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log-log blow up solutions blow up at exactly m points

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Abstract

We study the focusing mass-critical nonlinear Schrödinger equation, and construct certain solutions which blow up at exactly *m* points according to the log–log law.

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Résumé

Nous étudions l'équation de Schrödinger non linéaire focalisante de masse critique, et construisons certaines solutions avec exactement m points d'explosion d'après la loi de log-log.

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1. Introduction

We consider the Cauchy Problem for the mass-critical focusing nonlinear Schrödinger equation (NLS) on \mathbb{R}^d for d = 1, 2:

$$(NLS)\begin{cases} iu_t = -\Delta u - |u|^{\frac{4}{d}}u, \\ u(0) = u_0 \in H^1(\mathbb{R}^d). \end{cases}$$
(1.1)

Problem (1.1) has three conservation laws:

- Mass:

$$M(u(t,x)) := \int |u(t,x)|^2 dx = M(u_0), \tag{1.2}$$

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- Energy:

$$E(u(t,x)) := \frac{1}{2} \int |\nabla u(t,x)|^2 dx - \frac{1}{2 + \frac{4}{d}} \int |u(t,x)|^{2 + \frac{4}{d}} dx = E(u_0),$$
(1.3)

- Momentum:

$$P(u(t,x)) := \Im(\int \nabla u(t,x)\overline{u(t,x)})dx = P(u_0), \tag{1.4}$$

and the following symmetry:

- 1. Space–time translation: If u(t, x) solves (1.1), then $\forall t_0 \in \mathbb{R}, x_0 \in \mathbb{R}^d$, we have $u(t t_0, x x_0)$ solves (1.1). 2. Phase transformation: If u solves (1.1), then $\forall \theta_0 \in \mathbb{R}$, we have $e^{i\theta_0}u$ solves (1.1).
- 3. Galilean transformation: If u(t, x) solves (1.1), then $\forall \beta \in \mathbb{R}^d$, we have $u(t, x \beta t)e^{i\frac{\beta}{2}(x \frac{\beta}{2}t)}$ solves (1.1). 4. Scaling: If u(t, x) solves (1.1), then $\forall \lambda \in \mathbb{R}_+$, we have $u_{\lambda}(t, x) := \frac{1}{\lambda^{\frac{d}{2}}}u(\frac{t}{\lambda}, \frac{x}{\lambda})$ solves (1.1).
- 5. Pseudo-conformal transformation: If u(t, x) solves (1.1), then $\frac{1}{t^{\frac{d}{2}}} \bar{u}(\frac{1}{t}, \frac{x}{t}) e^{i\frac{|x|^2}{4t}}$ solves (1.1).

1.1. Setting of the problem and statement of the main result

The equation (1.1) is mass-critical since the conserved quantity given by the mass is invariant under scaling symmetry. It is called focusing since the conserved quantity given by the energy is not coercive.

The NLS (1.1) was proved to be locally well-posed (LWP) in H^1 by Ginibre and Velo in [5]. This means that for any initial data in H^1 , gives rise to a unique solution $u \in C([0, T]; H^1)$, for some time T and the solution depends continuously on the initial data. See also [1] for a L^2 local theory.

Since this equation is mass-critical, and H^1 is a norm subcritical with respect to the L^2 norm, one can take T = $T(||u_0||_{H^1}) > 0$. By LWP it is not hard to see that if the solution is defined in [0, T_0) and cannot be extend beyond T_0 , then it has to be that $\lim_{t\to T_0} \|u(t)\|_{H^1} = \infty$. In this case we say that *u* blows up in finite time in T_0 . A solution that does not blow up in finite time is called global.

It turns out that not all solutions of (1.1) are global. The classical virial identity ([6]) indicates the existence of solutions which blow up in finite time. By direct calculation one has:

$$\partial_t^2 \int |x|^2 |u|^2 = 4\partial_t \Im(\int x \nabla u \bar{u}) = 16E(u_0), \tag{1.5}$$

and (1.5) immediately indicates that if $xu_0 \in L^2$ and $E(u_0) < 0$, then the solution u must blow up in finite time.

Some questions are then natural: If a solution blows up in finite time, what is the mechanism for singularity formation, i.e. how does one describes the reason for a solution to blow up and how does one describe the behavior of the solution when it approaches the blow up time?

Virial identity (1.5) on its own does not give the answer to these questions, and up to the best of our knowledge, blow up solutions to (1.1) blow up way faster than the blow up time predicted by (1.5).

Before we continue, let's first introduce the ground state, which is one typical object appearing in the study of focusing equations. The ground state Q = Q(x) is the unique positive solution in $L^2(\mathbb{R}^d)$ that solves

$$-\Delta W + W = |W|^{\frac{4}{d}}W.$$
(1.6)

Q is an explicit function when d=1, $Q(x) = \left(\frac{3}{ch2x}\right)^{\frac{1}{4}}$, and Q is smooth and decays exponentially for d = 1, 2. It is a direct calculation to check that $Q(x)e^{it}$ is a solution of (1.1). In general we call solutions to (1.1) of the form $W(x)e^{it}$ soliton solutions. We have several quick remarks:

- $W(x)e^{it}$ solves (1.1) if and only if W(x) solves (1.6).
- A pure variational argument ([29]) shows that any H^1 solutions to (1.1) which initial data u_0 has mass below the mass of the ground state, does not blow up.

- To compare with (1.5), one can use variational arguments ([29]) to show that for any function $f \in H^1$ such that $||f||_2 \le ||Q||_2$, one has $E(f) \ge 0$, and in particular one can derive from (1.6) that E(Q) = 0.

By applying the pseudo-conformal transformation to the global soliton solution $Q(x)e^{it}$, one obtains the explicit blow up solution

$$S(t,x) := \frac{1}{|t|^{\frac{d}{2}}} Q(\frac{x}{t}) e^{-\frac{|x|^2}{4t} + \frac{i}{t}}.$$
(1.7)

We remark that $||S(t, x)||_2 = ||Q||_2$, and this is (up to symmetries) the minimal mass blow up solution, [15] (recall no H^1 solution with mass below $||Q||_2^2$ can blow up in finite time). We also remark by an easy scaling argument, that if a solution *u* to (1.1) blows up in finite time T, then its blow up rate for $||\nabla u(t)||_2$ has a lower bound:

$$\|\nabla u(t)\|_2 \gtrsim \frac{1}{\sqrt{T-t}},\tag{1.8}$$

and the explicit blow up solution S(t, x) has blow up rate $\frac{1}{|t|}$.

Despite the fact that the Schrödinger equation has the infinite speed of propagation, S(t, x) blows up locally in the physical space, and more precisely by looking at its explicit formula (1.7), we can say that the solution blows up at $x_0 := 0 \in \mathbb{R}^d$. Relying on the fact that the blow up behavior of S(t, x) is local, by a certain compactness argument and using S(t, x) as basic building blocks, Merle ([13]) constructed a solution u to (1.1) which blows up at k points in finite time T, and near blow up time T it has the following asymptotic:

$$u(t,x) \sim \sum_{i=1}^{k} S(\frac{t-T}{\lambda_i^2}, \frac{x-x_i}{\lambda_i}), \lambda_i > 0, x_i \neq x_j, \forall i \neq j.$$

$$(1.9)$$

The goal of this paper is to use the log–log blow up solutions, that we recall below, as basic building blocks instead of S(t, x) in (1.7) to construct blow up solutions to (1.1). The so-called log–log blow up solutions are solutions to (1.1) which blow up in finite time T with blow up rate $\|\nabla u(t)\|_2 \sim \left(\frac{\ln|\ln|T-t||}{T-t}\right)^{\frac{1}{2}}$. Such solutions had been suggested numerically by Landman, Papanocolaou, Sulem, Sulem, [11] and first constructed by Perelman, [22], and later intensively studied by Merle and Raphaël, [21,16,20,17,24,19]. Merle and Raphaël consider the solutions to (1.1) with initial data u_0 such that

$$\|Q\|_2 < \|u_0\|_2 < \|Q\|_2 + \alpha, \tag{1.10}$$

where α is a universal small positive constant.

Let's give a quick summary of the Merle and Raphaël's results relevant for this paper.

When (1.10) holds, they show that all solution with strict negative energy blow up in finite time *T* according to the log–log law: $\|\nabla u(t)\|_2 \sim \left(\frac{\ln|\ln|T-t||}{T-t}\right)^{\frac{1}{2}}$. And more precisely, near blow up time *T*, the solution *u* can be decomposed as:

$$u(t,x) = \frac{1}{\lambda(t)^{\frac{d}{2}}} Q(t, \frac{x - x(t)}{\lambda(t)}) + \Xi(t,x),$$

$$\lambda(t) \in \mathbb{R}^+, \frac{1}{\lambda(t)} \sim \sqrt{\frac{\ln|\ln T - t|}{T - t}},$$

$$\lim_{t \to T} \lambda(t) \|\Xi(t)\|_{H^1} = 0.$$
(1.11)

They also show that the log-log blow up solution is stable in the sense that the initial data, which generate log-log blow up solutions, is an open set in H^1 .

Remark 1.1. Merle and Raphaël's results are more general than what has been recalled above. For example the strict negative condition can be relaxed.

As mentioned above, the log-log blow up behavior is stable under H^1 perturbation. What's more, this log-log blow up behavior is local (in physical space) in a certain sense, in spite of infinity speed of propagation of the NLS equation. In [23], Planchon and Raphaël constructed a log-log blow up solutions on a bounded domain Ω of \mathbb{R}^d with $u(t) \in H_0^1(\Omega)$. In particular, if one looks at solutions of the form (1.11), they prove that x(t) has a limit as t approaches the blow up time, hence it makes sense to say that such solutions blow up at a certain point² $x_{\infty} := \lim_{t \to T} x(t)$ in \mathbb{R}^d . The key element in their proof is a very robust bootstrap argument. It turns out that the analysis of log-log blow up solutions has a very suitable bootstrap structure.

In this work, as an analogue of [13], we use $\log -\log$ blow up solutions as building blocks to construct solutions which blow up at (exactly) *m* points. To be precise, we show:

Theorem 1.1. For d = 1, 2, for each positive integer m, and given any m different points $x_{1,\infty}, \ldots x_{m,\infty}$ in \mathbb{R}^d , there exists a solution u to (1.1) such that u blows up in finite time T, and for t close enough to T,

$$u(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_{j}^{\frac{d}{2}}(t)} Q(\frac{x-x_{j,\infty}}{\lambda_{j}(t)}) e^{-i\gamma_{j}(t)} + \Xi(t,x),$$
(1.12)

where, for j = 1, ..., m,

$$\frac{1}{\lambda_j(t)} \sim \sqrt{\frac{\ln|\ln T - t|}{T - t}}, \quad and \quad \lambda_j \|\Xi(t)\|_{H^1} \xrightarrow{t \to T} 0, \tag{1.13}$$

i.e. Ξ can be viewed as an error term. In particular, since the m given points are arbitrary, the solutions do not necessarily have any symmetry restriction.

Remark 1.2. Essentially same construction as in the proof of Theorem 1.1, works on torus, \mathbb{T}^d , d = 1, 2.

Remark 1.3. The result of [23] already implies the existence of symmetric solutions which blow up at two points according to log–log law. In fact one first constructs a solution to NLS on the half line/plane $H := \mathbb{R}^d_+$, such that the solution blows up at one point according to log–log law and satisfies the Dirichlet condition $u \equiv 0$ on $\partial \mathbb{R}^{d+}$, then, by extending this solution symmetrically, one easily derives the solution which blows up at two points in the whole line/plane. Similarly, one can construct solutions that blow up at a even number of points according to log–log law, but they will have a very strong symmetry. See Corollary 1 in [23].

Remark 1.4. The work of [13] uses idea of "integrate from infinity", which means one needs to evolve the data backward. This does not seem to direct work in this setting because of the remainder term $\Xi(t, x)$. In this work, we will evolve the data forward.

Remark 1.5. We point out two applications of Theorem 1.1. First, it implies the existence of large mass log–log blow up solutions. More general results on this direction have been obtained in [18]. Second, for those who are familiar with standing ring blow up solutions, [25,26], our construction in 1 D case implies the existence of multiple-standing-ring blow up solutions for quintic NLS on dimension $N \ge 2$. To be precise, one can construct a radial solution *u* to the following Cauchy Problem:

$$\begin{cases} iu_t + \Delta u = -|u|^4 u, \\ u_0 \in H^{N_2(N)}(\mathbb{R}^N), \end{cases}$$
(1.14)

such that u blows up in finite time with log–log blow up rate and near blow up time T, u has the following asymptotic

$$u(t,x) = u(t,r) = \sum_{j=1}^{m} \frac{1}{\lambda_j(t)^{1/2}} P(\frac{r-r_j(t)}{\lambda_j(t)}) e^{-i\gamma_j(t)} + \Xi(t,x),$$
(1.15)

² Note that we already remarked this property for explicit blow up solutions S(t, x).

$$\frac{1}{\lambda_j(t)} \sim \sqrt{\frac{\ln|\ln T - t|}{T - t}}, \quad \text{and} \quad \lambda_j \|\Xi(t)\|_{H^1} \xrightarrow{t \to T} 0, \tag{1.16}$$

$$\lim_{t \to T} r_j(t) = r_{j,\infty} > 0, \tag{1.17}$$

where P is the unique positive L^2 solution which solves (1.6) for d = 1.

From this point on, we fix $m \in \mathbb{N}$ as the number of blow up points. We will do proof only for d = 2, the case d = 1 follows similarly.

To construct a solution satisfying Theorem 1.1 without relying on certain symmetry property of the initial data, the most intuitive process to follow is that one first prepares $m \log -\log b \log s$ solutions

$$u^{j}(t,x) = \frac{1}{\lambda_{i}^{\frac{d}{2}}} \mathcal{Q}(\frac{x-x_{j}(t)}{\lambda_{j}})e^{-i\gamma_{j}(t)} + \Xi_{j}, j = 1, \dots, m$$

satisfying (1.11). Then one shows that the solution to (1.1) with initial data $\sum_{j=1}^{m} u^{j}(0)$ evolves approximately as $\sum_{j=1}^{m} u^{j}(t)$, if one assumes that at the initial time all solutions are very close to blowing up, i.e. $\lambda_{j}(0)$, j = 1, ..., m is very small, and they are physically separated, i.e. $\min_{j \neq j'} |x_{j}(0) - x_{j'}(0)|$ is very large. To achieve this, one needs some mechanism to decouple u^{1} , u^{2} , ... u^{m} . Our choice is to require extra smoothness outside the (potential) singular points. Roughly speaking, we find neighborhoods U_{j} of $x_{j,0} \equiv x_{i}(0)$, j = 1, ..., m, and show that the solutions keep very high regularity outside these *m* neighborhoods. This approach is motivated by the work of Raphaël and Szeftel, [26]. They consider the focusing quintic NLS on \mathbb{R}^{d} and they require their data to be radial and in H^{d} . By the radial symmetry assumption, the problem can be understood in polar coordinates as a perturbation of the 1D-focusing-quintic NLS. The goal of Raphaël and Szeftel is to construct solutions that blow up at a sphere (or a ring). The crucial point in their paper is to understand the propagation of singularity/regularity. They show that all the singularities are kept around the sphere where the solution is supposed to blow up, and the solution is kept bounded in $H^{\frac{d-1}{2}}$ outside the sphere. Note that, thanks to the radial assumption, the authors are using the 1D NLS to model their solutions [26]. See also [25], which indeed has the same spirit as [26], but in the setting [25], one does not need to pursue high regularity.

We show the following theorem.

Theorem 1.2 (Propagation of regularity). For any given $K_1 > 1$, $K_2 > 2$, not necessarily integers, if $K_1 < \frac{K_2}{2}$, then we construct a solution u in H^{K_2} to (1.1) that blows up according to the log–log law as (1.11) at finite time T and such that

 $\begin{aligned} &- \sup_{t \in (0,T)} |x(t) - x(0)| < \frac{1}{1000}. \\ &- u(t,x) \text{ is bounded in } H^{K_1} \text{ when restricted in } |x - x_0| ≥ 1/2. \end{aligned}$

Remark 1.6. The choice of special numbers $\frac{1}{2}$, 1/1000 is (of course) just for concreteness and simplicity.

Remark 1.7. When d = 1, and K_2 is an integer, and one can prove Theorem 1.2 for $K_1 \le \frac{K_2-1}{2}$ by sightly modifying the language of [26]. Their method is based on a bootstrap argument and a certain pseudo-energy. When d = 2, Raphaël and Szeftel's method does not seem to directly work, and one should be able to use the argument in [30] to prove Theorem 1.2 for $K_1 \le \frac{K_2-1}{2}$ when K_2 is integer. Our proof improves the previous results in the two aspects. We can take $K_1 < \frac{K_2}{2}$ and we do not require K_2 to be integer. Our proof is written more in a harmonic analysis style, relying on the (upside-down) I-method, [2,27], interpolation and Strichartz estimate.

Remark 1.8. When $K_2 \ge K_1 = 1$, Theorem 1.2 is implied by the work of Holmer and Roudenko in [8].

One should understand Theorem 1.2 as a proof for the fact that the log–log blow up behavior is local, in the sense that it does not propagate singularity outside the blow up point. This will help us to decouple the m "solitons" (we sometimes also call them bubbles), in our construction of blow up solution.

Remark 1.9. Because of the good localized property of log–log blow up, one can even work on manifolds, see [7] for work in this direction.

Remark 1.10. It is not always true in this kinds of problems (if one doesn't put some restriction on the data), that the *m* "solitons" or *m* bubbles will be decoupled. Different bubbles may interact with each other in a strong way. See very recent work [12] for this direction.

Once one can somehow decouple the *m* soliton $u^1, \ldots u^m$, then we will use some topological argument to construct initial data, and balance those *m* bubbles and make them blow up in the same time. And, prescription of the blow up points is actually more subtle than making *m* bubbles blow up at the same time. Fortunately, by taking advantage of the sharp dynamic of log–log blow up, it can still be achieved by certain topological argument. We remark, it is typical that one may rely on soft topological argument rather than pure analysis to prove things like this, see [14,23,4], through one needs to find different topological argument in different settings.

1.2. Notation

Throughout this article, α is used to denote a universal small number, $\delta(\alpha)$ is a small number depending on α such that $\lim_{\alpha \to 0} \delta(\alpha) = 0$. We use $\delta_0, \delta_1, \ldots$ to denote universal constant (they are usually small, but don't depend on α). We use *C* to denote a large constant, it usually changes line by line. We also use *c*, η and *a* to denote small constants. For any constant *r*, we use $r \pm$ to mean $r \pm \delta$ where δ is a small positive constant.

We write $A \leq B$ when $A \leq CB$, for some universal constant C, we write $A \geq B$ if $B \leq A$. We write $A \sim B$ if $A \leq B$ and $B \leq A$. As usual, $A \leq_{\sigma} B$ means that $A \leq C_{\sigma} B$, where C_{σ} is a constant depending on σ . We use Λ to denote the operator $\frac{d}{2} + y\nabla$ on $H^1(\mathbb{R}^d)$. We use the notation

$$\epsilon_1 := \Re \epsilon, \epsilon_2 := \Im \epsilon, \text{ i.e. } \epsilon = \epsilon_1 + i\epsilon_2, \tag{1.18}$$

where \Re is the real part and \Im is the imaginary part.

We use usual functional spaces L^p, C^1, \ldots, C^k and C^∞ , we will also use Sobolev space $H^s, s \in \mathbb{R}$.

If not explicitly pointed out, L^p means $L^p(\mathbb{R}^d)$, so for the other spaces. We also use $L^q_t L^p_x$ to denote $L^q(\mathbb{R}; L^p(\mathbb{R}^d))$. When a certain function is only defined on $I \times \mathbb{R}^d$, we also use the notation $L^q(I; L^p(\mathbb{R}^d))$. Sometimes we use $||f||_p$ to denote $||f||_{L^p}$.

We use (,) to denote the usual L^2 (complex) inner product.

Finally, for a solution u(t, x), we use $(T^{-}(u), T^{+}(u))$ to denote its lifespan.

1.3. A quick review of Merle and Raphaël's work and heuristics for the localization of log-log blow up

Let us quickly review the work of Merle and Raphaël and highlight the bootstrap structure related to it. At the starting point of their series of work, in [21], they consider a solution u to (1.1) with initial data $u_0 \in H^1$ satisfies (1.10), with zero momentum and strictly negative energy. They rely on the following variational argument:

Lemma 1.1 (Lemma 1 in [21]). For an arbitrary function $f \in H^1$, with energy $E(f) \leq 0$, if also f satisfies (1.10), then one can find parameters $\lambda_0 \in \mathbb{R}^+$, $x_0 \in \mathbb{R}^d$, $\gamma_0 \in \mathbb{R}$ and $\epsilon \in H^1$, such that

$$e^{i\gamma_0}\lambda_0^{\frac{\alpha}{2}}f(\lambda_0 x + x_0) = Q + \epsilon, \qquad (1.19)$$

and

$$\|\epsilon\|_{H^1} \le \delta(\alpha). \tag{1.20}$$

This lemma implies that for the special solution u(t) to (1.1) considered by Merle and Raphaël, one has the geometric decomposition

$$u(t) = \frac{1}{\lambda(t)^{\frac{d}{2}}} \left(Q + \epsilon(t) \right) \left(\frac{x - x(t)}{\lambda(t)} \right) e^{-i\gamma(t)}, \tag{1.21}$$

$$\|\epsilon(t)\|_{H^1} \le \delta(\alpha). \tag{1.22}$$

Note that one has some freedom in choosing the three parameters $\lambda(t)$, x(t) and $\gamma(t)$. Because of this freedom one can further use the modulation theory to derive the next lemma.

Lemma 1.2 (Lemma 2 in [21]). Let u(t) be the solution to (1.1) with initial data u_0 , which has zero momentum, strictly negative energy and let u_0 satisfy (1.10). Then within the lifespan of u(t), there are three unique parameters $\lambda(t), x(t), \gamma(t)$ and $\epsilon = \epsilon_1 + i\epsilon_2 \in H^1$, such that

$$u(t) = \frac{1}{\lambda(t)^{\frac{d}{2}}} (Q + \epsilon(t)) \left(\frac{x - x(t)}{\lambda(t)}\right) e^{-i\gamma(t)},$$
(1.23)

$$\|\epsilon(t)\|_{H^1} \le \delta(\alpha),\tag{1.24}$$

$$(\epsilon_1, \Lambda Q) = (\epsilon_1, yQ) = (\epsilon_2, \Lambda^2 Q) = 0.$$
(1.25)

Now the study of (1.1) is transferred to the study of the evolution of system $\{\epsilon(t), x(t), \lambda(t), \gamma(t)\}$. We remark here that the blow up rate is determined by the parameter $\lambda(t)$.

In this setting, Merle and Raphaël are using the ground state Q to approximate the solution u(t) (up to space translation, scaling and phase transformation). It turns out that sharper results can be obtained by using \tilde{Q}_b , a modification of Q, [16,20]. Let's give a brief description of \tilde{Q}_b , see Proposition 1 in [20] for details.

Let $b \in \mathbb{R}$, $\eta \in \mathbb{R}_+$ be small enough, η is fixed. Let us define

$$R_b := \frac{2}{|b|} \sqrt{1 - \eta}, \quad R_b^- := \sqrt{1 - \eta} R_b, \tag{1.26}$$

and let ϕ_b be a smooth cut-off function which equals 1 on $|x| \le R_b^-$ and vanishes for $|x| \ge R_b$. Then the modified profile $\tilde{Q}_b := Q_b \phi_b$, where Q_b solves the equation

$$\begin{cases} \Delta Q_b - Q_b + ib\Lambda Q_b + |Q_b|^{4/d} Q_b = 0, \\ P_b \equiv Q_b e^{i\frac{b|y|^2}{4}} > 0 \text{ in } B_{R_b}, Q_b(0) \in (Q(0) - \epsilon^*(\eta), Q(0) + \epsilon^*(\eta)), Q_b(R_b) = 0. \end{cases}$$
(1.27)

Here we also define Ψ_b

$$\Psi_b = -\Delta \tilde{Q}_b + \tilde{Q}_b - ib\Lambda \tilde{Q}_b - \tilde{Q}_b |\tilde{Q}_b|^{4/d}, \qquad (1.28)$$

that will be used later.

We now list some useful estimates for \tilde{Q}_b :

1. \tilde{Q}_b is uniformly close to Q in the sense:

$$\|e^{(1-\eta)\frac{\theta(|b||y|)}{|b|}}(\tilde{Q}_{b}-Q)\|_{C^{3}}+\|e^{(1-\eta)\frac{\theta(|b||y|)}{|b|}}(\frac{\partial}{\partial_{b}}\tilde{Q}_{b}+i\frac{|y|^{2}}{4}Q)\|\xrightarrow{b\to 0}0,$$
(1.29)

where

$$\theta(r) = 1_{0 \le r \le 2} \int_{0}^{r} \sqrt{1 - \frac{z^2}{4}} dz + 1_{r>2} \frac{\theta(2)}{2} r.$$
(1.30)

- 2. \tilde{Q}_b is supported in $|y| \lesssim \frac{1}{|b|}$.
- 3. \tilde{Q}_b has strictly super critical mass:

$$0 < \frac{d}{d|b|^2} \|Q_b\|^2 < \infty, \tag{1.31}$$

i.e. $||Q_b||^2 - ||Q||_2^2 \sim b^2$.

4. \tilde{Q}_b is uniformly bound in H^s , $s \in R$ for all b small enough (recall we only consider d = 1, 2)

$$\|\tilde{Q}_b\|_{H^s} \lesssim_s 1. \tag{1.32}$$

Remark 1.11. $\theta(2) = \frac{\pi}{2}$.

Remark 1.12. Estimate (1.32) is implied by estimate (1.29) when $s \le 3$. However, Merle and Raphaël consider the C^3 rather than the general C^k convergence in (1.29) only due to the fact the nonlinearity $|Q|^{\frac{4}{d}}Q$ itself is not smooth enough when $d \ge 3$. Thus for d = 1, 2, since the nonlinearity is algebraic, (1.32) holds for all s. And indeed it is not hard to directly use standard elliptic estimates to prove (1.32) for $s \ge 3$ once we know this holds for $s \le 3$. This fact is already implicitly used in [26] for d = 1.

Remark 1.13. Later in this work many terms will involve $C\eta$ but we will be able to fix η , such that $C\eta$ is as small as we want.

Remark 1.14. Note that since Q decays exponentially and \tilde{Q}_b is uniformly close to Q, it is standard that for given $N \in \mathbb{N}$, terms of the form $(f, \tilde{Q}_b), (\nabla^N f, \tilde{Q}_b), (y^N f, \tilde{Q}_b)$ are controlled by $\int (|\nabla f|^2 + |f|^2 e^{-|y|})^{\frac{1}{2}}$. This is widely used in [16,20].

With \tilde{Q}_b , Merle and Raphaël modify the Lemma 1.2 to the following:

Lemma 1.3 (Lemma 2 in [16]). Let u(t) be the solution to (1.1) with initial data u_0 , which has zero momentum, strictly negative energy and satisfies (1.10). Then within the lifespan of the u(t), there are unique parameter $\{b(t), \lambda(t), x(t), \gamma(t)\} \in \mathbb{R} \times \mathbb{R}^+ \times \mathbb{R}^d \times \mathbb{R}$ such that

$$u(t,x) = \frac{1}{\lambda(t)^{\frac{1}{2}}} (\tilde{Q}_b + \epsilon) \left(\frac{x - x(t)}{\lambda(t)}\right) e^{-i\gamma(t)},\tag{1.33}$$

$$\|\epsilon\|_{H^1} + |b| \le \delta(\alpha), \tag{1.34}$$

$$(\epsilon_1, |y|^2 \Sigma_b) + (\epsilon_2), |y|^2 \Theta_{b(t)}) = 0,$$
(1.35)

$$(\epsilon_1, y\Sigma_b) + (\epsilon_2, y\Theta_b) = 0, \tag{1.36}$$
$$(1.36)$$
$$(1.37)$$

$$-(\epsilon_1, \Lambda^2 \Theta_b) + (\epsilon_2, \Lambda^2 \Sigma_b) = 0, \tag{1.37}$$

$$-(\epsilon_1, \Lambda \Theta_b) + (\epsilon_2, \Lambda \Sigma_b) = 0. \tag{1.38}$$

Here $\tilde{Q}_h := \Sigma_h + i \Theta_h$.

We note that (1.34) a priori assures that the whole analysis is of perturbative nature. Again, the study of (1.1)is transferred to the study of the system $\{\epsilon(t), \lambda(t), \gamma(t), b(t), x(t)\}$. To analyze this system it is essential that one considers the slowly varying time variable *s* rather than the *t*:

$$\frac{dt}{ds} = \frac{1}{\lambda^2}.$$
(1.39)

Note that this change of variable changes the lifespan of u (in t variable) to the whole \mathbb{R} (in s variable), no matter if the original solution *u* blows up in finite time or not.

Now u satisfying (1.1) is equivalent to $\epsilon_1, \epsilon_2, b(s), \lambda(s), x(s), \gamma(s)$ satisfying the system³:

$$b_{s}\frac{\partial\Sigma}{\partial b} + \partial_{s}\epsilon_{1} - M_{-}(\epsilon) + b\Lambda\epsilon_{1} = (\frac{\lambda_{s}}{\lambda} + b)\Lambda\Sigma + \tilde{\gamma}_{s}\Theta + \frac{x_{s}}{\lambda}\nabla\Sigma + (\frac{\lambda_{s}}{\lambda} + b)(\Lambda\epsilon_{1}) + \tilde{\gamma}_{s}\epsilon_{2} + \frac{x_{s}}{\lambda}(\nabla\epsilon_{1}) + \Im\Psi_{b} - R_{2}(\epsilon)$$
(1.40)
$$b_{s}\frac{\partial\Theta}{\partial b} + \partial_{s}\epsilon_{2} + M_{+}(\epsilon) + b\Lambda\epsilon_{2} = (\frac{\lambda_{s}}{\lambda} + b)\Lambda\Theta - \tilde{\gamma}_{s}\Sigma + \frac{x_{s}}{\lambda}\nabla\Theta + (\frac{\lambda_{s}}{\lambda} + b)\Lambda\epsilon_{2} - \tilde{\gamma}_{s}\epsilon_{1} + \frac{x_{s}}{\lambda}(\nabla\epsilon_{2}) - \Re\Psi_{b} + R_{1}(\epsilon).$$
(1.41)

³ Here we slightly abuse notation. For example, x(s) actually means x(t(s)).

Here⁴ we have $\tilde{\gamma}(s) = -s - \gamma(s)$, $\tilde{Q}_b = \Sigma_b + i\Theta_b$ and $M = (M_+, M_-)$ is the linearized operator near the profile \tilde{Q}_b and R_1 , R_2 are nonlinear terms. Interested readers may consult (2.31), (2.32) in [20] for more details.

Now if one plugs in the four orthogonality condition (1.35), (1.36), (1.37), (1.38), one can obtain a system for the parameters $\{\lambda, x, \tilde{\gamma}, b\}$, i.e. four ordinary equations involving $\{\lambda(s), x(s), \tilde{\gamma}(s), b(s)\}^5$:

$$\frac{d}{ds}\{(\epsilon_1(t), |y|^2 \Sigma_{b(t)}) + (\epsilon_2(t), |y|^2 \Theta_{b(t)})\} = 0,$$
(1.42)

$$\frac{d}{ds}\{(\epsilon_1(t), y\Sigma_{b(t)}) + (\epsilon_2(t), y\Theta_{b(t)})\} = 0,$$
(1.43)

$$\frac{d}{ds}\left\{-(\epsilon_1(t), \Lambda^2 \Theta_{b(t)}) + (\epsilon_2(t), \Lambda^2 \Sigma_{b(t)})\right\} = 0,$$
(1.44)

$$\frac{d}{ds}\left\{-(\epsilon_1(t), \Lambda\Theta_{b(t)}) + (\epsilon_2(t), \Lambda\Sigma_{b(t)})\right\} = 0.$$
(1.45)

To write down the above ODE system explicitly it requires elementary but involved algebraic computation (see (71), (72), (73), (74) in [24]). A more compact way of writing (1.42), (1.43), (1.44), (1.45) is

$$(b_s, \lambda_s, x_s, \gamma_s) = F(b_s, \lambda_s, x_s, \gamma_s, \epsilon_1, \epsilon_2), \tag{1.46}$$

which justifies the name of ODE system. We call (1.42), (1.43), (1.44), (1.45) modulational ODE.

Now assume that all conclusions in Lemma 1.3 hold. Then by applying them into the modulational ODE, one obtains the so-called modulational estimates in the following lemma.

Lemma 1.4 (Lemma 5 in [20]). Let the assumption of Lemma 1.3 hold, and let (1.42), (1.43), (1.44), (1.45) hold, then

$$\left|\frac{\lambda_s}{\lambda} + b\right| + \left|b_s\right| \le C\left(\int |\nabla\epsilon|^2 + \int |\epsilon|^2 e^{-|y|}\right) + \Gamma_b^{1-C\eta} + C\lambda^2 |E_0|, \tag{1.47}$$

$$\begin{aligned} |\tilde{\gamma}_{s} - \frac{1}{\|\Lambda Q\|_{2}^{2}} (\epsilon_{1}, L_{+}\Lambda^{2}Q)| + |\frac{x_{s}}{\lambda}| \leq \delta(\alpha) (\int |\nabla \epsilon|^{2} e^{-2(1-\eta)\frac{\theta(b|y|)}{|b|}} + \int |\epsilon|^{2} e^{-|y|})^{\frac{1}{2}} \\ + C \int |\nabla \epsilon|^{2} + \Gamma_{b}^{1-C\eta} + C\lambda^{2} |E_{0}|. \end{aligned}$$
(1.48)

The Γ_b term will naturally appear in the definition of the linear radiation term $\tilde{\zeta_b}$, which we will discuss later. However, most of the time one only needs to know that

$$e^{-(1+C\eta)\frac{\pi}{|b|}} \le \Gamma_b \le e^{-(1-C\eta)\frac{\pi}{|b|}}.$$
(1.49)

We point out that (1.49) is (2.17) in [20].

By applying the conservation laws (Energy and Momentum), one can obtain two more crucial estimates in the following lemma.

Lemma 1.5 (Lemma 5 in [20]). The following two estimates hold:

$$|2(\epsilon_1, \Sigma) + 2(\epsilon_2, \Theta)| \le C(\int |\nabla \epsilon|^2 + \int |\epsilon|^2 e^{-|y|}) + \Gamma_b^{1-C\eta} + C\lambda^2 |E_0|,$$
(1.50)

$$|(\epsilon_2, \nabla \Sigma)| \le C\delta(\alpha) \left(\int |\nabla \epsilon|^2 + \int |\epsilon|^2 e^{-|y|}\right)^{\frac{1}{2}}.$$
(1.51)

To derive the blow up rate for the blow up solution *u*, one performs the following three steps:

1. Based on (1.47), explore the fact that $b \sim -\frac{\lambda_s}{\lambda}$, and then transfer the evolution of λ (which determines the blow up rate) to the evolution of *b*.

⁵ Note
$$\frac{d}{ds} = \lambda^2 \frac{d}{dt}$$

⁴ The evolution of γ is of course equivalent to the evolution of $\tilde{\gamma}$.

- 2. Obtain a lower bound for b_s , which gives the upper bound for the blow up rate, [16].
- 3. Obtain an upper bound for b_s , which gives the lower bound for the blow up rate, [20].

The lower bound of b_s is given by the following lemma:

Lemma 1.6 (*Proposition 2 in [20]*). Let the results of Lemma 1.3 hold, let (1.42), (1.43), (1.44), (1.45) hold, let (1.47), (1.48), (1.50), (1.51) hold, then one has the estimate

$$b_s \ge \delta_0 \left(\int |\nabla \epsilon|^2 + \int |\epsilon|^2 e^{-|y|} \right) - C\lambda^2 E_0 - \Gamma_b^{1-C\eta}.$$

$$(1.52)$$

The inequality (1.52) is called local virial estimate, it is one of the key estimates in the work of Merle and Raphaël. The lower bound of b_s involves the construction of a certain Lyapounouv functional. For this construction one needs to introduce a certain tail term ζ_b or more precisely its cut-off version $\tilde{\zeta}_b$, [17,20].

Let us first quickly describe ζ_b , one may refer to Lemma 2 in [20] for more details. Let b, η , R_b , R_b^- , ϕ_b be as in (1.26). Let ζ_b be the unique radial solution to

$$\begin{cases} \Delta \zeta_b - \zeta_b + ib\Lambda \zeta_b = \Psi_b, \\ \int |\nabla \zeta_b|^2 < \infty. \end{cases}$$
(1.53)

Here Ψ_b is defined in (1.28). Note, as mentioned previously, that the Γ_b term appears naturally when one constructs $\tilde{\zeta}_b$, see (2.17) in [20].

What Merle and Raphaël actually use in [20] is the cut-off version⁶ of ζ_b , that here we denote by $\tilde{\zeta}_b$. See their discussion before formula (3.4) in [20]. Since later we will use $\tilde{\zeta}_b$ in several places, we write here the precise definition. Let $A = A(b) = e^{a\frac{\pi}{|b|}}$, where *a* is some universal small constant, we let χ_A be a smooth cut-off function that vanishes outside $|x| \ge A$. Then one defines $\tilde{\zeta}_b := \zeta_b \chi_A$. Clearly $\tilde{\zeta}_b$ is supported in $\{|y| \le A\}$. In the rest of this paper, the notation *A* means $A(b) = e^{a\frac{\pi}{|b|}}$ and note that $A \gg \frac{1}{|b|}$.

The tail term $\tilde{\zeta}_b$ is introduced to improve the local virial inequality (1.52). Essentially, one wants to change the term $-\Gamma_b^{1-C\eta}$ in (1.52) to $c\Gamma_b$. Let

$$f_1(s) := \frac{b}{4} \| y \tilde{\mathcal{Q}}_b \|_2^2 + \frac{1}{2} \Im(\int y \nabla \tilde{\zeta} \tilde{\tilde{\zeta}}) + (\epsilon_2, \Lambda(\tilde{\zeta}_{re})) - (\epsilon_1, \Lambda \tilde{\zeta}_{im}).$$
(1.54)

Then (this is highly nontrivial, and is one of the key point in [20])

$$\{f_{1}(s)\}_{s} \geq \delta_{1}(\int |\nabla \tilde{\epsilon}|^{2} + \int |\tilde{\epsilon}|^{2} e^{-|y|} + C\Gamma_{b} - C\lambda^{2}E_{0} - \frac{1}{\delta_{1}} \int_{A \leq |x| \leq 2A} |\epsilon|^{2}.$$
(1.55)

Here $\tilde{\epsilon} = \epsilon - \tilde{\zeta}_b$.

With this, one constructs the Lyapounov functional \mathcal{J} , [20] that we write explicitly as

$$\mathcal{J}(s) = \left(\int |\tilde{Q}_b|^2 - \int |Q|^2\right) + 2(\epsilon_1, \Sigma) + 2(\epsilon_2, \Theta) + \int (1 - \phi_A) |\epsilon|^2 - \frac{\delta_1}{800} \left(b\tilde{f}_1(b) - \int_0^b \tilde{f}_1(v)dv + b\{(\epsilon_2, (\tilde{\zeta}_{re})_1) - (\epsilon_1, (\tilde{\zeta}_{im})_1)\}\right),$$
(1.56)

where

$$\tilde{f}_{1}(b) = \frac{b}{4} \| y \tilde{Q}_{b} \|_{2}^{2} + \frac{1}{2} \Im(\int y \nabla \tilde{\zeta} \tilde{\zeta}).$$
(1.57)

One has the following inequality:

⁶ The main reason for the introduction of this cut-off is that ζ_b itself is not in L^2 .

$$\frac{d}{ds}J(s) \le -Cb\left(\Gamma_b + \int |\nabla\tilde{\epsilon}|^2 + \int |\tilde{\epsilon}|^2 e^{-|y|} + \int_A^{2A} |\epsilon|^2 - \lambda^2 E_0\right) + C\frac{\lambda^2}{|b|^2}E_0.$$
(1.58)

Inequality (1.58) finally leads to the lower bound of b_s .

In [23] and other related works, one can see that the analysis of the log–log blow up solutions can be decomposed into two stages:

- 1. At a certain time t_0 , the initial data evolves into some well prepared data.
- 2. The well-prepared data admits a suitable bootstrap structure, and analysis can be significantly simplified.

One can show for the solution u(t) considered by Merle and Raphaël, that there exists some t_0 such that $u(t_0)$ satisfies the following:

$$u(t_0, x) = \frac{1}{\lambda_0^{\frac{d}{2}}} (\tilde{Q}_{b_0} + \epsilon_0) (\frac{x - x_0}{\lambda_0}) e^{-i\gamma_0}.$$
(1.59)

Also $u(t_0, x)$ satisfies the following:

- 1. orthogonality conditions: (1.35), (1.36), (1.37), (1.38),
- 2. the sign condition of *b*:

 $b_0 := b(t_0) > 0, \tag{1.60}$

3. closeness to Q

$$\|\epsilon_0\|_{H^1} + b_0 < \alpha, \tag{1.61}$$

4. smallness condition of the error ϵ_0 :

$$\int |\nabla \epsilon_0|^2 + |\epsilon_0|^2 e^{-|y|} < \Gamma_{b_0}^{\frac{6}{7}},\tag{1.62}$$

5. renormalized energy/momentum control⁷:

$$\lambda_0^2 |E_0| + \lambda_0 |P_0| < \Gamma_{b_0}^{100}, \tag{1.63}$$

6. log-log regime

$$e^{-e^{\frac{2\pi}{b_0}}} < \lambda_0 < e^{-e^{\frac{2\pi}{b_0}}}.$$
(1.64)

Without loss of generality (by translation in time), we can assume $t_0 = 0$. Now let us focus on the initial data of the form $u(t_0)$, which from now on we denote with u_0 . It turns out that the evolution of the data after $t_0(=0)$ is described by the following lemma.

Lemma 1.7. Assume u solves (1.1) with initial data u_0 as (1.59). For all $T < T^+(u)$, the following bootstrap argument holds:

Let the rescaled time s be defined as $s = \int_0^t \frac{d\tau}{\lambda^2(\tau)} + s_0$, $s_0 = e^{\frac{3\pi}{4b_0}}$, if one assumes the bootstrap hypothesis for $t \in [0, T]$,

$$b(t) > 0, \quad b(t) + \|\epsilon(t)\|_{H^{1}} \le 10\alpha^{1/2}, \quad e^{-e^{\frac{10\pi}{b(t)}}} \le \lambda(t) \le e^{-e^{\frac{\pi}{10b(t)}}},$$

$$\frac{\pi}{10\ln s} \le b(t(s)) \le \frac{10\pi}{\ln s}, \quad \int \|\nabla\epsilon(t)\|^{2} + |\epsilon(t)|^{2}e^{-|y|} \le \Gamma_{b(t)}^{\frac{3}{4}},$$

$$\lambda(t_{2}) \le 3\lambda(t_{1}), \quad \forall T > t_{2} > t_{1} > 0 \ (almost \ monotonicity),$$

$$|x(t)| \le 1/1000,$$
(1.65)

 $^{^{7}}$ In this case, this condition is actually implied by the log-log regime condition below, we still keep it to make the notation consistent.

then one has the bootstrap estimate for $t \in [0, T]$:

$$b(t) > 0, \quad b(t) + \|\epsilon(t)\|_{H^{1}} \le 5\alpha^{1/2}, \quad e^{-e^{\frac{3\pi}{b(t)}}} \le \lambda(t) \le e^{-e^{\frac{\pi}{3b(t)}}},$$

$$\frac{\pi}{5\ln s} \le b(t(s)) \le \frac{5\pi}{\ln s}, \quad \int \|\nabla\epsilon(t)\|^{2} + |\epsilon(t)|^{2}e^{-|y|} \le \Gamma_{b(t)}^{\frac{4}{5}},$$

$$\lambda(t_{2}) \le 2\lambda(t_{1}), \quad \forall T > t_{2} > t_{1} > 0 \ (almost \ monotonicity),$$

$$|x(t)| \le 1/2000.$$
(1.66)

Remark 1.15. The special numbers 5, 4/5... appearing above, are of course only for technical reason, by sharpening the initial conditions at t_0 , one can push the bootstrap estimates to

$$\int |\nabla \epsilon|^2 + |\epsilon|^2 e^{-|y|} \le \Gamma_b^{1-c_1}, \quad e^{-e^{\frac{(1+c_1)\pi}{b}}} \le \lambda(t) \le e^{-e^{\frac{(1-c_1)\pi}{b}}}$$

for arbitrary $c_1 > 0$.

We refer to [25,26,23] for a proof. See in particular Proposition 1 in [26]. Under the bootstrap regime, the analysis is made easier since one can simplify (1.47), (1.48), (1.50), (1.51), (1.52), (1.58) following the observation below:

For any polynomial P, $b >> P(\Gamma_b), \Gamma_b >> P(\lambda)$. (1.67)

Now, the first step of the analysis leading to the log–log blow up solution listed above, that is $b \sim -\frac{\lambda_s}{\lambda}$, is quite clear since by (1.47), one has:

$$|\frac{\lambda_s}{\lambda} + b| \lesssim \Gamma_b^{\frac{1}{10}}.$$
(1.68)

The local virial inequality (1.52), which is used to show the lower bound of b_s , is simplified further to

$$b_s \ge -\Gamma_b^{1-C\eta},\tag{1.69}$$

and the Lyapounov functional (1.58), which is used to show the upper bound of b_s , is simplified to

$$\frac{d}{ds}\mathcal{J} \le -Cb\Gamma_b. \tag{1.70}$$

Basically, all those terms involving $\lambda^2 E_0$ can be neglected since λ is much smaller than *b*. This is actually one of the key observation in [23,3].

Recall from the geometric decomposition (1.34)

$$u(t,x) = \frac{1}{\lambda(t)^{\frac{d}{2}}} (\tilde{Q}_b + \epsilon) \left(\frac{x - x(t)}{\lambda(t)}\right) e^{-i\gamma(t)}$$

Since \tilde{Q}_b is supported in $|y| \lesssim \frac{2}{|b|}$, and the tail radiation term $\tilde{\zeta}_b$ is supported in $|y| \lesssim A \equiv \Gamma_b^{-a}$, then the analysis in work of Merle and Raphaël is taking place in $|x - x(t)| \lesssim \lambda(t)\Gamma_b^{-a} \ll 1$. The only part which connects the local dynamics and the information outside the potential blow up point x(t) is given by the conservation laws (energy and momentum). Thus, if one considers the local momentum and local energy, one could localize the dynamics, paying the prize that the local energy and momentum are no longer conserved. However, since in the analysis any term appearing with the energy or the momentum is multiplied by a power of λ , and λ is so small as mentioned above, ultimately we do not really need that the energy and momentum are conserved, but rather that they are bounded. This is the reason we need to investigate whether the log–log blow up regime may keep the solution u(t) very smooth away from the blow up point. This would of course imply that the local momentum and local energy are varying in a bounded (not necessarily small) way, thus finally totally localizing the dynamics and decoupling what happens in the region near the blow up point and away from it.

1.4. Strategy and structure of the paper

All arguments and results are valid for both dimensions d = 1 and 2. For simplicity, from now on, we work on d = 2. The main result in Theorem 1.1 is proven by describing the dynamics for well-prepared initial data of the form

$$u_0(x) = \sum_{j=1}^m \frac{1}{\lambda_{j,0}} \tilde{\mathcal{Q}}_{b_{j,0}}(\frac{x - x_{j,0}}{\lambda_{j,0}}) e^{-\gamma_{j,0}} + \Xi.$$
(1.71)

We call this data a "multi-solitions model". We sometimes call each "soliton" as bubble. We show that, under certain conditions, *m* bubbles $\frac{1}{\lambda_{j,0}}\tilde{Q}_{b_{j,0}}(\frac{x-x_{j,0}}{\lambda_{j,0}}) + \Xi$, j = 1, ..., m evolve as if they do not feel the existence of each other, then we use a topological argument to show the existence of a blow up solution which blows up at *m* prescribed points according to log-log law.

As mentioned previously, one key element in the proof is to show that solutions generated by the well prepared data in (1.71) will keep high regularity outside a small neighborhood of $x_{j,0}$, j = 1, ..., m.

To make the proof more accessible, we will first discuss an easier model (we call it "one soliton model"), i.e. we show that $\frac{1}{\lambda_0}\tilde{Q}_{b_0}(\frac{x}{\lambda_0}) + \Xi$ blows up according to the log-log law without propagating singularity outside a small neighborhood of the origin, which is essentially a restatement of Theorem 1.2.

We organize the rest of this paper as follows:

In Section 2, we introduce the well prepared data for the "one soliton model" and describe its dynamics in Lemma 2.2. In Section 3, we introduce the well prepared data for the "multi-solitons model" and describe its dynamic in Lemma 3.1. In Section 4, we prove Lemma 2.2 for the "one soliton model". In Section 5, we prove Lemma 3.1 for the "multi-solitons model". In Section 6, we use topological argument to prove the main Theorem 1.1.

2. Description of initial data and dynamic/modification of system: one soliton model

Let us recall that throughout the paper α will be used to denote a universal small number, though its exact value will be chosen at the very end of work. Also $\delta(\alpha)$ is used to denote small constant which depends on α and $\lim_{\alpha \to 0} \delta(\alpha) = 0.$

2.1. Description of initial data

We start with the "one soliton model". We define initial data u_0 in the following form:

$$u_0 = \frac{1}{\lambda_0} (\tilde{Q}_{b_0} + \epsilon_0) (\frac{x}{\lambda_0}), \tag{2.1}$$

which satisfies all the property of data described in (1.59), i.e. orthogonality conditions (1.35), (1.36), (1.37), (1.38), and the bounds (1.60), (1.61), (1.62), (1.63), (1.64). Moreover we assume that outside the origin the data is smooth in the sense that for some $K_2 > 1$, we have

$$\|u_0\|_{H^{K_2}(|x|\ge \frac{1}{3})} \lesssim 1,$$
(2.2)

and u_0 in H^{K_2} .

Remark 2.1. For non-integer values K_2 , formula (2.2) means that, there exists a smooth cut-off function χ equals 1 in $|x| \ge \frac{1}{3}$ such that $\|\chi u_0\|_{H^{K_2}} \le 1$.

Now let us restate Theorem 1.2.

Lemma 2.1. Let u solve (1.1), with initial data prepared as in (2.1). Then

$$\forall K_1 < \frac{K_2}{2}, \quad \|u(t)\|_{H^{K_1}(|x| \ge \frac{1}{2})} \lesssim 1$$

within the lifespan of u.

Remark 2.2. One should understand Theorem 1.2 and Lemma 2.1 in the following way: if the initial data is smooth, then partial smoothness will be kept during the log–log blow up process.

2.2. Modification of system

Let us point out here that Lemma 1.7 cannot directly be applied to u that solves (1.1) with the prepared data u_0 . In fact we cannot even say that u satisfies the geometric decomposition in (1.34) for $t \neq 0$.

Recall that previously a geometric decomposition was obtained via a variational argument, Lemma 1.1 (and a further modulation argument), all relying on a negative energy condition and on (1.10). In our case, for data described as (2.1) we don't even know if they have the negative energy condition.

This does not really matter in the "one soliton model", since one may modify Lemma 1.21 relying on the so-called orbital stability of the soliton and re-establish Lemma 1.3 without negative energy condition.

We do not use this approach here, since we will later deal with "multi-solitons model" and general orbital stability for multi-solitons is a very hard open problem.

Let us consider now a system for $\{u(t), b(t), \lambda(t), x(t), \gamma(t)\}$. We define $\epsilon(t) = \epsilon_1 + i\epsilon_2 := u(t) - \frac{1}{\lambda(t)}\tilde{Q}_{b(t)}(\frac{x-x(t)}{\lambda(t)})e^{i-\gamma(t)}$, and we consider the system as

$$\begin{cases}
iu_{t} = -\Delta u - |u|^{2}u, \\
\frac{d}{dt}\{(\epsilon_{1}(t), |y|^{2}\Sigma_{b(t)}) + (\epsilon_{2}(t), |y|^{2}\Theta_{b(t)})\} = 0, \\
\frac{d}{dt}\{(\epsilon_{1}(t), y\Sigma_{b(t)}) + (\epsilon_{2}(t), y\Theta_{b(t)})\} = 0, \\
\frac{d}{dt}\{-(\epsilon_{1}(t), \Lambda^{2}\Theta_{b(t)}) + (\epsilon_{2}(t), \Lambda^{2}\Sigma_{b(t)})\} = 0, \\
\frac{d}{dt}\{-(\epsilon_{1}(t), \Lambda\Theta_{b(t)}) + (\epsilon_{2}(t), \Lambda\Sigma_{b(t)})\} = 0, \\
\frac{d}{dt}\{-(\epsilon_{1}(t), \lambda(t), x(t), y(t))\}_{t=0} = \{u_{0}, b_{0}, \lambda_{0}, 0, 0\}.
\end{cases}$$
(2.3)

The local well posedness of (2.3) is straightforward, see Appendix A for more details. Note now the following geometric decomposition automatically holds due to the definition of ϵ .

$$u(t,x) = \frac{1}{\lambda(t)} (\tilde{Q}_b + \epsilon) \left(\frac{x - x(t)}{\lambda(t)} \right) e^{-i\gamma(t)}.$$
(2.4)

As pointed out by Merle and Raphaël, such system should be studied in rescaled time variable rather than in its original time variable. So we define the re-scaled time *s* as $\frac{ds}{dt} = \frac{1}{\lambda^2(t)}$, and if one rewrites (2.3) in rescaled time variable, ϵ automatically solves (1.40), (1.41). Since the orthogonality conditions (1.35), (1.36), (1.37), (1.38) hold at t = 0, by (2.3), they hold for all $t \in [0, T]$, the life span of (2.3).

What's more, the life span of (2.3) is exactly the lifespan of the problem (1.1) with initial data u_0 . This is not a trivial fact. Indeed (2.3) is an NLS coupled with four ODEs involving $\{b, \lambda, x, \gamma\}$ and there is the possibility that b may become large, and then this system no longer makes sense since \tilde{Q}_b is only defined for small b. However, as long as (2.3) holds, then the bootstrap Lemma 1.7 works, and this will ensure that b stays bounded and small, and thus the lifespan of (2.3) coincides with that of the NLS problem with initial data u_0 . There is of course another possibility, that the coupled ODE breakdown, i.e. λ becomes 0 and this of course, means that the NLS equation blows up.

To summarize, during the study of the NLS (1.1) with initial data u_0 as in (2.1), for any [0, T] in the lifespan of u(t), u(t) satisfies geometric decomposition:

$$u(t) = \frac{1}{\lambda(t)} (\tilde{Q}_{b(t)+\epsilon(t)}) \left(\frac{x-x(t)}{\lambda(t)}\right) e^{-i\gamma(t)},$$

and orthogonality conditions (1.35), (1.36), (1.37), (1.38). In particular the bootstrap Lemma 1.7 works for these kinds of data.

2.3. Description of the dynamic

Now, we can (equivalently) restate Lemma 2.1 in the following way.

Lemma 2.2. Consider the system (2.3) with initial data $\{u_0, b_0, \lambda_0, 0, 0\}$ described in subsection 2.1, then for any $T < T^+(u_0)$, we have the following bootstrap argument:

Let the rescaled time s be defined as $s = \int_0^t \frac{d\tau}{\lambda^2(\tau)} + s_0$, $s_0 = e^{\frac{3\pi}{4b_0}}$, if one assumes bootstrap hypothesis (1.65) for $t \in [0, T]$, then one has bootstrap estimate (1.66) for $t \in [0, T]$, and following regularity estimate holds for any fixed $K_1 < \frac{K_2}{2}$,

$$\|u(t)\|_{H^{K_1}(|x|\geq \frac{1}{2})} \lesssim 1, t \in [0, T].$$
(2.5)

3. Description of initial data and dynamic: multi-soliton model

Now we turn to "multi-solitons model".

We introduce some notations. Let $\chi_{0,loc}$, $\chi_{1,loc}$ be two smooth cut-off functions such that

$$\chi_{0,loc}(x) = \begin{cases} 1, |x| \le 3/4, \\ 0, |x| \ge 1, \end{cases}$$
(3.1)

$$\chi_{1,loc}(x) = \begin{cases} 1, |x| \le 2/3, \\ 0, |x| \ge 3/4. \end{cases}$$
(3.2)

Let's define the local version of energy and momentum. For any $f \in H^1(\mathbb{R}^2)$, $x_0 \in \mathbb{R}^2$, let

$$E_{loc}(f, x_0) := \int_{\mathbb{R}^2} \chi_{0, loc}(x - x_0) \left(\frac{1}{2} |\nabla f|^2 - \frac{1}{4} |f|^4 \right),$$
(3.3)

$$P_{loc}(f, x_0) := \int_{\mathbb{R}^2} \nabla f \, \bar{f} \, \chi_{0, loc}(x - x_0). \tag{3.4}$$

3.1. Description of the initial data

We define an initial data u_0 of the following form:

$$u_0 := \sum_{j=1}^m u_{j,0} + \Xi_0, \tag{3.5}$$

where

$$u_{j,0} = \frac{1}{\lambda_{j,0}} \tilde{Q}_{b_{j,0}} \left(\frac{x - x_{j,0}}{\lambda_{j,0}} \right) e^{-i\gamma_{j,0}}, j = 1, \dots, m$$

$$\Xi_0 \in H^1(\mathbb{R}^2),$$

$$(3.6)$$

and the properties of Ξ_0 are implicitly encoded below. We let

$$\epsilon_0^j := \lambda_{j,0} \Xi_0(\lambda_{j,0}x + x_{j,0}) e^{i\gamma_{j,0}},$$

$$\epsilon_0^j \equiv \epsilon_{1,0}^j + i\epsilon_{2,0}^j, j = 1, \dots, m.$$
(3.8)

We require the orthogonality condition:

$$(\epsilon_{1,0}^{j}, |y|^{2} \Sigma_{b_{j,0}}) + (\epsilon_{2,0}^{j}, |y|^{2} \Theta_{b_{j,0}}) = 0, j = 1, \dots, m$$
(3.9)

$$(\epsilon_{1,0}^{J}, y\Sigma_{b_{1,0}}) + (\epsilon_{2,0}^{J}, y\Theta_{b_{j,0}}) = 0, j = 1, \dots, m$$
(3.10)

$$-(\epsilon_{1,0}, \Lambda^2 \Theta_{b_{j,0}}) + (\epsilon_{j,0}^2, \Lambda^2 \Sigma_{b_{j,0}}) = 0, \ j = 1, \dots, m$$
(3.11)

$$-(\epsilon_{1,0}^{j}, \Lambda \Theta_{b_{j,0}}) + (\epsilon_{2,0}^{j}, \Lambda \Sigma_{b_{j,0}}) = 0, \, j = 1, \dots, m,$$
(3.12)

where $\tilde{Q}_b \equiv \Sigma_b + i \Theta_b$.

We further require the following:

1. Non-interaction of any two different solitons:

$$|x_{j,0} - x_{j',0}| \ge 10, \forall j \ne j', \tag{3.13}$$

2. Sign condition and smallness condition of $b_{i,0}$:

$$\alpha > b_{j,0} > 0, \, j = 1, \dots, m. \tag{3.14}$$

3. Smallness condition:

$$\sum_{j=1}^{m} \lambda_{j,0} + \sum_{j=1}^{m} \|\epsilon^{j}\|_{H^{1}} \le \alpha.$$
(3.15)

4. Log-log regime condition:

$$e^{-e^{\frac{2\pi}{b_{j,0}}}} < \lambda_{j,0} < e^{-e^{\frac{\pi}{2b_{j,0}}}}, j = 1, \dots, m.$$
 (3.16)

5. Smallness of the local error

$$\int_{|x| \le \frac{10}{\lambda_{j,0}}} |\nabla \epsilon_0^j|^2 + |\epsilon_0^j|^2 e^{-|y|} \le \Gamma_{b_{j,0}}^{\frac{6}{7}}, \, j = 1, \dots, m.$$
(3.17)

6. Tameness of local energy and local momentum

$$\lambda_j^2 |E_{loc}(u_0, x_{j,0})| \le \Gamma_{b_j}^{1000}, \, j = 1, \dots, m.$$
(3.18)

$$\lambda_j |E_{loc}(u_0, x_{j,0})| \le \Gamma_{b_j}^{1000}, \, j = 1, \dots, m.$$
(3.19)

7. Smoothness outside the singularity, here we fix a large N_2 , which would be chosen later.

$$|u_0|_{H^{N_2}(\min_{1\le j\le m}\{|x-x_{j,0}|\}\ge \frac{1}{3})} \le \alpha.$$
(3.20)

8. u_0 in H^{N_2} .

Remark 3.1. Such data can be constructed similarly as those constructed in [26].

3.2. Modification of system

Now let *u* be the solution to (1.1) with initial data u_0 described as in Subsection 3.1. We are expecting that throughout the lifespan of the evolution, we can find parameters $\{b_j(t), \lambda_j(t), x_j(t), \gamma_j(t)\}_{j=1}^m$ and a function $\Xi(t, x)$ such that the following geometric decomposition holds:

$$u(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_j(t)} \tilde{Q}_{b_j}(\frac{x - x_j(t)}{\lambda_j(t)}) e^{-i\gamma_j(t)} + \Xi(t,x),$$
(3.21)

$$|b_j(t)| \le \delta(\alpha), \|\epsilon^j\|_{H^1} \le \delta(\alpha), 1 \le j \le m, \epsilon^j = \lambda_j \xi(\lambda_j x + x_j) e^{i\gamma_j},$$
(3.22)

as well as the orthogonality condition:

$$(\epsilon_1^j, |y|^2 \Sigma_{b_j}) + (\epsilon_2^j, |y|^2 \Theta_{b_j}) = 0, \, j = 1, \dots, m,$$
(3.23)

$$(\epsilon_{1}^{j}, y\Sigma_{b_{j}}) + (\epsilon_{2}^{j}, y\Theta_{b_{j}}) = 0, j = 1, \dots, m,$$

$$(\epsilon_{1}^{j}, y\Sigma_{b_{j}}) + (\epsilon_{2}^{j}, y\Theta_{b_{j}}) = 0, j = 1, \dots, m,$$

$$(3.24)$$

$$(\epsilon_{1}^{j}, y\Sigma_{b_{j}}) + (\epsilon_{2}^{j}, y\Theta_{b_{j}}) = 0, j = 1, \dots, m,$$

$$(3.25)$$

$$-(\epsilon_1, \Lambda^2 \Theta_{b_j}) + (\epsilon_2^J, \Lambda^2 \Sigma_{b_j}) = 0, \, j = 1, \dots, m,$$
(3.25)

1445

$$-(\epsilon_1^j, \Lambda \Theta_{b_j}) + (\epsilon_2^j, \Lambda \Sigma_{b_j}) = 0, j = 1, \dots, m.$$
(3.26)

Again $\tilde{Q}_b = \Sigma_b + i\Theta_b$, $\epsilon^j = \epsilon_1^j + i\epsilon_2^j$. Since at this point we do not know that the general multi-solitons orbital stability holds, we consider the system for $\{u(t,x), \{b_j(t), \lambda_j(t), x_j(t), \gamma_j(t)\}_{j=1}^m\}$ as in Subsection 2.2:

$$\begin{cases} iu_t + \Delta u = -|u|^2 u, \\ \frac{d}{dt} \{ (\epsilon_1^j, |y|^2 \Sigma_{b_j}) + (\epsilon_2^j, |y|^2 \Theta_{b_j}) \} = 0, j = 1, \dots, m, \\ \frac{d}{dt} \{ (\epsilon_1^j, y \Sigma_{b_j}) + (\epsilon_2^j, y \Theta_{b_j}) \} = 0, j = 1, \dots, m, \\ \frac{d}{dt} \{ (\epsilon_1, \Lambda^2 \Theta_{b_j}) + (\epsilon_2^j, \Lambda^2 \Sigma_{b_j}) \} = 0, j = 1, \dots, m, \\ \frac{d}{dt} \{ -(\epsilon_1^j, \Lambda \Theta_{b_j}) + (\epsilon_2^j, \Lambda \Sigma_{b_j}) \} = 0, j = 1, \dots, m, \\ u(0, x) = u_0, \lambda_j(0) = \lambda_{j,0}, b_j(0) = b_{j,0}, x_j(0) = x_{j,0}, \gamma_j(0) = \gamma_{j,0}, j = 1, \dots, m, \end{cases}$$
(3.27)

where $\epsilon^j, \epsilon_1^j, \epsilon_2^j$ is defined by Ξ , as previously noted, and Ξ is defined by

$$\Xi(t,x) \equiv u(t,x) - \sum_{j=1}^{m} \frac{1}{\lambda_j} \tilde{Q}_{b_j}(\frac{x-x_j}{\lambda_j}) e^{-i\gamma_j(t)}.$$
(3.28)

By the orthogonality conditions (3.23), (3.24), (3.25), (3.26) of the initial data, one has the orthogonality condition within the lifespan of (3.27):

$$(\epsilon_1^j, |y|^2 \Sigma_{b_j}) + (\epsilon_2^j, |y|^2 \Theta_{b_j}) = 0, \, j = 1, \dots, m,$$
(3.29)

$$(\epsilon_1^j, y\Sigma_{b_1}) + (\epsilon_2^j, y\Theta_{b_j}) = 0, \ j = 1, \dots, m,$$
(3.30)

$$-(\epsilon_1, \Lambda^2 \Theta_{b_j}) + (\epsilon_2^j, \Lambda^2 \Sigma_{b_j}) = 0, j = 1, \dots, m,$$
(3.31)

$$-(\epsilon_1^j, \Lambda \Theta_{b_j}) + (\epsilon_2^j, \Lambda \Sigma_{b_j}) = 0, \ j = 1, \dots, m.$$

$$(3.32)$$

We discussed previously that the lifespan of (2.3) is the same as the lifespan of u. With the same argument one also has that the lifespan of (3.27) is the same as lifespan of u.

We perform one final simplification. Let us define⁸:

$$\Xi^{j}(x) := \Xi(x)\chi_{1,loc}(x-x_{j}), \varepsilon^{j}(t,y) := \lambda_{j}\Xi^{j}(\lambda_{j}y+x_{j})e^{-i\gamma_{j}},$$

$$\varepsilon^{j}(t,y) \equiv \varepsilon_{1}^{j} + i\varepsilon_{2}^{j}, j = 1, \dots, m.$$
(3.33)

We point out that our analysis will be performed under the condition $\lambda_j \ll b_j$ (see Lemma 3.1 below). Now, by definition, one has $\varepsilon^j = \epsilon^j \chi_{1,loc}(\frac{x}{\lambda_j})$, and recall \tilde{Q}_b is supported in $|y| \lesssim \frac{1}{b}$. It is not hard to see the following orthogonality condition for ε^{j} .

$$(\varepsilon_1^j, |y|^2 \Sigma_{b_j}) + (\varepsilon_2^j, |y|^2 \Theta_{b_j}) = 0, \, j = 1, \dots, m,$$
(3.34)

$$(\varepsilon_1^j, y\Sigma_{b_1}) + (\varepsilon_2^j, y\Theta_{b_j}) = 0, \, j = 1, \dots, m$$
(3.35)

$$-(\varepsilon_1, \Lambda^2 \Theta_{b_j}) + (\varepsilon_2^j, \Lambda^2 \Sigma_{b_j}) = 0, \, j = 1, \dots, m,$$
(3.36)

$$-(\varepsilon_1^j, \Lambda \Theta_{b_i}) + (\varepsilon_2^j, \Lambda \Sigma_{b_i}) = 0, j = 1, \dots, m.$$

$$(3.37)$$

Under the hypothesis $\lambda_i \ll b_i$, the system (3.27) is exactly:

⁸ Please note ε and ϵ are different notations.

$$\begin{aligned} & \{ u_{t} + \Delta u = -|u|^{2} u, \\ & \frac{d}{dt} \{ (\varepsilon_{1}^{j}, |y|^{2} \Sigma_{b_{j}}) + (\varepsilon_{2}^{j}, |y|^{2} \Theta_{b_{j}}) \} = 0, j = 1, \dots, m, \\ & \frac{d}{dt} \{ (\varepsilon_{1}^{j}, y \Sigma_{b_{1}}) + (\varepsilon_{2}^{j}, y \Theta_{b_{j}}) \} = 0, j = 1, 2, \dots, m, \\ & \frac{d}{dt} \{ - (\varepsilon_{1}^{j}, \Lambda^{2} \Theta_{b_{j}}) + (\varepsilon_{2}^{j}, \Lambda^{2} \Sigma_{b_{j}}) \} = 0, j = 1, \dots, m, \\ & \frac{d}{dt} \{ - (\varepsilon_{1}^{j}, \Lambda \Theta_{b_{j}}) + (\varepsilon_{2}^{j}, \Lambda \Sigma_{b_{j}}) \} = 0, j = 1, \dots, m, \\ & u(0, x) = u_{0}, \lambda_{j}(0) = \lambda_{j,0}, b_{j}(0) = b_{j,0}, x_{j}(0) = x_{j,0}, \gamma_{j}(0) = \gamma_{j,0}, j = 1, \dots, m. \end{aligned}$$

$$(3.38)$$

3.3. Description of the dynamic

We are now ready to present the main lemma which contains a bootstrap argument used to describe the dynamics of (3.38), where the initial data $\{u_0, \{b_{j,0}, \lambda_{j,0}, x_{j,0}, \gamma_{j,0}\}_{i=1}^m\}$ is described as in Subsection 3.1.

Lemma 3.1. Consider the system (3.38) with initial data described as in Subsection 3.1 (with the universal constant α is small enough), then $\forall T < T^+(u)$, the following bootstrap argument holds (and all the estimates are independent of T):

Let $s_{j,0} := e^{\frac{3\pi}{4b_{j,0}}}$, and define the re-scaled time $s_j := \int_0^t \frac{1}{\lambda_j^2} + s_{j,0}$, j = 1, ..., m. If for $t \in [0, T)$ the bootstrap hypotheses hold:

- smallness of b_i and ε :

$$\sum_{j=1}^{m} |b_j| + \sum_{j=1}^{m} \|\varepsilon^j\|_{H^1} \le 10\alpha^{1/2},$$
(3.39)

- log-log regime, part I:

$$e^{-e^{\frac{10\pi}{b_j(t)}}} \le \lambda_j(t) \le e^{-e^{\frac{10\pi}{10b_j(t)}}}, \ j = 1, \dots, m,$$
(3.40)

- log-log regime, part II:

$$\frac{\pi}{10\ln s_j} \le b_j \le \frac{10\pi}{\ln s_j}, \, j = 1, \dots, m,$$
(3.41)

- smallness of the local error:

$$\int |\nabla \varepsilon^{j}|^{2} + |\varepsilon^{j}|^{2} e^{-|y|} \le \Gamma_{b_{j}}^{\frac{3}{4}}, j = 1, \dots, m,$$
(3.42)

- almost monotonicity:

$$\forall T > t_2 \ge t_1 \ge 0, \, \lambda_j(t_2) \le 3\lambda_j(t_1), \, j = 1, \dots, m,$$
(3.43)

- control of translation parameters:

$$\forall 0 < t < T, |x_j(t) - x_{j,0}| \le \frac{1}{1000}, j = 1, \dots, m,$$
(3.44)

- local control of conserved quantity: for all 0 < t < T, j = 1, ..., m:

$$|E_{loc}(x_j(t), u(t)) - E_{loc}(x_{j,0}, u_0)| \le 1000, \tag{3.45}$$

$$|P_{loc}(x_j(t), u(t)) - P_{loc}(x_{j,0}, u_0)| \le 1000,$$
(3.46)

- outside-smoothness:

$$\|u\|_{H^{N_1}(\min_j\{|x-x_{j,0}| \ge \frac{1}{2}\})} \le \frac{1}{\max_j\{\lambda_{j,0}\}},\tag{3.47}$$

where N_1 is some fixed large constant, $N_1 < \frac{N_2}{2}$.

Then for $t \in [0, T)$ *, the bootstrap estimates hold:*

- smallness of b_i and ϵ :

$$\sum_{j=1}^{m} |b_j| + \sum_{j=1}^{m} \|\epsilon^j\|_{H^1} \le 5\alpha^{1/2},$$
(3.48)

- log-log regime, part I:

$$e^{-e^{\frac{5\pi}{b_j(t)}}} \le \lambda(t) \le e^{-e^{\frac{5\pi}{5b_j(t)}}}, j = 1, \dots, m,$$
(3.49)

- log-log regime, part II:

$$\frac{\pi}{5\ln s_j} \le b_j \le \frac{5\pi}{\ln s_j}, \, j = 1, \dots, m,$$
(3.50)

- smallness of the local error:

$$\int |\nabla \varepsilon^{j}|^{2} + |\varepsilon^{j}|^{2} e^{-|y|} \le \Gamma_{b_{j}}^{\frac{4}{5}}, j = 1, \dots, m,$$
(3.51)

- almost monotonicity:

$$\forall T > t_2 \ge t_1 \ge 0, \lambda_j(t_2) \le 2\lambda_j(t_1), j = 1, \dots, m,$$
(3.52)

- control of translation parameter:

$$\forall 0 < t < T, |x_{j,t} - x_{j,0}| \le \frac{1}{2000}, j = 1, \dots, m,$$
(3.53)

- local control of conserved quantity: for all 0 < t < T, j = 1, ..., m:

$$|E_{loc}(x_j(t), u(t)) - E_{loc}(x_{j,0}, u_0)| \le 500,$$
(3.54)

$$|P_{loc}((x_j(t), u(t)) - P_{loc}(x_{j,0}, u_0)| \le 500,$$
(3.55)

- outside smoothness:

$$\|u\|_{H^{N_1}(\min_j\{|x-x_{j,0}| \ge \frac{1}{2}\})} \lesssim 1.$$
(3.56)

Remark 3.2. According to our notation, (3.56) means that

$$\|u\|_{H^{N_1}(\min_i\{|x-x_{i,0}| \ge \frac{1}{2}\})} \le C$$

for some universal constant C. Also by (3.15) we have $\lim_{\alpha \to 0} \frac{1}{\max_j \{\lambda_{j,0}\}} = \infty$. Thus, when α is small enough, (3.56) is stronger than (3.47), i.e. this is a bootstrap estimate.

Remark 3.3. The Lemma 3.1 itself implies that $T^+(u) \le \delta_1(\alpha)$. Indeed, Lemma 3.1 implies that the solution blows up in finite time and according to the log–log law. Similarly for the solutions in Lemma 1.7 and Lemma 2.2. This is a standard argument, for details see for example [23]. One may also directly look at Subsection 6.2, see Remark 6.3.

Remark 3.4. Note that Lemma 3.1 implies that $\min_{j \neq j'} |x_j(t) - x_{j'}(t)| > 5$ for all $t < T^+(u)$.

3.4. Further remarks on Lemma 3.1

Lemma 3.1 means that all the bootstrap estimates hold within the lifespan of u generated by the initial data described in Subsection 3.1. In particular for initial data of the form

$$u(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_j} \tilde{Q}_{b_{j,0}}(\frac{x - x_{j,0}}{\lambda_j}) e^{-i\gamma_{j,0}} + \Xi_0(t,x)$$

with orthogonality conditions (3.9), (3.10), (3.11), (3.12) and with bounds (3.13), (3.14) ... (3.20), the associated solution is smooth in the area min_j $|x - x_{j,0}| \ge \frac{1}{2}$ (i.e. estimate (3.56) holds).

Furthermore, if one looks⁹ at $\chi_{1,loc}(x - x_j(t))u(t, x)$, j = 1, ..., m, then

$$\chi_{1,loc}(x-x_j(t))u(t,x) := \frac{1}{\lambda_j}(\tilde{\mathcal{Q}}_{b_j} + \varepsilon^j)(\frac{x-x_j}{\lambda_j})e^{-i\gamma_j}, \ j = 1,\dots,m$$
(3.57)

and the bootstrap estimates stated in Lemma 3.1 simply mean that $\chi(x - x_j(t))u(t, x)$ evolves according to the log–log law described in the series of work of Merle and Rapahel (until at least one soliton blows up).

4. Proof of Lemma 2.2: one soliton model

This section is devoted to the proof of Lemma 2.2, and we need only to show (2.5) since we already have (1.7). Fix [0, T], all the estimates below are independent of the choice of T.

4.1. Setting up

Recall that u is the solution to (1.1) with initial data u_0 as described in (2.1). Recall as discussed in Subsection 2.2, that we consider the system (2.3) for $\{u(t), b(t), \lambda(t), x(t), \gamma(t)\}$, and the geometric decomposition (2.4) holds. Also Lemma 1.7 holds. We finally point out that the bootstrap hypothesis (1.65) itself implies

$$\int_{0}^{l} \frac{1}{\lambda(\tau)^{\mu}} d\tau \leq \begin{cases} C(\mu) \text{ for } \mu < 2, \\ \frac{|\ln \lambda(t)|^{101}}{\lambda(t)^{\mu-2}} \text{ for } \mu \ge 2, \end{cases}$$

$$(4.1)$$

for any t < T. See (51) in [26].

Remark 4.1. Note here that we may also use the bootstrap estimate (1.66) to show (4.1), since we know that (1.65) implies (1.66). However, all the arguments in this section only rely on the bootstrap hypothesis (1.65), this will become important when we work on multi-solitons model.

4.2. An overview of the proof

Recall that for all $t \in [0, T^+(u))$, the geometric decomposition holds

$$u(t) = \frac{1}{\lambda(t)} (\tilde{Q}_b + \epsilon) (\frac{x - x(t)}{\lambda(t)}) e^{-i\gamma(t)},$$

and by the bootstrap hypothesis (1.65), we know $|x(t)| \le \frac{1}{1000}$ for $t \in [0, T]$. Now we fix χ_0 , such

$$\chi_0(x) := \begin{cases} 0, |x| \le a_0 \equiv \frac{1}{500}, \\ 1, |x| \ge d_0 \equiv \frac{1}{250}. \end{cases}$$
(4.2)

The key to proving that a log-log blow up solution can propagate some regularity outside its potential blow up point x(t) is the following control:

$$\forall T < T^{+}(u), \int_{0}^{1} \|\nabla(\chi_{0}u(t))\|_{2}^{2} \lesssim 1.$$
(4.3)

This is pointed out in [26], see formula (63) in [19] for a proof.

We first use the I-method to show the rough control,

⁹ We note that by the definition of $\chi_{1,loc}$ and the bootstrap estimate $|x_j(t) - x_{j,0}| \le \frac{1}{1000}$, then $u(t,x) = \chi_{1,loc}(x - x_j(t))u(t,x)$ in the region $|x - x_{j,0}| \le \frac{1}{2}$.

Lemma 4.1. In the setting of Lemma 2.2, the following estimates hold for $t \in [0, T]$, for any $\sigma > 0$,

$$\|u(t)\|_{H^{K_2}} \lesssim_{\sigma} \left(\frac{1}{\lambda(t)}\right)^{K_2 + \sigma}.$$
(4.4)

For a proof see Subsection 4.3.

Next, we introduce a sequence of cut off functions $\{\chi_l\}_{l=0}^L$, where *L* is a large but fixed number which will be chosen later:

$$\chi_l = \begin{cases} 0, |x| \le a_l, \\ 1, |x| \ge b_l. \end{cases}$$
(4.5)

Here we require $a_l < b_l < a_{l-1}, b_L \leq \frac{1}{2}$.

The idea is to retreat from $\chi_{l-1}u$ to $\chi_l \tilde{u}$ for each l, showing that $\chi_l u$ has higher regularity than $\chi_{l-1}u$. For convenience of notation we set $v_i := \chi_i u$. We use the crucial control (4.3), Strichartz estimates and interpolation techniques to show:

Lemma 4.2. *For all* 0 < v < 1 *and for* $t \in [0, T]$ *, we have*

$$\forall \sigma > 0, \|\nabla^{\nu} v_2(t)\|_{L^2} \lesssim_{\sigma} \frac{1}{\lambda(t)^{\sigma}}.$$
(4.6)

For a proof see Subsubsection 4.4.2.

This lemma gives a better L^{∞} estimate of v_2 , which of course implies a better L^{∞} estimate of $v_i, i \ge 2$. We also show:

Lemma 4.3. For any 0 < v < 1 and for $t \in [0, T]$, we have

$$\forall \sigma > 0, \|v_2\|_{L^{\infty}} \lesssim_{\sigma} \frac{1}{\lambda(t)^{\sigma}}.$$
(4.7)

For a proof see Subsubsection 4.4.3.

Remark 4.2. This lemma should be understood as an improvement of the L^{∞} estimate. Indeed, from the viewpoint of Sobolev embedding, $H^{1+}(\mathbb{R}^2) \hookrightarrow L^{\infty}$, thus the trivial L^{∞} estimate is $\|v_2\|_{L^{\infty}} \lesssim \|u(t)\|_{L^{\infty}} \lesssim \frac{1}{\lambda(t)^{1+}}$.

Lemma 4.3 improves Lemma 4.2 immediately:

Lemma 4.4. *For all* 0 < v < 1 *and for* $t \in [0, T]$ *, we have*

$$\|\nabla^{\nu}v_3(t)\|_{L^2} \lesssim 1. \tag{4.8}$$

For a proof see Subsubsection 4.4.4.

Finally, we have the following lemma (which is analogue to Lemma 4 in [26]) to iterate the gain of regularity for large K_2 .

Lemma 4.5. If the following estimate holds for some r > 0 and some $i, 1 \le i \le L - 1$,

$$\|v_i(t)\|_{H^r} \lesssim 1, t \in [0, T], \tag{4.9}$$

then there is a gain of regularity on v_{i+1} , that is

$$\forall \tilde{r} < \frac{2(K_2 - r)}{K_2} - 1 + r, \|v_{i+1}\|_{H^{\tilde{r}}} \lesssim_{\tilde{r}} 1, t \in [0, T].$$
(4.10)

For a proof see Subsubsection 4.4.5.

Now we are ready to end the proof of Lemma 2.2, i.e. to prove (2.5).

Proof of (2.5) **in Lemma 2.2.** Lemma 4.5 is enough for us to obtain the regularity estimate (2.5) for $K_1 < \frac{K_2}{2}$. To see this consider the end point case in (4.10), i.e. if one considers $r = \tilde{r} = \frac{2(K_2 - r)}{K_2} - 1 + r$, then one obtains $r = \frac{K_2}{2}$. Thus, when $K_1 < \frac{K_2}{2}$, by choosing the iteration time *L* large, the desired estimate (2.5) follows. \Box

4.3. Rough control

This subsection is devoted to the proof of Lemma 4.1, similar type of estimates are derived in [26] by considering some pseudo energy. Here, we rely on I-method. Without loss of generality, we only show:

$$\|u(T)\|_{H^{K_2}} \lesssim_{\sigma} \left(\frac{1}{\lambda(T)}\right)^{K_2 + \sigma}.$$
(4.11)

4.3.1. LWP interval

Let's introduce the so-called LWP interval as in [3]. To make the argument easier, we observe that, under the bootstrap hypothesis (1.65), we can show that $\lambda(t)$ is actually strictly decreasing. Indeed, by (1.65), b > 0, $\lambda \ll b$ in the sense of (1.67). Thus by (1.69) and (1.67), we obtain $-\frac{\lambda_s}{\lambda} \ge \frac{1}{2}b > 0$, thus $\lambda(t)$ is strictly decreasing. We define k_0, k_T as:

$$\frac{1}{2^{k_0}} \le \lambda(0) < \frac{1}{2^{k_0 - 1}}, \frac{1}{2^{k_T}} \le \lambda(T) < \frac{1}{2^{k_T - 1}},\tag{4.12}$$

and for $k_0 \le k \le k_T$, let t_k be the (unique) time such that

$$\lambda(t_k) = \frac{1}{2^k}.\tag{4.13}$$

Then as in [3], we can perform a bootstrap argument:

Lemma 4.6. In [0, T], assuming the bootstrap hypothesis

$$t_{k+1} - t_k \le k\lambda^2(t_k),\tag{4.14}$$

we obtain the bootstrap estimate

$$t_{k+1} - t_k \le \sqrt{k}\lambda^2(t_k). \tag{4.15}$$

This estimate is not sharp, indeed, morally one should have $t_{k+1} - t_k \sim (\ln \ln k)\lambda(t_k)^2$ according to the log–log law. We refer to (2.39) in Lemma 2.4 of [3].

As in [3], we divide all intervals $[t_k, t_{k+1}]$ into disjoint intervals $\bigcup_{j=0}^{J_k-1} [\tau_k^j, \tau_k^{j+1}]$, where

$$t_{k} = \tau_{k}^{0} < \tau_{k}^{1} \cdots < \tau_{k}^{J_{k}-1} < \tau_{k}^{J_{k}} = t_{k+1},$$

$$\tau_{k}^{j+1} - \tau_{k}^{j} = \frac{\delta_{1}}{4} \lambda(t_{k+1})^{2}, \forall j \leq J_{k} - 2,$$

$$\tau_{k}^{J_{k}} - \tau_{k}^{J_{k}-1} \leq \frac{\delta_{1}}{4} \lambda(t_{k+1})^{2}, \forall j \leq J_{k} - 2.$$
(4.16)

Here δ_1 is a fixed constant which will be chosen later. Now, by Lemma 4.6, using the bootstrap estimate (4.15) (indeed, we only use the bootstrap hypothesis (4.14), which is weaker), we have

$$J_k \le 10k^2, \sum_k J_k \le 10k_T^3 \lesssim |\ln \lambda(T)|^3.$$
(4.17)

4.3.2. A quick introduction of upside-down I-method

We point out without proof the following classical fact on any LWP interval:

$$\forall [\tau_k^j, \tau_k^{j+1}], \sup_{t \in [\tau_k^j, \tau_k^{j+1}]} \|u(t) - u(\tau_k^j)\|_{H^{\kappa_2}} \lesssim \|u(\tau_k^j)\|_{H^{\kappa_2}},$$
(4.18)

and directly iterate this over all the LWP intervals (recall we have about $|\ln\lambda(T)|^3$ such intervals), to obtain the estimate $||u(T)||_{H^{K_2}} \lesssim C^{|\ln\lambda(T)|^3}$ for some C > 1, which is clearly weaker than (4.11) but actually quite close. The idea of the original I-method [2] or up-side down I-method [27] are both to improve the estimate (4.18) on a

The idea of the original I-method [2] or up-side down I-method [27] are both to improve the estimate (4.18) on a LWP interval by working on certain a slowly varying/almost conserved quantity.

We introduce the upside-down I-operator \mathcal{D}_N :

$$\widehat{\mathcal{D}_N f}(\xi) := \mathcal{M}(\frac{\xi}{N}) \widehat{f}(\xi), \tag{4.19}$$

where $\mathcal{M}(\xi)$ is a smooth function such that

$$\mathcal{M}(\xi) := \begin{cases} (|\xi|)^{K_2}, |\xi| \ge 2, \\ 1, |\xi| \le 1. \end{cases}$$
(4.20)

It is easy to see that for any $f \in H^{K_2}$ we have

$$\|\mathcal{D}_N f\|_{L^2} \lesssim \|f\|_{H^{K_2}} \lesssim N^{K_2} \|\mathcal{D}_N f\|_{l^2}.$$
(4.21)

Now for some v that solves NLS, the idea of the upside-down I-method is to use $\|\mathcal{D}_N v\|_{L^2}^2$ to model $\|v\|_{H^K}^2$, and show that $E^1(v) := \|\mathcal{D}_N v\|_2^2$ is slow varying, see [27].

Lemma 4.7 (Proposition 3.4, Lemma 4.5 in [27]). There exists a higher modified energy E^2 (or to be more precise, E_N^2), such that for any f in H^{K_2} ,

$$|E^{2}(f) - ||\mathcal{D}_{N}f||_{2}^{2}| \lesssim \frac{1}{N^{1-}} ||\mathcal{D}_{N}f||_{2}^{2},$$
(4.22)

 $\forall M > 0$, there exists $\delta_0 = \delta_0(M) > 0$ such that if v solves (1.1) with initial data v_0 and $\|v_0\|_{H^1} \le M$, then $[0, \delta_0]$ is in the lifespan of v and $E^2(v)$ is slow varying in the following sense:

$$\sup_{0 \le \tau \le \delta_0} |E^2(v(\tau)) - E^2(v_0)| \lesssim \frac{1}{N^{\frac{7}{4}-}} E^2(v_0),$$
(4.23)

and

$$\sup_{0 \le \tau \le \delta_0} \|\mathcal{D}_N v(t)\|_{L^2}^2 \le (1 + C\frac{1}{N^{\frac{7}{4}}})(1 + C\frac{1}{N^{1-}})\|\mathcal{D}_N v(t_0)\|_{L^2(\mathbb{R}^2)}^2,$$
(4.24)

where C is some universal constant.

Actually in [27] one finds only (4.22) and (4.23), but (4.24) is directly implied. Formula (4.22) can be found in the last formula of the proof of the Proposition 3.4 in [27].

Remark 4.3. In [27], the equation is defocussing, while here we are working with a focusing equation. However, when one restricts analysis locally, the two problems actually have no real difference. In [27], the defocussing condition is only used to ensure that any H^s , s > 1 solution considered is global. Note [27] deals with both Euclidean case and Torus case.

4.3.3. Proof of Lemma 4.1

Now we are ready to finish the proof of Lemma 4.1.

Proof. Recall we need only to show (4.11). Let

$$N = (\frac{1}{\lambda(T)})^{1+\sigma_1}.$$
(4.25)

Here $\sigma_1 = \sigma_1(\sigma)$ is a small positive constant chosen later. We have the following estimate on LWP interval:

$$\sup_{t \in [\tau_k^j, \tau_k^{j+1}]} \|\mathcal{D}_N u(t)\|_{L^2}^2 \le (1 + C \frac{1}{(N\lambda(T))^{\frac{1}{2}}}) \|\mathcal{D}_N u(\tau_k^j)\|_{L^2}^2.$$
(4.26)

The constant C is independent of j, k. We now assume (4.26) temporarily and we finish the proof of (4.11). By our choice of N, we immediately obtain from (4.26) that

$$\sup_{\in [\tau_k^j, \tau_k^{j+1}]} \|\mathcal{D}_N u(t)\|_{L^2}^2 \le (1 + C \frac{1}{(\lambda(T)^{-\sigma_1})^{\frac{1}{2}}}) \|\mathcal{D}_N u(\tau_k^j)\|_{L^2}^2.$$
(4.27)

Then we iterate this estimate and we recall that the total number of LWP intervals is controlled by (4.17), thus we obtain

$$\|\mathcal{D}_N u(T)\|_2^2 \lesssim (1 + C\lambda(T)^{\frac{\sigma_1}{2}})^{C(|\ln\lambda(T)|)^3} \|u_0\|_{H^{K_2}} \lesssim 1.$$
(4.28)

Thus, we arrive to the estimate

t

$$\|u\|_{H^{K_{2}}} \lesssim N^{K_{2}} \|\mathcal{D}_{N}u\|_{2} \lesssim \frac{1}{\lambda(T)^{K_{2}(1+\sigma_{1})}}.$$
(4.29)

The desired estimate (4.4) clearly follows if we choose $\sigma_1 \le \sigma/K_2$. What is now left is to prove (4.26). We indeed show that

$$\sup_{t \in [\tau_k^j, \tau_k^{j+1}]} \|\mathcal{D}_N u(t)\|_{L^2}^2 \le (1 + C \frac{1}{(N\lambda(t_k^j))^{\frac{1}{2}}} \|\mathcal{D}_N u(\tau_k^j)\|_{L^2}^2.$$
(4.30)

Since, $\lambda(T) \le \lambda(t_{k+1})$, clearly (4.26) follows from (4.30). We now prove (4.30). Indeed, by scaling or direct computation, let

$$u_{j,k}(t,x) := \lambda(t_{k+1})u(\lambda(t_{k+1})^2(\tau_k^j + t), \lambda(t_{k+1}x)),$$

then $u_{j,k}$ solves (1.1) with initial data $\lambda(t_{k+1})u(\tau_k^j, \lambda(t_{k+1})x)$. A direct computation leads to

$$\|\lambda(t_{k+1})u(\tau_k^J,\lambda(t_{k+1})x)\|_{H^1} = \|Q+\epsilon\|_{H^1} \le 2\|Q\|_{H^1}.$$

To apply the upside-down I-method in Lemma 4.7, we first choose δ_0 in Lemma 4.7 as $\delta_0(2\|Q\|_{H^1}^2)$ and then we choose δ_1 in (4.16) as $\delta_1 = \frac{\delta_0}{2}$, and as a consequence for all *j*, *k*,

$$\frac{\tau_k^j - \tau_k^{j+1}}{\lambda(t_{k+1})^2} < \delta_0.$$

By the upside-down I-method (4.24), we have:

$$sup_{t \in [0, \frac{\tau_{k}^{j} - \tau_{k}^{j+1}}{\lambda(t_{k+1})^{2}}]} \|\mathcal{D}_{N\lambda(t_{k})}u_{j,k}(t)\|^{2}$$

$$\leq (1 + C \frac{1}{(N\lambda(t_{k}^{j}))^{\frac{7}{4}-}} (1 + C \frac{1}{(N\lambda(t_{k}^{j}))^{1-}} \|\mathcal{D}_{N\lambda(t_{k}^{j})}u_{j,k}(t)\|_{L^{2}}^{2}$$

$$\leq (1 + C \frac{1}{(N\lambda(t_{k}^{j}))^{\frac{1}{2}}})) \|\mathcal{D}_{N\lambda(t_{k}^{j})}u_{j,k}(t)\|_{L^{2}}^{2}.$$
(4.31)

1452

We now observe that for any $\lambda > 0$ and for any function *f* we have

$$\lambda(\mathcal{D}_N f)(\lambda x) = \mathcal{D}_{N\lambda} \left(\lambda f(\lambda x)\right),$$

$$\|f\|_2 = \|\lambda f(\lambda x)\|_2.$$

(4.32)

(4.30) clearly follows. \Box

Remark 4.4. It is not hard to see that the proof only relies on the fact that one can divide [0, T] into disjoint intervals $\bigcup_{j,k} I_{k,j}$ such that

1. $||u(t)||_{H^1} \sim 2^{-k}$ for $t \in I_{k,j}$ (this is equivalent to $\lambda(t) \sim 2^{-k}$ for $t \in I_{k,j}$);

2. $|I_{k,j}| \sim \frac{1}{2^{2k}}, \forall k, j;$

3. $\sharp \{I_{k,j}\} \lesssim |\ln \lambda(T)|^3$;

4. $\lambda(T) \lesssim \lambda(t), t \in [0, T]$

and $u(0) \in H^{K_2}$.

Note that condition 3 follows from the bootstrap Lemma 4.6.

We finally remark that by further dividing $I_{k,j}$, it is very easy to improve condition 2 to $|I_{k,j}| \le \frac{\delta}{2^{2k}}$ for any fixed $\delta > 0$.

4.4. Propagation of regularity

In this subsection, we give the proof of Lemma 4.2, Lemma 4.3, Lemma 4.4 and Lemma 4.5. Since our proof relies on the Strichartz estimates, for completeness we recall them below.

4.4.1. Strichartz estimates

Consider the equation:

$$\begin{cases} iU_t + \Delta U = F, \\ U(0, x) = U_0. \end{cases}$$

$$\tag{4.33}$$

We write it in the integral form using the Duhamel Formula:

$$U(t) = e^{it\Delta} U_0 - i \int_0^t e^{i(t-s)\Delta} F(s) ds,$$
(4.34)

where $e^{it\Delta}$ is the propagator of linear Schrödinger equation. For notation convenience, let:

$$\Gamma F := -i \int_{0}^{t} e^{i(t-s)\Delta} F(s) ds.$$

Then one has the following Strichartz estimates:

$$\begin{aligned} &\forall \frac{2}{q} + \frac{2}{r} = \frac{2}{2}, 2 < q \le \infty, \\ &\|e^{it\Delta} U_0\|_{L^q([0,\infty],L^r(\mathbb{R}^2))} \lesssim_{q,r} \|U_0\|_2, \\ &\|\Gamma F\|_{L^q([0,t],L^r(\mathbb{R}^2))} \lesssim_{q,r} \|F\|_{L^{q'}([0,t],L^{r'}(\mathbb{R}^2))}, \end{aligned}$$
(4.35)

where (q', r') is the dual of (q, r), i.e. $\frac{1}{q} + \frac{1}{q'} = 1$, $\frac{1}{r} + \frac{1}{r'} = 1$. We call (q, r) an admissible pair if and only if $\frac{2}{q} + \frac{2}{r} = \frac{2}{2}$, $2 < q \le \infty$. We refer to [28,9] and the references within for a proof.

4.4.2. Proof of Lemma 4.2

We fix $\nu < 1$, and we estimate $\|\nabla^{\nu} v_2\|_2$. Note that v_2 satisfies:

$$\left\{ i\partial_t v_2 + \Delta v_2 = \Delta \chi_2 v_1 + 2\nabla \chi_2 \nabla v_1 - v_2 |v_1|^2. \right.$$
(4.36)

As in previous section, we only need to show:

$$\|\nabla^{\nu}v_2(T)\|_2 \lesssim \left(\frac{1}{\lambda(T)}\right)^{\sigma}.$$
(4.37)

We explain some heuristics for this estimate. We view the system as a perturbation of the linear Schrödinger equation, and the dominating term in the perturbation is the last term on the right side of (4.36). We can view (4.36) as:

$$i\partial_t v_2 + \Delta v_2 \approx O(|v_1|^2 v_2). \tag{4.38}$$

From the viewpoint of persistence of regularity one gets:

$$\|\nabla_{\nu}v_{2}(T)\|_{2} \lesssim \|\nabla^{\nu}v_{2}(0)\|_{2}e^{\|v_{1}(t)\|_{L^{4}([0,t],L^{4})}^{*}}.$$
(4.39)

By estimating

$$\|v_1(t)\|_{L^4(\mathbb{R}^2)} \lesssim \|u(t)\|_{H^{\frac{1}{2}}} \lesssim \frac{1}{\lambda^{\frac{1}{2}}(t)},$$

one obtains by (4.1)

$$\|\nabla^{\nu}v_{2}(T)\|_{2} \lesssim e^{\int_{0}^{T} \frac{1}{\lambda^{2}(t)}} \lesssim e^{|\ln\lambda(T)|^{101}}.$$
(4.40)

This estimate though is too weak, of course. On the other hand we didn't even use the key estimate (4.3) in the log–log regime. To obtain the stronger estimate (4.37) instead of the $L_{t,x}^4$ norm in (4.39) we use a more flexible $L_t^q L_x^r$. More precisely we replace $||v_1(t)||_{L^4([0,t],L^4)}^4$ by some $||v_1(t)||_{L^q([0,T],L^r(\mathbb{R}^2))}^q$ such that (q, r) is an admissible pair. Now, we make the key observation that $L^2\dot{H}^1$ and $L_t^4 L_x^4$ have the same scaling as $L_t^q L_x^r$ whenever (q, r) is an admissible pair. Thus, by interpolating the two estimates

$$\int_{0}^{T} \|v_{1}(t)\|_{4}^{4} \lesssim |\ln \lambda(t)|^{101},$$

$$\int_{0}^{T} \|\nabla v_{1}(t)\|_{2}^{2} \lesssim 1,$$
(4.41)

we derive

$$\int_{0}^{1} \|u(t)\|_{L^{r}}^{q} \lesssim (|\ln\lambda(t)|)^{1/100},$$
(4.42)

for an admissible (q, r) carefully picked. Then we obtain

$$\|\nabla^{\nu}v_{2}(T)\|_{2} \lesssim \|\nabla^{\nu}v_{2}(0)\|_{2}e^{\|v_{1}(t)\|_{L^{q}([0,T],L^{r})}^{q}} \lesssim e^{|\ln(\lambda(T)|)\frac{1}{100}},$$
(4.43)

which implies the desired estimate (4.37).

We now go to the details. We need the following technical lemmata.

Lemma 4.8. Let $\frac{2}{q_0} + \frac{2}{r_0} = 1$, let q'_0 , r'_0 be their dual, let p_0 be defined by $\frac{1}{r'_0} = \frac{1}{2} + \frac{1}{p_0}$, let h_0 be defined by $h_0 = 1 - \frac{1}{p_0}$, then we have the following estimate uniform with respect to any time interval I

$$\||v_1|^2 \nabla^{\nu} v_2\|_{L^{r'_0}} \lesssim \|\nabla^{\nu} v_2\|_{L^2} \|v_1\|_{L^{2p_0}}^2, \tag{4.44}$$

$$\||v_{1}|^{2}\nabla^{\nu}v_{2}\|_{L^{q'_{0}}(I;L''_{0})} \lesssim \|\nabla^{\nu}v_{2}\|_{L^{\infty}(I;L^{2})}\|v_{1}\|_{L^{2q'_{0}}(I;L^{2p_{0}})}^{2},$$
(4.45)

$$\|v_1\|_{L^{2p_0}} \lesssim \|v_1\|_{H^{h_0}},\tag{4.46}$$

$$\|v_1\|_{L^{2q'_0}(I;H^{h_0})} \lesssim \|v_1\|_{L^2(I;H^1)}^{\frac{h_0-1/2}{1/2}} \|v_1\|_{L^2(I;H^1)}^{\frac{1-h_0}{1/2}} \|v_1\|_{L^4(I;H^{\frac{1}{2}})}^{\frac{1-h_0}{1/2}}.$$
(4.47)

Moreover we can choose q_0 large enough such that $1 - \frac{1}{10000} < h_0 < 1$, $q'_0 \le 2$.

Lemma 4.9. Let $\frac{2}{q_1} + \frac{2}{r_1} = 1$, let q'_1 , r'_1 be their dual, let g_1 , \tilde{g}_1 be defined by $v - 1 = -\frac{2}{g_1}$, $-v \equiv (1 - v) - 1 = -\frac{2}{\tilde{g}_1}$, let p_1 be defined as $\frac{1}{r'_1} = \frac{1}{2} + \frac{1}{p_1}$, let h_1 be defined as $h_1 - 1 = -\frac{2}{p_1}$ and let w_1 be defined as $\frac{1}{q'_1} = \frac{1}{2} + \frac{1}{w_1}$, then we have the following estimates uniform with respect to any time interval I

$$\||v_1||\nabla^{\nu}v_1|v_2\|_{r'_1} \lesssim \|\nabla^{\nu}v_1\|_{L^{\tilde{g}_1}} \|v_1\|_{L^{p_1}} \|v_2\|_{L^{g_1}}, \tag{4.48}$$

$$\|\nabla^{\nu}v_{1}\|_{L^{\tilde{g}_{1}}} \lesssim \|v_{1}\|_{H^{1}}, \|v_{2}\|_{L^{g_{1}}} \lesssim \|\nabla^{\nu}v_{2}\|_{L^{2}}, \tag{4.49}$$

$$\||v_1||\nabla^{\nu}v_1|v_2\|_{L^{q_1'}(I;L^{r_1'})} \lesssim \|v_2\|_{L^{\infty}(I;H^{\nu})} \|v_1\|_{L^2(I;H^1)} \|v_1\|_{L^{w_1}(I;L^{p_1})},$$
(4.50)

$$\|v_1\|_{L^{p_1}} \lesssim \|v_1\|_{H^{h_1}},\tag{4.51}$$

$$\|v_1\|_{L^{w_1}(I;H^{h_1})} \lesssim \|v_1\|_{L^2(I;H^1)}^{\frac{h_1-l/2}{1/2}} \|v_1\|_{L^2(I;H^1)}^{\frac{1-h_1}{1/2}}$$
(4.52)

Moreover we can choose q_1 large such that $1 - \frac{1}{10000} < h_1 < 1$, $w_1 \le 4$.

The proofs are a direct consequence of Sobolev and Hölder inequalities as well as standard interpolation

$$\|u\|_{H^s} \lesssim \|u\|_{H^{s_1}}^{\frac{s_2-s}{s-s_1}} \|u\|_{H^{s_2}}^{\frac{s-s_1}{s_2-s_1}}, s_1 < s < s_2.$$

We finally remark that all the indices happen to coincide because of scaling reasons. Indeed, we can check that $(2q'_0, 2p_0), (w_1, p_1)$ are both admissible pairs.

Now, we pick (q_0, r_0) , (q_1, r_1) as in Lemma 4.8 and Lemma 4.9, and we let all the other associated indices be determined as in these two lemmas. By estimating $||u(t)||_{H^{\frac{1}{2}}} \lesssim \frac{1}{\lambda(t)^{\frac{1}{2}}}$, we have $||v_1(t)||_{H^{\frac{1}{2}}} \lesssim \frac{1}{\lambda(t)^{\frac{1}{2}}}$. Thus, by (4.1) we have:

have:

$$\int_{0}^{T} \|v_{1}(t)\|_{4}^{4} \lesssim |\ln \lambda(t)|^{101}.$$
(4.53)

Also from (4.3), we have

$$\int_{0}^{1} \|\nabla v_{1}(t)\|_{2}^{2} \lesssim 1.$$
(4.54)

Using (4.53), (4.54) and (4.46), (4.47) in Lemma 4.8 we obtain

$$\|v_1\|_{L^{2q'_0}([0,T];L^{2p_0})}^{2q'_0} \le |\ln\lambda(T)|^{1/2},$$

and using (4.53), (4.54) and (4.51), (4.52) in Lemma 4.9 we obtain

$$\|v_1\|_{L^{w_1}([0,T];L^{p_1})}^{w_1} \le |\ln\lambda(T)|^{1/2}.$$

Thus, we are able to divide [0, T] into J_k disjoint intervals $[\tau_k, \tau_{k+1}], k = 1, \dots, J_k$, such that

$$J_{k} \sim_{\epsilon} |\ln(\lambda(T))|^{\frac{1}{2}},$$

$$\|v_{1}\|_{L^{2q'_{0}}([0,T];L^{2p_{0}})} \leq \epsilon,$$

$$\|v_{1}\|_{L^{w_{1}}([0,T]);L^{p_{1}}} \leq \epsilon,$$
(4.55)

where $\epsilon > 0$ is a fixed small constant, to be chosen later.

Now we use the Duhamel formula (4.34) for (4.36) in $[\tau_k, \tau_{k+1}]$, and we obtain for any $t \in [\tau_k, \tau_{k+1}]$,

$$\|v_{2}(t) - v_{2}(\tau_{k})\|_{\dot{H}^{\nu}}$$

$$\leq C \int_{0}^{T} \|v_{1}\|_{H^{\nu}} + C \int_{0}^{T} \|v(t)\|_{H^{1+\nu}} + \|\int_{\tau_{k}}^{t} e^{i(t-\tau)\Delta} \nabla^{\nu}(|v_{1}|^{2}v_{2})(\tau)d\tau\|_{2}$$

$$\leq C \int_{0}^{T} \|u(t)\|_{H^{1+\nu}} + C \||v_{1}|^{2} \nabla^{\nu}v_{2}||\|_{L^{q'_{0}}([t_{k},\tau];L^{r'_{0}})} + C \|\nabla^{\nu}(|v_{1}|^{2})v_{2}\|_{L^{q'_{1}}([t_{k},\tau]);L^{r'_{1}})},$$
(4.56)

where in the last step we have used Strichartz estimates and Fractional Leibniz rule (see Theorem A.8, [10]). We remark it is not hard to see that we can choose $(q_0, r_0) = (q_1, r_1)$. Now we plug in the estimates in (4.45), (4.49) into (4.56), to obtain

$$\|v_{2}(t) - v_{2}(\tau_{k})\|_{\dot{H}^{\nu}}$$

$$\leq C \int_{0}^{T} \|u(t)\|_{H^{1+\nu}}$$

$$+ C(\|v_{1}\|_{L^{2q'_{0}}([\tau_{k},\tau_{k+1}];L^{2p_{0}})}^{2} + \|v_{1}\|_{L^{2q'_{0}}([\tau_{k},\tau_{k+1}];L^{2p_{0}})}\|v_{1}\|_{L^{2}([0,T];H^{1})} \sup_{[\kappa \in [\tau_{k},t]]} \|v_{2}\|_{H^{\kappa}},$$

$$(4.57)$$

i.e.

$$\|v_{2}(t) - v_{2}(\tau_{k})\|_{H^{\nu}} \le C + C \int_{0}^{T} \|u(t)\|_{H^{1+\nu}} + C\epsilon sup_{t \in [t_{k}, \tau]} \|u(t)\|_{H^{\nu}}.$$
(4.58)

The term $\int_0^T \|u(t)\|_{H^{1+\nu}}$ is not hard to control. Interpolating between $\|u\|_{H^1}$ and $\|u\|_{H^{K^2}}$, we obtain

$$\|u(t)\|_{H^{1+\nu}} \lesssim \|u(t)\|_{H^{1}}^{\frac{K_2-1-\nu}{K_2-1}} \|u(t)\|_{H^{K_2}}^{\frac{\nu}{K_2-1}}.$$
(4.59)

Now we plug in the estimate $||u(t)||_{H^1} \lesssim \frac{1}{\lambda(t)}$ and $||u(t)||_{H^{K_2}} \lesssim_{\sigma} \frac{1}{\lambda(t)^{K_2+\sigma}}$ from Lemma 4.1, and we obtain

$$\|u(t)\|_{H^{1+\nu}} \lesssim_{\sigma} \frac{1}{\lambda^{1+\nu+C\sigma}(t)}.$$
(4.60)

By choosing σ small enough such that $1 + \nu + C\sigma < 2$, and using (4.1), we have

$$\int_{0}^{T} \|u(t)\|_{H^{1+\nu}} \lesssim 1.$$
(4.61)

We plug in (4.61) into (4.56), and we choose ϵ small enough to obtain

$$sup_{t\in[t_k,t_{k+1}]}\|u(t)\|_{H^{\nu}} \le 2\|u(t_k)\|_{H^{\nu}} + C.$$
(4.62)

Iterating this over the $J_k \sim |\ln \lambda(T)|^{\frac{1}{2}}$ intervals, we obtain the estimate

$$\|u(T)\|_{H^{\nu}} \lesssim e^{C|\ln\lambda(T)|^{\frac{3}{4}}},\tag{4.63}$$

which implies the desired estimate (4.37).

4.4.3. Proof of Lemma 4.3

We now prove Lemma 4.3. Indeed, by Lemma 4.1 and Lemma 4.2, we obtain for any $\nu < 1$, $\tilde{\sigma} > 0$ the estimates

$$\|v_2(t)\|_{H^{K_2}} \lesssim_{\tilde{\sigma}} \frac{1}{\lambda^{K^2 + \tilde{\sigma}}(t)},\tag{4.64}$$

$$\|v_2(t)\|_{H^{\nu}} \lesssim_{\tilde{\sigma}} \frac{1}{\lambda^{\tilde{\sigma}}(t)}.$$
(4.65)

Now, by Sobolev embedding and interpolation, we obtain

$$\|v_2(t)\|_{L^{\infty}(\mathbb{R}^2)} \lesssim_{\tilde{\sigma}} \|v_2(t)\|_{H^{1+\tilde{\sigma}}} \lesssim_{\tilde{\sigma}} \frac{1}{\lambda^{C\tilde{\sigma}}(t)} \frac{1}{\lambda^{\frac{K_2-K_2\nu}{K_2-\nu}}(t)}.$$
(4.66)

Since $\tilde{\sigma} > 0$ is arbitrary and we can choose v as close to 1 as we want, this clearly gives Lemma 4.3.

4.4.4. Proof of Lemma 4.4

Now we prove Lemma 4.4. We fix ν and we point out that all the constants in the proof will depend the choice of ν . By choosing σ small enough, Lemma 4.3 gives for some small $c = c_{\nu} > 0$

$$\|v_2(t)\|_{L^{\infty}} \lesssim \frac{1}{\lambda^{1-c\nu}(t)}.$$
 (4.67)

By choosing σ even smaller, by Lemma 4.2, we have

$$\|v_2(t)\|_{H^\nu} \lesssim \frac{1}{\lambda^{c\nu/10}(t)}.$$
(4.68)

Clearly v_3 also satisfies the estimate (4.67), (4.68), and also it satisfies the equation

$$i\partial_t v_3 + \Delta v_3 = \Delta \chi_3 v_2 + 2\nabla \chi_3 \nabla v_2 - v_3 |v_2|^2.$$
(4.69)

By the Duhamel's formula, we obtain

$$v_{3}(t) = e^{it\Delta}(v_{3}(0)) + i \int_{0}^{t} e^{i(t-s)\Delta}(\Delta\chi_{3}v_{2} + 2\nabla\chi_{3}\nabla v_{2} - v_{3}|v_{2}|^{2})(s)ds.$$
(4.70)

Thus,

$$\|v_{3}(t)\|_{H^{\nu}} \lesssim \|v_{3}(0)\|_{H^{\nu}} + \int_{0}^{t} \|v_{2}\|_{H^{\nu}} + \int_{0}^{t} \|v_{2}\|_{H^{1+\nu}} + \int_{0}^{t} \|v_{3}|v_{2}|^{2}\|_{H^{\nu}}.$$
(4.71)

As argued previously in (4.61), $\int_0^t \|v_2\|_{H^{1+\nu}} \lesssim 1$. Thus, to finish the proof of Lemma 4.4, we only need to show

$$\int_{0}^{T} \|v_{3}|v_{2}|^{2}(t)\|_{H^{\nu}} \lesssim 1.$$
(4.72)

Indeed,

$$|v_{3}|v_{2}|^{2}||_{H^{\nu}} \lesssim ||v_{3}||_{H^{\nu}} ||v_{2}||_{L^{\infty}}^{2} + ||v_{3}||_{L^{\infty}} ||v_{2}||_{L^{\infty}} ||v_{2}||_{H^{\nu}}.$$
(4.73)

Now we plug the estimates (4.67) and (4.68) into (4.73) and we recall that v_3 also satisfy estimate (4.67) and (4.68), we have

$$\|v_3|v_2|^2(t)\|_{H^{\nu}} \lesssim \frac{1}{\lambda^{2-c\nu}(t)}.$$
(4.74)

Thus, by (4.1), estimate (4.72) follows.

4.4.5. Proof of Lemma 4.5

Finally we prove Lemma 4.5. This is quite straightforward. First, one can directly check (again) that v_{j+1} satisfies:

$$i\partial_t v_{j+1} + \Delta v_j = \Delta \chi_{j+1} v_j + 2\nabla \chi_{j+1} \nabla v_j - v_{j+1} |v_j|^2.$$
(4.75)

Since we assume that $K_2 > 2$, then by Lemma 4.3, we have (for $j \ge 2$):

$$\|v_j(t)\|_{\infty} \lesssim_{\sigma} \frac{1}{\lambda^{\sigma}(t)}.$$
(4.76)

Clearly the same estimates hold for v_{j+1} . Again by the Duhamel's formula:

$$v_{j+1}(t) = e^{it\Delta}(\chi_{j+1}u_0) + i \int_0^t e^{i(t-\tau)\Delta}(\Delta\chi_{j+1}v_j + 2\nabla\chi_{j+1}\nabla v_j - v_{j+1}|v_j|^2)(\tau)d\tau.$$
(4.77)

We remark here that the main term to control is $2\nabla \chi_{j+1} \nabla v_j$. Directly from (4.77), we obtain:

$$\|v_{j+1}(t)\|_{H^{\tilde{r}}} \lesssim \|\chi_{j}u(0)\|_{H^{\tilde{r}}} + \int_{0}^{t} \|v_{j}(\tau)\|_{H^{\tilde{r}}} + \int_{0}^{t} \|v_{j}(\tau)\|_{H^{\tilde{r}+1}} + \int_{0}^{t} \|v_{j+1}|v_{j}|^{2}(\tau)\|_{H^{\tilde{r}}}.$$
(4.78)

Clearly, to finish the proof we only need to show

$$\int_{0}^{t} \|v_{j}(\tau)\|_{H^{\tilde{r}+1}} \lesssim 1,$$
(4.79)

and

$$\int_{0}^{t} \|v_{j+1}|v_{j}|^{2}\|_{H^{\tilde{r}}} \lesssim 1.$$
(4.80)

We first prove (4.80). By interpolating between the estimates

$$\|v_j\|_{H^r} \lesssim 1$$
 and $\|u(\tau)\|_{H^{K_2}} \lesssim_{\sigma} \frac{1}{\lambda^{K_2+\sigma}(\tau)}$

we obtain

$$\|v_j(\tau)\|_{H^{\tilde{r}}} \lesssim_{\sigma} \frac{1}{\lambda^{1+C\sigma}(\tau)}.$$
(4.81)

Note that the same estimates hold for v_{j+1} . Thus,

$$\|v_{j+1}|v_{j}|^{2}(\tau)\|_{H^{\tilde{r}}} \lesssim \|v_{j}\|_{H^{\tilde{r}}} \|v_{j}\|_{L^{\infty}}^{2} \lesssim_{\sigma} \frac{1}{\lambda(\tau)^{1+C\sigma}}.$$
(4.82)

By choosing σ small enough, such that $1 + C\sigma < 2$, (4.80) clearly follows from (4.1). Finally, we turn to the proof of (4.79). Again by interpolation, we obtain

$$\|v_{j}\|_{H^{\tilde{r}+1}} \lesssim \|v_{j}(\tau)\|_{H^{r}}^{\frac{K_{2}-\tilde{r}-1}{K_{2}-r}} \|v_{j}(\tau)\|_{H^{K^{2}}}^{\frac{\tilde{r}+1-r}{K_{2}}-r} \lesssim_{\sigma} \left(\frac{1}{\lambda(\tau)}\right)^{\frac{K_{2}(\tilde{r}+1-r)}{K_{2}-r}+C\sigma}.$$
(4.83)

The key point is our choice of \tilde{r} that ensures $\frac{K_2(\tilde{r}+1-r)}{K_2-r} < 2$, thus by choosing σ small enough, we have $\frac{K_2(\tilde{r}+1-r)}{K_2-r} + C\sigma < 2$. Finally by (4.1) estimate (4.79) follows.

1458

5. Proof of Lemma 3.1: multi solitons model

This section is devoted to the proof of Lemma 3.1.

The proof is involved, hence we first discuss some heuristics. Lemma 3.1 is the consequence of the following facts.

- 1. From previous work, [20,23], the solution u_j to (1.1) with initial data $\chi_{1,loc}(x-x_{j,0})u_0, j=1,\ldots,m$, blows up according to log-log law.
- 2. Assume some solution v to (1.1) blows up according to log-log law, and assume t_0 is close enough to the blow up time $T^+(v)$. Let F = F(t, x) be some smooth perturbation. Then, the solution \tilde{v} to the following Cauchy problem

$$\begin{cases} i\tilde{v}_t = -\Delta \tilde{v} - |\tilde{v}|^2 \tilde{v} + F, \\ \tilde{v}(0) = v(t_0), \end{cases}$$
(5.1)

still blows up according to the log-log law.

- 3. By our smoothness condition (3.20), we can show that solution to (1.1) with initial data $\chi_{1,loc}(x x_{i,0})u_0$ is smooth in the region $|x - x_j| \ge \frac{1}{2}$ for j = 1, ..., m. 4. For j = 1, ..., m the function $u_j := \chi_1(x - x_j(t))$ satisfies the following equation

$$i\partial_t u_j = -\Delta u_j - |u_j|^2 u_j + F_j, \tag{5.2}$$

where

$$F_{j} = -\nabla \chi_{1,loc}(x - x_{j}) \frac{dx_{j}}{dt} + \Delta \chi_{1,loc}(x - x_{j})u(t, x) + 2\nabla \chi_{1,loc}(x - x_{j})\nabla u + (\chi_{1,loc}^{3} - \chi_{1,loc})(x - x_{j})|u|^{2}u.$$
(5.3)

The idea is the following loop argument, which is false of course, but it can be made rigorous by bootstrap argument. If we assume u_i evolves according to the log-log law for each j, we basically know that u is smooth in the region $\min_{j}\{|x-x_{j}|\} \ge 1/2$. And since $\nabla \chi_{1,loc}(x-x_{j})$ and $\chi_{1,loc}(x-x_{j})$ are supported in $|x-x_{j}| \ge \frac{2}{3}$, this implies that F_i above is smooth. Thus fact 4 and fact 2 imply that u_i actually evolves according to the log-log law for each j. Thus, the assumption, that u_i evolves according to log–log law, is right.

Let us turn to the details and a rigorous mathematical proof.

5.1. Outline of the Proof

Recall that the re-scaled time s_j satisfies $\frac{dt}{ds_j} = \lambda_j^2$ for j = 1, ..., m. The system (3.38) implies the following:

$$\frac{d}{ds_{j}}\{(\varepsilon_{1}^{j},|y|^{2}\Sigma_{b_{j}}) + (\varepsilon_{2}^{j},|y|^{2}\Theta_{b_{j}})\} = 0, j = 1,...,m,
\frac{d}{ds_{j}}\{(\varepsilon_{1}^{j},y\Sigma_{b_{1}}) + (\varepsilon_{2}^{j},y\Theta_{b_{j}})\} = 0, j = 1,...,m,
\frac{d}{ds_{j}}\{-(\varepsilon_{1}^{j},\Lambda^{2}\Theta_{b_{j}}) + (\varepsilon_{2}^{j},\Lambda^{2}\Sigma_{b_{j}})\} = 0, j = 1,...,m,
\frac{d}{ds_{j}}\{-(\varepsilon_{1}^{j},\Lambda\Theta_{b_{j}}) + (\varepsilon_{2}^{j},\Lambda\Sigma_{b_{j}})\} = 0, j = 1,...,m.$$
(5.4)

Now, by pure algebraic computation as in [24], one is able to write down almost the same modulation ODE for $\{b_j, \lambda_j, x_j, \tilde{\gamma}_j\}$ for j = 1, ..., m as in formulas (71), (72), (73), (74) in [24], where $\tilde{\gamma}_j(s_j) := \gamma_j - s_j$. Basically, the only difference is that the $\lambda^2 E$ term in [24] is replaced¹⁰ by $E(\tilde{Q}_{b_i} + \varepsilon^j)$.

¹⁰ Because in [24], the solution u has the form $u := \frac{1}{\lambda(t)} (\tilde{Q}_b + \epsilon) (\frac{x - x(t)}{\lambda(t)}) e^{-i\gamma}$, thus all the term $E(\tilde{Q}_b + \epsilon)$ is equal to $\lambda^2 E$. In our model, we cannot make this substitution and have to keep the term of form $E(\tilde{Q}_{b_i} + \varepsilon^j)$.

5.1.1. Modulation estimates

Now, by exactly the same argument as the proof of Lemma 5 in [20] we have the analogue of Lemma 1.4

Lemma 5.1. In the setting of Lemma 3.1, the following modulation estimates hold for $t \in [0, T]$, and j = 1, ..., m:

$$\left|\frac{1}{\lambda_j}\frac{d\lambda_j}{ds_j} + b_j\right| + \left|\frac{db_j}{ds_j}\right| \le C\left(\int |\nabla\varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|}\right) + \Gamma_{b_j}^{1-C\eta} + C|E(\tilde{Q}_{b_j} + \varepsilon^j)|,\tag{5.5}$$

$$\begin{aligned} &|\frac{d\tilde{\gamma}_{j}}{ds_{j}} - \frac{1}{|Q_{1}|^{2}} (\varepsilon_{1}^{j}, L_{+}Q_{2}| + |\frac{1}{\lambda_{j}} \frac{dx_{j}}{ds_{j}}| \\ &\leq \delta(\alpha) (\int |\nabla \varepsilon^{j}|^{2} e^{-2(1-\eta)\frac{\theta(b_{j}|y|)}{b_{j}}} + \int |\varepsilon^{j}|^{2} e^{-|y|}|)^{\frac{1}{2}} + C \int |\nabla \epsilon^{j}|^{2} + \Gamma_{b_{j}}^{1-C\eta} \\ &+ C |E(\tilde{Q}_{b_{j}} + \varepsilon^{j})|. \end{aligned}$$
(5.6)

The proof follows exactly as the proof of (2.36), (2.37) in Lemma 5 of [20].

5.1.2. Estimates by the conservation law

Similarly, following the proof of Lemma 5 in [20], we have the analogue of Lemma 1.5,

Lemma 5.2. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, for j = 1, ..., m

$$|2(\varepsilon_1^j, \Sigma_{b_j}) + 2(\varepsilon_2^j, \Theta_{b_j})| \le C\left(\int |\nabla \varepsilon^j|^2 + |\varepsilon^j|^2 e^{-|y|}\right) + \Gamma_{b_j}^{1-C\eta} + C|E(\tilde{Q}_{b_j} + \varepsilon^j)|,$$
(5.7)

$$|(\epsilon_2^j, \nabla \Sigma_{b_j})| \le C\delta(\alpha) (\int |\nabla \varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|})^{\frac{1}{2}} + CP(\tilde{Q}_{b_j} + \varepsilon^j).$$
(5.8)

Proof. The proof of (5.7) is exactly the same as the proof of (2.35) in Lemma 5 in [20]. The proof of (5.8) is a little different, since in [20] the authors use the zero-momentum condition which is not used here. A direct simple algebraic computation shows that

$$P(\tilde{Q}_{b_j} + \varepsilon^j) = -(\nabla \Sigma_{b_j}, \varepsilon_2^j) + (\nabla \Theta_{b_j}, \varepsilon_1^j) - 2(\nabla \varepsilon_1^j, \varepsilon_2^j)$$

By a point-wise control $|\nabla \Theta_{b_i}|(y) \leq e^{-|K||y|}$ and Cauchy Schwartz, one has:

$$|(\nabla \Theta_{b_j}, \epsilon)| \leq C (\int |\varepsilon^j|^2 e^{-K|y|})^{\frac{1}{2}}$$

From here one has that $|(\nabla \varepsilon_1^j, \varepsilon_2^j)| \le ||\varepsilon||_2 ||\nabla \varepsilon||_2$.

Using the general functional analysis fact¹¹

$$\int |\varepsilon^j|^2 e^{-K|y|} \le C_K (\int |\nabla \varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|}),$$

the bootstrap hypothesis (3.39) and (3.42):

$$\|\varepsilon^{j}\|_{L^{2}} \leq \delta(\alpha), \int |\nabla \varepsilon^{j}|^{2} + \int |\varepsilon^{j}|^{2} e^{-|y|} \leq \delta(\alpha),$$

we have that (5.8) follows. \Box

¹¹ This estimate holds for all H^1 functions, see (2.38) in Lemma 5 in [20].

5.1.3. Control of local quantity

We use the bootstrap hypothesis (in particular, the control of local conserved quantity) to show the following

Lemma 5.3. In the setting of Lemma 3.1, for $t \in [0, T]$ and j = 1, ..., m the following estimates hold

$$|E(\varepsilon^{j} + \tilde{Q}_{bj})| \lesssim \lambda_{j}^{2-},$$

$$|P(\varepsilon^{j}) + \tilde{Q}_{bj}| \lesssim \lambda_{j}^{1-}.$$

$$(5.9)$$

Proof. We only prove (5.9), and the second inequality follows by a similar argument.

Direct computation shows that

$$|E(\varepsilon^{J} + \tilde{Q}_{b_{j}})(t) - \lambda_{j}^{2}(t)E_{loc}(x_{j}(t), u(t))|$$

= $\lambda_{j}^{2}|E(u\chi_{1,loc}(x - x_{j})) - E_{loc}(x_{j}(t), u(t))|$
 $\leq \lambda_{j}^{2} \int_{\frac{2}{3} \leq |x - x_{j}| \leq 1} |\nabla u|^{2} + |u|^{4}.$ (5.11)

Note that by a standard Sobolev imbedding and by the bootstrap hypothesis (3.47), we have:

$$\int_{\frac{2}{3} \le |x-x_j| \le 1} |\nabla u|^2 + |u|^4 \le ||u||_{H^{N_1}(\min_j\{|x-x_{j,0}|\} \ge \frac{1}{2}\})}^4 \le (\frac{1}{\max_j\{\lambda_{j,0}\}})^4.$$

(We will choose N_2 large enough, and thus $N_1 < \frac{N_2}{2}$ can be chosen large enough so that all the desired Sobolev embedding holds.)

On the other hand, by the bootstrap hypothesis (3.45), and the assumption on initial data (3.18), we obtain

$$\lambda_j^2 |E_{loc}(x_j(t), u(t))| \lesssim \lambda_j^2$$

We plug these two estimates into (5.11), and we obtain

$$|E(\varepsilon^j + \tilde{Q}_{b_j})| \lesssim \lambda_j^2 \left(1 + \left(\frac{1}{\max_j \{\lambda_{j,0}\}}\right)^4 \right).$$
(5.12)

Finally, by the bootstrap hypothesis (3.43), we have $\lambda_j(t) \leq \max_j \{\lambda_{j,0}\}$ for j = 1, ..., m. Thus, (5.9) clearly follows from (5.12). \Box

5.1.4. Modulation estimates and estimates by conservation law, restated

In this section we summarize what we have found above.

Lemma 5.4. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, and for j = 1, ..., m

$$\left|\frac{1}{\lambda_j}\frac{d\lambda_j}{ds_j} + b_j\right| + \left|\frac{db_j}{ds_j}\right| \le C\left(\int |\nabla\varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|}\right) + \Gamma_{b_j}^{1-C\eta},\tag{5.13}$$

$$\begin{aligned} &|\frac{d\tilde{\gamma}_{j}}{ds_{j}} - \frac{1}{|Q_{1}|^{2}} (\varepsilon_{1}^{j}, L_{+}Q_{2}| + |\frac{1}{\lambda_{j}} \frac{dx_{j}}{ds_{j}}| \\ &\leq \delta(\alpha) (\int |\nabla \varepsilon^{j}|^{2} e^{-2(1-\eta)\frac{\theta(b_{j}|y|)}{b_{j}}} + \int |\varepsilon^{j}|^{2} e^{-|y|}|)^{\frac{1}{2}} + C \int |\nabla \epsilon^{j}|^{2} + \Gamma_{b_{j}}^{1-C\eta}, \end{aligned}$$
(5.14)

$$|2(\varepsilon_1^j, \Sigma_{b_j}) + 2(\varepsilon_2^j, \Theta_{b_j})| \le C\left(\int |\nabla \varepsilon^j|^2 + |\varepsilon^j|^2 e^{-|y|}\right) + \Gamma_{b_j}^{1-C\eta},\tag{5.15}$$

$$|(\epsilon_2^j, \nabla \Sigma_{b_j})| \le C\delta(\alpha) (\int |\nabla \varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|}).$$
(5.16)

Proof. We just need to combine Lemma 5.1, Lemma 5.2, Lemma 5.3. □

5.1.5. Local virial estimate and Lyapounov functional control

Below we combine the orthogonality conditions (3.34), (3.35), (3.36), (3.37), modulation estimates (5.13), (5.14) and estimates induced by the (local) conservation laws (5.15), (5.16), following the work of Merle and Raphaël to obtain the analogue of Lemma 7 in [20].

Lemma 5.5 (Local virial). In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, and j = 1, ..., m

$$\frac{db_j}{ds_j} \ge \delta_1 \left(\int |\nabla \varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|} \right) - \Gamma_{b_j}^{1-C\eta},\tag{5.17}$$

where δ_1 is a universal constant.

Proof. This lemma is highly nontrivial, and it is actually one of the key elements in the work of Merle and Raphaël. However, by applying exactly the same argument, which they used to derive Lemma 7 in [20], one can derive

$$\frac{db_j}{ds_j} \ge \delta_1 \left(\int |\nabla \varepsilon^j|^2 + \int |\varepsilon^j|^2 e^{-|y|} \right) - \Gamma_{b_j}^{1-C\eta} - C|E(\tilde{Q}_{b_j} + \varepsilon^j)|, \ j = 1, \dots, m.$$
(5.18)

Now using estimate (5.9) in Lemma 5.3, one obtains $|E(\tilde{Q}_{b_j} + \varepsilon^j)| \leq \lambda_j^{2-}$, which is negligible compared to $\Gamma_{b_j}^{1-C\eta}$, since we have (3.40). Thus, (5.17) clearly follows from (5.18). \Box

Next, we recover the Lyaponouv functional, which is essential to establish the sharp log-log regime. This is the analogue of Proposition 8 in [20].

Lemma 5.6. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, and for j=1,...,m

$$\frac{d\mathcal{J}_j}{ds_j} \le -Cb_j \{\Gamma_{b_j} + \int |\nabla \tilde{\varepsilon}^j|^2 + \int |\tilde{\varepsilon}^j|^2 e^{-|y|} + \int_{A_j}^{2A_j} |\varepsilon^j|^2 \},\tag{5.19}$$

with

$$\mathcal{J}_{j} := \left(\int |\tilde{Q}_{b_{j}}|^{2} - \int Q^{2} \right) + 2(\varepsilon_{1}^{j}, \Sigma_{b_{j}}) + 2(\varepsilon_{2}^{j}, \Theta_{b_{j}}) + \int (1 - \phi_{A_{j}}) |\varepsilon^{j}|^{2} \\ - \frac{\delta_{1}}{800} \left(b \,\tilde{f}_{1}(b_{j}) - \int_{0}^{b_{j}} \tilde{f}_{1}(v) dv + b\{(\varepsilon_{2}^{j}, \Lambda \Re \tilde{\zeta}_{b_{j}}) - (\varepsilon_{1}^{j}, \Lambda \Im \tilde{\zeta}_{b_{j}})\} \right),$$
(5.20)

where

$$\tilde{f}_{1}(b) := \frac{b}{4} |y\tilde{Q}_{b}|_{2}^{2} + \frac{1}{2} \Im(\int y\nabla\tilde{\zeta}_{b}\tilde{\tilde{\zeta}}_{b}),$$

$$\tilde{\varepsilon}^{j} = \varepsilon^{j} - \tilde{\zeta}_{b_{j}},$$
(5.21)
(5.22)

and ϕ_{A_j} is a non-negative smooth cut-off function, j = 1, ..., m

$$\begin{cases} \phi_{A_j}(x) = 0, |x| \le \frac{A_j}{2}, \\ \frac{1}{4_{A_j}} \le |\nabla \phi_{A_j}| \le |\frac{1}{2A_j}|, A_j \le |x| \le 2A_j, \\ \phi_{A_j}(x) = 1, |x| \ge 3A_j, \end{cases}$$
(5.23)

where $A_j = A(b_j) = \Gamma_{b_j}^{-a}$.

Most parts of the proof follows directly the proof of Proposition 8 in [20], however, some extra technical elements need to be treated hence for completeness we will explain the proof in Subsection 5.2.

We have the following control for the scale of the Lyapounouv functional:

Lemma 5.7. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, and for j = 1, ..., m,

$$\frac{\mathcal{J}_j}{b_j^2} = C^* (1 + O(\delta(\alpha)), \tag{5.24}$$

where C^* is some fixed constant.

Proof. Following the proof of (5.15) in [20], one can derive $\mathcal{J}_j \sim b_j^2$. Further refined analysis, will give (5.24), see the formula between (5.24) and (5.25) in [20]. \Box

5.1.6. Bootstrap estimates except (3.56)

So far, we already have all the ingredients to prove most of the bootstrap estimates. In fact for (3.48), (3.49), (3.50), (3.51), (3.52), (3.53) one can follow the arguments of Planchon and Raphaël in [23], which we will review for completeness in Subsection 5.4. Here we prove instead (3.54), (3.55). Actually we only show the details for (3.54), since (3.55) is similar.

Proof of (3.54). A direct computation (using *u* that solves (1.1)) shows that for j = 1, ..., m

$$|E_{loc}(x_{j}(t), u(t)) - E(x_{j,0})| \leq \int_{0}^{t} |\frac{d}{d\tau} E_{loc}(x_{j}(\tau), u(t))| d\tau$$

$$\leq \int_{0}^{T} |\frac{d}{dt} E_{loc}(x_{j}(t), u(t))| dt$$

$$\leq E_{1} + E_{2},$$
(5.25)

where

T

$$E_1 = \int_0^1 \int_{\mathbb{R}^2} \left| \left(\frac{1}{2} |\nabla u|^2 - \frac{1}{4} |u|^4 \right) \nabla \chi_{0,loc}(x - x_j(t)) \frac{dx_j}{dt} \right|,$$
(5.26)

$$E_2 = \int_0^T \int \frac{1}{2} |\nabla \chi_{0,loc}(x - x_j(t)) \left(2i\Im \Delta u \nabla \bar{u} + 2i\Im |u|^2 u \nabla \bar{u} \right)|.$$
(5.27)

Recall that $\chi_{0,loc}(x - x_j)$ vanishes for $\frac{3}{4} \le |x - x_j| \le 1$. Using the bootstrap hypothesis (3.44), and Sobolev embedding, we obtain

$$\int \left| \left(\frac{1}{2} |\nabla u|^2 - \frac{1}{4} |u|^4 \right) \nabla \chi_0(x - x_j(t)) \frac{dx_j}{dt} \right| \lesssim \|u\|_{H^N(\min_j\{|x - x_{j,0}|\} \ge \frac{1}{2})}^4,$$

$$\int \frac{1}{2} |\nabla \chi_0(x - x_j) \left(\Delta u \nabla \bar{u} - \nabla u \Delta \bar{u} + |u|^2 u \nabla u - |u|^2 \bar{u} \nabla u \right) | \lesssim \|u\|_{H^N(\min_j\{|x - x_{j,0}| \ge \frac{1}{2})}^4.$$

$$(5.28)$$

Thus,

$$E_1 + E_2 \lesssim \int_0^T (1 + \frac{d}{dt} x_j) \|u\|_{H^N(\min_j\{|x - x_{j,0}|\} \ge \frac{1}{2})}^4.$$
(5.29)

By the bootstrap hypothesis (3.47), we have that

$$||u||_{H^N(\min_j |x-x_{j,0}| \ge \frac{1}{2})}^4 \le \frac{1}{\max_j \{\lambda_{j,0}\}}.$$

Note also that by the bootstrap hypothesis (3.43), also $\frac{1}{\max_j \{\lambda_{j,0}\}} \lesssim \frac{1}{\lambda_j(t)}^{\frac{1}{2}}$ for all *j*. By the modulation estimate (5.14), $|\frac{1}{\lambda_j} \frac{dx_j}{ds_j}| \lesssim 1$. Thus, $|\frac{dx_j}{dt}| = |\frac{1}{\lambda_j^2} \frac{dx_j}{ds_j}| \lesssim \frac{1}{\lambda_j}$, hence

$$E_1 + E_2 \lesssim \int_0^T \frac{1}{\lambda_j(t)^{1.5}} dt \lesssim \lambda_{j,0}^{0.1} \int_0^T \lambda_j(t)^{1.6}.$$
(5.30)

Here we use $\lambda_j(t) \lesssim \lambda_{j,0}$, i.e. (3.43).

Finally we have for j = 1, ..., m that

$$\int_{0}^{1} \frac{1}{\lambda_{j}^{\mu}} \lesssim C_{\mu}, \forall \mu < 2.$$
(5.31)

This is the analogue of (4.1), and it follows exactly the same proof that only relies on the bootstrap hypotheses (3.40), (3.41). Clearly (5.30) and (5.31) end the proof. \Box

5.1.7. Propagation of regularity and end of bootstrap estimate

To end the bootstrap estimate, we still have to show (3.56), and it will be done in Subsection 5.3.

5.2. Recovering the Lyapounouv functional under bootstrap hypothesis

This subsection is devoted to the proof of Lemma 5.6. We emphasize again that this proof follows the computation in [20] except for two points:

- We do not use the global energy, and all those $\lambda^2 E$ terms in [20] are replaced by $E(\varepsilon^j + \tilde{Q}_{b_i})$.
- In the definition of the original Lyapounov functional \mathcal{J} in (1.56), there is a term $\|\tilde{Q}_b\|_2^2 \|Q\|_2^2 + (\epsilon_1, \Sigma) + (\epsilon_2, \Theta) + \|\epsilon\|_2^2$, that is actually a constant in [20], thanks to the conservation of L^2 mass. In our case, since we analyze the dynamics locally, the natural substitution is the local mass, which is no longer a constant. So we need to show that the local mass is slowly varying.

We need the following two lemmas.

Lemma 5.8 (Analogue of Lemma 6 in [20]). In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$ for some universal constant δ_1 and j = 1, ..., m

$$\frac{d}{ds_j}f_1^j \ge \delta_1\left(\int |\nabla \tilde{\varepsilon}^j|^2 + \int |\tilde{\varepsilon}^j|^2 e^{-|y|}\right) + c\Gamma_{b_j} - \frac{1}{\delta_1}\int_{A_j}^{2A_j} |\varepsilon^j|^2,$$
(5.32)

with

$$f_1^j(s) = \frac{b_j}{4} \| y \tilde{\mathcal{Q}}_{b_j} \|_2^2 + \frac{1}{2} \Im \left(y \nabla \tilde{\zeta}_{b_j} \tilde{\tilde{\zeta}}_{b_j} \right) + (\varepsilon_2^j, \Lambda \mathfrak{R} \tilde{\zeta}_{b_j}) - (\varepsilon_1^j, \Lambda \mathfrak{R} \tilde{\zeta}_{b_j}).$$
(5.33)

Proof. This lemma is one of the most fundamental points in [20]. We quickly recall its proof. If $\frac{1}{\lambda(t)}(Q_b + \epsilon)(\frac{x-x(t)}{\lambda(t)})e^{-i\gamma}$ solves (1.1), then one is able to derive two equations for $\epsilon := \epsilon_1 + i\epsilon_2$, i.e. (1.40), (1.41).

Then one takes the inner product of (1.41) with $\Lambda(\Sigma_b + \Lambda \Re \tilde{\zeta})$ and of (1.41) with $-\Lambda(-\Theta_b + \Lambda \Im \tilde{\zeta})$ and sums. The detailed computation is displayed in Appendix B of [20].

1464

Here we follow exactly the same procedure. We pick j = 1, ..., m and we recall that by definition,

$$\frac{1}{\lambda_j}(\tilde{Q}_{b_j}+\varepsilon^j)(\frac{x-x_j}{\lambda_j})e^{-i\gamma_j}=\chi_{1,loc}(x-x_j)u,$$

and $\chi_{1,loc}(x - x_j)u$ almost solves (1.1). One may derive similar equations for $\varepsilon^j(y) = \varepsilon_1^j(y) + i\varepsilon_j(y)$ as in (1.40), (1.41) with some extra terms in right hand side. Since $\chi_{1,loc}(x - x_j) \equiv 1$ in $|x - x_j| \leq \frac{2}{3}$, these extra terms are supported in $|y| \geq \frac{1}{2\lambda_j}$ (i.e. $|x - x_j| \geq \frac{1}{2}$). Since $\tilde{Q}_{b_j}, \tilde{\zeta}_{b_j}$ is supported in $|y| \leq \Gamma_{b_j}^{-a}$ and by the bootstrap hypothesis (3.40), $\Gamma_{b_j}^{-a} \ll \frac{1}{\lambda_j}$, then when one pairs these equations with $\Lambda(\Sigma_{b_j} + \Lambda\mathfrak{R}\tilde{\zeta}_{b_j})$ or $-\Lambda(-\Theta_{b_j} + \Lambda\mathfrak{R}\tilde{\zeta}_{b_j})$, these extra terms automatically cancel.

As a consequence all the algebraic computations in the appendix B of [20] follow. There is one more difference compared to Appendix B of [20]. There the authors use the energy conservation (formula (4.17) in [20]), here instead we need to replace term $\lambda^2 E_0$ in that formula by the term $E(\varepsilon^j + \tilde{Q}_{b_j})$. Once this is done we follow the argument in [20] to recover the virial type estimate

$$\frac{d}{ds_j}f_1^{\,j} \ge \delta_1\left(\int |\nabla \tilde{\varepsilon}^j|^2 + \int |\tilde{\varepsilon}^j|^2 e^{-|y|}\right) + c\Gamma_{b_j} - \frac{1}{\delta_1}\int_{A_j}^{2A_j} |\varepsilon^j|^2 - CE(\varepsilon^j + \tilde{Q}_{b_j}), \, j = 1, \dots, m.$$
(5.34)

By Lemma 5.3 and the bootstrap hypothesis (3.40), we obtain

24

$$|E(\varepsilon^j + \tilde{Q}_{b_j})| \lesssim \lambda_j^{2-} \ll \Gamma_{b_j}^{10}, \tag{5.35}$$

and as a consequence formula (5.32). \Box

We have the following Lemma.

Lemma 5.9 (Analogue of Lemma 7 in [20]). In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$, and for j = 1, ..., m

$$\frac{d}{ds_j} \{ \int |\phi_{A_j}|\varepsilon^j|^2 \} \ge \frac{b_j}{400} \int_{A_j}^{2A_j} |\epsilon|^2 - \Gamma_{b_j}^{1+z_0} - \Gamma_{b_j}^{\frac{a}{2}} \int |\nabla \varepsilon^j|^2.$$
(5.36)

Proof. The proof follows as the proof of Lemma 7 in [20] except, as above, for replacing $\lambda^2 E$ with $E(\varepsilon^j + \tilde{Q}_{b_j})$, which is much smaller than $\Gamma_{b_j}^{100}$ by the bootstrap assumption, and thus negligible. \Box

The following lemma illustrates the fact that the local mass is slowly varying with respect to the rescaled time variable s_j .

Lemma 5.10 (Slow varying of local mass). With the same assumptions as in Lemma 3.1, the following estimate holds

$$\left|\frac{d}{ds_{j}}\left\{\int |\tilde{Q}_{b_{j}}|^{2} + 2(\varepsilon_{1}^{j}, \Sigma_{b_{j}}) + 2(\varepsilon_{2}^{j}, \Theta_{b_{j}}) + \int |\varepsilon^{j}|^{2}\right\}\right| \le \Gamma_{b_{j}}^{10}.$$
(5.37)

Proof. First we observe that

$$\{\int |\tilde{Q}_b|^2 + 2(\epsilon_1, \Sigma) + 2(\epsilon_2, \Theta) + \int |\epsilon|^2\} \equiv \|\tilde{Q}_{b_j} + \epsilon^j\|_2^2.$$
(5.38)

By (3.28), one has the geometric decomposition

$$u(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_j} \tilde{\mathcal{Q}}_{b_j}(\frac{x-x_j}{\lambda_j}) e^{-i\gamma_j} + \Xi.$$
(5.39)

Also recall that $\varepsilon^{j}(y) = \chi_{1,loc}(\lambda_{j}y)\epsilon^{j}(y)$ and that $\tilde{Q}_{b_{j}}$ is supported in $|x| \leq \frac{2}{b_{j}}$. Note that by the bootstrap assumption (3.40), $\frac{1}{\lambda_{i}} \gg \frac{1}{b_{i}}$. We then obtain

$$\|\tilde{Q}_{b_j} + \varepsilon^j\|_2^2 = \|\tilde{Q}_{b_j}\chi_{1,loc}(\lambda y) + \epsilon^j \chi_{1,loc}(\lambda y)\|_2^2 = \|u(\chi_{1,loc}(x - x_j))\|_2^2.$$
(5.40)

Thus,

$$\frac{d}{ds_{j}}\left\{\int |\tilde{Q}_{b_{j}}|^{2} + 2(\epsilon_{1}^{j}, \Sigma_{b_{j}}) + 2(\epsilon_{2}^{j}, \Theta_{b_{j}}) + \int |\epsilon|^{2}\right\} \\ \lesssim \frac{dx_{j}}{ds_{j}} ||u||_{2}^{2} + 2\lambda_{j}^{2} \int_{\frac{2}{3} \le |x - x_{j}| \le \frac{3}{4}} |\nabla \chi_{1, loc}||u|| \nabla u|.$$
(5.41)

Recall now that $\frac{d}{ds_j} = \lambda_j^2 \frac{d}{dt}$. The first term on left hand side is controlled by the modulation estimate (5.14), and the conservation of mass $||u(t)||_2 = ||u_0||_2$,

$$|\frac{dx_j}{ds_j}||u||_2^2 \lesssim |\frac{1}{\lambda_j}\frac{dx_j}{ds_j}|\lambda_j.$$
(5.42)

The second term on the left hand side is controlled by the bootstrap assumption (3.47),

$$2\lambda_j^2 \int_{\frac{2}{3} \le |x-x_j| \le \frac{3}{4}} |\nabla \chi_{1,loc}| |u| |\nabla u| \lesssim \lambda_j^2.$$
(5.43)

We finally use the bootstrap assumption (3.40) to control λ_i and the desired estimate easily follows. \Box

Now we are in good position to finish the proof of Lemma 5.19.

Proof of Lemma 5.19. With Lemma 5.32 and Lemma 5.9, we follow the proof of Proposition 4 in [20] and we obtain

$$\frac{d}{ds_j} \left(\int -\phi_A |\varepsilon^j|^2 - \frac{\delta_1}{800} \left(b \tilde{f}_1(b_j) - \int_0^{b_j} \tilde{f}_1(v) dv + b\{(\varepsilon_2^j, \Lambda \mathfrak{R} \tilde{\zeta}_{b_j}) - (\varepsilon_1^j, \Lambda \mathfrak{R} \tilde{\zeta}_{b_j})\} \right) \right) \leq -C \left(\int |\nabla \varepsilon^j|^2 + |\varepsilon^j|^2 e^{-|y|} \right) + \Gamma_{b_j}^{1-C\eta}.$$
(5.44)

Now plug in the estimate in Lemma 5.10 above and the desired estimate follows. \Box

5.3. Propagation of regularity under bootstrap hypothesis

This subsection mostly follows from the arguments in Section 4, indeed, the reason why we write Section 4 is to make this subsection more accessible. First, the analogue of (4, 1) holds.

First, the analogue of (4.1) holds,

Lemma 5.11. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$

$$\int_{0}^{t} \left(\sum_{j=1}^{m} \frac{1}{\lambda_{j}^{\mu}(\tau)} \right) d\tau \lesssim \begin{cases} C_{\mu}, \mu < 2, \\ \sum_{j=1}^{m} \frac{|\ln \lambda_{j}(t)|^{101}}{\lambda_{j}(t)^{\mu-2}}, \mu \ge 2. \end{cases}$$
(5.45)

Proof. The proof of (5.45) only relies on the bootstrap hypothesis (3.40), (3.41), and it is similar to the proof of (4.1). \Box

We again introduce a sequence of cut-off functions¹² $\{\chi_l\}_{l=0}^L$

$$\chi_l = \begin{cases} 0, |x - x_{j,0}| \le a_l, \ j = 1 \text{ or } 2, \\ 1, |x - x_{j,0}| \ge b_l, \ j = 1 \text{ and } 2, \end{cases}$$
(5.46)

such that $a_l < b_l < a_{l-1}$, $a_0 = \frac{1}{250} \ll 1$, $b_L = \frac{2}{5} < \frac{1}{2}$. The idea is (again) that we want to retreat from $\chi_{l-1}u$ to $\chi_l u$ for each *l*, showing that $\chi_l u$ has higher regularity than $\chi_{l-1}u$. We still use the notation $v_i = \chi_i u$.

The key to gain regularity in Section 4 is formula (4.3). Similar estimates also hold here:

Lemma 5.12. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$,

$$\int_{0}^{T} \|\nabla(\chi_{0}u(t))\|^{2} dt \lesssim 1.$$
(5.47)

Proof. First recall that by the bootstrap hypothesis (3.44), if $x \in supp \chi_0$, then $\min_i \{|x - x_i|\} \ge \frac{1}{500}$. Thus we need only to control

$$\int_{0}^{T} \int_{\min_{j}\{|x-x_{j}|\} \geq \frac{1}{500}} \|\nabla(\chi_{0}u)\|_{2}^{2}.$$

Secondly, since u is bounded in L^2 and T is bounded (indeed $T < \delta(\alpha)$, see Remark 3.3), we need only to control

$$\int_{0}^{I} \int_{\min_{j}\{|x-x_{j}|\} \ge \frac{1}{500}} \|\nabla u\|_{2}^{2}.$$

Indeed

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$$\int_{0}^{T} \int_{\min_{j}\{|x-x_{j}|\} \ge \frac{1}{500}} |\nabla u|^{2}$$

$$\leq \int_{0}^{T} \int_{\min_{j}\{|x-x_{j}|\} \ge \frac{3}{5}} |\nabla u|^{2} + \int_{0}^{T} \sum_{j=1}^{m} \int_{\frac{1}{500} \le |x-x_{j}| \le \frac{3}{5}} |\nabla u|^{2}$$

$$:= E_{1} + E_{2}.$$
(5.48)

The first term is controlled by the bootstrap hypotheses (3.47) and (3.43). In fact we easily have:

$$E_{1} \lesssim \int_{0}^{T} \min_{j} \{|\ln \lambda_{j,0}|\}$$

$$\lesssim \int_{0}^{T} \sum_{j} |\ln \lambda_{j}(t)|$$

$$\lesssim 1.$$
(5.49)

¹² We still use the notation of χ , but the definition of χ is different from Section 4.

In last step we have applies (5.45). The second term E_2 is estimated by using the Lyapounov functional \mathcal{J}_j in (5.19). This is actually one of the key estimates in [19]. We recall it here. Note that

$$\chi_{1,loc}(x) \equiv 1 \text{ for } |x| \le \frac{2}{3} \text{ and } \chi_{1,loc}(x-x_j) = \frac{1}{\lambda_j} (\tilde{Q}_{b_j} + \varepsilon^j) (\frac{x-x_j}{\lambda_j}) e^{-i\gamma_j}.$$

We make the observation that since $\lambda_j \lesssim e^{-e^{\frac{\pi}{10b_j}}}$ by the bootstrap hypothesis (3.40), when $|x - x_j| \ge \frac{1}{500}$ we have $\tilde{Q}_{b_j}(\frac{|x-x_j|}{\lambda_j}) \equiv 0$, since $supp \tilde{Q}_{b_j} \subset \{|y| \lesssim \frac{1}{b_j}\}$. Similarly, in this region the radiation term $\tilde{\zeta}_{b_j}(\frac{x-x_j}{\lambda_j})$ also vanishes. Thus, for j = 1, ..., m

$$E_{2} = \int_{\frac{1}{500} \le |x - x_{j}| \le \frac{3}{5}} |\nabla u|^{2} = \int_{\frac{1}{500} \le |x - x_{j}| \le \frac{3}{5}} |\nabla \frac{1}{\lambda_{j}} (\tilde{Q}_{b_{j}} + \varepsilon^{j}) (\frac{x - x_{j}}{\lambda_{j}}) e^{-i\gamma_{j}}|^{2}$$

$$= \lambda_{j}^{-2} \int_{\frac{1}{500\lambda_{j}} \le |x| \le \frac{3}{5\lambda_{j}}} |\nabla \varepsilon^{j}|^{2} = \lambda_{j}^{-2} \int_{\frac{1}{500\lambda_{j}} \le |x| \le \frac{3}{5\lambda_{j}}} |\nabla \tilde{\varepsilon}^{j}|^{2}.$$
(5.50)

As a consequence

$$E_2 \leq \int_{s_j(t_1)}^{s_j(t_2)} |\nabla \tilde{\varepsilon}^j|^2 \lesssim sup_{t \in [0,T)} \sqrt{\mathcal{J}_j} \sim b_j \lesssim 1.$$
(5.51)

In the last step, we use (5.19) and the fact that the Lyapounov functional $\mathcal{J}_j \sim b_j^2$, see Lemma 5.7. This is enough to end the proof. \Box

We now again use I-method to recover a rough control:

Lemma 5.13. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$,

$$\|u(t)\|_{H^{N_2}} \lesssim_{\sigma} \left(\sum_{j} \frac{1}{\lambda_j(t)}\right)^{N_2 + \sigma}.$$
(5.52)

Proof. Without loss of generality, we only show (5.52) for t = T. We remark that if one defines $\tilde{\lambda} = \min_{j} \{\lambda_{j}\}$, then $\|u(t)\|_{H^{1}} \sim \frac{1}{\tilde{\lambda}(T)}$, and (5.52) is equivalent to $\|u(t) \sim \left(\frac{1}{\tilde{\lambda}(t)}\right)^{N_{2}+\sigma}$.

The proof is almost the same as the proof of (4.37) in Lemma 4.1. Following Remark 4.4 we only need to show that we can divide [0, T] into disjoint intervals $\cup I_{k,h}$, such that

 $\begin{aligned} &- \|u(t)\|_{H^1} \sim 2^k \text{ for } t \in I_{k,j}. \\ &- |I_{k,h}| \sim \frac{1}{2^{2k}}, \forall k, h. \end{aligned}$

$$- \sharp \{I_{k,h}\} \lesssim (|\ln \lambda(T)|)^3.$$

These clearly follow from the facts below. Fix j, we can divide [0, T] into disjoint intervals $\cup I_{k,h}^{j}$, such that

$$-\lambda_j(t) \sim 2^{-k} \text{ for } t \in I_{k,j}$$

$$-|I_{k,h}^j| \sim \frac{1}{2^{2k}}$$

$$- \sharp \{I_{k,h}^j\} \lesssim (|\ln \lambda_j(T)|)^3$$

Now we prove these facts. Then by the bootstrap hypothesis (3.43), we can divide [0, T] into $0 = t_{k_0} < ... t_k < ... t_{K(T)} = T$ such that $\lambda_1(t) \sim 2^{-k}$ for $t \in [t_k, t_{k+1}]$, then relying on the bootstrap hypothesis (3.40), (3.41), similarly

as in the proof of Lemma 4.6, one can show that $t_{k+1} - t_k \leq \sqrt{k}\lambda_1^{2(t_k)} \sim \sqrt{k}2^{-2k}$. This estimate is enough for us to further divide $[t_k, t_{k+1}]$ into disjoint intervals $\bigcup_{h=1}^{J_h} [\tau_k^h, \tau_k^{h+1}]$ such that our desired $I_{k,h}$ can be chosen as $[\tau_k^h, \tau_k^{h+1}]$, see (4.16), (4.17). \Box

We now remark that the bootstrap estimate (3.56) follows from the lemma below.

Lemma 5.14. In the setting of Lemma 3.1, the following estimates hold for $t \in [0, T]$,

- 1. $\forall v < 1, \|v_1\|_{H^{\nu}} \lesssim 1.$
- 2. If there holds estimate for some r > 0, and some $1 \le i \le L 1$

$$\|v_i\|_{H^r} \lesssim 1,\tag{5.53}$$

then we will have a gain of regularity on v_{i+1} ,

$$\forall \tilde{r} < \frac{2(K_2 - r)}{K_2} - 1 + r, \|v_i\|_{H^{\tilde{r}}} \lesssim 1.$$
(5.54)

The proof follows as the proof of Lemma 4.2, Lemma 4.3, Lemma 4.4, Lemma 4.5, see also Proof of (2.5) in Lemma 2.2.

5.4. Proof of bootstrap estimate except (3.56), (3.54), (3.55)

The proof of these estimates can be found (up to a small modification) in [23,26] and the references therein. See in particular Proposition 1 in [26]. We quickly review those estimates for the convenience of the readers.

5.4.1. Proof of (3.48)

We bound the left hand side of (3.48) by $C\alpha$. The control of b_j directly follows from the bootstrap hypothesis (3.41), which implies $0 \le b_j \le \frac{10\pi}{\ln s_{j,0}} \le b_{j,0}$ and we note that by the initial condition (3.14), we have $b_{j,0} \le \alpha$. Thus $|b_j| \le \alpha \ll \alpha$. The control of $\|\nabla \varepsilon^j\|_2$ follows from the bootstrap hypothesis (3.42), therefore we have $\|\nabla \varepsilon^j\|_2 \le \gamma_{b_j}^{\frac{1}{100}} \ll \alpha$. The control of $\|\varepsilon^j\|_2$ comes from the L^2 conservation law. Indeed

$$\sum_{j=1}^{m} \|\varepsilon^{j}\|_{2}^{2} + \sum_{j=1}^{m} \|\tilde{Q}_{b_{j}}\|_{2}^{2} \le \|u\|_{2}^{2} + O\left(\sum_{j=1}^{m} (|\varepsilon^{j}|, |\tilde{Q}_{b_{j}}|)\right).$$

Note that

$$(|\varepsilon^{j}|, |\tilde{\mathcal{Q}}_{b_{j}}|) \leq \int |\nabla \varepsilon^{j}|^{2} + \int |\varepsilon^{j}|^{2} e^{-|y|}, \qquad (5.55)$$

$$\|\hat{Q}_{b_j}\|_2^2 = \|Q\|_2^2 + O(|b_j|^2), \tag{5.56}$$

$$\|u\|_2^2 = \|u_0\|_2^2. \tag{5.57}$$

Thus the control of $\|\nabla \varepsilon^{j}\|_{2}$ follows easily from the bootstrap hypothesis (3.42) and our choice of initial data.

5.4.2. Proof of (3.49), (3.50)

We first show (3.50). It follows from the local virial estimate (5.17) and the control of the Lyapounouv functional (5.19). We first show the lower bound of (3.50). Note that by (5.17), one obtains $\frac{db_j}{ds_j} \ge -\Gamma_{b_j}^{1-C\eta}$, which implies

$$\left(\frac{de^{\frac{\pi(1-C\eta)}{b_j}}}{ds_j}\right) \le 1.$$

Recall that $s_j \ge s_{j,0} \equiv e^{\frac{3\pi}{4b_{j,0}}}$, thus we have:

$$e^{\frac{\pi(1-C\eta)}{b_j}} \le s_j - s_{j,0} + e^{\frac{\pi(1-C\eta)}{b_{j,0}}} \le s_j + e^{\frac{\pi(1-C\eta)}{b_{j,0}}} \le s_j^{\frac{4}{3}}.$$
(5.58)

Therefore one obtains $b_{j,0} \ge \frac{\pi}{5 \ln s_j}$, which is the lower bound of (3.50). Now we turn to the upper bound of (3.50). The control (5.19) for the Lyapounov functional directly implies

$$\frac{d\mathcal{J}_j}{ds_j} \le -Cb_j \Gamma_{b_j}. \tag{5.59}$$

From Lemma 5.7, we know that there exists a C_* such that

$$\frac{\mathcal{J}_j}{b_j^2} = C^* (1 + O(\delta(\alpha))). \tag{5.60}$$

Let $g_j = \sqrt{\mathcal{J}_j}$, then by combining (5.59) and (5.60), one obtains:

$$\frac{dg_j}{ds_j} = \frac{1}{2} \frac{d\mathcal{J}_j}{s_j} \frac{1}{g_j} \le -Cg_j \times \frac{1}{2} C^* e^{\frac{\pi(1+C\eta)\sqrt{C^*(1-O(\delta(\alpha))})}{g_j}}.$$
(5.61)

One obtains

$$\frac{d}{d_{s_j}} \left(e^{\frac{\pi \sqrt{C^*}(1+O(\delta(\alpha)))}{g_j}} \right) \ge 1.$$
(5.62)

Thus

$$e^{\frac{\pi(1+O(\delta(\alpha)))}{b_j}} \ge s_j - s_{j,0} + e^{\frac{\pi(1+O(\delta(\alpha)))}{b_{j,0}}} \ge s_j,$$
(5.63)

where we recall that $s_{j,0} = e^{\frac{3\pi}{4b_{j,0}}}$. This implies the lower bound of (3.50).

Now we turn to the proof of (3.49). By (5.58), (5.63), we have

$$\frac{3\pi}{4\ln s_j} \le b_j \le \frac{4\pi}{3\ln s_j}.\tag{5.64}$$

We only show the upper bound of (3.49), the lower bound will be similar.

The computation follows from the proof of Lemma 6 in [25], in fact formula (2.68) mentioned in that lemma is exactly the upper bound we want.

The modulation estimate (5.5) plus the bootstrap hypothesis (3.42) imply that

$$\left|\frac{1}{\lambda_j}\frac{d\lambda_j}{ds_j} + b_j\right| \le C\Gamma_{b_j}^{\frac{1}{2}}.$$
(5.65)

Thus

$$\frac{d\lambda_j}{ds_j} = -(1+O(\delta(\alpha)))b_j,$$

and from here

$$\lambda_j((t)) \le \lambda_j(s_j(t)) \equiv \lambda_j(s_j) \le \lambda_{j,0} + \int_{s_{j,0}}^s \frac{1}{2\ln s}.$$

(Note that according to our notation $\lambda_{i,0} = \lambda_i(s_{i,0})$.)

A direct calculation implies that

$$-\ln\lambda_j(s_j) \ge -\frac{1}{2}\ln\lambda_j(s_{j,0}) + \frac{\pi}{3}\frac{s}{\ln s}$$

which further implies that $s_j \lambda(s_j) \le e^{-e^{\frac{3\pi}{8b_j}}}$. (See (3.106), (3.107) in [25] for more details.) This already implies our desired upper bound.

1470

5.4.3. Proof of (3.51)

This is exactly the step 3 of the proof of Proposition 5 in [20], called pointwise control of ϵ by b.

5.4.4. Proof of (3.52)

The modulation estimate (5.5) plus the bootstrap hypothesis (3.42) imply

$$\left|\frac{1}{\lambda_j}\frac{d\lambda_j}{ds_j} + b_j\right| \le C\Gamma_{b_j}^{\frac{1}{2}}.$$
(5.66)

Thus

$$\frac{d\lambda_j}{ds_j} = -(1 + O(\delta(\alpha)))b_j\lambda_j < 0$$

(since $b_i > 0$ due to hypothesis (3.41)). This clearly implies (3.52).

5.4.5. Proof of (3.53)

This easily follows from the modulation estimate (5.6) and the bootstrap hypotheses (3.40) and (3.41), as in [23]. We quickly review it for completeness.

For all $t \in [0, T]$,

$$|x_{j} - x_{j,0}| \leq \int_{0}^{t} |\frac{dx_{j}}{dt}| dt$$

$$= \int_{0}^{t} \frac{dx_{j}}{ds_{j}} \frac{1}{\lambda_{j}^{2}} dt = \int_{s_{j,0}}^{s_{j}(t)} |\frac{1}{\lambda_{j}} \frac{dx_{j}}{ds_{j}}| \lambda_{j} ds.$$
(5.67)

Note that by (5.6), $|\frac{1}{\lambda_j} \frac{dx_j}{ds_j}| \le \delta(\alpha)$, and by (3.40) and (3.41), one has $\lambda_j(s) \le s_j^{-100}$. Thus one clearly has $|x_j - x_{j,0}| \le \delta(\alpha)$, which easily implies (3.53).

Remark 5.1. The above computation indeed shows

$$\int_{0}^{T^{+}(u)} \left| \frac{dx_j}{dt} \right| dt < \infty, \tag{5.68}$$

which of course implies $\lim_{t\to T^+} x_i(t)$ exists.

6. Proof of Main Theorem

We will need several parameters throughout the whole section.

 $1 \gg a_0 \gg a_1 \gg a_2 > 0.$

Recall our goal is to construct log–log blow up solution u, which blows up at m prescribed points $x_{1,\infty}, \ldots x_{m,\infty}$, and has asymptotic near blow up time as (1.11). We will still focus on dimension d = 2. And since we have scaling symmetry, we assume without loss of generality

$$|x_{j,\infty} - x_{j',\infty}| \ge 20, \, j \ne j'. \tag{6.1}$$

It is clear that by Lemma 3.1, all initial data u_0 describe in Subsection 3.1, the associated solution u have the following geometric decomposition for $t < T^+(u)$.

$$u(t,x) = \sum_{j} \frac{1}{\lambda_{j}(t)} \tilde{Q}_{b_{j}}(\frac{x - x_{j}(t)}{\lambda_{j}}) e^{-i\gamma_{j}(t)} + \Xi(t,x),$$
(6.2)

such that for each i,

$$\chi_{1,loc}(x-x_j)u(t,x) \equiv \frac{1}{\lambda_j}(\tilde{Q}_{b_j}+\varepsilon^j)(\frac{x-x_j}{\lambda_j})e^{-i\gamma_j},$$

and $\Xi \equiv (1 - \sum_{j} \chi_{1,loc})u$ in the region $\{x \mid |x - x_{j,0}| \ge \frac{1}{2}, j = 1, \dots, m\}$ and bounded in H^{N_1} in this region. To finish the construction in Theorem 1.1, we need to construct initial data u_0 such that

- (a) The associated solution u blows up in finite time according to the log-log law.

- (b) The *m* points blows up simultaneously, i.e.

$$\lambda_j(t) \xrightarrow{t \to T^+} 0, \, j = 1, \dots, m. \tag{6.3}$$

- (c) The blow up points are as prescribed, i.e.

$$x_j(t) \xrightarrow{t \to T^+} x_{j,\infty}.$$
(6.4)

It is easy to check once we get (a), (b), (c), then other requirements in the Theorem 1.1 will be automatically satisfied. To see this, just observe

- 1. $\lambda_i \rightarrow 0 \Rightarrow b_i \rightarrow 0$ since we have (3.49).
- Thus, *Q̃_{bj}* converges strongly to *Q* in the sense of (1.29).
 ε^j is bounded in L² by (3.48) and converges to 0 in H¹ as b_j → 0 by (3.51).

Now let us turn to condition (a), (b), (c). Indeed, Lemma 3.1 already implies for all data u_0 , the associated solution u will blow up in finite time with log-log blow up rate

$$\|\nabla u(t)\|_2 \sim \sqrt{\frac{\ln|\ln T - t|}{T - t}}.$$
 (6.5)

Not all data described as in Subsection 3.1 will give a solution satisfy condition (b), (c). Morally speaking, the initial data in Subsection 3.1 has two types of parameters, $\lambda_{i,0}$ and $x_{i,0}$. $\lambda_{i,0}$ describes how concentrated the j-th bubble is and the $x_{j,0}$ describes the initial position of jth bubble. Lemma 3.1 says the m bubbles evolve with very weak interaction with each other. Thus, we need to choose λ_i carefully to make all the bubbles to blow up at the same time, and choose initial position $x_{i,0}$ to make the bubbles blow up at the prescribed points $x_{i,\infty}$. This is achieved by certain topological argument. If one does not want to prescribed the blow up points and only wants to get condition (b), an argument similar to the topological argument in [4] will suffice. However, if one wants to further prescribe the blow up points, then the problem is actually more tricky, since the choice of $\lambda_{i,0}$ and $x_{i,0}$ will be coupled.

Before we continue, we point out it is not hard to show $\lim_{t\to T^+} x_i(t)$ exists. Indeed, as remarked, we actually have

$$\int_{0}^{T^{+}} \left| \frac{dx_{j}}{dt} \right| dt < \infty.$$
(6.6)

The hard part here is to prescribe the limit, i.e. for given $x_{1,\infty}, \ldots x_{m,\infty}$, we want to construct solution such that

$$\lim_{t \to T^+} x_j(t) = x_{j,\infty}.$$
(6.7)

6.1. Preparation of data

Let us first fix a smooth-cut off function

$$\chi(x) = \begin{cases} 1, |x| \le 1, \\ 0, |x| \ge 2. \end{cases}$$
(6.8)

1472

Let $u_0 = \frac{1}{\lambda_0} \tilde{Q}_{b_0}(\frac{x}{\lambda_0}) + \epsilon_0$ as described in Subsection 2.1. Note Lemma 1.7 and Lemma 2.2 hold for those data. We need to choose K_2 large enough for later use. We also require u_0 be radial. The associated solution will have the geometric decomposition

$$u(t,x) = \frac{1}{\lambda(t)} (\tilde{Q}_b(t) + \epsilon(t)) (\frac{x}{\lambda(t)}) e^{-i\gamma(t)}$$
(6.9)

such that the conclusion of Lemma 1.7 and Lemma 2.2. And in the spirit of Remark 1.15, we may further sharpen the initial condition such that u further satisfies

$$\int |\nabla \epsilon|^2 + |\epsilon|^2 e^{-|y|} \le \Gamma_b^{1-a_2}, \quad e^{-e^{\frac{(1+a_2)\pi}{b}}} \le \lambda(t) \le e^{-e^{\frac{(1-a_2)\pi}{b}}}.$$

Recall a_2 is the parameter we have fixed at the beginning of this Section.

Note λ , b, γ , $\epsilon(t)$ depend on t in an continuous way. u(t, x) can be understood as a family of data continuously depending on t. By (1.69) and (1.67), we obtain $-\frac{\lambda_s}{\lambda} \ge \frac{1}{2}b > 0$, since $\lambda_t = \lambda^2 \lambda_s < 0$. Thus, the map from $t \to \lambda(t)$ is a homeomorphism. Thus u(t, x) in (6.9) can also be understood as a family of data indexed by λ . To summarize, we have a family of data index by λ small enough,

$$u_{\lambda}(x) = \frac{1}{\lambda} (\tilde{Q}_{b(\lambda)} + \epsilon(\lambda)) (\frac{x}{\lambda}) e^{-i\gamma(\lambda)}, \tag{6.10}$$

with

$$\begin{split} \lambda + b &< \alpha/10, \quad b > 0, \quad e^{e^{-\frac{(1+a_2)\pi}{b}}} < \lambda < e^{e^{-\frac{(1-a_2)\pi}{b}}}, \\ \int & \|\nabla\epsilon(t)\|^2 + |\epsilon(t)|^2 \le \Gamma_{b(t)}^{1-a_2}, \\ |x| \le \frac{10}{\lambda(t)} \\ \lambda^2 |\int \chi_{0,loc}(x - x(t)) \left(\frac{1}{2} |\nabla u(t)|^2 - \frac{1}{4} |u|^4\right)| \le \Gamma_{b(t)}^{10000}, \\ \lambda(t) |\Im \int \chi_{0,loc}(x - x(t)) \left(\nabla u(t)\bar{u}(t)\right)| \le \Gamma_{b(t)}^{10000}, \|u(t)\|_{H^{N_2}(|x - x_{1,0}| \ge \frac{1}{10000})} \le \alpha/10, \end{split}$$
(6.11)

and b, γ , ϵ depending on λ continuously. Recall $\chi_{0,loc}$ is defined in (3.1).

We will consider a family of data

$$u_{0,\lambda,x} \equiv u_{\lambda_{1,0},\lambda_{2,0},\dots\lambda_{m,0},x_{1,0},x_{2,0},\dots x_{m,0}}$$

$$\equiv \sum_{j=1}^{m} \chi(x - x_{j,0}) u_{\lambda_{j,0}}(x - x_{j,0}).$$
(6.12)

Here u_{λ} is defined as in (6.10). And we also require $|x_{j,0} - x_{j,\infty}| \le 1$. Note this implies $|x_{j,0} - x_{j',0}| \ge 10$, $j \ne j'$, since we have (6.1).

Now let us consider (1.1) with initial data $u_{0,\lambda,x}$. Lemma 3.1 will work for $u_{0,\lambda,x}$. Thus the associated solution $u_{\lambda,x}$ will satisfy the geometric decomposition in its lifespan

$$u_{\lambda,x}(t,x) = \frac{1}{\lambda_{j,\lambda,x}} \tilde{Q}_{b_{j,\lambda,x}}(\frac{x-x_{j,\lambda,x}}{\lambda_{j,\lambda,x}}) e^{-i\gamma_{j,\lambda,x}} + \Xi_{\lambda,x}$$
(6.13)

and all the bootstrap estimates in Lemma 3.1 holds. For notation convenience, we will write $u_{\lambda,x}$ as u, write $\lambda_{j,\lambda,x}$ as λ_j , write $x_{j,\lambda,x}$ as x_j , write $\gamma_{j,\lambda,x}$ as γ_j .

Again, in the spirit of Remark 1.15, we can further sharpen the condition on (6.11), i.e. make a_2 small enough, such that

$$e^{-e^{\frac{(1+a_1)\pi}{b_j}}} \le \lambda_j \le e^{-e^{\frac{(1-a_1)\pi}{b_j}}}.$$
 (6.14)

6.2. log-log blow up and almost sharp blow up dynamic

We first show the solution u with initial data $u_{\lambda,x}$ blows up according to the log–log law, indeed, we show every bubble itself evolves according to the log–log law, and for later use, we need the almost sharp dynamic, and keep in mind condition our data has already been sharpened so that (6.14) holds for the associated solution.

Lemma 6.1. Let u be the solution with initial data $u_{\lambda,x}$ as in (6.12), then for each $j = 1, \ldots m$, there is a T_j such that

$$\lambda_{i\,0}^2 \ln|\ln\lambda_{i,0}| = 2\pi \left(1 + O(a_1)\right) T_i,\tag{6.15}$$

$$\lambda_j(t)^2 \ln|\ln\lambda_j(t)| = 2\pi \left(1 + O(a_1)\right) (T_j - t).$$
(6.16)

In particular, since the blow up rate is modeled by $\min_i \lambda_i$, we have that the solution blow up in finite time T^+ and

$$\|\nabla u(t)\|_{2} \sim \sqrt{\frac{\ln|\ln(T^{+}-t)|}{T^{+}-t}}.$$
(6.17)

Remark 6.1. We implicitly require α to be small enough, as the whole paper.

Remark 6.2. Note Lemma 3.1 says the *m* bubbles evolve according to the log–log law without really seeing each other. T_j is the time that the *j*th bubble which is supposed to blow up. The solution will blow up at $T^+ = \min_j T_j$, and the dynamic will be stopped at T^+ . In particular, if $T_j > T_{j'}$, then it means *j* th bubble 'blows' up faster than *j*th bubble, though they may both not blow up.

Remark 6.3. The proof of Lemma 6.1 needs to use the bootstrap estimate rather than bootstrap hypothesis in Lemma 3.1 since we need to get control in term of $O(a_1)$. It is not hard to see one can argue as the proof below, and with bootstrap hypothesis rather than bootstrap estimate to show

$$\lambda_{j,0}^2 \ln|\ln\lambda_{j,0}| \sim T_j,\tag{6.18}$$

$$\lambda_j(t)^2 \ln|\ln\lambda_j(t)| \sim (T_j - t) \tag{6.19}$$

which in particular shows the associated solution *u* blows up in finite time $T^+ < \delta(\alpha)$.

Proof of Lemma 6.1. We will follow the computation in [20], which is used to show the exact log–log law.

$$\left|\frac{1}{\lambda_j}\frac{d\lambda_j}{ds_j} + b_j\right| \le \Gamma_{b_j}^{\frac{1}{2}-C\eta},\tag{6.20}$$

which immediately implies

$$(1 - \delta(\alpha))b_j \le -\frac{1}{\lambda_j} \frac{d\lambda_j}{ds_j} \le (1 + \delta(\alpha))b_j.$$
(6.21)

Note also (6.14) implies

$$b_j = \frac{(1+O(a_1))\pi}{\ln|\ln\lambda_j|}.$$
(6.22)

Thus

$$\frac{d}{dt}\lambda_j^2 \ln|\ln\lambda_j| = 2\frac{d\lambda_j}{\lambda_j^2 ds}\lambda_j \ln|\ln\lambda_j|(1+\delta(\alpha)) = 2(1+\delta(\alpha))\frac{d\lambda_j}{\lambda_j ds_j}\ln|\ln\lambda_j|.$$
(6.23)

Now, plug in (6.21) and (6.22), and choosing α small enough, we get

$$\frac{d}{dt}\lambda_j^2 \ln|\ln\lambda_j| = -2\pi \left(1 + O(a_1)\right).$$
(6.24)

which immediately implies (6.15), (6.16).

6.3. A quick discussion of blow up at same time

Now, we have a family of data $u_{\lambda,x}$, and the strategy is to adjust parameters to make the *m* bubbles to blow up at the same time at the prescribed position. If one only wants to make the *m* bubble to blow up at same time and does not track the final blow up points, then a topological argument similar to the one in [4] will be enough. We quickly illustrate this. We will fix $x_{1,0}, \ldots x_{m,0}$ and $\lambda_{1,0}$. And we will adjust $\lambda_{2,0}, \ldots \lambda_{m,0}$ to make the *m* bubbles blow up simultaneously.

Lemma 6.2. Fixed $x_{1,0}, \ldots x_{m,0}$ and $\lambda_{1,0} (|x_{j,0} - x_{j',0}| \ge 10, j \ne j')$, there exist $(\beta_{2,0}, \ldots, \beta_{m,0}) \in [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$ such that, the associated solution u to (1.1) with initial data $u_{\lambda_{1,0},\beta_{2,0},\ldots,\beta_{m,0},x_{1,0},\ldots,x_{m,0}}$ as in (6.9), will blow up simultaneously at m points, i.e.

$$\lambda_j(t) \sim \lambda_{j'}(t), t \in [0, T^+). \tag{6.25}$$

For notation convenience, we write $u_{\lambda_{1,0},\beta_{2,0},\dots,\beta_{m,0},x_{1,0},\dots,x_{m,0}}$ as $u_{\beta_2,\dots,\beta_m}$, and we further write $u_{\beta_{2,0},\dots,\beta_{m,0}}$ as u_{β} . We rewrite (6.13) as

$$u_{\boldsymbol{\beta}}(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_{j,\boldsymbol{\beta}}(t)} \tilde{Q}_{b_{j,\boldsymbol{\beta}}(t)} (\frac{x - x_{j,\boldsymbol{\beta}}(t)}{\lambda_{j,\boldsymbol{\beta}}}) e^{-i\gamma_{j,\boldsymbol{\beta}}(t)} + \Xi_{\boldsymbol{\beta}}(t,x).$$
(6.26)

Now we prove Lemma 6.2 by contradiction. Assume (6.25) is not true for any $\beta \in [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$. We consider several maps as following:

- Let $F(t, \boldsymbol{\beta}) \equiv F(t, \beta_2, \dots, \beta_m) \equiv (\frac{\lambda_{2,\boldsymbol{\beta}}(t)}{\lambda_{1,\boldsymbol{\beta}}(t)}, \frac{\lambda_{3,\boldsymbol{\beta}}(t)}{\lambda_{1,\boldsymbol{\beta}}(t)}, \dots, \frac{\lambda_{m,\boldsymbol{\beta}}(t)}{\lambda_{1,\boldsymbol{\beta}}}).$ - Let $T_{\boldsymbol{\beta}}$ be the first time $F(t, \boldsymbol{\beta})$ hits $\partial [(1 - a_0), (1 + a_0)]^{m-1}.$ - Let $G(\boldsymbol{\beta})$ be $F(T_{\boldsymbol{\beta}}, \boldsymbol{\beta}).$

Here $\boldsymbol{\beta} \in [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$.

Since we assume (6.25) is not true for any $\boldsymbol{\beta} \in [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$, then $T_{\boldsymbol{\beta}}$ is always well defined, i.e. $T_{\boldsymbol{\beta}} < \infty$. The key point here is that $T_{\boldsymbol{\beta}}$ depends on $\boldsymbol{\beta}$ continuously. Assume this for the moment and let us finish the proof by deriving a contradiction. Note that F is clearly continuous (by standard well posedness theory of NLS). Since we assume that $T_{\boldsymbol{\beta}}$ depends on $\boldsymbol{\beta}$ continuously, G is also continuous. Make the observation $T_{\boldsymbol{\beta}} = 0$ for $\boldsymbol{\beta} \in \partial [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$, and then it is easy to see $G|_{\partial [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}}$ is an homeomorphism from $\partial [(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$ to $\partial [(1 - a_0), (1 + a_0)]^{m-1}$. Then we have constructed a continuous map from $[(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$ to $\partial [(1 - a_0), (1 + a_0)]^{m-1}$ such that its restriction on $[(1 - a_0)\lambda_{1,0}, (1 + a_0)\lambda_{1,0}]^{m-1}$ is a homeomorphism, which is clearly false by classical Homology theory. A contradiction!

We need to check T_{β} does depend on β continuously. Note that by LWP of NLS, the map $F(t, \beta)$ is continuous and differentiable. Thus, to show that T_{β} depends on β continuously, we need only to show that $\partial_t F(T_{\beta}, \beta)$ points outward $\partial [1 - a_0, 1 + a_0]^{m-1}$. In order to show this, without loss of generality, we assume $\frac{\lambda_{2,\beta}}{\lambda_{1,\beta}}(T_{\beta}) = 1 + a_0$ and check

$$\frac{d}{dt}(\frac{\lambda_{2,\boldsymbol{\beta}}}{\lambda_{1,\boldsymbol{\beta}}})(T_{\boldsymbol{\beta}}) > 0.$$

This clearly follows from the fact that

$$\frac{1}{\lambda_{2,\boldsymbol{\beta}}}\frac{d}{dt}\lambda_{2,\boldsymbol{\beta}}-\frac{1}{\lambda_{1,\boldsymbol{\beta}}}\frac{d}{dt}\lambda_{1,\boldsymbol{\beta}}>0,$$

which is equivalent to

$$-\frac{1}{(\lambda_{1,\boldsymbol{\beta}})^3}\frac{d}{ds_1}\lambda_{1,\boldsymbol{\beta}} > -\frac{1}{(\lambda_{2,\boldsymbol{\beta}})^3}\frac{d}{ds_2}\lambda_{2,\boldsymbol{\beta}},\tag{6.27}$$

when $t = T_{\beta}$. Note that by (6.21) and (6.22) (again, we will choose α small enough),

$$-\frac{1}{(\lambda_{j,\boldsymbol{\beta}})^3}\frac{d}{ds_j}\lambda_{j,\boldsymbol{\beta}} = (1+O(a_1))\frac{1}{(\lambda_{j,\boldsymbol{\beta}})^2}\frac{1}{\ln|\ln\lambda_{j,\boldsymbol{\beta}}|}.$$
(6.28)

Note $\frac{\lambda_{2,\beta}}{\lambda_{1,\beta}}(T_{\beta}) = 1 + a_0$, and $a_0 \gg a_1$, (6.27) follows once one plugs in (6.28). This concludes the proof of Lemma 6.2.

6.4. Prescription of blow up points

To prescribe the blow up points, i.e., to make the solution blows up exactly at the given points $x_{1,\infty}, \ldots x_{m,\infty}$ is more tricky. We will need a topological argument inspired by [14], see also [23]. Morally speaking, we are dealing with a family of initial data with parameters $\lambda_{1,0}, \lambda_{2,0}, \ldots \lambda_{m,0}, x_{1,0}, \ldots x_{m,0}$, and the goal is to adjust those parameters to make the solution blow up at *m* points and the *m* points should be the given $x_{1,\infty}, \ldots x_{m,\infty}$. The analysis in Subsection 6.3 basically says for any given initial parameter $x_{1,0}, \ldots x_{m,0}$, one will be able to find $\lambda_{1,0}, \ldots \lambda_{m,0}$ to make the solution blow up at *m* points. If one can choose $\lambda_{1,0}, \ldots \lambda_{m,0}$ according to $x_{1,0}, \ldots x_{m,0}$ in a continuous way, then the argument in [14] will be able to help us adjust $x_{1,0}, \ldots x_{m,0}$ to make the *m* blow up points be exactly the prescribed $x_{1,\infty}, \ldots x_{m,\infty}$. However, with our arguments in Subsection 6.3, the choice of $\lambda_{1,0}, \ldots \lambda_{m,0}$ is not even uniquely determined by the $x_{1,0}, \ldots x_{m,0}$. Though this can be somehow fixed, it is very unclear whether one can choose $\lambda_{1,0}, \ldots \lambda_{m,0}$ in a continuous way¹³ according to $x_{1,0}, \ldots x_{m,0}$. Before we continue, let us first make several important observations

- The sharp dynamic of log-log blow up is known, we should make full use of it.
- The impact of parameters of $x_{1,0}, \ldots x_{m,0}$ is of lower order than $\lambda_{1,0}, \ldots \lambda_{m,0}$.

Our strategy is to choose all the parameters $x_{1,0}, \ldots x_{m,0}, \lambda_{1,0}, \ldots \lambda_{m,0}$ simultaneously to make the solutions blow up at exactly m points $x_{1,\infty}, \ldots x_{m,\infty}$. Finally, at the technique level, in [14] and [23], they rely on the following topological lemma (they call it index Theorem).

Lemma 6.3. Let f be a continuous map from \mathbb{R}^n to \mathbb{R}^n , let r > 0 and suppose

$$|f(y) - y| < |y|, \forall y \in \partial B_r$$

Then there is $y_0 \in B_r$ such $f(y_0) = 0$.

We will need a modified version

Lemma 6.4. Let f be a map from $\Omega \subset \mathbb{R}^n$ to \mathbb{R}^n . Let Ω be a convex domain and let $\partial \Omega$ be a closed surface which is homeomorphic to the sphere. We assume the original point is in the Ω . Let us further assume for each y in $\partial \Omega$, we have

$$\mathbf{0} \notin \{(1-t)y + tf(y) | t \in [0,1]\}.$$
(6.30)

(6.29)

Then $\mathbf{0} \in f(\bar{\Omega})$.

We will prove Lemma 6.4 in Appendix B, we point out Lemma Appendix B actually implies Lemma 6.3. Now, we turn to the prescription of blow up points. As previous mentioned in (6.1), we assume

$$|x_{j,\infty} - x_{j',\infty}| \ge 20, \, j \ne j'. \tag{6.31}$$

We will still consider the data as in (6.9) and we will fix $\lambda_{1,0}$ and adjust parameters $\lambda_{2,0}, \ldots, \lambda_{m,0}, x_{1,0}, \ldots, x_{m,0}$ to make the *m* bubbles blow up at the same time in $x_{1,\infty}, \ldots, x_{m,\infty}$.

1476

¹³ If one wants to direct borrow the arguments in [14], one will need an maximal principle type argument, which says the following: Fixed $x_{1,0}, \ldots x_{m,0}$, let $\lambda_{1,0}, \ldots \lambda_{m,0}$ be chosen such that the *m* bubbles blow up at the same time, then if one further adjusts $\lambda_{1,0}$ to be smaller and keep other parameters unchanged, then the first bubble will blow up first. This is not clear in our setting, and we even think this argument may not hold for log–log blow up solutions.

Lemma 6.5. Fix $\lambda_{1,0}$, there exists $(\beta_2, \ldots, \beta_m, d_1, \ldots, d_m) \in [-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m$ (here $B_1 \subset \mathbb{R}^2$ is the unit ball) such that the associated solution u to (1.1) with initial data

$$u_{\lambda,x} := u_{\lambda_{1,0},\lambda_{1,0}+\beta_{2,..\lambda_{1,0}+\beta_{m,0},x_{1,0}+d_{1,...,x_{m,0}+d_{m}}}$$

will blow up at m given prescribed points $x_{1,\infty}, ..x_{m,\infty}$, i.e.

$$\lim_{t \to T^+(u)} \lambda_j(0) = 0, \quad j = 1, \dots, m,$$

$$\lim_{t \to T^+(u)} x_j(t) = x_{j,0}, \quad j = 1, \dots, m.$$
(6.32)

Here, we use λ_j, x_j to denote $\lambda_{j,\lambda,x}, x_{j,\lambda,x}$ for notation convenience.

Proof of Lemma 6.5. As we previously did in the proof of Lemma 6.2, we write $u_{\lambda_{1,0},\lambda_{1,0}+\beta_2,..\lambda_{1,0}+\beta_{m,0},x_{1,0}+d_1,...x_{m,0}+d_0}$ as $u_{\beta_2,...\beta_m,d_1,...d_m}$. And we further write $u_{\beta_2,...\beta_m,d_1,...d_m}$ as $u_{\mathcal{A}}$, where $\mathcal{A} = (\beta_2,...\beta_m,d_1,...d_m) \in \mathbb{R}^{m-1} \times \mathbb{R}^{2m}$. We rewrite (6.13) as

$$u_{\mathcal{A}}(t,x) = \sum_{j=1}^{m} \frac{1}{\lambda_{j,\mathcal{A}}(t)} \tilde{Q}_{b_{j,\mathcal{A}}(t)} (\frac{x - x_{j,\mathcal{A}}(t)}{\lambda_{j,\mathcal{A}}}) e^{i\gamma_{j,\mathcal{A}}(t)} + \Xi_{\mathcal{A}}(t,x).$$
(6.33)

Let $T_{\mathcal{A}}$ be the blow up time of $u_{\mathcal{A}}$. We now consider the following map:

$$\mathcal{F}: [-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m \to \mathbb{R}^{m-1} \times \mathbb{R}^{2m}$$

$$\mathcal{F}(\mathcal{A}) := (y_{2,\mathcal{A}}, \dots, y_{m,\mathcal{A}}, z_{1,\mathcal{A}}, \dots, z_{m,\mathcal{A}})$$

$$y_{i,\mathcal{A}} = \lambda_{i,\mathcal{A}}(T_{\mathcal{A}}) - \lambda_{1,\mathcal{A}}(T_{\mathcal{A}}), z_{j,\mathcal{A}} = x_{j,\mathcal{A}}(T_{\mathcal{A}}) - x_{j,\infty}, i = 2, \dots, m, j = 1, \dots, m.$$
(6.34)

Here $B_1 \subset \mathbb{R}^2$ is the unit ball. And $\lambda_{j,\mathcal{A}}(T_{\mathcal{A}})$ and $x_{j,\mathcal{A}}(T_{\mathcal{A}})$ are defined as

$$\lambda_{j,\mathcal{A}}(T_{\mathcal{A}}) = \lim_{t \to T_{\mathcal{A}}} \lambda_{j,\mathcal{A}}(t), \ j = 1, \dots, m,$$

$$x_{j,\mathcal{A}}(T_{\mathcal{A}}) = \lim_{t \to T_{\mathcal{A}}} x_{j,\mathcal{A}}(t), \ j = 1, \dots, m.$$

(6.35)

 $\lambda_{j,\mathcal{A}}(T_{\mathcal{A}})$ is well defined as we have (6.21), which implies $\lambda_{j,\mathcal{A}}$ is strictly decreasing. $x_{j,\mathcal{A}}(T_{\mathcal{A}})$ is well defined as mentioned in Subsection 6.3, see also Remark 5.1. The point is

Lemma 6.6. The map $\mathcal{A} \to \lambda_{j,\mathcal{A}}(T_{\mathcal{A}})$ and the map $\mathcal{A} \to x_{j,\mathcal{A}}(T_{\mathcal{A}})$ is continuous.

We will prove Lemma 6.6 in Appendix C.

Note if $\mathcal{F}(\mathcal{A}) = 0$ for some \mathcal{A} , then $u_{\mathcal{A}}$ is the desired solution which blows up according to log–log law at exactly *m* prescribed points.

Lemma 6.6 implies the map \mathcal{F} is continuous and we will use Lemma 6.4 to show

$$\mathbf{0} \in \mathcal{F}([-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m).$$
(6.36)

To achieve this, we need to show if \mathcal{A} in $\partial \{[-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m\}$, then

$$\mathbf{0} \notin \{t\mathcal{A} + (1-t)\mathcal{F}(\mathcal{A}), t \in [0,1]\}.$$
(6.37)

Note if $\mathcal{A} \in \partial\{[-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m\}$, then at least one of the following holds (recall the notation $\mathcal{A} = (\beta_2, \dots, \beta_m)$)

- Case 1: $|\beta_j| = a_0 \lambda_{1,0}$, for some j = 2, ..., m.
- Case 2: $|d_j| = 1$, for some j = 1, ..., m.

We first show (6.37) holds in Case 2. Indeed, by bootstrap estimate (3.53) in Lemma 3.1, we have $\sup_{t < T_A} |x_{j,A}(t) - x_{j,0} - d_j| \le \frac{1}{2000}$, which implies $|d_j - z_{j,A} \le \frac{1}{2000}|$. (Recall our notation $\mathcal{F}(\mathcal{A}) := (y_{2,A}, \dots, y_{m,A}, z_{1,A}, \dots, z_{m,A})$.) Since $|d_j| = 1$, this implies

 $td_j + (1-t)z_{j,\mathcal{A}} \neq \mathbf{0}, \forall \in [0,1]$

In particular (6.37) holds in Case 2.

Next we show (6.37) also holds in Case 1. Without loss of generality, we assume $\beta_2 = (a_0)\lambda_{1,0}$. Using Lemma 6.1, we can find $T_{1,\mathcal{A}}$, $T_{2,\mathcal{A}}$ such that

$$\lambda_{j,\mathcal{A}}^{2}(t)\ln|\ln\lambda_{j,\mathcal{A}}(t)| = 2\pi(1+O(a_{1}))(T_{j,\mathcal{A}}-t), t < T_{\mathcal{A}},$$

$$T_{1,\mathcal{A}} = 2\pi(1+O(a_{1}))(\lambda_{1,0}^{2}\ln|\ln\lambda_{1,0}|),$$

$$T_{2,\mathcal{A}} = 2\pi(1+O(a_{1}))\left((1+a_{0})^{2}\lambda_{1,0}^{2}\ln|\ln(1+a_{0})\lambda_{1,0})|\right)$$
(6.38)

(Note this also implies $T_{\mathcal{A}} \leq \min(T_{1,\mathcal{A}}, T_{2,\mathcal{A}})$, and since $a_1 \ll a_0, T_{1,\mathcal{A}} < T_{2,\mathcal{A}}$.)

Since $a_0 \gg a_1$, it is easy to see

$$\inf_{t \le T_{1,\mathcal{A}}} \left(\lambda_{1,\mathcal{A}}^2(t) \ln |\ln \lambda_{1,\mathcal{A}}(t)| - \lambda_{2,\mathcal{A}}^2(t) \ln |\ln \lambda_{2,\mathcal{A}}(t)| \right) < 0$$
(6.39)

(Note it is important here we have < rather than \leq .)

In particular, we have $\lambda_{2,\mathcal{A}}(T_{\mathcal{A}}) > \lambda_{1,\mathcal{A}}(T_{\mathcal{A}})$ (since $T_{\mathcal{A}} \leq T_{1,\mathcal{A}}$). Thus, since $\beta_2 > 0$, and $y_{2,\mathcal{A}} = \lambda_{2,\mathcal{A}}(T_{\mathcal{A}}) - \lambda_{1,\mathcal{A}}(T_{\mathcal{A}}) > 0$,

 $t\beta_2 + (1-t)y_{2,\mathcal{A}} \neq 0, t \in [0, 1].$

which recovers (6.37). This concludes the proof. \Box

Conflict of interest statement

There is no conflict of interest.

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Appendix A. The local wellposedness of the modified system

We explain briefly why one can always locally solve (2.3), (3.27). We only explain the case about (2.3) here, (3.27) is similar. Indeed, the system is NLS couples with 4 ODEs. Since NLS is locally well posed and ODE is always locally well posed, it is no surprise (2.3) is locally well posed. To construct a solution, one first solves NLS $iu_t = -\Delta u - |u|^2 u$ in a time interval [0, T_1], and plugs this u(t, x) (which is not unknown in [0, T_1]) into the last 4 equations, we will obtain 4 ODEs about { $\lambda(t), b(t), x(t), \gamma(t)$ }. We just do some computation to illustrate this, say, the equation

$$\frac{d}{dt}\{(\epsilon_1(t), |y|^2 \Sigma_{b(t)}) + (\epsilon_2(t), |y|^2 \Theta_{b(t)})\} = 0$$

is now equivalent to

$$\frac{d}{dt}\Re(\frac{1}{\lambda(t)}\epsilon(\frac{x-x(t)}{\lambda(t)})e^{-i\gamma(t)},\frac{1}{\lambda(t)}(|y|^2\tilde{Q}_b)(\frac{x-x(t)}{\lambda(t)})e^{-i\gamma(t)})=0,$$

which is equivalent to

$$\frac{d}{dt}\Re(u-\frac{1}{\lambda(t)}\tilde{Q}_b(\frac{x-x(t)}{\lambda(t)}),\frac{1}{\lambda(t)}(|y|^2\tilde{Q}_b)(\frac{x-x(t)}{\lambda(t)}))e^{-i\gamma(t)})=0,$$

which is equivalent to

$$\begin{aligned} \Re(i\Delta u + i|u|^2, \frac{1}{\lambda(t)}(|y|^2\tilde{Q}_b)(\frac{x - x(t)}{\lambda(t)})e^{-i\gamma(t)}) \\ + \Re(u, \frac{d}{dt}\frac{1}{\lambda(t)}(|y|^2\tilde{Q}_b)(\frac{x - x(t)}{\lambda(t)})e^{-i\gamma(t)}) \\ - \frac{d}{dt}\Re(\frac{1}{\lambda(t)}\tilde{Q}_b(\frac{x - x(t)}{\lambda(t)}), \frac{1}{\lambda(t)}(|y|^2\tilde{Q}_b)(\frac{x - x(t)}{\lambda(t)})e^{-i\gamma(t)}) \\ = 0. \end{aligned}$$
(A.1)

Though (A.1) is complicate, it is an ODE about $\{\lambda(t), b(t), x(t), \gamma(t)\}$.

Similarly, the last 3 equations in (2.3) can also be transformed into ODEs about $\{\lambda(t), b(t), x(t), \gamma(t)\}$.

Thus the local well posedness theory about (2.3) is equivalent to the local well posedness theory about NLS.

Appendix B. Proof of Lemma 6.4

Let us turn to the proof of Lemma 6.4 now. This is very standard in algebraic topology. Note $\overline{\Omega} - \{0\}$ is the retract of $\partial\Omega$, i.e. there is a map

$$r: \mathbb{R}^n - \{\mathbf{0}\} \to \partial \Omega,$$

such that

$$r \circ \iota = i d_{\partial \Omega},$$

here $\iota: \partial \Omega \to \mathbb{R}^n - \{\mathbf{0}\}$ is the natural inclusion map.

Now we prove Lemma 6.4 by contradiction. Assume $\mathbf{0} \notin f(\bar{\Omega})$. Then $g := r \circ f$ is well defined and continuous. Note g is map from $\bar{\Omega}$ to $\partial \Omega$.

On the other hand $g|_{\partial \Omega}$ is homotopic to the $id|_{\partial \Omega}$. Indeed, we may write down the homotopy explicitly

 $r \circ (tf + (1-t)id).$

We emphasize here this homotopy is well defined since $tf(y) + (1-t)y \neq 0$ for $t \in [0, 1]$ and $y \in \partial \Omega$.

Now we have constructed a map g from $\overline{\Omega}$ to $\partial \Omega$, and $g|_{\partial \Omega}$ is homotopic to $id|_{\partial|_{\Omega}}$. Note Ω is convex domain and $\partial \Omega$ is homeomorphic to the sphere. This is a clear contradiction from standard homology theory.

Appendix C. Proof of Lemma 6.6

Before we go to the proof, let us point out Lemma 6.6 basically say the blow up point (model by $x_{j,A}(t)$) and the blow up time (modeled by $\lambda_{j,A}(t)$) depending on the initial data (modeled by A) in a continuous way. We remark here, in general, the problem whether blow up point and blow up time depend on the initial data in a continuous way is not an easy problem. Indeed, if one has a sequence initial data $u_{0,n}$, whose associated solution to (1.1) blows up according to the log–log law, and one assumes $u_{0,n}$ converges to u_0 in H^1 , it is not always right the associated solution u to (1.1) will blow up according to the log–log law. And for NLS, if we don't have some information about the dynamic near the blow up time, we cannot even define the blow up point. However, see [17] for this direction.

The proof of Lemma 6.6 is much easier, because we are only working on data with finite parameter A. And if $A_n \rightarrow A$, clearly u_A still blows up according the log–log law with dynamic described by Lemma 3.1. Let us turn to the proof. We recall

$$\mathcal{A} \in [-a_0\lambda_{1,0}, a_0\lambda_{1,0}]^{m-1} \times (B_1)^m.$$

(B.1)

Proof. Recall to understand the evolution of $x_{j,\mathcal{A}}(t)$, $\lambda_{j,\mathcal{A}}(t)$, one needs to consider system (3.27) or equivalently (3.38). As explained in Appendix A, the system is NLS coupled with 4 ODEs. Use the standard stability arguments for NLS and stability argument of ODEs, we have that for any $T < T_{\mathcal{A}}$ (recall $T_{\mathcal{A}}$ is the blow up time of $u_{\mathcal{A}}$) the map $\mathcal{A} \rightarrow x_{j,\mathcal{A}}(T)$, $\lambda_{j,\mathcal{A}}(T)$ is continuous. Now, Lemma 6.6 easily follows from the following lemma

Lemma 6.7. Given A, for any $\epsilon > 0$, there is $T < T_A$ and $\delta = \delta(A)$ such that for any A' with

$$|\lambda_{j,\mathcal{A}'}(T) - \lambda_{j,\mathcal{A}}(T)| < \delta, \tag{C.1}$$

then

$$\sup_{t \in [T, T'_A]} |\lambda_{j, \mathcal{A}'}(t) - \lambda_{j, \mathcal{A}'}(T)| < \epsilon, j = 1, ..., m,$$
(C.2)

$$\sup_{t \in [T, T_{\mathcal{A}'}]} |x_{j, \mathcal{A}'}(t) - x_{j, \mathcal{A}'}(T)| < \epsilon, j = 1, \dots m.$$
(C.3)

We now prove Lemma 6.7.

Note A is given, and we only need to prove (C.2) and (C.3) for every given j. We discuss the two cases.

- Case 1: the *j*th bubble of $u_{\mathcal{A}}$ blows up, i.e. $\lambda_{j,\mathcal{A}}(T_{\mathcal{A}}) = 0$.
- Case 2: the *j*th bubble of u_A blows up, i.e. $\lambda_{j,A}(T_A) \neq 0$.

We first discuss Case 1. In this case, we can choose T close to T_A enough, and δ small enough, such that $\lambda_{j,A'}(T)$ is small enough,

$$\lambda_{i,\mathcal{A}'}(T) \le \epsilon_0 \ll \epsilon^{100}. \tag{C.4}$$

Then (C.2) follows since $\frac{d}{dt}\lambda_{j,\mathcal{A}'} < 0$ and $\lambda_{j,\mathcal{A}'}(t) > 0$, $\forall t < T_{\mathcal{A}'}$. To prove (C.3), one needs modify a little bit the analysis in Subsubsection 5.4.5, as argued in Subsubsection 5.4.5.

$$|x_{j,\mathcal{A}'}(t) - x_{j,\mathcal{A}'}(T)| \le \int_{T}^{t} \frac{dx_j}{ds_j} \frac{1}{\lambda_j^2} d\tau \le \int_{T}^{T_{\mathcal{A}'}} \frac{1}{\lambda_{j,\mathcal{A}}}(\tau) d\tau$$
(C.5)

(We have use $\left|\frac{1}{\lambda_{j,\mathcal{A}'}}\frac{dx_{j,\mathcal{A}'}}{ds_j}\right| \lesssim 1$ by (5.6).) And recall (5.45), we further have

$$|x_{j,\mathcal{A}'}(t) - x_{j,\mathcal{A}'}(T)| \lesssim \lambda_{j,\mathcal{A}'}^{1/2}(T) \int_{T}^{T_{\mathcal{A}'}} \frac{1}{\lambda_{j,\mathcal{A}}^{3/2}} \lesssim \lambda_{j,\mathcal{A}'}^{1/2}(T) \ll \epsilon^{50}.$$
(C.6)

This gives (C.3).

Now we discuss Case 2. Without loss of generality, we assume $\lambda_{j,\mathcal{A}'}(T) \ge \epsilon_0 \ge \epsilon^{100}$, otherwise we just argue as Case 1. Note since $u_{\mathcal{A}}$ blows up at $T_{\mathcal{A}}$, at least one bubble blows up at $T_{\mathcal{A}}$, let us assume $\lim_{t\to T_{\mathcal{A}}} \lambda_{j_0,\mathcal{A}} = 0$. Thus, we are able to choose T close enough to $T_{\mathcal{A}}$ and δ small enough such that

$$\lambda_{j,\mathcal{A}'}(T) \gg \lambda_{j_0}(T), \tag{C.7}$$

and of course $\lambda_{j_0, \mathcal{A}'}(T) \ll \epsilon^{100}$.

We remark here we actually may need

$$\lambda_{j_0,\mathcal{A}'}(T) \ll \exp(-\exp\exp\exp\{\frac{1}{\lambda_{j,\mathcal{A}'}(T)}\}),\tag{C.8}$$

and don't worry about the special form because our arguments are kinds of soft.

The idea is $u_{\mathcal{A}}$ is going to blow up so fast that the *j*th bubble does not have much time to change.

Note it is not clear whether the j_0 th bubble will blow up when $u_{A'}$ blows up, however, we can still estimate blow up time by Lemma 6.1, i.e. the blow up time $T_{A'}$ of $u_{A'}$ is controlled by $T_{j_0,A}$ predicted by Lemma 6.1, we have¹⁴

$$T_{\mathcal{A}'} - T \lesssim \lambda_{j_0,\mathcal{A}'}^2(T) \ln |\ln \lambda_{j_0,\mathcal{A}'}(T)|.$$
(C.9)

And, using Lemma 6.1 again, we can ensure

$$\lambda_{j,\mathcal{A}'}(t) \sim \lambda_{j,\mathcal{A}'}(T), t \in [T, T_{\mathcal{A}'}]. \tag{C.10}$$

Now we use estimate (we need (6.21), (6.22))

$$0 < -\frac{d}{dt}\lambda_{j,\mathcal{A}'} \sim -\frac{1}{\lambda_{j,\mathcal{A}'}^3} \frac{d\lambda_{j,\mathcal{A}'}}{ds_j} \sim \frac{1}{\lambda_{j,\mathcal{A}'}^3} b_{j,\mathcal{A}'} \sim \frac{1}{\lambda_{j,\mathcal{A}'}^3} \ln|\ln\lambda_{j,\mathcal{A}}|$$
(C.11)

Combine (C.11) and (C.10), plug in (C.9), we have

$$\sup_{t \in [T, T_{\mathcal{A}'})} |\lambda_{j, \mathcal{A}'}(T) - \lambda_{j, \mathcal{A}'}(t)| \le \lambda_{j, \mathcal{A}'}^{-3}(T) \ln |\ln \lambda_{j, \mathcal{A}}(T)| \lambda_{j_0, \mathcal{A}'}^2(T) \ln |\ln \lambda_{j_0, \mathcal{A}'}(T)|$$
(C.12)

The desired estimate (C.2) follows since we have (C.8).

Similar arguments work for (C.3). \Box

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¹⁴ Lemma 6.1 is stated for solutions started at t = 0, clearly we can also set the start time as t = T.

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