ENERGY TRANSFER MODEL FOR THE DERIVATIVE NONLINEAR SCHröDINGER EQUATIONS ON THE TORUS

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Abstract. We consider the nonlinear derivative Schrödinger equation with a quintic nonlinearity, on the one dimensional torus. We exhibit that the nonlinear dynamic properties of the particular solution consisting of four frequency modes initially excited, whose frequencies include the resonant clusters and phase matched resonant interactions of nonlinearities. The proof is based on the analysis of resonant dynamics via a finite dimensional ordinary differential system.

1. Introduction. In this paper we consider the derivative nonlinear Schrödinger equation with quintic nonlinearity:

\[
\begin{cases}
    i\partial_t u + \partial_x^2 u = -i\lambda u^2 \partial_x u + \mu |u|^4 u, \quad (t, x) \in \mathbb{R} \times \mathbb{T}, \\
    u(0, x) = u_0(x), \quad x \in \mathbb{T},
\end{cases}
\]

(1)

where \( u = u(t, x) : \mathbb{R} \times \mathbb{T} \to \mathbb{C}, \lambda, \mu \in \mathbb{R} \), and \( \mathbb{T} = \mathbb{R}/2\pi \mathbb{Z} \) is the torus. Our aim is to prove that there exist solutions of (1) which initially oscillate, and illustrate the energy transfer from one coupled oscillator to another.

The derivative nonlinear Schrödinger equation in the completely integrable form is written as

\[
i\partial_t u + \partial_x^2 u = i\lambda \partial_x(|u|^2 u).
\]

(2)

As in [12], by using the gauge transform

\[
v(t, x) = u(t, x)e^{-i\mathcal{G}[u]},
\]

\[
\mathcal{G}[u](t, x) = \frac{\lambda}{2\pi} \int_0^{2\pi} \int_{\mathbb{T}} \left( |u(t, y)|^2 - \frac{1}{2\pi} |u(t)||^2 \right) dy d\theta,
\]

we derive an equivalent equation of (2) for \( v \) as follows:

\[
i\partial_t v + \partial_x^2 v - i\frac{\lambda}{\pi} ||v(t)||^2 \partial_x^2 v = -i\lambda \partial_x v - \frac{\lambda^2}{2} |v|^4 v
\]

\[
- \frac{\lambda}{\pi} \left( \int_0^{2\pi} \text{Im}(v \partial_x v)(t, \theta) d\theta \right) v
\]

(3)

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\[
+ \frac{\lambda^2}{2\pi} \|v(t)\|_{L_2^\frac{1}{2}}^2 |v|^2 v \\
+ \frac{\lambda^2}{4\pi} \left( \|v(t)\|_{L_4}^4 - \frac{1}{\pi} \|v(t)\|_{L_2^\frac{1}{2}}^4 \right) v.
\]

The equations (4) and (4) are completely integrable Hamiltonian equations, having infinitely many conservation laws (2).

The well-posedness results for the equation (2), in particular considering solutions in the Sobolev space \( H^s \), were studied in many works (e.g., \([12\ 17\ 18]\)). If \( u_0 \in H^s \) with \( s \geq 1/2 \), local well-posedness was established by Herr (12). This combined with the conservation laws, can often yield the corresponding global well-posedness result for data in \( H^s \) for \( s \geq 1 \), which is small in \( L^2 \). Extension of the global well-posedness with low regularity data was obtained in the case \( s > 1/2 \) by Su Win (15), by using the I-method developed in \([3\ 4\ 5]\). A further extension was given by \([8\ 14\ 17]\). For instance, in \([8\ 14\ 17]\) Grünrock and Herr proved the local well-posedness result for initial data in Fourier-Lebesgue space \( \hat{H}_r^s \) in the range \( s \geq 1/2, \ 4/3 < r < 2 \), where the space \( \hat{H}_r^s \) is defined by the norm

\[
\|f\|_{\hat{H}_r^s} = \|\langle \xi \rangle^s \hat{f}\|_{\ell_r^t}, \quad \langle \xi \rangle = (1 + \xi^2)^{1/2}, \quad \frac{1}{r} + \frac{1}{r'} = 1.
\]

Later on in \([14]\) Nahmod, Oh, Rey-Bellet and Staffilani established the almost sure global well-posedness of \([2]\) with random initial data in \( \hat{H}_r^s \) with \( s > 1/2 \) and certain \( 4/3 < r < 2 \). Below \( s = 1/2 \), in \([17]\) Takaoka obtained the existence and time continuity of solution in \( H^s \) with some \( s < 1/2 \). We recall that the initial value problem of the equation (2) is ill-posed in \( H^s \) provided \( s < 1/2 \) (see \([1]\)). In the non-periodic case, analogous results to the well-posedness of the equation (2) were taken from the references \([4\ 5\ 10\ 11\ 13\ 15\ 16]\).

The above technology can also be applied to the equation in (1) to obtain the local and global well-posedness in \( H^s \). Indeed, the initial value problem (1) is locally well-posed in \( H^s \) for all \( s \geq 1/2 \). In order to capture the global solution in the energy class \( s \geq 1 \), we need a priori estimate. Formally, solutions of (1) satisfy the conservation laws, that is, the following three quantities are constant in time:

\[
M[u](t) = \frac{1}{2\pi} \int_{0}^{2\pi} m[u](t,x) dx = M[u](0) \quad (L^2\text{-norm}), \\
E[u](t) = \frac{1}{2\pi} \int_{0}^{2\pi} e[u](t,x) dx = E[u](0) \quad (energy), \\
P[u](t) = \frac{1}{2\pi} \int_{0}^{2\pi} p[u](t,x) dx = P[u](0) \quad (momentum),
\]

where

\[
m[u] = |u|^2 \quad (mass\ density),
\]

\[
e[u] = \frac{1}{2} |\partial_x u|^2 + \frac{\lambda}{4} |u|^2 \text{Im}(\pi \partial_x u) + \frac{\lambda^2}{12} |u|^6 \quad (energy\ density),
\]

\[
p[u] = -\frac{1}{2} \text{Im}(\pi \partial_x u) - \frac{\lambda}{4} |u|^4 \quad (momentum\ density).
\]

Note that in the case \( \mu > -5\lambda^2/16 \), the energy \( E[u](t) \) assigns positive sign. In practice if \( \mu > -5\lambda^2/16 \), the absorption of the second term \( \frac{\lambda}{4} |u|^2 \text{Im}(\pi \partial_x u) \) into the first and third terms yields

\[
e[u] \geq c \left( |\partial_x u|^2 + |u|^6 \right),
\]
for some $c > 0$, and we conclude
\[ E[u](t) \geq c\|\partial_x u(t)\|^2_{L^2}. \] (5)

We refer to the cases $\mu > -5\lambda^2/16$ and $\mu < -5\lambda^2/16$ as that the nonlinearity are respectively defocusing and focusing. By (5) and $L^2$-norm conservation (4), in the case $\mu > -5\lambda^2/16$ we have the global well-posedness for any data in $H^s$ for any $s \geq 1$, whereas in the case $\mu \leq -5\lambda^2/16$ we require small in $L^2$.

In the present paper, we concentrate on defocusing case as prototype problem. The purpose of this paper is to investigate the dynamics on resonant clusters for (1) in the underlying nonlinear energy transfers. Notice that such energy exchange processes have been seen with the defocusing cubic nonlinear Schrödinger equation on two spacial dimension by Colliander, Keel, Staffilani, Takaoka and Tao [6], and with the quintic nonlinear Schrödinger equation on one spacial dimension by Grébert and Thomann [7]. We summarize these results in Remark 4 and Remark 5 respectively.

Now it is similar in [7], we begin with a resonant set associated to (1). We are interested in whether there exist solutions that have frequencies of the energy exchange observed between the Fourier modes when only four are initially excited.

**Definition 1.1.** For $(M, N) \in \mathbb{Z}^2$ with $\lambda(M + N) > 0$, we will use the notation
\[ \alpha_1 = M, \quad \alpha_2 = -3M - N, \quad \beta_1 = N, \quad \beta_2 = -3N - M, \]
and $M_* = \max\{|M|, |N|\}$. With numbers $\alpha_j, \beta_j (j = 1, 2)$ satisfying
\[ \gamma_1 + \gamma_3 \neq \gamma_2 + \gamma_4, \] (6)
for $\{\gamma_1, \gamma_3\} \neq \{\gamma_2, \gamma_4\}$, where each $\gamma_j$ is one of $\alpha_1, \alpha_2, \beta_1, \beta_2$, a resonant set $\Lambda = \Lambda(M, N)$ takes the form
\[ \Lambda = \{\alpha_1, \alpha_2, \beta_1, \beta_2\}. \]

Let $P_L$ and $P_H$ be frequency cut-off operators at low and high frequency modes, respectively, defined as follows:
\[ P_L f(\xi) = \begin{cases} f(\xi), & \text{if } |\xi| \lesssim |M + N|, \\ 0, & \text{if } |\xi| \gg |M + N|, \end{cases} \]
and $P_H f = (I - P_L)f$.

**Remark 1.** A typical sufficient condition for (6) is that $M + N = \text{sgn}\lambda$ and $|M| \gg 1$.

The main result of this paper is the following theorem.

**Theorem 1.2** (Case $\lambda = 20(M + N)$). Let $\lambda = 20(M + N)$, $-5\lambda^2/16 < \mu \neq 0$, $M_* \gg |\lambda|$ and $|\mu| \ll M_*^2$. There exist $c_* \in (0, 1/2)$ and $T > 0$ satisfying
\[ \frac{1}{c_*|\mu|} < T < \frac{1}{c_*^2}, \]
where the constant $c > 1$ is independent of $\lambda, \mu, M$ and $N$ such that the following holds: there exist $2T$-periodic function $K(t) : \mathbb{R} \to (0, 1)$ and a solution to (1) such that $K(0) \leq 1/2 - c_*$ and $K(T) \geq 1/2 + c_*$ and for all $t \in \mathbb{R}$
\[ u(t, x) = \sum_{\xi \in \Lambda} a_{\xi}(t)e^{ix\xi} + e(t, x), \quad e|_{t=0} = 0, \] (7)
where
\[ |a_{\alpha_1}(t)|^2 = 2|a_{\alpha_2}(t)|^2 = K(t), \quad |a_{\beta_1}(t)|^2 = 2|a_{\beta_2}(t)|^2 = 1 - K(t). \]
Moreover for $\delta \in [1/2, 1]$ and $|t| \ll \min\{\frac{1}{|\lambda| M^\varphi_\varepsilon}, \frac{1}{s^2 + |\mu|}\}$, 
\[
\sum_{\xi \notin \Lambda} \left| \left( M_\delta^\varepsilon P_L + |\xi|^\delta P_H \right) F_x e(t, \xi) \right| \lesssim \frac{1}{M^\delta_\varepsilon},
\] (8)

**Remark 2.** Notice that $\frac{1}{|\mu|} \gg \min\{\frac{1}{|\lambda| M^\varphi_\varepsilon}, \frac{1}{s^2 + |\mu|}\}$, provided $|\lambda| M^\varphi_\varepsilon \gg |\mu|$ or $|\lambda| \gg 1$. So it is unclear what the time trajectory of $e(t, x)$ in [7] describes beyond the time interval $|t| \ll \min\{\frac{1}{|\lambda| M^\varphi_\varepsilon}, \frac{1}{s^2 + |\mu|}\}$ in that setting.

**Remark 3.** There seems to be no complete bound on the error term in Theorem 1.2. Indeed, the frequencies in $\Lambda$ are missing in (8), and the error term $\varepsilon$ is only controlled for frequencies $\xi$ away from the resonance set $\Lambda$. However, the standard approach to analyzing the error in an approximation of the equation (1) (equivalent one in [16] below) only for frequencies $\xi \in \Lambda$ as truncated toy model (17) below, along with (8) and conservation laws (4) provides the good approximation
\[
u(t, x) \approx \sum_{\xi \in \Lambda} a_\xi(t) e^{ix\xi},
\]
over same time interval $|t| \ll \min\{\frac{1}{|\lambda| M^\varphi_\varepsilon}, \frac{1}{s^2 + |\mu|}\}$ as in Theorem 1.2 where masses $a_\beta_j$ are excited initially compared these of $a_\alpha_j$, then move decreasingly to that of $a_\alpha_j$, in periodically. The approximation lemma can be seen by a direct computation. We leave it to the readers.

**Remark 4.** For the defocusing cubic nonlinear Schrödinger equation on two spatial dimension:
\[i\partial_t u + \Delta u = |u|^2 u,
\]
one can obtain the low-to-high energy cascade scenario, see [9]. Indeed, for any $s > 1$ and $K \gg 1 \gg \varepsilon$ there exist a global solution $u(t)$ of 2d-NLS on $\mathbb{T}^2$ and time $T$ satisfying that
\[\|u(0)\|_{H^s} \leq \varepsilon, \quad \|u(T)\|_{H^s} \geq K.
\]

**Remark 5.** In [7], one striking example of a similar situation is considered for the quintic nonlinear Schrödinger equation on one spacial dimension:
\[i\partial_t u + \partial_x^2 u + \nu |u|^4 u.
\] (9)
For $(k, n) \in \mathbb{Z}^2$ with $k \neq 0$, a resonant set $\Lambda_1$ corresponding to [9] is the following:
\[\Lambda_1 = \{\alpha_1', \alpha_2', \beta_1', \beta_2'\},
\]
where
\[(\alpha_1', \alpha_2', \beta_1', \beta_2') = (n + 3k, n, n + k, n + 4k).
\]
Comparing to the result stated in Theorem 1.2 they obtained dynamical properties for [9] in stronger sense. More precisely, it is obtained that there exist $T > 0$, $\nu_0 > 0$, $2T$-periodic function $K : \mathbb{R} \to (0, 1)$ with $K(0) \leq 1/4$ and $K(T) \geq 3/4$ so that for all $0 < \nu < \nu_0$ and $s \in \mathbb{R}$, there exists a solution of [9] satisfying
\[u(t) = \sum_{\xi \in \Lambda_1} a_\xi(t) e^{ix\xi} + \left(\nu^{1/4} + t\nu^{3/2}\right) O_{H^s}(1)
\] (10)
for $|t| \lesssim \nu^{-3/2}$, where
\[|a_{n+3k}(t)|^2 = 2|a_n(t)|^2 = K(\nu t),
\]
\[|a_{n+k}(t)|^2 = 2|a_{n+4k}(t)|^2 = 1 - K(\nu t).
\]
The proof is centered on the resonant normal form approach. The formula (7) looks very similar to that of (10), but this conveys different information because of obstructions by the derivative nonlinearity in (1).

Recently, Haus and Procesi [9] proved via KAM theory, infinite-time existence of solutions of the type obtained for long times by [7].

We now turn to equation adjusting the amplitude scaling:

\[ u(t, x) \mapsto v(t, x) = \sqrt{\frac{20(M + N)}{\lambda'}} u(t, x), \quad \lambda'(M + N) > 0, \]

which transforms the equation (1) in the case \( \lambda = 20(M + N) \) to the following:

\[ i\partial_t v + \partial_x^2 v = -\lambda' v^2 \partial_x\nu + \mu'|v|^4v, \quad (11) \]

where

\[ \mu' = \mu \left( \frac{20(M + N)}{\lambda'} \right)^2. \]

Apply this to Theorem 1.2, we obtain the immediate corollary, which is equivalent to the counterpart of Theorem 1.2.

**Corollary 1** (Case \( \lambda'(M + N) > 0 \)). Let \(-5\lambda^2/16 < \mu' \neq 0, \quad M_* \gg |M + N| \) and \( M_*|M + N| \gg |M|\sqrt{\mu'}\). There exist \( c_* \in (0, 1/2) \) and \( T > 0 \) satisfying \( T \to 0 \) as \( \frac{1}{\mu'}(\frac{M+N}{\lambda'})^2 \to 0 \) such that the following holds: there exist 2T-periodic function \( K(t) : \mathbb{R} \to (0, 1) \) and a global solution to (11) such that \( K(0) \leq 1/2 - c_* \) and \( K(T) \geq 1/2 + c_* \) and for all \( t \in \mathbb{R} \)

\[ v(t, x) = \sqrt{\frac{20(M + N)}{\lambda'}} \left( \sum_{\xi \in \Lambda} a_\xi(t)e^{ix\xi} + e(t, x) \right), \quad e|_{t=0} = 0, \]

where

\[ |a_{\alpha_1}(t)|^2 = 2|a_{\alpha_2}(t)|^2 = K(t), \quad |a_{\beta_1}(t)|^2 = 2|a_{\beta_2}(t)|^2 = 1 - K(t). \]

Moreover for \( \delta \in [1/2, 1] \) and \( |t| \ll \min\left(\frac{1}{M_*\lambda'}, \frac{(M+M)^2}{(M+N)^2+\lambda^2\mu'}\right) \),

\[ \sum_{\xi \notin \Lambda} \left| (M_\delta P_L + |\xi|^\delta P_H) \mathcal{F}_x e(t, \xi) \right| \lesssim \frac{1}{M_*}. \]

Theorem 1.2 is obtained by analyzing an infinitely dimensional system of ordinary differential equations that include the resonant components on nonlinearity of (1). We first make the Fourier transform to the equation (1) and build a resonant approximation model (toy model), which is an ordinary differential system. This is accomplished by investigating the long time dynamics of toy model. The error estimate is proved to approximate well the solution of (1) by that of the toy model by using the bootstrap argument.

**Notation.** We prefer to use the notation \( \langle \cdot \rangle = (1 + |\cdot|^2)^{1/2} \). The over dot \( \dot{a}(t) \) denotes the time derivative of \( a(t) \).

The Fourier transform with respect to the space variable (discrete Fourier transform) is defined by

\[ \mathcal{F}_x f(\xi) = \frac{1}{\sqrt{2\pi}} \int_0^{2\pi} e^{-i\xi x} f(x) \, dx, \quad \xi \in \mathbb{Z}. \]

We abbreviate \( \mathcal{F}_x f \) as \( \hat{f} \).
For $1 \leq p, q \leq \infty$, we use the mixed norm notation $\|f\|_{L^q_t L^p_x}$ with norm
\[
\|f\|_{L^q_t L^p_x} = \left( \int_{-\infty}^{\infty} \|f(t)\|_{L^q_x(T)}^p \ dt \right)^{1/q},
\]
with the obvious modification when $q = \infty$. For $T > 0$, we also use $\|f\|_{L^q_T L^p_x}$ to denote the norm
\[
\|f\|_{L^q_T L^p_x} = \left( \int_0^T \|f(t)\|_{L^q_x(T)}^p \ dt \right)^{1/q},
\]
with the obvious modification when $q = \infty$. We also use $L^p_T L^q_x$ in a similar manner.

We use $c, C$ to denote various constants that does not depend on $\lambda, \mu, M$ and $N$. We use $A \lesssim B$ to denote $A \leq CB$ for some constant $C > 0$. Similarly, we write $A \ll B$ to mean $A \leq cB$ for some small constant $c > 0$.

For complex-valued $n$ functions $f_1, f_2, \ldots, f_n$ defined on the set $\mathbb{Z}$ of integers, we write the discrete convolution (convolution sum) $[f_1 * f_2 * \ldots * f_n](\xi)$ as
\[
[f_1 * f_2 * \ldots * f_n](\xi) = \sum_{\sum \lambda = \xi} \prod_{j=1}^n f_j(\xi_j),
\]
where $\sum$ denotes a summation over the set where $\xi_1 + \xi_2 + \ldots + \xi_n = \xi$.

The rest of this paper is organized as follows. In Section 2, we reduce the equation (1) to that for the Fourier coefficients, by taking the advantage of conservation laws (4). In Section 3, we construct an appropriate ordinary differential equation (ODE) system as theoretical toy model for (1). In Section 4, we solve the ODE system. In Section 5, the equation that we derived in the previous section should approximate the dynamics of the first differential equation (1).

2. Fourier ansatz. In this section we derive the resonant approximation model related to (1). We consider the smooth solutions thought the paper, and write
\[
u(t, x) = \sum_{\xi \in \mathbb{Z}} a_\xi(t) e^{i(\alpha x + \xi t^2)}.
\]
(12)

Substituting this into the equation (1) and using
\[
\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 = 2(\xi_1 - \xi_2)(\xi_1 - \xi) = 2(\xi_3 - \xi_2)(\xi_3 - \xi),
\]
(13)
for $\xi = \xi_1 - \xi_2 + \xi_3$, we obtain the equation for $a_\xi = a_\xi(t)$:
\[
i a_\xi = \lambda |a_\xi|^2 a_\xi - \mu \sum_{\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} = \emptyset} \xi_2 a_{\xi_1} a_{\xi_3} a_{\xi_2} e^{2i t(\xi_1 - \xi_2)(\xi_1 - \xi)}
\]
\[-2\lambda \sum_{\xi \in \mathbb{Z}} \xi |a_\xi|^2 a_\xi + \mu \sum_{\xi \in \mathbb{Z}} a_{\xi_1} a_{\xi_3} a_{\xi_2} e^{i t(\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2)},
\]
(14)
by splitting the cubic convolution sum on $\xi = \xi_1 - \xi_2 + \xi_3$ into two types of sums as same as in [17]: $\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} = \emptyset$ and $\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} \neq \emptyset$. 
Denote \( P_0 = P[u](0) \) and \( M_0 = M[u](0) \). Taking the momentum conservation via Fourier transformation, we obtain

\[
\sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 = -\frac{\lambda}{2} \sum_{\xi_1 + \xi_3 = \xi_2 + \xi_4} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)} - 2P_0.
\]

On substituting this into the second term on the right-hand side of (14), we have the following form

\[
i \dot{a}_{\xi} = \lambda |a_{\xi}|^2 a_{\xi} + 4\lambda P_0 a_{\xi}
- \lambda \sum_{\xi_1, \xi_3 \ni \xi_2, \xi} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)}
+ \lambda^2 |a_{\xi}|^2 \sum_{\xi_1 + \xi_4 = \xi_2 + \xi_3} a_{\xi_1} \overline{a_{\xi_4}} a_{\xi_2} \overline{a_{\xi_3}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)}
+ \mu \sum_{\xi_1, \xi_3 \ni \xi_2, \xi} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{it(\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2)}.
\]

Note that \((\xi_1 - \xi_2)(\xi_1 - \xi_3) \neq 0\) when \(\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} = \emptyset\) in the third term on the right-hand side of (15). When \((\xi_1 - \xi_2)(\xi_1 - \xi_4) \neq 0\), it oscillates in space-time at different frequencies, which illustrates hidden smoothing effect. Thus for the third term on the right-hand side of (15), we divide the sum into two cases: \(\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} = \emptyset\) and \(\{\xi_1, \xi_3\} \cap \{\xi_2, \xi\} \neq \emptyset\), then

\[
\sum_{\xi_1 + \xi_3 = \xi_2 + \xi_4} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)}
= 2 \left( \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 \right)^2 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4
+ \sum_{\xi_1 + \xi_3 = \xi_2 + \xi_4} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)}
= 2M_0^2 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 + \sum_{\xi_1 + \xi_3 = \xi_2 + \xi_4} a_{\xi_1} \overline{a_{\xi_3}} a_{\xi_2} \overline{a_{\xi_4}} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi_3)}.
\]

We also divide the sum in the last term on the right-hand side of (15) into two cases: \(\{\xi_1, \xi_3, \xi_5\} \neq \{\xi_2, \xi_4, \xi\}\) and \(\{\xi_1, \xi_3, \xi_5\} \neq \{\xi_2, \xi_4, \xi\}\). On the set \(\{\xi_1, \xi_3, \xi_5\} = \{\xi_2, \xi_4, \xi\}\), we split this case into three disjoint subsets: (i) \(\xi_1 = \xi, \xi_3 = \xi_5 = \xi_2, \xi_4\); (ii) \(\xi_1 \neq \xi, \xi_3 = \xi_5 = \xi_2, \xi_4\); (iii) \(\xi_1 \neq \xi, \xi_3 \neq \xi, \{\xi_1, \xi_3\} = \{\xi_2, \xi_3\}\). The contribution of the subset (i) to the last term on the right-hand side of (15) is equal to

\[
a_{\xi} \left( 2 \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 \sum_{\xi'' \neq \xi'} |a_{\xi''}|^2 + \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 \right)
= a_{\xi} \left( 2 \left( \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 \right)^2 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 \right)
\]
\[ a_\xi \left( 2M_0^2 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 \right). \]

The contribution of the subset (ii) to the last term on the right-hand side of (15) is equal to
\[
a_\xi \left( \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 \right)^2 - 2|a_\xi|^2 \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^2 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 + |a_\xi|^4.
\]

The contribution of the subset (iii) to the last term on the right-hand side of (15) is equal to
\[
a_\xi \left( 2M_0^2 - 2|a_\xi|^2 M_0 - \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 + |a_\xi|^4 \right).
\]

Collecting the above results, we have that
\[
\sum_{\ast} a_{\xi_1} \bar{a}_{\xi_2} a_{\xi_3} \bar{a}_{\xi_4} a_{\xi_5} e^{it(\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2)}
\]
\[
= \left( 6M_0^2 - 6M_0 |a_\xi|^2 - 3 \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 + 4|a_\xi|^4 \right) a_\xi
\]
\[
+ \sum_{\{\xi_1, \xi_2, \xi_3\} \neq \{\xi_2, \xi_3, \xi_1\}} a_{\xi_1} \bar{a}_{\xi_2} a_{\xi_3} \bar{a}_{\xi_4} a_{\xi_5} e^{it(\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2)}.
\]

From this, the formula (15) is described as
\[
M_0^2 (\lambda^2 - 6M_0 \mu) |a_\xi|^2 + 2 \left( M_0^2 (\lambda^2 + 3\mu) + 2\lambda P_0 \right) a_\xi
\]
\[
+ 4\mu|a_\xi|^4 a_\xi - (\lambda^2 + 3\mu) \sum_{\xi' \in \mathbb{Z}} |a_{\xi'}|^4 \right) a_\xi
\]
\[
- \lambda \sum_{\{\xi_1, \xi_2, \xi_3\} \neq \{\xi_2, \xi_3, \xi_1\}} \xi_2 a_{\xi_1} \bar{a}_{\xi_2} a_{\xi_3} e^{2it(\xi_1 - \xi_2)(\xi_1 - \xi)}.
\]

(16)
\[
+ \lambda^2 a_\xi \sum_{\xi_1 + \xi_3 = \xi_2 + \xi_4 \in \Lambda} a_{\xi_1} a_{\xi_2} a_{\xi_3} a_{\xi_4} e^{2i(\xi_1 - \xi_2)(\xi_1 - \xi_3)} + \mu \sum_{\{\xi_1, \xi_2, \xi_3\} \neq \{\xi_2, \xi_4, \xi_5\}} a_{\xi_1} a_{\xi_2} a_{\xi_3} a_{\xi_4} e^{i(\xi_1 - \xi_2 + \xi_3 - \xi_4 + \xi_5 - \xi_6)}.
\]

The cubic interactions through resonance in the first term on the right-hand side of (16) are self-interaction, which enhance the phase oscillations. In next section, we seek the contribution of quintic resonant interactions to the last term on the right-hand side of (16).

3. Toy model. In this section we shall construct resonant, finite dimensional truncation model to (16). Let

\[
\lambda = 20(\alpha_1 + \beta_1) = 20(M + N).
\]

We start out by referring to the observation on the resonant set used in [7], which can be helpful in predicting the resonant set related to the equation (16).

**Lemma 3.1** (Grébert-Thomann [7] Lemmas 2.1 and 2.2). Assume that \( \xi_j \in \mathbb{Z} \), (1 \( \leq \) \( j \leq \) 6) satisfy

\[
\xi_1 + \xi_3 + \xi_5 = \xi_2 + \xi_4 + \xi_6,
\]

\[
\xi_1^2 + \xi_3^2 + \xi_5^2 = \xi_2^2 + \xi_4^2 + \xi_6^2,
\]

and \( \{\xi_1, \xi_3, \xi_5\} \neq \{\xi_2, \xi_4, \xi_6\} \). Then \( \{\xi_1, \xi_3, \xi_5\} \cap \{\xi_2, \xi_4, \xi_6\} = \emptyset \). Moreover the cardinal of \( \{\xi_j \mid 1 \leq j \leq 6\} \) is greater than or equal to 4.

Considering the influence by the first term on the right-hand side of (16), we define the resonant interaction set

\[
\Lambda^* = \{\{\alpha_1, \alpha_2, \alpha_3\}, \{\beta_1, \beta_2, \beta_3\}\},
\]

where \( \alpha_j \) and \( \beta_j \) are defined in Definition [11]. We will use the notation

\[
\{\{\xi_1, \xi_3, \xi_5\}, \{\xi_2, \xi_4, \xi_6\}\} = \Lambda^*
\]

to represent that two of \( \xi_1, \xi_3, \xi_5 \) are \( \alpha_1 \), one of \( \xi_1, \xi_3, \xi_5 \) are \( \alpha_2 \), two of \( \xi_2, \xi_4, \xi_6 \) are \( \beta_1 \), one of \( \xi_2, \xi_4, \xi_6 \) is \( \beta_2 \), and its symmetrical objects.

Having discussed in Section 2, we define the resonant truncation ODE model of (16) by

\[
i\dot{c}_\xi = \left(\lambda \xi - 6M_0 \mu\right) |c_\xi|^2 c_\xi + 2 \left(M_0^2 (\lambda^2 + 3 \mu) + 2 \lambda P_0\right) c_\xi + 4\mu |c_\xi|^4 c_\xi - (\lambda^2 + 3 \mu) \sum_{\xi \in \Lambda} |c_\xi|^4 c_\xi + \mu \sum_{\{\xi_1, \xi_3, \xi_5\} \neq \{\xi_2, \xi_4, \xi_6\}} c_{\xi_1} c_{\xi_2} c_{\xi_3} c_{\xi_4} c_{\xi_5} c_{\xi_6} e^{i\left(\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2 - \xi_6^2\right)},
\]

with initial data

\[
c_\xi(0) = a_\xi(0),
\]

where \( \xi \in \Lambda \).

One of the most important features of the above ODE system is the following lemma.
Lemma 3.2. Let \( \{c_\xi\}_{\xi \in \Lambda} \) be a solution to (17)-(18) on the time interval \( I \). Then we have the relations:

\[
\frac{d}{dt} (|c_{\alpha_1}(t)|^2 + |c_{\beta_1}(t)|^2) = 0, \tag{19}
\]

\[
\frac{d}{dt} (|c_{\alpha_2}(t)|^2 + |c_{\beta_2}(t)|^2) = 0, \tag{20}
\]

\[
\frac{d}{dt} (|c_{\alpha_1}(t)|^2 - 2|c_{\alpha_2}(t)|^2) = 0, \tag{21}
\]

\[
\frac{d}{dt} (|c_{\beta_1}(t)|^2 - 2|c_{\beta_2}(t)|^2) = 0, \tag{22}
\]

for \( t \in I \).

Proof. The proof is straightforward calculation. We only give the proof of (19) and (21), as the proofs for (20) and (22) are analogous.

For (19), we multiply equation (17) for \( \xi = \alpha_1 \) by \( 2c_{\alpha_1} \) and take the imaginary part to get

\[
\frac{d}{dt} |c_{\alpha_1}|^2 = 2\mu \Im \sum_{\{\xi_1, \xi_2, \xi_3\} = \{\beta_1, \beta_2, \beta_3\}, \{\xi_4, \xi_5\} = \{\alpha_1, \alpha_2\}} c_{\xi_1} \xi_{\xi_2}^{c_{\xi_3}} c_{\xi_4} \xi_{\xi_5} c_{\alpha_1} e^{it(\xi_1^2 + \xi_4^2 - \xi_2^2 - \xi_5^2)}
\]

\[
= 2^3 \mu \Im c_{\beta_1} c_{\alpha_2} e^{it(2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2)},
\]

where we use that

\[
\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2 - \alpha_1^2 = 2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2,
\]

for \( \{\xi_1, \xi_3, \xi_5\} = \{\beta_1, \beta_1, \beta_2\} \) and \( \{\xi_2, \xi_4\} = \{\alpha_1, \alpha_2\} \). A similar argument as above but using equation (17) for \( \xi = \beta_1 \) yields

\[
\frac{d}{dt} |c_{\beta_1}|^2 = 2\mu \Im \sum_{\{\xi_1, \xi_3, \xi_5\} = \{\alpha_1, \alpha_2, \alpha_2\}, \{\xi_2, \xi_4\} = \{\beta_1, \beta_2\}} c_{\xi_1} \xi_{\xi_2}^{c_{\xi_3}} c_{\xi_4} c_{\alpha_2} c_{\beta_1} e^{it(\xi_1^2 + \xi_4^2 - \xi_2^2 - \xi_3^2 - \beta_1^2)}
\]

\[
= 2^3 \mu \Im c_{\beta_1} c_{\alpha_1} c_{\alpha_2} e^{it(2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2)}.
\]

Collecting the information in (23) and (24) we can obtain

\[
\frac{d}{dt} (|c_{\alpha_1}|^2 + |c_{\beta_1}|^2) = 2^3 \mu \Im \left( c_{\beta_1} c_{\alpha_2} c_{\beta_2} c_{\alpha_1} e^{it(2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2)}
\]

\[
+ c_{\alpha_1} c_{\beta_1} c_{\alpha_2} c_{\beta_2} e^{-it(2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2)} \right) = 0.
\]

On the other hand, to prove (21), we will see that

\[
\frac{d}{dt} |c_{\alpha_2}|^2 = 2\mu \Im \sum_{\{\xi_1, \xi_3, \xi_5\} = \{\beta_1, \beta_1, \beta_2\}, \{\xi_2, \xi_4\} = \{\alpha_1, \alpha\}} c_{\xi_1} \xi_{\xi_2}^{c_{\xi_3}} c_{\xi_4} c_{\alpha_2} c_{\alpha_1} e^{it(\xi_1^2 + \xi_4^2 - \xi_2^2 - \xi_3^2 - \alpha_1^2)}
\]

\[
= 2^3 \mu \Im c_{\beta_1} c_{\alpha_1} c_{\alpha_2} c_{\beta_1} c_{\alpha_2} e^{it(2\beta_1^2 + \beta_2^2 - 2\alpha_1^2 - \alpha_2^2)}.
\]

Hence the argument used in above gives

\[
\frac{d}{dt} (|c_{\alpha_1}|^2 - 2|c_{\alpha_2}|^2) = 0,
\]

which concludes the proof of (21). \(\square\)
Remark 6. The conservation laws in [19–21] are used to the study of properties of the system of [17] including the global existence results. In particular, the ODE system [17]–[18] has global solutions. We will show that the system has periodic orbits in time.

Let us set data \( \{a_\xi(0)\} \) in [18] where
\[
|a_{\alpha_1}(0)|^2 + |a_{\beta_1}(0)|^2 = 2|a_{\alpha_2}(0)|^2 + 2|a_{\beta_2}(0)|^2 = 1,
|a_{\alpha_1}(0)|^2 - 2|a_{\alpha_2}(0)|^2 = |a_{\beta_1}(0)|^2 - 2|a_{\beta_2}(0)|^2 = 0.
\]

By Lemma 3.2 and Remark 6, we have that for all \( t \in \mathbb{R} \),
\[
|c_{\alpha_1}(t)|^2 + |c_{\beta_1}(t)|^2 = 2|c_{\alpha_2}(t)|^2 + 2|c_{\beta_2}(t)|^2 = 1, \quad (25)
|c_{\alpha_1}(t)|^2 - 2|c_{\alpha_2}(t)|^2 = |c_{\beta_1}(t)|^2 - 2|c_{\beta_2}(t)|^2 = 0. \quad (26)
\]

We now write
\[
|c_{\alpha_1}(t)|^2 = K(t), \quad |c_{\alpha_2}|^2 = \frac{K(t)}{2},
|c_{\beta_1}(t)|^2 = 1 - K(t), \quad |c_{\beta_2}|^2 = \frac{1 - K(t)}{2}.
\]

Remark 7. From the relation in (25) and (26), we easily see that
\[
M_0 = \frac{3}{2} \sum_{\xi' \in \Lambda} |c_{\xi'}(t)|^4 = \frac{5}{4} - \frac{5}{2}K(t)(1 - K(t)). \quad (27)
\]

The equation [17] has a gauge freedom. On writing
\[
d_\xi(t) = c_\xi(t) \exp (iG_\xi(t)),
\]
where
\[
\dot{G}_\xi(t) = \lambda \xi |c_\xi(t)|^2 + 2 \left( M_0^2 (\lambda^2 + 3\mu) + 2\lambda P_0 \right) - (\lambda^2 + 3\mu) \sum_{\xi' \in \Lambda} |c_{\xi'}(t)|^4,
\]
\[
G_\xi(0) = 0,
\]
and using (27), we have that the equation [17] becomes the following equation for \( d_\xi \),
\[
\dot{d}_\xi = \mu |d_\xi|^2 (4|d_\xi|^2 - 9)d_\xi + \mu \sum_{\{(\xi_1, \xi_2, \xi_3), (\xi_4, \xi, \xi)\} = \Lambda^*} d_{\xi_1}d_{\xi_2}d_{\xi_3}d_{\xi_4}d_{\xi} e^{i\mathcal{I}}, \quad (28)
\]
where \( \mathcal{I}(0) = 0 \) and
\[
\mathcal{I}(t) = \xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2 - \xi^2
- \lambda (|d_{\xi_1}(t)|^2 - |d_{\xi_2}|^2|d_{\xi_3}(t)|^2 + \xi_3|d_{\xi_4}(t)|^2 - |d_{\xi_1}(t)|^2 - \xi_4|d_{\xi_4}(t)|^2 - \xi_5|d_{\xi_5}(t)|^2 - \xi|d_{\xi}(t)|^2).
\]

Let us discuss the contribution of the right-hand of (28).

Lemma 3.3. Let
\[
|d_{\alpha_1}(t)|^2 = K(t), \quad |d_{\alpha_2}(t)|^2 = \frac{K(t)}{2},
|d_{\beta_1}(t)|^2 = 1 - K(t), \quad |d_{\beta_2}(t)|^2 = \frac{1 - K(t)}{2},
\]
and \( \{\xi_1, \xi_3, \xi_5\}, \{\xi_2, \xi_4, \xi\} = \Lambda^* \). Then we have \( \dot{\mathcal{I}}(t) = 0. \)
Substituting these equations into (32) gives hypotheses in Corollary 1). This simply reflects the constraint on system of (32).

Proof. By symmetry, it suffices to consider the case \((\xi_1, \xi_3, \xi_5, \xi_2, \xi_4, \xi) = (\alpha_1, \alpha_1, \alpha_2, \beta_1, \beta_1, \beta_2)\).

First one notices that
\[
\frac{\xi_1[\xi_1^2 - \xi_2\xi_2^2 + \xi_3\xi_3^2 - \xi_4^2\xi_4^2 + \xi_5^2\xi_5^2 - \xi|\xi|^2]}{2K} + 2\alpha_1 + \alpha_2 + \frac{\beta_2}{2} - 2\beta_1 - \frac{\beta_2}{2}
\]
(29)

From the crucial fact that \((\alpha_1, \alpha_2, \beta_1, \beta_2) = (M, -3M - N, N, -3N - M)\), we see that the first term on the right-hand side of (29) is zero. From the identity
\[
\alpha_1 + \beta_1 = M + N, \ 2\alpha_1 + 2\beta_1 + \frac{\alpha_2 + \beta_2}{2} = 0,
\]
and
\[
2\alpha_1 - 2\beta_1 + \alpha_2 - \beta_2 = 0,
\]
thus \(\tilde{L}\) simplifies to
\[
2\alpha_1^2 - 2\beta_1^2 + \alpha_2^2 - \beta_2^2 + \lambda \left(2\beta_1 + \frac{\beta_2}{2}\right)
\]
\[
= 10(\alpha_1 + \beta_1)(\alpha_1 - \beta_1) + \frac{\lambda}{4}(\alpha_2 - \beta_2)
\]
\[
= 10(\alpha_1 - \beta_1) \left(\alpha_1 + \beta_1 - \frac{\lambda}{20}\right) = 0,
\]
since \(\lambda = 20(\alpha_1 + \beta_1)\). This completes the proof of Lemma 3.3.

Finally we reduce the form (28) to the following form:
\[
\frac{d\xi}{d\xi} = \mu[|d\xi|^2(4|d\xi|^2 - 9)]d\xi
\]
\[
+ \mu \sum_{\{(\xi_1, \xi_3, \xi_5), (\xi_2, \xi_4, \xi)\} = \Lambda^*} d\xi_i |d\xi_i|^2 d\xi_i, \quad (32)
\]
with initial data \(d\xi(0) = \alpha_d(0)\).

4. ODE dynamics. We now use the symplectic polar coordinates
\[
d\xi(t) = \sqrt{I_\xi(t)} \exp(i\theta_\xi(t)), \quad \xi \in \Lambda,
\]
where \(\theta_\xi(t) \in \mathbb{R}\). From (25) and (26), we may assume that \(0 < I_\xi(t) < 1\) (see the hypotheses in Corollary 1). This simply reflects the constraint on system of (32). Substituting these equations into (32) gives
\[
-\dot{\theta}_\xi = \mu I_\xi(4I_\xi - 9) + \mu \sum_{\{(\xi_1, \xi_3, \xi_5), (\xi_2, \xi_4, \xi)\} = \Lambda^*} \left(\prod_{j=1}^{3} I_{\xi_j}\right)^{1/2} \cos(\phi_\xi), \quad (33)
\]
\[
\frac{1}{2} \dot{I}_\xi = \sum_{\{(\xi_1, \xi_3, \xi_5), (\xi_2, \xi_4, \xi)\} = \Lambda^*} \left(\prod_{j=1}^{3} I_{\xi_j}\right)^{1/2} (I_\xi)^{1/2} \sin(\phi_\xi), \quad (34)
\]
where
\[
\phi_\xi = \begin{cases}
2\theta_\alpha_1 + \theta_\alpha_2 - 2\theta_\beta_1 - \theta_\phi_2, & \text{when } \xi \in \{\beta_1, \beta_2\}, \\
2\theta_\beta_1 + \theta_\beta_2 - 2\theta_\alpha_1 - \theta_\alpha_2, & \text{when } \xi \in \{\alpha_1, \alpha_2\}.
\end{cases}
\]
Remark 8. It is possible to consider the following Hamiltonian function associated to (33)-(34):

\[
H = -\mu \sum_{\xi \in \Lambda} \left( \frac{4}{3} I_3^3 - \frac{9}{2} I_2^2 \right) - 6\mu I_{\alpha_1} I_{\alpha_2} I_{\beta_1} I_{\beta_2} \cos(2\theta_{\alpha_1} + \theta_{\alpha_2} - 2\theta_{\beta_1} - \theta_{\beta_2}) \\
= \frac{33}{8} \mu (K^2 + (1 - K)^2) + \frac{3}{2} \mu K(1 - K) \\
-3\mu K^2 (1 - K)^2 \cos(2\theta_{\alpha_1} + \theta_{\alpha_2} - 2\theta_{\beta_1} - \theta_{\beta_2})
\]

(since by Lemma 3.3). In particular, we have that the system (33)-(34) takes the symplectic form

\[
\begin{align*}
\dot{\theta}_\xi &= -\frac{\partial H}{\partial I_\xi}, \\
\dot{I}_\xi &= \frac{\partial H}{\partial \theta_\xi},
\end{align*}
\]

for \(\xi \in \{\alpha_1, \alpha_2, \beta_1, \beta_2\}\).

By introducing the symplectic change of variables in a similar way to [7], we define

\[
\theta = \tau(\theta_{\alpha_1}, \theta_{\alpha_2}, \theta_{\beta_1}, \theta_{\beta_2}), \quad I = \tau(I_{\alpha_1}, I_{\alpha_2}, I_{\beta_1}, I_{\beta_2}),
\]

and define new variables

\[
\phi = \tau(\phi_1, \phi_2, \phi_3, \phi_4), \quad J = \tau(J_1, J_2, J_3, J_4),
\]

where

\[
\begin{align*}
\phi_1 &= 2\theta_{\alpha_1} + \theta_{\alpha_2} - 2\theta_{\beta_1} - \theta_{\beta_2}, \\
\phi_2 &= \theta_{\alpha_2}, \quad \phi_3 = \theta_{\beta_1}, \quad \phi_4 = \theta_{\beta_1}, \\
J_1 &= \frac{1}{2} I_{\alpha_1}, \quad J_2 = -\frac{1}{2} I_{\alpha_1} + I_{\alpha_2}, \quad J_3 = I_{\alpha_1} + I_{\beta_1}, \quad J_4 = \frac{1}{2} I_{\alpha_1} + I_{\beta_2}.
\end{align*}
\]

Then the Hamiltonian flow in the action-angle coordinates (33)-(34) satisfies

\[
\begin{pmatrix}
\phi \\
J
\end{pmatrix} = \begin{pmatrix}
A & 0 \\
0 & tA^{-1}
\end{pmatrix}
\begin{pmatrix}
\theta \\
I
\end{pmatrix},
\]

where

\[
A = \begin{pmatrix}
2 & 1 & -2 & -1 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{pmatrix}, \quad tA^{-1} = \begin{pmatrix}
\frac{1}{2} & 0 & 0 & 0 \\
-\frac{1}{2} & 0 & 0 & 0 \\
1 & 0 & 1 & 0 \\
1 & 0 & 0 & 1
\end{pmatrix}.
\]

By (25)-(26) and Remark 8 we have

\[
J_2 = 0, \quad J_3 = 1, \quad J_4 = \frac{1}{2}.
\]

We will interest in the periodic solutions \((\theta_j, I_j)\). We achieve this by finding a periodic solution \(K(t)\) for the ODE system on \((\phi_1, J_1)\).

For the purpose of obtaining the solutions \((\phi_1, J_1)\), one considers the equations for \(\phi_1\)

\[
\dot{\phi}_1 = -2 \frac{\partial H}{\partial I_{\alpha_1}} - \frac{\partial H}{\partial I_{\alpha_2}} + 2 \frac{\partial H}{\partial I_{\beta_1}} + \frac{\partial H}{\partial I_{\beta_2}} \\
= 4\mu \left( 2K^2 + \left( \frac{K}{2} \right)^2 - 2(1 - K)^2 - \left( \frac{1 - K}{2} \right)^2 \right)
\]
\[-9\mu \left(2K + \frac{K}{2} - 2(1 - K) - \frac{1 - K}{2}\right) + \mu \left(12 \left(\frac{K}{2}\right)^{\frac{1}{2}} (1 - K) \left(\frac{1 - K}{2}\right)^{\frac{1}{2}} + 3K(1 - K) \left(\frac{1 - K}{2}\right)^{\frac{1}{2}} \right) \\right. \\
-12K \left(\frac{K}{2}\right)^{\frac{1}{2}} \left(1 - K\right)^{\frac{1}{2}} \left(\frac{K}{2}\right)^{\frac{1}{2}} = 9\mu(1 - 2K) \left(\frac{7}{2} + K^{\frac{1}{2}} (1 - K)^{\frac{1}{2}} \cos \phi_1\right),
\]

and for \(J_1\)
\[
\dot{J}_1 = \frac{\partial H}{\partial \theta_{\alpha_1}} = 12\mu I_{\alpha_1} I_{\beta_1} I_{\beta_2} \sin \phi_1 = 6\mu K^{\frac{1}{2}} (1 - K)^{\frac{1}{2}} \sin \phi_1.
\]

Since from \(J_1 = I_{\alpha_1}/2 = K/2\), we conclude that \(\phi_1, K\) satisfy the following system:
\[
\begin{cases}
\dot{\phi}_1 = 9\mu(1 - 2K) \left(\frac{7}{2} + K^{\frac{1}{2}} (1 - K)^{\frac{1}{2}} \cos \phi_1\right), \\
\dot{K} = 12\mu K^{\frac{1}{2}} (1 - K)^{\frac{1}{2}} \sin \phi_1.
\end{cases}
\] (36)

We will apply the argument in [7, Proposition 4.2] for the dynamical system (36), and have the following proposition.

**Proposition 1.** There exists \(T > 0\) so that \((\phi_1, K)\) is a \(2T\)-periodic solution of (36) and
\[
K(0) + K(T) = 1.
\]

Moreover, \(\mu T\) is a constant independent of \(\mu\).

**Proof.** The existence of a periodic solution for (36) is similar to the one exposed in [7, Proposition 4.2]. The intuition is that in the region \((\phi_1, K) \in [-\pi, \pi] \times (0, 1)\), the equilibrium points associated with (36) are \(\omega_0 = (0, 1/2)\) and \(\omega_\pm = (\pm \pi, 1/2)\). At the point \(\omega_\pm\), we have that the flow obeys the Hamiltonian (35)
\[
H = \frac{45}{16} \mu,
\]
which constraints the equation (35) to be
\[
K(1 - K) \left(\frac{27}{4} + 3K^{\frac{1}{2}} (1 - K)^{\frac{1}{2}} \cos \phi_1\right) = \frac{21}{16}.
\] (37)

The equation (37) defines two heteroclinic orbits. It can be concluded that the equation (37) has periodic solutions \((\phi_1, K)\).

We use the rescaling theorem done by \(t \mapsto \mu t\) to concrete \(\mu T\) as a constant independent of \(\mu\). 

Once one has the periodic solution \(K(t)\) in Proposition 1, one can obtain the following proposition.
Proposition 2. There exist $c_\ast > 0$ and $c > 0$ that are independent of $\lambda$, $\mu$, $M$, $N$ so that the following holds: there exist $2T$-periodic function $K(t) : \mathbb{R} \to (0,1)$ and solutions $\{c_\xi \}_{\xi \in \Lambda}$ to the ODE system (17)-(18) such that $\mu T = c$. Moreover, $K(0) \leq 1/2 - c_\ast$ and $K(T) \geq 1/2 + c_\ast$ and for all $t \in \mathbb{R}$

$$|c_{\alpha_1}(t)|^2 = 2|c_{\alpha_2}(t)|^2 = K(t), \quad |c_{\beta_1}(t)|^2 = 2|c_{\beta_2}(t)|^2 = 1 - K(t).$$

5. Error estimates. In this section, we shall evaluate the residual terms $a_\xi(t)$, $\xi \notin \Lambda$ for (10).

Recall $M_\ast = \max\{|M|, |N|\}$. First we show the conservation laws.

Lemma 5.1. Consider a global solution $u(t)$ to (1) as formed in (13) satisfying $a_\xi(0) = 0$ for all $\xi \notin \Lambda$, and

$$|a_{\alpha_1}(0)|^2 = 2|a_{\alpha_2}(0)|^2 = K(0), \quad |a_{\beta_1}(0)|^2 = 2|a_{\beta_2}(0)|^2 = 1 - K(0).$$

Then we have that

$$\|\langle \xi \rangle a_\xi(t)\|_{L^2} \lesssim M_\ast. \quad (38)$$

Proof. As we shall see in (4) and $M_0 = 3/2$ by $L^2$-norm conservation, one can give

$$\|u(t)\|_{L^2}^2 = \sum_{\xi \in \Lambda} \left( |a_{\alpha_1}(0)|^2 + |a_{\alpha_2}(0)|^2 + |a_{\beta_1}(0)|^2 + |a_{\beta_2}(0)|^2 \right) = \frac{3}{2}. \quad (39)$$

Also by (5) and via the Sobolev inequality, we have

$$\|\partial_x u(t)\|_{L^2}^2 \lesssim E[u](0) \lesssim \|\partial_x u(0)\|_{L^2}(1 + \|u(0)\|_{L^2}^4) \lesssim M_\ast^2. \quad (40)$$

The estimate (38) follows from (39) and (40).

On the other hand, as in Section 3 setting

$$b_\xi(t) = a_\xi(t)e^{i\Phi(t, \xi)}, \quad (41)$$

where

$$\Phi(t, \xi) = \int_0^t \left( (\lambda \xi - 6M_0\mu)|a_\xi(s)|^2 + 4\mu|a_\xi(s)|^4 - (\lambda^2 + 3\mu) \sum_{\xi \in \mathbb{Z}} |a_{\xi'}(s)|^4 \right) ds$$

$$+ 2t \left( M_\ast^2 (\lambda^2 + 3\mu) + 2\lambda P_0 \right),$$

leads to the equation

$$ib_\xi = -\lambda \sum_{\{\xi_1, \xi_2, \xi_3, \xi_4\} \cap \{\xi_2, \xi_3\} = \emptyset} \xi_2 b_{\xi_1} \xi_3 b_{\xi_4} e^{i\Phi_1(t, \xi_1, \xi_2, \xi_3, \xi)}$$

$$+ \lambda^2 b_\xi \sum_{\{\xi_1, \xi_3 + \xi_4\} \cap \{\xi_2, \xi_4\} = \emptyset} b_{\xi_1} b_{\xi_3} \xi_4 b_{\xi_4} e^{i\Phi_1(t, \xi_1, \xi_2, \xi_3, \xi_4)}$$

$$+ \mu \sum_{\{\xi_1, \xi_3, \xi_4\} \neq \{\xi_2, \xi_4, \xi\}} b_{\xi_1} b_{\xi_3} b_{\xi_4} e^{i\Phi_2(t, \xi_1, \xi_2, \xi_3, \xi, \xi)}, \quad (42)$$

where

$$\Phi_1 = \Phi_1(t, \xi_1, \xi_2, \xi_3, \xi_4) = \sum_{j=1}^4 (-1)^{j-1} \left( \xi_j^2 t - \Phi(t, \xi_j) \right),$$

$$\Phi_2 = \Phi_2(t, \xi_1, \xi_2, \xi_3, \xi_4, \xi_5) = \sum_{j=1}^5 (-1)^{j-1} \left( \xi_j^2 t - \Phi(t, \xi_j) \right).$$
and

\[ \Phi_2 = \Phi_3(t, \xi_1, \xi_2, \xi_3, \xi_4, \xi_5, \xi_6) = \sum_{j=1}^{6} (-1)^{j-1} (\xi_j^2 t - \Phi(t, \xi_j)). \]

Notice that by (13), we have for \( \xi_1 + \xi_3 = \xi_2 + \xi_4 \),

\[ \sum_{j=1}^{4} \xi_j^2 = 2(\xi_1 - \xi_2)(\xi_1 - \xi_4). \]

We use (39) and \( \ell_2^{\xi} \hookrightarrow \ell_1^{\xi} \) to see that

\[ \|a_\xi(t)\|_{L^\infty_{\xi_1}} \lesssim 1, \]

as well as

\[ \|b_\xi(t)\|_{L^\infty_{\xi_1}} \lesssim 1. \] (43)

To consider the error estimates, we need to introduce the following definition.

**Definition 5.2.** For \( \xi \in \mathbb{Z} \), denote \( \Lambda^{**}_\xi = \Lambda^{**}(\xi) \) is the set of \( \mathbb{Z}^3 \) such that \((\xi_1, \xi_2, \xi_3) \not\in \Lambda^{**} \) if and only if \( \xi_1 + \xi_3 = \xi_2 + \xi_4 \) and either the following three cases holds:

(i) \( \xi_j \in \Lambda \) for all \( 1 \leq j \leq 3 \),
(ii) at least, two of \( \xi_j \) are not elements of \( \Lambda \),
(iii) two of \( \xi_j \) are elements of \( \Lambda \), one of \( \xi_j \) is not elements of \( \Lambda \), \( |\xi - \xi_1| \geq M_*/100 \) and \( |\xi - \xi_4| \geq M_*/100 \).

We introduce the notation:

\[ A_L(t) = \|P_L b_\xi(t)\|_{\ell_1^{\xi_1} \xi \not\in \Lambda}, \]

\[ A_H(t) = \|\langle \xi \rangle^4 P_H b_\xi(t)\|_{\ell_1^{\xi_1} \xi \not\in \Lambda}, \]

\[ B_L(t) = \left\| P_L \left( b_\xi(t) - \tilde{b}_\xi(t) \right) \right\|_{\ell_1^{\xi_1} \xi \not\in \Lambda}, \]

and

\[ B_H(t) = \left\| \langle \xi \rangle^4 P_H \left( b_\xi(t) - \tilde{b}_\xi(t) \right) \right\|_{\ell_1^{\xi_1} \xi \not\in \Lambda}, \]

where

\[ \tilde{b}_\xi(t) = i\lambda \int_0^t \sum_{\Lambda^{**}} \xi_2 b_{\xi_1}(s) b_{\xi_2}(s) b_{\xi_3}(s) e^{2is(\xi_1 - \xi_2)(\xi_1 - \xi)} ds. \]

Now we illustrate how the error term in \( (42) \) would stay bounded for a reasonable amount of time. More precisely we have the following lemma.

**Proposition 3.** Let \( 0 < \varepsilon \leq \frac{1}{M_*}, \delta \in [1/2, 1) \) and \( b_\xi(t) \) be solution to \( (42) \). Assume that \( M_* \gg |M + N|, |\mu| \ll M_*^2 \) and for all \( |t| \leq T \)

\[ A_L(t) \leq \varepsilon^2, \]

\[ A_H(t) \leq \varepsilon. \]

Then there exists a constant \( c > 0 \) such that for all \( |t| \leq T \),

\[ B_L(t) \leq \frac{1}{10} \varepsilon^2 + c \left( \lambda^2 + |\mu| \right) \varepsilon^2 t, \]
Lemma 5.3. Assume the same hypotheses of Proposition 3. Then

\[ b_\xi(t) = \sum_{j=1}^{3} N_j(t, \xi), \]

where

\[ N_1(t, \xi) = i\lambda \int_0^t \sum_{\{\xi, \xi_3, \xi_4\} \neq \emptyset} \xi_2 b_{\xi_1}(s) \overline{b_{\xi_3}(s)} b_{\xi_4}(s) e^{i\Phi_1(s, \xi_1, \xi_2, \xi_3, \xi_4)} ds, \]

\[ N_2(t, \xi) = -i\lambda^2 \int_0^t b_\xi(s) \sum_{\{\xi, \xi_3, \xi_4\} \neq \emptyset} \xi_2 b_{\xi_1}(s) \overline{b_{\xi_3}(s)} b_{\xi_4}(s) e^{i\Phi_1(s, \xi_1, \xi_2, \xi_3, \xi_4)} ds, \]

and

\[ N_3(t, \xi) = -i\mu \int_0^t \sum_{\{\xi, \xi_3, \xi_4\} \neq \emptyset} \xi_2 b_{\xi_1}(s) \overline{b_{\xi_3}(s)} b_{\xi_4}(s) b_{\xi_5}(s) e^{i\Phi_2(s, \xi_1, \xi_2, \xi_3, \xi_4, \xi_5)} ds. \]

Moreover, if one divides the sum of \( N_2 \) into two parts, \((\xi_1, \xi_2, \xi_3) \in \Lambda^{**}\) and otherwise, then

\[ b_\xi(t) = N_1(t, \xi) + N_2(t, \xi) + N_3(t, \xi), \]

where

\[ N_1(t, \xi) = i\lambda \int_0^t \sum_{\{\xi, \xi_3, \xi_4\} \neq \emptyset} \xi_2 b_{\xi_1}(s) \overline{b_{\xi_3}(s)} b_{\xi_4}(s) e^{i\Phi_1(s, \xi_1, \xi_2, \xi_3, \xi_4)} ds. \]

For this purpose, we have the several lemmas. Let us first consider the term \( N_2 \).

Lemma 5.3. Assume the same hypotheses of Proposition 3. Then

\[ \|\langle \xi \rangle^\delta N_2(t, \xi)\|_{\xi \in \Lambda} \lesssim \lambda^2 \varepsilon^2 t. \]

Proof of Lemma 5.3. One employs the constraint in (43). Because \( \xi_1 + \xi_3 = \xi_2 + \xi_4 \) and \( \{\xi_1, \xi_3\} \cap \{\xi_2, \xi_4\} = \emptyset \), there exists some \( 1 \leq j \leq 4 \) such that \( \xi_j \notin \Lambda \). Hence by (43) we have that

\[ \|\langle \xi \rangle^\delta N_2(t, \xi)\|_{\xi \in \Lambda} \lesssim \lambda^2 t \|\langle \xi \rangle^\delta b_\xi\|_{L^2_{\xi} L^2_{\xi} \xi \in \Lambda} \lesssim \lambda^2 \varepsilon^2 t, \]

since

\[ \|\langle \xi \rangle^\delta b_\xi\|_{L^2_{\xi} \xi \notin \Lambda} \lesssim \|P_L b_\xi\|_{L^2_{\xi} L^2_{\xi} \xi \in \Lambda} + \|\langle \xi \rangle^\delta P_H b_\xi\|_{L^2_{\xi} \xi \in \Lambda} \lesssim \varepsilon^2 + \varepsilon \lesssim 2 \varepsilon. \]

\[ \square \]

Remark 9. By Lemma 5.3 and its proof, we obtained the following estimates:

\[ \|P_L N_2(t, \xi)\|_{\xi \in \Lambda} + \|\langle \xi \rangle^\delta P_H N_2(t, \xi)\|_{\xi \in \Lambda} \lesssim \lambda^2 \varepsilon^2 t. \]
Next we illustrate the estimate for $N_3$. Before the estimate on $N_3$, we show the following lemma.

**Lemma 5.4.** Assume that $M_* \gg |M + N|$, $\xi_1 + \xi_3 + \xi_5 = \xi_2 + \xi_4 + \xi_6$ and 
$|\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5^2 - \xi_6^2| \ll M_*^2$, $\xi_j \in \Lambda$ for all $1 \leq j \leq 5$, then 
$\{\xi_1, \xi_3, \xi_5\} = \{\xi_2, \xi_4, \xi_6\}$ or $\xi_6 \in \Lambda$.

**Proof of Lemma 5.4.** By symmetry we assume $|M| \geq |N|$ which leads $M_* = |M|$. Put 
$$\xi_j \over M_* = \eta_j + o \left( \frac{M + N}{M_*} \right),$$
for $1 \leq j \leq 6$, where $\eta_j \in \mathbb{Z}$. By hypothesis, $\eta_1 + \eta_3 + \eta_5 = \eta_2 + \eta_4 + \eta_6$, $\eta_1^2 - \eta_2^2 + \eta_3^2 - \eta_4^2 + \eta_5^2 - \eta_6^2 = 0$. Moreover from 
$N = (M+N) - M$, $-3M-N = -(M+N)-2M$, $-3N-M = -3(M+N)+2M$, 
we have $\eta_j \in \{-1, -1, 2, -2\}$ for $1 \leq j \leq 5$.

When $\{\xi_1, \xi_3, \xi_5\} \cap \{\xi_2, \xi_4, \xi_6\} \neq \emptyset$, we may assume $\xi_1 = \xi_2$ by symmetry. In this case, since $\eta_3 + \eta_5 = \eta_4 + \eta_6$, we have 
$$0 = \eta_3^2 + \eta_5^2 - \eta_4^2 - \eta_6^2 = 2(\eta_3 - \eta_4)(\eta_5 - \eta_6),$$
which implies $\{\xi_3, \xi_5\} = \{\xi_4, \xi_6\}$.

When $\xi_6 \notin \Lambda$, $\{\xi_1, \xi_3, \xi_5\} \cap \{\xi_2, \xi_4, \xi_6\} = \emptyset$ and $\{\xi_1, \xi_3, \xi_5\} \neq \{\xi_2, \xi_4, \xi_6\}$, then there exists pairs $\{\xi_1, \xi_2\} \subset \{\xi_1, \xi_3, \xi_5\}$ or $\{\xi_1, \xi_2\} \subset \{\xi_2, \xi_4, \xi_6\}$ such that $i \neq j$, but $\xi_i = \xi_j$. By symmetry we may assume $i = 1$ and $j = 3$. Note 
$$2\eta_1^2 + \eta_3^2 = \eta_2^2 + \eta_4^2 + (2\eta_1 + \eta_5 - \eta_2 - \eta_4)^2.$$

If $\{\eta_1, \eta_5\} \subset \{\pm 1\}$, then $\{\eta_2, \eta_4\} \subset \{\pm 1\}$, $\eta_1 = \eta_5$ and $\eta_2 = \eta_4 = -\eta_1$, since 
$3 = 1^2 + 1^2 + 1^2$, which implies $|\eta_6| = 5$, but this contradicts with $\eta_6 \in \{\pm 1\}$.

If $\eta_1 \in \{\pm 1\}$ and $\eta_5 \in \{\pm 2\}$, then $\eta_2 = -\eta_1$ (by symmetry) and $\{\eta_4, \eta_6\} = \{-\eta_1, -\eta_5\}$, since $6 = 1^2 + 1^2 + 2^2$. Then $2\eta_1 + \eta_5 = 0$ holds, which implies $\xi_6 \in \Lambda$ by (30) and (31). This is a contradiction.

If $\eta_1 \in \{\pm 2\}$ and $\eta_5 \in \{\pm 1\}$, then $\eta_2 = -\eta_1$ (by symmetry) and $\{\eta_4, \eta_6\} = \{-\eta_1, -\eta_5\}$ since $9$ is expressed in sum of square by $1^2 + 2^2 + 2^2$, which contradicts with $2\eta_1 + \eta_5 = 0$ as same as above.

If $\{\eta_1, \eta_5\} \subset \{\pm 2\}$, then $\eta_1 = \eta_5$ and $\eta_2 = \eta_4 = -\eta_1$ since $12 = 2^2 + 2^2 + 2^2$, which implies $3\eta_1 = -3\eta_1$. This is a contradiction. Hence we conclude that $\{\xi_1, \xi_3, \xi_5\} = \{\xi_2, \xi_4, \xi_6\}$. \square

By using this lemma, we obtain the estimate on the term $N_3$.

**Lemma 5.5.** Assume the same hypotheses of Proposition 3. Then, 
$$\left\| |(\xi)^\delta P_H N_3(t, \xi)\|_{L^2_{|\xi|\Lambda}} \right\|_{L^2_{|\xi|\Lambda}} \leq \frac{1}{100} \varepsilon + c \left( \lambda^2 + |\mu| \right) \varepsilon t. \quad (44)$$

**Proof of Lemma 5.5.** We split the sum into two cases:

(i) at least, one of $\xi_j$ are not elements of $\Lambda$, 
(ii) $\xi_j \in \Lambda$ for all $1 \leq j \leq 5$.

In the first case (i), we see that the contribution to the left-hand side of (44) is bounded by 
$$c|\mu| M_*^\delta \left( \| P_L b \|_{L^\infty_{\xi} \ell^2_{|\xi|\Lambda}} + M_*^{-\delta} \| \langle \xi \rangle^\delta P_H b \|_{L^\infty_{\xi} \ell^2_{|\xi|\Lambda}} \right) \lesssim |\mu| M_*^\delta (\varepsilon^2 + M_*^{-\delta} \varepsilon) \lesssim |\mu| \varepsilon t,$$
where we use \( \| \langle \xi \rangle^\delta b_\xi \|_{L^2_x L^2_t} \lesssim M_*^\delta \) from \( (43) \).

In the second case (ii), we may assume \( \xi_j \in A \) for \( 1 \leq j \leq 5 \), without loss of generality. Observe that by Lemma 5.4 if \( \xi_1 + \xi_3 + \xi_5 = \xi_2 + \xi_4 + \xi \) and \( |\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5 - \xi^2| \ll M_*^2 \), then \( \{ \xi_1, \xi_3, \xi_5 \} = \{ \xi_2, \xi_4, \xi \} \), but it obviously contradicts with the restriction of the formula of \( N_{3j} \). This leads to the bound
\[
|\xi_1^2 - \xi_2^2 + \xi_3^2 - \xi_4^2 + \xi_5 - \xi^2| \gtrsim M_*^2.
\]
The main point is that by \( (43) \),
\[
\left| \left\langle \phi(t, \xi) \right\rangle \right| \lesssim (|\lambda|M_* + |\mu|),
\]
which in particular implies that
\[
\left| \left\langle \phi_2(t, \xi_1, \xi_2, \xi_3, \xi_4, \xi_5) \right\rangle \right| \gtrsim M_*^2,
\]
(since \( |\lambda| \ll M_* \) and \( |\mu| \ll M_*^2 \)). Using integration by parts along with the above inequality \( (45) \), it follows that the contribution of this case to the left-hand side of \( (44) \) is bounded by
\[
\frac{c|\mu|}{M_*^2} \sum_j \langle \xi \rangle^\delta \left[ \left| \left\langle b_{\xi_1}(s) \overline{b_{\xi_2}(s)} b_{\xi_3}(s) \overline{b_{\xi_4}(s)} b_{\xi_5}(s) \rangle \right\rangle \right]_{s=0}^{s=t} + \frac{c|\mu|}{M_*^2} \int_0^t |\partial_s \left( b_{\xi_1}(s) \overline{b_{\xi_2}(s)} b_{\xi_3}(s) \overline{b_{\xi_4}(s)} b_{\xi_5}(s) \right) | ds.
\]
(46)

We use the a priori bound \( (43) \) to obtain that the first term of \( (46) \) is bounded by
\[
c|\mu| \| M_*^{-2+\delta} \|_{L^2_x L^2_t} \lesssim |\mu| M_*^{-2+\delta} \lesssim \frac{1}{100} \varepsilon,
\]
where \( \delta < 1 \) and \( M_* \) is large enough. Substituting \( \partial_s b_{\xi_1} \) to \( (46) \) and using \( (38) \), we easily see that the contribution of the second term of \( (46) \) is bounded by
\[
c|\mu| \left( |\lambda|M_*^{1+\delta} \| b_\xi \|_{L^2_x L^2_t} + (\lambda^2 + |\mu|) M_*^{-2+\delta} \right) t \lesssim (\lambda^2 + |\mu|) \varepsilon t,
\]
which completes the proof of Lemma 5.5.

Analogously, we have the following lemma, by replace \( \delta = 0 \).

**Lemma 5.6.** Assume the same hypotheses of Proposition 3. Then
\[
\| P_{L^2} N_3(t, \xi) \|_{L^2_t L^2_x} \lesssim \frac{1}{100} \varepsilon^2 + c \left( \lambda^2 + |\mu| \right) \varepsilon^2 t.
\]

**Proof of Lemma 5.6** From the proof of Lemma 5.5, it follows that
\[
\| P_{L^2} N_3(t, \xi) \|_{L^2_t L^2_x} \lesssim \frac{1}{100} \varepsilon^2 + c|\mu| t \left( \varepsilon^2 + M_*^{-\delta} \varepsilon + M_*^{-2} + |\lambda|M_*^{-1} \| b_\xi \|_{L^2_x L^2_t} + (\lambda^2 + |\mu|) M_*^{-2} \right),
\]
which is bounded by
\[
\frac{1}{100} \varepsilon^2 + c \left( \lambda^2 + |\mu| \right) \varepsilon^2 t.
\]

It remains to consider the term \( N_{11} \).

**Lemma 5.7.** Assume the same hypotheses of Proposition 3. Then
\[
\| \langle \xi \rangle^\delta P_H N_{11}(t, \xi) \|_{L^2_t L^2_x} \lesssim \frac{1}{100} \varepsilon + c \left( \lambda^2 + |\mu| \right) \varepsilon t.
\]
Proof of Lemma 5.7. We divide the sum on \((\xi_1, \xi_2, \xi_3) \not\in \Lambda^*\) into three cases:

(i) \(\xi_j \in \Lambda\) for all \(1 \leq j \leq 3\),
(ii) at least, two of \(\xi_j\) are not elements of \(\Lambda\),
(iii) two of \(\xi_j\) are elements of \(\Lambda\), one of \(\xi_j\) is not elements of \(\Lambda\), \(|\xi - \xi_1| \geq M_*/100\)
and \(|\xi - \xi_3| \geq M_*/100\).

In the case (i), we see that

\[
|\xi - \xi_1|(|\xi - \xi_3|) = |(\xi_2 - \xi_1)(\xi_2 - \xi_3)| \geq M_*^2, \tag{48}
\]

since \(M_* \gg |\lambda| \sim |M + N|\) and the constrain \(\{\xi_1, \xi_3\} \cap \{\xi_2, \xi_3\} = \emptyset\). As we saw in (45), the following inequality holds:

\[
\left[\Phi_1(s, \xi_1, \xi_2, \xi_3, \xi)\right] \gg \left|\left(\xi_1 - \xi_2\right)(\xi_1 - \xi)\right|. \tag{50}
\]

Using integration by parts, it follows that the contribution of this case to the left-hand side of (47) is bounded by

\[
c|\lambda| \sum_{\Lambda^*} \left|\frac{\langle \xi \rangle^{\delta_2}}{(\xi_1 - \xi_2)(\xi_1 - \xi)} \right| \left|\left[\bar{b}_{\xi_1}(s)b_{\xi_2}(s)b_{\xi_3}(s)\right]_{s=0}^{s=1}\right|
+c|\lambda| \int_0^t \sum_{\Lambda^*} \left|\frac{\langle \xi \rangle^{\delta_2}}{(\xi_1 - \xi_2)(\xi_1 - \xi)} \right| \left|\partial_s \left(b_{\xi_1}(s)b_{\xi_2}(s)b_{\xi_3}(s)\right)\right| \, ds. \tag{49}
\]

By (43) and (48), we have that the first term of (49) has the bound

\[
c|\lambda| M_*^{-1-\delta} \leq \frac{1}{100} \varepsilon, \tag{50}
\]

where \(\delta < 1\) and \(M_*\) is large enough. To deal with the second term in (49), we use the equation (42). By symmetry we may only consider the case that the differential operator \(\partial_s\) applies to \(b_{\xi_1}(s)\), and so the contribution of this case to the second term of (49) is bounded by

\[
c \sum_{j=1}^3 \left|\lambda| \int_0^t \sum_{\Lambda^*} M_*^{5-1} \left|\tilde{N}_j(s, \xi_1)\bar{b}_{\xi_2}(s)b_{\xi_3}(s)\right| \, ds. \tag{51}
\]

Fortunately, except for \(j = 1\) in (51), it follows that the contribution of these cases to (51) is bounded by

\[
c|\lambda|^3 M_*^{5-1} t \varepsilon^2 + c|\lambda\mu| M_*^{5-1} t (\varepsilon^2 + M_*^{1-2} + |\lambda| M_*^{-1} + |\mu| M_*^{-2})
\lesssim (\lambda^2 + |\mu|) \varepsilon t,
\]

where \(\delta < 1\) and \(M_*\) is large enough. For \(j = 1\), \(\tilde{N}_1(s, \xi_1)\) is expressed as

\[
i\lambda \sum_{\{\xi'_j, \xi'_1\} \cap \{\xi'_1, \xi_1\} = \emptyset} \xi'_2 b_{\xi'_1}(s)\bar{b}_{\xi'_2}(s)b_{\xi'_3}(s)e^{i\Phi_4(s, \xi', \xi)}.
\]

Since \(\xi_1 \in \Lambda\) and constrain (40), we have that at least one of \(\xi'_j\) is elements of \(\Lambda\), thus we divide the sum into two sub-cases (using symmetry):

(i)’ two of \(\xi'_j\) are not elements of \(\Lambda\),
(ii)’ \(\xi'_2 \in \Lambda, \xi'_1 \in \Lambda, \xi'_3 \not\in \Lambda\),
(iii)’ \(\xi'_2 \not\in \Lambda\) and \(\xi'_j \in \Lambda\) for \(j = 1, 3\).
In the case (iii)', it easily follows that the contribution to (51) is bounded by
\[ c\lambda^2 M_*^2 \varepsilon^2 t \lesssim \lambda^2 \varepsilon t. \]
In the case (ii)', note that \( \xi_1, \xi_1', \xi_2', \xi_3 \) satisfy
\[ \xi_1 + \xi_2' = \xi_1' + \xi_3', \quad \xi_1, \xi_1' \in \Lambda, \quad \xi_2' \in \Lambda, \quad \xi_3 \notin \Lambda, \]
which implies that either \( |\xi_3'| \lesssim |M + N| \approx |\lambda| \) or \( |\xi_3'| \gtrsim M_* \) holds. Then we have that the contribution of this case to (51) is bounded by
\[ c\lambda^2 t \left( M_*^2 \| P_L b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} + \| (\xi_3')^\delta P_H b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} \right) \lesssim \lambda^2 \varepsilon t. \]
In the case (iii)', by
\[ |\xi_2' b_{\xi_3}| \lesssim |\lambda| \left( |P_L b_{\xi_3}| + M_*^{1-\delta} |(\xi_3')^\delta P_H b_{\xi_3}| \right), \]
we have that the contribution to (51) is bounded by
\[ c\lambda^2 t \left( |\lambda|M_*^{1-\delta} \varepsilon^2 + \varepsilon \right) \lesssim \lambda^2 \varepsilon t. \]

We now consider the case (ii). As in the case (i)' above, it follows that the contribution of this case to (51) is bounded by
\[ c\lambda^2 t \left( \| P_L b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} + \| (\xi_3')^\delta P_H b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} \right)^2 \lesssim |\lambda|M_*^{1+\delta} \varepsilon^2 t (\varepsilon + M_*^{-\delta})^2 \lesssim \lambda^2 \varepsilon t, \]
for \( \delta \geq 1/2 \).

It thus remains to establish the bound in the case (iii). As in the case (i) above, by using (48), it follows that the contribution of this case to (51) has the same bound as in (49). The estimate for the first term of (51) is same, namely we have the bound (50). For the second term in (51), we divide two cases (using symmetry):
(iii)” the differential operator \( \partial_s \) falls in \( b_{\xi_1} \) with \( \xi_j \in \Lambda \),
(ii)” the differential operator \( \partial_s \) falls in \( b_{\xi_3} \) with \( \xi_j \notin \Lambda \).

In the case (i)”, the sum in \( \tilde{N}_1(s, \xi_j) \) contains at least one frequency elements not in \( \Lambda \). Thus as same as in the case (ii), we have that the contribution of this case to (51) is bounded by
\[ c\lambda^2 t \left( \| P_L b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} + \| (\xi_3')^\delta P_H b_{\xi_3} \|_{L^\infty T^\infty_{\xi_3 \not\in \Lambda}} \right)^2 \lesssim \lambda^2 \varepsilon^2 t (\varepsilon + M_*^{-\delta})^2 \lesssim \lambda^2 \varepsilon t. \]
In the case (ii)”, we also use the formula (49), but we will take the absolute value at last. For simplicity, we may assume that
\[ \xi_1 \notin \Lambda, \quad \xi_2 \in \Lambda, \quad \xi_3 \in \Lambda, \]
and the second term is represented as follows:
\[ c\lambda^2 \left| \int_0^t \sum_{s} \frac{(\xi_3')^\delta \xi_2' s}{(\xi_1 - \xi_3)(\xi - \xi_3)} b_{\xi_1}'(s)b_{\xi_2}'(s)b_{\xi_3}(s) b_{\xi_3}(s) e^{i\theta_2(s, \xi_1', \xi_2', \xi_3', \xi_3, \xi_3)} ds \right|, \quad (52) \]
where \( \sum_{s} \) denotes the sum over
\[ \xi_1' + \xi_3' = \xi_2 + \xi_1, \quad \xi_1 + \xi_3 = \xi_2 + \xi, \]
and
\[ \{\xi_1' - \xi_2 + \xi_3', \xi_3\} \cap \{\xi_2, \xi\} = \{\xi_1', \xi_3'\} \cap \{\xi_2', \xi_1\} = \emptyset. \]
In the case that there exist at least one \( \xi'_j \notin \Lambda \), we will use a similar argument to the case (i) in the proof of Lemma 5.5. Then we should consider the case that
\[
\xi_j' \in \Lambda \text{ for all } 1 \leq j \leq 3, \quad \xi_2 \in \Lambda, \quad \xi_3 \in \Lambda, \quad \xi \notin \Lambda.
\]
(53)
In the case that \( |\Phi_2(s, \xi_1', \xi_2', \xi_3', \xi_2, \xi_3, \xi)| \gtrsim M_3' \), using Lemma 5.5 it follows that the contribution this case to (52) is bounded by
\[
\begin{align*}
&c \lambda^2 M_3' t \left( \varepsilon^2 + M_3'^{-2} + |\mu| M_3'^{-2} \right) \lesssim \lambda^2 \varepsilon t,
\end{align*}
\]
for \( \delta < 1 \).

In the case that \( |\Phi_2(s, \xi_1', \xi_2', \xi_3', \xi_2, \xi_3, \xi)| \ll M_3' \), by Lemma 5.4 we see that \( \{ \xi_1', \xi_2', \xi_3 \} = \{ \xi_2, \xi_2, \xi \} \). It follows that \( \{ \xi_1', \xi_3' \} = \{ \xi_2, \xi \} \) and \( \xi_3 = \xi_2' \) (also \( \Phi_2(s, \xi_1', \xi_2', \xi_3', \xi_2, \xi_3, \xi) = 0 \) holds), which contradicts with (53).

Hence it completes the proof of Lemma 5.7.

\[ \square \]

Lemma 5.8. Assume the same hypotheses of Proposition 3 Then
\[
\| P_L N_{11}(t, \xi) \|_{L^2_t L^2_{\xi, \eta}} \lesssim \frac{1}{100} \varepsilon^2 + c \lambda^2 \varepsilon^2 t.
\]

Proof of Lemma 5.8 The proof is the same as that of Lemma 5.7.

Proof of Proposition 3 We are now in the position to prove Proposition 3. The proof is a consequence of Lemma 5.3, Remark 9, Lemmas 5.5, 5.6, 5.7 and 5.8.

We apply the bootstrap argument to Proposition 3 and have the following proposition.

Proposition 4. Let \( \delta \in [1/2, 1] \) and \( a_\xi(t) \) be a global solution to (10). Then for all \( |t| \ll \min \{ \frac{1}{M_3'}, \frac{1}{100} \} \),
\[
M^0_\delta \| P_L a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}} + \| \langle \xi \rangle \delta P_H a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}} \lesssim \frac{1}{M_3'}.
\]

Proof. We shall use Proposition 3 with \( \varepsilon = \frac{1}{M_3'} \), which yields that for if
\[
\max \left\{ M^0_\delta \| P_L a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}}, \| \langle \xi \rangle \delta P_H a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}} \right\} \lesssim \varepsilon,
\]
for \( |t| \leq T \), then we have
\[
\max \left\{ M^0_\delta \| P_L a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}}, \| \langle \xi \rangle \delta P_H a_\xi(t) \|_{L^2_t L^2_{\xi, \eta}} \right\} \lesssim \frac{1}{100} + c (\lambda^2 + |\mu|) \varepsilon t,
\]
for \( |t| \leq T \). By bootstrap argument and notation (41) combined with a continuity argument, it will suffice to show that under the same hypothesis of Proposition 3
\[
\max \left\{ M^0_\delta \| P_L \tilde{b}_\xi(t) \|_{L^2_t L^2_{\xi, \eta}}, \| P_H \langle \xi \rangle \delta \tilde{b}_\xi(t) \|_{L^2_t L^2_{\xi, \eta}} \right\} \lesssim |\lambda| M_* \varepsilon t.
\]
(54)
Note that \( (\xi_1, \xi_2, \xi_3) \in \Lambda^{**}(\xi) \) implies that
\[
|\xi_1| \sim |\xi_2| \sim M_*,
\]
and there is one of \( \xi_j \notin \Lambda \) such that \( (\xi_1, \xi_2, \xi_3) \in \Lambda^{**}(\xi) \), which implies that the first term of (54) is eliminated. A direct computation allows us to obtain (54).

Proof of Theorem 1.2 The proof of the corollary follows by Propositions 1, 2 and 4. □
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