{n}+1, y+2 s q_{n}\right], \text { if } n_{j} \geq 0 .
\end{aligned}
$$

Lemma 4.3. Let $\theta_{j}=\theta+j \alpha$; then set $\left\{\theta_{j}\right\}_{j \in I_{1} \cup I_{2}}$ is $C_{4} \epsilon_{0}+C_{4} \xi$-uniform for some absolute constant $C_{4}$ and $y>y\left(\alpha, \epsilon_{0}, \xi\right)$.

Proof. Without loss of generality, we assume $n_{j}>0$. Take $x=\cos 2 \pi a$. Now it suffices to estimate

$$
\sum_{j \in I_{1} \cup I_{2}, j \neq i}\left(\ln \left|\cos 2 \pi a-\cos 2 \pi \theta_{j}\right|-\ln \left|\cos 2 \pi \theta_{i}-\cos 2 \pi \theta_{j}\right|\right) \triangleq \sum_{1}-\sum_{2}
$$

Lemma 2.8 reduces this problem to estimating the minimal terms.
First we estimate $\sum_{1}$ :

$$
\begin{aligned}
\sum_{1} & =\sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln \left|\cos 2 \pi a-\cos 2 \pi \theta_{j}\right| \\
& =\sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln \left|\sin \pi\left(a+\theta_{j}\right)\right|+\sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln \left|\sin \pi\left(a-\theta_{j}\right)\right|+\left(6 s q_{n}-1\right) \ln 2 \\
& \triangleq \sum_{1,+}+\sum_{1,-}+\left(6 s q_{n}-1\right) \ln 2 .
\end{aligned}
$$

We cut $\sum_{1,+}$ or $\sum_{1,-}$ into $6 s$ sums and then apply Lemma 2.8. We get that for some absolute constant $C_{1}$ :

$$
\sum_{1} \leq-6 s q_{n} \ln 2+C_{1} s \ln q_{n}
$$

Next, we estimate $\sum_{2}$ :

$$
\begin{aligned}
\sum_{2}= & \sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln \left|\cos 2 \pi \theta_{j}-\cos 2 \pi \theta_{i}\right| \\
= & \sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln |\sin \pi(2 \theta+(i+j) \alpha)| \\
& \quad+\sum_{j \in I_{1} \cup I_{2}, j \neq i} \ln |\sin \pi(i-j) \alpha|+\left(6 s q_{n}-1\right) \ln 2 \\
\triangleq & \sum_{2,+}+\sum_{2,-}+\left(6 s q_{n}-1\right) \ln 2
\end{aligned}
$$

We need to carefully estimate the minimal terms. For $\sum_{2,+}$, we use the property of resonant set; and for $\sum_{2,-}$, we use the Diophantine condition on $\alpha$.

For any $0<|j|<q_{n+1}$, we have $\|j \alpha\| \geq\left\|q_{n} \alpha\right\| \geq C_{\xi} e^{-\xi q_{n}}$. Therefore

$$
\max (\ln |\sin x|, \ln |\sin (x+\pi j \alpha)|) \geq-2 \xi q_{n} \text { for } y>y(\alpha, \xi) .
$$

This means in any interval of length $s q_{n}$, there can be at most one term which is less than $-2 \xi q_{n}$. Then there can be at most 6 such terms in total.

For the part $\sum_{2,-}$, since $\|(i-j) \alpha\| \geq C_{\xi} e^{-\xi|i-j|} \geq e^{-20 \xi s q_{n}}$, these 6 smallest terms must be bounded by $-20 \xi s q_{n}$ from below. Hence $\sum_{2,-} \geq-6 s q_{n} \ln 2-$ $C \xi s q_{n}-C s \ln q_{n}$ for $y>y(\xi)$ and some absolute constant $C$.

For the part $\sum_{2,+}$, notice $|i+j| \leq 2 y+4 s q_{n}<3 y<\left|n_{j+1}\right|$ and $i+j>0>-n_{j}$. Suppose $\left\|2 \theta+k_{0} \alpha\right\|=\min _{j \in I_{1} \cup I_{2}}\|2 \theta+(i+j) \alpha\| \leq e^{-100 \epsilon_{0} s q_{n}}<e^{-\epsilon_{0}\left|k_{0}\right|}$. Then for any $|k| \leq\left|k_{0}\right| \leq 40 s q_{n}$ (including $\left|n_{j}\right|$ ),

$$
\|2 \theta-k \alpha\| \geq\left\|\left(k+k_{0}\right) \alpha\right\|-\left\|2 \theta+k_{0} \alpha\right\|>\left\|2 \theta+k_{0} \alpha\right\| \text { for } y>y\left(\alpha, \epsilon_{0}, \xi\right)
$$

This means $-k_{0}$ must be a $\epsilon_{0}$-resonance, therefore $\left|k_{0}\right| \leq\left|n_{j-1}\right|$. Then

$$
\begin{aligned}
\left\|2 \theta-n_{j} \alpha\right\| & \geq\left\|\left(n_{j}+k_{0}\right) \alpha\right\|-\left\|2 \theta+k_{0} \alpha\right\| \\
& \geq C_{\xi} e^{-12 \xi s q_{n}}-e^{-100 \epsilon_{0} s q_{n}}>e^{-100 \epsilon_{0} s q_{n}} \geq\left\|2 \theta+k_{0} \alpha\right\|
\end{aligned}
$$

leads to a contradiction. Thus the smallest terms must be greater than $-100 \epsilon_{0} s q_{n}$. We can bound $\sum_{2,+}$ by $-6 s q_{n} \ln 2-600 \epsilon_{0} s q_{n}-12 \xi s q_{n}-C s \ln q_{n}$ from below. Therefore $\sum_{2} \geq-6 s q_{n} \ln 2-C \epsilon_{0} s q_{n}-C \xi s q_{n}-C s \ln q_{n}$. Thus the set $\left\{\theta_{j}\right\}_{j \in I_{1} \cup I_{2}}$ is $C_{4} \epsilon_{0}+C_{4} \xi$-uniform for $y>y\left(\alpha, \epsilon_{0}, \xi\right)$ and some absolute constant $C_{4}$.

Now let $C_{4}$ be the absolute constant in Lemma 4.3. Choose $0<1000 \xi<\epsilon_{0}<$ $\frac{\epsilon_{1}}{100 C_{4}}$. Combining Lemma 4.2 and Lemma 4.3, we know that when $y>y\left(\alpha, \epsilon_{0}, \xi\right)$, $\left\{\theta_{j}\right\}_{j \in I_{1} \cup I_{2}}$ cannot be inside the set $M_{6 s q_{n}-1, \tilde{L}-2 C_{4} \epsilon_{0}}$ at the same time. Therefore 0 and $y$ cannot be $\left(6 s q_{n}-1, \tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|-9 C_{4} \epsilon_{0}\right)$ at the same time. However 0 is $\left(6 s q_{n}-1, \tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|-9 C_{4} \epsilon_{0}\right)$-singular given $n$ large enough. Therefore

$$
\left\{\theta_{j}\right\}_{j \in I_{1}} \subset M_{6 s q_{n}-1, \tilde{L}-2 C_{4} \epsilon_{0}}
$$

Thus $y$ must be $\left(6 s q_{n}-1, \tilde{L}-\int \ln \left|c_{\hat{\lambda}}\right|-9 C_{4} \epsilon_{0}\right)$-regular. This implies

$$
|u(y)| \leq e^{-\left(\tilde{L}-\int \ln \left|c_{\lambda}\right|-9 C_{4} \epsilon_{0}\right) \frac{1}{4}|y|}<e^{-\frac{\epsilon_{1}}{5}|y|} \text { for }|y| \geq y\left(\lambda, \alpha, \epsilon_{0}, \xi\right) .
$$

Thus there exists $C_{3}=C_{\lambda, \alpha, \epsilon_{0}, \xi}$ such that $|u(y)| \leq C_{3} e^{-\frac{\epsilon_{1}}{5}|y|}$ for any $3\left|n_{j}\right| \leq|y| \leq$ $\frac{1}{3}\left|n_{j+1}\right|$ and $j \in \mathbb{N}$.

## 5. Almost reducibility in Region II

Proof of Theorem 3.3. For any $E \in \Sigma_{\lambda}$, take $\theta(E)$ and $\left\{u_{k}\right\}$ as in Theorem[2.6] Let $\epsilon_{1}$ be as in (2.4), $C_{4}$ be the absolute constant from Lemma 4.3 and $C_{2}$ be the absolute constant from Lemma 2.9, Fix $\max \left(32 C_{2} \xi, 1000 \xi\right)<\epsilon_{0}<\min \left(\frac{\epsilon_{1}}{200}, \frac{\epsilon_{1}}{100 C_{4}}\right)$. By Lemma 3.2, there exists $C$ depending on $\lambda$ and $\alpha$ such that for any $3\left|n_{j}\right|<$ $|k|<\frac{1}{3}\left|n_{j+1}\right|$, we have $\left|u_{k}\right| \leq C e^{-\frac{e_{1}}{5}|k|}$.

For any $n, 9\left|n_{j}\right|<n<\frac{1}{9}\left|n_{j+1}\right|$, of the form

$$
\begin{equation*}
n=r q_{m}-1<q_{m+1} \tag{5.1}
\end{equation*}
$$

Let $u(x)=u^{I}(x)=\sum_{k \in I} u_{k} e^{2 \pi i k x}$ with $I=\left[-\left[\frac{n}{2}\right],\left[\frac{n}{2}\right]\right]=\left[x_{1}, x_{2}\right]$. Define

$$
U(x)=\binom{e^{2 \pi i \theta} u(x)}{u(x-\alpha)}
$$

Let $A(\theta)=A_{\lambda, E}(\theta)$. By direct computation:

$$
A(x) U(x)=e^{2 \pi i \theta} U(x+\alpha)+\binom{g(x)}{0} \triangleq e^{2 \pi i \theta} U(x+\alpha)+G(x)
$$

The Fourier coefficients of $g(x)$ are possibly non-zero only at four points $x_{1}, x_{2}$, $x_{1}-1$ and $x_{2}+1$. Since $\left|u_{k}\right| \leq C_{1} e^{-\frac{1}{5}|k|}$ when $3\left|n_{j}\right|<|k|<\frac{1}{3}\left|n_{j+1}\right|$, we know that $\|G(x)\|_{\frac{\epsilon_{1}}{20 \pi}}^{20} \leq C_{1} e^{-\frac{\epsilon_{1}}{20} n}$.

Combining Lemmas A. 3 and 2.1 we have exponential control of the growth of the transfer matrix, for any $\delta>0$ there exists $C_{\delta}>0$ such that

$$
\left\|\tilde{A}_{k}(x)\right\|_{\frac{\epsilon_{1}}{2 \pi}} \leq C_{\delta} e^{\delta|k|}, \text { for any } k
$$

[^2]With some effort we are able to get the following significantly improved upper bound.

Theorem 5.1. For some $C>0$ depending on $\lambda$ and $\alpha$,

$$
\left\|\tilde{A}_{k}(x)\right\|_{\mathbb{T}} \leq C(1+|k|)^{C}
$$

Proof. Let $\tilde{U}(x)=Q(x) U(x), \tilde{G}(x)=Q(x+\alpha) G(x)$, where $Q=Q_{\lambda}$ is given in Lemma A. 2 Since

$$
\max \left(\|Q(x)\|_{\frac{\epsilon_{1}}{20 \pi}},\left\|Q^{-1}(x)\right\|_{\frac{\epsilon_{1}}{20 \pi}}\right) \leq C
$$

we have

$$
\tilde{A}(x) \tilde{U}(x)=e^{2 \pi i \theta} \tilde{U}(x+\alpha)+\tilde{G}(x)
$$

where $\|\tilde{G}(x)\|_{\frac{\epsilon_{1}}{20 \pi}} \leq C e^{-\frac{\epsilon_{1}}{20} n}$.
Lemma 5.2. Let $C_{2}$ be the constant from Lemma 2.9. Then for any $\delta, 2 C_{2} \xi<$ $\delta<\frac{\epsilon_{0}}{16}$, we have

$$
\inf _{|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{20 \pi}}\|\tilde{U}(x)\| \geq e^{-2 \delta n}
$$

for $n>n(\alpha, \delta)$.
Proof. We will prove the statement by contradiction. Suppose for some $x_{0} \in$ $\left\{|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{20 \pi}\right\}$ we have $\left\|\tilde{U}\left(x_{0}\right)\right\|<e^{-2 \delta n}$. Notice that for any $l \in \mathbb{N}$,

$$
\begin{aligned}
e^{2 \pi i l \theta} \tilde{U}\left(x_{0}+l \alpha\right)= & \tilde{A}_{l}\left(x_{0}\right) \tilde{U}\left(x_{0}\right) \\
& -\sum_{m=1}^{l} e^{2 \pi i(m-1) \theta} \tilde{A}_{l-m}\left(x_{0}+m \alpha\right) \tilde{G}\left(x_{0}+(m-1) \alpha\right) .
\end{aligned}
$$

This implies for $n>n(\delta)$ large enough and for any $0 \leq l \leq n,\left\|\tilde{U}\left(x_{0}+l \alpha\right)\right\| \leq e^{-\delta n}$, thus $\left\|u\left(x_{0}+l \alpha\right)\right\| \leq C_{\delta} e^{-\delta n}$. By Lemma 2.9, $\left\|u\left(x+i \operatorname{Im}\left(x_{0}\right)\right)\right\|_{\mathbb{T}} \leq C_{2} C_{\delta} e^{C_{2} \xi n} e^{-\delta n} \leq$ $e^{-\frac{\delta}{2} n}$. This contradicts with $\int_{\mathbb{T}} u\left(x+i \operatorname{Im}\left(x_{0}\right)\right) \mathrm{d} x=u_{0}=1$.

Lemma 5.3 (3). Let $V: \mathbb{T} \rightarrow \mathbb{C}^{2}$ be analytic in $|\operatorname{Im}(x)|<\eta$. Assume that $\delta_{1}<\|V(x)\|<\delta_{2}^{-1}$ holds on $|\operatorname{Im}(x)|<\eta$. Then there exists $M: \mathbb{T} \rightarrow S L(2, \mathbb{C})$ analytic on $|\operatorname{Im}(x)|<\eta$ with first column $V$ and $\|M\|_{\eta} \leq C \delta_{1}^{-2} \delta_{2}^{-1}\left(1-\ln \left(\delta_{1} \delta_{2}\right)\right)$.

Applying Lemma 5.3 let $M(x)$ be the matrix with first column $\tilde{U}(x)$. Then $e^{-2 \delta n} \leq\|\tilde{U}(x)\|_{\frac{\delta}{\pi}} \leq e^{\delta n}$ and hence $\|M(x)\|_{\frac{\delta}{\pi}} \leq C e^{6 \delta n}$. Therefore

$$
M^{-1}(x+\alpha) \tilde{A}(x) M(x)=\left(\begin{array}{cc}
e^{2 \pi i \theta} & 0 \\
0 & e^{-2 \pi i \theta}
\end{array}\right)+\left(\begin{array}{cc}
\beta_{1}(x) & b(x) \\
\beta_{3}(x) & \beta_{4}(x)
\end{array}\right)
$$

where $\left\|\beta_{1}(x)\right\|_{\frac{\delta}{\pi}},\left\|\beta_{3}(x)\right\|_{\frac{\delta}{\pi}},\left\|\beta_{4}(x)\right\|_{\frac{\delta}{\pi}} \leq C e^{-\frac{\epsilon_{1}}{40} n}$, and $\|b(x)\|_{\frac{\delta}{\pi}} \leq C e^{13 \delta n}$. Let

$$
\Phi(x)=M(x)\left(\begin{array}{cc}
e^{\frac{\epsilon_{1}}{100} n} & 0 \\
0 & e^{-\frac{\epsilon_{1}}{160} n}
\end{array}\right)
$$

Then we would have

$$
\Phi(x+\alpha)^{-1} \tilde{A}(x) \Phi(x)=\left(\begin{array}{cc}
e^{2 \pi i \theta} & 0 \\
0 & e^{-2 \pi i \theta}
\end{array}\right)+H(x)
$$

where $\|H(x)\|_{\frac{\delta}{\pi}} \leq C e^{-\frac{\epsilon_{1}}{160} n}$, and $\|\Phi(x)\|_{\frac{\delta}{\pi}} \leq C e^{\frac{\epsilon_{1}}{80} n}$. Thus

$$
\sup _{0 \leq s \leq e^{\frac{\epsilon}{320} n}}\left\|\tilde{A}_{s}(x)\right\|_{\mathbb{T}} \leq e^{\frac{\epsilon_{1}}{20} n}
$$

for $n \geq n(\lambda, \alpha)$ satisfying (5.1). For $s$ large, there always exists $9\left|n_{j}\right|<n<\frac{1}{9}\left|n_{j+1}\right|$ satisfying (5.1) such that $c n \leq \frac{320}{\epsilon_{1}} \ln s \leq n$ with some absolute constant $c$. Thus there exists $C$ depending on $\lambda$ and $\alpha$ such that $\left\|\tilde{A}_{k}(x)\right\|_{\mathbb{T}} \leq C(1+|k|)^{C}$.

Now we come back to the proof of Theorem 3.3. Fix some $n=\left|n_{j}\right|$, and $N=$ $\left|n_{j+1}\right|$. Let $u(x)=u^{I_{2}}(x)$ with $I_{2}=\left[-\left[\frac{N}{9}\right],\left[\frac{N}{9}\right]\right]$ and $U(x)=\binom{e^{2 \pi i \theta} u(x)}{u(x-\alpha)}$. Then

$$
A(x) U(x)=e^{2 \pi i \theta} U(x+\alpha)+G(x) \text { with }\|G(x)\|_{\frac{\epsilon_{1}}{20 \pi}} \leq C e^{-\frac{\epsilon_{1}}{90} N}
$$

Define $U_{0}(x)=e^{\pi i n_{j} x} U(x)$. Notice that if $n_{j}$ is even, then $U_{0}(x)$ is well-defined on $\mathbb{T}$, otherwise $U_{0}(x+1)=-U_{0}(x)$. Then

$$
\tilde{A}(x) \tilde{U}_{0}(x)=e^{2 \pi i \tilde{\theta}} \tilde{U}_{0}(x+\alpha)+H(x)
$$

where $\tilde{\theta}=\theta-\frac{n_{j}}{2} \alpha, \tilde{U}_{0}(x)=Q(x) U_{0}(x)$ and $\|H(x)\|_{\frac{\epsilon_{1}}{20 \pi}} \leq C e^{-\frac{\epsilon_{1}}{100} N}$. Consider the matrix $W(x)$ with $\tilde{U}_{0}(x)$ and $\overline{\tilde{U}_{0}(x)}$ being its two columns. Then

$$
\tilde{A}(x) W(x)=W(x+\alpha)\left(\begin{array}{cc}
e^{2 \pi i \tilde{\theta}} & 0 \\
0 & e^{-2 \pi i \tilde{\theta}}
\end{array}\right)+\tilde{H}(x)
$$

Theorem 5.4. Let $L^{-1}=\left\|2 \theta-n_{j} \alpha\right\|$. Then for $n>N_{0}(\lambda, \alpha)$ we have

$$
|\operatorname{det} W(x)| \geq L^{-4 C} \quad \text { for any } x \in \mathbb{T}
$$

where $C$ is the constant appeared in Theorem 5.1.
Proof. First, we fix $\xi_{1}<\frac{\epsilon_{0}}{1600}$ so that $\|k \alpha\| \geq C_{\xi_{1}} e^{-\xi_{1}|k|}$ for any $k \neq 0$. We have the following estimate about $L$ :
Lemma 5.5. $e^{\epsilon_{0} n} \leq L \leq e^{4 \xi_{1} N}$. This can be seen by the following inequality:

$$
e^{-2 \xi_{1} N} \leq\left\|\left(n_{j+1}-n_{j}\right) \alpha\right\| \leq 2\left\|n_{j} \alpha-2 \theta\right\|=2 L^{-1} \leq 2 e^{-\epsilon_{0} n} \text { for } n \geq N\left(\xi_{1}\right)
$$

Now we prove by contradiction. Suppose there exists $\kappa$ and $x_{0} \in \mathbb{T}$ such that $\left\|\tilde{U}_{0}\left(x_{0}\right)-\kappa \tilde{\tilde{U}_{0}\left(x_{0}\right)}\right\|<L^{-4 C}$. Then

$$
\begin{aligned}
\| \tilde{U}_{0}\left(x_{0}\right. & +l \alpha) e^{2 \pi i l \tilde{\theta}}-\kappa \overline{\tilde{U}_{0}\left(x_{0}+l \alpha\right)} e^{-2 \pi i l \tilde{\theta}} \| \\
\leq & \| \sum_{m=0}^{l-1} \tilde{A}_{l-m}\left(x_{0}+m \alpha\right) H\left(x_{0}+m \alpha\right) \\
& -\kappa \sum_{m=0}^{l-1} \tilde{A}_{l-m}\left(x_{0}+m \alpha\right) \overline{H\left(x_{0}+m \alpha\right)}\|+\| A_{l}\left(x_{0}\right) \| L^{-4 C} \\
\quad \leq & C L^{2 C} e^{-\frac{e_{1}}{100} N}+C L^{-2 C}<L^{-C}
\end{aligned}
$$

for $0 \leq|l| \leq L^{2}$. If we take $j=\frac{L}{4}$, then

$$
\begin{equation*}
\left\|\tilde{U}_{0}\left(x_{0}+\frac{L}{4} \alpha\right)+\kappa \overline{\tilde{U}_{0}\left(x_{0}+\frac{L}{4} \alpha\right)}\right\|<L^{-1} \tag{5.2}
\end{equation*}
$$

Next since $\left\|U_{0}(x)\right\|_{\mathbb{T}} \leq n$, we have $\left\|\tilde{U}_{0}(x)\right\|_{\mathbb{T}} \leq C n$. Thus

$$
\left\|\tilde{U}_{0}\left(x_{0}+l \alpha\right)-\kappa \overline{\tilde{U}_{0}\left(x_{0}+l \alpha\right)}\right\|<L^{-\frac{1}{3}} \text { for } 0 \leq|l| \leq L^{\frac{1}{2}} .
$$

For any analytic function $f(x)=\sum_{k \in \mathbb{Z}} \hat{f}_{k} e^{2 \pi i k x}$, define

$$
f_{[-m, m]}(x)=\sum_{|k| \leq m} \hat{f}_{k} e^{2 \pi i k x} .
$$

For any column vector $V(x)=\binom{v^{(1)}(x)}{v^{(2)}(x)}$, let $V_{[-m, m]}(x)=\binom{v_{[-m, m]}^{(1)}(x)}{v_{[-m, m]}^{(2)}(x)}$. Now let us define $\tilde{U}_{0}^{[9 n]}(x)=Q(x) e^{\pi i n_{j} x} U_{[-9 n, 9 n]}(x)$. Then

$$
\left\|\tilde{U}_{0}^{[9 n]}(x)-\tilde{U}_{0}(x)\right\|_{\mathbb{T}} \leq C e^{-\frac{9}{5} \epsilon_{1} n}
$$

Consider $\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}(x)\right]_{[-18 n, 18 n]}(x) e^{\pi i n_{j} x}$. This function differs from a polynomial with essential degree $36 n$ only by a multiple of $e^{\pi i n_{j} x}$. Notice that $Q(x)$ is analytic in $\left\{x:|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{4 \pi}\right\}$, thus $|\hat{Q}(k)| \leq C e^{-\frac{\epsilon_{1}}{2}|k|}$. Then

$$
\left|e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}(k)\right| \leq \sum_{|m| \leq 9 n}|\hat{Q}(k-m) \hat{U}(m)| \leq C n e^{-\frac{\epsilon_{1}}{2}(|k|-9 n)} \quad \text { for }|k| \geq 18 n
$$

Thus

$$
\begin{gathered}
\left\|e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}(x)-\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x)\right\|_{\mathbb{T}} \leq e^{-4 \epsilon_{1} n}, \\
\left\|\tilde{U}_{0}(x)-\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x) e^{\pi i n_{j} x}\right\|_{\mathbb{T}} \leq e^{-4 \epsilon_{1} n}
\end{gathered}
$$

Hence

$$
\begin{aligned}
& \|\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}\left(x_{0}+l \alpha\right) e^{2 \pi i n_{j}\left(x_{0}+l \alpha\right)} \\
& \quad-\kappa \overline{\left.\kappa e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}\left(x_{0}+l \alpha\right)} \|_{\mathbb{T}}<2 L^{-\frac{1}{3}}+e^{-4 \epsilon_{1} n}
\end{aligned}
$$

for $|l| \leq L^{\frac{1}{2}}$. Notice that

$$
\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x) e^{2 \pi i n_{j} x}-\kappa \overline{\kappa\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x)}
$$

is a polynomial whose essential degree is at most $37 n$. Thus by Lemma 2.9, we would have

$$
\begin{aligned}
& \|\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x) e^{\pi i n_{j} x} \\
& \quad-\kappa \overline{\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x) e^{\pi i n_{j} x}} \|_{\mathbb{T}}<L^{-\frac{1}{4}}+e^{-2 \epsilon_{1} n} .
\end{aligned}
$$

Hence $\left\|\tilde{U}_{0}(x)-\kappa \overline{\tilde{U}_{0}(x)}\right\|_{\mathbb{T}}<L^{-\frac{1}{4}}+2 e^{-2 \epsilon_{1} n}$. But combining with (5.2) we would get $\left\|\tilde{U}_{0}\left(x_{0}+\frac{L}{4} \alpha\right)\right\|<2 L^{-\frac{1}{4}}+2 e^{-2 \epsilon_{1} n}$, but this contradicts with $\inf _{x \in \mathbb{T}}\left\|\tilde{U}_{0}(x)\right\|>e^{-2 \delta n}$ since $\delta<\frac{\epsilon_{0}}{16}$.

Now for $n>N_{0}(\lambda, \alpha)$, take $S(x)=\operatorname{Re} \tilde{U}_{0}(x)$ and $T(x)=\operatorname{Im} \tilde{U}_{0}(x)$. Let $W_{1}(x)$ be the matrix with columns $S(x)$ and $T(x)$. Notice that $\operatorname{det} W_{1}(x)$ is well-defined on $\mathbb{T}$ and $\operatorname{det} W_{1}(x) \neq 0$ on $\mathbb{T}$, hence without loss of generality we could assume $\operatorname{det} W_{1}(x)>0$ on $\mathbb{T}$; otherwise we simply take $W_{1}(x)$ to be the matrix with columns $S(x)$ and $-T(x)$. Then

$$
\left\|\tilde{A}(x) W_{1}(x)-W_{1}(x+\alpha) R_{-\tilde{\theta}}\right\|_{\mathbb{T}} \leq C e^{-\frac{\epsilon_{1} N}{45}}
$$

By taking determinant, we get

$$
\operatorname{det} W_{1}(x)=\operatorname{det} W_{1}(x+\alpha)+O\left(e^{-\frac{\epsilon_{1}}{50} N}\right) \text { on } \mathbb{T} .
$$

Since $\operatorname{det} W_{1}(x)$ is analytic on $|\operatorname{Im} x| \leq \frac{\epsilon_{1}}{20 \pi}$, by considering the Fourier coefficients we could get

$$
\operatorname{det} W_{1}(x)=w_{0}+O\left(e^{-\frac{\epsilon_{1}}{100} N}\right) \text { on } \mathbb{T}
$$

where $w_{0} \geq L^{-5 C}$. Thus $\operatorname{det} W_{1}(x)$ is almost a positive constant.
Define $W_{2}(x)=\operatorname{det} W_{1}(x)^{-\frac{1}{2}} W_{1}(x)$. Then $W_{2}(x) \in C^{\omega}(\mathbb{T})$ and $\operatorname{det} W_{2}(x)=1$. We have

$$
\begin{gathered}
W_{2}^{-1}(x+\alpha) \tilde{A}(x) W_{2}(x)=\frac{\operatorname{det} W_{1}(x+\alpha)^{\frac{1}{2}}}{\operatorname{det} W_{1}(x)^{\frac{1}{2}}} R_{-\tilde{\theta}}+O\left(e^{-\frac{\epsilon_{1}}{100} N}\right) \text { on } \mathbb{T}, \\
W_{2}^{-1}(x+\alpha) \tilde{A}(x) W_{2}(x)=R_{-\tilde{\theta}}+O\left(e^{-\frac{\epsilon_{1}}{200} N}\right) \text { on } \mathbb{T} .
\end{gathered}
$$

Now let's prove $\operatorname{deg} W_{2}(x) \leq 36 n$. $\operatorname{deg} W_{2}(x)$ is the same as the degree of its columns. For $M: \mathbb{R} / 2 \mathbb{Z} \rightarrow \mathbb{R}^{2}$, we say $\operatorname{deg} M=k$ if $M$ is homotopic to $\binom{\cos k \pi x}{\sin k \pi x}$.

For some constant $c>0$, we obviously have

$$
\int_{\mathbb{T}}\|S(x)\| \mathrm{d} x+\int_{\mathbb{T}}\|T(x)\| \mathrm{d} x \geq \int_{\mathbb{T}}\|S(x)+i T(x)\| \mathrm{d} x=\int_{\mathbb{T}}\left\|\tilde{U}_{0}(x)\right\| \mathrm{d} x \geq c
$$

Without loss of generality we could assume $\int_{\mathbb{T}}\|S(x)\| \mathrm{d} x>\frac{c}{2}$. Also

$$
\tilde{A}(x) S(x)=S(x+\alpha) \cos 2 \pi \tilde{\theta}-T(x+\alpha) \sin 2 \pi \tilde{\theta}+O\left(e^{-\frac{\epsilon_{1}}{45} N}\right) \text { on } \mathbb{T}
$$

Then since $\|2 \tilde{\theta}\|=L^{-1}$,

$$
\tilde{A}(x) S(x)=S(x+\alpha)+O\left(L^{-\frac{1}{2}}\right) \text { on } \mathbb{T}
$$

First we prove $\inf _{x \in \mathbb{T}}\|S(x)\| \geq e^{-2 \epsilon_{1} n}$. Suppose otherwise. Then there exists $x_{0} \in \mathbb{T}$, so that $\left\|S\left(x_{0}\right)\right\|<e^{-2 \epsilon_{1} n}$. Then $\left\|\operatorname{Re} \tilde{U}_{0}\left(x_{0}+l \alpha\right)\right\|<e^{-\frac{\epsilon_{0}}{8} n}$ for $|l|<e^{\frac{\epsilon_{0}}{4 C} n}$, where $C$ is the constant that appeared in Theorem 5.1. We have already shown that

$$
\left\|\tilde{U}_{0}(x)-\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]} e^{\pi i n_{j} x}\right\|_{\mathbb{T}}<e^{-4 \epsilon_{1} n}
$$

Thus

$$
\left\|\operatorname{Re}\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}\left(x_{0}+l \alpha\right)\right\|<e^{-\frac{\epsilon_{0}}{16} n}
$$

for $|l|<e^{\frac{\epsilon_{0}}{4 C} n}$. However $\operatorname{Re}\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}$ is a polynomial with essential degree at most $36 n$. Using Lemma 2.9 we are able to get

$$
\left\|\operatorname{Re}\left[e^{-\pi i n x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]} e^{\pi i n_{j} x}\right\|_{\mathbb{T}}<e^{-\frac{\epsilon_{0}}{32} n}
$$

and thus $\left\|\operatorname{Re} \tilde{U}_{0}(x)\right\|_{\mathbb{T}}<e^{-\frac{\epsilon_{0}}{64} n}$ which is a contradiction to $\int_{\mathbb{T}}\left\|\operatorname{Re} \tilde{U}_{0}(x)\right\| \mathrm{d} x>\frac{c}{2}$. In the meantime, we also get

$$
\left\|S(x)-\operatorname{Re}\left[e^{-\pi i n_{j} x} \tilde{U}_{0}^{[9 n]}\right]_{[-18 n, 18 n]}(x) e^{\pi i n_{j} x}\right\|_{\mathbb{T}} \triangleq\|S(x)-h(x)\|_{\mathbb{T}} \leq e^{-4 \epsilon_{1} n}
$$

The first column of $W_{2}(x)$ is $\operatorname{det} W_{1}(x)^{-\frac{1}{2}} S(x)$. We have

$$
\begin{aligned}
& \left\|\frac{S(x)}{\operatorname{det} W_{1}(x)^{\frac{1}{2}}}-\frac{h(x)}{w_{0} \frac{1}{2}}\right\| \\
& \quad \leq \frac{1}{\left|\operatorname{det} W_{1}(x)^{\frac{1}{2}}\right|}\left\|S(x)-h(x)+\left(1-\frac{\operatorname{det} W_{1}(x)^{\frac{1}{2}}}{w_{0}{ }^{\frac{1}{2}}}\right) h(x)\right\| \\
& \quad \leq L^{2 C}\left(e^{-4 \epsilon_{1} n}+L^{8 C} e^{-\frac{\epsilon_{1}}{100} N}\right) \\
& \quad \leq e^{-3 \epsilon_{1} n}<\left\|\frac{S(x)}{\operatorname{det} W_{1}(x)^{\frac{1}{2}}}\right\| \text { on } \mathbb{T} .
\end{aligned}
$$

Thus by Rouché's theorem $\left|\operatorname{deg} W_{2}(x)\right|=|\operatorname{deg} h(x)| \leq 19 n$. Notice that

$$
\left|\rho\left(\alpha, W_{2}^{-1} \tilde{A} W_{2}\right)+\tilde{\theta}\right|<C e^{-\frac{\epsilon_{1}}{200} N}
$$

Then, by (2.2) for some $|m| \leq 19 n$ :

$$
\left|\rho(\alpha, \tilde{A})-\frac{m}{2} \alpha+\tilde{\theta}\right|<C e^{-\frac{\epsilon_{1}}{200} N} .
$$

## Appendix A

When $\lambda$ belongs to region II, let $\epsilon_{2}=\ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{\lambda_{1}+\lambda_{3}+\sqrt{\left(\lambda_{1}+\lambda_{3}\right)^{2}-4 \lambda_{1} \lambda_{3}}}>\epsilon_{1}$. Then $c(x)$ is analytic and non-zero on $|\operatorname{Im}(x)|<\frac{\epsilon_{2}}{2 \pi}$. Furthermore, the winding number of $c(\cdot+i \epsilon)$ is equal to zero when $|\epsilon|<\frac{\epsilon_{2}}{2 \pi}$.

Lemma A.1. When $\lambda$ belongs to region II, we can find an analytic function $f(x)$ on $|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{2 \pi}$ such that $c(x)=|c|(x) e^{f(x+\alpha)-f(x)}$ and $\tilde{c}(x)=|c|(x) e^{-f(x+\alpha)+f(x)}$.

Proof. Since the winding numbers of $c(x)$ and $\tilde{c}(x)$ are 0 on $|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{2 \pi}$, there exist analytic functions $g_{1}(x)$ and $g_{2}(x)$ on $|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{2 \pi}$, such that $c(x)=e^{g_{1}(x)}$ and $\tilde{c}(x)=e^{g_{2}(x)}$. Notice that

$$
\begin{aligned}
& \int_{\mathbb{T}} \ln |c(x)| \mathrm{d} x=\int_{\mathbb{T}} \ln |\tilde{c}(x)| \mathrm{d} x \\
& \int_{\mathbb{T}} \arg c(x) \mathrm{d} x=\int_{\mathbb{T}} \arg \tilde{c}(x) \mathrm{d} x
\end{aligned}
$$

so there exists an analytic function $f(x)$ such that $2 f(x+\alpha)-2 f(x)=g_{1}(x)-g_{2}(x)$. Then $c(x)=|c|(x) e^{f(x+\alpha)-f(x)}$.

Lemma A.2. When $\lambda$ belongs to region II, there exists an analytic matrix $Q_{\lambda}(x)$ defined on $|\operatorname{Im}(x)| \leq \frac{\epsilon_{1}}{2 \pi}$ such that

$$
Q_{\lambda}^{-1}(x+\alpha) \tilde{A}_{\lambda, E}(x) Q_{\lambda}(x)=A_{\lambda, E}(x)
$$

Proof.

$$
\begin{aligned}
\tilde{A}_{\lambda, E}(x) & =\frac{1}{\sqrt{|c|(x)|c|(x-\alpha)}}\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{\tilde{c}(x)}{c(x)}}
\end{array}\right)\left(\begin{array}{cc}
E-v(x) & -\tilde{c}(x-\alpha) \\
c(x) & 0
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{c(x-\alpha)}{\tilde{c}(x-\alpha)}}
\end{array}\right) \\
& =\frac{c(x)}{\sqrt{|c|(x)|c|(x-\alpha)}}\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{\tilde{c}(x)}{c(x)}}
\end{array}\right) A(x)\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{c(x-\alpha)}{\tilde{c}(x-\alpha)}}
\end{array}\right) \\
& =e^{f(x+\alpha)} \sqrt{|c|(x)}\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{\tilde{c}(x)}{c(x)}}
\end{array}\right) A(x)\left\{e^{f(x)} \sqrt{|c|(x-\alpha)}\left(\begin{array}{cc}
1 & 0 \\
0 & \sqrt{\frac{\tilde{c}(x-\alpha)}{c(x-\alpha)}}
\end{array}\right)\right\}^{-1} \\
& =Q_{\lambda}(x+\alpha) A_{\lambda, E}(x) Q_{\lambda}^{-1}(x) .
\end{aligned}
$$

Lemma A.3. If $\alpha$ is irrational, $\lambda$ belongs to region $I I, E \in \Sigma(\lambda)$, then

$$
L\left(\alpha, A_{\lambda, E}(\cdot+i \epsilon)\right)=L\left(\alpha, \tilde{A}_{\lambda, E}(\cdot+i \epsilon)\right)=0
$$

for $|\epsilon| \leq \frac{\epsilon_{1}}{2 \pi}$.
Proof. $L(A(\cdot+i \epsilon))=L(D(\cdot+i \epsilon))-\int \ln |c(x+i \epsilon)| \mathrm{d} x$.

$$
\begin{aligned}
D(x & +i \epsilon) \\
& =\left(\begin{array}{cc}
E-e^{2 \pi i(x+i \epsilon)}-e^{-2 \pi i(x+i \epsilon)} & -\lambda_{1} e^{2 \pi i\left(x-\frac{\alpha}{2}+i \epsilon\right)}-\lambda_{2}-\lambda_{3} e^{-2 \pi i\left(x-\frac{\alpha}{2}+i \epsilon\right)} \\
\lambda_{1} e^{-2 \pi i\left(x+\frac{\alpha}{2}+i \epsilon\right)}+\lambda_{2}+\lambda_{3} e^{2 \pi i\left(x+\frac{\alpha}{2}+i \epsilon\right)}
\end{array}\right) \\
& =e^{2 \pi \epsilon}\left(\begin{array}{cc}
-e^{2 \pi i x}+o(1) & -\lambda_{3} e^{-2 \pi i\left(x-\frac{\alpha}{2}\right)}+o(1) \\
\lambda_{1} e^{-2 \pi i\left(x+\frac{\alpha}{2}\right)}+o(1) & 0
\end{array}\right) .
\end{aligned}
$$

Thus the asymptotic behaviour of $L(D(\cdot+i \epsilon))$ is:

$$
\begin{aligned}
& L(D(\cdot+i \epsilon))=\ln \left|\frac{1+\sqrt{1-4 \lambda_{1} \lambda_{3}}}{2}\right|+2 \pi \epsilon \text { when } \epsilon \rightarrow \infty \\
& L(D(\cdot+i \epsilon))=\ln \left|\frac{1+\sqrt{1-4 \lambda_{1} \lambda_{3}}}{2}\right|-2 \pi \epsilon \text { when } \epsilon \rightarrow-\infty
\end{aligned}
$$

Then it suffices to calculate $\int \ln |c(x+i \epsilon)| \mathrm{d} x$ in region II. We have

$$
\begin{aligned}
& \int \ln |c(x+i \epsilon)| \mathrm{d} x \\
& \quad=\ln \lambda_{3}-2 \pi \epsilon+\int \ln \left|e^{2 \pi i x}-y_{1, \epsilon}\right|+\int \ln \left|e^{2 \pi i x}-y_{2, \epsilon}\right|,
\end{aligned}
$$

where $y_{1, \epsilon}=\frac{-\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{3}} e^{2 \pi \epsilon}$ and $y_{2, \epsilon}=\frac{-\lambda_{2}-\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{3}} e^{2 \pi \epsilon}$.

$$
\begin{aligned}
& \int \ln |c(x+i \epsilon)| \mathrm{d} x \\
& \quad=\left\{\begin{array}{lc}
2 \pi \epsilon+\ln \lambda_{1}, & \epsilon>\frac{1}{2 \pi} \ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{1}}, \\
\ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2}, & \frac{1}{2 \pi} \ln \frac{\lambda_{2}-\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{1}} \leq \epsilon \leq \frac{1}{2 \pi} \ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{1}}, \\
-2 \pi \epsilon+\ln \lambda_{3}, & \epsilon<\frac{1}{2 \pi} \ln \frac{\lambda_{2}-\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{2 \lambda_{1}}
\end{array}\right.
\end{aligned}
$$

Thus $L(A(\cdot+i \epsilon))=0$ when $|\epsilon| \leq \frac{1}{2 \pi} \ln \frac{\lambda_{2}+\sqrt{\lambda_{2}^{2}-4 \lambda_{1} \lambda_{3}}}{\max \left(1, \lambda_{1}+\lambda_{3}\right)+\sqrt{\max \left(1, \lambda_{1}+\lambda_{3}\right)^{2}-4 \lambda_{1} \lambda_{3}}}=\frac{\epsilon_{1}}{2 \pi}$.

Since $\tilde{A}_{\lambda, E}(x+i \epsilon)=Q_{\lambda}(x+\alpha+i \epsilon) A_{\lambda, E}(x+i \epsilon) Q_{\lambda}^{-1}(x+i \epsilon)$, the statement about $\tilde{A}_{\lambda, E}$ is also true.

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[^1]:    ${ }^{1}$ In general one cannot take $M \in C^{\omega}(\mathbb{T}, S L(2, \mathbb{R}))$.

[^2]:    ${ }^{2}$ The existence of such $n$ comes from (4.1).

