# PERSISTENCE OF HYPERBOLIC TORI IN HAMILTONIAN SYSTEMS 

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#### Abstract

We generalize the well-known result of Graff and Zehnder on the persistence of hyperbolic invariant tori in Hamiltonian systems by considering non-Floquet, frequency varying normal forms and allowing the degeneracy of the unperturbed frequencies. The preservation of part or full frequency components associated to the degree of non-degeneracy is considered. As applications, we consider the persistence problem of hyperbolic tori on a submanifold of a nearly integrable Hamiltonian system and the persistence problem of a fixed invariant hyperbolic torus in a non-integrable Hamiltonian system.


Dedicated to Professor George R. Sell on the occasion of his 65 th birthday

## 1. Introduction

In [14], Moser considered the following Hamiltonian system:

$$
\begin{equation*}
H=e+\left\langle\omega_{0}, y\right\rangle+\frac{1}{2}\langle y, A y\rangle+\frac{1}{2}\langle z, M z\rangle+P(x, y, z) \tag{1.1}
\end{equation*}
$$

where $(x, y, z) \in T^{n} \times R^{n} \times R^{2 m}, \omega_{0} \in R^{n}$ is a fixed Diophantine toral frequency, $A, M$ are $n \times n, 2 m \times 2 m$ non-singular, constant matrices, respectively, $J M$ is hyperbolic with all eigenvalues being real and distinct, and $P$ is a small perturbation. The persistence of the unperturbed Diophantine hyperbolic torus $T^{n} \times\{0\} \times\{0\}$ was shown along with the preservation of the toral frequency $\omega_{0}$. By considering a symplectic reduction of $M$ into the form

$$
M=\left(\begin{array}{cc}
O & B_{0}  \tag{1.2}\\
B_{0}^{\top} & O
\end{array}\right)
$$

Graff [9] generalized Moser's result by allowing multiple eigenvalues of $M$. An alternative proof of Graff's result was later given by Zehnder in [25] using implicit function technique. This result has played a fundamental role in analyzing global branches of invariant tori and Arnold diffusion in Hamiltonian systems ([10, 23]). For the Lindstedt series approach to the persistence of hyperbolic tori in Hamiltonian systems, we refer the readers to $[8,11]$.

The present paper concerns a generalization of the Graff-Zehnder result to the non-Floquet, frequency varying cases in which certain degeneracy of the unperturbed frequencies is allowed. More precisely, we consider the following Hamiltonian

[^0]systems
\[

$$
\begin{equation*}
H=e(\lambda)+\langle\omega(\lambda), y\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}(x, \lambda)\binom{y}{z}\right\rangle+h(x, y, z, \lambda)+P(x, y, z, \lambda) \tag{1.3}
\end{equation*}
$$

\]

where $(x, y, z) \in T^{n} \times R^{n} \times R^{2 m}, \lambda$ is a parameter in a bounded, closed, connected region $\mathcal{O} \subset R^{k}, e$ and $\omega$ are of class $C^{l_{0}}$ on $\mathcal{O}$ for some fixed $l_{0} \geq \max \{n, 2\}, \mathcal{M}$ is symmetric, real analytic in $x \in \mathcal{D}(r)=\left\{x \in C^{n} / Z^{n}:|\operatorname{Im} x|<r\right\}$ and $C^{l_{0}}$ in $\lambda \in \mathcal{O}$, $h(x, y, z, \lambda)=O\left(|(y, z)|^{3}\right)$ is real analytic, and, $P$, viewed as a perturbation, is real analytic in a complex neighborhood $D(r, s)=\{(x, y, z):|\operatorname{Im} x|<r,|y|<s,|z|<s\}$ of $T^{n} \times\{0\} \times\{0\}$ and $C^{l_{0}}$ in $\lambda \in \mathcal{O}$.

Throughout the paper, we shall use the same symbol $|\cdot|$ to denote the supnorm of vectors and its induced matrix norm, the standard $l_{1}$ norm in a lattice $Z^{p}$, absolute value of functions, and Lebesgue measure of sets etc. Thus, for any matrix $Q=\left(q_{i j}\right),|Q|=\max _{i} \sum_{j}\left|q_{i j}\right|$. Also, [•] will denote both the average of a matrix valued function on a torus and the integral part of a real number. For any (vector, matrix valued) function $f$ defined on a domain $D,|f|_{D}$ stands for $\sup _{D}|f|$, and, for any two complex column vectors $\xi, \zeta$ of the same dimension, $\langle\xi, \zeta\rangle$ means the transpose of $\xi$ times $\zeta$.

Write $\mathcal{M}$ in (1.3) into blocks:

$$
\mathcal{M}=\left(\begin{array}{cc}
A & B  \tag{1.4}\\
B^{\top} & M
\end{array}\right)
$$

where $A=A(x, \lambda), B=B(x, \lambda), M=M(x, \lambda)$ are $n \times n, n \times 2 m, 2 m \times 2 m$ blocks of $\mathcal{M}=\mathcal{M}(x, \lambda)$ respectively. With respect to the standard symplectic form

$$
\sum_{i=1}^{n} \mathrm{~d} x_{i} \wedge \mathrm{~d} y_{i}+\sum_{j=1}^{m} \mathrm{~d} z_{j} \wedge \mathrm{~d} z_{j+m}
$$

the equation of motion associated to (1.3) reads

$$
\left\{\begin{aligned}
\dot{x} & =\omega(\lambda)+A y+B z+\partial_{y} h+\partial_{y} P \\
\dot{y} & =-\frac{1}{2} \partial_{x}\left\langle\binom{ y}{z}, \mathcal{M}(x, \lambda)\binom{y}{z}\right\rangle-\partial_{x} h-\partial_{x} P \\
\dot{z} & =J M z+J B^{\top} y+J \partial_{z} h+J \partial_{z} P
\end{aligned}\right.
$$

where $J$ denotes the $2 m \times 2 m$ symplectic matrix. Thus, the unperturbed system associated to (1.3) admits a smooth family of invariant $n$-tori $T_{\lambda}=T^{n} \times\{0\} \times\{0\}$ with toral frequencies $\omega(\lambda)$ parameterized by $\lambda \in \mathcal{O}$. We first assume that $J[M]$ is hyperbolic on $\mathcal{O}$, i.e., if $\lambda_{i}(\lambda), i=1,2, \cdots, 2 m$, are eigenvalues of $J[M](\lambda)$, then
$\mathbf{H )}$ there is a $\sigma_{0}>0$ such that

$$
\left|\operatorname{Re} \lambda_{i}(\lambda)\right| \geq \sigma_{0}
$$

for all $\lambda \in \mathcal{O}$ and $i=1,2, \cdots, 2 m$.
We note that if both $|M-[M]|_{\mathcal{D}(r) \times \mathcal{O}}$ and $|B|_{\mathcal{D}(r) \times \mathcal{O}}$ are sufficiently small, then H) implies that the invariant tori $T_{\lambda}, \lambda \in \mathcal{O}$, are hyperbolic in the $z$-direction.

Next, we assume the Rüssmann condition on the frequency map, i.e.,
R)

$$
\max _{\lambda \in \mathcal{O}} \operatorname{rank}\left\{\partial^{\alpha} \omega(\lambda): \quad \forall|\alpha| \leq n-1\right\}=n
$$

The Rüssmann condition is known to be the weakest non-degenerate condition for the persistence of maximal dimensional invariant tori in nearly integrable Hamiltonian systems ([3, 19, 20, 22, 24]).

Define

$$
\eta_{0}=\frac{2}{\sqrt{\rho_{0}^{2}+4 \alpha_{0} \rho_{0}}+\rho_{0}}
$$

where,

$$
\begin{align*}
\alpha_{0} & =(1+2 m)\left|[M]^{-1}\right|_{\mathcal{O}} \\
\rho_{0} & =\frac{4 m}{\sigma_{0}}\left(1+\frac{2 m}{\sigma_{0}}|[M]|_{\mathcal{O}}\right)^{2 m-1} \tag{1.5}
\end{align*}
$$

Our main result is stated as follows.
Theorem 1. Consider (1.3). Assume conditions H), R) and that

$$
\begin{equation*}
|M-[M]|_{\mathcal{D}(r) \times \mathcal{O}},|B-[B]|_{\mathcal{D}(r) \times \mathcal{O}}<\eta_{0} \tag{1.6}
\end{equation*}
$$

Then there is a $\mu=\mu\left(r, s, l_{0}, \sigma_{0}\right)>0$ sufficiently small such that if

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} P\right|_{D(r, s) \times \mathcal{O}}<s^{2} \gamma^{3 l_{0}+4} \mu,|l| \leq l_{0} \tag{1.7}
\end{equation*}
$$

then there is a $0<r_{0}=r_{0}\left(r, \sigma_{0}\right) \leq r$ and a Cantor-like set $\mathcal{O}_{\gamma} \subset \mathcal{O}$, with $\left|\mathcal{O} \backslash \mathcal{O}_{\gamma}\right|=$ $O\left(\gamma^{\frac{1}{n_{*}-1}}\right)$, where $n_{*}=\max \{2, n\}$, for which the following holds. There is a $C^{l_{0}-1}$ Whitney smooth family of real analytic, symplectic transformations

$$
\Psi_{\lambda}: D\left(r_{0}, \frac{s}{2}\right) \rightarrow D\left(r_{0}, s\right), \quad \lambda \in \mathcal{O}_{\gamma}
$$

which are $C^{l_{0}}$ uniformly close to the identity such that

$$
H \circ \Psi_{\lambda}=e_{*}+\left\langle\Omega_{*}(\lambda), y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}_{*}(x, \lambda)\binom{y}{z}\right\rangle+h(x, y, z, \lambda)+P_{*}(x, y, z, \lambda),
$$

where

$$
\begin{aligned}
& \left|\partial_{\lambda}^{l} e_{*}-\partial_{\lambda}^{l} e\right|_{\mathcal{O}_{\gamma}}=O\left(\gamma^{2\left(l_{0}+1\right)} \mu^{\frac{1}{2}}\right) \\
& \left|\partial_{\lambda}^{l} \Omega_{*}-\partial_{\lambda}^{l} \omega\right|_{\mathcal{O}_{\gamma}}=O\left(\gamma^{2\left(l_{0}+1\right)} \mu^{\frac{1}{2}}\right) \\
& \left|\partial_{\lambda}^{l} \mathcal{M}_{*}-\partial_{\lambda}^{l} \mathcal{M}\right|_{\mathcal{D}\left(r_{0}\right) \times \mathcal{O}_{\gamma}}=O\left(\gamma^{l_{0}+1} \mu^{\frac{1}{4}}\right)
\end{aligned}
$$

for all $|l| \leq l_{0}$. Moreover,

$$
\left.\partial_{y}^{p} \partial_{z}^{q} P_{*}\right|_{(y, z)=(0,0)} \equiv 0,|p|+|q| \leq 2
$$

Thus, all unperturbed tori $T_{\lambda}$ with $\lambda \in \mathcal{O}_{\gamma}$ will persist and give rise to a $C^{l_{0}-1}$ Whitney smooth family of slightly deformed, analytic, quasi-periodic, invariant $n$ tori of the perturbed system with the Diophantine toral frequencies $\Omega_{*}(\lambda)$.

Our next result concerns the preservation of toral frequencies in connection with the degree of non-degeneracy of the matrix $[A]$. More precisely, we assume that
ND) there is an $1 \leq n_{0} \leq n$ such that both the $n_{0} \times n_{0}$ ordered principal block $U$ of $[A]$ and $Y \equiv[M]-[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right)[B]$ are non-singular on $\mathcal{O}$, where $O$ denotes the zero matrix.

It is clear that ND) holds automatically if $[A]$ is non-singular on $\mathcal{O}$ and $|[B]|_{\mathcal{O}}$ is sufficiently small (in particular when $[B] \equiv 0$ ).

Define

$$
\begin{equation*}
\eta=\frac{2}{\sqrt{\rho_{0}^{2}+4 \alpha \rho_{0}}+\rho_{0}} \tag{1.8}
\end{equation*}
$$

where $\rho_{0}$ is as in (1.5) and

$$
\begin{equation*}
\alpha=(1+2 m)\left(\left|Y^{-1}\right|+\left|U^{-1}\right|+\left(\left|Y^{-1}\right|\left|U^{-1}\right|\right)\left(1+|[B]|+|[B]|\left|U^{-1}\right|\right)\right)_{\mathcal{O}} \tag{1.9}
\end{equation*}
$$

Theorem 2. Consider (1.3). Assume conditions H), ND), the condition (1.6) with $\eta$ in place of $\eta_{0}$, and the condition (1.7) with respect to a sufficiently small positive number $\mu=\mu\left(r, s, l_{0}, \sigma_{0}, U\right)$.

1) If R$)$ holds, then there is a $r_{0}=r_{0}\left(r, \sigma_{0}, U\right)$ such that Theorem 1 holds with

$$
\left(\Omega_{*}(\lambda)\right)_{i}=\omega_{i}(\lambda), \quad \lambda \in \mathcal{O}_{\gamma}, \quad i=1,2, \cdots, n_{0}
$$

i.e., the first $n_{0}$ components of a perturbed toral frequency $\Omega_{*}(\lambda)$ coincide with the corresponding ones of the unperturbed toral frequency $\omega(\lambda)$.
2) If $n_{0}=n$, i.e., $U=[A]$ is non-singular on $\mathcal{O}$, then every hyperbolic Diophantine tori $T_{\lambda}$ with Diophantine type $(\gamma, \tau)$ for a fixed $\tau>n-1$ will persist with unchanged toral frequencies.

The roughness condition (1.6) in Theorem 1 and the similar roughness condition in Theorem 2 hold automatically if both $M$ and $B$ are small perturbations of $x$ independent matrices. Since both roughness conditions are independent of the size of the perturbation $P$ and there is no restriction on the smallness of $|A-[A]|$, the above theorems can be applied to Hamiltonians of form (1.3) which may be far from being integrable. In particular, if $M=[M], B=[B]$, then Theorem 1 holds for arbitrary $A(x, \lambda)$.

The above theorems extend the results of Graff and Zehnder even for the case of (1.1), by allowing the degeneracy of $A$, a general matrix $\mathcal{M}$, and the variation of toral frequencies. It is clear that if $B=0$ then Theorem 2 2) coincides with the results of Graff and Zehnder. We remark that, in the frequency varying case, a smoothly varying, symplectic reduction of (1.3) to (1.1), in particular with $M$ being reduced to the form (1.2), is generally not available especially when eigenvalues of $J[M](\omega)$ change their multiplicities.

The proof of the above theorems will use a linear iterative scheme which follows the traditional KAM framework but only deals with the elimination of all resonant terms in each KAM step.

The paper is organized as follows. In section 2, we consider applications of the above results to problems such as the persistence of hyperbolic lower dimensional tori on a submanifold in a partially integrable system and the persistence of a hyperbolic invariant torus in a non-integrable system. In Section 3, we describe the linear iterative scheme for one KAM step with respect to (1.3) and also give estimates for the symplectic transformation and the new Hamiltonian. In Section 4, we prove an iteration lemma which checks the validity of all KAM steps. Proofs of both theorems will be completed in Section 5 .

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## 2. Applications

2.1. Persistence of hyperbolic tori on submanifolds. A natural way to introduce parameters in the application of Theorems 1 and 2 is to consider persistence of hyperbolic tori on submanifolds in a partially integrable Hamiltonian system. To be more precise, let us consider the Hamiltonian

$$
\begin{equation*}
H(x, y, z)=H_{0}(y, z)+\varepsilon P(x, y, z), \quad x \in T^{n}, y \in R^{n}, z \in R^{2 m} \tag{2.1}
\end{equation*}
$$

endorsed with the standard symplectic structure. We assume that $H_{0}$ is partially integrable in the sense that the surface $\{z=0\}$ is invariant with respect to the Hamiltonian flow associated to $H_{0}$. Hence $H_{0}$ admits the following Taylor expansion

$$
H_{0}(y, z)=H_{0}(y, 0)+\frac{1}{2}\langle z, \tilde{M}(y) z\rangle+O\left(|z|^{3}\right)
$$

where $\tilde{M}(y)=\left(\partial^{2} H_{0} / \partial z^{2}\right)(y, 0)$. Now consider a $k$-dimensional $(1 \leq k \leq n)$ submanifold $\Xi$ of $R^{n}$ defined by

$$
\Xi=\{y=y(\lambda), z=0: \lambda \in \mathcal{O}\}
$$

where $\mathcal{O}$ is a closed bounded region in $R^{k}$ and $y: \mathcal{O} \rightarrow R^{n}$ is a smooth imbedding. Then with the translation $y-y(\lambda) \mapsto y, z=z$, the Hamiltonian $H_{0}$ becomes

$$
\begin{aligned}
H_{0}(y, z, \lambda)= & H_{0}(y(\lambda), 0)+\langle\tilde{\omega}(y(\lambda)), y\rangle \\
& +\frac{1}{2}\langle y, \tilde{A}(y(\lambda)) y\rangle+\frac{1}{2}\langle z, \tilde{M}(y(\lambda)) z\rangle+O\left(|z|^{3}\right)+O\left(|y||z|^{2}\right),
\end{aligned}
$$

where $\tilde{\omega}(y)=\left(\partial H_{0} / \partial y\right)(y, 0), \tilde{A}(y)=\left(\partial^{2} H_{0} / \partial y^{2}\right)(y, 0)$. Thus by letting $e(\lambda)=$ $H_{0}(y(\lambda), 0), \omega(\lambda)=\tilde{\omega}(y(\lambda)), A(\lambda)=\tilde{A}(y(\lambda)), M(\lambda)=\tilde{M}(y(\lambda))$, the Hamiltonian $H$ has the canonical form

$$
H=e(\lambda)+\langle\omega(\lambda), y\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}(\lambda)\binom{y}{z}\right\rangle+h(y, z, \lambda)+P(x, y, z, \lambda)
$$

with $\mathcal{M}=\operatorname{diag}(A, M)$, and,

$$
h(y, z, \lambda)=O\left(|z|^{3}\right)+O\left(|y||z|^{2}\right), P(x, y, z, \lambda)=\varepsilon P(y+y(\lambda), z, x)
$$

Since $P$ and its derivatives with respect to $\lambda$ will be sufficiently small in a complex neighborhood of $T^{n} \times\{0\} \times\{0\}$, an application of Theorems 1 and 2 yields the following.
Corollary 1. Consider (2.1) on the submanifold $\Xi$ and assume that $\tilde{M}$ is hyperbolic on $\Xi$ in the sense of H$)$.

1) If $\tilde{\omega}$ satisfies the Rüssmann condition on $\Xi$, i.e., $\omega$ satisfies $R$ ) on $\mathcal{O}$, then there is an $\varepsilon_{0}>0$ and a family of Cantor-like sets $\Xi_{\varepsilon} \subset \Xi, 0<\varepsilon \leq \varepsilon_{0}$, with $\left|\Xi \backslash \Xi_{\varepsilon}\right| \rightarrow 0$, as $\varepsilon \rightarrow 0$, such that for all $y \in \Xi_{\varepsilon}$, the unperturbed $n$-tori $T_{y}=\{y\} \times\{0\} \times T^{n}$ persist and give rise to a Whitney smooth family of hyperbolic, analytic, Diophantine n-tori $T_{y, \varepsilon}$ of the perturbed system.

Moreover, if for some $1 \leq n_{0} \leq n$, the $n_{0} \times n_{0}$ ordered principal block $U(y)$ of $\tilde{A}(y)$ is non-singular on $\Xi$, then the first $n_{0}$ components of the toral frequency of each $T_{y, \varepsilon}$ are the same as those of $T_{y}$.
2) If $\tilde{A}(y)$ itself is non-singular on $\Xi$, then all unperturbed hyperbolic Diophantine $n$-tori $T_{y}$ on $\Xi$ with Diophantine type $(\gamma, \tau)$, where $0<\gamma<\varepsilon^{\frac{1}{6 n+8}}$ is arbitrary and $\tau>n-1$ is fixed, will persist as $\varepsilon \rightarrow 0$ with unchanged toral frequencies.

To give an example, let us consider

$$
H(x, y, z)=\frac{1}{2} y_{1}^{2}+y_{2}+\frac{1}{2}\left(u^{2}-v^{2}\right)+\varepsilon P(x, y, z)
$$

where $x=\left(x_{1}, x_{2}\right)^{\top} \in T^{2}, y=\left(y_{1}, y_{2}\right)^{\top} \in R^{2}, z=(u, v)^{\top} \in R^{2}$, associated to the standard symplectic structure. Let $\Xi$ be the unit circle in $R^{2}$ which we parameterize by $\Xi=\{(\cos \lambda, \sin \lambda): 0 \leq \lambda<2 \pi\}$. With the notion above, it is clear that

$$
\begin{aligned}
\omega(\lambda) & =(\cos \lambda, 1)^{\top} \\
A(\lambda) & =\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) \\
M(\lambda) & =\left(\begin{array}{ll}
1 & 0 \\
0 & -1
\end{array}\right)
\end{aligned}
$$

Hence the condition H) is clearly satisfied on $\Xi$. Since

$$
\operatorname{rank}\left\{\partial^{\alpha} \omega(\lambda): \forall|\alpha| \leq 1\right\}=\operatorname{rank}\left(\begin{array}{ll}
\cos \lambda & -\sin \lambda \\
1 & 0
\end{array}\right)=2, \lambda \neq 0, \pi
$$

R) is also satisfied. An application of Corollary 2 yields the persistence of the majority of the unperturbed 2-tori on $\Xi$ after a small perturbation. Moreover, since the $1 \times 1$ principal block of $A$ equals 1 , the first component of the frequency of a perturbed torus $T_{\lambda, \varepsilon}$ is simply $\cos \lambda$.
2.2. Non-integrable KAM problem. Let $H_{0}(p, q),(p, q) \in R^{d} \times R^{d}$ be a real analytic Hamiltonian on $R^{2 d}$ endorsed with the standard symplectic structure. We assume that $H_{0}$ admits an analytic, invariant, Diophantine $n$-torus $T_{0}$ for some $1 \leq n<d$, i.e., there are real analytic functions $\hat{P}, \hat{Q}: T^{n} \rightarrow R^{d}$, and a Diophantine vector $\omega_{0} \in R^{n}$, such that $T_{0}=\operatorname{cl}\left\{\left(\hat{P}\left(\omega_{0} t\right), \hat{Q}\left(\omega_{0} t\right)\right): t \in R\right\}$. Let $x \in T^{n}$ be the standard coordinate on $T^{n}$. One can apply the Frobenius' theorem ([16]) to obtain a symplectic coordinate system $(x, y, z) \in T^{n} \times R^{n} \times R^{2 d-2 n}$ in the vicinity of $T_{0}$ such that $T_{0}=\{y=0, z=0\}$ and

$$
H_{0}=e+\left\langle\omega_{0}, y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}(x)\binom{y}{z}\right\rangle+O\left((|y|+|z|)^{3}\right) .
$$

Write $\mathcal{M}$ into blocks like (1.4) and assume that a) $[A]$ is non-singular; b) $[M]$ is hyperbolic in the sense of H$)$; c) $|M-[M]|,|B-[B]|$ are sufficiently small. Then an application of Theorem 2 2) yields the persistence as well as the frequency preservation of $T_{0}$ after a small perturbation of $H_{0}$. In the case that a smooth family of such invariant $n$-tori $T_{\lambda}$ of $H$ are given, one can show the persistence of the majority of invariant tori in the family under either the Rüssmann condition R) or the non-degeneracy of $\left[A_{\lambda}\right]$ according to Theorems 1 and 2.

The non-integrable persistence problem is even more significant when the persistence of a smooth family of maximal dimensional invariant tori $T_{\lambda} \simeq T^{d}, \lambda \in \mathcal{O}$, is considered for $H_{0}$, where $\mathcal{O} \subset R^{k}$ is a bounded closed region.

By introducing a symplectic coordinate $(x, y) \in D(r, s)=\{(x, y):|\operatorname{Im} x|<$ $r,|y|<s\}$ in the vicinity of the invariant tori, the perturbed Hamiltonian $H_{0}+\varepsilon P$ can be written in the form

$$
\begin{equation*}
H=e(\lambda)+\langle\omega(\lambda), y\rangle+\frac{1}{2}\langle y, A(x, \lambda) y\rangle+O\left(|y|^{3}\right)+\varepsilon P(x, y, \lambda) \tag{2.2}
\end{equation*}
$$

For a fixed Diophantine toral frequency $\omega=\omega_{0}$, persistence of an invariant $n$ torus of (2.2) for small $\varepsilon$ has been shown for the non-degenerate case (i.e., $[A]$ is a non-singular constant matrix) by Salamon and Zehnder ([21]) using the Lagrangian formalism and by Eliasson ([6]) using Lindstedt series. A generalization of these works is recently made in [5] without using action-angle variables.

In the frequency varying case with possible degeneracy of $\omega(\lambda)$, an immediate consequence of Theorems 1 and 2 is the following.

Corollary 2. Consider (2.2).

1) If $\omega(\lambda)$ satisfies the Rüssmann condition R ) for $n=d$, then there is a sufficiently small $\varepsilon_{0}>0$ and a family of Cantor-like sets $\mathcal{O}_{\varepsilon} \subset \mathcal{O}, 0<\varepsilon \leq$ $\varepsilon_{0}$, with $\left|\mathcal{O} \backslash \mathcal{O}_{\varepsilon}\right| \rightarrow 0$, as $\varepsilon \rightarrow 0$, such that for all $\lambda \in \mathcal{O}_{\varepsilon}$, the unperturbed $d$-tori $T_{\lambda}=T^{d} \times\{0\}$ persist and give rise to a Whitney smooth family of slightly deformed, analytic, Diophantine, invariant d-tori of the perturbed system.

Moreover, if for some $1 \leq d_{0} \leq d$, the $d_{0} \times d_{0}$ ordered principal block of $[A]$ is non-singular on $\mathcal{O}$, then the first $d_{0}$ components of the toral frequency of each perturbed torus coincide with those of the corresponding unperturbed toral frequency.
2) If $[A]$ itself is non-singular on $\mathcal{O}$, then all Diophantine tori $T_{\lambda}$ of Diophantine type $(\gamma, \tau)$, where $0<\gamma<\varepsilon^{\frac{1}{6 d+8}}$ is arbitrary and $\tau>n-1$ is fixed, will persist as $\varepsilon \rightarrow 0$ with unchanged toral frequencies.

Part 2) of Corollary 3 particularly holds when the persistence of a fixed Diophantine torus in a non-integrable Hamiltonian system is considered. Assume that the Hamiltonian $H_{0}$ admits a Diophantine invariant $d$-torus $T_{0}$ with toral frequency $\omega_{0}$. Let $(x, y) \in T^{d} \times R^{d}$ be a symplectic coordinate system in the vicinity of $T_{0}$ such that $T_{0}=\{y=0\}$. Then with respect to the new coordinate the Hamiltonian $H_{0}$ becomes

$$
H_{0}=e+\left\langle\omega_{0}, y\right\rangle+\frac{1}{2}\langle y, A(x) y\rangle+O\left(|y|^{3}\right)
$$

where $A(x)=\left(\partial^{2} H / \partial y^{2}\right)(x, 0)$. Applying part 2$)$ of Corollary 3 , it is clear that if $H_{0}$ is non-degenerate on $T_{0}$, i.e., $A$ is non-singular on $T^{d}$, then not only does $T_{0}$ persist under a small perturbation of $H_{0}$ but also the perturbed toral frequency is kept unchanged.

## 3. KAM Step

In this section, we describe our linear iterative scheme with respect to (1.3) for one KAM step, under the conditions of the first part of Theorem 2. As we shall see in the sequel, Theorem 1 and the second part of Theorem 2 can be more or less treated as special cases of the first part of Theorem 2 by taking $n_{0}=0$ and $n$ respectively (to unify the notation, $n_{0}=0$ means the omission of all $U$-related terms in the assumptions of Theorems 1 and 2). Below, we let $\tau>\max \{n(n-1)-1,0\}$ be fixed. Also, for simplicity, we set $l_{0}=n$.

Initially, set $e_{0}=e, \Omega_{0}=\omega, \mathcal{M}^{0}=\mathcal{M}, A^{0}=A, B^{0}=B, M^{0}=M, h_{0}=h$, $P_{0}=P, \mathcal{O}_{0}=\mathcal{O}, \beta_{0}=s, r_{*}=r, \gamma_{0}=\gamma, s_{0}=\gamma_{0}^{n+1+\frac{a_{0}}{2}} \mu^{\frac{1}{4}}, \mu_{0}=s^{2} \gamma_{0}^{n+1+\frac{a_{0}}{2}} \mu^{\frac{1}{2}}$,
$\mu_{*}=\mu$. We also write $\left[A^{0}\right](=[A])$ into blocks:

$$
\left[A^{0}\right]=\left(\begin{array}{ll}
U^{0} & E^{0} \\
\left(E^{0}\right)^{\top} & V^{0}
\end{array}\right)
$$

where $U^{0}=U$. Without loss of generality, assume that $0<r_{*}, \beta_{0}, \mu_{*} \leq 1, s_{0} \leq \beta_{0}$. By (1.7), we have that

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} P_{0}\right|_{D\left(r_{0}, s_{0}\right)} \leq \gamma_{0}^{n+1} s_{0}^{2} \mu_{0},|l| \leq n \tag{3.1}
\end{equation*}
$$

where $r_{0} \in\left(0, r_{*}\right]$ will be specified in Section 3.3.
Suppose at a KAM step, say the $\nu$ th step, we have arrived at a Hamiltonian

$$
\begin{align*}
& H=H_{\nu}=N+P  \tag{3.2}\\
& N=N_{\nu}(x, y, z, \lambda)=e+\langle\Omega(\lambda), y\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}(x, \lambda)\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda)
\end{align*}
$$

where $(x, y, z) \in D=D_{\nu}=D(r, s), 0<r=r_{\nu} \leq r_{0}, 0<s=s_{\nu} \leq s_{0}, \gamma=$ $\gamma_{\nu} \leq \gamma_{0}, \lambda \in \mathcal{O}=\mathcal{O}_{\nu} \subset \mathcal{O}_{0}, e(\lambda)=e_{\nu}(\lambda), \Omega(\lambda)=\Omega_{\nu}(\lambda)$ are smooth on $\mathcal{O}$ with $(\Omega(\lambda))_{i}=\omega_{i}(\lambda), 1 \leq i \leq n_{0}, \mathcal{M}(x, \lambda)=\mathcal{M}^{\nu}(x, \lambda)$ is real symmetric over $\mathcal{D} \times \mathcal{O}=\{x:|\operatorname{Im} x|<r\} \times \mathcal{O}$ which is smooth in $\lambda \in \mathcal{O}$ and real analytic in $x \in \mathcal{D}=\mathcal{D}_{\nu}=\mathcal{D}(r), P=P_{\nu}(x, y, z, \lambda)$ is real analytic in $(x, y, z) \in D$, smooth in $\lambda \in \mathcal{O}$, and moreover,

$$
\left|\partial_{\lambda}^{l} P\right|_{D \times \mathcal{O}} \leq \gamma^{n+1} s^{2} \mu,|l| \leq n
$$

for some $\mu=\mu_{\nu}>0$.
We shall construct a symplectic transformation $\Phi=\Phi_{\nu+1}$ which transforms the Hamiltonian (3.2), in smaller phase and frequency domains, to the desired Hamiltonian in the next KAM cycle (the $(\nu+1)$ th KAM step).

Below, for simplicity, quantities (domains, normal form, perturbation, etc.) in the next KAM cycle will be simply indexed by " + " $(=\nu+1)$ and we shall often suspend the dependence of functions on their arguments. Also, all constants $c_{1}-c_{7}$ in this section are positive and independent of the iteration process. We shall also use $c=c\left(r_{0}, \beta_{0}, l_{0}, \sigma_{0}\right)$ to denote any intermediate positive constant which is independent of the iteration process.

Let $b, \sigma, d$ be sufficiently small positive constants such that

$$
\begin{aligned}
& \sigma-(b+\sigma)(2 b+3 \sigma)>0, \delta(1+b+\sigma)>1 \\
& \left(n+1+\frac{a_{0}}{2}\right)(1-2 b-3 \sigma)>n+1
\end{aligned}
$$

where $\delta=1-d, a_{0} \in(0,1 / 3)$.

Define

$$
\begin{aligned}
\gamma_{+} & =\frac{\gamma_{0}}{4}+\frac{\gamma}{2} \\
r_{+} & =\delta r+d\left(1-\frac{\delta^{2}}{2}\right) r_{0}, \\
s_{+} & =s^{1+b+\sigma}, \\
\beta_{+} & =\frac{\beta}{2}+\frac{\beta_{0}}{4}, \\
K_{+} & =\left(\left[\log \frac{1}{s}\right]+1\right)^{3}, \\
D(a) & =D\left(r_{+}+\frac{6}{8}\left(r-r_{+}\right), a\right), a>0, \\
\mathcal{D}(a) & =\{x:|\operatorname{Im} x|<a\}, a>0, \\
\Gamma(a) & =e^{\frac{r_{0}(1-\delta) \delta^{2}}{16}} \sum_{0<|k| \leq K_{+}}|k|^{n+(n+1) \tau+6} e^{-|k| \frac{a}{8}}, a>0, \\
D_{+} & =D\left(r_{+}, s_{+}\right), \\
\mathcal{D}_{+} & =\mathcal{D}\left(r_{+}\right)=\left\{x:|\operatorname{Im} x|<r_{+}\right\}, \\
\tilde{D}_{+} & =D\left(r_{+}+\frac{5}{8}\left(r-r_{+}\right), \beta_{+}\right), \\
D_{i} & =D\left(r_{+}+\frac{i-1}{8}\left(r-r_{+}\right), i s_{+}\right), i=1,2, \cdots, 8 .
\end{aligned}
$$

3.1. Truncation. We express $P$ into Taylor-Fourier series

$$
P=\sum_{i \in Z_{+}^{n}, j \in Z_{+}^{2 m}, k \in Z^{n}} p_{k i j} y^{i} z^{j} e^{\sqrt{-1}\langle k, x\rangle}
$$

and consider the truncation

$$
\begin{align*}
& R= \sum_{|i|+|j|<3,|k| \leq K_{+}} p_{k i j} y^{i} z^{j} e^{\sqrt{-1}\langle k, x\rangle}=\sum_{|k| \leq K_{+}}\left(P_{k 00}+\left\langle P_{k 10}, y\right\rangle\right. \\
&\left.+\left\langle P_{k 01}, z\right\rangle+\left\langle y, P_{k 20} y\right\rangle+\left\langle y, P_{k 11} z\right\rangle+\left\langle z, P_{k 02} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} \tag{3.3}
\end{align*}
$$

Lemma 3.1. Assume that
H1) $s_{+} \leq \frac{s}{16}$;
H2) $\int_{K_{+}}^{\infty} \lambda^{n} e^{-\lambda \frac{r-r_{+}}{16}} \mathrm{~d} \lambda \leq s$.
Then there is a constant $c_{1}$ such that

$$
\left|\partial_{\lambda}^{l}(P-R)\right|_{D_{8}} \leq c_{1} \gamma^{n+1}\left(s^{3}+\frac{s_{+}^{3}}{s}\right) \mu,|l| \leq n
$$

Proof. Denote

$$
\begin{aligned}
I= & \sum_{|k|>K_{+}} p_{k i j} y^{i} z^{j} e^{\sqrt{-1}\langle k, x\rangle} \\
I I= & \sum_{|k| \leq K_{+},|i|+|j| \geq 3} p_{k i j} y^{i} z^{j} e^{\sqrt{-1}\langle k, x\rangle} \\
& =\int \frac{\partial^{(p, q)}}{\partial y^{p} \partial z^{q}} \sum_{|k| \leq K_{+},|i|+|j| \geq 3} p_{k i j} e^{\sqrt{-1}\langle k, x\rangle} y^{i} z^{j} \mathrm{~d} y \mathrm{~d} z \\
& =D\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right), s\right)
\end{aligned}
$$

where $\int$ is the obvious anti-derivative of $\frac{\partial^{(p, q)}}{\partial y^{p} \partial z^{q}}$ for $|p|+|q|=3$. Clearly,

$$
P-R=I+I I
$$

We note by H1) that $D_{8} \subset D_{*}$.
Since, by Cauchy's estimate,

$$
\left|\sum_{i \in Z_{+}^{n}, j \in Z_{+}^{2 m}} \partial_{\lambda}^{l} p_{k i j} y^{i} z^{j}\right| \leq\left|\partial_{\lambda}^{l} P\right|_{D(r, s)} e^{-|k| r} \leq \gamma^{n+1} s^{2} \mu e^{-|k| r},|l| \leq n
$$

$\mathrm{H} 2)$ implies that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} I\right|_{D_{*}} & \leq \sum_{|k|>K_{+}} \gamma^{n+1} s^{2} \mu e^{-|k| r} e^{|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \leq \gamma^{n+1} s^{2} \mu \sum_{\kappa=K_{+}}^{\infty} \kappa^{n} e^{-\kappa \frac{r-r_{+}}{8}} \leq \gamma^{n+1} s^{2} \mu \int_{K_{+}}^{\infty} \lambda^{n} e^{-\lambda \frac{r-r_{+}}{16}} \mathrm{~d} \lambda \\
& \leq \gamma^{n+1} s^{3} \mu,|l| \leq n
\end{aligned}
$$

It follows that

$$
\left|\partial_{\lambda}^{l}(P-I)\right|_{D_{*}} \leq\left|\partial_{\lambda}^{l} P\right|_{D(r, s)}+\left|\partial_{\lambda}^{l} I\right|_{D_{*}} \leq 2 \gamma^{n+1} s^{2} \mu,|l| \leq n
$$

By Cauchy's estimate of $\partial_{\lambda}^{l}(P-I)$ on $D_{*}$, we then have

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} I I\right|_{D_{8}} & \leq\left|\int \frac{\partial^{(p, q)}}{\partial y^{p} \partial z^{q}} \sum_{|k| \leq K_{+},|i|+|j| \geq 3} \partial_{\lambda}^{l} p_{k i j} e^{\sqrt{-1}\langle k, x\rangle} y^{i} z^{j} \mathrm{~d} y \mathrm{~d} z\right|_{D_{8}} \\
& \leq\left.\left.\left|\int\right| \frac{\partial^{(p, q)}}{\partial y^{p} \partial z^{q}} \partial_{\lambda}^{l}(P-I-R)\right|_{D_{*}} \mathrm{~d} y \mathrm{~d} z\right|_{D_{8}} \\
& \leq 2\left(\frac{1}{s-8 s_{+}}\right)^{3} \gamma^{n+1} s^{2} \mu\left|\int \mathrm{~d} y \mathrm{~d} z\right|_{D_{8}} \leq c \gamma^{n+1} \frac{s_{+}^{3}}{s} \mu,|l| \leq n .
\end{aligned}
$$

Thus,

$$
\left|\partial_{\lambda}^{l}(P-R)\right|_{D_{8}} \leq c \gamma^{n+1}\left(s^{3}+\frac{s_{+}^{3}}{s}\right) \mu,|l| \leq n
$$

3.2. The linear homological equations. Write $\mathcal{M}$ into blocks

$$
\mathcal{M}(x, \lambda)=\left(\begin{array}{cc}
A & B \\
B^{\top} & M
\end{array}\right)
$$

where

$$
\begin{aligned}
& A(x, \lambda)=\sum_{k \in Z^{n}} A_{k} e^{\sqrt{-1}\langle k, x\rangle} \\
& B(x, \lambda)=\sum_{k \in Z^{n}} B_{k} e^{\sqrt{-1}\langle k, x\rangle} \\
& M(x, \lambda)=\sum_{k \in Z^{n}} M_{k} e^{\sqrt{-1}\langle k, x\rangle}
\end{aligned}
$$

are $n \times n, n \times 2 m, 2 m \times 2 m$ blocks of $\mathcal{M}$ respectively.
To transform (3.2) into the Hamiltonian in the next KAM cycle, a symplectic transformation should at least eliminate all its first order resonant terms

$$
P_{k 00} e^{\sqrt{-1}\langle k, x\rangle},\left\langle P_{k 10}, y\right\rangle e^{\sqrt{-1}\langle k, x\rangle},\left\langle P_{k 01}, z\right\rangle e^{\sqrt{-1}\langle k, x\rangle},\left\langle P_{001}, z\right\rangle, 0<|k| \leq K_{+} .
$$

An essential idea of our linear iterative scheme is to find a Hamiltonian $F$ of the form

$$
\begin{equation*}
F=\sum_{0<|k| \leq K_{+}}\left(f_{k 0}+\left\langle f_{k 1}, y\right\rangle+\left\langle F_{k 1}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}+\left\langle F_{01}, z\right\rangle \tag{3.4}
\end{equation*}
$$

such that the time- 1 map $\phi_{F}^{1}$ of the flow generated by $F$, as a symplectic transformation, will precisely eliminate the above resonant terms. To be able to fix the first $n_{0}$ components of the toral frequencies as stated in Theorem 21 ), we shall also find a $Y_{*} \in R^{n_{0}}$ so that the translation of coordinate

$$
\phi: x \rightarrow x, y \rightarrow y+\binom{Y_{*}}{0}, z \rightarrow z
$$

removes all possible drifts among the first $n_{0}$ components of the new toral frequencies.

Denote

$$
\begin{align*}
{[A]=} & \left(\begin{array}{cc}
U & E \\
E^{\top} & V
\end{array}\right), \\
R^{\prime}= & {[R]+\sum_{0<|k| \leq K_{+}}\left(\left\langle y, P_{k 20} y\right\rangle+\left\langle y, P_{k 11} z\right\rangle+\left\langle z, P_{k 02} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} } \\
& -\left\langle P_{001}, z\right\rangle+\sum_{|k| \leq K_{+}}\left\langle B_{-k} J F_{k 1}, y\right\rangle,  \tag{3.5}\\
R_{t}= & (1-t)\{N, F\}+R,  \tag{3.6}\\
y_{*}= & \binom{Y_{*}}{0}
\end{align*}
$$

where $U, E, V$ are the $n_{0} \times n_{0}, n_{0} \times\left(n-n_{0}\right),\left(n-n_{0}\right) \times\left(n-n_{0}\right)$ blocks of $[A]$ respectively.

Let

$$
\Phi_{+}=\phi_{F}^{1} \circ \phi
$$

Then it is easy to see that

$$
\begin{aligned}
H_{+}= & H \circ \Phi_{+}=H \circ \phi_{F}^{1} \circ \phi=(N+R) \circ \phi_{F}^{1} \circ \phi+(P-R) \circ \phi_{F}^{1} \circ \phi \\
= & \left(N+R^{\prime}\right) \circ \phi-\left\langle y_{*},(A-[A]) y\right\rangle-\left\langle y_{*}, B z\right\rangle \\
& +\left(\{N, F\}+R-R^{\prime}\right) \circ \phi+\left\langle y_{*},(A-[A]) y\right\rangle+\left\langle y_{*}, B z\right\rangle \\
& +\int_{0}^{1}\left\{R_{t}, F\right\} \circ \phi_{F}^{t} \circ \phi \mathrm{~d} t+(P-R) \circ \phi_{F}^{1} \circ \phi .
\end{aligned}
$$

Since the Taylor-Fourier series of $R-R^{\prime}$ consists of terms of Fourier modes $e^{\sqrt{-1}\langle k, x\rangle}, 0<|k| \leq K_{+}$, but that of $\{N, F\}$ contains some high modes $e^{\sqrt{-1}\langle k, x\rangle}$ for $|k|>K_{+}$, we need to choose a function $Q$ of high order such that both equations

$$
\begin{align*}
& \left(\{N, F\}+R-R^{\prime}\right) \circ \phi-Q+\left\langle y_{*},(A-[A]) y\right\rangle+\left\langle y_{*}, B z\right\rangle=0,  \tag{3.7}\\
& \operatorname{diag}(U, O) y_{*}=\operatorname{diag}\left(I_{n_{0}}, O\right)\left(-P_{010}-\sum_{|j| \leq K_{+}} B_{-j} J F_{j 1}\right) \tag{3.8}
\end{align*}
$$

are solvable. If this is the case, then it is easy to see that

$$
\begin{aligned}
H_{+} & =N_{+}+P_{+} \\
N_{+} & =e_{+}+\left\langle\Omega_{+}(\lambda), y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}^{+}\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda) \\
& =e_{+}+\left\langle\Omega_{+}(\lambda), y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z},\left(\begin{array}{cc}
A^{+} & B^{+} \\
B^{+\top} & M^{+}
\end{array}\right)\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda),
\end{aligned}
$$

where

$$
\begin{align*}
& e_{+}=e+P_{000}+\left\langle\Omega, y_{*}\right\rangle,  \tag{3.9}\\
& \Omega_{+}=\Omega+\operatorname{diag}\left(O, I_{n-n_{0}}\right)\left([A] y_{*}+P_{010}+\sum_{|k| \leq K_{+}} B_{-k} J F_{k 1}\right),  \tag{3.10}\\
& A^{+}=A+\sum_{|k| \leq K_{+}} 2 P_{k 20} e^{\sqrt{-1}\langle k, x\rangle},  \tag{3.11}\\
& B^{+}=B+\sum_{|k| \leq K_{+}} P_{k 11} e^{\sqrt{-1}\langle k, x\rangle},  \tag{3.12}\\
& M^{+}=M+\sum_{|k| \leq K_{+}} 2 P_{k 02} e^{\sqrt{-1}\langle k, x\rangle},  \tag{3.13}\\
& P_{+}=\int_{0}^{1}\left\{R_{t}, F\right\} \circ \phi_{F}^{t} \circ \phi \mathrm{~d} t+(P-R) \circ \phi_{F}^{1} \circ \phi \\
& +\frac{1}{2}\left\langle y_{*},(A-2 \operatorname{diag}(U, O)) y_{*}\right\rangle+h_{0}\left(x, y+y_{*}, z, \lambda\right)-h_{0}(x, y, z, \lambda) \\
& +\sum_{|k| \leq K_{+}}\left(\left\langle y_{*}, P_{k 20} y_{*}\right\rangle+\left\langle y_{*}, 2 P_{k 20} y\right\rangle+\left\langle y_{*}, P_{k 11} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}+Q . \tag{3.14}
\end{align*}
$$

We now explore equations (3.7) and (3.8). Let

$$
\begin{aligned}
Q & =\left(\sum _ { 0 < | k | \leq K _ { + } } \left(\left\langle\frac{1}{2} \partial_{x}\langle y, A(x, \lambda) y\rangle+\partial_{x}\langle y, B(x, \lambda) z\rangle+\frac{1}{2} \partial_{x}\langle z, M(x, \lambda) z\rangle\right.\right.\right. \\
& \left.+\partial_{x} h_{0}(x, y, z, \lambda), f_{k 1}\right\rangle+(-\sqrt{-1}\langle k, A(x, \lambda) y+B(x, \lambda) z\rangle \\
& \left.\left.+\partial_{y} h_{0}(x, y, z, \lambda)\right)\left(f_{k 0}+\left\langle f_{k 1}, y\right\rangle+\left\langle F_{k 1}, z\right\rangle\right)\right) e^{\sqrt{-1}\langle k, x\rangle}
\end{aligned}
$$

$$
\begin{aligned}
& +\sum_{|k|>K_{+}}\left(\left\langle B_{k} J F_{01}, y\right\rangle+\left\langle M_{k} J F_{01}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} \\
& +\sum_{|k|>K_{+}, 0<|j| \leq K_{+}}\left(\left\langle B_{k-j} J F_{j 1}, y\right\rangle+\left\langle M_{k-j} J F_{j 1}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} \\
& \left.+\sum_{0<|k| \leq K_{+}}\left\langle\partial_{z} h_{0}(x, y, z, \lambda) J F_{k 1}, z\right\rangle e^{\sqrt{-1}\langle k, x\rangle}\right) \circ \phi \\
& +\sum_{0<|k| \leq K_{+}}\left(-\sqrt{-1}\left\langle f_{k 1}, y_{*}\right\rangle+\left\langle B^{\top}(x, \lambda) y_{*}, J F_{k 1}\right\rangle+\left\langle P_{k 10}, y_{*}\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} \\
& +\sum_{|k|>K_{+}}\left(\left\langle y_{*}, A_{k} y\right\rangle+\left\langle y_{*}, B_{k} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}+\left\langle B^{\top}(x, \lambda) y_{*}, J F_{01}\right\rangle .
\end{aligned}
$$

Substituting (3.3)-(3.5) and (3.15) into (3.7) yields

$$
\begin{aligned}
& -\sum_{0<|k| \leq K_{+}} \sqrt{-1}\langle k, \Omega(\lambda)\rangle\left(f_{k 0}+\left\langle f_{k 1}, y\right\rangle+\left\langle F_{k 1}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle} \\
& +\sum_{|j| \leq K_{+}}\left\langle z, M_{-j} J F_{j 1}\right\rangle+\sum_{0<|k| \leq K_{+},|j| \leq K_{+}}\left(\left\langle y, B_{k-j} J F_{j 1}\right\rangle\right. \\
& \left.+\left\langle z, M_{k-j} J F_{j 1}\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}+\sum_{0<|k| \leq K_{+}}\left(\left\langle y_{*}, A_{k} y\right\rangle+\left\langle y_{*}, B_{k} z\right\rangle+P_{k 00}+\left\langle P_{k 10}, y\right\rangle\right. \\
& \left.+\left\langle P_{k 01}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}+\left\langle P_{001}+[B]^{\top} y_{*}, z\right\rangle=0 .
\end{aligned}
$$

By comparing coefficients above, equations (3.7), (3.8) give rise to the following linear homological equations for all $0<|k| \leq K_{+}$:

$$
\begin{align*}
& \sqrt{-1}\langle k, \Omega(\lambda)\rangle f_{k 0}=P_{k 00},  \tag{3.16}\\
& \sqrt{-1}\langle k, \Omega(\lambda)\rangle f_{k 1}=P_{k 10}+A_{k} y_{*}+\sum_{|j| \leq K_{+}} B_{k-j} J F_{j 1},  \tag{3.17}\\
& \sqrt{-1}\langle k, \Omega(\lambda)\rangle F_{k 1}-[M] J F_{k 1}=\sum_{0<|j| \leq K_{+}, j \neq k} M_{k-j} J F_{j 1} \\
& \quad+P_{k 01}+B_{k}^{\top} y_{*}+M_{k} J F_{01},  \tag{3.18}\\
& {[M] J F_{01}=-P_{001}-\sum_{0<|j| \leq K_{+}} M_{-j} J F_{j 1}-[B]^{\top} y_{*},}  \tag{3.19}\\
& \operatorname{diag}(U, O) y_{*}=\operatorname{diag}\left(I_{n_{0}}, O\right)\left(-P_{010}-\sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}-[B] J F_{01}\right) . \tag{3.20}
\end{align*}
$$

Let

$$
\begin{align*}
\mathcal{O}_{+} & =\left\{\lambda \in \mathcal{O}:|\langle k, \Omega(\lambda)\rangle|>\frac{\gamma}{|k|^{\tau}}, 0<|k| \leq K_{+}\right\}  \tag{3.21}\\
Y & =[M] J-[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right)[B] J \\
Y^{0} & =\left[M^{0}\right] J-\left[B^{0}\right]^{\top} \operatorname{diag}\left(\left(U^{0}\right)^{-1}, O\right)\left[B^{0}\right] J
\end{align*}
$$

Then $\langle k, \Omega(\lambda)\rangle$ is invertible on $\mathcal{O}_{+}$for all $0<|k| \leq K_{+}$. If we assume that
H3) $\left|\partial_{\lambda}^{l}\left(\mathcal{M}-\mathcal{M}^{0}\right)\right|_{\mathcal{D}(r) \times \mathcal{O}} \leq \gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}},|l| \leq n$,
then as $\mu_{*}$ small, both $U$ and $Y$ are non-singular on $\mathcal{O}$. Thus, for all $0<|k| \leq K_{+}$, $\lambda \in \mathcal{O}_{+},(3.16)$ is immediately solvable, and, (3.17), (3.19) and (3.20) can be also solved in terms of $F_{k 1}, 0<|k| \leq K_{+}$. More precisely, we have

$$
\begin{align*}
& f_{k 0}=-\sqrt{-1}\langle k, \Omega(\lambda)\rangle^{-1} P_{k 00},  \tag{3.22}\\
& f_{k 1}=-\sqrt{-1}\langle k, \Omega(\lambda)\rangle^{-1}\left(P_{k 10}-A_{k} \operatorname{diag}\left(U^{-1}, O\right)\left(P_{010}\right.\right. \\
&\left.+\sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}\right)+\left(B_{k}-A_{k} \operatorname{diag}\left(U^{-1}, O\right)[B]\right) J Y^{-1}\left(-P_{001}\right. \\
&-\sum_{0<|j| \leq K_{+}} M_{-j} J F_{j 1}+[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010} \\
&\left.\left.+\operatorname{diag}\left(U^{-1}, O\right) \sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}\right)+\sum_{0<|j| \leq K_{+}} B_{k-j} J F_{j 1}\right),  \tag{3.23}\\
& F_{01}=Y^{-1}\left(-P_{001}-\sum_{0<|j| \leq K_{+}} M_{-j} J F_{j 1}+[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010}\right. \\
&\left.\quad+\operatorname{diag}\left(U^{-1}, O\right) \sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}\right),  \tag{3.24}\\
& y_{*}=\operatorname{diag}\left(U^{-1}, O\right)\left(-P_{010}-\sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}-[B] J Y^{-1}\left(-P_{001}\right.\right. \\
& \quad-\sum_{0<|j| \leq K_{+}} M_{-j} J F_{j 1}+[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010} \\
&\left.\left.\quad+\operatorname{diag}\left(U^{-1}, O\right) \sum_{0<|j| \leq K_{+}} B_{-j} J F_{j 1}\right)\right) . \tag{3.25}
\end{align*}
$$

Substituting (3.24), (3.25) into (3.18), it is easy to see that $F_{k 1}, 0<|k| \leq K_{+}$, satisfy the following equations:

$$
\begin{equation*}
L_{k} F_{k 1}=\sum_{0<|j| \leq K_{+}} M_{k j} F_{j 1}+\mathcal{P}_{k} \tag{3.26}
\end{equation*}
$$

where

$$
\begin{align*}
& L_{k}=\sqrt{-1}\langle k, \Omega(\lambda)\rangle I_{2 m}-[M] J,  \tag{3.27}\\
& M_{k j}= \begin{cases}\bar{M}_{k k}, & j=k, \\
M_{k-j} J+\bar{M}_{k j}, & j \neq k,\end{cases}  \tag{3.28}\\
& \mathcal{P}_{k}=P_{k 01}-B_{k}^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010}+B_{k}^{\top} \operatorname{diag}\left(U^{-1}, O\right)[B] J Y^{-1} P_{001} \\
& \quad-B_{k}^{\top} \operatorname{diag}\left(U^{-1}, O\right)[B] J Y^{-1}[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010} \\
& \quad-M_{k} J Y^{-1} P_{001}+M_{k} J Y^{-1}[B]^{\top} \operatorname{diag}\left(U^{-1}, O\right) P_{010}, \tag{3.29}
\end{align*}
$$

with

$$
\begin{aligned}
\bar{M}_{k j}= & M_{k} J Y^{-1}\left(-M_{-j} J+\operatorname{diag}\left(U^{-1}, O\right) B_{-j} J\right) \\
& +B_{k}^{\top} \operatorname{diag}\left(U^{-1}, O\right)\left(-B_{-j} J+[B] J Y^{-1}\left(M_{-j} J-\operatorname{diag}\left(U^{-1}, O\right) B_{-j} J\right)\right)
\end{aligned}
$$

Thus, equations (3.16)-(3.20) are solvable on $\mathcal{O}_{+}$if and only if (3.26) is. In fact, for the sake of convergence of the symplectic transformations to be constructed, not only do these equations need to be solvable, but also their solutions should satisfy
certain exponential decay properties. This motivates us to consider the following weighted functions:

$$
\begin{aligned}
\tilde{F}_{k 1} & =e^{|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} F_{k 1}, 0<|k| \leq K_{+} \\
\tilde{M}_{k j} & =e^{(|k|-|j|)\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} M_{k j}, 0<|k|,|j| \leq K_{+} \\
\tilde{\mathcal{P}_{k}} & =e^{|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \mathcal{P}_{k}, 0<|k| \leq K_{+}
\end{aligned}
$$

Let $\mathcal{T}="<"$ be a fixed ordering on $Z_{+}^{n}$ with the property that whenever $k, k_{*} \in Z_{+}^{n}$ with $|k|<\left|k_{*}\right|$, then $k<k_{*}$. We define

$$
\mathcal{F}=\left(\begin{array}{c}
\vdots \\
\tilde{F}_{j 1} \\
\vdots
\end{array}\right), \quad \mathcal{P}=\left(\begin{array}{c}
\vdots \\
\tilde{\mathcal{P}}_{j} \\
\vdots
\end{array}\right)
$$

as the column vectors which vertically line up all $\tilde{F}_{j 1}, \tilde{\mathcal{P}}_{j}, 0<|j| \leq K_{+}$, respectively, according to the ordering $\mathcal{T}$. We also form the matrices

$$
\begin{equation*}
\Lambda=\operatorname{diag}\left(L_{k}\right), \quad \mathcal{A}=\left(\tilde{M}_{k j}\right) \tag{3.30}
\end{equation*}
$$

according to the same ordering $\mathcal{T}$. Then it is clear that the equations in (3.26) can be putted into the following system form:

$$
\begin{equation*}
(\Lambda-\mathcal{A}) \mathcal{F}=\mathcal{P} \tag{3.31}
\end{equation*}
$$

3.3. Invertibility of $\Lambda-\mathcal{A}$. The invertibility of $\Lambda-\mathcal{A}$ will be shown by a perturbation argument.

Let $\left\{M_{k}^{0}\right\},\left\{B_{k}^{0}\right\}$ be the Fourier coefficients of $M^{0}, B^{0}$ respectively, and let $M_{k j}^{0}, 0<|k|,|j| \leq K_{+}$, be defined as in (3.28) with $M_{k}^{0}, B_{k}^{0}, U^{0}, Y^{0}$ in place of $M_{k}, B_{k}, U, Y$, for all $0<|k| \leq K_{+}$, respectively. Denote

$$
L_{k}^{0}=\sqrt{-1}\langle k, \omega\rangle-\left[M^{0}\right] J, \quad \tilde{M}_{k j}^{0}=e^{(|k|-|j|)\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} M_{k j}^{0}, 0<|j|,|k| \leq K_{+}
$$

and let $\Lambda^{0}, \mathcal{A}^{0}$ be defined as in (3.30) with $L_{k}^{0}, \tilde{M}_{k j}^{0}$ in place of $L_{k}, \tilde{M}_{k j}, 0<$ $|j|,|k| \leq K_{+}$, respectively. We first show the invertibility of $\Lambda^{0}-\mathcal{A}^{0}$ on $\mathcal{O}_{0}$, along with the estimate of an upper bound for its inverse.

By the hyperbolicity of $\left[M^{0}\right] J$, it is easy to see that all $L_{k}^{0}, 0<|k| \leq K_{+}$, are non-singular on $\mathcal{O}_{0}$. This implies that $\Lambda^{0}$ is non-singular on $\mathcal{O}_{0}$. To obtain an upper bound of $\left|\left(\Lambda^{0}\right)^{-1}\right|$ on $\mathcal{O}_{0}$, we need to estimate a uniform bound for all $\left|L_{k}^{0-1}\right|_{\mathcal{O}_{0}}$, $0<|k| \leq K_{+}$.

To do so, let $L=L(\lambda)$ be non-singular on $\mathcal{O}_{0}$ such that

$$
L^{-1}\left[M^{0}\right] J L=E
$$

where $E=E(\lambda)$ denotes the Jordan canonical form of $\left[M^{0}\right] J=\left[M^{0}(\lambda)\right] J$ for each $\lambda \in \mathcal{O}_{0}$. Due to the possible change of multiplicities of eigenvalues of $\left[M^{0}\right] J, L$ is in general not continuous on $\mathcal{O}_{0}$. However, by the standard Q-R decomposition procedure, for each $\lambda \in \mathcal{O}_{0}$, there is an upper triangular matrix $S=S(\lambda)$ such that $L S$ is orthogonal. Denote $\tilde{E}=\tilde{E}(\lambda) \equiv S^{-1} E S$, then $\tilde{E}$ is also upper triangular, and,

$$
(L S)^{\top}\left[M^{0}\right] J(L S)=(L S)^{-1}\left[M^{0}\right] J(L S)=S^{-1} E S=\tilde{E}
$$

Let $|\cdot|_{2}$ denote the matrix norm induced by the Euclidean vector norm. Then for any $n \times 2 m$ matrix $\left(a_{i j}\right)$,

$$
\frac{1}{\sqrt{2 m}}\left|\left(a_{i j}\right)\right| \leq\left|\left(a_{i j}\right)\right|_{2} \leq \sqrt{2 m}\left|\left(a_{i j}\right)\right|
$$

It follows that

$$
\begin{align*}
& |\tilde{E}| \leq \sqrt{2 m}|\tilde{E}|_{2}=\sqrt{2 m}\left|\left[M^{0}\right] J\right|_{2} \leq 2 m\left|\left[M^{0}\right]\right|  \tag{3.32}\\
& \left|\left(L_{k}^{0}\right)^{-1}\right| \leq \sqrt{2 m}\left|\left(L_{k}^{0}\right)^{-1}\right|_{2}=\sqrt{2 m}\left|\left(\sqrt{-1}\langle k, \omega\rangle I_{2 m}-\tilde{E}\right)^{-1}\right|_{2} \\
& \leq 2 m\left|\left(\sqrt{-1}\langle k, \omega\rangle I_{2 m}-\tilde{E}\right)^{-1}\right| \tag{3.33}
\end{align*}
$$

Since the eigenvalues of $\left[M^{0}\right] J$ coincide with those of $\tilde{E}$, eigenvalues of $\sqrt{-1}\langle k, \omega\rangle I_{2 m}-\tilde{E}$ are bounded below in absolute value by $\sigma_{0}$. Using H ), (3.32), (3.33) and the standard inverse formula for upper triangular matrices, we have that

$$
\begin{aligned}
\left|\left(L_{k}^{0}\right)^{-1}\right| & \leq 2 m\left|\left(\sqrt{-1}\langle k, \omega\rangle I_{2 m}-\tilde{E}\right)^{-1}\right| \leq \frac{2 m}{\sigma_{0}}\left(1+\frac{|\tilde{E}|_{\mathcal{O}_{0}}}{\sigma_{0}}\right)^{2 m-1} \\
& \leq \frac{2 m}{\sigma_{0}}\left(1+\frac{2 m}{\sigma_{0}}\left|\left[M^{0}\right]\right|_{\mathcal{O}_{0}}\right)^{2 m-1}=\frac{\rho_{0}}{2}
\end{aligned}
$$

for all $\lambda \in \mathcal{O}_{0}, 0<|k| \leq K_{+}$, i.e.,

$$
\begin{equation*}
\left|\left(\Lambda^{0}\right)^{-1}\right|_{\mathcal{O}_{0}} \leq \frac{\rho_{0}}{2} \tag{3.34}
\end{equation*}
$$

Next, we give an estimate for $\left|\mathcal{A}^{0}\right|_{\mathcal{O}_{0}}$.
Let $\mathcal{U}(x)=\left(u_{i j}(x)\right)$ be a real analytic, matrix valued function defined on $\mathcal{D}(a)$ $(a>0)$ and denote $\mathcal{U}_{k}=\left(u_{k i j}\right), k \in Z^{n} \backslash\{0\}$, as the $k$ th Fourier coefficient of $\mathcal{U}(x)$. Since

$$
\mathcal{U}_{k}=\int_{T^{n}}(\mathcal{U}-[\mathcal{U}]) e^{-\sqrt{-1}\langle k, x\rangle} \mathrm{d} x
$$

Cauchy's estimate yields that

$$
\begin{equation*}
\left|\mathcal{U}_{k}\right| \leq|\mathcal{U}-[\mathcal{U}]|_{\mathcal{D}(a)} e^{-|k| a} \tag{3.35}
\end{equation*}
$$

It also follows from Parseval's identity that

$$
\begin{align*}
\sum_{0<|k| \leq \tilde{K}}\left|\mathcal{U}_{k}\right| & =\sum_{0<|k| \leq \tilde{K}} \max _{i} \sum_{j}\left|u_{k i j}\right| \leq \sqrt{2 m} \sum_{0<|k| \leq \tilde{K}} \sum_{i, j}\left|u_{k i j}\right|^{2} \\
& =\sqrt{2 m} \sum_{i, j}\left(\sum_{0<|k| \leq \tilde{K}}\left|u_{k i j}\right|^{2}\right)^{\frac{1}{2}} \leq 2 m \max _{i} \sum_{j}\left\|u_{i j}-\left[u_{i j}\right]\right\|_{L^{2}(\mathcal{D}(a))} \\
& \leq 2 m \max _{i} \sum_{j}\left|u_{i j}-\left[u_{i j}\right]\right|_{\mathcal{D}(a)}=2 m|\mathcal{U}-[\mathcal{U}]|_{\mathcal{D}(a)} \tag{3.36}
\end{align*}
$$

for any $\tilde{K}>0$.

Applying (3.36) to the formula for $M_{k j}^{0}$ similar to (3.28), we see that, for any $0<|k| \leq K_{+}, \lambda \in \mathcal{O}_{0}$,

$$
\begin{align*}
& \sum_{0<|j| \leq K_{+}}\left|M_{k j}^{0}\right|=\left|M_{k k}^{0}\right|+\sum_{0<|j| \leq K_{+}, j \neq k}\left|M_{k j}^{0}\right| \\
& \quad \leq\left|M^{0}-\left[M^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}+(1+2 m)\left|\left(Y^{0}\right)^{-1}\right|\left|M^{0}-\left[M^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}^{2} \\
& \quad+(1+2 m)\left|\left(Y^{0}\right)^{-1}\right|\left|\left(U^{0}\right)^{-1}\right|\left|M^{0}-\left[M^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}\left|B^{0}-\left[B^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)} \\
& \quad+(1+2 m)\left|\left(Y^{0}\right)^{-1}\right|\left|\left(U^{0}\right)^{-1}\right|\left|B^{0}-\left[B^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}^{2} \\
& \quad+(1+2 m)\left|\left(U^{0}\right)^{-1}\right|\left|B^{0}-\left[B^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}^{2}+(1+2 m)\left|\left[B^{0}\right]\right|\left|\left(Y^{0}\right)^{-1}\right|\left|\left(U^{0}\right)^{-1}\right| \\
& \left|M^{0}-\left[M^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}\left|B^{0}-\left[B^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)} \\
& \quad+(1+2 m)\left|\left[B^{0}\right]\right|\left|\left(Y^{0}\right)^{-1}\right|\left|\left(U^{0}\right)^{-1}\right|^{2}\left|B^{0}-\left[B^{0}\right]\right|_{\mathcal{D}\left(r_{*}\right)}^{2}<\eta+\alpha \eta^{2}=\frac{1}{\rho_{0}}, \tag{3.37}
\end{align*}
$$

where $\eta, \alpha$ are as in (1.8), (1.9) respectively, i.e.,

$$
\begin{aligned}
\eta & =\frac{2}{\sqrt{\rho_{0}^{2}+4 \alpha \rho_{0}}+\rho_{0}} \\
\alpha & =(1+2 m)\left(\left|\left(Y^{0}\right)^{-1}\right|+\left|\left(U^{0}\right)^{-1}\right|\right. \\
& \left.+\left(\left|\left(Y^{0}\right)^{-1}\right|\left|\left(U^{0}\right)^{-1}\right|\right)\left(1+\left|\left[B^{0}\right]\right|+\left|\left[B^{0}\right]\right|\left|\left(U^{0}\right)^{-1}\right|\right)\right)_{\mathcal{O}_{0}}
\end{aligned}
$$

We now choose $r_{0}$ in (3.1). Define

$$
\mathcal{A}^{0}(a)=\left(e^{(|k-j|) a} M_{k j}^{0}\right)
$$

as the one parameter family of matrices which are of the same dimension as $\mathcal{A}$ and are formed according to the ordering $\mathcal{T}$. Then (3.37) clearly implies that

$$
\left|\mathcal{A}^{0}(0)\right|_{\mathcal{O}_{0}}<\frac{1}{\rho_{0}}
$$

Note that, by (3.35),

$$
\begin{aligned}
\sum_{|j|>0} e^{|k-j| a}\left|\partial_{\lambda}^{l} M_{k j}^{0}\right| & \leq c \sum_{|j|>0} e^{-|k-j|\left(r_{*}-a\right)} \leq c \sum_{|j| \geq 0} e^{-|j|\left(r_{*}-a\right)}<\infty \\
\sum_{|j|>0}|k-j| e^{|k-j| a}\left|\partial_{\lambda}^{l} M_{k j}^{0}\right| & \leq c \sum_{|j|>0}|k-j| e^{-|k-j|\left(r_{*}-a\right)} \\
& \leq c \sum_{|j| \geq 0}|j| e^{-|j|\left(r_{*}-a\right)}<\infty
\end{aligned}
$$

for all $|l| \leq n$. It follows that the family of functions

$$
\eta_{k, \tilde{K}}(a)=\sum_{0<|j| \leq \tilde{K}} e^{|k-j| a}\left|M_{k j}^{0}\right|
$$

is uniformly bounded and equi-continuous. Hence

$$
\left|\mathcal{A}^{0}(a)\right|_{\mathcal{O}_{0}}=\max _{0<|k| \leq K_{+}} \sum_{0<|j| \leq K_{+}} e^{|k-j| a}\left|M_{k j}^{0}\right| \mathcal{O}_{0}
$$

is continuous in $a \in\left[0, r_{*}\right)$ uniformly in $K_{+}$.
Let $r_{0}=r_{0}\left(r_{*}, \sigma_{0}, U^{0}\right) \in\left(0, \frac{r_{*}}{2}\right]$ be fixed such that

$$
\left|\mathcal{A}^{0}\left(\left(1-\frac{d \delta^{2}}{16}\right) r_{0}\right)\right|_{\mathcal{O}_{0}}<\frac{1}{\rho_{0}} .
$$

Denote

$$
\begin{equation*}
\xi(a)=\left(\frac{7}{8}+\frac{\delta}{8}\right) a+\frac{d}{8}\left(1-\frac{\delta^{2}}{2}\right) r_{0} \tag{3.38}
\end{equation*}
$$

Then

$$
\begin{aligned}
\xi(r) & =r_{+}+\frac{7}{8}\left(r-r_{+}\right) \\
\xi\left(r_{0}\right) & =\left(1-\frac{d \delta^{2}}{16}\right) r_{0}
\end{aligned}
$$

Since

$$
e^{(|k|-|j|) \xi(r)}\left|M_{k j}^{0}\right| \leq e^{|k-j| \xi(r)}\left|M_{k j}^{0}\right|<e^{|k-j| \xi\left(r_{0}\right)}\left|M_{k j}^{0}\right|, 0<|k|,|j| \leq K_{+},
$$

we have that

$$
\begin{equation*}
\left|\mathcal{A}^{0}\right|_{\mathcal{O}_{0}} \leq\left|\mathcal{A}^{0}\left(\xi\left(r_{0}\right)\right)\right|_{\mathcal{O}_{0}}<\frac{1}{\rho_{0}} \tag{3.39}
\end{equation*}
$$

Thus, by (3.34),

$$
\left|\left(\Lambda^{0}-\mathcal{A}^{0}\right)^{-1}\right|_{\mathcal{O}_{0}} \leq \frac{\left|\left(\Lambda^{0}\right)^{-1}\right|_{\mathcal{O}_{0}}}{1-\left|\mathcal{A}^{0}\right|_{\mathcal{O}_{0}}\left|\left(\Lambda^{0}\right)^{-1}\right|_{\mathcal{O}_{0}}}<2 \rho_{0}
$$

Lemma 3.2. Assume H3) and also that
H4) $\left|\partial_{\lambda}^{l} \mathcal{A}-\partial_{\lambda}^{l} \mathcal{A}^{0}\right|_{\mathcal{O}}<\mu_{*}^{\frac{1}{4}}$.
Then as $\mu_{*}$ sufficiently small, $\mathcal{L}=\Lambda-\mathcal{A}$ is non-singular on $\mathcal{O}$, and moreover, there is a constant $c_{2}$ such that

$$
\left|\partial_{\lambda}^{l} \mathcal{L}\right|_{\mathcal{O}} \leq c_{2} K_{+},|l| \leq n
$$

Proof. Given $\epsilon_{0}>0$ small. Since, by H3),

$$
\left|U-U^{0}\right|_{\mathcal{O}},\left|Y-Y^{0}\right|_{\mathcal{O}},\left|M-M^{0}\right|_{\mathcal{O}} \leq\left|\mathcal{M}-\mathcal{M}^{0}\right|_{\mathcal{O}} \leq \mu_{*}^{\frac{1}{4}}
$$

we can choose $\mu_{*}$ small, say, $\mu_{*} \leq \frac{\epsilon_{0}^{4}}{1+\epsilon_{0}} \min \left\{\frac{1}{\left|\left(U^{0}\right)^{-1}\right|_{\mathcal{O}_{0}}}, \frac{1}{\left|\left(Y^{0}\right)^{-1}\right| \mathcal{O}_{0}}\right\}$, such that

$$
\begin{align*}
\left|U^{-1}\right|_{\mathcal{O}} & \leq\left(1+\epsilon_{0}\right)\left|\left(U^{0}\right)^{-1}\right|_{\mathcal{O}_{0}}  \tag{3.40}\\
\left|Y^{-1}\right|_{\mathcal{O}} & \leq\left(1+\epsilon_{0}\right)\left|\left(Y^{0}\right)^{-1}\right|_{\mathcal{O}_{0}}  \tag{3.41}\\
|M|_{\mathcal{O}} & \leq\left(1+\epsilon_{0}\right)\left|M^{0}\right|_{\mathcal{O}_{0}}
\end{align*}
$$

Define $\rho_{*}$ similar to $\rho_{0}$ with $M, \mathcal{O}, \sigma_{*}=\left(1-\epsilon_{0}\right) \sigma_{0}$ in place of $M^{0}, \mathcal{O}_{0}, \sigma_{0}$ respectively. Then, $\rho_{*} \leq \frac{4}{3} \rho_{0}$ as $\epsilon_{0}$ small, and, as $\mu_{*}$ small the real parts of all eigenvalues of $[M] J$ are bounded below in absolute value by $\sigma_{0}$.

On one hand, by a similar argument as for (3.34), we have that

$$
\left|\Lambda^{-1}\right|_{\mathcal{O}} \leq \frac{1}{2} \rho_{*} \leq \frac{2}{3} \rho_{0}
$$

On the other hand, by (3.39) and H 4 ), we can make $\mu_{*}$ even smaller if necessary so that

$$
|\mathcal{A}|_{\mathcal{O}}<\frac{2}{3 \rho_{0}}
$$

It follows that $\mathcal{L}$ is non-singular on $\mathcal{O}$ and

$$
\left|\mathcal{L}^{-1}\right|_{\mathcal{O}} \leq 4 \rho_{0}
$$

Since, by H3) and H4),

$$
\left|\partial_{\lambda} \mathcal{L}\right|_{\mathcal{O}} \leq\left|\partial_{\lambda} \Lambda\right|_{\mathcal{O}}+\left|\partial_{\lambda} \mathcal{A}\right|_{\mathcal{O}} \leq c K_{+}
$$

we also have

$$
\left|\partial_{\lambda} \mathcal{L}^{-1}\right|_{\mathcal{O}}=\left|\mathcal{L}^{-1}\left(\partial_{\lambda} \mathcal{L}\right) \mathcal{L}^{-1}\right|_{\mathcal{O}} \leq\left|\partial_{\lambda} \mathcal{L}\right|_{\mathcal{O}}\left|\mathcal{L}^{-1}\right|_{\mathcal{O}}^{2} \leq c K_{+} .
$$

By induction,

$$
\left|\partial_{\lambda}^{l} \mathcal{L}^{-1}\right|_{\mathcal{O}} \leq c K_{+},|l| \leq n
$$

Above all, with the hypotheses H 3 ), H 4 ), the linear system (3.31) can be uniquely solved on $\mathcal{O}_{+}$to yield smooth functions $f_{k 0}, f_{k 1}, F_{k 1}, F_{01}, y_{*}, 0<|k| \leq K_{+}$. Hence, the Hamiltonian $F$ in (3.4) is well defined, smooth in $\lambda \in \mathcal{O}_{+}$, and real analytic in $(x, y, z) \in D$.
3.4. Estimates on the transformation. We first give some estimates on $F$ and its derivatives. Denote

$$
\zeta=K_{+} \Gamma\left(r-r_{+}\right)
$$

Lemma 3.3. Assume H3), H4) and also that
H5) $\left|\partial_{\omega}^{l}(\Omega-\omega)\right|_{O} \leq \mu_{*}^{\frac{1}{4}},|l| \leq n$.
Then there is a constant $c_{3}$ such that the following holds for all $|l| \leq n$.

1) $\left|\partial_{\lambda}^{l} y_{*}\right|_{\mathcal{O}_{+}} \leq c_{3} \gamma^{n+1} s \mu \zeta$.
2) $O n D(s) \times \mathcal{O}_{+}$,

$$
\left|\partial_{\lambda}^{l} F\right|,\left|\partial_{\lambda}^{l} F_{x}\right|, s\left|\partial_{\lambda}^{l} F_{y}\right|, s\left|\partial_{\lambda}^{l} F_{z}\right| \leq c_{3} s^{2} \mu \zeta
$$

3) $O n D(\beta) \times \mathcal{O}_{+}$,

$$
\left|\partial_{\lambda}^{l} D^{i} F\right| \leq c_{3} \mu \zeta,|i| \leq 4
$$

where $D=\partial_{(x, y, z)}$.
Proof. Let $|l| \leq n$. First, we observe by Cauchy's estimate that

$$
\begin{align*}
\left|\partial_{\lambda}^{l} P_{k i j}\right|_{\mathcal{O}} & \leq c s^{-(i+j)}\left|\partial_{\lambda}^{l} P\right|_{D(r, s) \times \mathcal{O}} e^{-|k| r} \\
& \leq c \gamma^{n+1} s^{2-i-j} \mu e^{-|k| r},|k| \geq 0, i, j=0,1,2 \tag{3.42}
\end{align*}
$$

and by (3.35) and H3) that

$$
\begin{align*}
\left|\partial_{\lambda}^{l} M_{k}\right|_{\mathcal{O}},\left|\partial_{\lambda}^{l} B_{k}\right|_{\mathcal{O}} & \leq\left|\partial_{\lambda}^{l} \mathcal{M}-\partial_{\lambda}^{l} \mathcal{M}^{0}\right|_{\mathcal{D}(r) \times \mathcal{O}} e^{-|k| r} \\
& \leq \gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}} e^{-|k| r} \leq 2^{n+1} \gamma^{n+1} \mu_{*}^{\frac{1}{4}} e^{-|k| r},|k|>0 \tag{3.43}
\end{align*}
$$

Let $\lambda \in \mathcal{O}_{+},|k|>0$. Then, by (3.29), (3.40)-(3.43), we have that

$$
\left|\partial_{\lambda}^{l} \mathcal{P}_{k}\right| \leq c \gamma^{n+1} s \mu e^{-|k| r}
$$

i.e.,

$$
\left|\partial_{\lambda}^{l} \mathcal{P}\right| \leq c \gamma^{n+1} s \mu
$$

Using (3.31) and Lemma 3.2, it follows that

$$
\left|\partial_{\lambda}^{l} \tilde{F}_{k 1}\right| \leq\left|\partial_{\lambda}^{l} \mathcal{F}\right|=\left|\partial_{\lambda}^{l}\left(\mathcal{L}^{-1} \mathcal{P}\right)\right| \leq c \gamma^{n+1} s \mu K_{+}
$$

Hence,

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} F_{k 1}\right| \leq c \gamma^{n+1} s \mu K_{+} e^{-|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \tag{3.44}
\end{equation*}
$$

By straightforward applications of (3.40)-(3.44) to (3.24) and (3.25), we then obtain 1) and also that

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} F_{01}\right| \leq c \gamma^{n+1} s \mu K_{+} . \tag{3.45}
\end{equation*}
$$

Next, we note by H5), (3.21) and a direct calculation that

$$
\left|\partial_{\lambda}^{l}\langle k, \Omega(\lambda)\rangle^{-1}\right| \leq c \frac{|k|^{|l|+(|k|+1) \tau}}{\gamma^{|l|+1}}
$$

This together with (3.40)-(3.44) implies that

$$
\begin{align*}
\left|\partial_{\lambda}^{l} f_{k 0}\right| & \leq c|k|^{|l|+(|l|+1) \tau} s^{2} \mu e^{-|k| r}  \tag{3.46}\\
\left|\partial_{\lambda}^{l} f_{k 1}\right| & \leq c|k|^{|l|+(|l|+1) \tau} s \mu K_{+} e^{-|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \tag{3.47}
\end{align*}
$$

Now, using the expression of $F, y_{*}$ in (3.4) and (3.25) respectively, parts 2), 3) of the lemma follow directly from (3.43)-(3.47).

Lemma 3.4. Assume H 1$), \mathrm{H} 3)-\mathrm{H} 5)$ and also that
H6) $c_{3} \mu \zeta<\frac{1}{8}\left(r-r_{+}\right)$;
H7) $c_{3} s \mu \zeta<s_{+}$;
H8) $c_{3} \mu \zeta<\frac{\beta-\beta_{+}}{2}$.
Let $\phi_{F}^{t}$ be the flow generated by $F$. Then the following holds.

1) For all $0 \leq t \leq 1$,

$$
\begin{array}{rll}
\phi_{F}^{t}: D_{3} & \rightarrow & D_{4} \\
\phi: D_{1} & \rightarrow & D_{3}
\end{array}
$$

are well defined, real analytic and depend smoothly on $\lambda \in \mathcal{O}_{+}$.
2) Let $\Phi_{+}=\phi_{F}^{1} \circ \phi$. Then for all $\lambda \in \mathcal{O}_{+}$,

$$
\begin{array}{ll}
\Phi_{+}: & D_{+} \rightarrow D \\
\tilde{D}_{+} \rightarrow D(r, \beta) .
\end{array}
$$

3) There is a constant $c_{4}$ such that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l}\left(\phi_{F}^{t}-i d\right)\right|_{D(s) \times \mathcal{O}_{+}} & \leq c_{4} s \mu \zeta \\
\left|\partial_{\lambda}^{l} D^{i}\left(\Phi_{+}-i d\right)\right|_{\tilde{D}_{+} \times \mathcal{O}_{+}} & \leq c_{4} \mu \zeta
\end{aligned}
$$

for all $|l| \leq n, 0 \leq i \leq 3,0 \leq t \leq 1$, where $D=\partial_{(x, y, z)}$.
Proof. Let $\lambda \in \mathcal{O}_{+}$.

1) We note that

$$
\begin{equation*}
\phi_{F}^{t}=\mathrm{id}+\int_{0}^{t} X_{F} \circ \phi_{F}^{\xi} \mathrm{d} \xi \tag{3.48}
\end{equation*}
$$

where $X_{F}=\left(F_{y},-F_{x}, J F_{z}\right)^{\top}$.
Denote $\phi_{F 1}^{t}, \phi_{F 2}^{t}, \phi_{F 3}^{t}$ as components of $\phi_{F}^{t}$ in $x, y, z$ planes respectively. For any $(x, y, z) \in D_{3}$, let $t_{*}=\sup \left\{t \in[0,1]: \phi_{F}^{t}(x, y, z) \in D_{4}\right.$. By H1), we have that
$D_{4} \subset D(s)$. It follows from H6), H7) and Lemma 3.3 that

$$
\begin{aligned}
\left|\phi_{F 1}^{t}(x, y, z)\right| & =|x|+\left|\int_{0}^{t} F_{y} \circ \phi_{F}^{\xi} \mathrm{d} \xi\right| \leq|x|+\left|F_{y}\right|_{D(s)} \leq r_{+}+\frac{2}{8}\left(r-r_{+}\right)+c_{3} s \mu \zeta \\
& <r_{+}+\frac{3}{8}\left(r-r_{+}\right) \\
\left|\phi_{F 2}^{t}(x, y, z)\right| & =|y|+\left|-\int_{0}^{t} F_{x} \circ \phi_{F}^{\xi} \mathrm{d} \xi\right| \leq|y|+\left|F_{x}\right|_{D(s)} \leq 3 s_{+}+c_{3} s^{2} \mu \zeta<4 s_{+} \\
\left|\phi_{F 3}^{t}(x, y, z)\right| & =|z|+\left|\int_{0}^{t} J F_{z} \circ \phi_{F}^{\xi} \mathrm{d} \xi\right| \leq|z|+\left|F_{z}\right|_{D(s)} \leq 3 s_{+}+c_{3} s \mu \zeta<4 s_{+}
\end{aligned}
$$

i.e., $\phi_{F}^{t}(x, y, z) \in D_{4}$ for all $0 \leq t \leq t_{*}$. Thus, $t_{*}=1$ and 1$)$ holds.
2) follows from Lemma 3.3, H8) and a similar argument as 1).
3) Using Lemma 3.3 and (3.48), we immediately have

$$
\left|\phi_{F}^{t}-i d\right|_{D(s)} \leq c_{3} s \mu \zeta
$$

Differentiating (3.48) with respect to $\lambda$ yields

$$
\partial_{\lambda} \phi_{F}^{t}=\int_{0}^{t} D X_{F} \circ \phi_{F}^{\xi} \partial_{\lambda} \phi_{F}^{\xi} \mathrm{d} \xi+\int_{0}^{t}\left(\partial_{\lambda} F_{y},-\partial_{\lambda} F_{x}, J \partial_{\lambda} F_{z}\right)^{\top} \circ \phi_{F}^{\xi} \mathrm{d} \xi
$$

It follows from Lemma 3.3 and Gronwall's inequality that

$$
\left|\partial_{\lambda} \phi_{F}^{t}\right|_{D(s)} \leq c s \mu \zeta .
$$

By induction, we have

$$
\left|\partial_{\lambda}^{l} \phi_{F}^{t}\right|_{D(s)} \leq c s \mu \zeta,|l| \leq n
$$

The estimates for $\Phi_{+}$follow from a similar application of Lemma 3.3 and Gronwall's inequality, and the identity

$$
\Phi_{+}-i d=\left(\phi_{F}^{1}-i d\right) \circ \phi+\left(\begin{array}{l}
0 \\
y_{*} \\
0
\end{array}\right)
$$

We omit the details.
3.5. Estimates on the new Hamiltonian. We first estimate the new normal form.

Lemma 3.5. Assume H3), H5). Then there is a constant $c_{5}$ such that the following holds for all $|l| \leq n$ :

$$
\begin{aligned}
\left|\partial_{\lambda}^{l}\left(e_{+}-e\right)\right|_{\mathcal{O}_{+}} & \leq c_{5} \gamma^{n+1} s \mu \zeta \\
\left|\partial_{\lambda}^{l}\left(\Omega_{+}-\Omega\right)\right|_{\mathcal{O}_{+}} & \leq c_{5} \gamma^{n+1} s \mu \zeta \\
\left|\partial_{\lambda}^{l}\left(\mathcal{M}^{+}-\mathcal{M}\right)\right|_{\mathcal{D}_{+} \times \mathcal{O}_{+}} & \leq c_{5} \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right)
\end{aligned}
$$

Proof. The estimates for $\left|\partial_{\lambda}^{l}\left(e^{+}-e\right)\right|_{\mathcal{O}_{+}}$follows immediately from (3.9), (3.42), H5) and Lemma 3.3 1). Also, it follows from (3.10), (3.42)-(3.44) and Lemma 3.3 1) that

$$
\left|\partial_{\lambda}^{l}\left(\Omega_{+}-\Omega\right)\right|_{\mathcal{O}_{+}} \leq c \gamma^{n+1} s \mu \zeta+c \gamma^{n+1} s \mu K_{+} \sum_{|k|>0} e^{-\frac{|k| r_{0}}{2}} \leq c \gamma^{n+1} s \mu \zeta
$$

and from (3.11)-(3.13) and (3.42) that

$$
\begin{align*}
& \left|\partial_{\lambda}^{l}\left(\mathcal{M}^{+}-\mathcal{M}\right)\right|_{\mathcal{D}\left(r_{+}+\frac{1}{8}\left(r-r_{+}\right)\right) \times \mathcal{O}_{+}} \leq c \sum_{|k|>0}|k| \sum_{i+j=2}\left|\partial_{\lambda}^{l} P_{k i j}\right|_{\mathcal{O}} e^{|k|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \quad \leq c \gamma^{n+1} \mu \sum_{|k|>0}|k| e^{-|k| \frac{r-r_{+}}{8}} \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right),|l| \leq n \tag{3.49}
\end{align*}
$$

This proves the lemma.

Lemma 3.6. Assume that
H9) $c_{5} \gamma^{n+1} s \mu \zeta K_{+}^{\tau+1}<\gamma-\gamma_{+}$.
Then

$$
\left|\left\langle k, \Omega_{+}(\lambda)\right\rangle\right|>\frac{\gamma_{+}}{|k|^{\tau}}
$$

for all $\lambda \in \mathcal{O}_{+}$and $0<|k| \leq K_{+}$.
Proof. By H9) and Lemma 3.5, we have

$$
\begin{aligned}
\left|\left\langle k, \Omega_{+}(\lambda)\right\rangle\right| & \geq \mid\langle k, \Omega(\lambda)|-c_{5} \gamma^{n+1} s \mu \zeta K_{+}^{\tau+1} \\
& \geq \frac{\gamma}{|k|^{\tau}}-c_{5} \gamma^{n+1} s \mu \zeta K_{+}^{\tau+1}>\frac{\gamma_{+}}{|k|^{\tau}}
\end{aligned}
$$

as desired.
Let $U^{+}$be the $n_{0} \times n_{0}$ ordered principal block of $\left[A^{+}\right]$. By (3.49) and H7), we see that

$$
\begin{aligned}
\left|U^{+}-U\right|_{\mathcal{O}_{+}} & \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) \leq c \mu \zeta \leq \frac{c}{c_{3}} \frac{s_{+}}{s} \\
& =\frac{c}{c_{3}} s^{b+\sigma} \leq \frac{c}{c_{3}} s_{0}^{b+\sigma} \leq \frac{c}{c_{3}} \mu_{*}^{\frac{b+\sigma}{4}}
\end{aligned}
$$

It follows that $U^{+}$is non-singular on $\mathcal{O}_{+}$as long as $\mu_{*}$ is small. A similar argument shows the same for

$$
Y^{+}=\left[M^{+}\right] J-\left[B^{+}\right]^{\top} \operatorname{diag}\left(\left(U^{+}\right)^{-1}, O\right)\left[B^{+}\right] J
$$

Now, let $M_{k}^{+}, B_{k}^{+}$denote the Fourier coefficients of $M^{+}, B^{+}$, respectively, and define $M_{k j}^{+},|k|,|j| \neq 0$, as in (3.28) with $M_{k}^{+}, B_{k}^{+}, U^{+}, Y^{+}$in place of $M_{k}, B_{k}, U, Y$, for all $|k| \neq 0$, respectively. We have the following.

Lemma 3.7. Assume H 3$)$, H 5 ) and H 7 ). Then there is a constant $c_{6}$ such that

$$
\max _{|k|>0} \sum_{|j|>0} e^{(|k|-|j|) \xi\left(r_{+}\right)}\left|\partial_{\lambda}^{l} M_{k j}^{+}-\partial_{\lambda}^{l} M_{k j}\right| \mathcal{O}_{+} \leq c_{6} \gamma^{n+1} \mu \Gamma^{2}\left(r-r_{+}\right),|l| \leq n
$$

where $\xi$ is as in (3.38).
Proof. First of all, we note by H3) and (3.49) that $M_{k}^{+}, B_{k}^{+}, M_{k}, B_{k},\left(U^{+}\right)^{-1}$, $\left(Y^{+}\right)^{-1}, U^{-1}, Y^{-1}$ are uniformly bounded on $\mathcal{O}_{+}$by a constant which is independent of the iteration process. Secondly, by (3.49) and Cauchy's estimate, we have
that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} M_{j}^{+}-\partial_{\lambda}^{l} M_{j}\right|_{\mathcal{O}_{+}} & \leq\left|\partial_{\lambda}^{l} M^{+}-\partial_{\lambda}^{l} M\right|_{\mathcal{D}\left(r_{+}+\frac{1}{8}\left(r-r_{+}\right)\right) \times \mathcal{O}_{+}} e^{-|j|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \left.\leq \mid \partial_{\lambda}^{l} \mathcal{M}^{+}-\partial_{\lambda}^{l} \mathcal{M}\right)\left.\right|_{\mathcal{D}\left(r_{+}+\frac{1}{8}\left(r-r_{+}\right)\right) \times \mathcal{O}_{+}} e^{-|j|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) e^{-|j| \xi(r)}, \\
\left|\partial_{\lambda}^{l} B_{j}^{+}-\partial_{\lambda}^{l} B_{j}\right|_{\mathcal{O}_{+}} & \leq\left|\partial_{\lambda}^{l} B^{+}-\partial_{\lambda}^{l} B\right|_{\mathcal{D}\left(r_{+}+\frac{1}{8}\left(r-r_{+}\right)\right) \times \mathcal{O}_{+}} e^{-|j|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \left.\leq \mid \partial_{\lambda}^{l} \mathcal{M}^{+}-\partial_{\lambda}^{l} \mathcal{M}\right)\left.\right|_{\mathcal{D}\left(r_{+}+\frac{1}{8}\left(r-r_{+}\right)\right) \times \mathcal{O}_{+}} e^{-|j|\left(r_{+}+\frac{7}{8}\left(r-r_{+}\right)\right)} \\
& \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) e^{-|j| \xi(r)}
\end{aligned}
$$

It follows from (3.28) and a similar formular for $M_{k j}^{+}$that

$$
\left|\partial_{\lambda}^{l} M_{k j}^{+}-\partial_{\lambda}^{l} M_{k j}\right|_{\mathcal{O}_{+}} \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) e^{-|k-j| \xi(r)}
$$

for all $|k|,|j|>0$.
Hence,

$$
\begin{aligned}
& \sum_{|j|>0} e^{(|k|-|j|) \xi\left(r_{+}\right)}\left|\partial_{\lambda}^{l} M_{k j}^{+}-\partial_{\lambda}^{l} M_{k j}\right|_{\mathcal{O}_{+}} \leq \sum_{|j|>0} e^{(|k-j|) \xi\left(r_{+}\right)}\left|\partial_{\lambda}^{l} M_{k j}^{+}-\partial_{\lambda}^{l} M_{k j}\right|_{\mathcal{O}_{+}} \\
& \quad \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) \sum_{|j|>0} e^{-|k-j|\left(\xi(r)-\xi\left(r_{+}\right)\right)} \\
& \quad \leq c \gamma^{n+1} \mu \Gamma\left(r-r_{+}\right) \sum_{|j|>0} e^{-|k-j| \frac{r-r_{+}}{8}} \leq c \gamma^{n+1} \mu \Gamma^{2}\left(r-r_{+}\right)
\end{aligned}
$$

Let

$$
\begin{equation*}
\Delta=\frac{\zeta^{3}}{r-r_{+}}\left(s^{3} \mu+\gamma^{n+1} s^{2} \mu^{2}+\gamma^{n+1} \frac{s_{+}^{3}}{s} \mu\right) \tag{3.50}
\end{equation*}
$$

Lemma 3.8. Assume H 2$)-\mathrm{H} 7)$. Then there is a constant $c_{7}$ such that

$$
\left|\partial_{\lambda}^{l} P_{+}\right|_{D_{+}} \leq c_{7} \Delta,|l| \leq n .
$$

Thus, if
H10) $c_{7} \Delta \leq \gamma_{+}^{n+1} s_{+}^{2} \mu_{+}$, and $2 s \mu K_{+}^{\tau+1}<\gamma-\gamma_{+}$,
then

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} P_{+}\right|_{D_{+}} \leq \gamma_{+}^{n+1} s_{+}^{2} \mu_{+} \tag{3.51}
\end{equation*}
$$

Proof. Let $|l| \leq n, \lambda \in \mathcal{O}_{+}$. By Lemma 3.1 and Cauchy's estimate, we have

$$
\left|\partial_{\lambda}^{l} D(P-R)\right|_{D_{4}} \leq \frac{c}{\left(r-r_{+}\right) s_{+}}\left|\partial_{\lambda}^{l}(P-R)\right|_{D_{8}} \leq \frac{c \gamma^{n+1} \mu}{r-r_{+}}\left(\frac{s^{3}}{s_{+}}+\frac{s_{+}^{2}}{s}\right)
$$

where $D=\partial_{(x, y, z)}$. This together with Lemma 3.31$)$ implies that

$$
\begin{align*}
\left|\partial_{\lambda}^{l}(P-R) \circ \Phi_{+}\right|_{D_{+}} & \leq c \frac{\gamma^{n+1} \mu^{2}}{r-r_{+}}\left(\frac{s^{4}}{s_{+}}+s_{+}^{2}\right) \zeta+c \mu\left(s^{3}+\frac{s_{+}^{3}}{s}\right) \\
& \leq c \frac{\gamma^{n+1}}{r-r_{+}}\left(s^{2} \mu^{2}+\frac{s_{+}^{3}}{s} \mu\right) \zeta \tag{3.52}
\end{align*}
$$

Also, a direct estimate using (3.42) yields that

$$
\begin{align*}
\left|\partial_{\lambda}^{l} D^{\alpha} R\right|_{D(s)} & \leq \sum_{0<|k| \leq K_{+}, 0 \leq i, j \leq 2}\left|\partial_{\lambda}^{l} P_{k i j}\right| s^{i+j-1}|k| e^{|k|\left(r_{+}+\frac{6}{8}\left(r-r_{+}\right)\right.} \\
(3.53) & \leq c \gamma^{n+1} s \mu \sum_{0<|k| \leq K_{+}}|k| e^{-|k| \frac{r-r_{+}}{4}} \leq c \gamma^{n+1} s \mu \Gamma\left(r-r_{+}\right), \alpha=0,1 \tag{3.53}
\end{align*}
$$

Similarly, by (3.5), (3.42)-(3.44),

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} D^{\alpha} R^{\prime}\right|_{D(s)} \leq c \gamma^{n+1} s \mu K_{+} \Gamma\left(r-r_{+}\right)=c \gamma^{n+1} s \mu \zeta, \alpha=0,1 \tag{3.54}
\end{equation*}
$$

We now estimate $Q$ in (3.15). Let

$$
\begin{aligned}
& Q_{1}=\sum_{0<|k| \leq K_{+}}\left(\left\langle\frac{1}{2} \partial_{x}\langle y, A(x, \lambda) y\rangle+\partial_{x}\langle y, B(x, \lambda) z\rangle\right.\right. \\
& \left.+\frac{1}{2} \partial_{x}\langle z, M(x, \lambda) z\rangle+\partial_{x} h_{0}(x, y, z, \lambda), f_{k 1}\right\rangle \\
& +\left(-\sqrt{-1}\langle k, A(x, \lambda) y+B(x, \lambda) z\rangle+\partial_{y} h_{0}(x, y, z, \lambda)\right)\left(f_{k 0}+\left\langle f_{k 1}, y\right\rangle+\left\langle F_{k 1}, z\right\rangle\right) \\
& \left.+\left\langle\partial_{z} h_{0}(x, y, z, \lambda) J F_{k 1}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}, \\
& Q_{2}=\sum_{|k|>K_{+}}\left(\left\langle B_{k} J F_{01}, y\right\rangle+\left\langle M_{k} J F_{01}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}, \\
& Q_{3}=\sum_{|k|>K_{+},|j| \leq K_{+}}\left(\left\langle B_{k-j} J F_{j 1}, y\right\rangle+\left\langle M_{k-j} J F_{j 1}, z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}, \\
& q_{0}=\left\langle B^{\top}(x) y_{*}, J F_{01}\right\rangle, \\
& \left.q_{1}=-\sum_{0<|k| \leq K_{+}} \sqrt{-1}\left\langle f_{k 1}, y_{*}\right\rangle+\left\langle B^{\top}(x) y_{*}, J F_{k 1}\right\rangle+\left\langle P_{k 10}, y_{*}\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}, \\
& q_{2}=\sum_{|k|>K_{+}}\left(\left\langle y_{*}, A_{k} y\right\rangle+\left\langle y_{*}, B_{k} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}, \\
& q=h_{0}\left(x, y+y_{*}, z, \lambda\right)-h_{0}(x, y, z, \lambda) .
\end{aligned}
$$

Then

$$
Q=\left(Q_{1}+Q_{2}+Q_{3}\right) \circ \phi+q_{0}+q_{1}+q_{2}
$$

By H3), it is clear that
$\left|\partial_{\lambda}^{l} A\right|_{\mathcal{D}(r)},\left|\partial_{\lambda}^{l} B\right|_{\mathcal{D}(r)},\left|\partial_{\lambda}^{l} M\right|_{\mathcal{D}(r)}, s^{-3}\left|\partial_{\lambda}^{l} \partial_{x} h_{0}\right|_{\mathcal{D}(r)}, s^{-2}\left|\partial_{\lambda}^{l} \partial_{(y, z)} h_{0}\right|_{\mathcal{D}(r)},|U| \leq c$.
It follows from the above and (3.44)-(3.47) that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} Q_{1}\right|_{D_{3}} & \leq c \sum_{0<|k| \leq K_{+}}\left(s_{+}^{2}\left|\partial_{\lambda}^{l} f_{k 1}\right|+s_{+}|k|\left(\left|\partial_{\lambda}^{l} f_{k 0}\right|+s_{+}\left|\partial_{\lambda}^{l} f_{k 1}\right|\right.\right. \\
& \left.+s_{+}\left|\partial_{\lambda}^{l} F_{k 1}\right|\right) e^{|k|\left(r_{+}+\frac{1}{4}\left(r-r_{+}\right)\right)} \\
& \leq c s^{2} s_{+} \mu K_{+} \sum_{0<|k| \leq K_{+}}|k|^{n+(n+1) \tau+6} e^{-|k| \frac{5\left(r-r_{+}\right)}{8}} \\
& \leq c s^{2} s_{+} \mu K_{+} \Gamma\left(r-r_{+}\right)=c s^{2} s_{+} \mu \zeta
\end{aligned}
$$

and from (3.42), (3.47) and Lemma 3.3 1) that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} q_{0}\right|_{D_{1}} \leq & c \sum_{0<|k| \leq K_{+}}\left(\left|\partial_{\lambda}^{l} y_{*}\right|\left|F_{k 01}\right|+\left|y_{*}\right|\left|\partial_{\lambda}^{l} F_{k 01}\right|\right) \leq c \gamma^{n+1} s^{2} \mu^{2} K_{+} \zeta \\
\left|\partial_{\lambda}^{l} q_{1}\right|_{D_{1}} \leq & \sum_{0<|k| \leq K_{+}}\left(\left|\partial_{\lambda}^{l} y_{*}\right|\left(\left|f_{k 1}\right|+\left|F_{k 1}\right|+\left|P_{k 10}\right|\right)\right. \\
& \left.+\left|y_{*}\right|\left(\left|\partial_{\lambda}^{l} f_{k 1}\right|+\left|\partial_{\lambda}^{l} F_{k 1}\right|+\left|\partial_{\lambda}^{l} P_{k 10}\right|\right)\right) e^{|k|\left(r_{+}+\frac{1}{4}\left(r-r_{+}\right)\right)} \\
\leq & c \gamma^{n+1} s^{2} \mu^{2} \zeta K_{+} \sum_{0<|k| \leq K_{+}}|k|^{1+2 \tau} e^{-|k| \frac{3\left(r-r_{+}\right)}{4}} \leq c \gamma^{n+1} s^{2} \mu^{2} \zeta^{2}
\end{aligned}
$$

Similarly,

$$
\begin{align*}
\left|\partial_{\lambda}^{l} Q_{1}\right|_{D(s)} & \leq c s^{3} \mu \zeta \\
\left|\partial_{\lambda}^{l} q_{0}\right|_{D(s)} & \leq c \gamma^{n+1} s^{2} \mu^{2} K_{+} \zeta \\
\left|\partial_{\lambda}^{l} q_{1}\right|_{D(s)} & \leq c s^{2} \mu^{2} \zeta^{2} \\
\left|\partial_{\lambda}^{l} q\right|_{D(s)} & \leq c \gamma^{n+1} s_{+} s^{3} \mu \zeta \tag{3.55}
\end{align*}
$$

By (3.43)-(3.45) and H2), we also have

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} Q_{2}\right|_{D_{3}} \leq & c \sum_{|k|>K_{+}} s_{+}\left(\left(\left|\partial_{\lambda}^{l} B_{k}\right|+\left|\partial_{\lambda}^{l} M_{k}\right|\right)\left|F_{01}\right|\right. \\
& \left.+\left(\left|B_{k}\right|+\left|M_{k}\right|\right)\left|\partial_{\lambda}^{l} F_{01}\right|\right) e^{|k|\left(r_{+}+\frac{1}{4}\left(r-r_{+}\right)\right)} \\
\leq & c \gamma^{n+1} s s_{+} \mu K_{+} \sum_{|k|>K_{+}} e^{-|k| \frac{3\left(r-r_{+}\right)}{4}} \leq c \gamma^{n+1} s s_{+} \mu^{2} K_{+} \\
\left|\partial_{\lambda}^{l} Q_{3}\right|_{D_{3}} \leq & c \sum_{|k|>K_{+}, 0<|j| \leq K_{+}} s_{+}\left(\left(\left|\partial_{\lambda}^{l} B_{k-j}\right|+\left|\partial_{\lambda}^{l} M_{k-j}\right|\right)\left|F_{j 1}\right|\right. \\
& \left.+\left(\left|B_{k-j}\right|+\left|M_{k-j}\right|\right)\left|\partial_{\lambda}^{l} F_{j 1}\right|\right) e^{|k|\left(r_{+}+\frac{1}{4}\left(r-r_{+}\right)\right)} \\
\leq & c \gamma^{n+1} s s_{+} \mu K_{+}\left(\sum_{|k|>K_{+}} e^{-|k| \frac{\left(r-r_{+}\right)}{2}}\right)\left(\sum_{0<|j| \leq K_{+}} e^{-|j| \frac{\left(r-r_{+}\right)}{8}}\right) \\
\leq & c \gamma^{n+1} s s_{+} \mu^{2} K_{+} \Gamma\left(r-r_{+}\right)=c \gamma^{n+1} s s_{+} \mu^{2} \zeta .
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} Q_{2}\right|_{D(s)} & \leq c \gamma^{n+1} s^{2} \mu^{2} K_{+} \\
\left|\partial_{\lambda}^{l} Q_{3}\right|_{D(s)} & \leq c \gamma^{n+1} s^{2} \mu^{2} \zeta
\end{aligned}
$$

Note that, similar to (3.43),

$$
\left|\partial_{\lambda}^{l} A_{k}\right|_{\mathcal{O}} \leq\left|\partial_{\lambda}^{l} \mathcal{M}-\partial_{\lambda}^{l} \mathcal{M}_{0}\right|_{\mathcal{D}(r) \times \mathcal{O}_{0}} e^{-|k| r} \leq \gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}} e^{-|k| r}
$$

This together with (3.43), H2) and Lemma 3.3 1) implies that

$$
\begin{aligned}
\left|\partial_{\lambda}^{l} q_{2}\right|_{D_{3}} & \leq c \sum_{|k|>K_{+}} s_{+}\left(\left(\left|\partial_{\lambda}^{l} A_{k}\right|+\left|\partial_{\lambda}^{l} B_{k}\right|\right)\left|y_{*}\right|\right. \\
& \left.+\left(\left|A_{k}\right|+\left|B_{k}\right|\right)\left|\partial_{\lambda}^{l} y_{*}\right|\right) e^{|k|\left(r_{+}+\frac{1}{4}\left(r-r_{+}\right)\right)} \\
& \leq c \gamma_{0}^{n+1} s_{+} s \mu K_{+} \sum_{|k|>K_{+}} e^{-|k| \frac{3\left(r-r_{+}\right)}{4}} \leq 2^{n+1} c \gamma^{n+1} s_{+} s \mu^{2} \zeta .
\end{aligned}
$$

Similarly,

$$
\left|\partial_{\lambda}^{l} q_{2}\right|_{D(s)} \leq c \gamma^{n+1} s^{2} \mu^{2} \zeta
$$

It now follows from Lemma 3.3 1) that

$$
\begin{align*}
\left|\partial_{\lambda}^{l} Q\right|_{D_{+}} & \leq c\left(s^{2} s_{+} \mu+\gamma^{n+1} s^{2} \mu^{2}\right) \zeta^{2}  \tag{3.56}\\
\left|\partial_{\lambda}^{l} Q \circ \phi^{-1}\right|_{D(s)} & \leq c\left(s^{3} \mu+\gamma^{n+1} s^{2} \mu^{2}\right) \zeta^{2} \tag{3.57}
\end{align*}
$$

Applying Cauchy's estimate to (3.57), we then have

$$
\begin{align*}
\left|\partial_{\lambda}^{l} D\left(Q \circ \phi^{-1}\right)\right|_{D_{3}} & \leq \frac{c}{\left(r-r_{+}\right)\left(s-3 s_{+}\right)}\left|\partial_{\lambda}^{l}\left(Q \circ \phi^{-1}\right)\right|_{D(s)} \\
& \leq \frac{c}{r-r_{+}}\left(s^{2} \mu+\gamma^{n+1} s \mu^{2}\right) \zeta^{2} \tag{3.58}
\end{align*}
$$

Denote

$$
\begin{aligned}
& W_{0}=\int_{0}^{1}\left\{R_{t}, F\right\} \circ \phi_{F}^{t} \mathrm{~d} t \\
& W_{1}=\frac{1}{2}\left\langle y_{*},(A-2 \operatorname{diag}(U, O)) y_{*}\right\rangle \\
& +\sum_{|k| \leq K_{+}}\left(\left\langle y_{*}, P_{k 20} y_{*}\right\rangle+\left\langle y_{*}, 2 P_{k 20} y\right\rangle+\left\langle y_{*}, P_{k 11} z\right\rangle\right) e^{\sqrt{-1}\langle k, x\rangle}
\end{aligned}
$$

Then by (3.14),

$$
\begin{equation*}
P_{+}=W_{0} \circ \phi+W_{1}+Q+q+(P-R) \circ \Phi_{+} . \tag{3.59}
\end{equation*}
$$

Since by (3.6) and (3.7)

$$
\begin{aligned}
R_{t} & =(1-t)\{N, F\}+R \\
& =(1-t)\left(\left(Q-\left\langle y_{*},(A-[A]) y\right\rangle-\left\langle y_{*}, B z\right\rangle\right) \circ \phi^{-1}-\left(R-R^{\prime}\right)\right)+R
\end{aligned}
$$

we have
$W_{0}=\int_{0}^{1}\left\{(1-t)\left(\left(Q-\left\langle y_{*},(A-[A]) y\right\rangle-\left\langle y_{*}, B z\right\rangle\right) \circ \phi^{-1}-R+R^{\prime}\right)+R, F\right\} \circ \phi_{F}^{t} \mathrm{~d} t$.
Using Lemma 3.3 1), 2) and Lemma 3.4 3), we have by (3.58) that

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} \int_{0}^{1}\left\{Q \circ \phi^{-1}, F\right\} \circ \phi_{F}^{t} \mathrm{~d} t \circ \phi\right|_{D_{+}} \leq \frac{c}{r-r_{+}}\left(s^{3} \mu^{2}+\gamma^{n+1} s^{2} \mu^{3}\right) \zeta^{3}, \tag{3.60}
\end{equation*}
$$

and, by (3.53) and (3.54) that
(3.61) $\left(\left|\partial_{\lambda}^{l} \int_{0}^{1}\{R, F\} \circ \phi_{F}^{t} \mathrm{~d} t \circ \phi\right|+\left|\partial_{\lambda}^{l} \int_{0}^{1}\left\{R^{\prime}, F\right\} \circ \phi_{F}^{t} \mathrm{~d} t \circ \phi\right|\right)_{D_{+}} \leq c \gamma^{n+1} s^{2} \mu^{2} \zeta^{3}$.

By Lemma 3.3, it is also clear that

$$
\left|\partial_{\lambda}^{l} \int_{0}^{1}\left\{\left(\left\langle y_{*},(A-[A]) y\right\rangle+\left\langle y_{*}, B z\right\rangle\right) \circ \phi^{-1}, F\right\} \circ \phi_{F}^{t} \mathrm{~d} t \circ \phi\right|_{D_{+}} \leq c \gamma^{n+1} s^{2} \mu^{2} \zeta^{3} .
$$

This together with (3.60) and (3.61) yields that

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} W_{0} \circ \phi\right|_{D_{+}} \leq c \frac{\mu^{2}}{r-r_{+}}\left(s^{3}+\gamma^{n+1} s^{2}\right) \zeta^{3} \tag{3.62}
\end{equation*}
$$

Using the above arguments along with (3.42) and Lemma 3.3, we also have

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} W_{1}\right|_{D_{+}} \leq c \gamma^{n+1}\left(s_{+}+s\right) s \mu^{2} \zeta^{2} \tag{3.63}
\end{equation*}
$$

Above all, it follows from (3.52), (3.55), (3.56), (3.59), (3.62), (3.63) that

$$
\left|\partial_{\lambda}^{l} P_{+}\right|_{D_{+}} \leq c \frac{\zeta^{3}}{r-r_{+}}\left(s_{+} s^{2} \mu+\gamma^{n+1} s^{2} \mu^{2}+s^{3} \mu^{2}+\gamma^{n+1} \frac{s_{+}^{3}}{s} \mu\right) \leq c \Delta
$$

## 4. Iteration Lemma

Let $r_{0}, s_{0}, \mu_{0}, \gamma_{0}, \mathcal{O}_{0}, H_{0}, N_{0}, e_{0}, \Omega_{0}, \mathcal{M}^{0}, A^{0}, B^{0}, M^{0}, \mathcal{A}^{0}, h_{0}, P_{0}$ be defined as in Section 3 and let $\tilde{D}_{0}=D\left(r_{0}, \beta_{0}\right), D_{0}=D\left(r_{0}, s_{0}\right), \mathcal{D}_{0}=\left\{x:|\operatorname{Im} x|<r_{0}\right\}, K_{0}=0$, $\Phi_{0}=i d$. For $\nu=1,2 \cdots$, we index all index-free quantities in Section 3 by $\nu$ and index all" + "-indexed quantities in Section 3 by $\nu+1$. This yields the following sequences with the properties stated in Section 3:

$$
\begin{aligned}
H_{\nu} & =H_{\nu}(x, y, z, \lambda)=N_{\nu}+P_{\nu} \\
N_{\nu} & =e_{\nu}+\left\langle\Omega_{\nu}, y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}^{\nu}\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda) \\
\mathcal{M}^{\nu} & =\left(\begin{array}{cc}
A^{\nu} & B^{\nu} \\
\left(B^{\nu}\right)^{\top} & M^{\nu}
\end{array}\right) \\
r_{\nu} & =r_{0}\left(1-\frac{1}{2}(1-\delta) \sum_{i=1}^{\nu} \delta^{i+1}\right) \\
s_{\nu} & =s_{\nu-1}^{1+b+\sigma}, \\
\beta_{\nu} & =\beta_{0}\left(1-\sum_{i=1}^{\nu} \frac{1}{2^{i+1}}\right) \\
\gamma_{\nu} & =\gamma_{0}\left(1-\sum_{i=0}^{\nu-1} \frac{1}{2^{i+2}}\right) \\
\mu_{\nu} & =c_{0} s_{\nu-1}^{\sigma} \mu_{\nu-1}, c_{0}=\max \left\{1, c_{1}, \cdots, c_{7}\right\} \\
K_{\nu} & =\left(\left[\log \frac{1}{s_{\nu-1}}\right]+1\right)^{3}, \nu \geq 1 \\
\mathcal{O}_{\nu} & =\left\{\lambda \in \mathcal{O}_{\nu-1}:\left|\left\langle k, \Omega_{\nu-1}(\lambda)\right\rangle\right|>\frac{\gamma_{\nu-1}}{|k|^{\tau}}, 0<|k| \leq K_{\nu}\right\}, \nu \geq 1 \\
D_{\nu} & =D\left(r_{\nu}, s_{\nu}\right) \\
\tilde{D}_{\nu} & =D\left(r_{\nu}+\frac{7}{8}\left(r_{\nu-1}-r_{\nu}\right), \beta_{\nu}\right) \\
\mathcal{D}_{\nu} & =\left\{x:|\operatorname{Im} x|<r_{\nu}\right\}
\end{aligned}
$$

$\nu=0,1, \cdots$.
We now prove an iteration lemma which checks the validity of all KAM steps.
Lemma 4.1. (Iteration Lemma) If $\mu_{*}=\mu_{*}\left(r_{*}, \beta_{0}, \sigma_{0}, U^{0}\right)$ is sufficiently small, then the following holds for all $|l| \leq n ; \nu=1,2, \cdots$.
1)

$$
\begin{align*}
& \left|\partial_{\lambda}^{l}\left(e_{\nu}-e_{0}\right)\right|_{\mathcal{O}_{\nu}} \leq \gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}  \tag{4.1}\\
& \left|\partial_{\lambda}^{l}\left(e_{\nu}-e_{\nu-1}\right)\right|_{\mathcal{O}_{\nu}} \leq \frac{\gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}}{2^{\nu}},  \tag{4.2}\\
& \left|\partial_{\lambda}^{l}\left(\Omega_{\nu}-\Omega_{0}\right)\right|_{\mathcal{O}_{\nu}} \leq \gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}  \tag{4.3}\\
& \left|\partial_{\lambda}^{l}\left(\Omega_{\nu}-\Omega_{\nu-1}\right)\right|_{\mathcal{O}_{\nu}} \leq \frac{\gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}}{2^{\nu}}  \tag{4.4}\\
& \left|\partial_{\lambda}^{l}\left(\mathcal{M}^{\nu}-\mathcal{M}^{0}\right)\right|_{\mathcal{D}_{\nu} \times \mathcal{O}_{\nu}} \leq \gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}}  \tag{4.5}\\
& \left|\partial_{\lambda}^{l}\left(\mathcal{M}^{\nu}-\mathcal{M}^{\nu-1}\right)\right|_{\mathcal{D}_{\nu} \times \mathcal{O}_{\nu}} \leq \frac{\gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}}}{2^{\nu}}  \tag{4.6}\\
& \left|\partial_{\lambda}^{l} P_{\nu}\right|_{D_{\nu} \times \mathcal{O}_{\nu}} \leq \gamma_{\nu}^{n+1} s_{\nu}^{2} \mu_{\nu} \tag{4.7}
\end{align*}
$$

2) $\left(\Omega_{\nu}(\lambda)\right)_{i}=\omega_{i}(\lambda), i=1,2, \cdots, n_{0}$.
3) There is a transformation $\Phi_{\nu}: \tilde{D}_{\nu} \times \mathcal{O}_{\nu} \longrightarrow \tilde{D}_{\nu-1}, D_{\nu} \times \mathcal{O}_{\nu} \longrightarrow D_{\nu-1}$, which is symplectic and analytic in $(x, y, z) \in \tilde{D}_{\nu+1}$, and smooth in $\lambda \in$ $\mathcal{O}_{\nu+1}$, such that

$$
H_{\nu}=H_{\nu-1} \circ \Phi_{\nu}
$$

and

$$
\begin{equation*}
\left|\partial_{\lambda}^{l} D^{i}\left(\Phi_{\nu}-i d\right)\right|_{\tilde{D}_{\nu} \times \mathcal{O}_{\nu}} \leq \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu}}, 0 \leq i \leq 3 \tag{4.8}
\end{equation*}
$$

4) 

$$
\mathcal{O}_{\nu+1}=\left\{\lambda \in \mathcal{O}_{\nu}:\left|\left\langle k, \Omega_{\nu}(\lambda)\right\rangle\right|>\frac{\gamma_{\nu}}{|k|^{\tau}}, K_{\nu}<|k| \leq K_{\nu+1}\right\} .
$$

Proof. The lemma will be proved by induction. We first verify the conditions H1), H2), H5)-H10) in Section 3 for all $\nu=0,1, \cdots$.

By definitions of $\mu_{\nu}, s_{\nu}$, we clearly have

$$
\begin{align*}
\mu_{\nu} & =c_{0}^{\nu} \mu_{0} s_{0}^{\frac{\sigma}{b+\sigma}\left((1+b+\sigma)^{\nu}-1\right)}  \tag{4.9}\\
s_{\nu} & =s_{0}^{(1+b+\sigma)^{\nu}}  \tag{4.10}\\
s_{0} & =\gamma^{n+1+\frac{a_{0}}{2}} \mu_{*}^{\frac{1}{4}} . \tag{4.11}
\end{align*}
$$

Define

$$
\begin{aligned}
E_{\nu} & =\frac{r_{\nu}-r_{\nu+1}}{8}=\frac{1}{16} r_{0}(1-\delta) \delta^{\nu+2} \\
\chi & =n+6+(n+1)[\tau]
\end{aligned}
$$

and let $\varepsilon>0$ be fixed such that

$$
\begin{aligned}
& \left(n+1+a_{0}\right)(1-2 b-3 \sigma-2 \varepsilon)>(n+1) \\
& \sigma>(b+\sigma)(2 b+3 \sigma+2 \varepsilon)+\varepsilon \\
& 3 \varepsilon<b
\end{aligned}
$$

By (4.9)-(4.11), it is clear that if $\mu_{*}$ (hence $s_{0}$ ) is small, then

$$
\begin{align*}
c_{0} \frac{\mu_{\nu}}{s_{\nu}^{2 b+3 \sigma+2 \varepsilon}} & =\frac{c_{0}^{\nu+1} \mu_{0} s_{0}^{\frac{\sigma}{b+\sigma}\left((1+b+\sigma)^{\nu}-1\right)}}{s_{0}^{(1+b+\sigma)^{\nu}(2 b+3 \sigma+2 \varepsilon)}} \\
& =c_{0} s_{0}^{-\frac{\sigma}{b+\sigma}} \mu_{0} c_{0}^{\nu} s_{0}^{\left(\frac{\sigma}{b+\sigma}-(2 b+3 \sigma+2 \varepsilon)\right)(1+b+\sigma)^{\nu}} \\
& \leq s_{0}^{-\frac{\sigma}{b+\sigma}} \mu_{0}\left(c_{0} s_{0}^{\varepsilon}\right)^{\nu} \leq s_{0}^{-\frac{\sigma}{b+\sigma}} \mu_{0} \leq \mu_{*}^{\frac{1}{4}} \tag{4.12}
\end{align*}
$$

Since

$$
\begin{aligned}
\Gamma_{\nu} & =\Gamma\left(r_{\nu}-r_{\nu+1}\right)=e^{E_{0}} \sum_{0<|k| \leq K_{n u+1}}|k|^{n+(n+1) \tau+6} e^{-|k| E_{\nu}} \\
& \leq e^{E_{0}} \int_{1}^{\infty} \lambda^{\chi} e^{-\lambda \frac{E_{\nu}}{2}} \mathrm{~d} \lambda \leq \frac{e^{E_{0}} 2^{\chi}(\chi+1)!}{E_{\nu}^{\chi}}
\end{aligned}
$$

we can make $\mu_{*}$ (hence $s_{0}$ ) small such that

$$
c_{0} \frac{s_{\nu}^{\varepsilon} \Gamma_{\nu}^{3}}{E_{\nu}} \leq c_{0}\left(e^{E_{0}} 2^{\chi}(\chi+1)!\right)^{3} \frac{s_{\nu}^{\varepsilon}}{E_{\nu}^{3 \chi+1}} \leq \frac{\left.c_{0} e^{3 E_{0}} 2^{3 \chi}(\chi+1)!\right)^{3}}{E_{0}^{3 \chi+1}}\left(\frac{s_{0}^{(b+\sigma) \varepsilon}}{\delta}\right)^{\nu} s_{0} \leq 1
$$

On the other hand, we can also make $\mu_{*}$ (hence $s_{\nu}$ ) small such that

$$
s_{\nu}^{\varepsilon} K_{\nu+1}^{3}=s_{\nu}^{\varepsilon}\left(\left[\log \frac{1}{s_{\nu}}\right]+1\right)^{9} \leq 1
$$

Therefore,

$$
\begin{equation*}
c_{0} \frac{s_{\nu}^{2 \varepsilon} \zeta_{\nu}^{3}}{E_{\nu}}=\left(\frac{c_{0} s_{\nu}^{\varepsilon} \Gamma_{\nu}^{3}}{E_{\nu}}\right)\left(s_{\nu}^{\varepsilon} K_{\nu+1}^{3}\right) \leq 1 \tag{4.13}
\end{equation*}
$$

By (4.9) and (4.13), we can again make $\mu_{*}$ (hence $s_{0}$ ) small such that

$$
\begin{align*}
c_{0} \mu_{\nu} \zeta_{\nu}^{3} & =c_{0} \frac{\mu_{\nu}}{s_{\nu}^{2 \varepsilon}}\left(s_{\nu}^{2 \varepsilon} \zeta_{\nu}^{3}\right) \leq c_{0} \frac{\mu_{\nu}}{s_{\nu}^{2 \varepsilon}} \leq \mu_{*}^{\frac{1}{4}} s_{\nu}^{2 b+3 \sigma} \\
& =\mu_{*}^{\frac{1}{4}} s_{0}^{(2 b+3 \sigma)(1+b+\sigma)^{\nu}} \leq \mu_{*}^{\frac{1}{4}}\left(s_{0}^{(2 b+3 \sigma)(b+\sigma)}\right)^{\nu} \leq \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu+1}} \tag{4.14}
\end{align*}
$$

Using the definition of $s_{\nu}$, we clearly have

$$
s_{\nu+1} \leq s_{0}^{b+\sigma} s_{\nu} \leq \frac{s_{\nu}}{16}
$$

i.e., H1) holds.

By making $\mu_{*}$ small, we can apply (4.10) and (4.11) to make $s_{\nu}$ small such that

$$
\begin{aligned}
& \log (n+1)!+3 n \log \left(\left[\log \frac{1}{s_{\nu}}\right]+1\right)-\frac{E_{\nu}}{2}\left(\left[\log \frac{1}{s_{\nu}}\right]+1\right)^{3}+(n+1)\left(\log 2+\log E_{\nu}\right) \\
& \leq \log (n+1)!+3 n \log \left(\log \frac{1}{s_{\nu}}+2\right)-\left(\log \frac{1}{s_{\nu}}\right)^{2} \\
& +(n+1)\left(\log 2+\log \left(2^{-4} r_{0}(1-\delta) \delta^{2}\right)+\nu \log \delta\right) \\
& \leq-\log \frac{1}{s_{\nu}}
\end{aligned}
$$

Hence,

$$
\int_{K_{\nu+1}}^{\infty} \lambda^{n} e^{-\lambda \frac{E_{\nu}}{2}} \mathrm{~d} \lambda \leq 2^{n}(n+1)!\left(\frac{K_{\nu+1}}{E_{\nu}}\right)^{n} e^{-K_{\nu+1} \frac{E_{\nu}}{2}} \leq s_{\nu}
$$

This verifies H2).

Applying (4.12) and (4.13), we have that

$$
\begin{aligned}
\frac{c_{0} \mu_{\nu} \zeta_{\nu}}{E_{\nu}} & =\frac{\mu_{\nu}}{s_{\nu}^{2 \varepsilon}} \frac{c_{0} s_{\nu}^{2 \varepsilon} \zeta_{\nu}}{E_{\nu}} \leq \frac{\mu_{\nu}}{s_{\nu}^{2 \varepsilon}} \leq s_{\nu}^{2 b+3 \sigma}<1 \\
\frac{c_{0} s_{\nu} \mu_{\nu} \zeta_{\nu}}{s_{\nu+1}} & =\frac{\mu_{\nu}}{s_{\nu}^{b+\sigma+2 \varepsilon}}\left(c_{0} s_{\nu}^{2 \varepsilon} \zeta_{\nu}\right) \leq \frac{\mu_{\nu}}{s_{\nu}^{b+\sigma+2 \varepsilon}}<1
\end{aligned}
$$

which verify H6) and H7) respectively. H9) can be verified similarly.
Since

$$
\beta_{\nu}-\beta_{\nu+1}=\frac{\beta_{0}}{2^{\nu+2}}
$$

it follows from (4.14) that

$$
\frac{c_{0} \mu_{\nu} \zeta_{\nu}}{\beta_{\nu}-\beta_{\nu+1}} \leq 2 \frac{\mu_{*}^{\frac{1}{4}}}{\beta_{0}}<\frac{1}{2}
$$

as $\mu_{*}<\frac{\beta_{0}^{4}}{2^{8}}$. This verifies H8).
Applying (4.12) and (4.13) again, we have that

$$
\begin{aligned}
\frac{c_{0} \Delta_{\nu}}{\gamma_{\nu+1}^{n+1} s_{\nu+1}^{2} \mu_{\nu+1}} & =c_{0} \frac{\zeta_{\nu}^{3}}{4 E_{\nu}}\left(\frac{s_{\nu}^{3} \mu_{\nu}}{\gamma_{\nu}^{n+1} s_{\nu+1}^{2} \mu_{\nu+1}}+\frac{s_{\nu}^{2} \mu_{\nu}^{2}}{s_{\nu+1}^{2} \mu_{\nu+1}}+\frac{s_{\nu+1} \mu_{\nu}}{s_{\nu} \mu_{\nu+1}}\right) \\
& =\frac{1}{4}\left(\frac{c_{0} s_{\nu}^{2 \varepsilon} \zeta_{\nu}^{3}}{E_{\nu}}\right)\left(\frac{s_{\nu}^{1-2 b-3 \sigma-2 \varepsilon}}{c_{0} \gamma_{\nu}^{n+1}}+\frac{\mu_{\nu}}{c_{0} s_{\nu}^{2 b+3 \sigma+2 \varepsilon}}+\frac{s_{\nu}^{b-2 \varepsilon}}{c_{0}}\right) \\
& \leq \frac{1}{4}\left(\frac{c_{0} s_{\nu}^{2 \varepsilon} \zeta_{\nu}^{3}}{E_{\nu}}\right)\left(\frac{s_{0}^{1-2 b-3 \sigma-2 \varepsilon}}{\gamma_{\nu}^{n+1}}+2\right)<1
\end{aligned}
$$

i.e., the first part of H10) holds. The second part of H10) is obvious.

Next, we verify H3)-H5) by induction. For each $\nu=0,1, \cdots$, we define $M_{k j}^{\nu}$, $|k|,|j|>0$, as in (3.28) with $M_{k}^{\nu}, B_{k}^{\nu}, U^{\nu}, Y^{\nu}$ in place of $M_{k}, B_{k}, U, Y$, respectively, for all $|k|>0$. Using the same ordering $\mathcal{T}$ as in Section 3, we also have the matrices

$$
\mathcal{A}^{\nu}=\left(e^{(|k|-|j|) \xi\left(r_{\nu}\right)} M_{k j}^{\nu}\right), \mathcal{A}^{0}=\left(e^{(|k|-|j|) \xi\left(r_{\nu}\right)} M_{k j}^{0}\right), 0<|k|,|j| \leq K_{\nu+1}
$$

where $\xi$ is as in (3.38).
Clearly, H3)-H5) trivially hold for $\nu=0$. We now assume that for some positive integer $\left.\nu_{*} \mathrm{H} 3\right)-\mathrm{H} 5$ ) hold for all $\nu=1, \cdots, \nu_{*}$. Then the KAM step described in Section 3 is valid for all $\nu=1, \cdots, \nu_{*}$. In particular, Lemmas 3.5, 3.6 hold for all
$\nu=1, \cdots, \nu_{*}$. It follows from (4.14) that

$$
\begin{aligned}
& \left|\partial_{\lambda}^{l}\left(\Omega_{\nu_{*}+1}-\Omega_{0}\right)\right| \leq \sum_{\nu=0}^{\nu_{*}}\left|\partial_{\lambda}^{l}\left(\Omega_{\nu+1}-\Omega_{\nu}\right)\right|_{\mathcal{O}_{\nu+1}} \leq \sum_{\nu=0}^{\nu_{*}} c_{0} \gamma_{0}^{n+1} s_{\nu} \mu_{\nu} \zeta_{\nu} \\
& \leq \gamma_{0}^{n+1} s_{0} \sum_{\nu=0}^{\nu_{*}} \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu+1}} \leq \gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}<\mu_{*}^{\frac{1}{4}}, \\
& \left|\partial_{\lambda}^{l}\left(\mathcal{M}^{\nu_{*}+1}-\mathcal{M}^{0}\right)\right|_{D_{\nu_{*}+1} \times \mathcal{O}_{\nu_{*}+1}} \leq \sum_{\nu=0}^{\nu_{*}}\left|\partial_{\lambda}^{l}\left(\mathcal{M}_{\nu+1}-\mathcal{M}_{\nu}\right)\right|_{D_{\nu+1} \times \mathcal{O}_{\nu+1}} \\
& \leq \sum_{\nu=0}^{\nu_{*}} c_{0} \gamma_{0}^{n+1} \mu_{\nu} \zeta_{\nu} \leq \gamma_{0}^{n+1} \sum_{\nu=0}^{\nu_{*}} \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu+1}} \leq \gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}}<\mu_{*}^{\frac{1}{4}}, \\
& \max _{|k|>0} \sum_{|j|>0} e^{(|k|-|j|) \xi\left(r_{\nu_{*}+1}\right)}\left|M_{k j}^{\nu_{*}+1}-M_{k j}^{0}\right|_{\mathcal{O}_{\nu_{*}+1}} \\
& \leq \sum_{\nu=0}^{\nu_{*}} \max _{|k|>0} \sum_{|j|>0} e^{(|k|-|j|) \xi\left(r_{\nu+1}\right)}\left|M_{k j}^{\nu+1}-M_{k j}^{\nu}\right| \mathcal{O}_{\nu+1} \\
& \leq \sum_{\nu=0}^{\nu_{*}} c_{0} \gamma_{0}^{n+1} \mu_{\nu} \zeta_{\nu}^{2} \leq \gamma_{0}^{n+1} \sum_{\nu=0}^{\nu_{*}} \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu+1}}<\mu_{*}^{\frac{1}{4}} .
\end{aligned}
$$

Thus, H3)-H5) also hold for $\nu=\nu_{*}+1$.
Above all, H 1$)-\mathrm{H} 10$ ) hold for all $\nu=0,1, \cdots$, i.e., the KAM step described in Section 3 is valid for all $\nu=0,1, \cdots$. Now, part 4) of the lemma easily follows from Lemma 3.6. Also, using (4.14), part 3) of the lemma follows from Lemma 3.4 ; (4.7) follows from Lemma 3.8; (4.2) and (4.4) follow from Lemma 3.51 ); (4.6) follows from Lemma 3.7. Finally, (4.1), (4.3) and (4.5) follow from (4.2), (4.4) and (4.6) respectively, and, part 2) of the lemma follows from an inductive application of (3.10). This completes the proof.

## 5. Proof of Theorems 1 and 2

We first consider Theorem 2 1). By making $\mu_{*}=\mu_{*}(r, s)$ small in Theorem 2, we obtain the following sequences

$$
\begin{aligned}
& \Psi^{\nu}=\Phi_{0} \circ \Phi_{1} \circ \cdots \circ \Phi_{\nu}: \tilde{D}_{\nu+1} \times \mathcal{O}_{\nu+1} \rightarrow \tilde{D}_{0} \\
& H \circ \Psi^{\nu}=H_{\nu}=N_{\nu}+P_{\nu} \\
& N_{\nu}=e_{\nu}+\left\langle\Omega_{\nu}, y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}^{\nu}\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda)
\end{aligned}
$$

$\nu=0,1, \cdots$, which satisfy all properties described in Lemma 4.1.
Let

$$
\mathcal{O}_{*}=\prod_{\nu=0}^{\infty} \mathcal{O}_{\nu}, G_{*}=D\left(\frac{r_{0}}{2}, \frac{\beta_{0}}{2}\right) \times \mathcal{O}_{*}
$$

Then $\mathcal{O}_{*}$ is a Cantor-like set consisting of non-resonant frequencies, and moreover, a measure estimate similar to that in [24] (also [4, 13]) yields that $\left|\mathcal{O} \backslash \mathcal{O}_{*}\right|=O\left(\gamma_{0}^{\frac{1}{n_{*}-1}}\right)$.

By Lemma 4.11 ), it is clear that $e_{\nu}$ and $\Omega_{\nu}$ converge uniformly on $\mathcal{O}_{*}$, and, $\mathcal{M}^{\nu}$ converges uniformly on $\mathcal{D}\left(\frac{r_{0}}{2}\right) \times \mathcal{O}_{*}$. We denote their limits by $e_{\infty}, \Omega_{\infty}, \mathcal{M}^{\infty}$,
respectively. Then $\mathcal{M}^{\infty}$ is real analytic in $x$, and, it follows from the Whitney's extension theorem ([17]) that these limits are also Hölder continuous in $\omega$. Moreover, by Lemma 4.1 1), we have that

$$
\begin{aligned}
\left|e_{\infty}-e_{0}\right|_{\mathcal{O}_{*}} & =O\left(\gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}\right), \\
\left|\Omega_{\infty}-\Omega_{0}\right|_{\mathcal{O}_{*}} & =O\left(\gamma_{0}^{2(n+1)} \mu_{*}^{\frac{1}{2}}\right), \\
\left.\left|\mathcal{M}^{\infty}-\mathcal{M}^{0}\right|_{\mathcal{D}\left(\frac{r_{0}}{2}\right.}\right) \times \mathcal{O}_{*} & =O\left(\gamma_{0}^{n+1} \mu_{*}^{\frac{1}{4}}\right)
\end{aligned}
$$

Thus, $N_{\nu}$ converges uniformly on $G_{*}$ to

$$
N_{\infty}=e_{\infty}+\left\langle\Omega_{\infty}, y\right\rangle+\frac{1}{2}\left\langle\binom{ y}{z}, \mathcal{M}^{\infty}\binom{y}{z}\right\rangle+h_{0}(x, y, z, \lambda) .
$$

To show the convergence of $\Psi^{\nu}$ on $G_{*}$, we note that

$$
\begin{aligned}
\Psi^{\nu}-\Psi^{\nu-1} & =\Phi_{0} \circ \cdots \circ \Phi_{\nu}-\Phi_{0} \circ \cdots \circ \Phi_{\nu-1} \\
& =\int_{0}^{1} D\left(\Phi_{0} \circ \cdots \circ \Phi_{\nu-1}\right)\left(i d+\theta\left(\Phi_{\nu}-i d\right)\right) \mathrm{d} \theta\left(\Phi_{\nu}-i d\right)
\end{aligned}
$$

It follows from Lemma 4.1 3) that

$$
\begin{aligned}
& \left|D\left(\Phi_{1} \circ \cdots \circ \Phi_{\nu-1}\right)\left(i d+\theta\left(\Phi_{\nu}-i d\right)\right)\right| \\
\leq & \left.\mid D \Phi_{1}\left(\Phi_{2} \circ \cdots \circ \Phi_{\nu-1}\right)\left(i d+\theta\left(\Phi_{\nu}-i d\right)\right)\right)|\cdots| D \Phi_{\nu-1}\left(i d+\theta\left(\Phi_{\nu}-i d\right)\right) \mid \\
\leq & \left(1+\frac{\mu_{*}^{\frac{1}{4}}}{2}\right) \cdots\left(1+\frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu-1}}\right) \leq e^{\frac{\mu_{*}^{\frac{1}{4}}}{2}+\cdots+\frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu-1}}} \leq e^{\mu_{*}^{\frac{1}{4}}}
\end{aligned}
$$

and

$$
\left|\Phi_{\nu}-i d\right|_{G_{*}} \leq \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu}}
$$

Hence,

$$
\left|\Psi^{\nu}-\Psi^{\nu-1}\right|_{G_{*}} \leq e \frac{\mu_{*}^{\frac{1}{4}}}{2^{\nu}}
$$

which implies the uniform convergence of $\Psi^{\nu}$. Let $\Psi^{\infty}$ be the limit of $\Psi^{\nu}$. Then, $\Psi^{\infty}$ is uniformly continuous in $\lambda \in \mathcal{O}_{*}$ and analytic in $(x, y, z) \in D\left(\frac{r_{0}}{2}, \frac{\dot{\beta_{0}}}{2}\right)$, and moreover,

$$
\left|\Psi^{\infty}-i d\right|_{G_{*}}=O\left(\mu_{*}^{\frac{1}{4}}\right)
$$

Using Lemma 4.1 3) and a similar argument as above, we can further show the uniform convergence of $D \Psi^{\nu}, D^{2} \Psi^{\nu}$ to $D \Psi^{\infty}, D^{2} \Psi^{\infty}$, respectively, on $G_{*}$. Hence, on $G_{*}$,

$$
P_{\nu}=H \circ \Psi^{\nu}-N_{\nu}, D P_{\nu}, D^{2} P_{\nu}
$$

converge uniformly to

$$
P_{\infty}=H \circ \Psi^{\infty}-N_{\infty}, D P_{\infty}, D^{2} P_{\infty}
$$

respectively. Clearly, these limits above are uniformly continuous in $\lambda \in \mathcal{O}_{*}$ and analytic in $(x, y, z) \in D\left(\frac{r_{0}}{2}, \frac{\beta_{0}}{2}\right)$.

Note that

$$
\left|P_{\nu}\right|_{D_{\nu}} \leq \gamma_{\nu}^{n+1} s_{\nu}^{2} \mu_{\nu}
$$

It follows from Cauchy's estimate that, for any $\lambda \in \mathcal{O}_{*}, j \in Z_{+}^{n}, k \in Z_{+}^{2 m}$ with $|j|+|k| \leq 2$,

$$
\left|\partial_{y}^{j} \partial_{z}^{k} P_{\nu}\right|_{D\left(r_{\nu+i}, \frac{1}{2} s_{\nu}\right)} \leq \gamma_{\nu}^{n+1} \mu_{\nu}
$$

Since, by (4.9), the right hand side of the above converges to 0 as $\nu \rightarrow 0$, we have that

$$
\left.\partial_{y}^{j} \partial_{z}^{k} P_{\infty}\right|_{(y, z)=0}=0
$$

for all $x \in T^{n}, \lambda \in \mathcal{O}_{*}, j \in Z_{+}^{n}, k \in Z_{+}^{2 m}$ with $|j|+|k| \leq 2$.
Thus, for each $\lambda \in \mathcal{O}_{*}$, the perturbed system (1.3) possesses an analytic, quasiperiodic, invariant torus with the Diophantine toral frequency $\Omega_{\infty}(\lambda)$. Since, by Lemma 4.1 2),

$$
\left(\Omega_{\nu}(\lambda)\right)_{i}=\omega_{i}(\lambda), \quad \lambda \in \Omega_{\nu}, i=1,2, \cdots, n_{0}
$$

we have that

$$
\left(\Omega_{\infty}(\lambda)\right)_{i}=\omega_{i}(\lambda), \lambda \in \mathcal{O}_{*}, i=1,2, \cdots, n_{0}
$$

i.e., the perturbed toral frequencies also preserve the first $n_{0}$ components of their corresponding ones. This proves part 1) of Theorem 2.

Theorem 1 and part 2) of Theorem 2 are almost immediate consequences of part 1) of Theorem 2 with respect to $n_{0}=0$ and $n_{0}=n$ respectively. In the case that $n_{0}=0$, we can modify the arguments in Section 3 by choosing $y_{*}=0$, $U^{0}=U=\emptyset, \operatorname{diag}\left(I_{n_{0}}, O\right)=O$, and, $\operatorname{diag}\left(O, I_{n-n_{0}}\right)=I_{n}$. In the case that $n_{0}=n$, we let $\hat{\mathcal{O}}_{\gamma}$ be the set of $\lambda \in \mathcal{O}_{0}$ such that $\Omega_{0}(\lambda)$ is Diophantine of Diophantine type $(\gamma, \tau)$ for a fixed $\tau>n-1$. It is clear that $U^{0}=\left[A^{0}\right], U=[A], \operatorname{diag}\left(I_{n_{0}}, O\right)=I_{n}$, and, $\operatorname{diag}\left(O, I_{n-n_{0}}\right)=O$. Hence, $\Omega_{\nu} \equiv \Omega_{0}$ and $\mathcal{O}_{\nu}=\hat{\mathcal{O}}_{\gamma}$ for all $\nu=0,1, \cdots$, i.e., $\Omega_{\infty} \equiv \Omega_{0}$ and $\mathcal{O}_{*}=\hat{\mathcal{O}}_{\gamma}$.

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