MULTI-VORTEX TRAVELING WAVES FOR THE GROSS-PITAEVSKII EQUATION AND THE ADLER-MOSER POLYNOMIALS

YONG LIU AND JUNCHENG WEI

Abstract. For $N \leq 34$, we construct traveling waves with small speed for the Gross-Pitaevskii equation, by gluing $N(N+1)/2$ pairs of degree $\pm 1$ vortices of the Ginzburg-Landau equation. The location of these vortices is symmetric in the plane and determined by the Adler-Moser polynomials, which has its origin in the study of Calogero-Moser system and rational solutions of the KdV equation. The construction still works for $N > 34$, under the additional assumption that the corresponding Adler-Moser polynomial has no repeated root. It is expected that this assumption holds for any $N \in \mathbb{N}$.

1. Introduction

The Gross-Pitaevskii equation (GP equation) arises as a model equation in Bose-Einstein condensate. It reads as

$$i\partial_t \Phi = \Delta \Phi + \Phi \left(1 - |\Phi|^2\right), \text{ in } \mathbb{R}^2,$$

where $\Phi$ is complex valued. Throughout the paper, $i$ will represent the imaginary unit. For traveling wave solutions of the form $U(x, y - \varepsilon t)$, GP equation becomes

$$-i\varepsilon \partial_y U = \Delta U + U \left(1 - |U|^2\right), \text{ in } \mathbb{R}^2. \tag{1.2}$$

We would like to construct multi-vortex type solutions of (1.2) when the speed $\varepsilon$ is close to zero. If $\varepsilon = 0$, then the above equation reduces to the well-known Ginzburg-Landau equation

$$\Delta U + U \left(1 - |U|^2\right) = 0, \text{ in } \mathbb{R}^2. \tag{1.3}$$

We shall use $(r, \theta)$ to denote the polar coordinate of $\mathbb{R}^2$. For each $d \in \mathbb{Z}\setminus\{0\}$, it is known that the Ginzburg-Landau equation has a degree $d$ vortex solution, of the form $e^{id\theta} S_d (r)$. The function $S_d$ is real valued and vanishes exactly at $r = 0$. It satisfies

$$-S''_d - \frac{1}{r} S'_d + \frac{d^2}{r^2} S_d = S_d \left(1 - S^2_d\right), \text{ in } (0, +\infty).$$

This equation has a unique solution $S_d$ satisfying $S_d(0) = 0$ and $S_d(+\infty) = 1$ and $S'(r) > 0$. See [29,33] for a proof. The “standard” degree $\pm 1$ solutions $S_1 (r) e^{i\theta}$ are global minimizers of the energy functional (For uniqueness of the global minimizer, see [34,37]). When $|d| > 1$, these standard vortices are unstable([9,33]). It is also worth mentioning that for $|d| > 1$, the uniqueness of degree $d$ vortex in the class of solutions with degree $d$ is still an open problem. We refer to [6,35,36] and the references therein for more discussion on the Ginzburg-Landau equation.
The constant 1 is a solution to the equation (1.2). We are interested in the solution $U$ which satisfies

$$U(z) \to 1, \text{ as } |z| \to +\infty.$$ 

The existence or nonexistence of solutions to (1.2) with this asymptotic behavior has been extensively studied in the literature. Jones, Putterman, Roberts ([24, 25]) studied it from the physical point of view, both in dimension two and three. It turns out the existence of solutions is related to the traveling speed $\varepsilon$. When $\varepsilon \geq \sqrt{2}$ (the sound speed in this context), nonexistence of traveling wave with finite energy is proved by Gravejat in [21, 22]. On the other hand, for $\varepsilon \in (0, \sqrt{2})$, the existence of travelling waves as constrained minimizer is studied by Bethuel, Gravejat, Saut [9, 11], by variational arguments. For $\varepsilon$ close to 0, these solutions have two vortices. The existence issue in higher dimension is studied [10, 14, 15]. We also refer to [8] for a review on this subject. Recently, Chiron-Scheid [13] performed numerical simulation on this equation. Among other things, their results indicate the existence of higher energy traveling waves. We also mention that as $\varepsilon$ tends to $\sqrt{2}$, a suitable rescaled traveling waves will converge to solutions of the KP-I equation ([7]), which is classical integrable system. In a forthcoming paper, we will construct transonic traveling waves based on the lump solution of the KP-I equation.

Another motivation for studying (1.2) arises in the study of superfluids passing an obstacle. Equation (1.2) is the limiting equation in the search of vortex nucleation solution. We refer to recent paper [31] for references and derivations.

To simplify notations, we write the degree $\pm 1$ vortex solutions of the Ginzburg-Landau equation as

$$v_+ = e^{i\theta} S_1(r), v_- = e^{-i\theta} S_1(r).$$

In this paper, we construct new traveling waves for $\varepsilon$ close to 0, using $v_+, v_-$ as basic blocks. Our main result is

**Theorem 1.1.** For each $N \leq 34$, there exists $\varepsilon_0 > 0$, such that for all $\varepsilon \in (0, \varepsilon_0)$, the equation (1.2) has a solution $U_\varepsilon$ which has the form

$$U_\varepsilon = \prod_{k=1}^{N(N+1)/2} \left( v_+ (z - \varepsilon^{-1} p_k) v_- (z - \varepsilon^{-1} q_k) \right) + o(1),$$

where $p_k, k = 1, \ldots, N(N+1)/2$ are the roots of the Adler-Moser polynomial $A_n$ defined in the next section ($n = \frac{N(N+1)}{2}$), and $q_k = -\overline{p_k}$.

**Remark 1.2.** $N = 1$ corresponds to the two vortices solutions constructed by variational method ([11]) as well as reduction method ([30]). For large $N$ this corresponds to higher energy solutions which have been observed numerically [13].

**Remark 1.3.** The condition $N \leq 34$ is only technical. In this case, we can verify numerically that the Adler-Moser polynomial $A_n$ has no repeated roots. (This is equivariant to $A_n$ and $A_{n-1}$ have no common roots.) Possibly, there are other numerically ways to verify this for large $N$ (using the recursive identity to compute the Adler-Moser polynomial, instead of computing the Wronskian), but we will not pursue this here. We conjecture that the Adler-Moser polynomial (as constructed in this paper) has only simple roots for all $N$. 


Our method is based on finite dimensional Lyapunov-Schmidt reduction. We show that the existence of multi-vortex solutions is essentially reduced to the study of the non-degeneracy of a symmetric vortex-configuration. To show this non-degeneracy, we use the theory of Adler-Moser polynomials and the Darboux transformation. An interesting feature of the solutions in Theorem 1.1 is that the vortex location has a ring-like structure. The emergence of this remarkable property still remains mysterious.

In Section 2, we introduce the Adler-Moser polynomials and prove the non-degeneracy of the symmetric configuration. In Section 3, we recall the linear theory of the degree one vortex of the Ginzburg-Landau equation. In Section 4, we use Lyapunov-Schmidt reduction to glue the vortices together and get a traveling wave solution for $\varepsilon$ small enough.

Acknowledgement The research of J. Wei is partially supported by NSERC of Canada. Part of this work is finished while the first author is visiting the University of British Columbia in 2017. He thanks the institute for the financial support. Both authors thank Professor Fanghua Lin for stimulating discussions and suggestions.

2. Vortex location and the Adler-Moser polynomials

Adler-Moser [1] have studied a set of polynomials corresponding to rational solutions of the KdV equation. It turns out that these polynomials have deep connections to the vortex dynamics with logarithmic interaction energy. This connection is first observed in [5], and later studied in [2–4,16,26]. (See the reference therein.) Surprisingly, to authors’ knowledge, except the papers mentioned above, it seems that this relation has not been much explored for the corresponding PDEs. One of our aims in this paper is to fill this gap.

In this section, we will first recall some basic facts about these polynomials and then analyze some of their properties, which will be used in our construction of the traveling wave for the GP equation.

Let $p_1,...,p_n$ be the position of the positive vortices and $q_1,...,q_m$ be that of the negative ones. Let $\mu \in \mathbb{R}$ be a fixed parameter. As we will see later, the vortex location of the traveling waves will be determined by the following systems of equations

$$\left\{ \begin{array}{l}
\sum_{j \neq \alpha} \frac{1}{p_\alpha - p_j} - \sum_{j} \frac{1}{p_\alpha - q_j} = \mu, \text{ for } \alpha = 1,...,n, \\
\sum_{j \neq \alpha} \frac{1}{q_\alpha - q_j} - \sum_{j} \frac{1}{q_\alpha - p_j} = -\mu, \text{ for } \alpha = 1,...,m.
\end{array} \right. \tag{2.1}$$

Adding all these equation together, we find that if $\mu \neq 0$, then $m = n$. (In the case of $\mu = 0$, this is no longer true). That is, the number of positive vortices must equal to that of the negative vortices. Solutions of this system (see for instances [4]) are related to the Adler-Moser polynomials. To explain this, let us define the generating polynomials

$$P(z) = \prod_j (z - p_j), \quad Q(z) = \prod_j (z - q_j).$$

If $p_j, q_j$ satisfy (2.1), then we have (see equation (68) of [4])

$$P''Q - 2P'Q' + P Q'' = 2\mu (P'Q - PQ'). \tag{2.2}$$
This equation is usually called generalized Tkachenko equation. Setting \( \psi (z) = \frac{P}{Q} e^{\mu z} \), we derive from (2.2) that
\[
\psi'' + 2 (\ln Q)'' \psi = \mu^2 \psi.
\]

This is a one dimensional Schrödinger equation with the potential \( 2 (\ln Q)'' \). It is well known that this equation appears in the Lax pair of the KdV equation. Hence equation (2.2) is naturally related to the theory of integrable systems.

For any \( z \in \mathbb{C} \), we use \( \bar{z} \) to denote its complex conjugate. To simplify the notation, we also write \( -\bar{z} \) as \( z^* \). Note that this is just the reflection of \( z \) across the \( y \) axis. Let \( K = (k_i, \ldots) \), where \( k_i \) are complex parameters. Following [16], we define functions \( \theta_n \), depending on \( K \), by
\[
\theta_n (z; K) := \sum_{j=0}^{n} \frac{k_j \lambda^{2j-1}}{2j-1},
\]
where \( \lambda = z \lambda - \sum_{j=2}^{\infty} \frac{k_j \lambda^{2j-1}}{2j-1} \). (2.3)

Note that \( \theta_n \) is a degree \( n \) polynomial in \( z \) and \( \theta'_n + 1 = \theta_n \). Explicitly, \( \theta_1 (z; K) = z \),
\[
\theta_2 (z; K) = \frac{\lambda}{6} - \frac{k_2}{3},
\]
\[
\theta_3 (z; K) = \frac{\lambda^3}{5} - \frac{k_2}{6} z^2 + \frac{1}{120} z^5.
\]

Let \( c_n = \prod_{j=1}^{n} (2j + 1)^{n-j} \). For each \( n \in \mathbb{N} \), the Adler-Moser polynomials are then defined by
\[
\Theta_n (z, K) := c_n W (\theta_1, \theta_3, \ldots, \theta_{2n-1}), \quad (2.4)
\]
where \( W (\theta_1, \theta_3, \ldots, \theta_{2n-1}) \) is the Wronskian of \( \theta_1, \ldots, \theta_{2n-1} \). In particular, the degree of \( \Theta_n \) is \( n(n+1)/2 \). The constant \( c_n \) is chosen such that the leading coefficient of \( \Theta_n \) is 1. The first three Adler-Moser polynomials are \( \Theta_1 (z, K) = z \), \( \Theta_2 (z, K) = z^3 + k_2 \), and
\[
\Theta_3 (z, K) = z^6 + 5k_2 z^3 - 9k_3 z - 5k_4^2.
\]
Note that this definition is slightly different from that of Adler-Moser [1]. (The parameter \( \tau_i \) in that paper is different from \( k_i \) here.)

Let \( \mu \) be another parameter. The modified Adler-Moser polynomial \( \tilde{\Theta} \) is defined by
\[
\tilde{\Theta}_n (z, \mu, K) := c_n e^{-\mu z} W (\theta_1, \theta_3, \ldots, \theta_{2n-1}, e^{\mu z}) \quad (2.5)
\]
where \( \tilde{K} = (k_2 + \mu^{-3}, k_3 + \mu^{-5}, \ldots, k_n + \mu^{-2n+1}) \). It is still a polynomial in \( z \) with degree \( n(n+1)/2 \). We observe that for a given \( \mu \), \( \Theta_n \) depends on \( n-1 \) complex parameters \( k_2, \ldots, k_n \). This together with the translation in \( z \) gives us a total of \( n \) complex parameters.

The following result, stated without proof in [16], will play an important role in our later analysis.

**Lemma 2.1.** The Adler-Moser and modified Adler-Moser polynomials are related by
\[
\tilde{\Theta}_n (z, \mu, K) = \mu^n \Theta_n \left( z - \mu^{-1}, \tilde{K} \right),
\]
Proof. We sketch the proof for completeness. First of all, direction computation shows that
\[ \sum_{n=0}^{+\infty} \theta_n(z; K)\lambda^n = \sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} \sum_{n=0}^{+\infty} \theta_n(z - \mu^{-1}; \tilde{K})\lambda^n. \]
From this we obtain
\[ \mu^{-1}\sum_{n=0}^{+\infty} \theta_{n-1}(z; K)\lambda^n = \mu^{-1}\lambda \sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} \sum_{n=0}^{+\infty} \theta_n(z - \mu^{-1}; \tilde{K})\lambda^n. \]
Hence using the fact that \( \theta_n' = \theta_{n-1} \), we get
\[
\sum_{n=0}^{+\infty} \left( \theta_n(z; K) - \mu^{-1}\theta_n'(z; K) - \theta_n(z - \mu^{-1}; \tilde{K}) \right)\lambda^n = \left( \sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} - 1 - \mu^{-1}\lambda \sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} \right) \sum_{n=0}^{+\infty} \theta_n(z - \mu^{-1}; \tilde{K})\lambda^n.
\]
We observe that
\[
\sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} - 1 - \mu^{-1}\lambda \sqrt{\frac{1 + \mu^{-1}\lambda}{1 - \mu^{-1}\lambda}} = \sqrt{1 - \mu^{-2}\lambda^2} - 1.
\]
The Taylor expansion of this function contains only even powers of \( \lambda \). Hence for odd \( \alpha \), \( \theta_n(z; K) - \mu^{-1}\theta_n'(z; K) - \theta_n(z - \mu^{-1}; \tilde{K}) \) can be written as a linear combination of \( \theta_k(z - \mu^{-1}; \tilde{K}) \) with \( k \) being odd. The desired identity then follows. \( \square \)

The next result, which essentially follows from Crum type theorem, reveals the relation of the Adler-Moser polynomial with the vortex dynamics \( \text{[4], see also Theorem 3.3 in [16]} \).

Lemma 2.2. The functions \( Q = \Theta_n(z, K), P = \tilde{\Theta}_n(z, \mu, K) \) satisfy (2.2).

Note that a general degree \( m \) term in \( \theta_n \) has the form \( k_j^2 \cdots k_j^j z^m \). We define its index to be \( (-1)^{j^2 + \cdots + j^j + m} \). We now prove the following

Lemma 2.3. For each term of \( \theta_{2n+1} \), its index is \(-1\).

Proof. Let \( k_j^2 \cdots k_j^j z^m \) be a degree \( m \) term in \( \theta_{2n+1} \). By Taylor expansion of the generating function and using the fact that \( 2n + 1 \) is odd, this term comes from functions of the form,
\[
\frac{1}{\alpha!} \left( z\lambda - \sum_{j=2}^{\infty} \frac{k_j\lambda^{2j-1}}{2j-1} \right)^\alpha,
\]
where \( \alpha \) is an odd integer. Hence the total degree of \( k_j \) is \( \alpha - m \). Then the index is \( (-1)^{\alpha-2m} = -1 \). \( \square \)

Lemma 2.4. For each term of \( \Theta_n \), its index is equal to \( (-1)^{\frac{n(n+1)}{2}} \).
Proof. Let us consider a typical term of $\Theta_n$, say $\theta_1\theta_2^{(n-1)}$, where the notation $(\cdot)$ represents taking derivatives. By Lemma 2.6, terms in $\theta_k^{(j)}$ have index $(-1)^{1+j}$. Hence the index of terms in $\theta_1\theta_2^{(n-1)}$ is $(-1)^{1+2+\cdots+n} = (-1)^{\frac{n(n+1)}{2}}$. This finishes the proof.

Now we introduce the notation

$$\Theta_n(t,z) := \Theta_n(z-t,K).$$

For any polynomial $\phi(z)$, we use $R(\phi)$ to denote the set of roots of $\phi$. We have the following

Lemma 2.5. Suppose $\mu$ is a real number. Assume $t = -\frac{\mu}{2}$ and $k_j = -\frac{1}{2} \mu^{2j-1}$ for $j = 2,\ldots$. Then

$$(\Theta_n,t(z,K))^* = (-1)^{\frac{a(n+1)}{2}} \Theta_n(t,z,K).$$

As a consequence, in this case, the reflection of $R(\Theta_n,t(z,K))$ across the $y$ axis is $R(\Theta_n,t(z,K))$, and $R(\Theta_n,t(z,K))$ is invariant respect to the reflection across the $x$ axis.

Proof. By Lemma 2.4, for each term $f = k_1^{i_1} \cdots k_j^{i_j} (z-t)^m$ of the function $\Theta_n(t,z,K)$, we have in $\tilde{\Theta}_n(t,z,K)$ a corresponding term $\tilde{k}_1^{i_1} \cdots \tilde{k}_j^{i_j} (z^* - t - \mu)^m$, denoted by $g$. By the choice of $k_j$, we know that

$$\tilde{k}_j = -k_j.$$

By Lemma 2.4, the index of $k_1^{i_1} \cdots k_j^{i_j} z^m$ is $(-1)^{\frac{a(n+1)}{2}}$. Hence using the fact that $\mu$ is real, we get

$$f^* = -k_1^{i_1} \cdots k_j^{i_j} (-z^* - t)^m$$

$$= (-1)^{1+i_1+\cdots+i_j+m} \tilde{k}_1^{i_1} \cdots \tilde{k}_j^{i_j} (z^* + t)^m$$

$$= (-1)^{\frac{a(n+1)}{2}+1} g.$$

This completes the proof. \qed

Taking for example $\mu = 2, t = -1, k_2 = -4, k_3 = -16$, we get $\Theta_2,t(z,K) = (z+1)^3 - 4$. It has one real root and a pair of conjugate roots, forming a regular triangle, and given numerically by

$$p_1 = 0.5874, p_2 = 1.7937 - 1.3747i, p_3 = 1.7937 + 1.3747i. \quad (2.7)$$

In the sequel, for simplicity, we shall choose $\mu = 1$. Then $t = k_j = -\frac{1}{2}$. Let us then denote the corresponding polynomial $\Theta_n(t,z,K)$ by $A_n(z)$. Then $\Theta_n(z,K)$ is equal to $A_n(-z)$, which we denote by $B_n(z)$.

Since our traveling wave solutions will roughly speaking have vortices at the roots of $A_n$, it is natural to ask that whether all the roots of $A_n$ are simple. This question seems to be nontrivial.

Lemma 2.6. Let $P(z), Q(z)$ be two polynomials satisfying

$$P''Q - 2P'Q' + P Q'' = 2\mu (P'Q - PQ'),$$

or

$$P''Q - 2P'Q' + PQ'' = 0.$$
Suppose \( P(\xi) = 0, Q(\xi) \neq 0, \) for some \( \xi. \) Then \( \xi \) is a simple root of \( P. \)

**Proof.** We prove the lemma assuming (2.8). The case of (2.9) is similar.

Suppose \( \xi \) is root of \( P \) with multiplicity \( k \geq 2. \) We have

\[
P''Q = 2P'Q' - PQ'' + 2\mu (P'Q - PQ').
\]

Then \( \xi \) is a root of the right hand side polynomial with multiplicity at least \( k - 1. \)
But its multiplicity in \( P''Q \) is \( k - 2. \) This is a contradiction. \( \square \)

**Lemma 2.7.** Suppose \( P(z), Q(z) \) are two polynomials satisfying (2.8) or (2.9). Let \( \xi \) be a common root of \( P \) and \( Q. \) Assume \( \xi \) is a simple root of \( Q. \) Then \( \xi \) can not be a simple root of \( P. \)

**Proof.** We prove this lemma assuming (2.9). The case of (2.8) is similar.

Assume to the contrary that \( \xi \) is a simple root of \( P. \) Then

\[
2P'(\xi) Q'(\xi) \neq 0.
\]

But this contradicts with the equation (2.9). This finishes the proof. \( \square \)

We introduce the following assumption:

(A). The polynomials \( A_n(z) \) and \( A_{n-1}(z) \) have no common roots.

**Lemma 2.8.** Suppose the assumption (A) holds. Then \( A_n \) has no repeated roots.
Moreover, \( A_n(\xi) \) and \( A_n(-\xi) \) have no common roots.

**Proof.** We know that the sequence of Adler-Moser polynomials satisfy the following recursion relation

\[
A''_n A_{n-1} - 2A'_n A'_{n-1} + A_n A''_{n-1} = 0, \text{ for any } n. \tag{2.10}
\]

By Lemma 2.6, any root of \( A_n \) is a simple root. Similarly, any root of \( A_n(-z) \) is a simple root.

Now suppose to the contrary that \( \xi \) is a common root of \( A_n(z) \) and \( A_n(-z). \)
Letting \( B_n(x) = A_n(-x), \) we have

\[
A''_n B_n - 2A'_n B'_n + A_n B''_n = 2\mu (A'_n B_n - A_n B'_n).
\]

Then by Lemma 2.7, either \( \xi \) is a repeated root of \( A_n(z), \) or it is a repeated root of \( A_n(-z). \) This is a contradiction. \( \square \)

2.1. **Linearization of the symmetric configuration.** Our construction of traveling wave requires that the vortex configuration we found is nondegenerate in the symmetric setting (in the sense of Lemma 2.5). For small number of vortices, this can be verified directly. To explain this, we now consider the case of \( n = 2. \) Let \( p_1, p_2, p_3 \) be the three roots of the Adler-Moser polynomial appeared in Lemma 2.5. Here \( p_1 \) is the real root and \( p_3 = \bar{p}_2. \) Let \( q_i = p_i^*. \) For \( z_1 \in \mathbb{R}, z_2 \in \mathbb{C}, \) we define the force map

\[
F_1(z_1, z_2) = \frac{1}{z_1 - z_2} + \frac{1}{z_1 - \bar{z}_2} - \frac{1}{2z_1} - \frac{1}{z_1 + z_2} - \frac{1}{z_1 - z_2}, \tag{2.11}
\]

\[
F_2(z_1, z_2) = \frac{1}{z_2 - z_1} + \frac{1}{z_2 - \bar{z}_2} - \frac{1}{z_2 + z_1} - \frac{1}{z_2 - z_2}. \tag{2.12}
\]

We have in mind that \( z_1 \) represents the vortex on the real axis and \( z_2 \) represents the one lying in the second quadrant. Note that by symmetry, \( F_1(z_1, z_2) \in \mathbb{R}. \) Writing \( z_1 = a_1, z_2 = a_2 + b_2i, \) where \( a_i, b_i \in \mathbb{R}, \) we can define

\[
F(a_1, a_2, b_2) := (F_1, \Re F_2, \Im F_2).
\]
The configuration \((p_1, p_2, p_3, q_1, q_2, q_3)\) is called nondegenerate, if
\[
\det DF (p_1, \Re p_2, \Im p_2) \neq 0.
\]

Numerical computation shows that \(\det DF (p_1, \Re p_2, \Im p_2) \neq 0\). Hence it is non-degenerate. It turns out for \(n\) large, this procedure is tedious and we have to find other ways to overcome this difficulty.

In the general case, let \(p = (p_1, \ldots, p_{n(n+1)/2})\), \(q = (q_1, \ldots, q_{n(n+1)/2})\). Define the map \(F :\)
\[
(p, q) \rightarrow (F_1, \ldots, F_{n(n+1)/2}, G_1, \ldots, G_{n(n+1)/2}),
\]
where
\[
F_k = \sum_{j \neq k} \frac{1}{p_k - p_j} - \sum_{j} \frac{1}{p_k - q_j},
\]
(2.13)
\[
G_k = \sum_{j \neq k} \frac{1}{q_k - q_j} - \sum_{j} \frac{1}{q_k - p_j}.
\]
(2.14)

Let \(a = (a_1, \ldots, a_{n(n+1)/2})\) be the roots of \(A_n\) and \(b = -(\bar{a}_1, \ldots, \bar{a}_{n(n+1)/2})\). Moreover, we assume that for \(i = 1, \ldots, i_0,\)
\[
a_{2i-1} = \bar{a}_{2i},
\]
while for \(i = 2i_0 + 1, \ldots, n(n+1)/2, \Im a_i = 0.\) Consider the linearization of \(F\) at \((p, q) = (a, b)\). Denote it by \(DF|_{(a,b)}\). This is a map from \(\mathbb{C}^{n(n+1)}\) to \(\mathbb{C}^{n(n+1)}\).

We remark that the points in \(a \cup b\) lie “approximately” on \(n\) circles (not exactly on these circles), and “approximately” on a certain number of straight lines.

The map \(DF|_{(a,b)}\) always has kernel. Indeed, for any parameter \(K = (k_2, \ldots, k_n), \Theta_n(z, K)\) and \(\tilde{\Theta}_n(z, K)\) satisfy
\[
\Theta_n'' \tilde{\Theta}_n - 2\Theta_n' \tilde{\Theta}_n' + \Theta_n \tilde{\Theta}_n'' = 2\mu \left(\Theta_n' \tilde{\Theta}_n - \Theta_n \tilde{\Theta}_n'\right).
\]

Differentiating this equation with respect to the translation and the parameters \(k_i, i = 2, \ldots, n,\) we get \(n\) (complex) dimensional kernel. Denote them by \(\varpi_1, \ldots, \varpi_n.\)

Let \(\xi = (\xi_1, \ldots, \xi_{n(n+1)/2}) \in \mathbb{C}^{n(n+1)/2}, \eta = (\eta_1, \ldots, \eta_{n(n+1)/2}) \in \mathbb{C}^{n(n+1)/2},\) the vector \((\xi, \eta),\) with \(\eta = \xi^*,\) is called symmetric if for \(i = 1, \ldots, i_0,\)
\[
\xi_{2i-1} = \bar{\xi}_{2i},
\]
while for \(i = 2i_0 + 1, \ldots, n(n+1)/2, \Im \xi_i = 0.\)

The main result of this section is the nondegeneracy of the vortex configuration given by \(A_n,\) i.e.e all the kernels are given by the above:

**Proposition 2.9.** Suppose \(DF|_{(a,b)} (\xi, \eta) = 0\) and \((\xi, \eta)\) is symmetric. Then \(\xi = \eta = 0.\)

Proposition 2.9 is proved by the linearization of Darboux transformation and recursive relations.
2.2. Darboux transformation and the nondegeneracy of the configuration. We first recall the following classical Darboux transformation theorem (Theorem 2.1, [32]).

**Theorem 2.10.** Suppose

\[-\Psi'' + u\Psi = \lambda \Psi,
-\Psi_1'' + u\Psi_1 = \lambda_1 \Psi_1.\]

Then the function \( \Phi := W(\Psi_1, \Psi) / \Psi_1 \) satisfies

\[-\Phi'' + \tilde{u}\Phi = \lambda \Phi,\]

where \( \tilde{u} := u - 2 (\ln \Psi_1)'' \).

The function \( \Phi \) is called the Darboux transformation of \( \Psi \).

Let \( \phi_n = \frac{A_{n+1}}{A_n} \). Consider the equation

\[
\phi'' + 2 (\ln A_n)'' \phi = 0. \tag{2.15}
\]

The second order recursive relation (2.10) is equivariant to the fact \( \phi_n \) is solution of (2.15). From [1], we know that \( \phi^{-1}_{n-1} \) is also a solution.

On the other hand, we have the following Darboux transformation relation

\[
(2n + 1) \phi^{-1}_n = \frac{W(\phi^{-1}_{n-1}, \phi_n)}{\phi_n}. \tag{2.16}
\]

Indeed, this is equivalent to the relation

\[
(2n + 1) A^2_n = A_{n-1} A'_{n+1} - A'_{n-1} A_{n+1}. \tag{2.17}
\]

Note that the constant \( 2n + 1 \) makes the coefficient of the leading order term of \( A_n \) to be 1. We also have the reversed transformation

\[
(2n + 3) \phi_n = \frac{W(\phi^{-1}_{n-2}, \phi_{n+1})}{\phi_{n+1}}.
\]

Indeed, this is equivalent to

\[
(2n + 3) A^2_{n+1} = A_{n-2} A'_{n+3} - A'_{n-2} A_{n+3}.
\]

We now recall that the function \( \psi_n(z) = \frac{A_n}{A_n e^{\mu z}} \) satisfying

\[
\psi'' + 2 (\ln A_n)'' \psi = \mu^2 \psi_n.
\]

Note that

\[
\psi_n = \frac{W(\theta_1, \ldots, \theta_{2n-1}, e^{\mu z})}{W(\theta_1, \ldots, \theta_{2n-1})}.
\]

Then the Darboux transformation [32] between \( \psi_n \) and \( \psi_{n+1} \) is given by

\[
\psi_{n+1} = \frac{W(\psi_n, \phi_n)}{\phi_n}. \tag{2.16}
\]

Explicitly,

\[
\psi_{n+1} = \psi_n (\ln \phi_n)' - \psi_n'. \tag{2.17}
\]

Let us verify this for the \( n = 1 \) case. We have

\[
\phi_1 = \frac{A_2}{A_1} = \frac{z^3 + k_2}{z},
\]

\[
\psi_1 = \frac{W(\theta_1, e^{\mu z})}{W(\theta_1)} = \frac{(\mu z - 1)e^{\mu z}}{z}.
\]
We also have

\[ \psi_2 = \frac{W(\theta_1, \theta_2, e^{\mu z})}{z^3 + k_2} \]

\[ = \frac{e^{\mu z}}{z^3 + k_2} \left( z - z^2 \mu + \frac{1}{3} \mu^2 k_2 + \frac{1}{3} z^3 \mu^2 \right). \]

Therefore,

\[ \psi_1 (\ln \phi_1)' - \psi_1' \]

\[ = \left( \frac{\mu z - 1}{z} \right) e^{\mu z} \frac{d}{dz} \left( \ln \frac{z^3 + k_2}{z} \right) - \frac{d}{dz} \left( \frac{(\mu z - 1) e^{\mu z}}{z} \right) \]

\[ = \psi_2. \]

Next we would like to analyze the linearized Darboux transformation. First of all, we linearize the equation (2.10) at \((A_n, A_{n+1})\). We obtain

\[ \xi''_{n+1} A_{n+1} - 2 \xi''_n A'_{n+1} + \xi_n A''_{n+1} + A''_n \xi_{n+1} - 2 A'_{n+1} \xi'_{n+1} + A_n \xi''_{n+1} = 0. \]

Let \(\xi_n = A_n \xi_n\). Then

\[ \begin{align*}
A''_n \xi_n + 2 A'_n \xi'_n + A_n \xi''_n & \quad A_{n+1} - 2 \left( A'_{n+1} \xi'_n + A_n \xi'_n + A_{n+1} \xi''_n \right) A_{n+1} + A_n \xi_{n+1} A''_{n+1} \\
A''_{n+1} \xi_{n+1} - 2 A'_{n+1} \left( A'_{n+1} \xi_{n+1} + A_{n+1} \xi''_{n+1} \right) + A_n \left( A''_{n+1} \xi_{n+1} + 2 A'_{n+1} \xi''_{n+1} + A_{n+1} \xi''_{n+1} \right) & = 0.
\end{align*} \]

Introducing \(f_n = \xi'_n\), we get

\[ \begin{align*}
A_n A_{n+1} f'_n + (2 A'_n A_{n+1} - 2 A_n A_{n+1}') f_n & \\
+ A_n A_{n+1} f'_{n+1} + (2 A_n A_{n+1}' - 2 A'_n A_{n+1}) f_{n+1} & = 0.
\end{align*} \]

This equation can be written as

\[ f'_n + 2 \left( \ln \frac{A_n}{A_{n+1}} \right)' f_n + f'_{n+1} + 2 \left( \ln \frac{A_{n+1}}{A_n} \right)' f_{n+1} = 0. \]

Hence for any given function \(f_{n+1}\), we can solve this equation and get

\[ f_n = \frac{A^2_{n+1}}{A^2_n} \int_0^z \frac{A^2_n}{A^2_{n+1}} \left( f'_n + 2 \left( \ln \frac{A_{n+1}}{A_n} \right)' f_{n+1} \right) ds \]

\[ = -f_{n+1} + 2 \frac{A^2_{n+1}}{A^2_n} \int_0^z \frac{A^2_n}{A^2_{n+1}} f'_{n+1} ds. \]

The last equality follows from integrating by parts for the second term.

Next, we linearize the equation (2.17) at \((A_n, A_{n+1})\) and obtain

\[ \sigma_{n+1} = \sigma_n (\ln \phi_n)' - \sigma'_n + \psi_n \left( \frac{\xi_{n+1}}{A_{n+1}} - \frac{\xi_n}{A_n} \right)' \]

We recall that \(\left( \frac{\xi_n}{A_n} \right)' = f_n\). Hence we get the equation

\[ \sigma'_n - \sigma_n (\ln \phi_n)' = \psi_n (f_{n+1} - f_n) - \sigma_{n+1}. \]
From this we get

\[ \sigma_n = \phi_n \int_0^z \phi_n^{-1} (\psi_n (f_{n+1} - f_n) - \sigma_{n+1}) \, ds. \]  \hspace{1cm} (2.19)

We are lead to the system

\[
\begin{aligned}
& f_n' + 2 \left( \ln \frac{A_n}{A_{n+1}} \right)' f_n + f_{n+1}' + 2 \left( \ln \frac{A_{n+1}}{A_n} \right)' f_{n+1} = 0, \\
& \sigma_n' - \sigma_n \left( \ln \phi_n \right) = \psi_n (f_{n+1} - f_n) - \sigma_{n+1}.
\end{aligned}
\]  \hspace{1cm} (2.20)

For given function \( f_{n+1} \) and \( \sigma_{n+1} \), we can solve this system and get a solution \( (f_n, \sigma_n) \) from \((2.18), (2.19)\).

Let \( K_0 = (-\frac{1}{2}, -\frac{1}{2}, \ldots) \). For each fixed \( n \), and \( j = 2, \ldots, n \), we define the polynomials

\[ \omega_{n,j} = \left( \frac{\partial_{k_j} \Theta_n (z + 1, K) |_{K=K_0}}{A_n} \right), \]

\[ \hat{\omega}_{n,j} = \left( \frac{\partial_{k_j} \hat{\Theta}_n (z + 1, K) |_{K=K_0}}{B_n} \right). \]

Let

\[ \beta_{n,j} = e^{\mu z} \left( -\frac{B_n \omega_{n,j}}{A_n^2} + \hat{\omega}_{n,j} \right). \]

**Lemma 2.11.** \( f_n = \omega_{n,j}, \sigma_n = \beta_{n,j} \) satisfy the system \((2.20)\). We also need the following uniqueness lemma on the symmetric configuration.

**Lemma 2.12.** Suppose \( \hat{K} \) is a \( n - 1 \) dimensional vector and \( \left| \hat{K} - K \right| + t + \frac{1}{2} < \delta \) for some small \( \delta \), with \( \hat{K} \not= K \). Then

\[ \Theta_n \left( -z - t, \hat{K} \right) \not= \hat{\Theta}_n \left( z - t, \hat{K} \right). \]

**Proof.** We prove this by induction. This is true for \( n = 1 \). Suppose this is true for \( n = k \), we prove that it is also true for \( n = k + 1 \). Indeed, suppose to be contrary that \( \Theta_n \left( -z - t, \hat{K} \right) \not= \hat{\Theta}_n \left( z - t, \hat{K} \right) \).

\[
\begin{align*}
\Theta_n'' \left( z - t, \hat{K} \right) \Theta_{n+1} \left( z - t, \hat{K} \right) - 2 \Theta_n' \left( z - t, \hat{K} \right) \Theta_{n+1}' \left( z - t, \hat{K} \right) \\
+ \Theta_n \left( z - t, \hat{K} \right) \Theta_{n+1}'' \left( z - t, \hat{K} \right) &= 0.
\end{align*}
\]

Replacing \( z \) be \( -z \), we get

\[
\begin{align*}
\hat{\Theta}_n'' \left( -z - t, \hat{K} \right) \hat{\Theta}_{n+1} \left( -z - t, \hat{K} \right) - 2 \hat{\Theta}_n' \left( -z - t, \hat{K} \right) \hat{\Theta}_{n+1}' \left( -z - t, \hat{K} \right) \\
+ \hat{\Theta}_n \left( -z - t, \hat{K} \right) \hat{\Theta}_{n+1}'' \left( -z - t, \hat{K} \right) &= 0.
\end{align*}
\]

This then implies that

\[ \Theta_{n+1} \left( -z - t, \hat{K} \right) = \hat{\Theta}_{n+1} \left( z - t, \hat{K} \right). \]

Hence \( t = -\frac{1}{2} \), and the first \( n - 2 \) components of \( \hat{K} \) is \(-\frac{1}{2}\). It then follows that the last component of \( \hat{K} \) is also \(-\frac{1}{2}\). This finishes the proof. \( \square \)
Since $\psi_n = \frac{B_n}{A_n} e^{\mu z}$, we have the relation
\[
\sigma_n e^{-\mu z} = -\frac{B_n \xi_n}{A_n^2} + \frac{\eta_n}{A_n}.
\]
Hence the function $\eta_n$ is given in terms of $\sigma_n, \xi_n$ by
\[
\eta_n = A_n \sigma_n e^{-\mu z} + \frac{B_n}{A_n} \xi_n.
\] (2.21)
Note that $(A_n, B_n)$ satisfies
\[
A_n'' B_n - 2 A_n' B_n' + A_n B_n'' - 2 \mu (A_n' B_n - A_n B_n') = 0.
\]
Linearizing this equation we get
\[
\xi_n'' B_n - 2 \xi_n' B_n' + \xi_n B_n'' - 2 \mu (\xi_n' B_n - \xi_n B_n') \\
+ A_n'' \eta_n - 2 A_n' \eta_n' + A_n \eta_n'' - 2 \mu (A_n' \eta_n - A_n \eta_n') = 0.
\]
In the case $n = 0$, we have $A_n = B_n = 1$, the above equation reads
\[
\xi_0'' - 2 \mu \xi_0' + \eta_0'' + 2 \mu \eta_0' = 0.
\] (2.22)
That is, $(\xi_0 + \eta_0)' = 2 \mu (\xi_0 - \eta_0)$. By (2.21),
\[
\eta_0 = \sigma_0 e^{-\mu z} + \xi_0.
\]
Hence
\[
\xi_0 + \eta_0 = \sigma_0 e^{-\mu z} + 2 \xi_0,
\]
\[
\xi_0 - \eta_0 = -\sigma_0 e^{-\mu z}.
\]
It follows that
\[
(\sigma_0 e^{-\mu z} + 2 \xi_0)' + 2 \mu (\sigma_0 e^{-\mu z}) = 0.
\] (2.23)
Now suppose the Adler-Moser polynomial $A_N$ satisfies assumption (A). Given functions $\xi_N$ and $\eta_N$, we have corresponding functions $f_N, \sigma_N$. Using (2.20), we can define recursively $(\xi_{N-1}, \eta_{N-1}, f_{N-1}, \sigma_{N-1}), \ldots, (\xi_0, \eta_0, f_0, \sigma_0)$. Linearizing the Darboux transformation, we find that $\xi_0, \eta_0$ satisfy (2.22).

**Proposition 2.13.** Suppose $\xi_N, \eta_N$ are polynomials with degree less than $N(N + 1)/2$. Then $\xi_0 = \eta_0 = 0$.

**Proof.** We first consider the case that for any $n \neq j \leq N$, $A_n$ and $A_j$ have no common roots. (This assumption is true for $N = 34$, as can be verified by Maple.) The idea for the general case is similar but notations are more involved.

Since $f_N = \left(\frac{\xi_N}{A_N}\right)$, $f_N$ is a rational function with possible poles at the roots of $A_N$. We know that for each $n \leq N - 1$, $f_n$ and $f_{n+1}$ are related by
\[
f_n = -f_{n+1} + 2 \frac{A_{n+1}^2}{A_n^2} \int_0^z \frac{A_{n+1}^2}{A_n^2} f_{n+1}' ds.
\] (2.24)
Hence $f_n$ has possible poles at the roots of $A_n, A_{n+1}, \ldots, A_N$. We remark that as a complex valued function with poles, $f_n$ may be multiple-valued. By (2.24),

\[
    f_n - f_{n+1} = -2f_{n+1} + 2\frac{A_{n+1}^2}{A_n^2} \int_0^z \frac{A_{n+1}^2}{A_n^2} f'_{n+1} ds
    = -2\frac{A_{n+1}^2}{A_n^2} \int_0^z f_{n+1} \left( \frac{A_{n+1}^2}{A_{n+1}^2} \right)' ds.
\]  

(2.25)

In particular,

\[
    f_0 = f_1 + 4(z - 1)^2 \int_0^z \frac{f_1}{(s - 1)^2} ds.
\]  

(2.26)

On the other hand,\n
\[
    \sigma_n = \phi_n \int_0^z \phi_n^{-1} (\psi_n (f_{n+1} - f_n) - \sigma_{n+1}) ds.
\]

Recall that $\phi_0 = z - 1, \psi_0 = 1$. Hence

\[
    \sigma_0 = (z - 1) \int_0^z \frac{1}{z - 1} e^{\mu z} (f_1 - f_0 - \sigma_1) ds
    = \frac{A_{n+1}}{A_n} \int_0^z \frac{B_n}{A_{n+1}} e^{\mu s} (f_{n+1} - f_n) ds - \frac{A_{n+1}}{A_n} \int_0^z \frac{A_n}{A_{n+1}} \sigma_{n+1} ds.
\]

By (2.23),

\[
    (\sigma_0 e^{\mu s})' + 2e^{2\mu s} f_0 = 0.
\]

Hence using the fact that $\mu = 1$, we obtain

\[
    (\sigma_0 e^{\mu s})' + 2e^{2\mu s} f_0
    = e^{2\mu s} (f_0 + f_1 - \sigma_1)
    + ze^{\mu s} \int_0^z \frac{1}{s - 1} e^{\mu s} (f_1 - f_0 - \sigma_1) ds
    = 0.
\]  

(2.27)

Our next aim is to show that $f_1$ has no singularity except the point $z = 1$. Let us consider the term $\frac{A_{n+1}}{A_n} \int_0^z A_n \frac{A_{n+1}}{A_{n+1}} \sigma_{n+1} ds$. Let $z = d_0$ be a singularity of $f_n$ which is not the root of $A_n$. Then loosely speaking, the degree of singularity $\sigma_n$ is smaller than that of $\sigma_{n+1}$ and $f_n$. By (2.25), $f_n$ and $f_{n+1}$ has essentially the same degree of singularity at $d_0$. But this contradicts with the identity (2.27). Hence $f_1$ can only have singularity at $z = 1$.

Now we show that $f_0 = 0$. To see this, we observe that since $f_1$ has no other singularities, by the recursive relation, we deduce that $f_1$ is actually single valued and $f_1 = c_1 \frac{1}{z - 1} + c_2 \frac{1}{(z - 1)^2}$, and

\[
    \sigma_1 = \phi_1 \int_0^z \phi_1^{-1} (\psi_1 (f_2 - f_1) - \sigma_2) ds.
\]

Putting this into (2.27), we find that $c_1 = 0$. Hence $f_0 = 0$ and $\sigma_0 = 0$. \qed

Now we can prove Proposition 2.9. By Proposition 2.13, the kernel of the map $\partial \mid_{(a, b)}$ is given by linear combination of $\varphi_i$. For $\mu = 1, k_i = -\frac{1}{2}$ are the only parameters for which $\Theta_n$ and $\tilde{\Theta}_n$ give rise to symmetric configuration. Hence the configuration is nondegenerate.
3. Preliminaries on the Ginzburg-Landau equation

In this section, we recall some results on the Ginzburg-Landau equation. Most of the materials in this section can be found in [35] (possibly with different notations though).

Stationary solutions of the GP equation (1.1) solve the following Ginzburg-Landau equation

$$-\Delta \Phi = \Phi \left(1 - |\Phi|^2\right)$$

in \(\mathbb{R}^2\), \((3.1)\)

where \(\Phi\) is a complex valued function. As we mentioned before, it has degree \(\pm d\) vortices of the form \(S_d(r)e^{\pm id\theta}\). The asymptotic behavior of \(S_d\) can be described.

It is known that as \(r \to +\infty\),

$$S_d(r) = 1 - \frac{d^2}{2r^2} + O\left(r^{-4}\right). \quad (3.2)$$

On the other hand, as \(r \to 0\), there is a constant \(\kappa = \kappa_d > 0\) such that

$$S_d(r) = \kappa r \left(1 - \frac{r^2}{8} + O\left(r^4\right)\right). \quad (3.3)$$

See [19] for the proof of these facts.

Let \(\varepsilon > 0\) be small. For technical reasons, we need to modify \(S\) in the region \(r > C_0\varepsilon^{-1}\), where \(C_0\) is a fixed large constant, such that \(S(r) = 1\) for \(r > C_0\varepsilon^{-1} + 1\). We still denote it by \(S\) for notational simplicity.

The linearized operator of the Ginzburg-Landau equation around \(v_+\) will be denoted by \(L\):

$$\eta \to \Delta \eta + \left(1 - |v_+|^2\right)\eta - 2v_+ \text{Re}(\eta \bar{v}_+)$$

It turns out to be more convenient to study the operator

$$\mathcal{L}\eta := e^{-i\theta} L(e^{i\theta} \eta)$$

If we write the complex function \(\eta\) as \(w_1 + iw_2\) with \(w_1, w_2\) being real valued functions, then explicitly

$$\mathcal{L}\eta = e^{-i\theta} \Delta (e^{i\theta} \eta) + (1 - S^2) \eta - 2S^2w_1 = \Delta w_1 + (1 - 3S^2) w_1 - \frac{1}{r^2} w_1 - \frac{2}{r^2} \partial_\theta w_2$$

$$+ i \left(\Delta w_2 + (1 - S^2) w_2 - \frac{1}{r^2} w_2 + \frac{2}{r^2} \partial_\theta w_1\right).$$

Invariance of the equation (3.1) under rotation and translation gives us three linearly independent kernels of the operator \(\mathcal{L}\), called Jacobi fields. Rotational invariance yields the solution

$$\Phi^0 := ie^{-i\theta} v_+ = iS,$$  \((3.4)\)

while the translational invariance along \(x\) and \(y\) direction leads to the solutions

$$\Phi^+ := S' \cos \theta - \frac{S}{r} \sin \theta,$$

$$\Phi^- := S' \sin \theta + \frac{S}{r} \cos \theta.$$
bounded solutions of $\mathcal{L} \eta = 0$ are given by linear combinations of $\Phi^0$ and $\Phi^+, \Phi^-$. (See [Theorem 3.2, [35]]. Another proof can be found in [18].) Similar results hold for the degree $-1$ vortex $v$. It is worth mentioning that the nondegeneracy of those higher degree vortices $e^{in\theta} S_d(r)$, $|d| > 1$, is still an open problem. Actually this is the main reason why we only deal with the degree $\pm 1$ vortices in this paper.

The analysis of the asymptotic behavior of the kernels of $\mathcal{L}$ near 0 and $\infty$ is crucial in understanding the mapping property of the linearized operator $\mathcal{L}$. In doing this, the main strategy is to decompose the kernel into different Fourier modes. Let us now briefly describe the results in the sequel. Lemma 3.1, Lemma 3.2 and Lemma 3.3 below can be found in Section 3.3 of [35].

We start the discussion with the lowest Fourier mode, which is the simplest case.

**Lemma 3.1.** Suppose $a$ is a complex valued solution of the equation $\mathcal{L} a = 0$, depending only on $r$.

(I) As $r \to 0$, either $|a|$ blows up at least like $r^{-1}$, or $a$ can be written as a linear combination of two linearly independent solutions $w_{0,1}, w_{0,2}$, with

\[
\begin{align*}
w_{0,1}(r) &= r \left( 1 + O(r^2) \right), \\
w_{0,2}(r) &= ir \left( 1 + O(r^2) \right).
\end{align*}
\]

(II) As $r \to +\infty$, if $a$ is an imaginary valued function, then $a = c_1 + c_2 \ln r + O(r^{-2})$; if $a$ is real valued, then it either blows up or decays exponentially.

**Proof.** We sketch the proof for completeness.

If $L a = 0$ and the complex function $a$ depends only on $r$, then $a$ will satisfy

\[
a'' + \frac{1}{r} a' - \frac{1}{r^2} a = S^2 a - \left( 1 - 2S^2 \right) a. \tag{3.5}
\]

Note that this equation is not complex linear and its solution space is a 4-dimensional real vector space. The Jacobi field $\Phi^0$ defined by (3.4) is a purely imaginary solution of (3.5). Writing $a = a_1 + a_2i$, where $a_i$ are real valued functions, we get from (3.5) two decoupled equations:

\[
\begin{align*}
a_1'' + \frac{1}{r} a_1' - \frac{1}{r^2} a_1 + \left( 1 - 3S^2 \right) a_1 &= 0, \\
a_2'' + \frac{1}{r} a_2' - \frac{1}{r^2} a_2 + \left( 1 - S^2 \right) a_2 &= 0. \tag{3.6}
\end{align*}
\]

Observe that due to (3.2), as $r \to +\infty$,

\[
\begin{align*}
1 - 3S^2 - r^{-2} &= 0 + O(r^{-2}), \\
1 - S^2 - r^{-2} &= O(r^{-4}).
\end{align*}
\]

While due to (3.3), as $r \to 0$, $1 - S^2 = 1 + O(r^2)$. The results of this lemma then follow from a perturbation argument. \hfill $\square$

For each integer $n \geq 1$, we consider kernels of $\mathcal{L}$ the form $a(r) e^{in\theta} + b(r) e^{-in\theta}$. The complex valued functions $a, b$ will satisfy the following coupled ODE system in $(0, +\infty)$:

\[
\begin{align*}
a'' + \frac{1}{r} a' - \frac{(n+1)^2}{r^2} a &= S^2 b - \left( 1 - 2S^2 \right) a, \\
b'' + \frac{1}{r} b' - \frac{(n-1)^2}{r^2} b &= S^2 a - \left( 1 - 2S^2 \right) b. \tag{3.7}
\end{align*}
\]
By analyzing this coupled ODE system, one gets the precise asymptotic behavior of its solutions. The next lemma deals with the $n = 1$ case.

**Lemma 3.2.** Suppose $w = a(r) e^{i\theta} + b(r) e^{-i\theta}$ solves $Lw = 0$.

(I) As $r \to 0$, either $|w|$ blows up at least like $-\ln r$, or $w$ can be written as a linear combination of 4 linearly independent solutions $w_{1,i}, i = 1, \ldots, 4$, satisfying: As $r \to 0$,

$$
w_{1,1} = r^{2} \left( 1 + O \left( r^{2} \right) \right) e^{i\theta} + O \left( r^{6} \right) e^{-i\theta},$$
$$w_{1,2} = i r^{2} \left( 1 + O \left( r^{2} \right) \right) e^{i\theta} + O \left( r^{6} \right) e^{-i\theta},$$
$$w_{1,3} = (1 + O \left( r^{2} \right) ) e^{-i\theta} + O \left( r^{4} \right) e^{i\theta},$$
$$w_{1,4} = i (1 + O \left( r^{2} \right) ) e^{-i\theta} + O \left( r^{4} \right) e^{i\theta}.$$

(II) As $r \to +\infty$, either $|w|$ is unbounded (blows up exponentially or like $r$), or $|w|$ decays to zero (exponentially or like $r^{-1}$).

For the $n \geq 2$ case, we have the following

**Lemma 3.3.** Suppose $w = a(r) e^{i\theta} + b(r) e^{-i\theta}$ solves $Lw = 0$.

(I) As $r \to 0$, either $|w|$ blows up at least like $r^{1-n}$, or $w$ can be written as a linear combination of 4 linearly independent solutions $w_{1,i}, i = 1, \ldots, 4$, satisfying: As $r \to 0$,

$$w_{n,1} = r^{n+1} \left( 1 + O \left( r^{2} \right) \right) e^{i\theta} + O \left( r^{n+5} \right) e^{-i\theta},$$
$$w_{n,2} = i r^{n+1} \left( 1 + O \left( r^{2} \right) \right) e^{i\theta} + O \left( r^{n+5} \right) e^{-i\theta},$$
$$w_{n,3} = r^{n-1} \left( 1 + O \left( r^{2} \right) \right) e^{-i\theta} + O \left( r^{n+3} \right) e^{i\theta},$$
$$w_{n,4} = i r^{n-1} \left( 1 + O \left( r^{2} \right) \right) e^{-i\theta} + O \left( r^{n+3} \right) e^{i\theta}.$$

(II) As $r \to +\infty$, either $|w|$ is unbounded (blows up exponentially or like $r^{n}$), or $|w|$ decays to zero (exponentially or like $r^{-n}$).

By Lemma 3.3, for $n \geq 3$, if $Lw = 0$ and $w$ is bounded near 0, then decays at least like $r^{2}$ as $r \to 0$, hence decaying faster than the vortex solution itself. For $n \leq 2$, solutions of $Lw = 0$ bounded near 0 behaves like $O \left( r \right)$ or $O \left( 1 \right)$. Note that $\Phi_{0}, \Phi_{1}, \Phi_{-1}$ have this property. Now define $\Psi_{0} = \kappa w_{0,2}$,

$$\Psi_{+1} = \kappa w_{1,3} + \frac{\kappa}{8} w_{1,1}, \Psi_{-1} = \kappa w_{1,4} - \frac{\kappa}{8} w_{1,2},$$
$$\Psi_{+2} = w_{2,3}, \Psi_{-2} = w_{2,4}.$$

They behave like $O \left( r \right)$ or $O \left( 1 \right)$ near 0, but blows up as $r \to +\infty$.

**4. Construction of multi-vortex solutions**

4.1. **Approximate solutions and estimate of the error.** We would like to construct traveling wave solutions by gluing together $N \left( N + 1 \right)/2$ pairs of degree ±1 vortices. Let us simply choose $N = 2$, the general case is similar, but notations will be much more involved.
For $k = 1, 2, 3$, let $p_k, q_k \in \mathbb{C}$. We have in mind that $p_k$ are close to roots of the Adler-Moser polynomial $A_2$. We define the translated vortices

$$u_k = v_+ \left( z - \varepsilon^{-1} p_k \right), u_{3+k} = v_- \left( z - \varepsilon^{-1} q_k \right).$$

We then define the approximate solution

$$u := \prod_{k=1}^{6} u_k.$$ 

Note that as $r \to +\infty$, $u \to 1$. Hence the degree of $u$ is 0. Let us denote the function $z \to u(z)$ by $\bar{u}$. The next lemma states that the real part of $u$ is even both in the $x$ and $y$ variables, while the imaginary part is even in $x$ and odd in $y$.

**Lemma 4.1.** The approximate solution $u$ has the following symmetry:

$$u(\bar{z}) = \bar{u}(z), \quad u(z^*) = u(z).$$

**Proof.** Observe that the standard vortex $v_+ = S(r) e^{i\theta}$ satisfies

$$v_+ (\bar{z}) = \bar{v}_+ (z), \quad v_+ (z^*) = (v_+ (z))^*.$$ 

The oppositive vortex $v_-$ has similar properties. Hence using the fact that the set $\{p_1, p_2, p_3\}$ is invariant with respect to the reflection across the $x$ axis, we get

$$u(\bar{z}) = \prod_{k=1}^{3} \left( v_+ (\bar{z} - \varepsilon^{-1} p_k) v_- (\bar{z} - \varepsilon^{-1} q_k) \right) = \prod_{k=1}^{3} \left( \bar{v}_+ (z - \varepsilon^{-1} \bar{p}_k) \bar{v}_- (z - \varepsilon^{-1} \bar{q}_k) \right) = \bar{u}(z).$$

Moreover, since $v_- = \bar{v}_+$, we have

$$u(z^*) = \prod_{k=1}^{3} \left( v_+ (z^* - \varepsilon^{-1} p_k) v_- (z^* - \varepsilon^{-1} q_k) \right) = \prod_{k=1}^{3} \left( (v_+ (z - \varepsilon^{-1} q_k))^* (v_- (z - \varepsilon^{-1} p_k))^* \right) = \prod_{k=1}^{3} \left( \bar{v}_+ (z - \varepsilon^{-1} q_k) (\bar{v}_- (z - \varepsilon^{-1} p_k)) \right) = u(z).$$

This finishes the proof. $\square$

We use $E(u)$ to denote the error of the approximate solution:

$$E(u) := \varepsilon i \partial_y u + \Delta u + u \left( 1 - |u|^2 \right).$$

We have

$$\Delta u = \Delta (u_1 \ldots u_6)$$

$$= \sum_k \left( \Delta u_k \prod_{j \neq k} u_j \right) + \sum_{k \neq j} \left( \nabla u_k \star \nabla u_j \prod_{l \neq k,j} u_l \right).$$
where symbol \( * \) denotes \( \nabla u_k \ast \nabla u_j := \partial_x u_k \partial_x u_j + \partial_y u_k \partial_y u_j \). On the other hand, writing \( |u_k|^2 - 1 = \rho_k \), we obtain

\[
|u|^2 - 1 = \prod_k (1 + \rho_k) - 1 = \sum_k \rho_k + \sum_{k=2}^6 Q_k,
\]

where \( Q_k = \sum_{i_1 < i_2 < \cdots < i_k} (\rho_{i_1} \cdots \rho_{i_k}) \). Using the fact that \( u_k \) solves the Ginzburg-Landau equation, we get

\[
E(u) = \varepsilon \sum_k \left( \partial_y u_k \prod_{j \neq k} u_j \right) + \sum_{k,j,k \neq j} \left( (\nabla u_k \ast \nabla u_j) \prod_{l \neq i,j} u_l \right) - u \sum_{k=2}^6 Q_k.
\]

We have in mind that the main order terms are \( \partial_y u_k \prod_{j \neq k} u_j \) and \( (\nabla u_k \ast \nabla u_j) \prod_{l \neq i,j} u_l \).

Throughout the paper \((r_j, \theta_j)\) will denote the polar coordinate with respect to the point \(p_j\). Note that

\[
\partial_x (e^{i\theta}) = \frac{-yi e^{i\theta}}{r^2}, \quad \partial_y (e^{i\theta}) = \frac{xie^{i\theta}}{r^2}.
\]

We compute

\[
\partial_x u_k = -\frac{i y_k e^{i\theta_k}}{r_k^2} S(r_k) + \frac{x_k}{r_k} S'(r_k) e^{i\theta_k},
\]

\[
\partial_y u_k = \frac{i x_k e^{i\theta_k}}{r_k} S(r_k) + \frac{y_k}{r_k} S'(r_k) e^{i\theta_k}.
\]

**Lemma 4.2.** *In the region* \(|z| > C\varepsilon^{-1}\),

\[
|E(u)| \leq C |z|^{-2},
\]

\[
\text{Im}(E(u)) \leq C |z|^{-3}.
\]

**Proof.** We first estimate, for \( z > C\varepsilon^{-1}\),

\[
|\theta_1 - \theta_4| \leq C r^{-1},
\]

\[
|\partial_x (\theta_1 - \theta_4)| \leq C r^{-2}.
\]

Hence \(|\partial_x u| \leq C r^{-2}\). Next,

\[
|\nabla u_k \ast \nabla u_j| = |\partial_x u_k \partial_x u_j + \partial_y u_k \partial_y u_j| \\
\leq |\partial_x u_k| |\partial_x u_j| + |\partial_y u_k| |\partial_y u_j| \\
\leq C |z|^{-2}.
\]

Finally, since \( \rho_k \leq C |z|^{-2} \), we have \( Q_k \leq C |z|^{-4} \). This finishes the proof. \( \square \)
4.2. **Projection of the error on the kernel.** Now we study the projection of the error of the approximate solution on the kernels. We have, in the region where 
\[|z - \varepsilon^{-1}p_k| \leq \frac{1}{2}\varepsilon^{-1},\]

\[
\nabla u_k \ast \nabla u_j = \partial_x u_k \partial_x u_j + \partial_y u_k \partial_y u_j \\
\sim \partial_x u_k \left(-\frac{y_j i\epsilon^{\iota_0}}{r^2_j}\right) + \partial_y u_k \left(-\frac{x_j i\epsilon^{\iota_0}}{r^2_j}\right).
\]

We have
\[
\text{Im} \left(\partial_y u_k (\partial_x u_k)\right) = \frac{SS'}{r_k}.
\]

It follows that
\[
\text{Re} \int_{|z - \varepsilon^{-1}p_k| \leq r_j} e^{-\iota \theta_0} (\nabla u_k \ast \nabla u_j) \overline{(\partial_x u_k)} \, dx \, dy \\
= -\text{Re} \left(\frac{\varepsilon}{p_k - p_j}\right) \text{Im} \int_{|z - \varepsilon^{-1}p_k| \leq r_j} \partial_y u_k (\partial_x u_k) \\
= -\text{Re} \left(\frac{\varepsilon}{p_k - p_j}\right) \int_{|z - \varepsilon^{-1}p_k| \leq r_j} \frac{SS'}{r_k} \\
= -\pi \text{Re} \left(\frac{\varepsilon}{p_k - p_j}\right) + O \left(\frac{r^{-2}}{r_k}\right).
\]

Similarly,
\[
\text{Re} \int_{|z - \varepsilon^{-1}p_k| \leq r_j} e^{-\iota \theta_0} (\nabla u_k \ast \nabla u_j) \overline{(\partial_y u_k)} \\
= \pi \text{Im} \left(\frac{\varepsilon}{p_k - p_j}\right) + O \left(\frac{r^{-2}}{r_k}\right).
\]

Next, if \(l, j \neq k\), we estimate that for \(|z - \varepsilon^{-1}p_k| \leq C\varepsilon^{-1},\)
\[
(\nabla u_l \ast \nabla u_j) (\overline{(\partial_x u_k)}) \sim e^{-\iota \theta_0} \left(\frac{y_j i\epsilon^{\iota_0}}{r^2_j} + \frac{x_j i\epsilon^{\iota_0}}{r^2_j}\right) \left(\frac{y_j i\epsilon^{\iota_0}}{r^2_j} + \frac{x_j i\epsilon^{\iota_0}}{r^2_j}\right) \\
= O \left(\frac{1}{rk}\right).
\]

Finally,
\[
\text{Re} \int_{\mathbb{R}^2} i\varepsilon \partial_y u_k (\overline{(\partial_y u_k)}) = O \left(\varepsilon^2\right), \\
\text{Re} \int_{\mathbb{R}^2} i\varepsilon \partial_y u_k (\overline{(\partial_x u_k)}) = \pi\varepsilon + O \left(\varepsilon^3\right).
\]

Combining these estimates, we find that the projected equation at the main order is (2.1) with \(\mu = 1\).

4.3. **The nonlinear scheme.** We search a traveling wave solution \(U\) of GP equation:
\[-i\varepsilon \partial_y U = \Delta U + U \left(1 - |U|^2\right).\]

After a rescaling, we get the equation
\[-i\partial_y U = \Delta U + \varepsilon^{-2} U \left(1 - |U|^2\right).\]
Then \( \phi \) should satisfy

\[
\varepsilon i \partial_y (u + \phi) = \Delta (u + \phi) + (u + \phi) \left( 1 - |u + \phi|^2 \right).
\]

We write this equation as

\[
- \Delta \phi + \left( 1 - |u|^2 \right) \phi - 2u \text{Re} (\bar{u}\phi) = E(u) + \varepsilon i \partial_x \phi + N(\phi).
\]

(4.1)

Here \( N(\phi) \) is a higher order perturbation term and equals \( 2u|\phi|^2 + \bar{u}\phi^2 + |\phi|^2 \).

Denote the left hand side by \( G\phi \).

Instead of analyze the operator \( G \) directly, we will study its conjugate operator, possibly with different forms in different regions of \( \mathbb{R}^2 \).

This technique has already appeared in Section 2. We write \( u \) as \( \tilde{\rho} e^{i\tilde{\theta}} \), where \( \tilde{\rho}, \tilde{\theta} \) are real valued. Observe that

\[
L \eta := e^{-i\tilde{\theta}} G \left( e^{i\tilde{\theta}} \eta \right) = -e^{-i\tilde{\theta}} \Delta \left( e^{i\tilde{\theta}} \eta \right) + \left( 1 - \tilde{\rho}^2 \right) \eta - 2\tilde{\rho} \text{Re} (\tilde{\rho} \eta).
\]

If we write \( \eta = \eta_1 + i\eta_2 \), with \( \eta_k \) being real function, then

\[
L \eta = -\Delta \eta - 2e^{-i\tilde{\theta}} \nabla \left( e^{i\tilde{\theta}} \right) \nabla \eta - e^{-i\tilde{\theta}} \Delta \left( e^{i\tilde{\theta}} \right) \eta
\]

\[
+ \left( 1 - \tilde{\rho}^2 \right) (\eta_1 + i\eta_2) - 2\tilde{\rho} \text{Re} (\tilde{\rho} \eta_1).
\]

This tells us the in the region \(|z| > C\varepsilon^{-1}\), the real part of \( L\eta \) is a well behaved operator like \( -\Delta \eta_1 + 2\eta_1 \). As for the imaginary part, we recall that by our definition of \( S \),

\[
1 - \tilde{\rho}^2 = 0, \text{ for } |z| \geq C_1\varepsilon^{-1}.
\]

Hence the imaginary part behaves like the Laplacian operator \(-\Delta \eta_2 \). Writing \( \phi = e^{i\tilde{\theta}} \eta \), equation (4.1) becomes

\[
L \eta = (E(u) + \varepsilon i \partial_x \phi + N(\phi)) e^{-i\tilde{\theta}}.
\]

Let us denote the right hand side by \( M \).

**Lemma 4.3.** \( M \) has the following symmetry

\[
M(\bar{z}) = \overline{M(z)}, M(z^*) = M(z).
\]

**Proof.** This follows from the symmetry of the approximate solution \( u \). \( \square \)

We also recall the following result from Lemma 4.2 in [30]:

**Lemma 4.4.** Let \( h \) satisfy

\[
\Delta h + f(z) = 0, h(\bar{z}) = -h(z), |h| \leq C,
\]

where \( f \) satisfies

\[
|f(z)| \leq \frac{C}{(1 + |z|)^{2+\sigma}}, \sigma \in (0,1).
\]

Then

\[
|h(z)| \leq \frac{C}{(1 + |z|)^{\eta}}.
\]
Now we introduce the functional framework adapted to the mapping property of the linearized operator $\mathbb{L}_\gamma$. Following [30], we fix $\gamma \in (0,1)$ and $\sigma \in \left(0, \frac{1}{2}\right)$. For $\eta = \eta_1 + \eta_2 i$, we define
\[
\|\eta\|_* := \sum_j \|u \eta\|_{C^{2,\gamma}(r_j < 2)} + \sum_j \|u \eta\|_{C^{2,\gamma}(r_j < 3)} + \sum_j \|r_j^{1+\sigma} \eta_1\|_{L^\infty(r_j > 2)} + \sum_j \|r_j^{2+\sigma} \partial_x \eta_1\|_{L^\infty(r_j > 2)} + \sum_j \|r_j^{1+\sigma} \eta_2\|_{L^\infty(r_j > 2)} + \sum_j \|r_j^{2+\sigma} \partial_x \eta_2\|_{L^\infty(r_j > 2)}.
\]

For $h = h_1 + i h_2$, we define
\[
\|h\|_{**} := \sum_j \|h\|_{C^{0,\gamma}(r_j < 2)} + \sum_j \|r_j^{1+\sigma} h_1\|_{L^\infty(r_j < 2)} + \sum_j \|r_j^{2+\sigma} h_2\|_{L^\infty(r_j < 2)}.
\]

We have the following projected linear theory.

**Proposition 4.5.** Let $\varepsilon$ be small. Suppose
\[
\begin{align*}
\mathbb{L}_\gamma \eta &= h, \\
\Re \left( \int_{|z| \leq \epsilon^{-1} p_k} \bar{\eta} e^{-i \theta} \partial_x u \right) &= 0, \quad k = 1, \ldots, 6, \\
\Re \left( \int_{|z| \leq \epsilon^{-1} p_k} \bar{\eta} e^{-i \theta} \partial_y u \right) &= 0, \quad k = 1, \ldots, 6 \\
\eta e^{i \theta} &\text{ has the symmetry as $M$.}
\end{align*}
\]

Then $\|\eta\|_* \leq C |\ln \varepsilon| \|h\|_{**}$.

**Proof.** This can be proved by using the linear theory of the standard $\pm 1$ vortex described in Section 3 and Lemma 4.4. It can also be proved along the same ideas as that of Lemma 5.1 in [30], using blow up and contradiction arguments, which is in the spirit similar as that of [17]. Since this type of results are by now more or less standard, we omit the details. \(\square\)

Now we are ready to prove our main theorem in this paper. Since technically the method is quite similar to that of [30], we only sketch the main steps.

Setting $\phi = e^{i \theta} \eta$, we write the nonlinear problem as
\[
\mathbb{L}_\gamma \eta = e^{-i \theta} \left( E(u) + \varepsilon i \partial_x \phi + N(\phi) \right).
\]

The error $E(u)$ can be estimated by $\|E(u)\|_{**} \leq C \varepsilon^{1-\sigma}$. By Proposition 4.5, the equation (4.2) can be solved modulo projection on the kernel $\partial_x u, \partial_y u$, using contradiction argument. More precisely, let $\eta_k$ be cutoff functions supported in the region $|z - \varepsilon^{-1} p| \leq c \varepsilon^{-1}$, for a fixed small constant $c$ less than the distances between any two roots of the Adler-Moser polynomials $A_n, B_n$. We can find $c_k, d_k, \eta$ such that
\[
\mathbb{L}_\gamma \eta = e^{-i \theta} \left( E(u) + c \varepsilon^{-1} \partial_x \phi + N(\phi) \right) + \sum_k \left( c_k e^{-i \theta} \partial_x u + d_k e^{-i \theta} \partial_y u \right) \eta_k.
\]

Moreover, $\|\eta\|_{**} \leq C |\ln \varepsilon| \varepsilon^{1-\sigma}$. Projecting both sides on $\partial_x u, \partial_y u$ and using the estimate of $\eta$, we find that $c_k, d_k$ equal zero, is equivariant to that $p_k, q_k$ satisfy the
system
\[
\begin{align*}
\sum_{j \neq \alpha} \frac{1}{p_{\alpha} - p_j} - \sum_{j \neq \alpha} \frac{1}{p_{\alpha} - q_j} &= 1 + O(\varepsilon^\alpha), \quad \text{for } \alpha = 1, \ldots, N, \\
\sum_{j \neq \alpha} \frac{1}{q_{\alpha} - q_j} - \sum_{j \neq \alpha} \frac{1}{q_{\alpha} - p_j} &= O(\varepsilon^\alpha), \quad \text{for } \alpha = 1, \ldots, N,
\end{align*}
\]
for some $\alpha > 0$. Now using the nondegeneracy of the roots of the Adler-Moser polynomial and the Lipschitz dependence of the $O(\varepsilon^\alpha)$ term on $p_k$, we can solve this system using contraction mapping principle and get a solution $p_k, q_k$, close to roots of the Adler-Moser polynomial $a, b$.

References
[1] Adler, M.; Moser, J. On a class of polynomials connected with the Korteweg-de Vries equation. *Comm. Math. Phys.* 61 (1978), no. 1, 1–30.
[2] Aref, Hassan. Vortices and polynomials. *Fluid Dyn. Res.* 39 (2007), no. 1-3, 5–23.
[3] Aref, Hassan. Point vortex dynamics: a classical mathematics playground. *J. Math. Phys.* 48 (2007), no. 6, 065401, 23 pp.
[4] Hassan Aref, Paul K. Newton, Mark A. Sremler, Tadashi Tokieda, Dmitri L. Vainchtein, Vortex crystals. *Adv. Appl. Mech.* 39, 1–79., 2002.
[5] Bartman, A.B., A new interpretation of the Adler–Moser KdV polynomials: interaction of vortices. In: Sagdeev, R.Z. (Ed.), Nonlinear and Turbulent Processes in Physics, vol. 3.,1983, Harwood Academic Publishers, NewYork, pp. 1175–1181.
[6] F. Bethuel; H. Brezis; F. Helein, Ginzburg-Landau vortices. Progress in Nonlinear Differential Equations and their Applications, 13. Birkhauser Boston, Inc., Boston, MA, 1994.
[7] F. Bethuel; P. Gravejat; J.-C. Saut, On the KP I Transonic Limit of Two-Dimensional ross-Pitaevskii Travelling Waves, *Dynamics of PDE*, Vol.5, No.3, 241-280, 2008.
[8] F. Bethuel; P. Gravejat; J.-C. Saut. Existence and properties of travelling waves for the Gross-Petaevskii equation, *Contemporary Mathematics*, Vol. 473, 2008, 55–103.
[9] Bethuel, F.; Gravejat, P.; Saut, J.-C., Travelling waves for the Gross-Pitaevskii equation. II. *Comm. Math. Phys.* 285 (2009), no. 2, 567–651.
[10] Bethuel, F.; Orlandi, G.; Smets, D. Vortex rings for the Gross-Pitaevskii equation. *J. Eur. Math. Soc. (JEMS)* 6 (2004), no. 1, 17–94.
[11] Bethuel, F.; Saut, J.-C., Travelling waves for the Gross-Pitaevskii equation. I. *Ann. Inst. H. Poincare Phys. Theor.* 70 (1999), no. 2, 147–238.
[12] Burchann, J. L.; Chaundy, T. W. A Set of Differential Equations which can be Solved by Polynomials. *Proc. London Math. Soc.* (2) 30 (1930), no. 6, 401–414.
[13] D. Chiron, C. Scheid, Multipe branches of travelling waves for the Gross Pitaevskii equation, Preprint.
[14] D. Chiron. Travelling waves for the Gross-Pitaevskii equation in dimension larger than two. *Nonlinear Anal., Theory, Methods, Appl.*, 58(1-2):175–204, 2004.
[15] D. Chiron and M. Maris. Traveling Waves for Nonlinear Schrodinger Equations with Nonzero Conditions at Infinity. *Arch. Rational Mech. Anal.* 226(1):143–242, 2017.
[16] Clarkson, P. A. Vortices and polynomials. *Stud. Appl. Math.* 123 (2009), no. 1, 37–62.
[17] del Pino, Manuel; Kowalczyk, Michal; Musso, Monica. Variational reduction for Ginzburg-Landau vortices. *J. Funct. Anal.* 239 (2006), no. 2, 497–541.
[18] del Pino, Manuel; Felmer, P.; Kowalczyk, Michal, Minimality and nondegeneracy of degree-one Ginzburg-Landau vortex as a Hardy’s type inequality. *Int. Math. Res. Not.* (2004), no. 30, 1511–1527.
[19] Fife, P. C.; Peletier, L. A. On the location of defects in stationary solutions of the Ginzburg-Landau equation in $R^2$, *Quart. Appl. Math.* 54 (1996), no. 1, 85–104.
[20] P. Gravejat. First order asymptotics for the travelling waves in the Gross-Pitaevskii equation. *Adv. Differential Equations* 11(3):259–280, 2006.
[21] Gravejat, P., A non-existence result for supersonic travellingwaves in the Gross-Pitaevskii equation. *Commun.Math. Phys.* 243(1), 93–103 (2003).
[22] Gravejat, P., Limit at infinity and nonexistence results for sonic travelling waves in the Gross-Pitaevskii equation. *Differ. Int. Eqs.* 17(11–12), 1213–1232 (2004).
[23] Herve, R.-M.; Herve, M.; Etude qualitative des solutions reelles d’une equation differentielle liee ‘a l’équation de Ginzburg–Landau. *Ann. Inst. H. Poincare Anal. Non Lineaire* 11 (1994), no. 4, 427–440.

[24] C. Jones and P. Roberts. Motion in a Bose condensate IV. Axisymmetric solitary waves. *J. Phys. A: Math. Gen.* 15:2599–2619, 1982.

[25] Jones, C.A., Putterman, S.J., Roberts, P.H.: Motions in a Bose condensate V. Stability of solitarywave solutions of nonlinear Schrodinger equations in two and three dimensions. *J. Phys. A, Math.Gen.* 19, 2991–3011, 1986.

[26] J. B. Kaddeke; L. J. Campbell, Method for finding stationary states of point vortices, *Phys. Rev.* A36:4360–4370, 1987.

[27] Y. N. Ovchinnikov and I. M. Sigal. Ginzburg-Landau equation III. Vortex dynamics. *Nonlinearity* 11(5):1277–1294, 1998.

[28] Y. N. Ovchinnikov and I. M. Sigal. The energy of Ginzburg-Landau vortices, *Euro. Journal of Applied Mathematics* (2002), Vol. 13, 153–178.

[29] Lin, Tai-Chia. The stability of the radial solution to the Ginzburg-Landau equation. *Comm. Partial Differential Equations* 22 (1997), no. 3-4, 619–632.

[30] Lin, Fanghua; Wei, Juncheng. Traveling wave solutions of the Schrodinger map equation. *Comm. Pure Appl. Math.* 63 (2010), no. 12, 1585–1621.

[31] Lin, Fanghua; Wei, Juncheng, Superfluids passing an obstacle and vortex nucleation, preprint 2018.

[32] V. B. Matveev; M. A. Salle, Darboux transformations and solitons, Springer-Verlag, 1991.

[33] Mironescu, Petru. On the stability of radial solutions of the Ginzburg-Landau equation. *J. Funct. Anal.* 130 (1995), no. 2, 334–344.

[34] Mironescu, P. Les minimiseurs locaux pour l’équation de Ginzburg–Landau ont ‘a symetrie radiale. *C. R. Acad. Sci. Paris, Ser. I* 323 (1996), no 6, 593–598.

[35] Pacard, Frank; Riviere, Tristan. Linear and nonlinear aspects of vortices. The Ginzburg-Landau model. Progress in Nonlinear Differential Equations and their Applications, 39. Birkhauser Boston, Inc., Boston, MA, 2000.

[36] Sandier, Etienne; Serfaty, Sylvia. Vortices in the magnetic Ginzburg-Landau model. Progress in Nonlinear Differential Equations and their Applications, 70. Birkhauser Boston, Inc., Boston, MA, 2007.

[37] Shafrir, I. Remarks on solutions of $-\Delta u = \left( 1 - |u|^2 \right) y$ in $\mathbb{R}^2$, *C. R. Acad. Sci. Paris Ser. I Math.* 318 (1994), no. 4, 327–331.

**School of Mathematics and Physics, North China Electric Power University, Beijing, China**

*E-mail address: liuyong@ncepu.edu.cn*

**Department of Mathematics, University of British Columbia, Vancouver, B.C., Canada, V6T 1Z2**

*E-mail address: jcweimath.ubc.ca*