

# ANALYSIS & PDE

Volume 11

No. 4

2018

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## ON RANK-2 TODA SYSTEMS WITH ARBITRARY SINGULARITIES: LOCAL MASS AND NEW ESTIMATES

CHANG-SHOU LIN, JUN-CHENG WEI, WEN YANG AND LEI ZHANG

For all rank-2 Toda systems with an arbitrary singular source, we use a unified approach to prove:

(1) The pair of local masses  $(\sigma_1, \sigma_2)$  at each blowup point has the expression

$$\sigma_i = 2(N_{i1}\mu_1 + N_{i2}\mu_2 + N_{i3}),$$

where  $N_{ij} \in \mathbb{Z}$ ,  $i = 1, 2$ ,  $j = 1, 2, 3$ .

(2) At each vortex point  $p_t$  if  $(\alpha_t^1, \alpha_t^2)$  are integers and  $\rho_i \notin 4\pi\mathbb{N}$ , then all the solutions of Toda systems are uniformly bounded.

(3) If the blowup point  $q$  is a vortex point  $p_t$  and  $\alpha_t^1, \alpha_t^2$  and 1 are linearly independent over  $\mathcal{Q}$ , then

$$u^k(x) + 2 \log |x - p_t| \leq C.$$

The Harnack-type inequalities of 3 are important for studying the bubbling behavior near each blowup point.

### 1. Introduction

Let  $(M, g)$  be a Riemann surface without boundary and  $\mathbf{K} = (k_{ij})_{n \times n}$  be the Cartan matrix of a simple Lie algebra of rank  $n$ . For example, for the Lie algebra  $\mathfrak{sl}(n+1)$  (the so-called  $A_n$ ) we have

$$\mathbf{K} = \begin{pmatrix} 2 & -1 & 0 & \cdots & 0 \\ -1 & 2 & -1 & \cdots & 0 \\ \vdots & & & & \vdots \\ 0 & \cdots & -1 & 2 & -1 \\ 0 & \cdots & 0 & -1 & 2 \end{pmatrix}. \tag{1-1}$$

In this paper we consider the solution  $u = (u_1, \dots, u_n)$  of the following system defined on  $M$ :

$$\Delta_g u_i + \sum_{j=1}^n k_{ij} \rho_j \left( \frac{h_j e^{u_j}}{\int_M h_j e^{u_j} dV_g} - 1 \right) = \sum_{p_t \in S} 4\pi \alpha_t^i (\delta_{p_t} - 1), \tag{1-2}$$

where  $\Delta_g$  is the Laplace–Beltrami operator ( $-\Delta_g \geq 0$ ),  $S$  is a finite set on  $M$ ,  $h_1, \dots, h_n$  are positive and smooth functions on  $M$ ,  $\alpha_t^i > -1$  is the strength of the Dirac mass  $\delta_{p_t}$  and  $\rho = (\rho_1, \dots, \rho_n)$  is a constant vector with nonnegative components. Here for simplicity we just assume that the total area of  $M$  is 1.

*MSC2010:* primary 35J47; secondary 35J60, 35J55.

*Keywords:*  $SU(n+1)$ -Toda system, asymptotic analysis, a priori estimate, classification theorem, topological degree, blowup solutions, Riemann–Hurwitz theorem.

Obviously, (1-2) remains the same if  $u_i$  is replaced by  $u_i + c_i$  for any constant  $c_i$ . Thus we might assume that each component of  $u = (u_1, \dots, u_n)$  is in

$$\mathring{H}^1(M) := \{v \in L^2(M), \nabla v \in L^2(M) \text{ and } \int_M v \, dV_g = 0\}.$$

Then (1-2) is the Euler–Lagrange equation for the following nonlinear functional  $J_\rho(u)$  in  $\mathring{H}^1(M)$ :

$$J_\rho(u) = \frac{1}{2} \int_M \sum_{i,j=1}^n k^{ij} \nabla_g u_i \nabla_g u_j \, dV_g - \sum_{i=1}^n \rho_i \log \int_M h_i e^{u_i} \, dV_g,$$

where  $(k^{ij})_{n \times n} = \mathbf{K}^{-1}$ .

It is hard to overestimate the importance of system (1-2), as it covers a large number of equations and systems deeply rooted in geometry and physics. Even if (1-2) is reduced to a single equation with Dirac sources, it is a mean-field equation that describes metrics with conic singularities. Finding metrics with constant curvature with prescribed conic singularity is a classical problem in differential geometry and extensive references can be found in [Bartolucci and Tarantello 2002; Battaglia and Malchiodi 2014; Eremenko et al. 2014; Lin et al. 2012; 2015; Lin and Zhang 2010; 2013; 2016; Troyanov 1989; 1991; Yang 1997]. Recently profound relations among mean-field equations, the classical Lamé equation, hyperelliptic curves, modular forms and the Painlevé equation have been discovered and developed in [Chai et al. 2015; Chen et al. 2016].

The general form of (1-2) has close ties with algebraic geometry and integrable systems. Here we just briefly explain the relation between the  $sl(n+1)$ -Toda system and the holomorphic curves in projective spaces: Let  $f$  be a holomorphic curve from a domain  $D$  of  $\mathbb{R}^2$  into  $\mathbb{C}\mathbb{P}^n$ . Then  $f$  can be lifted locally to  $\mathbb{C}^{n+1}$  and we use  $v(z) = [v_0(z), \dots, v_n(z)]$  to denote the lift and  $f_k$  the  $k$ -th associated curve,

$$f_k : D \rightarrow G(k, n+1) \subset \mathbb{C}\mathbb{P}^n(\Lambda^k \mathbb{C}^{n+1}), \quad f_k(z) = [v(z) \wedge v'(z) \wedge \dots \wedge v^{(k-1)}(z)],$$

where  $v^{(j)}$  is the  $j$ -th derivative of  $v$  with respect to  $z$ . Let

$$\Lambda_k(z) = v(z) \wedge \dots \wedge v^{(k-1)}(z).$$

Then the well-known infinitesimal Plücker formula gives

$$\frac{\partial^2}{\partial z \partial \bar{z}} \log \|\Lambda_k(z)\|^2 = \frac{\|\Lambda_{k-1}(z)\|^2 \|\Lambda_{k+1}(z)\|^2}{\|\Lambda_k(z)\|^4} \quad \text{for } k = 1, 2, \dots, n, \tag{1-3}$$

where we put  $\|\Lambda_0(z)\|^2 = 1$  as convention and the norm  $\|\cdot\|^2 = \langle \cdot, \cdot \rangle$  is defined by the Fubini–Study metric in  $\mathbb{C}\mathbb{P}(\Lambda^k \mathbb{C}^{n+1})$ . Here we observe that (1-3) holds only for  $\|\Lambda_k(z)\| > 0$ , i.e., for all the unramified points  $z \in M$ . Now we set  $\|\Lambda_{n+1}(z)\| = 1$  by normalization (analytically extended at the ramification points) and

$$U_k(z) = -\log \|\Lambda_k(z)\|^2 + k(n-k+1) \log 2, \quad 1 \leq k \leq n.$$

For every ramified point  $p$  we use  $\{\gamma_{p,1}, \dots, \gamma_{p,n}\}$  to denote the total ramification index at  $p$  and set

$$u_i^* = \sum_{j=1}^n k_{ij} U_j, \quad \alpha_{p,i} = \sum_{j=1}^n k_{ij} \gamma_{p,j},$$

Then we have

$$\Delta u_i^* + \sum_{j=1}^n k_{ij} e^{u_j^*} - K_0 = 4\pi \sum_{p \in S} \alpha_{p,i} \delta_p, \quad i = 1, \dots, n, \tag{1-4}$$

where  $K_0$  is the Gaussian curvature of the metric  $g$ .

Therefore any holomorphic curve from  $M$  to  $\mathbb{C}\mathbb{P}^n$  is associated with a solution  $u^* = (u_1^*, \dots, u_n^*)$  of (1-4). Conversely, given any solution  $u^* = (u_1^*, \dots, u_n^*)$  of (1-4) in  $\mathbb{S}^2$ , it is possible to construct a holomorphic curve of  $\mathbb{S}^2$  into  $\mathbb{C}\mathbb{P}^n$  which has the given ramification index  $\gamma_{p,i}$  at  $p$  if  $\gamma_{p,i} \in \mathbb{N}$ . One can see [Lin et al. 2012] for the details of this construction. Therefore, (1-4) is related to the following problem in more general setting: given a set of ramified points on  $M$  and its ramification indices at these points, can we find holomorphic curves into  $\mathbb{C}\mathbb{P}^n$  that satisfy the given ramification information?

Equation (1-2) is also related to many physical models from gauge field theory. For example, to describe the physics of high critical temperature superconductivity, a model related to the Chern–Simons model was proposed, which can be reduced to an  $n \times n$  system with exponential nonlinearity if the gauge potential and the Higgs field are algebraically restricted. The Toda system with (1-1) is one of the limiting equations if a coupling constant tends to zero. For extensive discussions on the relationship between the Toda system and its background in Physics we refer the readers to [Bennett 1934; Ganoulis et al. 1982; Lee 1991; Mansfield 1982; Yang 2001].

In this article we are concerned with rank-2 Toda systems. There are three types of Cartan matrices of rank 2:

$$A_2 = \begin{pmatrix} 2 & -1 \\ -1 & 2 \end{pmatrix}, \quad B_2 (= C_2) = \begin{pmatrix} 2 & -1 \\ -2 & 2 \end{pmatrix}, \quad G_2 = \begin{pmatrix} 2 & -1 \\ -3 & 2 \end{pmatrix}.$$

One of our main theorems is the following estimate:

**Theorem 1.1.** *Let  $(k_{ij})_{2 \times 2}$  be one of the matrices above,  $h_i$  be positive  $C^1$  functions on  $M$ ,  $\alpha_i^t \in \mathbb{N} \cup \{0\}$ ,  $t \in \{1, 2, \dots, N\}$  and  $K$  be a compact subset of  $M \setminus S$ . If  $\rho_i \notin 4\pi\mathbb{N}$ , then there exists a constant  $C(K, \rho_1, \rho_2)$  such that for any solution  $u = (u_1, u_2)$  of (1-2)*

$$|u_i(x)| \leq C \quad \text{for all } x \in K, \quad i = 1, 2.$$

Our proof of Theorem 1.1 is based on the analysis of the behavior of solutions  $u^k = (u_1^k, u_2^k)$  near each blowup point. A point  $p \in M$  is called a blowup point if, along a sequence of points  $p_k \rightarrow p$ ,

$$\max_{i=1,2} \{ \tilde{u}_1^k(p_k), \tilde{u}_2^k(p_k) \} \rightarrow +\infty,$$

where

$$\tilde{u}_i^k(x) = u_i^k(x) + 4\pi \sum_t \alpha_t^k G(x, p_t),$$

and  $G(x, y)$  is the Green’s function of the Laplacian operator on  $M$ .

Suppose  $u^k$  is a sequence of solutions of (1-2). When  $n = 1$ , it has been proved that if  $u^k$  blows up somewhere, the mass distribution  $\rho h e^{u^k} / (\int_M h e^{u^k})$  will concentrate; that is, for a set of finite points

$p_1, p_2, \dots, p_L$  and positive numbers  $m_1, \dots, m_L$

$$\frac{\rho h e^{u^k}}{\int_M h e^{u^k}} \rightarrow \sum_{i=1}^L m_i \delta_{p_i} \quad \text{as } k \rightarrow \infty.$$

In other words, “ $u_k$  concentrates” means  $u^k(x) \rightarrow -\infty$  if  $x$  is not a blowup point. This “blowup implies concentration” was first noted by Brezis and Merle [1991] and was later proved by Li [1999], Li and Shafirir [1994] and Bartolucci and Tarantello [2002]. But for  $n \geq 2$ , this phenomenon might fail in general. A component  $u_i^k$  is called not concentrating if  $u_i^k \not\rightarrow -\infty$  away from blowup points, or equivalently,  $\tilde{u}_i^k$  converges to some smooth function  $w_i$  away from blowup points. It is natural to ask whether it is possible to have all components not concentrating. For  $n = 2$ , we prove it is impossible.

**Theorem 1.2.** *Suppose  $u^k$  is a sequence of blowup solutions of a rank-2 Toda system (1-2). Then at least one component of  $u^k$  satisfies  $u_i^k(x) \rightarrow -\infty$  if  $x$  is not contained in the blowup set.*

The first example of such nonconcentration phenomenon was proved by Lin and Tarantello [2016]. The new phenomenon makes the study of systems ( $n \geq 2$ ) much more difficult than the mean-field equation ( $n = 1$ ). Recently, Battaglia [2015] and Lin, Yang and Zhong [Lin et al. 2017] independently proved the result of Theorem 1.2 for  $n \geq 3$ .

As mentioned before, our proofs of Theorems 1.1 and 1.2 are based on the asymptotic behavior of local bubbling solutions. For simplicity we set up the situation as follows:

Let  $u^k = (u_1^k, u_2^k)$  be a sequence of solutions of

$$\Delta u_i^k + \sum_{j=1}^2 k_{ij} h_j^k e^{u_j^k} = 4\pi \alpha_i \delta_0 \quad \text{in } B(0, 1), \quad i = 1, 2, \tag{1-5}$$

where  $\alpha_i > -1$ .  $B(0, 1)$  is the unit ball in  $\mathbb{R}^2$  (we use  $B(p, r)$  to denote the ball centered at  $p$  with radius  $r$ ) and  $(k_{ij})_{2 \times 2}$  is  $A_2, B_2$  or  $G_2$ . Throughout the paper,  $h_1^k, h_2^k$  are smooth functions satisfying  $h_1^k(0) = h_2^k(0) = 1$  and

$$\frac{1}{C} \leq h_i^k \leq C, \quad \|h_i^k\|_{C^1(B(0,1))} \leq C \quad \text{in } B(0, 1), \quad i = 1, 2. \tag{1-6}$$

For solutions  $u^k = (u_1^k, u_2^k)$  we assume

$$\left\{ \begin{array}{l} 0 \text{ is the only blowup point of } u^k, \\ |u_i^k(x) - u_i^k(y)| \leq C \quad \text{for all } x, y \in \partial B(0, 1), \quad i = 1, 2, \\ \int_{B(0,1)} h_i^k e^{u_i^k} \leq C, \quad i = 1, 2. \end{array} \right. \tag{1-7}$$

For this sequence of blowup solutions we define the local mass by

$$\sigma_i = \lim_{r \rightarrow 0} \lim_{k \rightarrow \infty} \frac{1}{2\pi} \int_{B(0,r)} h_i^k e^{u_i^k}, \quad i = 1, 2. \tag{1-8}$$

It is known that 0 is a blowup point if and only if  $(\sigma_1, \sigma_2) \neq (0, 0)$ . The proof is to use ideas from [Brezis and Merle 1991] and has become standard now. We refer the readers to [Lee et al. 2017] for a

complete proof. One important property of  $(\sigma_1, \sigma_2)$  is the so-called Pohozaev identity (P.I. in short)

$$k_{21}\sigma_1^2 + k_{12}k_{21}\sigma_1\sigma_2 + k_{12}\sigma_2^2 = 2k_{21}\mu_1\sigma_1 + 2k_{12}\mu_2\sigma_2, \tag{1-9}$$

where  $\mu_i = 1 + \alpha_i$ . Take  $A_2$  as an example; the P.I. is

$$\sigma_1^2 - \sigma_1\sigma_2 + \sigma_2^2 = 2\mu_1\sigma_1 + 2\mu_2\sigma_2.$$

The proof of (1-9) was given in [Lin et al. 2015] where we initiated an algorithm to calculate all the possible (finitely many) values of local masses and (1-9) played an essential role. But the argument there seems not very efficient. In this work we add major new ingredients to our approach and improve the classification of  $(\sigma_1, \sigma_2)$  to the following sharper form:

**Theorem 1.3.** *Let  $u^k$  be a sequence of blowup solutions of (1-5) which also satisfies (1-6) and (1-7). Suppose  $\sigma_1$  and  $\sigma_2$  are local masses defined by (1-8). Then  $\sigma_i$  can be written as*

$$\sigma_i = 2(N_{i,1}\mu_1 + N_{i,2}\mu_2 + N_{i,3}), \quad i = 1, 2,$$

for some  $N_{i,1}, N_{i,2}, N_{i,3} \in \mathbb{Z}$  ( $i = 1, 2$ ).

**Theorem 1.3** is proved in Sections 5 and 6. In Section 5, we give an explicit procedure to calculate the local masses. Take the  $A_2$  system as an example; we start with  $\sigma_1 = 0$  and the P.I. gives  $\sigma_2 = 2\mu_2$ . With  $\sigma_2 = 2\mu_2$ , the P.I. gives  $\sigma_1 = 2\mu_1 + 2\mu_2$  and so on. Let  $\Gamma(\mu_1, \mu_2)$  be the set obtained by the above algorithm. Then  $\Gamma(\mu_1, \mu_2)$  is equal to:

- (i)  $(2\mu_1, 0), (2\mu_1, 2\mu_1 + 2\mu_2), (2\mu_1 + 2\mu_2, 2\mu_1 + 2\mu_2), (2\mu_1 + 2\mu_2, 2\mu_2), (0, 2\mu_2)$  for  $A_2$ ,
- (ii)  $(2\mu_1, 0), (2\mu_1, 4\mu_1 + 2\mu_2), (4\mu_1 + 2\mu_2, 4\mu_1 + 2\mu_2), (4\mu_1 + 2\mu_2, 4\mu_1 + 4\mu_2),$   
 $(0, 2\mu_2), (2\mu_1 + 2\mu_2, 2\mu_2), (2\mu_1 + 2\mu_2, 4\mu_1 + 4\mu_2)$  for  $B_2$ ,
- (iii)  $(2\mu_1, 0), (2\mu_1, 6\mu_1 + 2\mu_2), (6\mu_1 + 2\mu_2, 6\mu_1 + 2\mu_2), (6\mu_1 + 2\mu_2, 12\mu_1 + 6\mu_2),$   
 $(8\mu_1 + 4\mu_2, 12\mu_1 + 6\mu_2), (8\mu_1 + 4\mu_2, 12\mu_1 + 8\mu_2), (0, 2\mu_2), (2\mu_1 + 2\mu_2, 2\mu_2),$   
 $(2\mu_1 + 2\mu_2, 6\mu_1 + 6\mu_2), (6\mu_1 + 4\mu_2, 6\mu_1 + 6\mu_2), (6\mu_1 + 4\mu_2, 12\mu_1 + 8\mu_2)$  for  $G_2$ .

**Definition 1.4.** A pair of local masses  $(\sigma_1, \sigma_2) \in \Gamma(\mu_1, \mu_2)$  is called special if

$$(\sigma_1, \sigma_2) = \begin{cases} (2\mu_1 + 2\mu_2, 2\mu_1 + 2\mu_2) & \text{for } A_2, \\ (4\mu_1 + 2\mu_2, 4\mu_1 + 4\mu_2) & \text{for } B_2, \\ (8\mu_1 + 4\mu_2, 12\mu_1 + 8\mu_2) & \text{for } G_2. \end{cases}$$

The analysis of local solutions in [Lin et al. 2015] describes a method to pick a family of points  $\Gamma_k = \{0, x_1^k, \dots, x_N^k\}$  (if 0 is a singular point, otherwise 0 can be deleted from  $\Gamma_k$ ) such that a tiny ball  $B(x_i^k, l_j^k)$  contributes an amount of mass (which is quantized), and the following Harnack-type inequality holds:

$$u_i^k(x) + 2 \log \text{dist}(x, \Sigma_k) \leq C \quad \text{for all } x \in B(0, 1). \tag{1-10}$$

When  $\alpha_1 = \alpha_2 = 0$ , we can use **Theorem 1.3** to calculate all the pairs of even positive integers satisfying (1-9) and the set is exactly the same as  $\Gamma(1, 1)$ .

It is interesting to see whether any pair of the above really consists of the local masses of some sequence of blowup solutions of (1-2). For  $K = A_2$  the existence of such a local blowup sequence has been obtained; see [Musso et al. 2016; Lin and Yan 2013].

After  $\Sigma_k$  is picked, the difficulty at the next step is how to calculate the mass contributed from outside  $B(x_j^k, l_j^k)$   $j = 1, 2, \dots, N$ . In Section 6, we see that the mass outside this union could be very messy. However, the picture is very clean if  $(\alpha_1, \alpha_2)$  satisfies the following  $Q$ -condition:

$$\alpha_1, \alpha_2 \text{ and } 1 \text{ are linearly independent over } Q.$$

**Theorem 1.5.** *Suppose  $(\alpha_1, \alpha_2)$  satisfies the  $Q$ -condition. Then  $(\sigma_1, \sigma_2) \in \Gamma(\mu_1, \mu_2)$ . Furthermore, for any sequence of solutions of (1-5) satisfying (1-6) and (1-7), the following Harnack-type inequality holds:*

$$u_i^k(x) + 2 \log |x| \leq C \quad \text{for } x \in B(0, 1).$$

For (1-2), let  $\mu_{1,t} = \alpha_t^1 + 1$  and  $\mu_{2,t} = \alpha_t^2 + 1$  at a vortex point  $p_t \in S$ , and define

$$\Gamma_i = \{2\pi(\sum_{t \in J} \sigma_{i,t} + 2n) \mid (\sigma_{1,t}, \sigma_{2,t}) \in \Gamma(\mu_{1,t}, \mu_{2,t}), J \subseteq S, n \in \mathbb{N} \cup \{0\}\}. \tag{1-11}$$

Based on Theorem 1.5, Theorem 1.1 can be extended to the following version:

**Theorem 1.6.** *Let  $h_i$  be positive  $C^1$  functions on  $M$ , and  $K$  be a compact set in  $M$ . For every point  $p_t \in S$ , if either both  $\alpha_t^1, \alpha_t^2 \in \mathbb{N} \cup \{0\}$  or  $(\alpha_t^1, \alpha_t^2)$  satisfies the  $Q$ -condition, then for  $\rho_i \notin \Gamma_i$  and  $u = (u_1, u_2)$  a solution of (1-2), there exists a constant  $C$  such that*

$$|u_i(x)| \leq C \quad \text{for all } x \in K.$$

The organization of this article is as follows. In Section 2 we establish the global mass for the entire solutions of some singular Liouville equation defined in  $\mathbb{R}^2$ . Then in Section 3 we review some fundamental tools proved in the previous work [Lin et al. 2015]. In Section 4 we present two crucial lemmas, which play the key role in the proof of main results. In Sections 5 and 6 we discuss the local mass on each bubbling disk centered at 0 and not at 0 respectively, and then all the main results are established based on previous discussions.

### 2. Total mass for Liouville equation

The main purpose of this section is to prove an estimate of the total mass for the solutions of the equation

$$\begin{cases} \Delta u + e^u = \sum_{i=1}^N 4\pi \alpha_i \delta_{p_i} & \text{in } \mathbb{R}^2, \\ \int_{\mathbb{R}^2} e^u < \infty, \end{cases} \tag{2-1}$$

where  $p_1, \dots, p_N$  are distinct points in  $\mathbb{R}^2$  and  $\alpha_i > -1$  for all  $1 \leq i \leq N$ .

**Theorem 2.1.** *Suppose  $u$  is a solution of (2-1) and  $\alpha_1, \dots, \alpha_N$  are positive integers. Then  $\frac{1}{4\pi} \int_{\mathbb{R}^2} e^u$  is an even integer.*

*Proof.* It is known that any solution  $u$  of (2-1) has, at infinity, the asymptotic behavior

$$u(z) = -2\alpha_\infty \log |z| + O(1), \quad \alpha_\infty > 1, \tag{2-2}$$

and  $u$  satisfies

$$\frac{1}{2\pi} \int_{\mathbb{R}^2} e^u dx = 2 \sum_{i=1}^N \alpha_i + 2\alpha_\infty. \tag{2-3}$$

We shall prove that  $\alpha_\infty + \sum_{i=1}^N \alpha_i$  is an even integer. A classical Liouville theorem (see [Chou and Wan 1994]) says that  $u$  can be written as

$$u = \log \frac{8|f'(z)|^2}{(1 + |f(z)|^2)^2}, \quad z \in \mathbb{R}^2, \tag{2-4}$$

for some meromorphic function  $f$ . In general,  $f(z)$  is multivalued and any vertex  $p_i$  is a branch point. However if  $\alpha_i \in \mathbb{N} \cup \{0\}$ , then  $f(z)$  is single-valued. Furthermore (2-2) implies that  $f(z)$  is meromorphic at infinity. Hence for any solution  $u$  of (2-1) there is a meromorphic function  $f$  on  $\mathbb{S}^2 = \mathbb{C} \cup \{\infty\}$  such that (2-4) holds. Then

$$\begin{aligned} 4\pi \left( \sum_{j=1}^N \alpha_j + \alpha_\infty \right) &= \int_{\mathbb{R}^2} e^u = 8 \int_{\mathbb{R}^2} \frac{|f'(z)|^2}{(1 + |f(z)|^2)^2} dx dy \\ &= 8(\deg f) \int_{\mathbb{R}^2} \frac{d\tilde{x} d\tilde{y}}{(1 + |w|^2)^2} = 8\pi(\deg f), \end{aligned}$$

where  $\deg(f)$  is the degree of  $f$  as a map from  $\mathbb{S}^2 = \mathbb{C} \cup \{\infty\}$  onto  $\mathbb{S}^2$ , and  $w = f(z) = \tilde{x} + i\tilde{y}$ . Thus we have

$$\sum_{j=1}^N \alpha_j + \alpha_\infty = 2 \deg(f). \tag{2-5}$$

**Theorem 2.2.** *Suppose  $u$  is a solution of*

$$\begin{cases} \Delta u + e^u = 4\pi\alpha_0\delta_{p_0} + \sum_{i=1}^N 4\pi\alpha_i\delta_{p_i} & \text{in } \mathbb{R}^2, \\ \int_{\mathbb{R}^2} e^u < \infty, \end{cases} \tag{2-5}$$

where  $p_0, p_1, \dots, p_N$  are distinct points in  $\mathbb{R}^2$  and  $\alpha_1, \dots, \alpha_N$  are positive integers,  $\alpha_0 > -1$ . Then  $\frac{1}{4\pi} \int_{\mathbb{R}^2} e^u$  is equal to  $2(\alpha_0 + 1) + 2k$  for some  $k \in \mathbb{Z}$  or  $2k_1$  for some  $k_1 \in \mathbb{N}$ .

*Proof.* As in Theorem 2.1, there is a developing map  $f(z)$  of  $u$  such that

$$u(z) = \log \frac{8|f'(z)|^2}{(1 + |f(z)|^2)^2}, \quad z \in \mathbb{C}. \tag{2-6}$$

On one hand by (2-5),  $u_{zz} - \frac{1}{2}u_z^2$  is a meromorphic function in  $\mathbb{C} \cup \{\infty\}$  because away from the Dirac masses

$$4(u_{zz} - \frac{1}{2}u_z^2)_{\bar{z}} = -(e^u)_z + u_z e^u = 0.$$

By  $u(z) = 2\alpha_i \log |z - p_i| + O(1)$  near  $p_i$  we have

$$u_{zz} - \frac{1}{2}u_z^2 = -2 \left( \sum_{j=0}^N \frac{1}{2}\alpha_j (\frac{1}{2}\alpha_j + 1)(z - p_j)^{-2} + A_j(z - p_j)^{-1} + B \right),$$



where  $A_0, \dots, A_N, B \in \mathbb{C}$  are some constants. On the other hand by (2-6), a straightforward computation shows that

$$u_{zz} - \frac{1}{2}u_z^2 = \frac{f'''}{f'} - \frac{3}{2}\left(\frac{f''}{f'}\right)^2. \tag{2-7}$$

Using the Schwarz derivative of  $f$ ,

$$\{f; z\} = \frac{f'''(z)}{f'(z)} - \frac{3}{2}\left(\frac{f''(z)}{f'(z)}\right)^2,$$

and letting

$$I(z) = \sum_{j=0}^N \frac{1}{2}\alpha_j(\frac{1}{2}\alpha_j + 1)(z - p_j)^{-2} + A_j(z - p_j)^{-1} + B,$$

we can write the equation for  $f$  as

$$\{f, z\} = -2I(z). \tag{2-8}$$

A well-known classic theorem (see [Whittaker and Watson 1927]) says that for any two linearly independent solutions  $y_1$  and  $y_2$  of

$$y''(z) = I(z)y(z), \tag{2-9}$$

the ratio  $y_2/y_1$  always satisfies

$$\{y_2/y_1; z\} = -2I(z).$$

By (2-8) and a basic result of the Schwarz derivative,  $f(z)$  can be written as the ratio of two linearly independent solutions. This is how (2-1) is related to the complex ODE (2-9). We refer the readers to [Chai et al. 2015] for the details.

For (2-9), there is an associated monodromy representation  $\rho$  from  $\pi_1(\mathbb{C} \setminus \{p_0, p_1, \dots, p_N\}; q)$  to  $GL(2; \mathbb{C})$ , where  $q$  is a base point. Note that at any singular point  $p_j$ , the local exponents are  $\frac{1}{2}\alpha_j + 1$  and  $-\frac{1}{2}\alpha_j$ . It is known from [Lin et al. 2012, Section 7] that  $e^{-u}$  can be locally written as

$$e^{-u} = |v_1|^2 + |v_2|^2 = \langle (v_1, v_2)^t, (v_1, v_2)^t \rangle,$$

where  $v_1, v_2$  are the two fundamental solutions of (2-9). After encircling the singular point  $p_j$  once, we have  $e^{-u} = \langle \rho_j(v_1, v_2)^t, \rho_j(v_1, v_2)^t \rangle$  and the value does not change. Therefore, we conclude that  $\rho_j$  is unitary and

$$\rho_j = \rho(\gamma_j) = C_j \begin{pmatrix} e^{\pi i \alpha_j} & 0 \\ 0 & e^{-\pi i \alpha_j} \end{pmatrix} C_j^{-1},$$

where  $\gamma_j \in \pi_1(\mathbb{C} \setminus \{p_0, \dots, p_N\}, q)$  encircles  $p_j$  only once,  $0 \leq j \leq N$ , while the monodromy at  $\infty$  is  $\rho_\infty$ . Then we have

$$\rho_\infty \rho_N \cdots \rho_0 = I_{2 \times 2}.$$

Note that  $\rho_j = \pm I_{2 \times 2}$  for  $1 \leq j \leq N$ . Hence

$$\rho_\infty^{-1} = D_0 \begin{pmatrix} e^{\pi i \sum_{j=0}^N \alpha_j} & 0 \\ 0 & e^{-\pi i \sum_{j=0}^N \alpha_j} \end{pmatrix} D_0^{-1}$$

for some constant invertible matrix  $D_0$ .

On the other hand, the local exponents at  $\infty$  can be computed as follows. Let  $\hat{y}(z) = y(\frac{1}{z})$ , where  $y$  is a solution of (2-9). Then we have

$$\hat{y}''(z) + \frac{2}{z}\hat{y}'(z) = \hat{I}(z)\hat{y}(z), \tag{2-10}$$

where  $\hat{I}(z) = I(\frac{1}{z})z^{-4}$ . Since  $I(z)$  is the Schwarz derivative of  $f(z)$ , by direct computation  $\hat{I}(z)$  is the Schwarz derivative of  $f(\frac{1}{z})$ . As before we let  $\hat{u}(z) = u(\frac{1}{z}) - 4 \log |z|$ . Then  $f(\frac{1}{z})$  is the developing map of  $\hat{u}(z)$ . Since

$$\hat{u}(z) = 2(\alpha_\infty - 2) \log |z| + O(1) \quad \text{near } 0,$$

(because  $u(z) = -2\alpha_\infty \log |z| + O(1)$  at infinity), we have

$$\hat{I}(z) = \frac{1}{2}\alpha_\infty(\frac{1}{2}\alpha_\infty - 1)z^{-2} + \text{higher-order terms of } z \quad \text{near } 0.$$

By (2-10) we see that the local exponents of (2-9) are  $-\frac{1}{2}\alpha_\infty$  and  $\frac{1}{2}\alpha_\infty - 1$ . Hence  $e^{\pi i \alpha_\infty}$  equals either  $e^{\pi i \sum_{j=0}^N \alpha_j}$  or  $e^{-\pi i \sum_{j=0}^N \alpha_j}$ , which yields

$$\alpha_\infty = -\sum_{j=0}^N \alpha_j + 2k \quad \text{or} \quad \alpha_\infty = \sum_{j=0}^N \alpha_j + 2k \tag{2-11}$$

for some  $k \in \mathbb{Z}$ . Since

$$\frac{1}{4\pi} \int_{\mathbb{R}^2} e^u = \sum_{j=0}^N \alpha_j + \alpha_\infty,$$

we either have  $\frac{1}{4\pi} \int_{\mathbb{R}^2} e^u = 2k$  for some  $k \in \mathbb{N}$  if the first case holds or  $\frac{1}{4\pi} \int_{\mathbb{R}^2} e^u = 2(\alpha_0 + 1) + 2k'$  for  $k' = \sum_{i=1}^N \alpha_i + k - 1$  if the second case holds. □

**Remark 2.3.** After proving Theorems 2.1 and 2.2, we found a stronger version of both theorems in [Eremenko et al. 2014]. Because we only need the present form of both theorems, we include our proofs here to make the paper more self-contained.

### 3. Review of bubbling analysis from a selection process

Let  $u^k = (u_1^k, u_2^k)$  be solutions of (1-5) such that (1-6) and (1-7) hold. In this section we review the process to select a set  $\Sigma_k = \{0, x_1^k, \dots, x_n^k\}$  and balls  $B(x_i^k, l_k)$  such that  $u^k$  has nonzero local masses in  $B(x_i^k, l_k)$ . This selection process was first carried out in [Lin et al. 2015]. We briefly review it below.

The set  $\Sigma_k$  is constructed by induction. If (1-5) has no singularity, we start with  $\Sigma_k = \emptyset$ . If (1-5) has a singularity, we start with  $\Sigma_k = \{0\}$ . By induction suppose  $\Sigma_k$  consists of  $\{0, x_1^k, \dots, x_{m-1}^k\}$ . Then we consider

$$\max_{x \in B_1} \max_{i=1,2} (u_i^k(x) + 2 \log \text{dist}(x, \Sigma_k)). \tag{3-1}$$

If the maximum is bounded from above by a constant independent of  $k$ , the process stops and  $\Sigma_k$  is exactly equal to  $\{0, x_1^k, \dots, x_{m-1}^k\}$ . However if the maximum tends to infinity, let  $q_k$  be where (3-1) is achieved and we set

$$d_k = \frac{1}{2} \text{dist}(q_k, \Sigma_k)$$

and

$$S_i^k(x) = u_i^k(x) + 2 \log(d_k - |x - q_k|) \quad \text{in } B(q_k, d_k), \quad i = 1, 2.$$

Suppose  $i_0$  is the component that attains

$$\max_i \max_{x \in \bar{B}(q_k, d_k)} S_i^k \tag{3-2}$$

at  $p_k$ . Then we set

$$\tilde{l}_k = \frac{1}{2}(d_k - |p_k - q_k|)$$

and scale  $u_i^k$  by

$$v_i^k(y) = u_i^k(p_k + e^{-\frac{1}{2}u_{i_0}^k(p_k)}y) - u_{i_0}^k(p_k) \quad \text{for } |y| \leq R_k \doteq e^{\frac{1}{2}u_{i_0}^k(p_k)}\tilde{l}_k. \tag{3-3}$$

It can be shown that  $R_k \rightarrow \infty$  and  $v_i^k$  is bounded from above over any fixed compact subset of  $\mathbb{R}^2$ . Thus by passing to a subsequence,  $v_i^k$  satisfies one of the following two alternatives:

(a)  $(v_1^k, v_2^k)$  converges in  $C_{\text{loc}}^2(\mathbb{R}^2)$  to  $(v_1, v_2)$  which satisfies

$$\Delta v_i + \sum_{j=1}^2 k_{ij} e^{v_j} = 0 \quad \text{in } \mathbb{R}^2, \quad i = 1, 2. \tag{3-4}$$

(b) Either  $v_1^k$  converges to

$$\Delta v_1 + 2e^{v_1} = 0 \quad \text{in } \mathbb{R}^2 \tag{3-5}$$

and  $v_2^k \rightarrow -\infty$  over any fixed compact subset of  $\mathbb{R}^2$  or  $v_2^k$  converges to  $\Delta v_2 + 2e^{v_2} = 0$  in  $\mathbb{R}^2$  and  $v_1^k \rightarrow -\infty$  over any fixed compact subset of  $\mathbb{R}^2$ .

Therefore in either case, we could choose  $l_k^* \rightarrow \infty$  such that

$$v_i^k(y) + 2 \log |y| \leq C \quad \text{for } i = 1, 2 \quad \text{and } |y| \leq l_k^* \tag{3-6}$$

and

$$\int_{B(0, l_k^*)} h_i^k e^{v_i^k} dy = \int_{\mathbb{R}^2} e^{v_i(y)} + o(1).$$

By scaling back to  $u_i^k$ , we add  $p_k$  in  $\Sigma_k$  with

$$l_k = e^{-\frac{1}{2}u_{i_0}^k(p_k)}l_k^*.$$

We can continue in this way until the Harnack-type inequality (1-10) holds.

We summarize what the selection process has done in the following proposition (a detailed proof for a more general case can be found in [Lin et al. 2015, Proposition 2.1]):

**Proposition 3A.** *Let  $u^k$  be described as above. Then there exist a finite set  $\Sigma_k := \{0, x_1^k, \dots, x_m^k\}$  (if 0 is not a singular point, then 0 can be deleted from  $\Sigma_k$ ) and positive numbers  $l_1^k, \dots, l_m^k \rightarrow 0$  as  $k \rightarrow \infty$  such that the following hold:*

- (1) *There exists  $C > 0$  independent of  $k$  such that (1-10) holds and all the components have fast decay on  $\partial B(x_j^k, l_j^k)$ ,  $j = 1, \dots, m$ . (The definition of fast decay can be found in Definition 3.1 below).*

(2) In  $B(x_j^k, l_j^k)$  ( $j = 1, \dots, m$ ), let  $R_{j,k} = e^{\frac{1}{2}u_{i_0}^k(x_j^k)}l_j^k$ ,  $u_{i_0}^k(x_j^k) = \max_i u_i^k(x_j^k)$  and

$$v_i^k(y) = u_i^k(x_j^k + e^{-\frac{1}{2}u_{i_0}^k(x_j^k)}y) - u_{i_0}^k(x_j^k) \tag{3-7}$$

for  $|y| \leq R_{j,k}$ ; then  $v^k = (v_1^k, v_2^k)$  satisfies either (a) or (b).

(3)  $B(x_j^k, l_j^k) \cap B(x_i^k, l_i^k) = \emptyset$ ,  $i \neq j$ .

The inequality (1-10) is a Harnack-type inequality, because it implies the following result.

**Proposition 3B.** Suppose  $u^k$  satisfies (1-5)–(1-7) and

$$u_i^k(x) + 2 \log |x - x_0| \leq C \quad \text{for } x \in B(x_0, 2r_k).$$

Then

$$|u_i^k(x_1) - u_i^k(x_2)| \leq C_0 \quad \text{for } \frac{1}{2} \leq \frac{|x_1 - x_0|}{|x_2 - x_0|} \leq 2 \quad \text{and } x_1, x_2 \in B(x_0, r_k). \tag{3-8}$$

The proof of Proposition 3B is standard, see [Lin et al. 2015, Lemma 2.4], so we omit it here. Let  $x_l^k \in \Sigma_k$  and  $\tau_l^k = \frac{1}{2} \text{dist}(x_l^k, \Sigma_k \setminus \{x_l^k\})$ ; then (3-8) implies

$$u_i^k(x) = \bar{u}_{x_l^k, i}^k(r) + O(1), \quad x \in B(x_l^k, \tau_l^k), \tag{3-9}$$

where  $r = |x_l^k - x|$  and  $\bar{u}_{x_l^k, i}^k$  is the average of  $u_i^k$  on  $\partial B(x_l^k, r)$ ,

$$\bar{u}_{x_l^k, i}^k(r) = \frac{1}{2\pi r} \int_{\partial B(x_l^k, r)} u_i^k dS, \tag{3-10}$$

and  $O(1)$  is independent of  $r$  and  $k$ .

Next we introduce the notions of slow decay and fast decay in our bubbling analysis.

**Definition 3.1.** We say  $u_i^k$  has fast decay on  $\partial B(x_0, r_k)$  if along a subsequence

$$u_i^k(x) + 2 \log |x - x_0| \leq -N_k \quad \text{for all } x \in \partial B(x_0, r_k),$$

for some  $N_k \rightarrow \infty$  and we say  $u_i^k$  has slow decay if there is a constant  $C$  independent of  $k$  such that

$$u_i^k(x) + 2 \log |x - x_0| \geq -C \quad \text{for all } x \in \partial B(x_0, r_k).$$

Furthermore, we say  $u_i^k$  is fast-decaying in  $B(x_0, s_k) \setminus B(x_0, r_k)$  if  $u_i^k$  has fast decay on  $\partial B(x_0, l_k)$  for any  $l_k \in [r_k, s_k]$ .

The concept of fast decay is important for evaluating the Pohozaev identities. The following proposition is a direct consequence of [Lin et al. 2015, Proposition 3.1] and it says if both components are fast-decaying on the boundary, the Pohozaev identity holds for the local masses.

In the following proposition, we let  $B = B(x^k, r_k)$ . If  $x^k \neq 0$ , we assume  $0 \notin B(x^k, 2r_k)$ .

**Proposition 3C.** Suppose both  $u_1^k, u_2^k$  have fast decay on  $\partial B$ , where  $B$  is given above. Then  $(\sigma_1, \sigma_2)$  satisfies the P.I. (1-9), where

$$\sigma_i = \lim_{k \rightarrow 0} \frac{1}{2\pi} \int_B h_i^k e^{u_i^k}, \quad i = 1, 2.$$

We refer the readers to [Lin et al. 2015, Proposition 3.1] for the proof.

### 4. Two lemmas

In this section, we prove two crucial lemmas which play the key role in Sections 5 and 6. For Lemma 4.1, we assume:

(i) The Harnack inequality

$$u_i^k(x) + 2 \log |x| \leq C \quad \text{for } \frac{1}{2}l_k \leq |x| \leq 2s_k, \quad i = 1, 2,$$

holds for both components.

(ii) Both components  $u_i^k$  have fast decay on  $\partial B(0, l_k)$  and  $\sigma_i^k(B(0, l_k)) = \sigma_i + o(1)$  for  $i = 1, 2$ , where

$$\sigma_i = \lim_{r \rightarrow 0} \lim_{k \rightarrow \infty} \frac{1}{2\pi} \int_{B(0, rs_k)} h_i^k e^{u_i^k}, \quad i = 1, 2.$$

(iii) One of  $u_i^k$ ,  $i = 1, 2$ , has slow decay on  $\partial B(0, s_k)$ .

**Lemma 4.1.** (a) Assume (i) and (ii). If  $u_i^k$  has slow decay on  $\partial B(0, s_k)$ , then

$$2\mu_i - \sum_{j=1}^2 k_{ij}\sigma_j > 0.$$

(b) Assume (i), (ii) and (iii). Let  $u_i^k$  be a slow-decaying component on  $\partial B(0, s_k)$ . Then the other component has fast decay on  $\partial B(0, s_k)$ .

*Proof.* (a) Suppose that  $u_i^k$  has slow decay on  $\partial B(0, s_k)$ . Then the scaling

$$v_j^k(y) = u_j^k(s_k y) + 2 \log s_k, \quad j = 1, 2 \text{ for } y \in B_2$$

gives

$$\Delta v_j^k(y) + \sum_{l=1}^2 k_{jl} h_l^k(s_k y) e^{v_l^k(y)} = 4\pi \alpha_j \delta_0 \quad \text{in } y \in B_2.$$

If the other component also has slow decay on  $\partial B(0, s_k)$ , then  $(v_1^k, v_2^k)$  converges to  $(v_1, v_2)$  which satisfies

$$\Delta v_j(y) + \sum_{l=1}^2 k_{jl} e^{v_l} = 0 \quad \text{in } B_2 \setminus \{0\}, \quad j = 1, 2. \tag{4-1}$$

If the other component has fast decay on  $\partial B(0, s_k)$ , then  $v_i^k(y)$  converges to  $v_i(y)$  and  $v_j(y) \rightarrow -\infty$ ,  $j \neq i$ . Furthermore,  $v_i(y)$  satisfies

$$\Delta v_i(y) + 2e^{v_i} = 0 \quad \text{in } B_2 \setminus \{0\}. \tag{4-2}$$

For any  $r > 0$ ,

$$\begin{aligned} \int_{\partial B(0,r)} \frac{\partial v_i(y)}{\partial \nu} dS &= \lim_{k \rightarrow \infty} \left( 4\pi \alpha_i - \sum_{j=1}^2 \int_{B(0,r)} k_{ij} h_j^k e^{v_j^k} dy \right) \\ &= 4\pi \alpha_i - 2\pi \sum_{j=1}^2 k_{ij} \sigma_j + o(1) \doteq 4\pi \beta_i + o(1), \end{aligned}$$

which implies the right-hand sides of both (4-1) and (4-2) should be replaced by  $4\pi\beta_i\delta_0$ . If  $\beta_i \leq -1$ , we can use the finite energy assumption (see the bottom assumption in (1-7)) to conclude that either (4-1) or (4-2) has no solutions. Hence  $\alpha_i - \frac{1}{2} \sum_{j=1}^2 k_{ij}\sigma_j > -1$  and then (a) is proved.

(b) Since both components have fast decay on  $\partial B(0, l_k)$ , the pair  $(\sigma_1, \sigma_2)$  satisfies the P.I. (1-9). By a simple manipulation, the P.I. (1-9) can be written as

$$k_{21}\sigma_1(4\mu_1 - k_{12}\sigma_2 - k_{11}\sigma_1) + k_{12}\sigma_2(4\mu_2 - k_{21}\sigma_1 - k_{22}\sigma_2) = 0. \tag{4-3}$$

Note by (a),

$$4\mu_i - \sum_{l=1}^2 k_{il}\sigma_l > 2\mu_i - \sum_{l=1}^2 k_{il}\sigma_l \geq 0.$$

Hence for  $j \neq i$

$$2\mu_j - \sum_{l=1}^2 k_{jl}\sigma_l < 4\mu_j - \sum_{l=1}^2 k_{jl}\sigma_l < 0,$$

where the last inequality is due to (4-3). By (a) again,  $u_j^k$  does not have slow decay on  $\partial B(0, s_k)$ . □

Our second lemma says that a fast-decaying component does not change its energy more than  $o(1)$ , regardless of the behavior of the other component.

**Lemma 4.2.** *Suppose the Harnack-type inequality holds for both components over  $r \in [\frac{1}{2}l_k, 2s_k]$ . If  $u_i^k$  is fast-decaying on  $r \in [l_k, s_k]$ , then*

$$\sigma_i^k(B(0, s_k)) = \sigma_i^k(B(0, l_k)) + o(1).$$

*Proof.* Obviously the conclusion holds if  $s_k/l_k \leq C$ . So we assume  $s_k/l_k \rightarrow +\infty$ . The Harnack-type inequality implies  $u_i^k(x) = \bar{u}_i^k(r) + o(1)$  for  $\frac{1}{2}l_k \leq |x| \leq 2s_k$ . Thus we obtain from (1-5) that

$$\frac{d}{dr}(\bar{u}_i^k(r) + 2 \log r) = \frac{2\mu_i - \sum_{j=1}^2 k_{ij}\sigma_j^k(r)}{r}, \quad l_k \leq r \leq s_k, \quad i = 1, 2,$$

where  $\sigma_j^k(r) = \sigma_j^k(B(0, r))$  and  $\sigma_j = \lim_{k \rightarrow +\infty} \sigma_j^k(l_k)$ ,  $j = 1, 2$ .

Without loss of generality, we assume that  $u_j^k$ ,  $j \neq i$ , is fast-decaying on  $\partial B(0, l_k)$ . Otherwise, we may choose  $\tilde{l}_k$  such that  $l_k \ll \tilde{l}_k$ ,  $u_i^k$  remains fast-decaying for  $r \in [l_k, \tilde{l}_k]$  and  $\sigma_i^k(B(0, r))$  does not change more than  $o(1)$ , while  $u_j^k$  is fast-decaying on  $\partial B(0, \tilde{l}_k)$ . If  $s_k/\tilde{l}_k \leq C$ , we get the conclusion as explained above. If  $s_k/\tilde{l}_k \rightarrow +\infty$ , by a little abuse of notation, we may replace  $\tilde{l}_k$  by  $l_k$ . Then both  $u_1^k, u_2^k$  have fast decay on  $\partial B(0, l_k)$ , and the P.I. holds at  $l_k$ , which implies that at least one component (say  $l$ ) satisfies

$$4\mu_l - \sum_{j=1}^2 k_{lj}\sigma_j^k(l_k) < 0.$$

Thus,

$$\frac{d}{dr}(\bar{u}_l^k(r) + 2 \log r) \leq -\frac{2\mu_l + o(1)}{r} \quad \text{at } r = l_k. \tag{4-4}$$

Suppose  $r_k \in [l_k, s_k]$  is the largest  $r$  such that

$$\frac{d}{dr}(\bar{u}_l^{(k)}(r) + 2 \log r) \leq -\frac{\mu_l}{r} \quad \text{for } r \in [l_k, r_k]. \tag{4-5}$$

Thus, either the equality holds at  $r = r_k$  or  $r_k = s_k$ . For simplicity, we let  $\varepsilon = \mu_l$ . By integrating (4-4) from  $l_k$  up to  $r \leq r_k$ , we have

$$\bar{u}_l^{(k)}(r) + 2 \log r \leq \bar{u}_l^{(k)}(l_k) + 2 \log(l_k) + \varepsilon \log\left(\frac{l_k}{r}\right);$$

that is for  $|x| = r$ ,

$$e^{u_l^k(x)} \leq O(1)e^{\bar{u}_l^{(k)}(r)} \leq O(1)e^{-N_k} l_k^\varepsilon r^{-(2+\varepsilon)},$$

where we used  $\bar{u}_l^{(k)}(l_k) + 2 \log l_k \leq -N_k$  by the assumption of fast decay. Thus

$$\int_{l_k \leq |x| \leq r_k} e^{u_l^k(x)} dx \leq O(1)e^{-N_k} l_k^\varepsilon \int_{l_k}^{r_k} r^{-(1+\varepsilon)} dr = O(1)\frac{e^{-N_k}}{\varepsilon} \rightarrow 0$$

as  $k \rightarrow +\infty$ . Hence

$$\sigma_l^k(r_k) = \sigma_l^k(l_k) + o(1). \tag{4-6}$$

If both components are fast-decaying on  $r \in [l_k, r_k]$ , then  $\lim_{k \rightarrow +\infty}(\sigma_1^k(r_k), \sigma_2^k(r_k)) = (\hat{\sigma}_1, \hat{\sigma}_2)$  also satisfies the P.I. (1-9). If  $\hat{\sigma}_j > \sigma_j$ , then  $j \neq l$  by (4-6). We choose  $r_k^* \leq r_k$  such that  $\sigma_j(r_k^*) = \sigma_j^k(l_k) + \varepsilon_0$  for small  $\varepsilon_0$ , and let  $\sigma_j^* = \lim_{k \rightarrow 0} \sigma_j(r_k^*)$ . Then  $\sigma_j^*$  and  $\sigma_l$  satisfies the P.I. (1-9) and it yields a contradiction provided  $\varepsilon_0$  is small. Thus, we have  $\sigma_m^k(r_k) = \sigma_m^k(l_k) + o(1)$ ,  $m = 1, 2$ . Then (4-4) holds at  $r = r_k$ , which implies  $r_k = s_k$ , and Lemma 4.2 is proved in this case.

If one of the components does not have fast decay on  $[l_k, r_k]$ , then we have  $l = i$  and  $u_j^k$ ,  $j \neq i$ , has slow decay on  $\partial B(0, r_k^*)$  for some  $r_k^* \leq r_k$ . If  $s_k/r_k \leq C$ , then (4-6) implies the lemma. If  $s_k/r_k \rightarrow +\infty$ , then by the scaling of  $u_j^k$  at  $r = r_k^*$ , the standard argument implies that there is a sequence of  $r_k^* \ll \tilde{r}_k = R_k r_k^* \ll s_k$  such that both components have fast decay on  $\tilde{r}_k$  and

$$\sigma_i^k(\tilde{r}_k) = \sigma_i(r_k^*) + o(1) = \sigma_i(l_k) + o(1) \quad \text{and} \quad \sigma_j^k(\tilde{r}_k) \geq \sigma_j^k(l_k) + \varepsilon_0$$

for  $j \neq i$  and  $\varepsilon_0 > 0$ . Therefore the assumption of Lemma 4.2 holds at  $r \in [\tilde{r}_k, s_k]$ . Then we repeat the argument starting from (4-4) and the lemma can be proved in a finite steps. □

**Remark 4.3.** Both lemmas will be used in Section 6 (and Section 5) for the case with singularity at 0 (and without singularity at 0).

### 5. Local mass on the bubbling disk centered at $x_l^k \neq 0$

**5A.** In this subsection we study the local behavior of  $u^k$  near  $x_l^k$ , where  $x_l^k \neq 0$ . For simplicity, we use  $x^k$  instead of  $x_l^k$  and  $\bar{u}_i^k(r)$  rather than  $\bar{u}_{x_l^k, i}^k(r)$ . Let

$$\tau^k = \frac{1}{2} \text{dist}(x^k, \Sigma_k \setminus \{x^k\}), \quad \sigma_i^k(r) = \frac{1}{2\pi} \int_{B(x^k, r)} h_i^k e^{u_i^k}, \quad i = 1, 2.$$

By Proposition 3A,  $l_k \leq \tau^k$ . Clearly  $u^k = (u_1^k, u_2^k)$  satisfies

$$\Delta u_i^k + \sum_{j=1}^2 k_{ij} h_j^k e^{u_j^k} = 0 \quad \text{in } B(x^k, \tau^k).$$

For a sequence  $s_k$ , we define

$$\hat{\sigma}_i(s_k) = \begin{cases} \lim_{k \rightarrow +\infty} \sigma_i^k(s_k) & \text{if } u_i^k \text{ has fast decay on } \partial B(x^k, s_k), \\ \lim_{r \rightarrow 0} \lim_{k \rightarrow +\infty} \sigma_i^k(rs_k) & \text{if } u_i^k \text{ has slow decay on } \partial B(x^k, s_k). \end{cases} \tag{5-1}$$

Recall that both components of  $u^k$  have fast decay on  $\partial B(x^k, l_k)$ . This is the starting point of the following proposition, which is a special case of Proposition 5.2 below.

In Proposition 5.1,  $(\mu_1, \mu_2)$  will be  $(1, 1)$  in both lemmas of Section 4.

**Proposition 5.1.** *Let  $u^k = (u_1^k, u_2^k)$  be the solution of (1-5) satisfying (1-7) and  $\hat{\sigma}_i(s_k)$  be defined in (5-1). The following holds:*

- (1) *At least one component  $u^k$  has fast decay on  $\partial B(x^k, \tau^k)$ .*
- (2)  *$(\hat{\sigma}_1(\tau^k), \hat{\sigma}_2(\tau^k))$  satisfies the P.I. (1-9) with  $\mu_1 = \mu_2 = 1$ .*
- (3)  *$(\hat{\sigma}_1(\tau^k), \hat{\sigma}_2(\tau^k)) \in \Gamma(1, 1)$ .*

*Proof.* If  $\tau_k/l_k \leq C$ , (1)–(3) hold obviously for  $\tau^k$ . So we assume  $\tau^k/l_k \rightarrow +\infty$ . First we remark that if  $u^k$  is fully bubbling in  $B(x^k, l_k)$  (i.e., (1) in Proposition 3A holds),  $(\hat{\sigma}_1(l_k), \hat{\sigma}_2(l_k))$  is special (see Definition 1.4) and satisfies

$$2\mu_i - \sum_{j=1}^2 k_{ij} \hat{\sigma}_j(l_k) < 0, \quad i = 1, 2.$$

Then by Lemma 4.1, both  $u_i^k$  have fast decay on  $\partial B(0, \tau^k)$  and Proposition 5.1 follows immediately.

Now we assume  $v_i^k$  defined in (3-7) and satisfies case (2) in Proposition 3A. We already know that both components have fast decay at  $r = l_k$ . If both components remain fast-decaying as  $r$  increases from  $l_k$  to  $\tau^k$ , Lemma 4.2 implies

$$\sigma_1^k(\tau^k) = \sigma_1^k(l_k) + o(1), \quad \sigma_2^k(\tau^k) = \sigma_2^k(l_k) + o(1)$$

and we are done. So we only consider the case that at least one component changes to a slow-decaying component. For simplicity, we assume that  $u_1^k$  changes to a slow-decaying component for some  $r_k \gg l_k$ . By Lemma 4.2,

$$\sigma_1^k(B(x^k, r_k)) \geq \sigma_1^k(B(x^k, l_k)) + c_0 \quad \text{for some } c_0 > 0.$$

We might choose  $s_k \leq r_k$  such that

$$\sigma_1^k(B(x^k, s_k)) = \sigma_1^k(B(x^k, l_k)) + \varepsilon_0,$$

and

$$\sigma_1^k(B(x^k, r)) < \sigma_1^k(B(x^k, l_k)) + \varepsilon_0 \quad \text{for all } r < s_k,$$

where  $\varepsilon_0 < \frac{1}{2}c_0$  is small.



Then Lemmas 4.1 and 4.2 together imply that  $u_1^k$  has slow decay on  $\partial B(x^k, s_k)$  and  $u_2^k$  has fast decay on  $\partial B(x^k, s_k)$  with

$$\hat{\sigma}_1(s_k) = \sigma_1^k(l_k) + o(1) \quad \text{and} \quad \hat{\sigma}_2(s_k) = \sigma_2^k(l_k) + o(1).$$

Let  $v_i^k(y) = u_i^k(x^k + s_k y) + 2 \log s_k$ . If  $\tau^k/s_k \leq C$ , there is nothing to prove. So we assume  $\tau^k/s_k \rightarrow \infty$ . Then  $v_1^k(y)$  converges to  $v_1(y)$  and  $v_2^k(y) \rightarrow -\infty$  in any compact set of  $\mathbb{R}^2$  as  $k \rightarrow +\infty$  and  $v_1(y)$  satisfies

$$\Delta v_1 + 2e^{v_1} = -2\pi \sum_{j=1}^2 (k_{1j} \hat{\sigma}_j(l_k)) \delta(0) \quad \text{in } \mathbb{R}^2. \tag{5-2}$$

Hence there is a sequence  $N_k^* \rightarrow +\infty$  as  $k \rightarrow +\infty$  that satisfies

- (1)  $N_k^* s_k \leq \tau^k$ ,
- (2)  $\int_{B(0, N_k^*)} e^{v_1} dy = \int_{\mathbb{R}^2} e^{v_1} dy + o(1)$ ,
- (3)  $v_i^k(y) + 2 \log |y| \leq -N_k$ ,  $i = 1, 2$ , for  $|y| = N_k^*$ .

Scaling back to  $u_i^k$ , we obtain that  $u_i^k$ ,  $i = 1, 2$ , have fast decay on  $\partial B(x^k, N_k^* s_k)$ .

We could use the classification theorem of [Prajapat and Tarantello 2001] to calculate the total mass of  $v_1$ , but instead we use the P.I. (1-9) to compute it. We know that both  $(\hat{\sigma}_1(l_k), \hat{\sigma}_2(l_k))$  and  $(\hat{\sigma}_1(N_k^* s_k), \hat{\sigma}_2(N_k^* s_k))$  satisfy the P.I. and  $\hat{\sigma}_2(N_k^* s_k) = \hat{\sigma}_2(l_k)$  by Lemma 4.2. With a fixed  $\sigma_2 = \hat{\sigma}_2(l_k)$ , P.I. (1-9) is a quadratic polynomial in  $\sigma_1$ ; then  $\hat{\sigma}_1(l_k)$  and  $\hat{\sigma}_1(N_k^* s_k)$  are two roots of the polynomial. From it, we can easily calculate  $\hat{\sigma}_1(N_k^* s_k)$ .

By a direct computation, we have

$$(\hat{\sigma}_1(N_k^* s_k), \hat{\sigma}_2(N_k^* s_k)) \in \Gamma(1, 1) \quad \text{if} \quad (\hat{\sigma}_1(l_k), \hat{\sigma}_2(l_k)) \in \Gamma(1, 1).$$

Thus (1)–(3) hold at  $r = N_k^* s_k$ . By denoting  $N_k^* s_k$  as  $l_k$ , we can repeat the same argument until  $\tau^k/l_k \leq C$ . Hence Proposition 5.1 is proved. □

**5B. Local mass in a group that does not contain 0.** In this subsection we collect some  $x_i^k \in \Sigma_k$  into a group  $S$ , a subset of  $\Sigma_k$  satisfying the following  $S$ -conditions:

- (1)  $0 \notin S$  and  $|S| \geq 2$ .
- (2) If  $|S| \geq 3$  and  $x_i^k, x_j^k, x_l^k$  are three distinct elements in  $S$ , then

$$\text{dist}(x_i^k, x_j^k) \leq C \text{dist}(x_j^k, x_l^k)$$

for some constant  $C$  independent of  $k$ .

- (3) For any  $x_m^k \in \Sigma_k \setminus S$ , we have  $\text{dist}(x_m^k, S) / \text{dist}(x_i^k, x_j^k) \rightarrow \infty$  as  $k \rightarrow \infty$ , where  $x_i^k, x_j^k \in S$ .

We write  $S$  as  $S = \{x_1^k, \dots, x_m^k\}$  and let

$$l^k(S) = 2 \max_{1 \leq j \leq m} \text{dist}(x_1^k, x_j^k). \tag{5-3}$$

Recall  $\tau_i^k = \frac{1}{2} \text{dist}(x_i^k, \Sigma_k \setminus \{x_i^k\})$ ; by (2) and (3) above we have  $l^k(S) \sim \tau_i^k$  for  $1 \leq i \leq m$ . Let

$$\tau_S^k = \frac{1}{2} \text{dist}(x_1^k, \Sigma_k \setminus S).$$

Then by (3) above we have  $\tau_S^k / \tau_i^k \rightarrow \infty$  for any  $x_i^k \in S$ .

By Proposition 5.1, we know that at least one of  $u_i^k$  has fast decay on  $\partial B(x_1^k, \tau_1^k)$ . Suppose  $u_1^k$  has fast decay on  $\partial B(x_1^k, \tau_1^k)$ . Then

$$u_1^k \text{ has fast decay on } \partial B(x_1^k, l^k(S)), \tag{5-4}$$

and we get

$$\begin{aligned} \sigma_1^k(B(x_1^k, l^k(S))) &= \frac{1}{2\pi} \int_{B(x_1^k, l^k(S))} h_1^k e^{u_1^k} dx \\ &= \frac{1}{2\pi} \int_{\bigcup_{j=1}^m B(x_j^k, \tau_j^k)} h_1^k e^{u_1^k} + \frac{1}{2\pi} \int_{B(x_1^k, l^k(S)) \setminus (\bigcup_{j=1}^m B(x_j^k, \tau_j^k))} h_1^k e^{u_1^k}. \end{aligned}$$

Since  $u_1^k$  has fast decay outside of  $B(x_j^k, \tau_j^k)$ , we have

$$e^{u_1^k(x)} \leq o(1) \max_j \{|x - x_j^k|^{-2}\} \quad \text{for } x \notin \bigcup_{j=1}^k B(x_j^k, \tau_j^k)$$

and the second integral is  $o(1)$ . Hence by Proposition 5.1,

$$\sigma_1^k(B(x_1^k, l^k(S))) = 2m_1 + o(1) \quad \text{for some } m_1 \in \mathbb{N} \cup \{0\}. \tag{5-5}$$

Similarly if  $u_2^k$  has fast decay on  $\partial B(x_1^k, \tau_1^k)$ , we have

$$\sigma_2^k(B(x_1^k, l^k(S))) = 2m_2 + o(1) \quad \text{for some } m_2 \in \mathbb{N} \cup \{0\}. \tag{5-6}$$

If  $u_2^k$  has slow decay on  $\partial B(x_1^k, \tau_1^k)$ , then it is easy to see that  $u_2^k$  has slow decay on  $\partial B(x_j^k, \tau_j^k)$ . By Proposition 5.1 we denote  $n_{i,j} \in \mathbb{N}$  by

$$2n_{i,j} = \lim_{r \rightarrow 0} \lim_{k \rightarrow \infty} \sigma_i^k(B(x_j^k, r\tau_j^k)), \quad 1 \leq j \leq m, i = 1, 2.$$

Define  $\hat{n}_{i,j}$  by

$$\hat{n}_{i,j} = - \sum_{l=1}^2 k_{il} n_{l,j}.$$

Then the slow decay of  $u_2^k$  on  $\partial B(x_j^k, \tau_j^k)$  implies  $1 + \hat{n}_{2,j} > 0$ . Since  $\hat{n}_{2,j} \in \mathbb{Z}$  we have  $\hat{n}_{2,j} \geq 0$ .

Furthermore, if we scale  $u^k$  by

$$v_i^k(y) = u_i^k(x_1^k + l^k(S)y) + 2 \log l^k(S), \quad i = 1, 2,$$

the sequence  $v_2^k$  converges to  $v_2(y)$  and  $v_1^k$  tends to  $-\infty$  over any compact subset of  $\mathbb{R}^2 \setminus \{0\}$ . Then  $v_2$  satisfies

$$\Delta v_2(y) + 2e^{v_2(y)} = 4\pi \sum_{j=1}^m \hat{n}_{2,j} \delta_{p_j} \quad \text{in } \mathbb{R}^2, \tag{5-7}$$

where  $p_j = \lim_{k \rightarrow \infty} (x_j^k - x_1^k) / l^k(S)$ . By [Theorem 2.1](#)

$$\frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v^2} = 2N \quad \text{for some } N \in \mathbb{N}.$$

Thus using the argument in [Proposition 5.1](#), we conclude that there is a sequence of  $N_k^* \rightarrow \infty$  such that both  $u_i^k$  ( $i = 1, 2$ ) have fast decay on  $\partial B(x_1^k, N_k^* l^k(S))$  and  $\sigma_i^k(B(x_1^k, N_k^* l^k(S))) = 2m_i + o(1)$ . Denote  $N_k^* l^k(S)$  by  $l_k$  for simplicity; we see that [\(5-5\)](#) and [\(5-6\)](#) hold at  $l_k$ . Then by using [Lemmas 4.1](#) and [4.2](#) we continue this process to obtain the following conclusion:

$$\text{At least one component of } u^k \text{ has fast decay on } \partial B(x_1^k, \tau_S^k). \tag{5-8}$$

Let  $\hat{\sigma}_i^k(B(x_1^k, \tau_S^k))$  be defined as in [\(5-1\)](#). Then

$$\hat{\sigma}_i^k(B(x_1^k, \tau_S^k)) = 2m_i(S), \quad \text{where } m_i(S) \in \mathbb{N} \cup \{0\}, \tag{5-9}$$

and the pair  $(2m_1(S), 2m_2(S))$  satisfies the P.I. [\(1-9\)](#).

Denote the group  $S$  by  $S_1$ . Based on this procedure, we can continue to select a new group  $S_2$  such that the  $S$ -conditions holds except we have to modify condition (2). In (2), we consider  $S_1$  as a single point as long as we compare the distance of distinct elements in  $S_2$ .

Set

$$\tau_{S_2}^k = \frac{1}{2} \text{dist}(x_1^k, \Sigma_k \setminus S_2) \quad \text{for } x_1^k \in S_2.$$

Then we follow the same argument as above to obtain the same conclusion as [\(5-8\)–\(5-9\)](#).

If [\(1-5\)](#) does not contain a singularity, the final step is to collect all the  $x_i^k$  into the single biggest group and [\(5-8\)–\(5-9\)](#) hold. Then we get  $(\sigma_1, \sigma_2) = (2m_1, 2m_2)$  (which satisfies the Pohozaev identity), where

$$\sigma_i = \lim_{r \rightarrow 0} \lim_{k \rightarrow \infty} \frac{1}{2\pi} \int_{B(0,r)} h_i^k e^{u_i^k}, \quad i = 1, 2.$$

By a direct computation, we can prove that the set of all the pairs of even integers solving [\(1-9\)](#) is exactly  $\Gamma(1, 1)$ . This proves [Theorem 1.3](#) if [\(1-5\)](#) has no singularities.

If 0 is a singularity of [\(1-5\)](#) then  $\Sigma_k$  can be written as a disjoint union of  $\{0\}$  and  $S_j$  ( $j = 1, \dots, m$ ). Here each  $S_j$  is collected by the process described above and is maximal in the following sense:

- (i)  $0 \notin S$ ,  $|S| \geq 2$  and for any two distinct points  $x_i^k, x_j^k$  in  $S$  we have

$$\text{dist}(x_i^k, x_j^k) \ll \tau^k(S),$$

where  $\tau^k(S) = \text{dist}(S, \Sigma_k \setminus S)$ .

- (ii) For any  $0 \neq x_i^k \in \Sigma_k \setminus S$ ,

$$\text{dist}(x_i^k, 0) \leq C \text{dist}(x_i^k, S)$$

for some constant  $C$ .

For  $S_j$  we define

$$\tau_{S_j}^k = \frac{1}{2} \text{dist}(S_j, \Sigma_k \setminus S_j).$$

Then the process described above proves the main result of this section:

**Proposition 5.2.** *Let  $S_j$  ( $j = 1, \dots, m$ ) be described as above. Then (5-8)–(5-9) hold, where  $B(x_1^k, \tau_S^k)$  is replaced by  $B(x_i^k, \tau_{S_j}^k)$  and  $x_i^k$  is any element in  $S_j$ .*

**6. Proofs of Theorems 1.2, 1.3, 1.5 and 1.6**

In Proposition 5.2, we write  $\Sigma_k = \{0\} \cup S_1 \cup \dots \cup S_N$ . From the construction, the ratio  $|x^k|/|\tilde{x}^k|$  is bounded for any  $x^k, \tilde{x}^k \in S_j$ . Let

$$\|S_j\| = \min_{x^k \in S_j} |x^k|$$

and arrange  $S_j$  by

$$\|S_1\| \leq \|S_2\| \leq \dots \leq \|S_N\|.$$

Assume  $l$  is the largest number such that  $\|S_l\| \leq C \|S_1\|$ . Then  $\|S_l\| \ll \|S_{l+1}\|$ .

We recall the local mass contributed by  $x_j^k \in S_j$  is

$$(\hat{\sigma}_1(B(x_j^k, \tau_j^k)), \hat{\sigma}_2(B(x_j^k, \tau_j^k))) = (m_{1,j}, m_{2,j}), \quad \text{where } m_{1,j}, m_{2,j} \in 2\mathbb{N} \cup \{0\}.$$

Let

$$r_1^k = \frac{1}{2} \|S_1\|.$$

Then we have

$$u_i^k(x) + 2 \log |x| \leq C \quad \text{for } 0 < |x| \leq r_1^k, \quad i = 1, 2.$$

*Proof of Theorem 1.3.* Let

$$\tilde{u}_i^k(x) = u_i^k(x) + 2\alpha_i \log |x|, \quad i = 1, 2.$$

Then (1-5) becomes

$$\Delta \tilde{u}_i^k(x) + \sum_{j=1}^2 k_{ij} |x|^{2\alpha_j} h_j^k(x) e^{\tilde{u}_j^k(x)} = 0, \quad |x| \leq r_1^k, \quad i = 1, 2.$$

Let

$$-2 \log \delta_k = \max_{i \in I} \max_{x \in \bar{B}(0, r_1^k)} \frac{\tilde{u}_i^k}{1 + \alpha_i}, \tag{6-1}$$

and

$$\tilde{v}_i^k(y) = \tilde{u}_i^k(\delta_k y) + 2(1 + \alpha_i) \log \delta_k, \quad |y| \leq r_1^k / \delta_k, \quad i = 1, 2. \tag{6-2}$$

Then  $\tilde{v}_i^k$  satisfies

$$\Delta \tilde{v}_i^k(y) + \sum_{j=1}^2 k_{ij} |y|^{2\alpha_j} h_j^k(\delta_k y) e^{\tilde{v}_j^k(y)} = 0, \quad |y| \leq r_1^k / \delta_k, \quad i = 1, 2. \tag{6-3}$$

We have either

- (a)  $\lim_{k \rightarrow \infty} r_1^k / \delta_k = \infty$ , or
- (b)  $r_1^k / \delta_k \leq C$ .

For case (a), our purpose is to prove a result similar to [Proposition 5.1](#):

(1) At most one component of  $u^k$  has slow decay on  $\partial B(0, r_1^k)$ . As in [Section 5](#), we define

$$\hat{\sigma}_{i,1} = \begin{cases} \lim_{k \rightarrow +\infty} \sigma_i^k(B(0, r_1^k)) & \text{if } u_i^k \text{ has fast decay on } \partial B(0, r_1^k), \\ \lim_{r \rightarrow 0} \lim_{k \rightarrow +\infty} \sigma_i^k(B(0, r r_1^k)) & \text{if } u_i^k \text{ has slow decay on } \partial B(0, r_1^k), \end{cases}$$

(2)  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1})$  satisfies the Pohozaev identity (1-9), and

(3)  $\hat{\sigma}_{i,1} = 2 \sum_{j=1}^2 n_{i,j} \mu_j + 2n_{i,3}$ ,  $n_{i,j} \in \mathbb{Z}$ ,  $i = 1, 2$ ,  $j = 1, 2, 3$ .

We carry out the proof in the discussion of the following two cases.

Case 1: If both  $\tilde{v}_i^k(y)$  converge in any compact set of  $\mathbb{R}^2$ , then  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1})$  can be obtained by the classification theorem in [[Lin et al. 2012](#)]:

$$(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) = \begin{cases} (2\mu_1 + 2\mu_2, 2\mu_1 + 2\mu_2) & \text{for } A_2, \\ (4\mu_1 + 2\mu_2, 4\mu_1 + 4\mu_2) & \text{for } B_2, \\ (8\mu_1 + 4\mu_2, 12\mu_1 + 8\mu_2) & \text{for } G_2. \end{cases}$$

By [Lemma 4.1](#), both  $u_i^k$  have fast decay on  $\partial B(0, r_1^k)$ . So this proves (1)–(3) in this case.

Case 2: Only one  $\tilde{v}_i^k$  converges to  $v_i(y)$  and the other tends to  $-\infty$  uniformly in any compact set. Then it is easy to see that there is  $l_k \ll r_1^k$  such that both  $u_i^k$  have fast decay on  $\partial B(0, l_k)$  and

$$(\sigma_1(B(0, l_k)), \sigma_2(B(0, l_k))) = (2\mu_1, 0) \quad \text{or} \quad (\sigma_1(B(0, l_k)), \sigma_2(B(0, l_k))) = (0, 2\mu_2).$$

So this is the same situation as in the starting point for [Proposition 5.1](#). Then the same argument of [Proposition 5.1](#) leads to the conclusion (1)–(3).

The pair  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1})$  can be calculated by the same method in [Proposition 5.1](#). Then  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) \in \Gamma(\mu_1, \mu_2)$ , which is given in [Section 2](#).

To continue for  $r \in [r_1^k, r_2^k]$ , where  $r_2^k = \frac{1}{2} \|S_{l+1}\|$ , we separate our discussion into two cases also.

Case 1: One component has slow decay on  $\partial B(0, r_1^k)$ , say  $u_1^k$ . Then we scale

$$v_i^k(y) = u_i^k(r_1^k y) + 2 \log r_1^k.$$

By our assumption,  $v_1^k(y)$  converges to  $v_1(y)$  and  $v_2^k(y) \rightarrow -\infty$  in any compact set. Let  $x_j^k \in S_j$  and  $y_j^k = (r_1^k)^{-1} x_j^k \rightarrow p_j$  for  $j \leq l$ . Then  $v_1(y)$  satisfies

$$\Delta v_1 + 2e^{v_1} = 4\pi \tilde{\alpha}_1 \delta_0 + 4\pi \sum_{j=1}^l \tilde{n}_{1,j} \delta_{p_j}, \tag{6-4}$$

where

$$\tilde{n}_{1,j} = -\frac{1}{2} \sum_{i=1}^2 k_{1i} m_{i,j} \quad \text{for some } m_{i,j} \in \mathbb{Z} \quad \text{and} \quad \tilde{\alpha}_1 = \alpha_1 - \frac{1}{2} \sum_{i=1}^2 k_{1i} \hat{\sigma}_{i,1}. \tag{6-5}$$

The finiteness of  $\int_{\mathbb{R}^2} e^{v_1}$  implies

$$\tilde{\alpha}_1 > -1 \quad \text{and} \quad \tilde{n}_{1,j} \geq 0.$$

By [Theorem 2.2](#), we have

$$\frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v_1} dy = 2(\tilde{\alpha}_1 + 1) + 2k_1, \quad \frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v_1} dy = 2k_2, \quad \text{where } k_1, k_2 \in \mathbb{Z}. \quad (6-6)$$

As before, we can choose  $l_k, r_1^k \ll l_k \ll r_2^k$ , such that both  $u_i^k$  have fast decay on  $\partial B(0, l_k)$ . Then the new pair  $(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2})$ , which is defined by

$$\hat{\sigma}_{t,2} = \frac{1}{2\pi} \lim_{k \rightarrow 0} \int_{B(0, l_k)} h_t^k e^{u_t^k}, \quad t = 1, 2,$$

becomes

$$(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2}) = \left( \hat{\sigma}_{1,1} + \frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v_1} + \sum_{j=1}^l m_{1,j}, \hat{\sigma}_{2,1} + \sum_{j=1}^l m_{2,j} \right) \quad (6-7)$$

for  $m_{1j}, m_{2j} \in 2\mathbb{N} \cup \{0\}$ . Using [\(6-6\)](#), we get

$$\hat{\sigma}_{1,2} = \begin{cases} \hat{\sigma}_{1,1} + 2k_2 + \sum_{j=1}^l m_{1,j} & \text{if } \frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v_1} dy = 2k_2, \\ 2\mu_1 + \hat{\sigma}_{1,1} - \sum_{i=1}^2 k_{1i} \hat{\sigma}_{i,1} + 2k_1 + \sum_{j=1}^l m_{1,j} & \text{if } \frac{1}{2\pi} \int_{\mathbb{R}^2} e^{v_1} dy = 2(\tilde{\alpha}_1 + 1) + 2k_1. \end{cases} \quad (6-8)$$

We note that if  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) \in \Gamma(\mu_1, \mu_2)$  and

$$2\mu_1 + \hat{\sigma}_{1,1} - \sum_{i=1}^2 k_{1i} \hat{\sigma}_{i,1} > 0,$$

then

$$\left( 2\mu_1 + \hat{\sigma}_{1,1} - \sum_{i=1}^2 k_{1i} \hat{\sigma}_{i,1}, \hat{\sigma}_{2,1} \right) \in \Gamma(\mu_1, \mu_2).$$

Let  $(\sigma_1^*, \sigma_2^*) = (2\mu_1 + \hat{\sigma}_{1,1} - \sum_{i=1}^2 k_{1i} \hat{\sigma}_{i,1}, \hat{\sigma}_{2,1})$ . We can write

$$(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2}) = (\sigma_1^* + m_1, \sigma_2^* + m_2), \quad (6-9)$$

with  $(\sigma_1^*, \sigma_2^*) \in \Gamma(\mu_1, \mu_2)$  and  $m_1, m_2 \in 2\mathbb{Z}$ .

Case 2: If both  $u_i^k$  have fast decay on  $\partial B(0, r_1^k)$ , then they have fast decay on  $\partial B(0, cr_1^k)$ , where we choose  $c$  bounded such that  $\bigcup_{j=1}^l S_j \subset B(0, \frac{1}{2} cr_1^k)$ . Then the new pair  $(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2})$  becomes

$$(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2}) = \left( \hat{\sigma}_{1,1} + \sum_{j=1}^l m_{1,j}, \hat{\sigma}_{2,1} + \sum_{j=1}^l m_{2,j} \right) \quad \text{for } m_{1,j}, m_{2,j} \in 2\mathbb{Z}. \quad (6-10)$$

Hence, in this case we can also write

$$(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2}) = (\sigma_1^* + m_1, \sigma_2^* + m_2), \quad (6-11)$$

with  $(\sigma_1^*, \sigma_2^*) = (\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) \in \Gamma(\mu_1, \mu_2)$  and  $m_1, m_2 \in 2\mathbb{Z}$ . Set  $cr_1^k = l_k$ . Then we can continue our process starting from  $l_k$ . After finitely many steps, we can prove that at most one component of  $u^k$  has slow decay on  $\partial B(0, 1)$  and their local masses have the expression in [\(3\)](#).

For case (b), i.e.,  $r_1^k/\delta_k \leq C$ , first  $\tilde{v}_i^k \leq 0$  implies  $|y|^{2\alpha_j} h_j^k(\delta_k y) e^{\tilde{v}_j^k} \leq C$  on  $B(0, r_1^k/\delta_k)$ . Then the fact that  $\tilde{v}_i^k$  has bounded oscillation on  $\partial B(0, r_1^k/\delta_k)$  further gives

$$\tilde{v}_i^k(x) = \bar{v}_i^k(\partial B(0, r_1^k/\delta_k)) + O(1) \quad \text{for all } x \in B(r_1^k/\delta_k),$$

where  $\bar{v}_i^k(\partial B(0, r_1^k/\delta_k))$  stands for the average of  $\tilde{v}_i^k$  on  $\partial B(0, r_1^k/\delta_k)$ . Direct computation shows that

$$\int_{B(0, r_1^k)} h_i^k e^{u_i^k} dx = \int_{B(0, r_1^k/\delta_k)} |y|^{2\alpha_i} h_i^k(\delta_k y) e^{\tilde{v}_i^k(y)} dy = O(1) e^{\bar{v}_i^k(\partial B(0, r_1^k/\delta_k))}.$$

Thus if  $\bar{v}_i^k(\partial B(0, r_1^k/\delta_k)) \rightarrow -\infty$ , we get  $\int_{B(0, r_1^k)} h_i^k e^{u_i^k} dx = o(1)$ . On the other hand, we note that  $\bar{v}_i^k(\partial B(0, r_1^k/\delta_k)) \rightarrow -\infty$  is equivalent to  $u_i^k$  having fast decay on  $\partial B(0, r_1^k)$ . Consequently  $\hat{\sigma}_{i,1} = 0$  if  $u_i^k$  has fast decay on  $\partial B(0, r_1^k)$ . So if both components have fast decay on  $\partial B(0, r_1^k)$  we have  $(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) = (0, 0)$ .

If some component of  $u^k$  has slow decay, say  $u_2^k$ , according to the definition of  $\hat{\sigma}_{2,1}$ , we have

$$\begin{aligned} \hat{\sigma}_{2,1} &= \lim_{r \rightarrow 0} \lim_{k \rightarrow +\infty} \sigma_2^k(B(0, rr_1^k)) = \frac{1}{2\pi} \lim_{r \rightarrow 0} \lim_{k \rightarrow +\infty} \int_{B(0, rr_1^k)} h_2^k e^{u_2^k} dx \\ &= \frac{1}{2\pi} \lim_{r \rightarrow 0} \lim_{k \rightarrow +\infty} \int_{B(0, rr_1^k/\delta_k)} |y|^{2\alpha_2} h_2^k(\delta_k y) e^{\tilde{v}_2^k(y)} dy = 0, \end{aligned} \tag{6-12}$$

where we used  $|y|^{2\alpha_2} h_2^k(\delta_k y) e^{\tilde{v}_2^k} \leq C$  on  $B(0, r_1^k/\delta_k)$ . Then we still get

$$(\hat{\sigma}_{1,1}, \hat{\sigma}_{2,1}) = (0, 0).$$

Now we can continue our discussion as in case (a) and [Theorem 1.3](#) is proved completely. □

Next, we shall prove [Theorem 1.5](#), that is,  $\Sigma_k = \{0\}$ , by way of contradiction. Suppose  $\Sigma_k$  has points other than 0. Using the notation from the beginning of this section, we have

$$\Sigma_k = \{0\} \cup S_1 \cup \dots \cup S_N.$$

Now suppose  $r_1^k/\delta_k \rightarrow \infty$  as  $k \rightarrow \infty$ . Let  $(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2})$  be the local masses defined by [\(6-7\)](#) for one of the components  $u_i^k$  having slow decay on  $\partial B(0, r_1^k)$  or by [\(6-10\)](#) for both components having fast decay on  $\partial B(0, r_1^k)$ . We summarize the results in the following:

- (i)  $\hat{\sigma}_{i,2} = \sigma_i^* + m_i$ , where  $(\sigma_1^*, \sigma_2^*) \in \Gamma(\mu_1, \mu_2)$  and  $m_i, i = 1, 2$ , are even integers.
- (ii) Both pairs  $(\sigma_1^*, \sigma_2^*)$  and  $(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2})$  satisfy the Pohozaev identity.

Based on the description above, we now present the proof of [Theorem 1.5](#).

*Proof of Theorem 1.5.* From the discussion above, we have

$$(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2}) = (\sigma_1^* + m_1, \sigma_2^* + m_2).$$

We note that the conclusion of [Theorem 1.5](#) is equivalent to proving  $m_i = 0, i = 1, 2$ . In order to prove this we first observe that both  $(\hat{\sigma}_{1,2}, \hat{\sigma}_{2,2})$  and  $(\sigma_1^*, \sigma_2^*)$  satisfy the P.I.

$$k_{21}\sigma_1^2 + k_{12}k_{21}\sigma_1\sigma_2 + k_{12}\sigma_2^2 = 2k_{21}\mu_1\sigma_1 + 2k_{12}\mu_2\sigma_2. \tag{6-13}$$

Thus we can write

$$k_{21}(\sigma_1^*)^2 + k_{12}k_{21}\sigma_1^*\sigma_2^* + k_{12}(\sigma_2^*)^2 = 2k_{21}\mu_1\sigma_1^* + 2k_{12}\mu_2\sigma_2^*, \tag{6-14}$$

and

$$\begin{aligned} k_{21}(\sigma_1^* + m_1)^2 + k_{12}k_{21}(\sigma_1^* + m_1)(\sigma_2^* + m_2) + k_{12}(\sigma_2^* + m_2)^2 \\ = 2k_{21}\mu_1(\sigma_1^* + m_1) + 2k_{12}\mu_2(\sigma_2^* + m_2). \end{aligned} \tag{6-15}$$

It is easy to obtain the following from (6-15) and (6-14):

$$\begin{aligned} 2k_{21}m_1\sigma_1^* + k_{12}k_{21}m_2\sigma_1^* + k_{12}k_{21}m_1\sigma_2^* + 2k_{12}m_2\sigma_2^* \\ = 2k_{21}m_1\mu_1 + 2k_{12}m_2\mu_2 - (k_{21}m_1^2 + k_{12}k_{21}m_1m_2 + k_{12}m_2^2). \end{aligned} \tag{6-16}$$

Since  $(\sigma_1^*, \sigma_2^*) \in \Gamma(\mu_1, \mu_2)$ , we set

$$\sigma_1^* = l_{1,1}\mu_1 + l_{1,2}\mu_2, \quad \sigma_2^* = l_{2,1}\mu_1 + l_{2,2}\mu_2.$$

Then we can rewrite (6-16) as

$$\begin{aligned} (2k_{21}l_{1,1}m_1 + k_{12}k_{21}l_{2,1}m_1 - 2k_{21}m_1 + 2k_{12}l_{2,1}m_2 + k_{12}k_{21}l_{1,1}m_2)\mu_1 \\ + (2k_{21}l_{1,2}m_1 + k_{12}k_{21}l_{2,2}m_1 + 2k_{12}l_{2,2}m_2 + k_{12}k_{21}l_{1,2}m_2 - 2k_{12}m_2)\mu_2 \\ + (k_{21}m_1^2 + k_{12}k_{21}m_1m_2 + k_{12}m_2^2) = 0. \end{aligned} \tag{6-17}$$

Since  $\mu_1, \mu_2$  and 1 are linearly independent, the coefficients of  $\mu_1$  and  $\mu_2$  must vanish. Equivalently we have

$$\begin{pmatrix} 2k_{21}l_{1,1} + k_{12}k_{21}l_{2,1} - 2k_{21} & 2k_{12}l_{2,1} + k_{12}k_{21}l_{1,1} \\ 2k_{21}l_{1,2} + k_{12}k_{21}l_{2,2} & 2k_{12}l_{2,2} + k_{12}k_{21}l_{1,2} - 2k_{12} \end{pmatrix} \begin{pmatrix} m_1 \\ m_2 \end{pmatrix} = 0. \tag{6-18}$$

Let  $M_K$  be the coefficient matrix

$$M_K = \begin{pmatrix} 2k_{21}l_{1,1} + k_{12}k_{21}l_{2,1} - 2k_{21} & 2k_{12}l_{2,1} + k_{12}k_{21}l_{1,1} \\ 2k_{21}l_{1,2} + k_{12}k_{21}l_{2,2} & 2k_{12}l_{2,2} + k_{12}k_{21}l_{1,2} - 2k_{12} \end{pmatrix}.$$

Our goal is to show that  $M_k$  is nonsingular, which immediately implies  $m_1 = m_2 = 0$  and completes the proof of [Theorem 1.5](#). The proof of the nonsingularity of  $M_k$  is divided into the following three cases.

Case 1:  $K = A_2$ . Then we can write (6-18) as

$$\begin{pmatrix} 2l_{1,1} - l_{2,1} - 2 & 2l_{2,1} - l_{1,1} \\ 2l_{1,2} - l_{2,2} & 2l_{2,2} - l_{1,2} - 2 \end{pmatrix} \begin{pmatrix} m_1 \\ m_2 \end{pmatrix} = 0. \tag{6-19}$$

We note that

$$(l_{1,1}, l_{1,2}, l_{2,1}, l_{2,2}) \in \{(2, 0, 0, 0), (0, 0, 0, 2), (2, 2, 0, 2), (2, 0, 2, 2), (2, 2, 2, 2)\}.$$

Then it is easy to see that  $M_K$  is nonsingular when  $(l_{1,1}, l_{1,2}, l_{2,1}, l_{2,2})$  belongs the above set.

Case 2:  $K = B_2$ . Then we can write (6-18) as

$$\begin{pmatrix} 2l_{1,1} - l_{2,1} - 2 & l_{2,1} - l_{1,1} \\ 2l_{1,2} - l_{2,2} & l_{2,2} - l_{1,2} - 1 \end{pmatrix} \begin{pmatrix} m_1 \\ m_2 \end{pmatrix} = 0. \tag{6-20}$$



We note that

$$(l_{1,1}, l_{1,2}, l_{2,1}, l_{2,2}) \in \{(2, 0, 0, 0), (2, 0, 4, 2), (4, 2, 4, 2), (0, 0, 0, 2), (2, 2, 0, 2), (2, 2, 4, 4), (4, 2, 4, 4)\}$$

From the above set, we can see that  $4 \mid (l_{2,1} - l_{1,1})(2l_{1,2} - l_{2,2})$ . As a result, if the determinant of  $M_K$  is 0, we have to make  $4 \mid (2l_{1,1} - l_{2,1} - 2)$ , which forces  $l_{2,1} \equiv 2 \pmod{4}$ . However, this is impossible according to the above list. Thus  $M_k$  is nonsingular in this case.

Case 3:  $K = G_2$ . Then we can write (6-18) as

$$\begin{pmatrix} 6l_{1,1} - 3l_{2,1} - 6 & 2l_{2,1} - 3l_{1,1} \\ 6l_{1,2} - 3l_{2,2} & 2l_{2,2} - 3l_{1,2} - 2 \end{pmatrix} \begin{pmatrix} m_1 \\ m_2 \end{pmatrix} = 0. \tag{6-21}$$

We note that

$$(l_{1,1}, l_{1,2}, l_{2,1}, l_{2,2}) \in \{(2, 0, 0, 0), (2, 0, 6, 2), (6, 2, 6, 2), (6, 2, 12, 6), (8, 4, 12, 6), (8, 4, 12, 8), (0, 0, 0, 2), (2, 2, 0, 2), (2, 2, 6, 6), (6, 4, 6, 6), (6, 4, 12, 8)\}.$$

From the above list, we have  $3 \mid l_{2,1}$ ; then we get  $9 \mid (2l_{2,1} - 3l_{1,1})(6l_{1,2} - 3l_{2,2})$ . On the other hand, we see that

$$l_{1,1} \equiv 0, 2 \pmod{3} \quad \text{and} \quad l_{2,2} \equiv 0, 2 \pmod{3},$$

which implies  $(6l_{1,1} - 3l_{2,1} - 6)(2l_{2,2} - 3l_{1,2} - 2)$  is not multiple of 9; therefore we have the determinant of  $M_K$  is not zero. Thus  $M_k$  is nonsingular when  $K = G_2$ .

**Theorem 1.5** is established. □

Finally we prove Theorems 1.2 and 1.6.

*Proof of Theorems 1.2 and 1.6.* Suppose there exists a sequence of blowup solutions  $(u_1^k, u_2^k)$  of (1-2) with  $(\rho_1, \rho_2) = (\rho_1^k, \rho_2^k)$ . First, we prove **Theorem 1.2**. From the previous discussion of this section, we get that at least one component (say  $u_1^k$ ) of  $u^k$  has fast decay on a small ball  $B$  near each blowup point  $q$ , which means  $u_1^k(x) \rightarrow -\infty$  if  $x \notin S$  and  $x$  is not a blowup point. Hence **Theorem 1.2** holds.

Because the mass distribution of  $u_1^k$  concentrates as  $k \rightarrow +\infty$ , we get that  $\lim_{k \rightarrow +\infty} \rho_1^k$  is equal to the sum of the local mass  $\sigma_1$  at a blowup point  $q$ , which implies  $\rho_1 \in \Gamma_1$ , a contradiction to the assumption. Thus, we finish the proof of **Theorem 1.6**. □

### Acknowledgement

We would like to thank the referee for perusing the whole article and for many excellent suggestions.

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Received 3 Nov 2016. Revised 17 Aug 2017. Accepted 5 Dec 2017.

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Analysis & PDE (ISSN 1948-206X electronic, 2157-5045 printed) at Mathematical Sciences Publishers, 798 Evans Hall #3840, c/o University of California, Berkeley, CA 94720-3840, is published continuously online. Periodical rate postage paid at Berkeley, CA 94704, and additional mailing offices.

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APDE peer review and production are managed by EditFlow® from MSP.

PUBLISHED BY

 **mathematical sciences publishers**  
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Volume 11 No. 4 2018

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