Localization in One-dimensional Quasi-periodic Nonlinear Systems

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Abstract

To investigate localization in one-dimensional quasi-periodic nonlinear systems, we consider the Schrödinger equation

$$\mathrm{i}\dot{q}_n + \epsilon(q_{n+1} + q_{n-1}) + V(n\tilde{\alpha} + x)q_n + |q_n|^2 q_n = 0, \quad n \in \mathbb{Z},$$

as a model, with V a nonconstant real-analytic function on \mathbb{R}/\mathbb{Z} , and $\tilde{\alpha}$ satisfying a certain Diophantine condition. It is shown that, if ϵ is sufficiently small, then for a.e. $x \in \mathbb{R}/\mathbb{Z}$, dynamical localization is maintained for "typical" solutions in a quasi-periodic time-dependent way.

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1 Introduction

Physical motivations. Localization of particles and waves in disordered media is one of the most intriguing phenomena in solid-state physics. This phenomenon was first analyzed by P.W.Anderson[1]. He studied the transport of non-interacting electrons in a crystal lattice, described by a single particle with random on-site energy. In his model, he showed that when the amplitude of the disorder becomes higher than a critical value, the diffusion in the lattice of an initially localized wavepacket is suppressed. An Anderson localized state is characterized by an exponential decay of the amplitude of the wave function.

In many physics experiments, a relatively weak disorder on the structure of the lattice is introduced by a quasi-periodic potential. This kind of system corresponds to an experimental realization of the so-called Aubry-André[2] or Harper[19] model. It is important in the study of Bose-Einstein condensation and nonlinear optics. Anderson localization in such linear systems, especially in the one-dimensional case, has been thoroughly studied[31], and rigorous mathematical results have been established[22]. As a well-known model in mathematical physics, the almost Mathieu operator $H_{x,\lambda,\tilde{\alpha}}$ acting on $\ell^2(\mathbb{Z})$ is defined by

$$(H_{x,\lambda,\tilde{\alpha}}\psi)_n = (\psi_{n+1} + \psi_{n-1}) + \lambda \cos 2\pi (x + n\tilde{\alpha})\psi_n, \quad n \in \mathbb{Z},$$

where n is the primary lattice site index, $\tilde{\alpha}$ is some ratio between the wavenumbers of two lattices, $x \in \mathbb{R}/\mathbb{Z}$ is an arbitrary phase, and ψ_n is a complex variable whose modulus square gives the probability of finding a particle at the lattice site n. With $\tilde{\alpha}$ a fixed Diophantine number, for a.e. x and λ large enough, $H_{x,\lambda,\tilde{\alpha}}$ exhibits dynamical localization[16, 17], i.e., for any $\psi \in \ell^2(\mathbb{Z})$ with compact support and arbitrary d > 0,

$$\sup_{t} r^{(d)}(t) := \sup_{t} \sum_{n \in \mathbb{Z}} n^{2d} |(e^{iH_{x,\lambda,\tilde{\alpha}}t}\psi)_n|^2 < \infty.$$

In particular, there exists a transparent transition between diffusion and localization for the almost Mathieu operator. From the perspective of spectrum theory, it is shown by Jitomirskaya[22] that, for a.e. x, $H_{x,\lambda,\tilde{\alpha}}$ has

- 1. $\lambda > 2$: only pure point spectrum with exponentially decaying eigenfunctions;
- 2. $\lambda = 2$: purely singular-continuous spectrum;
- 3. $\lambda < 2$: purely absolutely continuous spectrum.

There is a perfect agreement with this conclusion in some experiments(e.g., [20]). For $\tilde{\alpha} = \frac{\sqrt{5}-1}{2}$, with an initial δ -function wavepacket, the asymptotic spreading of the wavepacket width $r^{(1)}(t)$ can be parametrized as $r^{(1)}(t) \sim t^{\gamma}$, and one finds three different regimes

1. $\lambda > 2$: localized regime, $\gamma = 0$;

- 2. $\lambda = 2$: sub-diffusive, $\gamma \sim 0.5$;
- 3. $\lambda < 2$: ballistic regime, $\gamma = 1$.

However, the situation is much less clear in the presence of interactions (nonlinearities). It strongly influences the possibility to observe the localization induced by disorder. One can start from the Gross-Pitaevskii (GP) equation [18, 27] in Hartree-Fock theory, and get a generalized Aubry-André model which includes an additional nonlinear term that represents the mean-field interaction. The Hamiltonian is

$$H = \sum_{n \in \mathbb{Z}} \left[(\psi_{n+1} \bar{\psi}_n + \bar{\psi}_{n+1} \psi_n) + \lambda \cos 2\pi (n\tilde{\alpha} + x) |\psi_n|^2 + \frac{1}{2} \beta |\psi_n|^4 \right],$$

and the equation of motion is generated by $i\dot{\psi}_n = -\frac{\partial H}{\partial \bar{\psi}_n}$, yielding the nonlinear Schrödinger equation

$$i\dot{\psi}_n + (\psi_{n+1} + \psi_{n-1}) + \lambda \cos 2\pi (n\tilde{\alpha} + x)\psi_n + \beta |\psi_n|^2 \psi_n = 0, \quad n \in \mathbb{Z},$$
 (1.1)

that can be considered as the GP equation on a discretized lattice. Similar versions of a discretized GP equation have been already used to investigate the dynamics of condensates in different situations(see, for instance, [33]).

It is shown experimentally in [25] that, if the condensate initially occupies a single lattice site, i.e., a δ -function $\psi_n(0) = \delta_{n,0}$, the dynamics of the interacting gas is dominated by self-trapping in a wide range of parameters, even for weak interaction. Conversely, if the diffusion starts from a Gaussian wavepacket of width σ , $\psi_n(0) = ce^{-\frac{n^2}{2\sigma^2}}$, then self-trapping is significantly suppressed and the destruction of localization by interaction is more easily observable.

The aim of the present work is to analyze localization in the quasi-periodic nonlinear dynamical systems, which are modeled by the discrete one-dimensional disordered non-linear Schrödinger equations of the same form as (1.1). When the disorder is sufficiently large, a rigorous mathematical argument for the maintainability of localization is given in this paper, corresponding to the experimental conclusion in [25].

Related mathematical works. In the theory of mathematical physics, localization in disordered, nonlinear dynamical systems was initiated by Fröhlich-Spencer-Wayne[12](Similar work was also done by Pöschel[28] and Vittot-Bellissard[34]), who constructed infinite-dimensional, compact invariant tori for a large class of non-coupling systems

$$i\dot{q}_n + V_n q_n + \sum_{m \in \mathbb{Z}} \epsilon_{mn} (q_m + \bar{q}_m)^2 q_n = 0, \quad n \in \mathbb{Z},$$

via the KAM techniques, where $\{V_n\}_{n\in\mathbb{Z}}$ are i.i.d. random variables, ϵ_{mn} are sufficiently small and vanish for |m-n| large enough. Solutions which lie on such tori are localized for all times. Besides the conclusion, they raised the following conjecture in that paper.

Conjecture. [12] Consider the equation

$$i\dot{q}_n + \epsilon(q_{n+1} + q_{n-1}) + V_n q_n + \delta |q_n|^2 q_n = 0, \quad n \in \mathbb{Z},$$
 (1.2)

with $\{V_n\}_{n\in\mathbb{Z}}$ i.i.d. random variables. If ϵ and δ are small enough, with the equation in a large class, then for "most" initial conditions ("Most", e.g., with respect to the uniform measure on finite-dimensional unit spheres.), $q(0) = (q_n(0))_{n\in\mathbb{Z}}$, of finite support, the solutions $q(t) = (q_n(t))_{n\in\mathbb{Z}}$ of (1.2) satisfy

$$\overline{\lim}_{t \to \infty} t^{-1} \sum_{n \in \mathbb{Z}} n^2 |q_n(t)|^2 = 0.$$

Recently, there are several breakthroughs on such problem. For a large class of equation in (1.2), Bourgain-Wang[7] constructed a quasi-periodic solution when ϵ , δ are sufficiently small. They also considered the slightly tempered equations[8], with the nonlinearity replaced by $\lambda_n |q_n|^2 q_n$, $|\lambda_n| < \epsilon (1+|n|)^{-\tau}$ for some small $\tau > 0$. By constructing symplectic transformations to create energy barriers, they proved that, if ϵ is sufficiently small, then the diffusion bound (i.e., the H^1 norm) of the solution grows polynomially with t almost surely. Moreover, a Nekhoroshev-type result about Equation(1.2) was given by Wang-Zhang [35], who proved the long time Anderson localization for arbitrary ℓ^2 initial data.

By establishing an abstract KAM theorem, Geng-Zhao[15] constructed small-amplitude time quasi-periodic solutions of the lattice Schrödinger equation

$$i\dot{q}_n + \epsilon(q_{n+1} + q_{n-1}) + \tan \pi(x + n\tilde{\alpha})q_n + \epsilon|q_n|^2 q_n = 0, \quad n \in \mathbb{Z},$$
(1.3)

for most of $x \in \mathbb{R}/\mathbb{Z}$ if ϵ is sufficiently small and $\tilde{\alpha}$ is Diophantine. This is based on the work by Bellissard-Lima-Scoppola[4], which have studied the linear operator corresponding to Equation (1.3), the well-known Maryland model. The operator, which has dense point spectrum, describes media with no resonance, and this provides convenience for the KAM iteration.

Statement of the main result. Inspired by the conjecture in [12], we try to establish a nonlinear version of "dynamical localization" in the quasi-periodic potential case.

Consider the one-dimensional nonlinear Schrödinger equation

$$i\dot{q}_n + \epsilon(q_{n+1} + q_{n-1}) + V(n\tilde{\alpha} + x)q_n + |q_n|^2 q_n = 0, \quad n \in \mathbb{Z},$$
 (1.4)

for $0 < \epsilon \ll 1$, with V a nonconstant real-analytic function on \mathbb{R}/\mathbb{Z} , and $\tilde{\alpha} \in \mathbb{R}$ is a Diophantine number, i.e., there exist $\tilde{\tau} > 1$ and $\tilde{\gamma} > 0$ such that $(|x|_1)$ is the absolute value of x modulo 1 defined so that $0 \le |x|_1 \le \frac{1}{2}$.)

$$|n\tilde{\alpha}|_1 \ge \frac{\tilde{\gamma}}{|n|^{\tilde{\tau}}}, \quad n \ne 0.$$
 (1.5)

The nonconstant real-analytic potential V, as in [11], is a smooth function in the Gevrey class

$$\sup_{x \in \mathbb{R}/\mathbb{Z}} |\partial^m V(x)| \le CL^m (m!)^2, \quad m \ge 0, \tag{1.6}$$

for some C, L > 0, and satisfying the transversality condition

$$\max_{0 \le m \le \tilde{s}} |\partial_{\varphi}^{m}(V(x+\varphi) - V(x))| \ge \tilde{\xi} > 0, \quad \forall x, \forall \varphi,$$
(1.7)

$$\max_{0 \le m \le \tilde{s}} |\partial_x^m (V(x+\varphi) - V(x))| \ge \tilde{\xi} |\varphi|_1, \quad \forall x, \forall \varphi, \tag{1.8}$$

for some $\tilde{\xi}$, $\tilde{s} > 0$. Clearly, the case $V(x) = \cos 2\pi x$ is included.

Based on an earlier KAM mechanism which was introduced by Eliasson[11], we construct an abstract KAM theorem, and apply this theorem to prove well-localization of Equation (1.4) for typical initial data. From the KAM perspective, the main technical challenges in this work are the following:

- i) Unlike the model in [12], we need to tackle with the second order perturbation in the Hamiltonian;
- ii) Different from the method in [7], our proof is developed from the traditional KAM method;
- iii) Compared with the work in [15], the main difficulty is that the corresponding linear operator has dense point spectrum with *infinitely many resonances*.

The main result can be stated as follows.

Theorem 1 Consider the one-dimensional nonlinear Schrödinger equation

$$i\dot{q}_n + \epsilon(q_{n+1} + q_{n-1}) + V(n\tilde{\alpha} + x)q_n + |q_n|^2 q_n = 0, \quad n \in \mathbb{Z},$$
 (1.9)

with V a nonconstant real-analytic function on \mathbb{R}/\mathbb{Z} , and $\tilde{\alpha} \in \mathbb{R}$ a Diophantine number. Given an integer b > 1, and any $\mathcal{J} = \{n_1, \dots, n_b\} \subset \mathbb{Z}$, there exists a sufficiently small $\epsilon_* = \epsilon_*(V, \tilde{\alpha}, \mathcal{J})$, such that if $0 < \epsilon < \epsilon_*$, then the following holds for a.e. $x \in \mathbb{R}/\mathbb{Z}$.

There exists a Cantor set $\mathcal{O}_{\epsilon} = \mathcal{O}_{\epsilon}(x) \subset [0,1]^b$ with $|[0,1]^b \setminus \mathcal{O}_{\epsilon}| \to 0^{-1}$ as $\epsilon \to 0$ such that the solution $q(t) = (q_n(t))_{n \in \mathbb{Z}}$ of Equation (1.9), with initial datum $q(0) \in \mathcal{O}_{\epsilon}$ supported on \mathcal{J} , satisfies, for any fixed d > 0,

$$\sup_{t} \sum_{n \in \mathbb{Z}} n^{2d} |q_n(t)|^2 < \infty.$$

Moreover, for each $n \in \mathbb{Z}$, $q_n(t)$ is quasi-periodic in time.

Remark 1.1 The quasi-periodic solutions we obtained are not necessarily small-amplitude, since the nonlinearity $|q_n|^2 q_n$ is integrable. Moreover, if the nonlinearity is "diagonal dominant" with some short-range decay, e.g.,

$$|q_n|^2 q_n + \epsilon \sum_{m \neq n} e^{-\varrho|m-n|} |q_m|^2 q_n,$$

the theorem above can also be obtained.

Remark 1.2 Smallness assumption on ϵ is necessary, otherwise the result is not true even for the linear problem. This is different from the random potential case.

2 Preliminaries

From now on, in the formulations and proofs of various assertions, we shall encounter absolute constants depending on the Hamiltonian, the dimension and so on. All such constants will be denoted by c, c_1, c_2, \dots , and sometimes even different constants will be denoted by the same symbol.

¹Hereafter, we use the symbol $|\mathcal{O}|$ to denote the Lebesgue measure of $\mathcal{O} \subset \mathbb{R}^b$.

2.1 The quasi-periodic Schrödinger operator

Consider the Schrödinger operator $T = T(x) : \ell^2(\mathbb{Z}) \to \ell^2(\mathbb{Z})$ defined as

$$(Tq)_n := \epsilon(q_{n+1} + q_{n-1}) + V(x + n\tilde{\alpha})q_n, \quad n \in \mathbb{Z}, \tag{2.1}$$

with V and $\tilde{\alpha}$ as in Equation (1.9). It is well-known from [11] that if ϵ is sufficiently small, then for a.e. $x \in \mathbb{R}/\mathbb{Z}$, the spectrum of T(x) is pure point. We refer to [6, 9, 13, 21, 23, 32] for other works on the pure point spectrum and Anderson localization of quasi-periodic Schrödinger operators.

Now, we are going to represent the main idea of [11], which is critical for the KAM iteration in this paper. Let us start with some notations for infinite-dimensional matrices. Given an infinite-dimensional matrix D, with $D_{mn} \in \mathbb{R}$ the $(m, n)^{\text{th}}$ entry, for a subset $\Lambda \subset \mathbb{Z}$, we define $\Lambda^{\perp} := \mathbb{Z} \setminus \Lambda$,

$$\mathbb{R}^{\Lambda} := \{ n \in \mathbb{R}^{\mathbb{Z}} : n_i = 0 \text{ if } i \notin \Lambda \}, \quad D_{\Lambda} := \left\{ \begin{array}{ll} D_{mn}, & m, n \in \Lambda \\ \delta_{mn}, & \text{otherwise} \end{array} \right..$$

Then $D_{\Lambda}: \mathbb{R}^{\Lambda} + \mathbb{R}^{\Lambda^{\perp}} \to \mathbb{R}^{\Lambda} + \mathbb{R}^{\Lambda^{\perp}}$, acts as $\mathbb{R}^{\Lambda} \hookrightarrow \mathbb{R}^{\mathbb{Z}} \xrightarrow{D} \mathbb{R}^{\mathbb{Z}} \xrightarrow{\text{iproj}} \mathbb{R}^{\Lambda}$ on the first component and as the identity on the second component. (When there is no risk for confusion, we will use D_{Λ} also to denote its first component.)

Let $D_0 = \operatorname{diag}\{V(x+n\tilde{\alpha})\}_{n\in\mathbb{Z}}$ and $Z_0 = \epsilon\Delta$ with Δ the discrete Laplacian. With $\varepsilon_0 = \epsilon^{\frac{1}{4}}$, $\sigma_0 = 1$ and any

$$M_0 \ge \max \left\{ 2^{\tilde{s}+4} C \frac{L^{\tilde{s}+1} ((\tilde{s}+1)!)^2}{\tilde{\xi}}, 2\tilde{\tau}, 8 \right\}, \quad N_0 \ge 1, \quad \rho_0 = N_0^{-1},$$

one can define the following sequences as in [11],

$$M_{\nu+1} = M_{\nu}^{\tilde{s}M_{\nu}^{3}}, \quad a_{\nu} = \frac{1}{\tilde{\tau}} M_{\nu}^{-3\tilde{s}M_{\nu}^{3}}, \quad \varepsilon_{\nu+1} = \varepsilon_{\nu}^{\frac{1}{2}\varepsilon_{\nu}^{-a_{\nu}/2}},$$

$$N_{\nu+1} = \varepsilon_{\nu}^{-a_{\nu}}, \qquad \rho_{\nu+1} = \varepsilon_{\nu}^{a_{\nu}}, \qquad \sigma_{\nu+1} = \frac{1}{3}\rho_{\nu}.$$
(2.2)

These sequences of parameters will be applied in the KAM iteration in this paper.

Proposition 1 There exists a constant $\epsilon_0 = \epsilon_0(C, L, \tilde{\xi}, \tilde{s}, \tilde{\gamma}, \tilde{\tau})$ such that the following holds for the operator (2.1) if $0 < \epsilon < \epsilon_0$.

Fix any $x \in \mathbb{R}/\mathbb{Z}$. There exists a sequence of orthogonal matrices U_{ν} , $\nu = 1, 2, \dots$, with

$$|(U_{\nu} - I_{\mathbb{Z}})_{mn}| \le \varepsilon_0^{\frac{1}{2}} e^{-\frac{3}{2}\sigma_{\nu}|m-n|},$$

such that $U_{\nu}^*(D_0 + Z_0)U_{\nu} = D_{\nu} + Z_{\nu}$, where Z_{ν} is a symmetric matrix satisfying

$$|(Z_{\nu})_{mn}| \leq \varepsilon_{\nu} e^{-\rho_{\nu}|m-n|}$$

and D_{ν} is a symmetric matrix which can be block-diagonalized via an orthogonal matrix Q_{ν} with

$$(Q_{\nu})_{mn} = 0 \ if \ |m-n| > N_{\nu}.$$

More precisely, there is a disjoint decomposition $\bigcup_i \Lambda_i^{\nu} = \mathbb{Z}$ such that

$$\tilde{D}^{\nu} = Q_{\nu}^* D_{\nu} Q_{\nu} = \prod_{j} \tilde{D}_{\Lambda_{j}^{\nu}}^{\nu} \quad with \quad \sharp \Lambda_{j}^{\nu} \leq M_{\nu}, \quad \text{diam} \Lambda_{j}^{\nu} \leq M_{\nu} N_{\nu}, \quad \forall j.^2$$

Moreover, there exists a full-measure subset $\tilde{\mathcal{X}} \subset \mathbb{R}/\mathbb{Z}$ such that if we fix $x \in \tilde{\mathcal{X}}$, then for each $k \in \mathbb{Z}$, there is a $\nu_0(k)$ such that

$$\Lambda^{\nu+1}(k) = \Lambda^{\nu}(k), \quad \forall \nu \ge \nu_0(k).$$

Proposition 1 shows the pure point spectrum of T. In Appendix A.1, we shall give an outline of the proof.

2.2 Decay property of matrices

Lemma 2.1 Given two matrices $G = (G_{mn})_{m,n \in \mathbb{Z}}$ and $F = (F_{mn})_{m,n \in \mathbb{Z}}$. Let K = GF.

(1) If $|G_{mn}| \le c_G e^{-\sigma_G |m-n|}$, $|F_{mn}| \le c_F e^{-\sigma_F |m-n|}$ for some positive c_G , c_F , σ_G , $\sigma_F > 0$, then we have

$$|K_{mn}| < c_K e^{-\sigma_K |m-n|}$$

for any $0 < \sigma_K < \min\{\sigma_G, \sigma_F\}$ and $c_K = c \cdot c_G c_F (\min\{\sigma_G, \sigma_F\} - \sigma_K)^{-1}$.

(2) If $|G_{mn}| \le c_G e^{-\sigma_G \max\{|m|,|n|\}}$, $|F_{mn}| \le c_F e^{-\sigma_F |m-n|}$, then

$$|K_{mn}| \le c_K e^{-\sigma_K \max\{|m|, |n|\}}.$$

(3) If $|G_{mn}| \le c_G e^{-\sigma_G |m-n|}$, $|F_{mn}| \le c_F e^{-\sigma_F \max\{|m|, |n|\}}$, then

$$|K_{mn}| \le c_K e^{-\sigma_K \max\{|m|,|n|\}}.$$

(4) If $|G_{mn}| \le c_G e^{-\sigma_G \max\{|m|, |n|\}}$, $|F_{mn}| \le c_F e^{-\sigma_F \max\{|m|, |n|\}}$, then

$$|K_{mn}| < c_K e^{-\sigma_K \max\{|m|,|n|\}}$$
.

In particular, if $\sigma_G \neq \sigma_F$, then the conclusions above hold with $\sigma_K = \min\{\sigma_G, \sigma_F\}$ and $c_K = c \cdot c_G c_F |\sigma_G - \sigma_F|^{-1}$.

Proof: Since the matrix element of K = GF can be formulated as $K_{mn} = \sum_{l \in \mathbb{Z}} G_{ml} F_{ln}$, we have that, in Case (1), for any $0 < \sigma_K < \min\{\sigma_G, \sigma_F\}$,

$$\begin{aligned} |(GF)_{mn}| &\leq \sum_{l \in \mathbb{Z}} |G_{ml}| |F_{ln}| \\ &\leq c_G c_F \sum_{l \in \mathbb{Z}} e^{-\sigma_G |m-l|} e^{-\sigma_F |l-n|} \\ &\leq c_G c_F e^{-\sigma_K |m-n|} \sum_{l \in \mathbb{Z}} e^{-(\sigma_G - \sigma_K) |m-l|} e^{-(\sigma_F - \sigma_K) |l-n|} \\ &\leq c \cdot c_G c_F (\min\{\sigma_G, \sigma_F\} - \sigma_K)^{-1} e^{-\sigma_K |m-n|}. \end{aligned}$$

The disjoint decomposition defines an equivalence relation $m \sim n$ on the integers and, for each $n \in \mathbb{Z}$, we denote its equivalence class by $\Lambda^{\nu}(n)$.

Here we have applied the basic triangular inequality $|m-l|+|l-n| \ge |m-n|$. Moreover, if $\sigma_G \ne \sigma_F$, assume that $0 < \sigma_G < \sigma_F$ without loss of generality, then

$$|(GF)_{mn}| \leq c_G c_F \sum_{l \in \mathbb{Z}} e^{-\sigma_G |m-l|} e^{-\sigma_F |l-n|}$$

$$\leq c_G c_F e^{-\sigma_G |m-n|} \sum_{l \in \mathbb{Z}} e^{-(\sigma_F - \sigma_G)|l-n|}$$

$$\leq c \cdot c_G c_F (\sigma_F - \sigma_G)^{-1} e^{-\sigma_G |m-n|}.$$

As for Case (2)-(4), the corresponding conclusions can also be obtained by using the trivial facts

 $|m-l| + \max\{|l|, |n|\} \ge \max\{|m|, |n|\}, \quad \max\{|m|, |l|\} + \max\{|l|, |n|\} \ge \max\{|m|, |n|\}.$

Thus Lemma 2.1 has been proved.

Remark 2.1 If we replace the matrix F satisfying $|F_{mn}| \leq c_F e^{-\sigma_F \max\{|m|,|n|\}}$ with a vector $f = (f_n)_{n \in \mathbb{Z}}$ satisfying $|f_n| \leq c_f e^{-\sigma_f |n|}$ in Case (3) and (4), then for the vector Gf, we can obtain the conclusion that $|(Gf)_n| \leq c_K e^{-\sigma_K |n|}$, where c_K and σ_K are the same as in the lemma.

3 Abstract KAM theorem

3.1 Function space norms

Given $d, \rho > 0$, let $\ell^1_{d,\rho}(\mathbb{Z})$ be the space of summable complex valued sequences $q = (q_n)_{n \in \mathbb{Z}}$, with the norm

$$||q||_{d,\rho} := \sum_{n \in \mathbb{Z}} |q_n| \langle n \rangle^d e^{\rho|n|} < \infty,$$

where $\langle n \rangle := \sqrt{1+n^2}$. For r, s > 0, let $\mathcal{D}_{d,\rho}(r,s)$ be the complex b-dimensional neighborhood of $\mathbb{T}^b \times \{0\} \times \{0\}$ in $\mathbb{T}^b \times \mathbb{R}^b \times \ell^1_{d,\rho}(\mathbb{Z}) \times \ell^1_{d,\rho}(\mathbb{Z})$, i.e.,

$$\mathcal{D}_{d,\rho}(r,s) := \{ (\theta, I, q, \bar{q}) : |\text{Im}\theta| = |\text{Im}(\theta_1, \dots, \theta_b)| < r, |I| < s^2, ||q||_{d,\rho} = ||\bar{q}||_{d,\rho} < s \},$$

where $|\cdot|$ is the ℓ^1 -norm of complex vectors.

Given a real-analytic function $F(\theta, I, q, \bar{q}; \xi)$ on $\mathcal{D} = \mathcal{D}_{d,\rho}(r, s)$, C_W^1 (i.e., C^1 in the sense of Whitney) parametrized by $\xi \in \mathcal{O}^3$, a closed region in \mathbb{R}^b . We expand F into the Taylor-Fourier series with respect to θ, I, q, \bar{q} :

$$F(\theta, I, q, \bar{q}; \xi) = \sum_{\alpha, \beta} F_{\alpha\beta}(\theta, I; \xi) q^{\alpha} \bar{q}^{\beta},$$

where, for multi-indices $\alpha := \sum_{n \in \mathbb{Z}} \alpha_n e_n$, $\beta := \sum_{n \in \mathbb{Z}} \beta_n e_n$, $\alpha_n, \beta_n \in \mathbb{N}$, with finitely many non-vanishing components,

$$F_{\alpha\beta}(\theta,I;\xi) = \sum_{k \in \mathbb{Z}^b, \, l \in \mathbb{N}^b} F_{kl\alpha\beta}(\xi) I^l e^{\mathrm{i}\langle k,\theta \rangle}, \quad q^\alpha \bar{q}^\beta = \prod_{(\alpha_n,\beta_n) \neq (0,0)} q_n^{\alpha_n} \bar{q}_n^{\beta_n}.$$

³In this paper, all dependencies on the parameter $\xi \in \mathcal{O}$ are assumed of class C_W^1 , thus all derivatives with respective to ξ will be interpreted in this sense.

(Here e_m denotes the vector with the m^{th} component being 1 and the other components being zero.)

Definition 3.1 For each non-zero multi-index $(\alpha, \beta) = (\cdots, \alpha_n, \beta_n, \cdots)_{n \in \mathbb{Z}}, \alpha_n, \beta_n \in \mathbb{N}$, with finitely many non-vanishing components, we define

$$n_{\alpha\beta}^{+} := \max\{n \in \mathbb{Z} : (\alpha_{n}, \beta_{n}) \neq (0, 0)\},\$$

$$n_{\alpha\beta}^{-} := \min\{n \in \mathbb{Z} : (\alpha_{n}, \beta_{n}) \neq (0, 0)\},\$$

$$n_{\alpha\beta}^{*} := \max\{|n_{\alpha\beta}^{+}|, |n_{\alpha\beta}^{-}|\},\$$

and $|\alpha| := \sum_n \alpha_n$, $|\beta| := \sum_n \beta_n$. In particular, for $|\alpha| = |\beta| = 0$, we define $n_{\alpha\beta}^+ = n_{\alpha\beta}^- = n_{\alpha\beta}^* = 0$.

In particular, for $|\alpha| = |\beta| = 0$, we define $n_{\alpha\beta} = n_{\alpha\beta} = n_{\alpha\beta} = 0$.

With
$$|\partial_{\xi} F_{kl\alpha\beta}| := \sum_{i=1}^{b} |\partial_{\xi_i} F_{kl\alpha\beta}|$$
 and $|F_{kl\alpha\beta}|_{\mathcal{O}} := \sup_{\xi \in \mathcal{O}} (|F_{kl\alpha\beta}| + |\partial_{\xi} F_{kl\alpha\beta}|)$, let

$$||F_{\alpha\beta}||_{\mathcal{O}} := \sum_{k,l} |F_{kl\alpha\beta}|_{\mathcal{O}} |I^l| e^{|k||\mathrm{Im}\theta|}, \quad ||F||_{\mathcal{O}} := \sum_{k,l,\alpha,\beta} |F_{kl\alpha\beta}|_{\mathcal{O}} |I^l| e^{|k||\mathrm{Im}\theta|} |q^{\alpha}||\bar{q}^{\beta}|.$$

Define the weighted norm of F as

$$||F||_{\mathcal{D},\mathcal{O}} := \sup_{\mathcal{D}} ||F||_{\mathcal{O}}.^4$$

For the Hamiltonian vector field $X_F = (\partial_I F, -\partial_{\theta} F, (-i\partial_{q_n} F)_{n \in \mathbb{Z}}, (i\partial_{\bar{q}_n} F)_{n \in \mathbb{Z}})$ associated F on $\mathcal{D} \times \mathcal{O}$, its norm is defined by

$$||X_F||_{\mathcal{D},\mathcal{O}} := ||\partial_I F||_{\mathcal{D},\mathcal{O}} + \frac{1}{s^2} ||\partial_\theta F||_{\mathcal{D},\mathcal{O}} + \sup_{\mathcal{D}} \frac{1}{s} \sum_{n \in \mathbb{Z}} (||\partial_{q_n} F||_{\mathcal{O}} + ||\partial_{\bar{q}_n} F||_{\mathcal{O}}) \langle n \rangle^d e^{|n|\rho}.$$

Sometimes, for the sake of notational simplification, we shall not write the subscript \mathcal{O} in the norms defined above if it is obvious enough.

Given two real-analytic functions F and G on \mathcal{D} , let $\{\cdot,\cdot\}$ denote the Poisson bracket of such functions, i.e.,

$$\{F,G\} = \langle \partial_I F, \partial_\theta G \rangle - \langle \partial_\theta F, \partial_I G \rangle + \mathrm{i} \sum_{n \in \mathbb{Z}} \left(\partial_{q_n} F \cdot \partial_{\bar{q}_n} G - \partial_{\bar{q}_n} F \cdot \partial_{q_n} G \right).$$

Some basic estimates about the Hamiltonian vector field and the Poisson bracket are given in Appendix A.2.

3.2 Statement of the abstract KAM theorem

Now, we consider the perturbed Hamiltonian

$$H = \mathcal{N} + \breve{P} + P$$

$$= e(x,\xi) + \langle \omega(x,\xi), I \rangle + \langle \Omega(x,\xi)q, \bar{q} \rangle + \breve{P}(q,\bar{q};x) + P(\theta,I,q,\bar{q};x,\xi), \quad (3.1)$$

defined on the domain $\mathcal{D} = \mathcal{D}_{d,\rho}(r,s)$. Our goal is to prove that, for a.e. $x \in \mathbb{R}/\mathbb{Z}$, the Hamiltonian H admits invariant tori for "most" of the parameter $\xi \in \mathcal{O} = \mathcal{O}(x)$, provided that $\|X_{\check{P}+P}\|_{\mathcal{D},\mathcal{O}}$ is sufficiently small.

⁴In the case of a vector-valued function $F: \mathcal{D} \times \mathcal{O} \to \mathbb{C}^n (n < \infty)$, the norm is defined as $||F||_{\mathcal{D}, \mathcal{O}} := \sum_{i=1}^n ||F_i||_{\mathcal{D}, \mathcal{O}}$.

- **Remark 3.1** From now on, we shall not report x for convenience if it is irrelevant. To this end, we need to impose some conditions on ω , Ω , and the perturbations $\check{P} + P$.
- (A1) Nondegeneracy of tangential frequencies: The map $\xi \to \omega$ is a C_W^1 diffeomorphism between \mathcal{O} and its image.
- **(A2)** Regularity of Ω : $\Omega = T + A + W$.
 - T is the symmetric matrix defined in (2.1), independent of ξ . More precisely,

$$T = \operatorname{diag}\{V(x + n\tilde{\alpha})\}_{n \in \mathbb{Z}} + \epsilon \Delta,$$

with V and $\tilde{\alpha}$ as in Equation (1.9).

- A is Hermitian, independent of ξ , satisfying

$$|A_{mn}| \le \begin{cases} c, & |m|, |n| \le \hat{N} \\ 0, & \text{otherwise} \end{cases}$$
 (3.2)

for some positive \hat{N} .

- W is C_W^1 parametrized by $\xi \in \mathcal{O}$, with

$$|W_{mn}|_{\mathcal{O}} \le \begin{cases} pe^{-\sigma \max\{|m|,|n|\}}, & |m|,|n| \le N \\ 0, & \text{otherwise} \end{cases}$$
 (3.3)

for some positive $p \ll 1$, $\sigma \gg \rho$ and sufficiently large N.

Moreover, there exists a subset $\mathcal{J} \subset \mathbb{Z}$ such that

$$\Omega_{mn} \equiv 0 \quad \text{if} \quad m \quad \text{or} \quad n \in \mathcal{J}.$$
(3.4)

(A3) Short range of \check{P} : $\check{P}(q,\bar{q}) = \sum_{|\alpha|=|\beta|\geq 2} \check{P}_{\alpha\beta} q^{\alpha} \bar{q}^{\beta}$ is real-analytic in q,\bar{q} , and independent of ξ , with

$$|\check{P}_{\alpha\beta}| \le e^{-\rho(n_{\alpha\beta}^+ - n_{\alpha\beta}^-)}, \quad |\alpha| = |\beta| \ge 2,$$
 (3.5)

$$\partial_{q_n} \check{P} = \partial_{\bar{q}_n} \check{P} \equiv 0, \quad \forall n \in \mathcal{J}.$$
 (3.6)

(A4) Decay property of $P: P = \sum_{\alpha,\beta} P_{\alpha\beta}(\theta,I;\xi) q^{\alpha} \bar{q}^{\beta}$ is real-analytic in $\theta,I,q,\bar{q}, C_W^1$ parametrized by $\xi \in \mathcal{O}$, and, with $\varepsilon = \epsilon^{\frac{1}{4}}$,

$$||P_{\alpha\beta}||_{\mathcal{D},\mathcal{O}} \le \begin{cases} \varepsilon e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \le 2\\ e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \ge 3 \end{cases}, \tag{3.7}$$

$$\partial_{q_n} P = \partial_{\bar{q}_n} P \equiv 0, \quad \forall n \in \mathcal{J}.$$
 (3.8)

(A5) Gauge invariance of P: For $P = \sum_{k \in \mathbb{Z}^b, l \in \mathbb{N}^b \atop \alpha, \beta} P_{kl\alpha\beta} I^l e^{\mathrm{i}\langle k, \theta \rangle} q^{\alpha} \bar{q}^{\beta}$, we have

$$P_{kl\alpha\beta} \equiv 0$$
 if $\sum_{j=1}^{b} k_j + |\alpha| - |\beta| \neq 0$.

Our abstract KAM theorem can be stated as follows.

Theorem 2 Consider the Hamiltonian H in (3.1), with $(\mathbf{A1}) - (\mathbf{A5})$ satisfied. There is a positive constant $\varepsilon_* = \varepsilon_*(\omega, V, \tilde{\alpha}, \hat{N}, p, \sigma, N, r, s, d, \rho)$ such that if $\|X_{\tilde{P}+P}\|_{\mathcal{D}, \mathcal{O}} \leq \varepsilon \leq \varepsilon_*$, then for every $x \in \tilde{\mathcal{X}}$, there exists a Cantor set $\mathcal{O}_{\varepsilon} = \mathcal{O}_{\varepsilon}(x) \subset \mathcal{O}(x)$ with $|\mathcal{O} \setminus \mathcal{O}_{\varepsilon}| \to 0$ as $\varepsilon \to 0$, such that the following holds.

- (a) There exists a C_W^1 map $\tilde{\omega}: \mathcal{O}_{\varepsilon} \to \mathbb{R}^b$, such that $|\tilde{\omega} \omega|_{\mathcal{O}_{\varepsilon}} \to 0$ as $\varepsilon \to 0$.
- (b) There exists a map $\Psi : \mathbb{T}^b \times \mathcal{O}_{\varepsilon} \to \mathcal{D}_{d,0}(r/2,0)$, real-analytic in $\theta \in \mathbb{T}^b$ and C_W^1 parametrized by $\xi \in \mathcal{O}$, such that $\|\Psi \Psi_0\|_{\mathcal{D}_{d,0}(r/2,0),\mathcal{O}_{\varepsilon}} \to 0$ as $\varepsilon \to 0$, where Ψ_0 is the trivial embedding: $\mathbb{T}^b \times \mathcal{O} \to \mathbb{T}^b \times \{0\} \times \{0\}$.
- (c) For any $\theta \in \mathbb{T}^b$ and $\xi \in \mathcal{O}_{\varepsilon}$, $\Psi(\theta + \tilde{\omega}(\xi)t, \xi) = (\theta + \tilde{\omega}(\xi)t, I(t), q(t), \bar{q}(t))$ is a b-frequency quasi-periodic solution of equations of motion associated with the Hamiltonian (3.1).
- (d) For each t, $q(t) = (q_n(t))_{n \in \mathbb{Z}} \in \ell^1_{d,0}(\mathbb{Z})$.

Remark 3.2 The statement (d) of Theorem 2 implies that

$$\sup_{t} \sum_{n \in \mathbb{Z}} n^{2d} |q_n(t)|^2 < c \left(\sup_{t} \sum_{n \in \mathbb{Z}} \langle n \rangle^d |q_n(t)| \right)^2 < \infty,$$

which is exactly the conclusion of Theorem 1.

Remark 3.3 In case that H satisfies $(\mathbf{A1}) - (\mathbf{A5})$ at the first step, all assumptions hold for the Hamiltonian at each KAM step (with suitable parameters).

3.3 Application to Equation (1.9)

The Hamiltonian associated with Equation (1.9) is

$$H = \sum_{n \in \mathbb{Z}} V(x + n\tilde{\alpha}) q_n \bar{q}_n + \epsilon \sum_{n \in \mathbb{Z}} \bar{q}_n (q_{n+1} + q_{n-1}) + \frac{1}{2} \sum_{n \in \mathbb{Z}} |q_n|^4.$$
 (3.9)

Fix $\mathcal{J} = \{n_1, \dots, n_b\} \subset \mathbb{Z}$, and $\mathbb{Z}_1 = \mathbb{Z} \setminus \mathcal{J}$. Let $\varepsilon = \epsilon^{\frac{1}{4}}$, with ϵ sufficiently small such that

$$|n_i| \le |\ln \varepsilon| = \frac{1}{4} |\ln \epsilon|, \quad i = 1, \dots, b.$$

We introduce action-angle variables and amplitude parameters to the Hamiltonian (3.9),

$$q_n = \sqrt{I_n + \xi_n} e^{\mathrm{i}\theta_n}, \quad \bar{q}_n = \sqrt{I_n + \xi_n} e^{-\mathrm{i}\theta_n}, \quad n \in \mathcal{J},$$

where $(I,\theta)=(I_{n_1},\cdots,I_{n_b},\theta_{n_1},\cdots,\theta_{n_b})$ are the standard action-angle variables in the $(q_n,\bar{q}_n)_{n\in\mathcal{J}}$ -space around ξ , with $\xi=(\xi_{n_1},\cdots,\xi_{n_b})\in\mathcal{O}=[\epsilon^{\frac{1}{12}},1]\subset[0,1]^b$ a parameter,

and $(q, \bar{q}) = (q_n, \bar{q}_n)_{n \in \mathbb{Z}_1}$. Then the Hamiltonian (3.9) becomes $H = \mathcal{N}(\theta, I, q, \bar{q}; x, \xi) + \check{P}(q, \bar{q}) + P(\theta, I, q, \bar{q}; \xi)$, with

$$\mathcal{N}(\theta, I, q, \bar{q}; x, \xi) := \sum_{n \in \mathcal{J}} (V(x + n\tilde{\alpha})\xi_n + \frac{1}{2}\xi_n^2) + \sum_{n \in \mathcal{J}} (V(x + n\tilde{\alpha}) + \xi_n)I_n
+ \sum_{n \in \mathbb{Z}_1} V(x + n\tilde{\alpha})|q_n|^2 + \epsilon \sum_{\substack{n \notin \mathcal{J} \\ n+1 \notin \mathcal{J},}} (\bar{q}_n q_{n+1} + q_n \bar{q}_{n+1}),
\check{P}(q, \bar{q}) := \frac{1}{2} \sum_{n \in \mathbb{Z}_1} |q_n|^4,
P(\theta, I, q, \bar{q}; \xi) := \frac{1}{2} \sum_{n \in \mathcal{J}} I_n^2 + \epsilon \sum_{\substack{m \in \mathcal{J}, n \notin \mathcal{J} \\ |m-n|=1}} \sqrt{I_m + \xi_m} (e^{-i\theta_m} q_n + e^{i\theta_m} \bar{q}_n)
+ \epsilon \sum_{\substack{m, n \in \mathcal{J} \\ |m-n|=1}} \sqrt{I_m + \xi_m} \sqrt{I_n + \xi_n} (e^{-i(\theta_m - \theta_n)} + e^{i(\theta_m - \theta_n)}).$$

After introducing the action-angle variables, we find that the structure of the linear operator T in (3.9) has been destroyed. To overcome this disadvantage, we need to add b variables $q'_{n_1}, \dots, q'_{n_b}$ and the corresponding conjugates $\bar{q}'_{n_1}, \dots, \bar{q}'_{n_b}$ into this system. For convenience, omit the prime of the newly-added variables and still use q to denote $(q_n)_{n \in \mathbb{Z}}$, since there is no confusion. We then rewrite \mathcal{N} as

$$\mathcal{N} = \sum_{n \in \mathcal{J}} (V(x + n\tilde{\alpha})\xi_n + \frac{1}{2}\xi_n^2) + \sum_{n \in \mathcal{J}} (V(x + n\tilde{\alpha}) + \xi_n)I_n
+ \left[\sum_{n \in \mathbb{Z}_1} V(x + n\tilde{\alpha})|q_n|^2 + \sum_{n \in \mathcal{J}} V(x + n\tilde{\alpha})|q_n|^2 + \epsilon \sum_{n \in \mathbb{Z}} (\bar{q}_n q_{n+1} + q_n \bar{q}_{n+1}) \right]
- \sum_{n \in \mathcal{J}} V(x + n\tilde{\alpha})|q_n|^2 - \epsilon \sum_{n \text{ or } n+1 \in \mathcal{J}} (\bar{q}_n q_{n+1} + q_n \bar{q}_{n+1})
= e(x, \xi) + \langle \omega(x, \xi), I \rangle + \langle T(x)q, \bar{q} \rangle + \langle A(x)q, \bar{q} \rangle,
\text{with } e(x, \xi) := \sum_{n \in \mathcal{J}} V(x + n\tilde{\alpha})\xi_n + \frac{1}{2} \sum_{n \in \mathcal{J}} \xi_n^2,
\omega(x, \xi) := (V(x + n_1\tilde{\alpha}) + \xi_{n_1}, \dots, V(x + n_b\tilde{\alpha}) + \xi_{n_b}),$$

$$T_{mn}(x) := \begin{cases} V(x + m\tilde{\alpha}), & m = n \\ \epsilon, & m - n = \pm 1 \\ 0, & \text{otherwise} \end{cases}$$

$$A_{mn}(x) := \begin{cases} -V(x + m\tilde{\alpha}), & m = n, m \in \mathcal{J} \\ -\epsilon, & m - n = \pm 1, m \text{ or } n \in \mathcal{J} \\ 0, & \text{otherwise} \end{cases}$$

$$(3.12)$$

Now, on some domain $\mathcal{D}_{d,\rho}(r,s)$, the regularity of $\check{P}+P$ holds true:

Lemma 3.1 For $\epsilon > 0$ sufficiently small and $s = \frac{1}{8} \epsilon^{\frac{1}{4}}$, if $|I| < s^2$ and $||q||_{d,\rho} < s$, then

$$||X_{\breve{P}+P}||_{\mathcal{D}_{d,o}(r,s),\mathcal{O}} \le \epsilon^{\frac{1}{4}} = \varepsilon.$$

We need to show that the Hamiltonian $H = \mathcal{N} + P + P$ satisfies the assumptions (A1) - (A5) of the KAM theorem, in which (A3) and (A5) are obviously satisfied.

(A1): Since $\{V(x+n\tilde{\alpha})\}_{n\in\mathbb{Z}}$ is independent of ξ , we have that $\frac{\partial\omega}{\partial\xi}\equiv I_{\mathcal{J}}$ in view of (3.10). Thus $(\mathbf{A1})$ holds.

(A2): Here, $W \equiv 0$. Then it is evident that (A2) holds with $\hat{N} = \frac{1}{4} |\ln \epsilon|$, by (3.12).

(A4): We focus on the formulation of P. Note that terms of P merely correspond to the normal variables q_n , \bar{q}_n , $n \notin \mathcal{J}$, n-1 or $n+1 \in \mathcal{J}$, with the coefficients no more than ϵ , and $\mathcal{J} \subset [-\hat{N}, \hat{N}] = [-\frac{1}{4}|\ln \epsilon|, \frac{1}{4}|\ln \epsilon|]$. Then, with $\rho \leq \frac{1}{6}\hat{N}^{-1}$, (3.7) is verified since

$$c\epsilon^{1-\frac{1}{24}} \le \epsilon^{\frac{1}{4}}e^{-\rho\hat{N}}.$$

Hence, Theorem 1 is a corollary of Theorem 2.

4 KAM step

To start the KAM iteration for the Hamiltonian (3.1), let $\mathcal{D}_0 = \mathcal{D}_{d,\rho_0}(r_0,s_0)$, \mathcal{O}_0 , \mathcal{N}_0 (including $e_0, \omega_0, W_0, p_0, \sigma_0, N_0$), $P_0, \varepsilon_0 = \epsilon^{\frac{1}{4}}$ denote the initial quantities given in the assumptions $(\mathbf{A1}) - (\mathbf{A5})$ respectively, and require that ϵ smaller than the ϵ_0 given in Proposition 1.

Suppose we have arrived at the ν^{th} step of the KAM iteration, $\nu = 0, 1, 2, \cdots$, recalling that several sequences have been given in (2.2). We consider the Hamiltonian on \mathcal{D}_{ν} := $\mathcal{D}_{d,\rho_{\nu}}(r_{\nu},s_{\nu})$ and \mathcal{O}_{ν} ,

$$H_{\nu} = \mathcal{N}_{\nu} + \breve{P} + P_{\nu}$$

= $e_{\nu} + \langle \omega_{\nu}, I \rangle + \langle \Omega_{\nu} q, \bar{q} \rangle + \breve{P} + P_{\nu},$ (4.1)

where $\Omega_{\nu} = T + A + W_{\nu}$, and $(\mathbf{A1}) - (\mathbf{A5})$ are satisfied, including (3.2), (3.5), (3.6) and

$$(\Omega_{\nu})_{mn} \equiv 0, \quad m \text{ or } n \in \mathcal{J},$$
 (4.2)

$$(\Omega_{\nu})_{mn} \equiv 0, \quad m \text{ or } n \in \mathcal{J},$$

$$|(W_{\nu})_{mn}|_{\mathcal{O}_{\nu}} \leq \begin{cases} p_{\nu}e^{-\sigma_{\nu}\max\{|m|,|n|\}}, & |m|,|n| \leq N_{\nu} \\ 0, & \text{otherwise} \end{cases},$$

$$|(P_{\nu})_{\alpha\beta}|_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \leq \begin{cases} \varepsilon_{\nu}e^{-\rho_{\nu}n_{\alpha\beta}^{*}}, & |\alpha|+|\beta| \leq 2 \\ e^{-\rho_{\nu}n_{\alpha\beta}^{*}}, & |\alpha|+|\beta| \geq 3 \end{cases},$$

$$(4.2)$$

$$\|(P_{\nu})_{\alpha\beta}\|_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \le \begin{cases} \varepsilon_{\nu} e^{-\rho_{\nu} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \le 2\\ e^{-\rho_{\nu} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \ge 3 \end{cases}, \tag{4.4}$$

$$\partial_{q_n} P_{\nu} = \partial_{\bar{q}_n} P_{\nu} \equiv 0, \quad n \in \mathcal{J}.$$
 (4.5)

Moreover, $||X_{\check{P}+P_{\nu}}||_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \leq \varepsilon_{\nu}$.

Choose some $r_{\nu+1}$ such that $0 < r_{\nu+1} < r_{\nu}$, and let $J_{\nu} := \left\lceil \frac{5}{2} \varepsilon_{\nu}^{-\frac{a_{\nu}}{2}} \right\rceil$. For $j = 0, 1, \dots, J_{\nu}$, we define the quantities at each KAM sub-step as

$$\rho_{\nu}^{(j)} = (1 - \frac{j}{2J_{\nu}})\rho_{\nu}, \quad r_{\nu}^{(j)} = r_{\nu} - \frac{j(r_{\nu} - r_{\nu+1})}{J_{\nu}}, \quad s_{\nu}^{(j)} = 2^{-3j}\varepsilon_{\nu}^{\frac{j}{5}}s_{\nu},$$

and $\mathcal{D}_{\nu}^{(j)} = \mathcal{D}_{d,\rho_{\nu+1}}(r_{\nu}^{(j)}, s_{\nu}^{(j)}), \ \varepsilon_{\nu}^{(j)} = \varepsilon_{\nu}^{\frac{j}{5}+1}$. Our goal is to construct a set $\mathcal{O}_{\nu+1} \subset \mathcal{O}_{\nu}$ and a finite sequence of maps

$$\Phi_{\nu}^{(j)}: \mathcal{D}_{\nu}^{(j)} \to \mathcal{D}_{\nu}^{(j-1)}, \quad j = 1, 2, \cdots, J_{\nu},$$

so that the Hamiltonian transformed into the $(\nu+1)^{\rm th}$ KAM cycle

$$H_{\nu+1} = H_{\nu} \circ \Phi_{\nu}^{(1)} \circ \cdots \circ \Phi_{\nu}^{(J_{\nu})}$$

$$= \mathcal{N}_{\nu+1} + \check{P} + P_{\nu+1}$$

$$= e_{\nu+1} + \langle \omega_{\nu+1}, I \rangle + \langle \Omega_{\nu+1}q, \bar{q} \rangle + \check{P} + P_{\nu+1}$$

satisfies all the above iterative assumptions $(\mathbf{A1}) - (\mathbf{A5})$ on $\mathcal{D}_{\nu+1} = \mathcal{D}_{\nu}^{(J_{\nu})}$ and C_W^1 parametrized by $\xi \in \mathcal{O}_{\nu+1}$, with new suitable parameters. Moreover,

$$\|X_{P+P_{\nu+1}}\|_{\mathcal{D}_{\nu+1},\mathcal{O}_{\nu+1}} \le \varepsilon_{\nu}^{(J_{\nu})} \le \varepsilon_{\nu}^{\frac{1}{2}\varepsilon_{\nu}^{-a_{\nu}/2}} = \varepsilon_{\nu+1}.$$

In the remaining part of this paper, all constants labeled with c, c_0, c_1, \cdots are positive and independent of the iteration step.

4.1 Construction of $\mathcal{O}_{\nu+1}$

As described in Proposition 1, there exists an orthogonal matrix U_{ν} with

$$|(U_{\nu} - I_{\mathbb{Z}})_{mn}| \le \varepsilon_0^{\frac{1}{2}} e^{-\frac{3}{2}\sigma_{\nu}|m-n|},$$
 (4.6)

such that $U_{\nu}^*TU_{\nu} = D_{\nu} + Z_{\nu}$, where Z_{ν} is a symmetric matrix satisfying

$$|(Z_{\nu})_{mn}| \le \varepsilon_{\nu} e^{-\rho_{\nu}|m-n|},\tag{4.7}$$

and D_{ν} is a symmetric matrix which can be block-diagonalized via an orthogonal matrix Q_{ν} with

$$(Q_{\nu})_{mn} = 0 \text{ if } |m-n| > N_{\nu}.$$
 (4.8)

More precisely,

$$\tilde{D}_{\nu} = Q_{\nu}^* D_{\nu} Q_{\nu} = \prod_{j} \tilde{D}_{\Lambda_{j}^{\nu}}^{\nu} \text{ with } \sharp \Lambda_{j}^{\nu} \le M_{\nu}, \text{ diam} \Lambda_{j}^{\nu} \le M_{\nu} N_{\nu}, \quad \forall j.$$

To describe $U_{\nu}^*\Omega_{\nu}U_{\nu}$, we need furthermore to consider $U_{\nu}^*AU_{\nu}$ and $U_{\nu}^*W_{\nu}U_{\nu}$. In view of (3.2), (4.3) and (4.6), there exists a constant $c_1 > 0$ such that

$$|(U_{\nu}^*(A+W_{\nu})U_{\nu})_{mn}|_{\mathcal{O}_{\nu}} \le c_1 \max\{\hat{N}^2 e^{3\sigma_{\nu}\hat{N}}, p_{\nu}\sigma_{\nu}^{-2}\} \cdot e^{-\sigma_{\nu}\cdot \max\{|m|, |n|\}},$$

by a simple application of Lemma 2.1. Define the truncation \hat{A}_{ν} as

$$(\hat{A}_{\nu})_{mn} := \begin{cases} (U_{\nu}^{*}(A + W_{\nu})U_{\nu})_{mn}, & |m|, |n| \leq N_{\nu} \\ 0, & \text{otherwise} \end{cases}$$
 (4.9)

It follows that

$$\left| \left(U_{\nu}^* (A + W_{\nu}) U_{\nu} - \hat{A}_{\nu} \right)_{mn} \right|_{\mathcal{O}_{\nu}} \le \varepsilon_{\nu} e^{-\rho_{\nu} \max\{|m|, |n|\}}$$

$$\tag{4.10}$$

under the assumption

(C1) $c_1 \max\{\hat{N}^2 e^{3\sigma_{\nu}\hat{N}}, p_{\nu}\sigma_{\nu}^{-2}\} \cdot e^{-(\sigma_{\nu}-\rho_{\nu})N_{\nu}} \le \varepsilon_{\nu}.$

Let $K_{\nu+1} := N_{\nu+1} - (M_{\nu} + 1)N_{\nu}$ with the sequences M_{ν} , N_{ν} , $\nu = 0, 1, \dots$, defined in (2.2) and

$$\tilde{D}^{\nu}_{\Lambda^{\nu}} := \prod_{\Lambda^{\nu}_{j} \subset \Lambda^{\nu}} \tilde{D}^{\nu}_{\Lambda^{\nu}_{j}}, \quad \tilde{A}_{\nu} := Q^{*}_{\nu} \hat{A}_{\nu} Q_{\nu}, \tag{4.11}$$

where $\Lambda^{\nu} := \bigcup \{\Lambda_{j}^{\nu} : \Lambda_{j}^{\nu} \cap [-(K_{\nu+1} + N_{\nu}), K_{\nu+1} + N_{\nu}] \neq \emptyset\} \subset [-N_{\nu+1}, N_{\nu+1}].$ In view of (4.8) and (4.9), we have

$$(\tilde{A}_{\nu})_{mn} \equiv 0$$
 if $|m|$ or $|n| > 2N_{\nu}$.

Since both of $\tilde{D}^{\nu}_{\Lambda^{\nu}}$ and \tilde{A}_{ν} are Hermitian, there is an orthogonal matrix O_{ν} such that

$$O_{\nu}^*(\tilde{D}_{\Lambda^{\nu}}^{\nu} + \tilde{A}_{\nu})O_{\nu} = \operatorname{diag}\{\mu_{j}^{\nu}\}_{j \in \Lambda^{\nu}},$$

where $\{\mu_j^{\nu}\}_{j\in\Lambda^{\nu}}$ are eigenvalues of $\tilde{D}_{\Lambda^{\nu}}^{\nu} + \tilde{A}_{\nu}$. Due to the block-diagonal structure of $\tilde{D}_{\Lambda^{\nu}}^{\nu} + \tilde{A}_{\nu}$, we also have

$$(O_{\nu})_{mn} \equiv 0 \text{ if } |m-n| > 2(M_{\nu}+2)N_{\nu}.$$
 (4.12)

Indeed, $\tilde{D}^{\nu}_{\Lambda^{\nu}} + \tilde{A}_{\nu}$ can be expressed as

$$\tilde{D}^{\nu}_{\Lambda^{\nu}} + \tilde{A}_{\nu} = (\tilde{D}^{\nu}_{\Lambda^{\nu}_{\nu}} + \tilde{A}_{\nu}) \cdot \prod_{\Lambda^{\nu}_{j} \cap [-2N_{\nu}, 2N_{\nu}] = \emptyset} \tilde{D}^{\nu}_{\Lambda^{\nu}_{j}}$$

where $\Lambda'_{\nu} := \bigcup \{\Lambda^{\nu}_{j} : \Lambda^{\nu}_{j} \cap [-2N_{\nu}, 2N_{\nu}] \neq \emptyset, \ \Lambda^{\nu}_{j} \subset \Lambda^{\nu} \}$ with $\operatorname{diam} \Lambda'_{\nu} \leq 2(M_{\nu} + 2)N_{\nu}$. The diagonalization of $\tilde{D}^{\nu}_{\Lambda^{\nu}} + \tilde{A}_{\nu}$ is just the diagonalization of blocks $(\tilde{D}^{\nu}_{\Lambda^{\nu}_{j}} + \tilde{A}_{\nu})$ and $\tilde{D}^{\nu}_{\Lambda^{\nu}_{j}}$.

As for the eigenvalues of $\tilde{D}_{\Lambda\nu}^{\nu} + \tilde{A}_{\nu}$, it is well-known that $\{\mu_{n}^{\nu}\}_{n \in \Lambda^{\nu}} C_{W}^{1}$ -smoothly depend on ξ and there exist orthonormal eigenvectors ψ_{n}^{ν} corresponding to μ_{n}^{ν} , C_{W}^{1} -smoothly depending on ξ (see e.g. [10]). In fact, $\mu_{n}^{\nu} = \langle (\tilde{D}_{\Lambda^{\nu}} + \tilde{A}_{\nu})\psi_{n}^{\nu}, \bar{\psi}_{n}^{\nu} \rangle$ and

$$\partial_{\xi_j} \mu_n^{\nu} = \langle (\partial_{\xi_j} (\tilde{D}_{\Lambda^{\nu}} + \tilde{A}_{\nu})) \psi_n^{\nu}, \bar{\psi}_n^{\nu} \rangle, \quad j = 1, \cdots, b.$$

By the construction of \tilde{A}_{ν} , we have $\partial_{\xi_{j}}\tilde{A}_{\nu}=Q_{\nu}^{*}(\partial_{\xi_{j}}\hat{A}_{\nu})Q_{\nu}$, with \hat{A}_{ν} the truncation of $U_{\nu}^{*}(A+W_{\nu}(\xi))U_{\nu}$. Since D_{ν} , A, U_{ν} and Q_{ν} are all independent of ξ ,

$$\sup_{\xi \in \mathcal{O}_{\nu}} |\partial_{\xi} \mu_n^{\nu}| \le c \sup_{\xi \in \mathcal{O}_{\nu} \atop m,n} |\partial_{\xi} (W_{\nu})_{mn}| \le c p_{\nu}. \tag{4.13}$$

Now we defined the new parameter set $\mathcal{O}_{\nu+1} \subset \mathcal{O}_{\nu}$ as

$$\mathcal{O}_{\nu+1} := \left\{ \xi \in \mathcal{O}_{\nu} : \begin{array}{l} |\langle k, \omega_{\nu} \rangle| > \frac{\gamma_{\nu}}{|k|^{\tau}}, & k \neq 0, \\ |\langle k, \omega_{\nu} \rangle + \mu_{n}^{\nu}| > \frac{\gamma_{\nu}}{|k|^{\tau} N_{\nu+1}^{2}}, & k \neq 0, & n \in \Lambda^{\nu}, \\ |\langle k, \omega_{\nu} \rangle + \mu_{m}^{\nu} \pm \mu_{n}^{\nu}| > \frac{\gamma_{\nu}}{|k|^{\tau} N_{\nu+1}^{4}}, & k \neq 0, & m, n \in \Lambda^{\nu}. \end{array} \right\}$$
(4.14)

for some $0 < \gamma_{\nu} \ll 1$, $\tau \geq b$. These inequalities are famous small-divisor conditions for controlling the solutions of the linearized equations.

From now on, to simplify notations, the subscripts (or superscripts) " ν " of quantities at the ν^{th} step are neglected, and the corresponding quantities at the $(\nu + 1)^{\text{th}}$ step are labeled with "+". In addition, we still use the superscript (j) to distinguish quantities at various sub-steps.

4.2 Homological equation and its approximate solution

For $P = \sum_{k,l,\alpha,\beta} P_{kl\alpha\beta}(\xi) I^l e^{i\langle k,\theta \rangle} q^{\alpha} \bar{q}^{\beta}$, according to (4.4) and the definition of norm in subsection 3.1, we have

$$|P_{kl\alpha\beta}|_{\mathcal{O}} \le \varepsilon e^{-\rho n_{\alpha\beta}^*} e^{-|k|r}, \quad 2|l| + |\alpha| + |\beta| \le 2, \quad \forall k \in \mathbb{Z}^b.$$
 (4.15)

Decompose P = R + (P - R) with

$$R := \sum_{\substack{k \ 2|l|+|\alpha|+|\beta| \le 2}} P_{kl\alpha\beta} e^{\mathrm{i}\langle k,\theta\rangle} I^l q^\alpha \bar{q}^\beta, \quad P - R = \sum_{\substack{k \ 2|l|+|\alpha|+|\beta| \ge 3}} P_{kl\alpha\beta} e^{\mathrm{i}\langle k,\theta\rangle} I^l q^\alpha \bar{q}^\beta.$$

It follows $||X_R||_{\mathcal{D},\mathcal{O}} \leq ||X_P||_{\mathcal{D},\mathcal{O}} \leq \varepsilon$. Recalling that $\check{P}(q,\bar{q})$ is a sum of high-order terms, there exists a constant $c_2 > 0$ such that

$$\|X_{\breve{P}}\|_{\mathcal{D}_{d,\rho}(r,\eta s),\mathcal{O}}, \|X_{P-R}\|_{\mathcal{D}_{d,\rho}(r,\eta s),\mathcal{O}} \le c_2 \eta s \le \frac{1}{8} \varepsilon^{\frac{6}{5}},$$
 (4.16)

with $\eta := \varepsilon^{\frac{1}{5}}$, provided

(C2) $c_2 s \leq \frac{1}{8} \varepsilon$.

Let $e' := P_{0000}$ and $\omega' := \int \frac{\partial P}{\partial I} d\theta|_{q=\bar{q}=0,I=0}$. With \mathcal{O}_+ defined as in (4.14), we have

Proposition 2 There exist two real-analytic Hamiltonians

$$F = \sum_{\substack{k \neq 0 \\ 1 \leq 2|l| + |\alpha| + |\beta| \leq 2}} F_{kl\alpha\beta} q^{\alpha} \bar{q}^{\beta} I^l e^{i\langle k, \theta \rangle}, \quad \grave{P} = \sum_{\substack{k \\ 1 \leq |\alpha| + |\beta| \leq 2}} \grave{P}_{k0\alpha\beta} q^{\alpha} \bar{q}^{\beta} e^{i\langle k, \theta \rangle},$$

and a Hermitian matrix W', all of which are C_W^1 parametrized by $\xi \in \mathcal{O}_+$, such that

$$\{\mathcal{N}, F\} + R = e' + \langle \omega', I \rangle + \langle W'q, \bar{q} \rangle + \hat{P}. \tag{4.17}$$

Moreover, both of F and \dot{P} have gauge invariance, and for ε sufficiently small,

$$|F_{kl\alpha\beta}|_{\mathcal{O}_{+}} \le \varepsilon^{\frac{4}{5}} |k|^{2\tau+1} e^{-|k|r} e^{-\rho n_{\alpha\beta}^{*}},$$
(4.18)

$$|\dot{P}_{k0\alpha\beta}|_{\mathcal{O}_{+}} \le \varepsilon^{\frac{7}{5}} |k|^{2\tau+1} e^{-|k|r} e^{-\rho^{(1)} n_{\alpha\beta}^{*}},$$

$$(4.19)$$

$$|W'_{mn}|_{\mathcal{O}_{+}} \leq \begin{cases} \varepsilon e^{-\rho \max\{|m|,|n|\}}, & |m|,|n| \leq N_{+}, & m,n \notin \mathcal{J} \\ 0, & otherwise \end{cases}, \tag{4.20}$$

$$\partial_{q_n} F = \partial_{\bar{q}_n} F = \partial_{q_n} \dot{P} = \partial_{\bar{q}_n} \dot{P} \equiv 0, \quad n \in \mathcal{J}.$$
 (4.21)

Proof of Proposition 2: We decompose the proof into the following parts.

• Truncation and approximate linearized equations

At first, we rewrite R as

$$R = \sum_{\substack{k \ |l| < 1}} P_{kl00} e^{\mathrm{i}\langle k, \theta \rangle} I^l + \sum_{k} (\langle P^{k10}, q \rangle + \langle P^{k01}, \bar{q} \rangle + \langle P^{k20}q, q \rangle + \langle P^{k11}q, \bar{q} \rangle + \langle P^{k02}\bar{q}, \bar{q} \rangle) e^{\mathrm{i}\langle k, \theta \rangle},$$

where P^{k10} , P^{k01} , P^{k20} , P^{k11} , P^{k02} respectively denote

The gauge invariance (A5) implies that P^{010} , P^{001} , P^{020} , $P^{002} \equiv 0$.

We try to construct a Hamiltonian F, of the same form as R, such that

$$\{\mathcal{N}, F\} + R = e' + \langle \omega', I \rangle + \langle P^{011}q, \bar{q} \rangle. \tag{4.22}$$

By a straightforward calculation and simple comparison of coefficients, Equation (4.22) is equivalent to the following equations for $k \neq 0$ and $|l| \leq 1$,

$$\langle k, \omega \rangle F_{kl00} = i P_{kl00}, \tag{4.23}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} - \Omega) F^{k10} = i P^{k10}, \tag{4.24}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} + \Omega) F^{k01} = i P^{k01}, \tag{4.25}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} - \Omega) F^{k20} - F^{k20} \Omega = i P^{k20}, \tag{4.26}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} - \Omega) F^{k11} + F^{k11} \Omega = i P^{k11}, \tag{4.27}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} + \Omega) F^{k02} + F^{k02} \Omega = i P^{k02}. \tag{4.28}$$

In view of the definition of \mathcal{O}_+ , we know that (4.23) is solved on \mathcal{O}_+ , with

$$|F_{kl00}|_{\mathcal{O}_+} \le \gamma^{-2}|k|^{2\tau+1}\varepsilon e^{-|k|r}.$$

As for (4.24) - (4.28), consider the equations

$$\left(\langle k, \omega \rangle I_{\mathbb{Z}} - (D + \hat{A})\right) \hat{F}^{k10} = i\hat{R}^{k10}, \tag{4.29}$$

$$\left(\langle k, \omega \rangle I_{\mathbb{Z}} + (D + \hat{A})\right) \hat{F}^{k01} = i\hat{R}^{k01}, \tag{4.30}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} - (D + \hat{A})) \hat{F}^{k20} - \hat{F}^{k20}(D + \hat{A}) = i\hat{R}^{k20}, \tag{4.31}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} - (D + \hat{A})) \hat{F}^{k11} + \hat{F}^{k11}(D + \hat{A}) = i\hat{R}^{k11}, \tag{4.32}$$

$$(\langle k, \omega \rangle I_{\mathbb{Z}} + (D + \hat{A})) \hat{F}^{k02} + \hat{F}^{k02}(D + \hat{A}) = i\hat{R}^{k02}$$
(4.33)

instead, where D and A are defined in the previous subsection, and for $k \neq 0$,

$$\hat{R}_n^{kx} = \begin{cases} (U^* P^{kx})_n, & |n| \le K_+ \\ 0, & \text{otherwise} \end{cases}, \quad x = \text{``10"}, \text{``01"},$$
 (4.34)

$$\hat{R}_{mn}^{kx} = \begin{cases} (U^* P^{kx} U)_{mn}, & |m|, |n| \le K_+ \\ 0, & \text{otherwise} \end{cases}, \quad x = \text{"20", "11", "02"}. \tag{4.35}$$

By (4.6) and (4.15), combining with Lemma 2.1, there exists $c_3 > 0$ such that

$$|(U^*P^{kx})_n|_{\mathcal{O}} \le c_3(\sigma - \rho)^{-1} \varepsilon e^{-\rho|n|} e^{-|k|r}, \tag{4.36}$$

$$|(U^*P^{kx}U)_{mn}|_{\mathcal{O}} \le c_3(\sigma - \rho)^{-2} \varepsilon e^{-\rho \max\{|m|, |n|\}} e^{-|k|r}. \tag{4.37}$$

This means

$$|(U^*P^{kx} - \hat{R}^{kx})_n|_{\mathcal{O}} \le \frac{1}{4}\varepsilon^{\frac{7}{5}}e^{-\rho^{(1)}|n|}e^{-|k|r},\tag{4.38}$$

$$|(U^*P^{kx}U - \hat{R}^{kx})_{mn}|_{\mathcal{O}} \le \frac{1}{4}\varepsilon^{\frac{7}{5}}e^{-\rho^{(1)}\max\{|m|,|n|\}}e^{-|k|r}$$
(4.39)

under the assumption that

(C3)
$$c_3(\sigma-\rho)^{-4}e^{-(\rho-\rho^{(1)})K_+} \leq \frac{1}{4}\varepsilon^{\frac{2}{5}}.$$

Equation (4.29) - (4.33) provide us with approximate solutions to (4.24) - (4.28), with the error estimated later.

\bullet Block-diagonalization and construction of F

Consider the equations

$$\left(\langle k, \omega \rangle I_{\Lambda} - (\tilde{D}_{\Lambda} + \tilde{A})\right) \tilde{F}^{k10} = i\tilde{R}^{k10}, \tag{4.40}$$

$$(\langle k, \omega \rangle I_{\Lambda} + (\tilde{D}_{\Lambda} + \tilde{A})) \tilde{F}^{k01} = i\tilde{R}^{k01}, \tag{4.41}$$

$$(\langle k, \omega \rangle I_{\Lambda} - (\tilde{D}_{\Lambda} + \tilde{A})) \tilde{F}^{k20} - \tilde{F}^{k20} (\tilde{D}_{\Lambda} + \tilde{A}) = i\tilde{R}^{k20}, \tag{4.42}$$

$$\left(\langle k, \omega \rangle I_{\Lambda} - (\tilde{D}_{\Lambda} + \tilde{A})\right) \tilde{F}^{k11} + \tilde{F}^{k11}(\tilde{D}_{\Lambda} + \tilde{A}) = i\tilde{R}^{k11}, \tag{4.43}$$

$$\left(\langle k, \omega \rangle I_{\Lambda} + (\tilde{D}_{\Lambda} + \tilde{A})\right) \tilde{F}^{k02} + \tilde{F}^{k02} (\tilde{D}_{\Lambda} + \tilde{A}) = i\tilde{R}^{k02}, \tag{4.44}$$

where \tilde{D}_{Λ} , \tilde{A} are defined as in (4.11) via the orthogonal matrix Q, and

$$\tilde{R}^{kx} := \begin{cases} Q^* \hat{R}^{kx}, & x = \text{``10"}, \text{``01"} \\ Q^* \hat{R}^{kx} Q, & x = \text{``20"}, \text{``11"}, \text{``02"} \end{cases}.$$

Note that $Q_{mn} = 0$ if |m - n| > N, then by (4.34) and (4.35), we have

$$\begin{split} \tilde{R}_{n}^{kx} &\equiv 0, \ \ \text{if} \ \ |n| > K_{+} + N, \quad x = \text{``10"}, \text{``01"}, \\ \tilde{R}_{mn}^{kx} &\equiv 0, \ \ \text{if} \ \ |m| \ \ \text{or} \ \ |n| > K_{+} + N, \quad x = \text{``20"}, \text{``11"}, \text{``02"}. \end{split}$$

Thus, recalling that $\Lambda := \bigcup \{\Lambda_j : \Lambda_j \cap [-(K_+ + N), K_+ + N] \neq \emptyset\}$, solutions of these finite-dimensional equations satisfy

$$\begin{split} \tilde{F}_n^{kx} &\equiv 0, &\text{if} \quad n \not\in \Lambda, \quad x = \text{``10"}, \text{``01"}, \\ \tilde{F}_{mn}^{kx} &\equiv 0, &\text{if} \quad m \text{ or } n \not\in \Lambda, \quad x = \text{``20"}, \text{``11"}, \text{``02"}. \end{split}$$

Then, in view of the facts

$$\left(\langle k, \omega \rangle I_{\mathbb{Z}} \pm (\tilde{D} + \tilde{A}) \right) \tilde{F}^{kx} = \left(\langle k, \omega \rangle I_{\Lambda} \pm (\tilde{D}_{\Lambda} + \tilde{A}) \right) \tilde{F}^{kx}, \quad x = \text{``10"}, \text{``01"},$$

$$\left(\langle k, \omega \rangle I_{\mathbb{Z}} \pm (\tilde{D} + \tilde{A}) \right) \tilde{F}^{kx} = \left(\langle k, \omega \rangle I_{\Lambda} \pm (\tilde{D}_{\Lambda} + \tilde{A}) \right) \tilde{F}^{kx}, \quad x = \text{``20"}, \text{``11"}, \text{``02"},$$

$$\tilde{F}^{kx} (\tilde{D} + \tilde{A}) = \tilde{F}^{kx} (\tilde{D}_{\Lambda} + \tilde{A}), \quad x = \text{``20"}, \text{``11"}, \text{``02"},$$

they are also solutions of

$$\begin{split} \left(\langle k,\omega\rangle I_{\mathbb{Z}} - (\tilde{D} + \tilde{A})\right)\tilde{F}^{k10} &= \mathrm{i}\tilde{R}^{k10},\\ \left(\langle k,\omega\rangle I_{\mathbb{Z}} + (\tilde{D} + \tilde{A})\right)\tilde{F}^{k01} &= \mathrm{i}\tilde{R}^{k01},\\ \left(\langle k,\omega\rangle I_{\mathbb{Z}} - (\tilde{D} + \tilde{A})\right)\tilde{F}^{k20} - \tilde{F}^{k20}(\tilde{D} + \tilde{A}) &= \mathrm{i}\tilde{R}^{k20},\\ \left(\langle k,\omega\rangle I_{\mathbb{Z}} - (\tilde{D} + \tilde{A})\right)\tilde{F}^{k11} + \tilde{F}^{k11}(\tilde{D} + \tilde{A}) &= \mathrm{i}\tilde{R}^{k11},\\ \left(\langle k,\omega\rangle I_{\mathbb{Z}} + (\tilde{D} + \tilde{A})\right)\tilde{F}^{k02} + \tilde{F}^{k02}(\tilde{D} + \tilde{A}) &= \mathrm{i}\tilde{R}^{k02},\\ \end{split}$$

which are respectively equivalent to Equation (4.29) - (4.33) since D can be block-diagonalized by the orthogonal matrix Q.

Now we focus on the following equations

$$\begin{split} &(\langle k,\omega\rangle - \mu_n) \check{F}_n^{k10} = \mathrm{i}(O^* \tilde{R}^{k10})_n, \\ &(\langle k,\omega\rangle + \mu_n) \check{F}_n^{k01} = \mathrm{i}(O^* \tilde{R}^{k01})_n, \\ &(\langle k,\omega\rangle - \mu_m - \mu_n) \check{F}_{mn}^{k20} = \mathrm{i}(O^* \tilde{R}^{k20}O)_{mn}, \\ &(\langle k,\omega\rangle - \mu_m + \mu_n) \check{F}_{mn}^{k11} = \mathrm{i}(O^* \tilde{R}^{k11}O)_{mn}, \\ &(\langle k,\omega\rangle + \mu_m + \mu_n) \check{F}_{mn}^{k02} = \mathrm{i}(O^* \tilde{R}^{k02}O)_{mn}, \end{split}$$

for $k \neq 0$ and $m, n \in \Lambda$, which is transformed from (4.40) - (4.44) by diagonalizing $\tilde{D}_{\Lambda} + \tilde{A}$ via the orthogonal matrix O. Obviously, these equations can be solved in \mathcal{O}_+ . Hence, (4.29) - (4.33) are solved with

$$\hat{F}^{kx} = \left\{ \begin{array}{ll} QO\check{F}^{kx}, & x = \text{``}10\text{''}, \text{``}01\text{''} \\ QO\check{F}^{kx}O^*Q^*, & x = \text{``}20\text{''}, \text{``}11\text{''}, \text{``}02\text{''} \end{array} \right..$$

Let

$$F^{kx} := \left\{ \begin{array}{ll} U\hat{F}^{kx}, & x = \text{``10"}, \text{``01"} \\ U\hat{F}^{kx}U^*, & x = \text{``20"}, \text{``11"}, \text{``02"} \end{array} \right.,$$

then we obtain a Hamiltonian

$$F = \sum_{k \neq 0 \atop |l| \leq 1} F_{kl00} e^{\mathrm{i}\langle k, \theta \rangle} I^l + \sum_{k \neq 0} (\langle F^{k10}, q \rangle + \langle F^{k01}, \bar{q} \rangle + \langle F^{k20}q, q \rangle + \langle F^{k11}q, \bar{q} \rangle + \langle F^{k02}\bar{q}, \bar{q} \rangle) e^{\mathrm{i}\langle k, \theta \rangle}.$$

It is easy to see that $\bar{F} = F$, by noting

$$\frac{\overline{F_{(-k)l00}}}{F^{(-k)20}} = F_{kl00}, \quad \overline{F^{(-k)10}} = F^{k01}, \quad \overline{F^{(-k)01}} = F^{k10}, F^{(-k)20} = F^{k02}, \quad (F^{(-k)11})^* = F^{k11}, \quad \overline{F^{(-k)02}} = F^{k20}.$$

• Estimates for coefficients of F

Let us consider F_{mn}^{k20} for instance, and the other terms can be treated in an analogous way. By the construction above, one sees that

$$F_{mn}^{k20} = i \sum_{\mathcal{F}_0} \frac{U_{mn_1} Q_{n_1 n_2} O_{n_2 n_3} O_{n_3 n_4}^* Q_{n_4 n_5}^* \hat{R}_{n_5 n_6}^{k20} Q_{n_6 n_7} O_{n_7 n_8} O_{n_8 n_9}^* Q_{n_9 n_{10}}^* U_{n_1 0n}^*}{\langle k, \omega \rangle - \mu_{n_3} - \mu_{n_8}}, \quad (4.45)$$

where the summation notation \mathcal{F}_0 denotes

$$\left\{ \begin{array}{l} n_1 \in \mathbb{Z}, & |n_2 - n_1| \le N, & |n_3 - n_2|, & |n_4 - n_3| \le 2(M+2)N, & |n_5 - n_4| \le N, \\ n_{10} \in \mathbb{Z}, & |n_9 - n_{10}| \le N, & |n_8 - n_9|, & |n_7 - n_8| \le 2(M+2)N, & |n_6 - n_7| \le N \end{array} \right\}$$

by virtue of the structure of Q and O, i.e, (4.8) and (4.12). Then, by (4.37) and Lemma

$$\sup_{\xi \in \mathcal{O}_+} |F_{mn}^{k20}(\xi)| \le c(\gamma^{-1}|k|^{\tau}N_+^4)(\sigma - \rho)^{-4}M^4N^8e^{(4M+10)N\rho}\varepsilon e^{-|k|r}e^{-\rho\max\{|m|,|n|\}}.$$

Here we have applied the property of the orthogonal matrices Q and O, and used the factor $e^{(4M+10)N\rho}$ to recover the exponential decay.

To estimate $|\partial_{\xi_i} F_{mn}^{k20}|$, we need to differentiate both sides of (4.42) with respect to ξ_j , $j=1,2,\cdots,b$. Then we obtain the equation about $\partial_{\xi_i}\tilde{F}^{k20}$

$$(\langle k,\omega\rangle I_{\Lambda}-(\tilde{D}_{\Lambda}+\tilde{A}))(\partial_{\xi_{j}}\tilde{F}^{k20})-(\partial_{\xi_{j}}\tilde{F}^{k20})(\tilde{D}_{\Lambda}+\tilde{A})=\breve{R}_{\xi_{j}}^{k20},$$

which can also be solved by diagonalizing $\tilde{D}_{\Lambda} + \tilde{A}$ via O as above, where

$$\breve{R}_{\xi_{i}}^{k20}:=\mathrm{i}\partial_{\xi_{i}}\tilde{R}^{k20}+\tilde{F}^{k20}(\partial_{\xi_{i}}\tilde{A})-(\partial_{\xi_{i}}(\langle k,\omega\rangle I-\tilde{A}))\tilde{F}^{k20}.$$

We get the formulation

$$\partial_{\xi_j} F_{mn}^{k20} = \sum_{\mathcal{F}_1} \frac{U_{mn_1} Q_{n_1 n_2} O_{n_2 n_3} O_{n_3 n_4}^* (\breve{R}_{\xi_j}^{k20})_{n_4 n_5} O_{n_5 n_6} O_{n_6 n_7}^* Q_{n_7 n_8}^* U_{n_8 n_5}^*}{\langle k, \omega \rangle - \mu_{n_3} - \mu_{n_6}}$$

with \mathcal{F}_1 denotes

$$\left\{\begin{array}{ll} n_1 \in \mathbb{Z}, & |n_2 - n_1| \le N, & |n_3 - n_2|, |n_4 - n_3| \le 2(M+2)N, \\ n_8 \in \mathbb{Z}, & |n_7 - n_8| \le N, & |n_6 - n_7|, |n_5 - n_6| \le 2(M+2)N \end{array}\right\}.$$

By the decay property of \hat{R}^{k20} and $\partial_{\xi_i}\hat{A}$, we have that

$$\sup_{\xi \in \mathcal{O}_+} |(\breve{R}_{\xi_j}^{k20})_{mn}| \le c(\gamma^{-1}|k|^{\tau+1}N_+^4)(\sigma-\rho)^{-4}M^4N^8e^{(4M+11)N\rho}\varepsilon e^{-|k|r}e^{-\rho\max\{|m|,|n|\}}.$$

Thus there exists $c_4 > 0$ such that

$$\sup_{\xi \in \mathcal{O}_{+}} (|F_{mn}^{k20}| + |\partial_{\xi} F_{mn}^{k20}|)$$

$$\leq c_{4} (\gamma^{-2} |k|^{2\tau+1} N_{+}^{8}) (\sigma - \rho)^{-6} M^{8} N^{14} e^{(8M+20)N\rho} \varepsilon e^{-\rho \max\{|m|,|n|\}} e^{-|k|r}$$

$$\leq \varepsilon^{\frac{4}{5}} |k|^{2\tau+1} e^{-|k|r} e^{-\rho \max\{|m|,|n|\}},$$

under the assumption

(C4)
$$c_4 \gamma^{-2} (\sigma - \rho)^{-6} N_+^8 M^8 N^{14} e^{(8M+20)N\rho} \varepsilon^{\frac{1}{5}} \le 1.$$

Suppose that $\sum_{i=1}^b k_i + 2 \neq 0$, which means $P^{k20} \equiv 0$. Then $\hat{R}^{k20} \equiv 0$, since it is a truncation of $U^*P^{k20}U$. By the formulation of F^{k20}_{mn} in (4.45), $F^{k20} \equiv 0$. Doing the same thing for F^{k11} , F^{k02} , F^{k10} , F^{k01} as above, we obtain the gauge invari-

ance of F and the inequality (4.18).

ullet Estimates for coefficients of \dot{P}

Let W' be the truncation of P^{011} , satisfying

$$W'_{mn} = \left\{ \begin{array}{ll} P_{mn}^{011}, & |m|, |n| \leq N_+ \\ 0, & \text{otherwise} \end{array} \right.,$$

and

$$\grave{P} = \langle \grave{P}^{011}q, \bar{q} \rangle + \sum_{k \neq 0} (\langle \grave{P}^{k10}, q \rangle + \langle \grave{P}^{k01}, \bar{q} \rangle + \langle \grave{P}^{k20}q, q \rangle + \langle \grave{P}^{k11}q, \bar{q} \rangle + \langle \grave{P}^{k02}\bar{q}, \bar{q} \rangle) e^{\mathrm{i}\langle k, \theta \rangle}$$

with

$$\begin{array}{lll} \grave{P}^{011} & := & P^{011} - W', \\ \grave{P}^{k10} & := & (P^{k10} - U\hat{R}^{k10}) - \mathrm{i}(\grave{A} + \grave{Z})F^{k10}, \\ \grave{P}^{k01} & := & (P^{k01} - U\hat{R}^{k01}) + \mathrm{i}(\grave{A} + \grave{Z})F^{k01}, \\ \grave{P}^{k20} & := & (P^{k20} - U\hat{R}^{k20}U^*) - \mathrm{i}(\grave{A} + \grave{Z})F^{k20} - \mathrm{i}F^{k20}(\grave{A} + \grave{Z}), \\ \grave{P}^{k11} & := & (P^{k11} - U\hat{R}^{k11}U^*) - \mathrm{i}(\grave{A} + \grave{Z})F^{k11} + \mathrm{i}F^{k11}(\grave{A} + \grave{Z}), \\ \grave{P}^{k02} & := & (P^{k02} - U\hat{R}^{k02}U^*) + \mathrm{i}(\grave{A} + \grave{Z})F^{k02} + \mathrm{i}F^{k02}(\grave{A} + \grave{Z}), \end{array}$$

where $\dot{A} := (A + W) - U\hat{A}U^*$, and $\dot{Z} := UZU^*$. Then we obtain

$$\{\mathcal{N}, F\} + R = e' + \langle \omega', I \rangle + \langle W'q, \bar{q} \rangle + \dot{P}. \tag{4.46}$$

By (4.4) and (4.5), we have (4.20) holds and

$$|\dot{P}_{mn}^{011}|_{\mathcal{O}_{\perp}} \leq \varepsilon e^{-\rho \max\{|m|,|n|\}} \leq \varepsilon^{\frac{7}{5}} e^{-\rho^{(1)} \max\{|m|,|n|\}},$$

under the assumption

(C5)
$$e^{-(\rho-\rho^{(1)})N_+} \leq \varepsilon^{\frac{2}{5}}$$
.

As for the case $k \neq 0$ in (4.19), we only estimate \dot{P}^{k20} , with the others entirely analogous. By (4.39) and (C3), combining with Lemma 2.1,

$$\left| \left(P^{k20} - U\hat{R}^{k20}U^* \right)_{mn} \right|_{\mathcal{O}} = \left| \left(U(U^*P^{k20}U - \hat{R}^{k20})U^* \right)_{mn} \right|_{\mathcal{O}} \le \frac{1}{4} \varepsilon^{\frac{7}{5}} e^{-\rho^{(1)} \max\{|m|, |n|\}} e^{-|k|r}.$$
(4.47)

In view of (4.7) and (4.10),

$$|\dot{A}_{mn}|_{\mathcal{O}} \le c(\sigma - \rho)^{-2} \varepsilon e^{-\rho \max\{|m|,|n|\}}, \quad |\dot{Z}_{mn}| \le c(\sigma - \rho)^{-2} \varepsilon e^{-\rho|m-n|}.$$

Then, by applying Lemma 2.1 again, there exists $c_6 > 0$ such that

$$\left| \left(F^{k20} (\dot{A} + \dot{Z}) \right)_{mn} \right|_{\mathcal{O}_{+}}, \quad \left| \left((\dot{A} + \dot{Z}) F^{k20} \right)_{mn} \right|_{\mathcal{O}_{+}} \\
\leq c_{6} (\sigma - \rho)^{-2} (\rho - \rho^{(1)})^{-1} \varepsilon^{\frac{9}{5}} |k|^{2\tau + 1} e^{-|k|r} e^{-\rho^{(1)} \max\{|m|, |n|\}} \\
\leq \frac{1}{4} \varepsilon^{\frac{7}{5}} |k|^{2\tau + 1} e^{-|k|r} e^{-\rho^{(1)} \max\{|m|, |n|\}}, \tag{4.48}$$

provided that

(C6)
$$c_6(\sigma-\rho)^{-2}(\rho-\rho^{(1)})^{-1}\varepsilon^{\frac{2}{5}} \leq \frac{1}{4}$$
.

Thus, we can obtain the estimate for \hat{P}^{k20} by putting (4.47) and (4.48) together. By the construction of \hat{P} , the gauge invariance is easily verified.

• Verification of (4.21)

In view of the construction of R and W' above, the objects in (4.46) that may depend on the variables $(q_n, \bar{q}_n)_{n \in \mathcal{J}}$ are F and \dot{P} . Let

$$\dot{F} = \sum_{k \neq 0} \left(\sum_{n \in \mathcal{J}} (F_n^{k10} q_n + F_n^{k01} \bar{q}_n) + \sum_{m \text{ or } n \in \mathcal{J}} (F_{mn}^{k20} q_m q_n + F_{mn}^{k11} q_m \bar{q}_n + F_{mn}^{k02} \bar{q}_m \bar{q}_n) \right) e^{i\langle k, \theta \rangle}$$

$$=: \sum_{k \neq 0} \left(\langle \acute{F}^{k10}, q \rangle + \langle \acute{F}^{k01}, \bar{q} \rangle + \langle \acute{F}^{k20} q, q \rangle + \langle \acute{F}^{k11} q, \bar{q} \rangle + \langle \acute{F}^{k02} \bar{q}, \bar{q} \rangle \right) e^{i\langle k, \theta \rangle}.$$

For m or $n \in \mathcal{J}$, by (4.2), we have

$$\begin{split}
&\left((\langle k,\omega\rangle I_{\mathbb{Z}}-\Omega)\acute{F}^{k20}-\acute{F}^{k20}\Omega\right)_{mn} &= \langle k,\omega\rangle \acute{F}^{k20}_{mn}-\sum_{l\not\in\mathcal{J}}\Omega_{ml}\acute{F}^{k20}_{ln}-\sum_{l\not\in\mathcal{J}}\acute{F}^{k20}_{ml}\Omega_{ln} \\
&= \begin{cases}
\langle k,\omega\rangle F^{k20}_{mn}, & m,n\in\mathcal{J} \\
\langle k,\omega\rangle F^{k20}_{mn}-\sum_{l\not\in\mathcal{J}}\Omega_{ml}F^{k20}_{ln}, & m\not\in\mathcal{J}, & n\in\mathcal{J} \\
\langle k,\omega\rangle F^{k20}_{mn}-\sum_{l\not\in\mathcal{J}}F^{k20}_{ml}\Omega_{ln}, & m\in\mathcal{J}, & n\notin\mathcal{J}
\end{cases} \\
&= \left((\langle k,\omega\rangle I_{\mathbb{Z}}-\Omega)F^{k20}-F^{k20}\Omega\right)_{mn}.$$

This means, by comparing the coefficients in both side of Equation (4.46),

$$\left((\langle k,\omega\rangle I_{\mathbb{Z}}-\Omega)\acute{F}^{k20}-\acute{F}^{k20}\Omega\right)_{mn}=-\mathrm{i}\grave{P}_{mn}^{k20}\quad m\ \mathrm{or}\ n\in\mathcal{J}.$$

Similarly,

$$\begin{split} &\left((\langle k,\omega\rangle I_{\mathbb{Z}}-\Omega) \acute{F}^{k10}\right)_{n}=-\mathrm{i} \grave{P}_{n}^{k10}, \quad n\in\mathcal{J},\\ &\left((\langle k,\omega\rangle I_{\mathbb{Z}}+\Omega) \acute{F}^{k01}\right)_{n}=-\mathrm{i} \grave{P}_{n}^{k01}, \quad n\in\mathcal{J},\\ &\left((\langle k,\omega\rangle I_{\mathbb{Z}}-\Omega) \acute{F}^{k11}+\acute{F}^{k11}\Omega\right)_{mn}=-\mathrm{i} \grave{P}_{mn}^{k11}, \quad m \ \text{or} \ n\in\mathcal{J},\\ &\left((\langle k,\omega\rangle I_{\mathbb{Z}}+\Omega) \acute{F}^{k02}+\acute{F}^{k02}\Omega\right)_{mn}=-\mathrm{i} \grave{P}_{mn}^{k02}, \quad m \ \text{or} \ n\in\mathcal{J}. \end{split}$$

Thus, $\{\mathcal{N}, \acute{F}\}$ equals to

$$\sum_{k \neq 0} \left(\sum_{n \in \mathcal{J}} (\grave{P}_n^{k10} q_n + \grave{P}_n^{k01} \bar{q}_n) + \sum_{m \text{ or } n \in \mathcal{J}} (\grave{P}_{mn}^{k20} q_m q_n + \grave{P}_{mn}^{k11} q_m \bar{q}_n + \grave{P}_{mn}^{k02} \bar{q}_m \bar{q}_n) \right) e^{i\langle k, \theta \rangle}.$$

Hence, if we substitute F with $F - \acute{F}$, which is independent of the variables $(q_n, \bar{q}_n)_{n \in \mathcal{J}}$, then \grave{P} in the homological equation (4.46) is replaced correspondingly, independent of $(q_n, \bar{q}_n)_{n \in \mathcal{J}}$. (4.21) is satisfied.

4.3 Verification of assumptions after one sub-step

We proceed to estimate the norm of X_F , and to study properties of Φ_F^1 on smaller domains $\mathcal{D}_i := \mathcal{D}_{d,\rho_+}(r^{(1)} + \frac{i}{4}(r - r^{(1)}), \frac{i}{4}s), i = 1, 2, 3, 4.$

Lemma 4.1 For ε sufficiently small, we have $\|X_F\|_{\mathcal{D}_3,\mathcal{O}_+} \leq \varepsilon^{\frac{3}{4}}$ and $\|X_{\grave{P}}\|_{\mathcal{D}_3,\mathcal{O}_+} \leq \varepsilon^{\frac{5}{4}}$.

Proof: In view of the decay property of F in Proposition 2, it follows that

$$\frac{1}{s^2} \|\partial_{\theta} F\|_{\mathcal{D}_3, \mathcal{O}_+}, \|\partial_I F\|_{\mathcal{D}_3, \mathcal{O}_+} \le c(r - r^{(1)})^{-(2\tau + b + 1)} \varepsilon^{\frac{4}{5}},$$

and

$$\sup_{\mathcal{D}_{3}} \frac{1}{s} \sum_{n \in \mathbb{Z}} \left(\|\partial_{q_{n}} F\|_{\mathcal{O}_{+}} + \|\partial_{\bar{q}_{n}} F\|_{\mathcal{O}_{+}} \right) \langle n \rangle^{d} e^{\rho_{+} |n|} \\
\leq \sup_{\mathcal{D}_{3}} \frac{c}{s} \sum_{n \in \mathbb{Z}} \sum_{k \neq 0} \left(|F_{n}^{k10}|_{\mathcal{O}_{+}} + |F_{n}^{k01}|_{\mathcal{O}_{+}} \right) e^{|k|(r - \frac{1}{4}(r - r^{(1)}))} \langle n \rangle^{d} e^{\rho_{+} |n|} \\
+ \sup_{\mathcal{D}_{3}} \frac{c}{s} \sum_{n \in \mathbb{Z}} \sum_{\substack{k \neq 0 \\ m \in \mathbb{Z}}} \left(|F_{mn}^{k20}|_{\mathcal{O}_{+}} + |F_{mn}^{k11}|_{\mathcal{O}_{+}} + |F_{mn}^{k02}|_{\mathcal{O}_{+}} \right) |q_{m}| e^{|k|(r - \frac{1}{4}(r - r^{(1)}))} \langle n \rangle^{d} e^{\rho_{+} |n|} \\
\leq c(r - r^{(1)})^{-(2\tau + b + 1)} (\rho - \rho_{+})^{-2} \varepsilon^{\frac{4}{5}}.$$

Putting together the estimates above, there is a constant $c_7 > 0$ such that

$$||X_F||_{\mathcal{D}_3,\mathcal{O}_+} \le c_7(r-r^{(1)})^{-(2\tau+b+1)}(\rho-\rho_+)^{-2}\varepsilon^{\frac{4}{5}}.$$

In an entirely analogous way, we have

$$||X_{\dot{\rho}}||_{\mathcal{D}_3,\mathcal{O}_+} \le c_7 (r - r^{(1)})^{-(2\tau + b + 1)} (\rho^{(1)} - \rho_+)^{-2} \varepsilon^{\frac{7}{5}}.$$

Moreover, if

(C7)
$$c_7(r-r^{(1)})^{-(2\tau+b+1)}(\rho^{(1)}-\rho_+)^{-2}\varepsilon^{\frac{1}{20}} \leq \frac{1}{3}$$

then Lemma 4.1 follows.

Let
$$\mathcal{D}_{i\eta} = \mathcal{D}_{d,\rho_+}(r^{(1)} + \frac{i}{4}(r - r^{(1)}), \frac{i}{4}\eta s), i = 1, 2, 3, 4.$$

Lemma 4.2 For ε sufficiently small, we have $\Phi_F^t: \mathcal{D}_{2\eta} \to \mathcal{D}_{3\eta}, -1 \leq t \leq 1$ and moreover,

$$||D\Phi_F^t - I||_{\mathcal{D}_{1\eta}} < 2\varepsilon^{\frac{3}{4}}.$$

Let $F^{(1)}$, $e^{(1)}$, $\omega^{(1)}$, $W^{(1)}$, $\dot{P}^{(1)}$ be the corresponding quantities in (4.17) respectively, which means that we are in the 1st sub-step. Define $H^{(1)}$ as

$$\begin{split} H^{(1)} &:= H \circ \Phi_{F^{(1)}}^1 \\ &= (\mathcal{N} + \check{P} + R) \circ \Phi_{F^{(1)}}^1 + (P - R) \circ \Phi_{F^{(1)}}^1 \\ &= \mathcal{N} + \check{P} + \{\mathcal{N}, F^{(1)}\} + R + \int_0^1 (1 - t) \{\{\mathcal{N}, F^{(1)}\}, F^{(1)}\} \circ \Phi_{F^{(1)}}^t \, dt \\ &+ \int_0^1 \{\check{P} + R, F^{(1)}\} \circ \Phi_{F^{(1)}}^t \, dt + (P - R) \circ \Phi_{F^{(1)}}^1 \\ &= \mathcal{N} + \check{P} + e^{(1)} + \langle \omega^{(1)}, I \rangle + \langle W^{(1)}q, \bar{q} \rangle + P^{(1)}, \end{split}$$

where

$$P^{(1)} := \grave{P}^{(1)} + \int_0^1 \{(1-t)\{\mathcal{N}, F^{(1)}\} + \check{P} + R, F^{(1)}\} \circ \Phi_{F^{(1)}}^t dt + (P-R) \circ \Phi_{F^{(1)}}^1.$$

Let $R(t) := (1-t)(e^{(1)} + \langle \omega^{(1)}, I \rangle + \langle W^{(1)}q, \bar{q} \rangle + \dot{P}^{(1)}) + tR$, which satisfies $||X_{R(t)}||_{\mathcal{D}_3} \le c\varepsilon$. Then $P^{(1)}$ can be written as

$$P^{(1)} = \grave{P}^{(1)} + \int_0^1 \{R(t) + \check{P}, F^{(1)}\} \circ \Phi_{F^{(1)}}^t dt + (P - R) \circ \Phi_{F^{(1)}}^1.$$

Hence,

$$X_{P^{(1)}-\grave{P}^{(1)}} \,=\, \int_0^1 (\Phi_{F^{(1)}}^t)^* X_{\{R(t)+\check{P},\,F^{(1)}\}} \,dt \,+\, (\Phi_{F^{(1)}}^1)^* X_{(P-R)}.$$

By Lemma A.4,

$$\|X_{\{R(t)+\check{P},F^{(1)}\}}\|_{\mathcal{D}_{2\eta}} \le c\eta^{-2}\varepsilon^{\frac{7}{4}} = \varepsilon^{\frac{27}{20}}.$$

Then, combining with (4.16), recalling the conclusion of Lemma 4.1 and 4.2,

$$\|X_{P^{(1)}}\|_{\mathcal{D}^{(1)},\mathcal{O}_+} \leq \frac{1}{2}\varepsilon^{\frac{6}{5}} + 2\varepsilon^{\frac{5}{4}} + 2c\varepsilon^{\frac{27}{20}} \leq \varepsilon^{\frac{6}{5}} = \varepsilon^{(1)}.$$

Now we need to show $P^{(1)}$ satisfies assumptions (A4) and (A5). Note that

$$P^{(1)} = \dot{P}^{(1)} + P - R + \{ \breve{P}, F^{(1)} \} + \{ P, F^{(1)} \}$$

$$+ \frac{1}{2!} \{ \{ \mathcal{N}, F^{(1)} \}, F^{(1)} \} + \frac{1}{2!} \{ \{ \breve{P}, F^{(1)} \}, F^{(1)} \} + \frac{1}{2!} \{ \{ P, F^{(1)} \}, F^{(1)} \} + \cdots$$

$$+ \frac{1}{n!} \{ \cdots \{ \mathcal{N}, \underbrace{F^{(1)} \} \cdots, F^{(1)}}_{n} \} + \frac{1}{n!} \{ \cdots \{ \breve{P}, \underbrace{F^{(1)} \} \cdots, F^{(1)}}_{n} \}$$

$$+ \frac{1}{n!} \{ \cdots \{ P, \underbrace{F^{(1)} \} \cdots, F^{(1)}}_{n} \} + \cdots .$$

Since all of \mathcal{N} , \check{P} , P, $F^{(1)}$, $\grave{P}^{(1)}$ have gauge invariance, independent of variables $(q_n, \bar{q}_n)_{n \in \mathcal{J}}$, so does $P^{(1)}$ due to Lemma A.5 and A.6 in Appendix.

so does
$$P^{(1)}$$
 due to Lemma A.5 and A.6 in Appendix.
For $P - R = \sum_{2|l|+|\alpha|+|\beta| \ge 3} P_{kl\alpha\beta} e^{\mathrm{i}\langle k,\theta\rangle} I^l q^{\alpha} \bar{q}^{\beta}$, we have

$$||P_{\alpha\beta}||_{\mathcal{D}^{(1)}} \le \begin{cases} \frac{1}{4}\varepsilon^{(1)}e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \le 2\\ e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \ge 3 \end{cases}.$$

Here we applied the estimate $|I| \le s^{(1)} \le \frac{1}{4}\varepsilon^{(1)}$ to handle the case that $|\alpha| + |\beta| \le 2$ and $2|l| + |\alpha| + |\beta| \ge 3$.

The decay property of remaining terms, which are made up of several Poisson brackets, is covered by the following lemmas.

Lemma 4.3 For ε sufficiently small, $\{P, F^{(1)}\}$ satisfies

$$\|\{P, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_{3\eta}, \mathcal{O}_{+}} \leq \begin{cases} \varepsilon^{\frac{5}{4}} e^{-\rho^{(1)} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \leq 2\\ \varepsilon^{\frac{1}{4}} e^{-\rho^{(1)} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \geq 3 \end{cases}.$$

Proof: A straightforward calculation yields that

$$\{P, F^{(1)}\}_{\alpha\beta} = i \sum_{\substack{n \in \mathbb{Z} \\ (\check{\alpha}, \check{\beta}) + (\hat{\alpha}, \hat{\beta}) = (\alpha, \beta)}} \left(P_{\check{\alpha} + e_n, \check{\beta}} F_{\hat{\alpha}, \hat{\beta} + e_n}^{(1)} - P_{\check{\alpha}, \check{\beta} + e_n} F_{\hat{\alpha} + e_n, \hat{\beta}}^{(1)} \right)$$
(4.49)

$$+ \sum_{(\check{\alpha},\check{\beta})+(\hat{\alpha},\hat{\beta})=(\alpha,\beta)} \left\{ P_{\check{\alpha}\check{\beta}}, F_{\hat{\alpha}\hat{\beta}}^{(1)} \right\}. \tag{4.50}$$

• Terms in (4.49) Let us consider terms $P_{\check{\alpha}+e_n,\check{\beta}}F_{\hat{\alpha},\hat{\beta}+e_n}^{(1)}$ first.

i)
$$|\alpha| + |\beta| \le 2$$

Since $|\hat{\alpha}| + |\hat{\beta} + e_n| = 1$ or 2 in view of the construction of $F^{(1)}$, we have that

$$|\check{\alpha} + e_n| + |\check{\beta}| = |\alpha| + |\beta| + 1 - (|\hat{\alpha}| + |\hat{\beta}|) \le 3.$$
 (4.51)

If $|\check{\alpha} + e_n| + |\check{\beta}| \le 2$, then, noting that $n^*_{\alpha\beta} \le \max\{n^*_{\check{\alpha} + e_n, \check{\beta}}, n^*_{\hat{\alpha}, \hat{\beta} + e_n}\}$, we have

$$\|P_{\check{\alpha}+e_n,\check{\beta}}F_{\hat{\alpha},\hat{\beta}+e_n}^{(1)}\|_{\mathcal{D}_3,\mathcal{O}_+} \le \varepsilon e^{-\rho n_{\check{\alpha}+e_n,\check{\beta}}^*} \cdot \varepsilon^{\frac{3}{4}} e^{-\rho n_{\hat{\alpha},\hat{\beta}+e_n}^*} \le \varepsilon^{\frac{7}{4}} e^{-\rho n_{\alpha\beta}^*}. \tag{4.52}$$

If $|\check{\alpha} + e_n| + |\check{\beta}| = 3$, then, by (4.51), $(\hat{\alpha}, \hat{\beta}) = (0, 0)$, $(\check{\alpha}, \check{\beta}) = (\alpha, \beta)$. By the definition of norm $||X_P||_{\mathcal{D},\mathcal{O}}$ and the construction of $F^{(1)}$,

$$||P_{\alpha+e_n,\beta}||_{\mathcal{D}_3,\mathcal{O}} \le e^{-\rho n_{\alpha+e_n,\beta}^*}, \quad ||F_{0,e_n}^{(1)}||_{\mathcal{D}_3,\mathcal{O}_+} \le s\varepsilon^{\frac{3}{4}}e^{-\rho|n|}.$$

Thus, noting that $n_{\alpha\beta}^* \leq \max\{n_{\alpha+e_n,\beta}^*, |n|\},\$

$$||P_{\alpha+e_n,\beta}F_{0,e_n}^{(1)}||_{\mathcal{D}_3,\mathcal{O}_+} \le s\varepsilon^{\frac{3}{4}}e^{-\rho n_{\alpha\beta}^*} \le \frac{1}{4}\varepsilon^{\frac{7}{4}}e^{-\rho n_{\alpha\beta}^*}.$$
 (4.53)

ii) $|\alpha| + |\beta| \ge 3$

By the same argument as above,

$$\|P_{\check{\alpha}+e_{n},\check{\beta}}F_{\hat{\alpha},\hat{\beta}+e_{n}}^{(1)}\|_{\mathcal{D}_{3},\mathcal{O}_{+}} \leq e^{-\rho n_{\check{\alpha}+e_{n},\check{\beta}}^{*}} \cdot \varepsilon^{\frac{3}{4}}e^{-\rho n_{\check{\alpha},\hat{\beta}+e_{n}}^{*}} \leq \varepsilon^{\frac{3}{4}}e^{-\rho n_{\alpha\beta}^{*}}.$$
 (4.54)

Doing the same for $P_{\check{\alpha},\check{\beta}+e_n}F_{\hat{\alpha}+e_n,\hat{\beta}}^{(1)}$, we finish estimates for terms in (4.49).

• Terms in (4.50) By Lemma A.3 and the inequality $n_{\alpha\beta}^* \leq \max\{n_{\tilde{\alpha}\tilde{\beta}}^*, n_{\hat{\alpha}\hat{\beta}}^*\}$, we have

$$\|\{P_{\check{\alpha}\check{\beta}}, F_{\hat{\alpha}\hat{\beta}}^{(1)}\}\|_{\mathcal{D}_{3\eta}} \le c(r - r^{(1)})^{-1}\eta^{-2} \begin{cases} \varepsilon^{\frac{7}{4}}e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \le 2\\ \varepsilon^{\frac{3}{4}}e^{-\rho n_{\alpha\beta}^*}, & |\alpha| + |\beta| \ge 3 \end{cases}$$
(4.55)

Combining (4.52) - (4.55), there exists $c_8 > 0$ such that

$$\|\{P, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_{3\eta}} \le c_8(r - r^{(1)})^{-1}\eta^{-2}(\rho - \rho^{(1)})^{-2} \begin{cases} \varepsilon^{\frac{7}{4}}e^{-\rho^{(1)}n_{\alpha\beta}^*}, & |\alpha| + |\beta| \le 2\\ \varepsilon^{\frac{3}{4}}e^{-\rho^{(1)}n_{\alpha\beta}^*}, & |\alpha| + |\beta| \ge 3 \end{cases},$$

applying the fact that $|\hat{\alpha}| + |\hat{\beta}| \leq 2$. Moreover, if

(C8)
$$c_8(r-r^{(1)})^{-1}\eta^{-2}(\rho-\rho^{(1)})^{-2}\varepsilon^{\frac{1}{2}} \leq \frac{1}{4}$$

Lemma 4.3 is proved.

By (4.15), (4.19) and (4.20), it is evident that the coefficients of

$$\{\mathcal{N}, F^{(1)}\} = e^{(1)} + \langle \omega^{(1)}, I \rangle + \langle W^{(1)}q, \bar{q} \rangle + \dot{P}^{(1)} - R$$

satisfies $\|\{\mathcal{N}, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_3,\mathcal{O}_+} \leq c\varepsilon e^{-\rho^{(1)}n_{\alpha\beta}^*}$. Then we have the following lemma, whose proof is analogous to that of Lemma 4.3.

Lemma 4.4 For ε sufficiently small, $\{\{\mathcal{N}, F^{(1)}\}, F^{(1)}\}$ satisfies

$$\|\{\{\mathcal{N}, F^{(1)}\}, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_{3\eta}, \mathcal{O}_{+}} \leq \frac{1}{4} \varepsilon^{\frac{6}{5}} e^{-\rho^{(1)} n_{\alpha\beta}^{*}}.$$

Lemma 4.5 For ε sufficiently small, $\{\breve{P}, F^{(1)}\}$ satisfies

$$\|\{\breve{P}, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_3, \mathcal{O}_+} \le \varepsilon^{\frac{1}{4}} e^{-\rho^{(1)} n_{\alpha\beta}^*}, \quad |\alpha| + |\beta| \ge 3.$$

Proof: It can be calculated that

$$\{\breve{P}, F^{(1)}\}_{\alpha\beta} = i \sum_{\substack{n \in \mathbb{Z} \\ (\check{\alpha}, \check{\beta}) + (\hat{\alpha}, \hat{\beta}) = (\alpha, \beta)}} \left(\breve{P}_{\check{\alpha} + e_n, \check{\beta}} F^{(1)}_{\hat{\alpha}, \hat{\beta} + e_n} - \breve{P}_{\check{\alpha}, \check{\beta} + e_n} F^{(1)}_{\hat{\alpha} + e_n, \hat{\beta}}\right). \tag{4.56}$$

For $\check{P}_{\check{\alpha}+e_n,\check{\beta}}F_{\hat{\alpha},\hat{\beta}+e_n}^{(1)}$ in (4.56), since $|\hat{\alpha}|+|\hat{\beta}+e_n|=1$ or 2 and $|\check{\alpha}+e_n|+|\check{\beta}|\geq 4$ here, it is obvious that $|\alpha|+|\beta|=|\check{\alpha}|+|\check{\beta}|+|\hat{\alpha}|+|\hat{\beta}|\leq 3$.

Note that $n_{\alpha\beta}^* \leq \max\{n_{\check{\alpha}+e_n,\check{\beta}}^*, n_{\hat{\alpha},\hat{\beta}+e_n}^*\}$, and

$$n^*_{\check{\alpha} + e_n, \check{\beta}} = \max\{n^+_{\check{\alpha} + e_n, \check{\beta}}, -n^-_{\check{\alpha} + e_n, \check{\beta}}\}, \quad n^*_{\hat{\alpha}, \hat{\beta} + e_n} = \max\{n^+_{\hat{\alpha}, \hat{\beta} + e_n}, -n^-_{\hat{\alpha}, \hat{\beta} + e_n}\}.$$

Then $n^+_{\check{\alpha}+e_n,\check{\beta}}-n^-_{\check{\alpha}+e_n,\check{\beta}}+n^*_{\hat{\alpha},\hat{\beta}+e_n}\geq n^*_{\alpha\beta}$, and hence

$$\left\| \check{P}_{\check{\alpha} + e_n, \check{\beta}} F_{\hat{\alpha}, \hat{\beta} + e_n}^{(1)} \right\|_{\mathcal{D}_3, \mathcal{O}_+} \leq e^{-\rho(n_{\check{\alpha} + e_n, \check{\beta}}^+ - n_{\check{\alpha} + e_n, \check{\beta}}^-)} \cdot \varepsilon^{\frac{3}{4}} e^{-\rho n_{\check{\alpha}, \hat{\beta} + e_n}^*} \leq \varepsilon^{\frac{3}{4}} e^{-\rho n_{\alpha\beta}^*}.$$

Doing the estimate for $\check{P}_{\check{\alpha},\check{\beta}+e_n}F^{(1)}_{\hat{\alpha}+e_n,\hat{\beta}}$ in (4.56) similarly, we have that

$$\|\{\breve{P}, F^{(1)}\}_{\alpha\beta}\|_{\mathcal{D}_3, \mathcal{O}_+} \le c_8(\rho - \rho^{(1)})^{-2} \varepsilon^{\frac{3}{4}} e^{-\rho^{(1)} n_{\alpha\beta}^*} \le \varepsilon^{\frac{1}{4}} e^{-\rho^{(1)} n_{\alpha\beta}^*}, \quad |\alpha| + |\beta| \ge 3,$$

if **(C8)** holds.

Summarize the analysis above, then the decay property for $P^{(1)}$ can be expressed as

Proposition 3 For ε sufficiently small, $P^{(1)} = \sum_{\alpha,\beta} P_{\alpha\beta}^{(1)}(\theta,I;\xi) q^{\alpha} \bar{q}^{\beta}$ satisfies

$$||P_{\alpha\beta}^{(1)}||_{\mathcal{D}^{(1)},\mathcal{O}_{+}} \leq \begin{cases} \varepsilon^{(1)} e^{-\rho^{(1)} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \leq 2 \\ e^{-\rho^{(1)} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \geq 3 \end{cases}.$$

4.4 A succession of symplectic transformations

With the verification of assumptions (A4) and (A5) completed, we finish one sub-step of KAM iteration. Suppose that we have arrived at the j^{th} sub-step, $j=1,\cdots,J$, with $J=\left\lceil\frac{5}{2}\varepsilon^{\frac{a}{2}}\right\rceil$, then we encounter the Hamiltonian

$$\begin{split} H^{(j-1)} &= H \circ \Phi^1_{F^{(1)}} \circ \dots \circ \Phi^1_{F^{(j-1)}} \\ &= \mathcal{N} + \check{P} + \sum_{i=1}^{j-1} \left(e^{(i)} + \langle \omega^{(i)}, I \rangle + \langle W^{(i)}q, \bar{q} \rangle \right) + P^{(j-1)}, \end{split}$$

with the superscript "(0)" labeling quantities before the 1st sub-step in particular. Let

$$R^{(j-1)} := \sum_{\substack{k \ 2|l|+|\alpha|+|\beta| \le 2}} P_{kl\alpha\beta}^{(j-1)} e^{i\langle k,\theta \rangle} I^l q^{\alpha} \bar{q}^{\beta}. \tag{4.57}$$

As demonstrated in Proposition 2, on \mathcal{O}_+ , the following homological equation

$$\{\mathcal{N}, F^{(j)}\} + R^{(j-1)} = e^{(j)} + \langle \omega^{(j)}, I \rangle + \langle W^{(j)}q, \bar{q} \rangle + \dot{P}^{(j)},$$
 (4.58)

can be solved, with $F^{(j)}$, $e^{(j)}$, $\omega^{(j)}$, $W^{(j)}$, $\dot{P}^{(j)}$ having properties similar to $F^{(1)}$, $e^{(1)}$, $\omega^{(1)}$, $W^{(1)}$, $\dot{P}^{(1)}$ respectively. Then we obtain

$$H^{(j)} = H^{(j-1)} \circ \Phi^1_{F^{(j)}} = \mathcal{N} + \check{P} + \sum_{i=1}^j \left(e^{(i)} + \langle \omega^{(i)}, I \rangle + \langle W^{(i)}q, \bar{q} \rangle \right) + P^{(j)}.$$

The estimates for $F^{(j)}$ and the verification of assumptions for $P^{(j)}$ can be done similarly as in subsection 4.3.

The process above can be summarized as

Proposition 4 Consider the Hamiltonian H in (4.1). There exist J symplectic transformations $\Phi^{(j)} = \Phi^1_{F^{(j)}}$, $j = 1, \dots, J$, generated by the corresponding real-analytic Hamiltonians $F^{(j)}$ respectively, such that

$$H^{(j)} = H \circ \Phi^{(1)} \circ \cdots \circ \Phi^{(j)} = \mathcal{N} + \check{P} + G_j + P^{(j)},$$

is real-analytic on $\mathcal{D}^{(j)} = \mathcal{D}_{d,\rho_+}(r^{(j)},s^{(j)})$, with $G_j = \sum_{i=1}^j \left(e^{(i)} + \langle \omega^{(i)},I \rangle + \langle W^{(i)}q,\bar{q} \rangle\right)$. For $i=1,2,3,4, \ \eta = \varepsilon^{\frac{1}{5}}$, let

$$\begin{split} \mathcal{D}_{i}^{(j)} &= \mathcal{D}_{d,\rho_{+}}(r^{(j+1)} + \frac{i}{4}(r^{(j)} - r^{(j+1)}), \frac{i}{4}s^{(j)}), \\ \mathcal{D}_{i\eta}^{(j)} &= \mathcal{D}_{d,\rho_{+}}(r^{(j+1)} + \frac{i}{4}(r^{(j)} - r^{(j+1)}), \frac{i}{4}\eta s^{(j)}). \end{split}$$

(a) With $R^{(j-1)}$ defined in (4.57), $F^{(j)}$ satisfies the homological equation (4.58) on \mathcal{O}_+ , $\|X_{F^{(j)}}\|_{\mathcal{D}_3^{(j-1)},\mathcal{O}_+} \leq \varepsilon^{-\frac{1}{4}}\varepsilon^{(j-1)}, \ \Phi_{F^{(j)}}^t: \mathcal{D}_{2\eta}^{(j-1)} \to \mathcal{D}_{3\eta}^{(j-1)}, -1 \leq t \leq 1, \ and$

$$\begin{split} & \| D\Phi_{F^{(j)}}^t - I \|_{\mathcal{D}_{1\eta}^{(j-1)}} < 2\varepsilon^{-\frac{1}{4}}\varepsilon^{(j-1)}, \\ & \| F_{\alpha\beta}^{(j)} \|_{\mathcal{D}_{3}^{(j-1)},\mathcal{O}_{+}} \leq \left\{ \begin{array}{cc} \varepsilon^{-\frac{1}{4}}\varepsilon^{(j-1)}e^{-\rho^{(j-1)}}n_{\alpha\beta}^*, & |\alpha| + |\beta| \leq 2 \\ 0, & |\alpha| + |\beta| \geq 3 \end{array} \right., \\ & \partial_{q_{n}}F^{(j)} = \partial_{\bar{q}_{n}}F^{(j)} \equiv 0, \quad \forall n \in \mathcal{J}. \end{split}$$

(b)
$$G_j$$
 satisfies that $\|X_{G_j}\|_{\mathcal{D}_2^{(j)},\mathcal{O}_+} \leq c\varepsilon$ and for $i=1,2,\cdots,j$,

$$|\omega^{(i)}|_{\mathcal{O}_{+}} \leq \varepsilon^{(i-1)},$$

$$|W_{mn}^{(i)}|_{\mathcal{O}_{+}} \leq \begin{cases} \varepsilon^{(i-1)} e^{-\rho^{(i-1)} \max\{|m|,|n|\}}, & |m|,|n| \leq N_{+}, & m,n \notin \mathcal{J} \\ 0, & otherwise \end{cases}$$

(c) $P^{(j)}$ satisfies $\|X_{\check{P}+P^{(j)}}\|_{\mathcal{D}^{(j)},\mathcal{O}_+} \leq \varepsilon^{(j)}$ and assumptions (A4), (A5), which include

$$\begin{split} \|P_{\alpha\beta}^{(j)}\|_{\mathcal{D}^{(j)},\,\mathcal{O}_{+}} &\leq \left\{ \begin{array}{l} \varepsilon^{(j)}e^{-\rho^{(j)}n_{\alpha\beta}^{*}}, & |\alpha|+|\beta| \leq 2 \\ e^{-\rho^{(j)}n_{\alpha\beta}^{*}}, & |\alpha|+|\beta| \geq 3 \end{array} \right., \\ \partial_{q_{n}}P^{(j)} &= \partial_{\bar{q}_{n}}P^{(j)} \equiv 0, \quad \forall n \in \mathcal{J}. \end{split}$$

Let
$$s_+ = s^{(J)} = 2^{-3J} \varepsilon^{\frac{J}{5}} s$$
, $\Phi = \Phi^{(1)} \circ \cdots \circ \Phi^{(J)}$, and
$$\mathcal{N}_+ = e_+ + \langle \omega_+, I \rangle + \langle \Omega_+ q, \bar{q} \rangle,$$

with $\Omega_+ = T + A + W_+$, and

$$e_{+} = e + \sum_{j=1}^{J} e^{(j)}, \quad \omega_{+} = \omega + \sum_{j=1}^{J} \omega^{(j)}, \quad W_{+} = W + \sum_{j=1}^{J} W^{(j)}.$$

Then $\Phi: \mathcal{D}_+ \times \mathcal{O}_+ \to \mathcal{D} \times \mathcal{O}$. From the estimates of $\omega^{(j)}$ and $W^{(j)}$, we have

$$|\omega_{+} - \omega|_{\mathcal{O}_{+}} \le c\varepsilon,\tag{4.59}$$

$$|(W_{+} - W)_{mn}|_{\mathcal{O}_{+}} \le \begin{cases} \varepsilon^{\frac{1}{2}} e^{-\frac{\rho}{2} \max\{|m|, |n|\}}, & |m|, |n| \le N_{+}, & m, n \notin \mathcal{J} \\ 0, & \text{otherwise} \end{cases}$$
 (4.60)

Since $W^* = W$ and $(W^{(i)})^* = W^{(i)}$, W_+ is still a Hermitian matrix. Then, by (4.59) and (4.60), (A1) and (A2) hold with $p_+ = p + \varepsilon^{\frac{1}{2}}$ and $\sigma_+ := \frac{1}{3}\rho$.

Let $P_+ = P^{(J)}$. It has been verified that the assumptions (A4) and (A5) for $P^{(J)}$ hold, which is an analogue to the process in subsection 4.3.

This completes one step of KAM iterations.

5 Proof of the KAM theorem

With $\varepsilon_0 = \epsilon^{\frac{1}{4}}$, $\sigma_0 = 1$, $\hat{N} = |\ln \varepsilon_0|$, and

$$M_0 = \max \left\{ 2^{\tilde{s}+4} C \frac{L^{\tilde{s}+1}((\tilde{s}+1)!)^2}{\tilde{\xi}}, \ 2\tilde{\tau}, \ 8, \ \frac{12(2\tau+b+3)}{\tilde{\tau}} \right\}, \quad N_0 = 6|\ln \varepsilon_0|, \quad \rho_0 = N_0^{-1},$$

one can define the following sequences as in [11],

$$\begin{split} M_{\nu+1} &= M_{\nu}^{\tilde{s}M_{\nu}^{3}}, \quad a_{\nu} = \frac{1}{\tilde{\tau}} M_{\nu}^{-3\tilde{s}M_{\nu}^{3}}, \quad \varepsilon_{\nu+1} = \varepsilon_{\nu}^{\frac{1}{2}\varepsilon_{\nu}^{-a_{\nu}/2}}, \\ N_{\nu+1} &= \varepsilon_{\nu}^{-a_{\nu}}, \qquad \rho_{\nu+1} = \varepsilon_{\nu}^{a_{\nu}}, \qquad \sigma_{\nu+1} = \frac{1}{3}\rho_{\nu}. \end{split}$$

Given $p_0 = \varepsilon_0^{\frac{1}{2}}$, $r_0 = r$, $s_0 = s$, the other sequences are defined as

$$p_{\nu+1} = p_{\nu} + \varepsilon_{\nu}^{\frac{1}{2}}, \qquad K_{\nu+1} = N_{\nu+1} - (M_{\nu} + 1)N_{\nu}, \quad J_{\nu} = \left[\frac{5}{2}\varepsilon_{\nu}^{-\frac{a_{\nu}}{2}}\right],$$

$$r_{\nu} = r_{0}\left(1 - \sum_{i=2}^{\nu+1} 2^{-i}\right), \quad s_{\nu+1} = 2^{-3J_{\nu}}\varepsilon_{\nu}^{\frac{J_{\nu}}{5}}s_{\nu}, \qquad \gamma_{\nu} = \varepsilon_{\nu}^{\frac{1}{80}}.$$

Let \mathcal{D}_{ν} and \mathcal{O}_{ν} be as defined in Section 4.

5.1 Iteration lemma

The preceding analysis can be summarized as follows.

Lemma 5.1 There exists ε_0 sufficiently small such that the following holds for all $\nu = 0, 1, \cdots$.

(a) $H_{\nu} = \mathcal{N}_{\nu} + \check{P} + P_{\nu}$ is real-analytic on \mathcal{D}_{ν} , and C_W^1 parametrized by $\xi \in \mathcal{O}_{\nu}$, where $\mathcal{N}_{\nu} = e_{\nu} + \langle \omega_{\nu}, I \rangle + \langle \Omega_{\nu} q, \bar{q} \rangle, \quad P_{\nu} = \sum_{\alpha, \beta} (P_{\nu})_{\alpha\beta}(\theta, I) q^{\alpha} \bar{q}^{\beta},$

with $\Omega_{\nu} = T + A + W_{\nu}$ satisfying

$$(\Omega_{\nu})_{mn} \equiv 0 \quad if \quad m \quad or \quad n \in \mathcal{J},$$

$$|(W_{\nu})_{mn}|_{\mathcal{O}_{\nu}} \leq \begin{cases} p_{\nu}e^{-\sigma_{\nu}\max\{|m|,|n|\}}, & |m|,|n| \leq N_{\nu} \\ 0, & otherwise \end{cases},$$

$$|\omega_{\nu+1} - \omega_{\nu}|_{\mathcal{O}_{\nu+1}} \leq \varepsilon_{\nu},$$

$$|(W_{\nu+1} - W_{\nu})_{mn}|_{\mathcal{O}_{\nu+1}} \leq \begin{cases} \varepsilon_{\nu}^{\frac{1}{2}}e^{-\frac{\rho_{\nu}}{2}\max\{|m|,|n|\}}, & |m|,|n| \leq N_{\nu+1}, & m,n \notin \mathcal{J} \\ 0, & otherwise \end{cases}.$$

Moreover, P_{ν} has gauge invariance and $\|X_{\check{P}+P_{\nu}}\|_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \leq \varepsilon_{\nu}$,

$$\begin{split} & \| (P_{\nu})_{\alpha\beta} \|_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \leq \left\{ \begin{array}{l} \varepsilon_{\nu} e^{-\rho_{\nu} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \leq 2 \\ e^{-\rho_{\nu} n_{\alpha\beta}^{*}}, & |\alpha| + |\beta| \geq 3 \end{array} \right., \\ & \partial_{q_{n}} P_{\nu} = \partial_{\bar{q}_{n}} P_{\nu} \equiv 0, \quad \forall n \in \mathcal{J}. \end{split}$$

(b) For each ν , there is a symplectic transformation $\Phi_{\nu}: \mathcal{D}_{\nu+1} \to \mathcal{D}_{\nu}$ with

$$||D\Phi_{\nu} - Id||_{\mathcal{D}_{\nu+1}, \mathcal{O}_{\nu+1}} \le \varepsilon_{\nu}^{\frac{1}{2}},$$

such that $H_{\nu+1} = H_{\nu} \circ \Phi_{\nu}$.

Proof: Let $c_0 := 8e^{20} \max\{c_1, \dots, c_8\}$. We need to verify the assumptions (C1) – (C8) for $\nu = 0, 1, \dots$. By noting that

$$N_{\nu+1} = \varepsilon_{\nu}^{a_{\nu}} = \rho_{\nu+1}^{-1}, \quad \sigma_{\nu+1} = \frac{1}{3}\rho_{\nu}, \quad r_{\nu}^{(j)} - r_{\nu}^{(j+1)} = \frac{r_{\nu} - r_{\nu+1}}{2J_{\nu}}, \quad \rho_{\nu}^{(j)} - \rho_{\nu}^{(j+1)} = \frac{\rho_{\nu} - \rho_{\nu+1}}{2J_{\nu}},$$

it is sufficient for us to check:

(D1) $c_0 s_{\nu} \leq \varepsilon_{\nu}$,

(D2)
$$c_0 \left(\frac{r_{\nu} - r_{\nu+1}}{2J_{\nu}}\right)^{-(2\tau + b + 1)} \left(\frac{\rho_{\nu} - \rho_{\nu+1}}{2J_{\nu}}\right)^{-2} \le \varepsilon_{\nu}^{-\frac{1}{20}},$$

(D3)
$$c_0 N_{\nu+1}^8 M_{\nu}^8 N_{\nu}^{20} e^{8M_{\nu}N_{\nu}\rho_{\nu}} \le \varepsilon_{\nu}^{-\frac{7}{40}},$$

(D4)
$$e^{-\frac{\rho_{\nu}K_{\nu+1}}{2J_{\nu}}} \le \varepsilon_{\nu}^{\frac{2}{5}},$$

for all $\nu = 0, 1, \cdots$.

By the choice of s_0 , the condition (**D1**) clearly holds for $\nu = 0$. Suppose that it holds for some ν , then it is easy to see that

$$c_0 s_{\nu+1} = 2^{-3J_{\nu}} \varepsilon_{\nu}^{\frac{J_{\nu}}{5}} \cdot c_0 s_{\nu} < 2^{-3J_{\nu}} \varepsilon_{\nu}^{\frac{J_{\nu}}{5}} \cdot \varepsilon_{\nu} < \varepsilon_{\nu+1}.$$

Hence (**D1**) holds for all ν .

Let us first take ε_0 sufficiently small such that

$$\varepsilon_0^{\frac{1}{20} - \frac{1}{2}a_0(2\tau + b + 3)} \leq \frac{1}{c_0} \left(\frac{r_0}{20}\right)^{2\tau + b + 1} \left(\frac{1 - \varepsilon_0^{a_0}}{5}\right)^2.$$

Here we have applied $M_0 \ge \frac{12}{\tilde{\tau}}(2\tau + b + 3)$ and $a_0 = M_0^{-3\tilde{s}M_0^3}$ such that $\frac{1}{20} - \frac{1}{2}a_0(2\tau + b + 3) > 0$. Then, recalling that $r_{\nu} - r_{\nu+1} = \frac{r_0}{2^{2+\nu}}$ and $J_{\nu} = \left[\frac{5}{2}\varepsilon_{\nu}^{-\frac{a_{\nu}}{2}}\right]$,

$$c_0 \left(\frac{r_0 - r_1}{2J_0} \right)^{-(2\tau + b + 1)} \left(\frac{\rho_0 - \rho_1}{2J_0} \right)^{-2} \le \varepsilon_0^{-\frac{1}{20}},$$

i.e., (**D2**) holds for $\nu = 0$. Since for $\nu \ge 1$ and for ε_0 sufficiently small,

$$\varepsilon_{\nu}^{\frac{1}{40} - \frac{1}{2}a_{\nu}(2\tau + b + 3)} \ll \varepsilon_{0}^{\left(\frac{6}{5}\right)^{\nu}} \ll \frac{1}{2^{\nu(2\tau + b + 1)}c_{0}} \left(\frac{r_{0}}{20}\right)^{2\tau + b + 1}, \quad \varepsilon_{\nu}^{\frac{1}{40}} \ll \left(\frac{\varepsilon_{\nu - 1}^{a_{\nu - 1}} - \varepsilon_{\nu}^{a_{\nu}}}{5}\right)^{2},$$

we have

$$c_0 \left(\frac{r_{\nu} - r_{\nu+1}}{2J_{\nu}} \right)^{-(2\tau + b + 1)} \left(\frac{\rho_{\nu} - \rho_{\nu+1}}{2J_{\nu}} \right)^{-2} \le \varepsilon_{\nu}^{-\frac{1}{20}}.$$

Thus, (**D2**) holds true.

In Section 6 of [11], the basic smallness assumption of ε_{ν} , i.e., the inequality (A.1) in Lemma A.1, has been verified, then all other assumptions are immediate, including the inequality

$$\Gamma_{\nu} N_{\nu}^2 e^{6M_{\nu}N_{\nu}\rho_{\nu}} \leq \varepsilon_{\nu}^{-\frac{1}{8}},$$

where Γ_{ν} increases superexponentially in M_{ν} . Since all of M_{ν} , N_{ν} , ρ_{ν} and ε_{ν} here are defined in the same way as [11], we can apply this inequality. So (**D3**) has been verified.

By the definition of ρ_{ν} , a_{ν} and ε_{ν} , we have

$$\rho_{\nu} \varepsilon_{\nu}^{-\frac{1}{2}a_{\nu}} > \ln \frac{1}{\varepsilon_{\nu}}.$$

Then we see that (**D4**) holds for $\nu = 0, 1, \cdots$

Convergence

Now we fix $x \in \mathcal{X}$, with \mathcal{X} defined as in Proposition 1. This means that the blocks mentioned in Proposition 1 are eventually stationary after some step, i.e., for each $n \in \mathbb{Z}$, there is a $\nu_0(n)$ such that

$$\Lambda^{\nu+1}(n) = \Lambda^{\nu}(n), \quad \forall \nu \ge \nu_0(n).$$

In this case, the local decay rate for n may not shrink with ν necessarily (ρ_{ν} is the global upper bound of the rates for all $n \in \mathbb{Z}$).

Define $\Psi^{\nu} = \Phi_0 \circ \Phi_1 \circ \cdots \circ \Phi_{\nu-1}, \ \nu = 1, 2, \cdots$. An induction argument shows that $\Psi^{\nu}: \mathcal{D}_{\nu+1} \to \mathcal{D}_0$, and

$$H_0 \circ \Psi^{\nu} = H_{\nu} = \mathcal{N}_{\nu} + \breve{P} + P_{\nu}.$$

Let $\mathcal{O}_{\varepsilon_0} = \bigcap_{\nu=0}^{\infty} \mathcal{O}_{\nu}$. As in standard arguments (e.g. [24, 29]), thanks to Lemma 4.2, it concludes that H_{ν} , \mathcal{N}_{ν} , P_{ν} , Ψ^{ν} , e_{ν} , ω_{ν} and W_{ν} converge uniformly on $\mathcal{D}_{d,0}(\frac{1}{2}r_0,0)\times\mathcal{O}_{\varepsilon_0}$ to, say, H_{∞} , \mathcal{N}_{∞} , P_{∞} , Ψ^{∞} , e_{∞} , ω_{∞} and W_{∞} respectively, in which case it is clear that

$$\mathcal{N}_{\infty} = e_{\infty} + \langle \omega_{\infty}, I \rangle + \langle (T + A + W_{\infty})q, \bar{q} \rangle,$$

with $\Omega_{\infty} = T + A + W_{\infty}$ satisfying $(\Omega_{\infty})_{mn} \equiv 0$ if m or $n \in \mathcal{J}$. Since $||X_{P_{\nu}}||_{\mathcal{D}_{\nu},\mathcal{O}_{\nu}} \leq \varepsilon_{\nu}$

with $\varepsilon_{\nu} \to 0$, it follows that $\|X_{P_{\infty}}\|_{\mathcal{D}_{d,0}(\frac{1}{2}r_0,0),\mathcal{O}_{\varepsilon_0}} = 0$. Since $H_0 \circ \Psi^{\nu} = H_{\nu}$, we have $\Phi^t_{H_0} \circ \Psi^{\nu} = \Psi^{\nu} \circ \Phi^t_{H_{\nu}}$, with $\Phi^t_{H_0}$ denoting the flow of the Hamiltonian vector field X_{H_0} . The uniform convergence of Ψ^{ν} and $X_{H_{\nu}}$ implies that one can pass the limit in the above and conclude that

$$\Phi^t_{H_0} \circ \Psi^\infty = \Psi^\infty \circ \Phi^t_{H_\infty}, \quad \Psi^\infty : \mathcal{D}_{d,0}(\frac{1}{2}r_0,0) \to \mathcal{D}_0.$$

Hence,

$$\Phi_{H_0}^t(\Psi^{\infty}(\mathbb{T}^b \times \{\xi\})) = \Psi^{\infty}\Phi_{\mathcal{N}_{\infty}}^t(\mathbb{T}^b \times \{\xi\}) = \Psi^{\infty}(\mathbb{T}^b \times \{\xi\}), \quad \forall \xi \in \mathcal{O}_{\varepsilon_0}.$$

This means that $\Psi^{\infty}(\mathbb{T}^b \times \{\xi\})$ is an embedded invariant torus of the original perturbed Hamiltonian system at $\xi \in \mathcal{O}_{\varepsilon_0}$. Moreover, the frequencies $\omega_{\infty}(\xi)$ associated with $\Psi^{\infty}(\mathbb{T}^b \times \mathbb{T}^b)$ $\{\xi\}$) are slightly deformed from the unperturbed ones, $\omega(\xi)$.

5.3Measure estimate

At the ν^{th} step of KAM iteration, we need to exclude the following resonant parameter set

$$\mathcal{R}_k^{\nu} := \mathcal{R}_k^{\nu 1} \bigcup \left(\bigcup_{n \in \Lambda^{\nu}} \mathcal{R}_{kn}^{\nu 2} \right) \bigcup \left(\bigcup_{m,n \in \Lambda^{\nu}} \mathcal{R}_{kmn}^{\nu 3} \right) \bigcup \left(\bigcup_{m,n \in \Lambda^{\nu}} \mathcal{R}_{kmn}^{\nu 4} \right), \quad k \neq 0$$

for any fixed $x \in \mathcal{X}$, where

$$\begin{split} \mathcal{R}_k^{\nu 1} &:= \left\{ \xi \in \mathcal{O}_{\nu} : |\langle k, \omega_{\nu} \rangle| < \frac{\gamma_{\nu}}{|k|^{\tau}} \right\}, \\ \mathcal{R}_{kn}^{\nu 2} &:= \left\{ \xi \in \mathcal{O}_{\nu} : |\langle k, \omega_{\nu} \rangle + \mu_n^{\nu}| < \frac{\gamma_{\nu}}{|k|^{\tau} N_{\nu+1}^2} \right\}, \end{split}$$

$$\mathcal{R}_{kmn}^{\nu 3} := \left\{ \xi \in \mathcal{O}_{\nu} : |\langle k, \omega_{\nu} \rangle + \mu_{m}^{\nu} + \mu_{n}^{\nu}| < \frac{\gamma_{\nu}}{|k|^{\tau} N_{\nu+1}^{4}} \right\},
\mathcal{R}_{kmn}^{\nu 4} := \left\{ \xi \in \mathcal{O}_{\nu} : |\langle k, \omega_{\nu} \rangle + \mu_{m}^{\nu} - \mu_{n}^{\nu}| < \frac{\gamma_{\nu}}{|k|^{\tau} N_{\nu+1}^{4}} \right\},$$

with $\{\mu_i^{\nu}\}_{j\in\Lambda^{\nu}}$ eigenvalues of $\tilde{D}_{\Lambda^{\nu}}^{\nu}+\tilde{A}_{\nu}$. It is clear $\mathcal{O}_0\setminus\mathcal{O}_{\varepsilon_0}\subseteq\bigcup_{\nu>0}\bigcup_{k\neq 0}\mathcal{R}_k^{\nu}$.

Recalling that ω_0 is a diffeomorphism of ξ , together with the estimates in (4.13), (4.59) and (4.60), we have

$$|\partial_{\xi}(\langle k, \omega_{\nu} \rangle + \mu_{m}^{\nu} - \mu_{n}^{\nu})| \ge |\partial_{\xi}\langle k, \omega_{0} \rangle| - \varepsilon_{0}^{\frac{1}{4}}|k| - p = O(|k|)$$

for the set $\mathcal{R}_{kmn}^{\nu 4}$. The cases for $\mathcal{R}_{k}^{\nu 1}$, $\mathcal{R}_{kn}^{\nu 2}$, $\mathcal{R}_{kmn}^{\nu 3}$ can be handled in an entirely analogous way. Thus

$$\left| \mathcal{R}_k^{\nu 1} \bigcup \left(\bigcup_{n \in \Lambda^{\nu}} \mathcal{R}_{kn}^{\nu 2} \right) \bigcup \left(\bigcup_{m,n \in \Lambda^{\nu}} \mathcal{R}_{kmn}^{\nu 3} \right) \bigcup \left(\bigcup_{m,n \in \Lambda^{\nu}} \mathcal{R}_{kmn}^{\nu 4} \right) \right| \leq \frac{c \gamma_{\nu}}{|k|^{\tau + 1}}.$$

Since $\tau \geq b$, we have that

$$|\mathcal{O}_0 \setminus \mathcal{O}_{\varepsilon_0}| \le \left| \bigcup_{\nu \ge 0} \bigcup_{k \ne 0} \mathcal{R}_k^{\nu} \right| \le c \sum_{\nu \ge 0} \sum_{k \ne 0} \frac{\gamma_{\nu}}{|k|^{\tau+1}} = c \sum_{\nu \ge 0} \gamma_{\nu} \sim \gamma_0 = \varepsilon_0^{\frac{1}{80}}.$$

A Appendix

A.1 Outline of the proof of Proposition 1

For any smooth function f defined on $\mathcal{I} \subset \mathbb{R}/\mathbb{Z}$, let $|f|_{C^j} := \max_{0 \le k \le j} \sup_{x \in \mathcal{I}} \frac{1}{k!} |\partial_x^k f(x)|$.

The operator T in (2.1) can be viewed as a sum of two infinite-dimensional matrices, i.e., $\operatorname{diag}\{V(x+n\tilde{\alpha})\}_{n\in\mathbb{Z}}+\epsilon\Delta$ with Δ denoting the discrete Laplacian. It is natural to define an abstract normal form containing $\operatorname{diag}\{V(x+n\tilde{\alpha})\}_{n\in\mathbb{Z}}$.

Definition A.1 Given a symmetric matrix D, smoothly parametrized by $x \in \mathbb{R}/\mathbb{Z}$ and satisfying the shift condition

$$D_{m+k,n+k}(x) = D_{mn}(x + k\tilde{\alpha}), \quad \forall k \in \mathbb{Z},$$

where $\tilde{\alpha}$ is a Diophantine number, i.e., for some $\tilde{\gamma} > 0$ and $\tilde{\tau} > 1$,

$$|n\tilde{\alpha}|_1 \ge \frac{\tilde{\gamma}}{|n|^{\tilde{\tau}}}, \quad n \ne 0.$$

We say that D is in **normal form** if the following conditions hold.

(a) Short-range.

$$|D_{mn}|_{C^k} \le \begin{cases} Ce^{-\rho|m-n|}L^k, & |m-n| \le N \\ 0, & otherwise \end{cases}, \forall k \ge 0.$$

(b) Block diagonalization. Fix any $x_* \in \mathbb{R}/\mathbb{Z}$. There exist an interval \mathcal{I} centered in x_* , a disjoint decomposition $\bigcup_i \Lambda_i = \mathbb{Z}$ and a smooth orthogonal matrix Q on \mathcal{I} such that

- (b1) $\sharp \Lambda_j \leq M$ and $\operatorname{diam} \Lambda_j \leq MN$ for each j.
- (b2) the conjugated matrix $\tilde{D} = Q^*DQ$ is a product of commuting blocks

$$\prod_{j} \tilde{D}_{\Lambda_{j}}(x), \quad \forall x \in \mathcal{I}.$$

- (b3) $Q_{mn} \equiv 0$ if |m-n| > N. Moreover, for all m, $Q_{mn} \not\equiv 0$ for at most M different n.
- (b4) $|Q|_{C^k} \leq L^k$ for each $k \geq 0$.
- (c) Eigenvalues. There is a piecewise smooth function E(x) such that for each j,

$$\{E(x_* + n\tilde{\alpha})\}_{n \in \Lambda_i}$$
 are the eigenvalues of $\tilde{D}_{\Lambda_i}(x_*)$,

and there are sets $\Omega_j \supset \Lambda_j$ such that

(c1) for each n, if $\inf_{l \in \Lambda_j} |E(x_* + l\tilde{\alpha}) - E(x_* + n\tilde{\alpha})| < \kappa$, then

$$x_* + n\tilde{\alpha} \in x_* + m\tilde{\alpha} + \frac{1}{2}(\mathcal{I} - x_*)$$
 for some $m \in \Omega_j$,

$$Q(x)(\mathbb{R}^{\Lambda(n)}) \subset \mathbb{R}^{\Omega_j + n - m}, \quad \forall x \in \mathcal{I}.$$

(c2) the resultant

$$u_{\Omega_j}(\varphi, x) = \operatorname{Res}\left(\det(D(x + \varphi)_{\Omega_j} - tI_{\Omega_j}), \det(D(x)_{\Omega_j} - tI_{\Omega_j})\right)^{5}$$

satisfies

$$|u_{\Omega_j}|_{C^k} < (4MC)^{2M^2} B^k, \quad \forall k \le \tilde{s} M^2 + 1,^6$$

$$\max_{0 \le k \le \tilde{s} M^2} \left| \frac{1}{\nu! B^k} \partial_{\varphi}^k u_{\Omega_j}(\varphi, x) \right| \ge \vartheta, \quad \forall \varphi, \, \forall x \in \mathbb{R}/\mathbb{Z}.$$

- (c3) $\sharp \Omega_j \leq M$ and $\operatorname{diam} \Omega_j \leq \left(\frac{1}{\lambda}\right)^{\tilde{\tau}+2}$.
- (c4) the intervals $\{n\tilde{\alpha} + \mathcal{I}\}_{\text{dist}(n,\Omega_j) < N}$ are pairwise disjoint.
- (c5) for each $\varphi \in \mathcal{I}$, $u_{\Omega_j}(\varphi, x)$ satisfies

$$|u_{\Omega_j}|_{C^k} < (2MC)^{2M^2} B^k, \quad \forall k \le \tilde{s}M^2 + 1,^7$$

$$\max_{0 \leq k \leq \tilde{s}M^2} \left| \frac{1}{\nu!B^k} \partial_x^k u_{\Omega_j}(\varphi, x) \right| \geq \vartheta \left(\prod_{m,n \in \Omega_j} |\varphi + (m-n)\tilde{\alpha}|_1 \right), \quad \forall x \in \mathbb{R}/\mathbb{Z}.$$

⁵The resultant of two monic polynomials P and Q is defined as the product $\operatorname{Res}(P,Q) = \prod_{\substack{P(x)=0 \ Q(y)=0}} (x-y)$.

⁶The norm is with respect to the variable φ .

⁷The norm is with respect to the variable x.

Remark A.1 Condition (a) implies an estimate of D in the operator norm on $\ell^2(\mathbb{Z}) \to \ell^2(\mathbb{Z})$:

$$||D||_{C^k} \le C \frac{e^{\rho} + 1}{e^{\rho} - 1} L^k \le C \frac{4}{\rho} L^k, \quad \forall k \ge 0,$$

if $\rho \leq 1$.

Lemma A.1 (The inductive lemma of [11]) Let D be in normal form on an interval $\mathcal{I} \subset \mathbb{R}/\mathbb{Z}$ with parameters C, L, ρ , M, N, κ , B, ϑ , λ , and let a < g < h be numbers restricted by

$$\frac{1}{\tilde{\tau}M^{3\tilde{s}M^3}} \leq a < \frac{g}{20\tilde{s}\tilde{\tau}M^4} < \frac{h}{100\tilde{s}^2\tilde{\tau}M^8}, \quad h \leq \frac{1}{5\tilde{s}M^{2\tilde{s}M^3}}.$$

Assume, as simplification, that

$$1 \le B \le L$$
, $M \ge 8$, $1 < C < 2$, $\rho, \kappa, \vartheta \le 1$.

Let Z be a symmetric matrix, smoothly parametrized on \mathbb{R}/\mathbb{Z} , satisfying the shift condition. Assume that

$$\lambda \le |\mathcal{I}| \le \vartheta/B,$$
$$|Z_{mn}|_{C^k} < \varepsilon e^{-\rho|m-n|} L^k, \quad k \ge 0.$$

If there is a constant $\Gamma = \Gamma(\tilde{\gamma}, \tilde{\tau}, \tilde{s}, M)$, super-exponentially decaying in M, such that

$$|\varepsilon| < \Gamma \left[\frac{\rho^{\tilde{\tau}} \kappa \vartheta \lambda^{\tilde{\tau}^2}}{LN^{\tilde{\tau}}} e^{-N\rho} \right]^{e^{e^{\tilde{s}}M^4}},$$
 (A.1)

then there is a smooth orthogonal matrix \tilde{U} , satisfying the shift condition, such that

$$|(\tilde{U}-I)_{mn}|_{C^k} < \varepsilon^{\frac{1}{2}} e^{-\rho'|m-n|} L'^k$$

and

$$\tilde{U}^*(D+Z)\tilde{U} = D' + Z',$$

with Z' a symmetric matrix, smoothly parametrized on \mathbb{R}/\mathbb{Z} , satisfying the shift condition, and D' in normal form on an interval $\mathcal{I}' \subset \mathcal{I}$, with parameters

$$\begin{split} C' &= (1+\varepsilon^{\frac{1}{2}})C, \quad L' = \varepsilon^{-h}L, \qquad \rho' = \frac{1}{2}\rho, \\ \lambda' &= 9^{-M'}\lambda, \qquad M' = M^{\tilde{s}M^3}, \quad N' = \varepsilon^{-a}, \\ \kappa' &= \varepsilon^h, \qquad B' = L, \qquad \vartheta' = \varepsilon^g L, \end{split}$$

and

$$2\lambda' \le |\mathcal{I}'| \le \varepsilon^g,$$
$$|Z'_{mn}|_{C^k} < \varepsilon^{\frac{1}{2}\varepsilon^{-a/2}} e^{-\rho'|m-n|} L'^k.$$

In addition,

$$|E(x_* + m\tilde{\alpha}) - E(x_* + n\tilde{\alpha})| < M' \frac{L}{\rho} \varepsilon^g, \quad \forall m \in \Lambda'(n),$$

$$Q'(x)(\mathbb{R}^{\Lambda'(n)}) \subset \sum_{m \in \Lambda'(n)} Q(x)(\mathbb{R}^{\Lambda(m)}), \quad \forall x \in \mathcal{I}',$$

D' is in normal form with the same parameters also on $x_* + \frac{1}{2}(\mathcal{I}' - x_*)$.

Finally, if $M \geq 2\tilde{\tau}$ then the closure of the sets

$$\{x_* + m\tilde{\alpha} : |E(x_* + m\tilde{\alpha}) - E(x_* + (m+n)\tilde{\alpha})| < 2M' \frac{L}{\rho} \varepsilon^g\}, \quad \forall \, 4(1/\lambda)^{\tilde{\tau}+2} < |n| < M'N',$$

$$\{x_* + m\tilde{\alpha} : |E(x_* + m\tilde{\alpha}) - E(x_* + (m+n)\tilde{\alpha})| < 2\varepsilon^{\frac{1}{8}}\}, \ \forall M'N' < |n| < 4(1/\lambda')^{\tilde{\tau}+2}$$

are unions of, respectively, at most $\varepsilon^{-\frac{g}{5\tilde{s}M^2}}$ and ε^{-M^4g} many components, each component being of length, respectively, at most $\varepsilon^{\frac{g}{4\tilde{s}M^2}}$ and ε^{2M^4g} .

For the detail of proof, which contains the construction of new blocks Λ'_i , i.e., the new equivalence relation on \mathbb{Z} , and the new orthogonal transformation Q', see Section 5 of Reference [11].

Recall that (1.5)-(1.8) have defined quantities $\tilde{\gamma}$, $\tilde{\tau}$, C, L, \tilde{s} , $\tilde{\xi}$ associated with the Diophantine number $\tilde{\alpha}$ and the nonconstant real-analytic function V. It has been proved by Eliasson in Section 6 of [11] that $D_0 = \text{diag}\{V(x+n\tilde{\alpha})\}_{n\in\mathbb{Z}}$ is in normal form with $C_0 = C$, $L_0 = L$, any

$$M_0 \ge \max \left\{ 2^{\tilde{s}+4} C \frac{L^{\tilde{s}+1}((\tilde{s}+1)!)^2}{\tilde{\xi}}, 2\tilde{\tau}, 8 \right\}, \quad N_0 \ge 1, \quad \rho = N_0^{-1},$$

and other suitable parameters κ_0 , B_0 , λ_0 , ϑ_0 . With $\varepsilon_0 = \epsilon^{\frac{1}{4}}$, $Z_0 = \epsilon \Delta$ satisfies

$$|(Z_0)_{mn}|_{C^k} < \varepsilon_0 e^{-\rho_0|m-n|} L_0^k.$$

For $\nu = 0, 1, 2, \dots$, let $M_{\nu+1} = M_{\nu}^{\tilde{s}M_{\nu}^3}$, and

$$a_{\nu} = \frac{1}{\tilde{\tau}} \left(\frac{1}{M_{\nu}} \right)^{3\tilde{s}M_{\nu}^{3}}, \quad g_{\nu} = 20\tilde{s}\tilde{\tau}M_{\nu}^{4}a_{\nu}, \quad h_{\nu} = \frac{1}{5\tilde{s}} \left(\frac{1}{M_{\nu}} \right)^{2\tilde{s}M_{\nu}^{3}}.$$

The other sequences can be defined as

$$\begin{split} \varepsilon_{\nu+1} &= \varepsilon_{\nu}^{\frac{1}{2}\varepsilon_{\nu}^{-a_{\nu}/2}}, \quad C_{\nu+1} &= (1+\varepsilon_{\nu}^{1/2})C_{\nu}, \quad L_{\nu+1} &= \varepsilon_{\nu}^{-h_{\nu}}L_{\nu}, \\ N_{\nu+1} &= \varepsilon_{\nu}^{-a_{\nu}}, \qquad \rho_{\nu+1} &= \varepsilon_{\nu}^{a_{\nu}}, \qquad \kappa_{\nu+1} &= \varepsilon_{\nu}^{h_{\nu}}, \\ B_{\nu+1} &= L_{\nu}, \qquad \lambda_{\nu+1} &= 9^{-M_{\nu}}\varepsilon_{\nu}^{g_{\nu}}, \qquad \vartheta_{\nu+1} &= \varepsilon_{\nu}^{g_{\nu}}L_{\nu}. \end{split}$$

The inequality (A.1), about parameters at the ν^{th} step, has been verified in Section 6 of [11], so we can apply Lemma A.1 iteratively. For each $\nu \geq 0$, there is an orthogonal matrix \tilde{U}_{ν} satisfying the shift condition, such that

$$|(\tilde{U}_{\nu} - I_{\mathbb{Z}})_{mn}|_{C^k} < \varepsilon_{\nu}^{\frac{1}{2}} e^{-\frac{\rho_{\nu}}{2}|m-n|} L_{\nu+1}^k$$

and

$$(\tilde{U}_0 \cdots \tilde{U}_{\nu})^* (D_0 + Z_0) (\tilde{U}_0 \cdots \tilde{U}_{\nu}) = D_{\nu+1} + Z_{\nu+1},$$

where $D_{\nu+1}$ is in normal form with parameters $C_{\nu+1}$, $L_{\nu+1}$, $\rho_{\nu+1}$, $M_{\nu+1}$, $N_{\nu+1}$, $\kappa_{\nu+1}$, $B_{\nu+1}$, $\vartheta_{\nu+1}$, $\lambda_{\nu+1}$, and $Z_{\nu+1}$ is a symmetric matrix, smoothly parametrized on \mathbb{R}/\mathbb{Z} , satisfying the shift condition and

$$|(Z_{\nu+1})_{mn}|_{C^k} \le \varepsilon_{\nu+1} e^{-\rho_{\nu+1}|m-n|} L_{\nu+1}^k.$$

Hence, in the operator norm $\|\cdot\|_{C^k}$,

$$\tilde{U}_0 \cdots \tilde{U}_{\nu} \to U, \quad Z_{\nu} \to 0, \quad D_{\nu} \to D_{\infty}.$$

Let $U_{\nu+1} = \tilde{U}_0 \cdots \tilde{U}_{\nu}$, by a simple calculation, we have

$$|(U_{\nu+1} - I_{\mathbb{Z}})_{mn}|_{C^k} < \varepsilon_0^{\frac{1}{2}} e^{-\frac{\rho_{\nu}}{2}|m-n|} L_{\nu+1}^k.$$

Clearly there is a uniform limit $E^{\nu}(x) \to E^{\infty}(x)$ which describes the spectrum of $D_{\infty}(x)$ -it is the closure of the image of E^{∞} . Consider now the closure S_{ν} of the set of all x such that

$$|E_{\infty}(x) - E_{\infty}(x + n\tilde{\alpha})| < \frac{3}{2}M_{\nu+1}\frac{L_{\nu}}{\rho_{\nu}}\varepsilon_{\nu}^{g_{\nu}}$$
 for some $4(1/\lambda_{\nu})^{\tilde{\tau}+2} < |n| < M_{\nu+1}N_{\nu+1}$

or

$$|E_{\infty}(x) - E_{\infty}(x + n\tilde{\alpha})| < \frac{3}{2} \varepsilon_{\nu}^{\frac{1}{8}}$$
 for some $M_{\nu+1} N_{\nu+1} < |n| < 4(1/\lambda_{\nu+1})^{\tilde{\tau}+2}$.

According to the final statement of Lemma A.1, this set is of measure less than $c\varepsilon_{\nu}^{g_{\nu}/20\tilde{s}M_{\nu}^2}$. By Borel-Cantelli Lemma, we conclude that there is a full-measure subset $\tilde{\mathcal{X}}$ of \mathbb{R}/\mathbb{Z} such that for any $x \in \tilde{\mathcal{X}}$, each $x + n\tilde{\alpha}$ will belong to only finitely many S_{ν} 's. Choose $x = x_*$ of this sort, i.e., for all $n \in \mathbb{Z}$ there is a $\nu_0(n)$ such that $x_* + n\tilde{\alpha} \notin S_{\nu}$ for $\nu \geq \nu_0(n)$. Hence for such ν 's,

$$|E^{\nu}(x_* + n\tilde{\alpha}) - E^{\nu}(x_* + n\tilde{\alpha} + m\tilde{\alpha})| \ge 2M_{\nu+1} \frac{L_{\nu}}{\rho_{\nu}} \varepsilon_{\nu}^{g_{\nu}}, \quad \forall 4(1/\lambda_{\nu})^{\tilde{\tau}+2} < |m| < M_{\nu+1}N_{\nu+1},$$

$$|E^{\nu}(x_* + n\tilde{\alpha}) - E^{\nu}(x_* + n\tilde{\alpha} + m\tilde{\alpha})| \ge 2\varepsilon_{\nu}^{\frac{1}{8}}, \quad \forall M_{\nu+1}N_{\nu+1} < |m| < 4(1/\lambda_{\nu+1})^{\tilde{\tau}+2}$$

This implies that $\Lambda^{\nu}(n) \subset [n - 4(1/\lambda_{\nu_0(n)})^{\tilde{\tau}+2}, n + 4(1/\lambda_{\nu_0(n)})^{\tilde{\tau}+2}]$ for $\nu \geq \nu_0(n)$. The blocks $\Lambda^{\nu}(n)$ therefore become eventually stationary:

$$\Lambda^{\nu+1}(n) = \Lambda^{\nu}(n), \quad \forall \nu \ge \nu_0(n).$$

A.2 Hamiltonian vector field and Poisson bracket

For $d, \rho, r, s > 0$, let F, G be two real-analytic functions on $\mathcal{D} = \mathcal{D}_{d,\rho}(r,s)$, both of which C_W^1 depend on the parameter $\xi \in \mathcal{O}$.

Lemma A.2 The norm $\|\cdot\|_{\mathcal{D},\mathcal{O}}$ has the Banach algebraic property, i.e.,

$$||FG||_{\mathcal{D},\mathcal{O}} \le ||F||_{\mathcal{D},\mathcal{O}} ||G||_{\mathcal{D},\mathcal{O}}.$$

Proof: Since
$$(FG)_{kl\alpha\beta} = \sum_{\substack{\check{k}+\hat{k}=k,\ \check{l}+\hat{l}=l\\ \check{\alpha}+\check{\alpha}=\alpha,\ \check{\beta}+\hat{\beta}=\beta}} F_{\check{k}\check{l}\check{\alpha}\check{\beta}} G_{\hat{k}\hat{l}\hat{\alpha}\hat{\beta}}$$
, we have that

$$\begin{split} \|FG\|_{\mathcal{D},\mathcal{O}} &= \sup_{\mathcal{D}} \sum_{k,l,\alpha,\beta} |(FG)_{kl\alpha\beta}|_{\mathcal{O}} |q^{\alpha}| |\bar{q}^{\beta}| |I^{l}| e^{|k||\operatorname{Im}\theta|} \\ &\leq \sup_{\mathcal{D}} \sum_{k,l,\alpha,\beta} \sum_{\substack{\tilde{k}+\tilde{k}=k,\ \tilde{l}+\tilde{l}=l\\ \tilde{\alpha}+\tilde{\alpha}=\alpha,\ \beta+\hat{\beta}=\beta}} |F_{\tilde{k}\tilde{l}\check{\alpha}\check{\beta}}G_{\hat{k}\hat{l}\hat{\alpha}\check{\beta}}|_{\mathcal{O}} |q^{\alpha}| |\bar{q}^{\beta}| |I^{l}| e^{(|\check{k}|+|\hat{k}|)|\operatorname{Im}\theta|} \\ &\leq \|F\|_{\mathcal{D},\mathcal{O}} \|G\|_{\mathcal{D},\mathcal{O}}. \end{split}$$

Lemma A.3 (Generalized Cauchy Inequalities) The various components of the Hamiltonian vector field X_F satisfy: for any 0 < r' < r, $0 < \rho' < \rho$,

$$\begin{split} &\|\partial_{\theta}F\|_{\mathcal{D}_{d,\rho}(r',s)} \leq \frac{c}{r-r'}\|F\|_{\mathcal{D}}, \\ &\|\partial_{I}F\|_{\mathcal{D}_{d,\rho}(r,\frac{s}{2})} \leq \frac{c}{s^{2}}\|F\|_{\mathcal{D}}, \\ &\sup_{\mathcal{D}_{d,\rho}(r,\frac{s}{2})} \sum_{n \in \mathbb{Z}} \left(\|\partial_{q_{n}}F\|_{\mathcal{O}} + \|\partial_{\bar{q}_{n}}F\|_{\mathcal{O}}\right) \langle n \rangle^{d} e^{|n|\rho'} \leq \frac{c}{s(\rho-\rho')}\|F\|_{\mathcal{D}}. \end{split}$$

Proof: We only prove the third inequality, with others shown analogously. Given $\omega \in \ell^1_{d,\rho}(\mathbb{Z}) \setminus \{0\}, \ f(t) = F(\cdot,\cdot,q+t\omega,\cdot)$ is an analytic function on the the complex disc $\{z \in \mathbb{C} : |z| < \frac{s}{\|\omega\|_{d,\rho}}\}$. Hence

$$|f'(0)| = \left| \sum_{n \in \mathbb{Z}} \omega_n \cdot \partial_{q_n} F \right| \le \frac{c}{s} ||F||_{\mathcal{D}} \cdot ||\omega||_{d,\rho},$$

by the usual Cauchy inequality. As a linear operator on $\ell^1_{d,\rho}(\mathbb{Z})$, $\partial_q F$ satisfies

$$\|\partial_q F\|_{\text{op}} := \sup_{\omega \neq 0} \frac{\left| \sum_{n \in \mathbb{Z}_1} \omega_n \cdot \partial_{q_n} F \right|}{\|\omega\|_{d,\rho}} \le \frac{c}{s} \|F\|_{\mathcal{D}}.$$

Let $\|\omega\|_{d,\rho} = \frac{s}{2}$, then

$$|\partial_{q_n} F| \leq \sup_{\|\omega\|_{d,\rho} = \frac{s}{2}} \frac{|\partial_{q_n} F| \cdot |\omega_n|}{\|\omega\|_{d,\rho}} \leq \frac{\|\partial_q F\|_{\mathrm{op}} |\omega_n|}{\frac{s}{2}} \leq \frac{c}{s} \|F\|_{\mathcal{D}} \langle n \rangle^{-d} e^{-|n|\rho}.$$

Hence, for any $0 < \rho' < \rho$,

$$\sum_{n\in\mathbb{Z}}|\partial_{q_n}F|\langle n\rangle^d e^{|n|\rho'}\leq \sum_{n\in\mathbb{Z}_1}\frac{c}{s}\|F\|_{\mathcal{D}}e^{-|n|(\rho-\rho')}\leq \frac{c}{s(\rho-\rho')}\|F\|_{\mathcal{D}}.$$

With $\tilde{F} = \sum_{k,l,\alpha,\beta} (\partial_{\xi} F_{kl\alpha\beta}) I^l e^{i\langle k,\theta \rangle} q^{\alpha} \bar{q}^{\beta}$, it can be proved similarly that

$$\sum_{n\in\mathbb{Z}} |\partial_{q_n} \tilde{F}| e^{|n|\rho'} \le \frac{c}{s(\rho - \rho')} ||F||_{\mathcal{D}}.$$

Since in the process above, $\xi \in \mathcal{O}$ and $(\theta, I, q, \bar{q}) \in \mathcal{D}_{d,\rho}(r, \frac{s}{2})$ are arbitrarily chosen, this inequality is proved.

Remark A.2 These inequalities can be seen as a generalization of the standard Cauchy estimates, which is similar to Lemma A.3 in [30].

Let $\{\cdot,\cdot\}$ denote the Poisson bracket of smooth functions, i.e.,

$$\{F,G\} = \langle \partial_I F, \partial_\theta G \rangle - \langle \partial_\theta F, \partial_I G \rangle + i \sum_{n \in \mathbb{Z}_1} (\partial_{q_n} F \cdot \partial_{\bar{q}_n} G - \partial_{\bar{q}_n} F \cdot \partial_{q_n} G).$$

Lemma A.4 If $||X_F||_{\mathcal{D}} < \varepsilon'$, $||X_G||_{\mathcal{D}} < \varepsilon''$, then

$$||X_{\{F,G\}}||_{\mathcal{D}_{d,\rho}(r-\sigma,\,\eta s)} < c\sigma^{-1}\eta^{-2}\varepsilon'\varepsilon'',$$

for any $0 < \sigma < r$ and $0 < \eta \ll 1$.

The proof is similar to that of Lemma 7.3 in [14].

Remark A.3 For more information about the norm $\|\cdot\|_{\mathcal{D},\mathcal{O}}$, see references [3, 5, 26].

Lemma A.5 If both of F and G have gauge invariance, then $\{F,G\}$ has gauge invariance.

Proof: F and G can be written as

$$F = \sum_{k,\alpha,\beta} F_{k\alpha\beta}(I;\xi) e^{i\langle k,\theta\rangle} q^{\alpha} \bar{q}^{\beta}, \quad G = \sum_{k,\alpha,\beta} G_{k\alpha\beta}(I;\xi) e^{i\langle k,\theta\rangle} q^{\alpha} \bar{q}^{\beta},$$

with $F_{k\alpha\beta} = G_{k\alpha\beta} \equiv 0$ if $\sum_{i=1}^{b} k_i + |\alpha| - |\beta| \neq 0$. By a simple calculation, we have

$$\{F,G\}_{k\alpha\beta} = i \sum_{\substack{\tilde{k}+\hat{k}=k\\ \tilde{\alpha}+\hat{\alpha}=\alpha\\ \hat{\beta}+\hat{\beta}=\beta}} \left(\langle \partial_I F_{\check{k}\check{\alpha}\check{\beta}}, \hat{k} \rangle G_{\hat{k}\hat{\alpha}\hat{\beta}} - \langle \check{k}, \partial_I G_{\hat{k}\hat{\alpha}\hat{\beta}} \rangle F_{\check{k}\check{\alpha}\check{\beta}} \right) \tag{A.2}$$

$$+\mathrm{i} \sum_{\substack{\tilde{k}+\hat{k}=k\\ \tilde{\alpha}+\hat{\alpha}=\alpha\\ \tilde{\beta}+\hat{\beta}=\beta}} \sum_{m\in\mathbb{Z}} \left(F_{\tilde{k}(\tilde{\alpha}+e_m)\tilde{\beta}} G_{\hat{k}\hat{\alpha}(\hat{\beta}+e_m)} - F_{\tilde{k}\tilde{\alpha}(\tilde{\beta}+e_m)} G_{\hat{k}(\hat{\alpha}+e_m)\hat{\beta}} \right). \tag{A.3}$$

Assume $\sum_{j=1}^{b} k_i + |\alpha| - |\beta| \neq 0$. Then, in the summation above, it is impossible that

$$\sum_{j=1}^{b} \check{k}_j + |\check{\alpha}| - |\check{\beta}| = \sum_{j=1}^{b} \hat{k}_j + |\hat{\alpha}| - |\hat{\beta}| = 0,$$

or

$$\sum_{j=1}^{b} \check{k}_j + |\check{\alpha} + e_m| - |\check{\beta}| = \sum_{j=1}^{b} \hat{k}_j + |\hat{\alpha}| - |\hat{\beta} + e_m| = 0,$$

$$\sum_{i=1}^{b} \check{k}_{j} + |\check{\alpha}| - |\check{\beta} + e_{m}| = \sum_{i=1}^{b} \hat{k}_{j} + |\hat{\alpha} + e_{m}| - |\hat{\beta}| = 0.$$

This means, in (A.2) and (A.3), each term $\equiv 0$. Thus Lemma A.5 is obtained.

Lemma A.6 If there exists $n_* \in \mathbb{Z}$ such that

$$\partial_{q_{n_*}} F = \partial_{\bar{q}_{n_*}} F = \partial_{q_{n_*}} G = \partial_{\bar{q}_{n_*}} G \equiv 0,$$

then $\partial_{q_{n_*}}\{F,G\} = \partial_{\bar{q}_{n_*}}\{F,G\} \equiv 0.$

Proof: Since

$$\begin{array}{lcl} \partial_{qn_*}\{F,G\} & = & \partial_{qn_*}\left(\langle\partial_IF,\,\partial_\theta G\rangle - \langle\partial_\theta F,\,\partial_IG\rangle + \mathrm{i}\sum_{m\in\mathbb{Z}}\left(\partial_{q_m}F\cdot\partial_{\bar{q}_m}G - \partial_{\bar{q}_m}F\cdot\partial_{q_m}G\right)\right) \\ & = & \left\langle\partial_I(\partial_{qn_*}F),\,\partial_\theta(\partial_{qn_*}G)\right\rangle - \left\langle\partial_\theta(\partial_{qn_*}F),\,\partial_I(\partial_{qn_*}G)\right\rangle \\ & & + \mathrm{i}\sum_{m\in\mathbb{Z}}\left(\partial_{q_m}(\partial_{qn_*}F)\cdot\partial_{\bar{q}_m}(\partial_{qn_*}G) - \partial_{\bar{q}_m}(\partial_{qn_*}F)\cdot\partial_{q_m}(\partial_{qn_*}G)\right) \\ & \equiv & 0 \end{array}$$

and similarly, $\partial_{\bar{q}_{n_*}}\{F,G\}\equiv 0$, this lemma is proved.

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