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REAL SUBMANIFOLDS OF MAXIMUM COMPLEX TANGENT SPACE AT A CR SINGULAR POINT, I

XIANGHONG GONG AND LAURENT STOLOVITCH

ABSTRACT. We study a germ of real analytic n -dimensional submanifold of \mathbf{C}^n that has a complex tangent space of maximal dimension at a CR singularity. Under some assumptions, we show its equivalence to a normal form under a local biholomorphism at the singularity. We also show that if a real submanifold is formally equivalent to a quadric, it is actually holomorphically equivalent to it, if a small divisors condition is satisfied. Finally, we investigate the existence of a complex submanifold of positive dimension in \mathbf{C}^n that intersects a real submanifold along two totally and real analytic submanifolds that intersect transversally at a possibly non-isolated CR singularity.

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1. INTRODUCTION AND MAIN RESULTS

1.1. Introduction. We are concerned with the local holomorphic invariants of a real analytic submanifold M in \mathbf{C}^n . The tangent space of M at a point x contains a maximal complex subspace of dimension d_x . When d_x is constant, M is called a Cauchy-Riemann (CR) submanifold. The CR submanifolds have been extensively studied since E. Cartan [Car32], [Car33], Tanaka [Tan62], and Chern-Moser [CM74].

We say that a point x_0 in a real submanifold M in \mathbf{C}^n is a CR singularity, if the complex tangent spaces $T_x M \cap J_x T_x M$ do not have a constant dimension in any neighborhood of x_0 . A real submanifold with a CR singularity must have codimension at least 2. The study

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of real submanifolds with CR singularities was initiated by E. Bishop in his pioneering work [Bis65]. He investigated a C^∞ real submanifold M of which the complex tangent space at a CR singularity is minimal, that is exactly one-dimensional. The very elementary models of this kind of manifolds are the Bishop quadrics Q_γ that depends on the Bishop invariant $0 \leq \gamma \leq \infty$, given by

$$Q_\gamma \subset \mathbf{C}^2: z_2 = |z_1|^2 + \gamma(z_1^2 + \bar{z}_1^2), \quad 0 \leq \gamma < \infty; \quad Q_\infty: z_2 = z_1^2 + \bar{z}_1^2.$$

The complex tangent at the origin is said to be *elliptic* if $0 \leq \gamma < 1/2$, *parabolic* if $\gamma = 1/2$, or *hyperbolic* if $\gamma > 1/2$. In [MW83], Moser and Webster studied the normal form problem of a real analytic surface M in \mathbf{C}^2 which is the higher order perturbation of Q . They showed that when $0 < \gamma < 1/2$, M is holomorphically equivalent, near the origin, to a normal form which is an algebraic surface that depends only on γ and two discrete invariants. We mention that the Moser-Webster normal form theory, as in Bishop's work, actually deals with an n -dimensional real submanifold M in \mathbf{C}^n , of which the complex tangent space has (minimum) dimension 1 at a CR singularity.

The main purpose of this paper is to investigate an n -dimensional real analytic submanifold M in \mathbf{C}^n , which is totally real outside a proper analytic subset and of which the complex tangent space has the *largest* possible dimension at a given CR singularity. We shall say that the singularity is a (maximal) complex tangent. The dimension must be $p = n/2$. Therefore, $n = 2p$ is even. We are interested in the normal form problem, the rigidity property, and the local analytic geometry of such real analytic manifolds.

In suitable holomorphic coordinates, a $2p$ -dimensional real analytic submanifold M in \mathbf{C}^{2p} that has a complex tangent space of maximum dimension at the origin is given by

$$(1.1) \quad M: z_{p+j} = E_j(z', \bar{z}'), \quad 1 \leq j \leq p,$$

where $z' = (z_1, \dots, z_p)$ and

$$E_j(z', \bar{z}') = h_j(z', \bar{z}') + q_j(\bar{z}') + O(|(z', \bar{z}')|^3).$$

Moreover, each $h_j(z', \bar{z}')$ is a homogeneous quadratic polynomial in z', \bar{z}' without holomorphic or anti-holomorphic terms, and each $q_j(\bar{z}')$ is a homogeneous quadratic polynomial in \bar{z}' . One of our goals is to seek suitable normal forms of perturbations of quadrics at the CR singularity (the origin).

1.2. Basic invariants. To study M , we consider its complexification in $\mathbf{C}^{2p} \times \mathbf{C}^{2p}$ defined by

$$\mathcal{M}: \begin{cases} z_{p+i} = E_i(z', w'), & i = 1, \dots, p, \\ w_{p+i} = \bar{E}_i(w', z'), & i = 1, \dots, p. \end{cases}$$

It is a complex submanifold of complex dimension $2p$ with coordinates $(z', w') \in \mathbf{C}^{2p}$. Let π_1, π_2 be the restrictions of the projections $(z, w) \rightarrow z$ and $(z, w) \rightarrow w$ to \mathcal{M} , respectively. Note that $\pi_2 = C\pi_1\rho_0$, where ρ_0 is the restriction to \mathcal{M} of the anti-holomorphic involution $(z, w) \rightarrow (\bar{w}, \bar{z})$ and C is the complex conjugate.

Our basic assumption is the following condition.

Condition B. $q(z') = (q_1(z'), \dots, q_p(z'))$ satisfies $q^{-1}(0) = \{0\}$.

When $p = 1$, condition B corresponds to the non-vanishing of the Bishop invariant γ . When $\gamma = 0$, Moser [Mos85] obtained a formal normal form that is still subject to further formal changes of coordinates. In [HY09a], Huang and Yin obtained a formal normal form and a complete holomorphic classification for real analytic surfaces with $\gamma = 0$. The formal normal forms for co-dimension two real submanifolds in \mathbf{C}^n have been further studied by Huang-Yin [HY12] and Burcea [Bur13]. Coffman [Cof06] showed that any m dimensional real analytic submanifold in \mathbf{C}^n of one-dimensional complex tangent space at a CR singularity satisfying certain non-degeneracy conditions is locally holomorphically equivalent to a unique algebraic submanifold, provided $2(n+1)/3 \leq m < n$.

When M is a *quadric*, i.e. all E_j in (1.1) are quadratic polynomials, our basic condition B is equivalent to π_1 being a 2^p -to-1 branched covering. Since $\pi_2 = C\pi_1\rho_0$, then π_2 is also a 2^p -to-1 branched covering. We will see that the CR singularities of the real submanifolds are closely connected with these branched coverings and their deck transformations.

We now introduce our main results. Some of them are analogous to the Moser-Wester theory. We will also describe new situations which arise with the maximum complex tangency.

1.3. Branched coverings and deck transformations. In section 2, we study the existence of deck transformations for π_1 . We show that they must be involutions generating an abelian group of order 2^k for some $0 \leq k \leq p$. The latter is a major difference with the case $p = 1$. Indeed, in the Moser-Webster theory, the group of deck transformations is generated by a unique non-trivial involution τ_1 . Therefore, we will impose the following condition.

Condition D. *M satisfies condition B and the branched covering π_1 of \mathcal{M} admits the maximum 2^p deck transformations.*

Condition D gives rise to two families of commuting involutions $\{\tau_{i1}, \dots, \tau_{i2^p}\}$ intertwined by the anti-holomorphic involution $\rho_0: (z', w') \rightarrow (\bar{w}', \bar{z}')$ such that $\tau_{2j} = \rho_0\tau_{1j}\rho_0$ ($1 \leq j \leq 2^p$) are deck transformations of π_2 . We will call $\{\tau_{11}, \dots, \tau_{12^p}, \rho_0\}$ the set of *Moser-Webster involutions*. We will show that there is a unique set of p generators for the deck transformations of π_1 , denoted by $\tau_{11}, \dots, \tau_{1p}$, such that each τ_{1j} fixes a hypersurface in \mathcal{M} pointwise. Then

$$\tau_1 = \tau_{11} \circ \dots \circ \tau_{1p}$$

is the unique deck transformation of which the fixed-point set has the smallest dimension p . Let $\tau_2 = \rho_0\tau_1\rho_0$ and

$$\sigma = \tau_1\tau_2.$$

Then σ is *reversible* by τ_j and ρ_0 , i.e. $\sigma^{-1} = \tau_j\sigma\tau_j^{-1}$ and $\sigma^{-1} = \rho_0\sigma\rho_0$.

As in the Moser-Webster theory, we will show that the existence of such 2^p deck transformations transfers the normal form problem for the real submanifolds into the normal form problem for the sets of involutions $\{\tau_{11}, \dots, \tau_{1p}, \rho_0\}$.

In this paper we will make the following assumption.

Condition J. *M satisfies condition D and M is diagonalizable, i.e. $\sigma'(0)$ is diagonalizable.*

Note that the condition excludes the higher dimensional analogous complex tangency of *parabolic* type, i.e. of $\gamma = 1/2$. The normal form problem for the parabolic complex

tangents has been studied by Webster [Web92], and in [Gon96] where the normalization is divergent in general. In [AG09], Ahern and Gong constructed a moduli space for real analytic submanifolds that are formally equivalent to the Bishop quadric with $\gamma = 1/2$.

1.4. Product quadrics. In this paper, the basic model for quadric manifolds with a CR singularity satisfying condition J is a product of 3 types of quadrics defined by

$$(1.2) \quad Q_{\gamma_e} \subset \mathbf{C}^2: z_2 = (z_1 + 2\gamma_e \bar{z}_1)^2;$$

$$(1.3) \quad Q_{\gamma_h} \subset \mathbf{C}^2: z_2 = (z_1 + 2\gamma_h \bar{z}_1)^2, \quad 1/2 < \gamma_h < \infty; \quad Q_\infty: z_2 = z_1^2 + \bar{z}_1^2;$$

$$(1.4) \quad Q_{\gamma_s} \subset \mathbf{C}^4: z_3 = (z_1 + 2\gamma_s \bar{z}_2)^2, \quad z_4 = (z_2 + 2(1 - \bar{\gamma}_s) \bar{z}_1)^2.$$

Here $\gamma_s \in \mathbf{C}$ and

$$(1.5) \quad 0 < \gamma_e < 1/2, \quad 1/2 < \gamma_h \leq \infty, \quad \operatorname{Re} \gamma_s \leq 1/2, \quad \operatorname{Im} \gamma_s \geq 0, \quad \gamma_s \neq 0, 1/2.$$

Note that $Q_{\gamma_e}, Q_{\gamma_h}$ are elliptic and hyperbolic Bishop quadrics, respectively. Realizing a type of pairs of involutions introduced in [Sto07], we will say that the complex tangent of Q_{γ_s} at the origin is *complex*. We emphasize that this last type of quadric is new as it is not holomorphically equivalent to a product of two Bishop surfaces. A product of the above quadrics will be called a *product of quadrics*, or a *product quadric*. We denote by $e_*, h_*, 2s_*$ the number of elliptic, hyperbolic and complex coordinates, respectively. We remark that the complex tangent of complex type has another basic model Q_{γ_s} with $\gamma_s = 1/2$, which is excluded by condition J (see Proposition 2.10).

This is the first part of two papers devoted to the local study of real analytic manifold at maximal complex tangent point. To limit its scope, we have to leave the complete classification of quadratic submanifolds of maximum deck transformations to the second paper [GS15] (see Theorem 1.1 therein), showing that there are quadratic manifolds which are not holomorphically equivalent to a product quadric. In [GS15], we also show that all Poincaré-Dulac normal forms of the σ of a general higher order perturbation of a product quadric are divergent when $p > 1$. With the divergent Poincaré-Dulac normal forms at our disposal, we seek types of CR singularities that ensure the convergent normalization and the analytic structure of the hull of holomorphy associated with the types of CR singularities.

We now introduce our main geometrical and dynamical results for analytic higher order perturbations of product quadric. We first turn to a holomorphic normalization of a real analytic submanifold M with the so-called abelian CR singularity. This will be achieved by studying an integrability problem on a general family of commuting biholomorphisms described below. The holomorphic normalization will be used to construct the local hull of holomorphy of M . We will also study the rigidity problem of a quadric under higher order analytic perturbations, i.e. the problem if such a perturbation remains holomorphically equivalent to the quadric if it is formally equivalent to the quadric. The rigidity problem is reduced to a theorem of holomorphic linearization of one or several commuting diffeomorphisms that was devised in [Sto15]. Finally, we will study the existence of holomorphic submanifolds attached to the real submanifold M . These are complex submanifolds of dimension p intersecting M along two totally real analytic submanifolds that intersect transversally at a CR singularity. Attaching complex submanifolds has less constraints

than finding a convergent normalization. A remarkable feature of attached complex submanifolds is that their existence depends only on the existence of suitable (convergent) invariant submanifolds of σ .

1.5. Normal form of commuting biholomorphisms.

Definition 1.1. Let $\mathcal{F} = \{F_1, \dots, F_\ell\}$ be a finite family of germs of biholomorphisms of \mathbf{C}^n fixing the origin. Let D_m be the linear part of F_m at the origin. We say that the family \mathcal{F} is (resp. formally) *completely integrable*, if there is a (resp. formal) biholomorphic mapping Φ such that $\{\Phi^{-1}F_m\Phi: 1 \leq m \leq \ell\} = \{\hat{F}_m: 1 \leq m \leq \ell\}$ satisfies

- (i) $\hat{F}_m(z) = (\mu_{m1}(z)z_1, \dots, \mu_{mn}(z)z_n)$ where μ_{mj} are germs of holomorphic (resp. formal) functions such that $\mu_{mj} \circ D_{m'} = \mu_{mj}$ for $1 \leq m, m' \leq \ell$ and $1 \leq j \leq n$. In particular, \hat{F}_m commutes with $D_{m'}$ for all $1 \leq m, m' \leq \ell$.
- (ii) For each j and each $Q \in \mathbf{N}^n$ with $|Q| > 1$, $\mu_m^Q(0) = \mu_{mj}(0)$ hold for all m if and only if $\mu_m^Q(z) = \mu_{mj}(z)$ hold for all m .

A necessary condition for \mathcal{F} to be formally completely integrable is that F_1, \dots, F_ℓ commute pairwise. The main result of section 4 is the following.

Theorem 1.2. *Let \mathcal{F} be a family of finitely many germs of biholomorphisms at the origin. If \mathcal{F} is formally completely integrable and its linear part \mathcal{D} has the Poincaré type, then it is holomorphically completely integrable.*

The definition of Poincaré type is in Definition 4.11. Such a formal integrability condition can hold under some geometrical properties. For instance, for a single germ of real analytic hyperbolic area-preserving mapping, the result was due to Moser [Mos56], and for a single germ of reversible hyperbolic holomorphic mapping $\sigma = \tau_1\tau_2$ of which τ_1 fixes a hypersurface, this result was due to Moser-Webster [MW83]. Such results for commuting germs of vector fields were obtained in [Sto00, Sto05] under a collective small divisors Brjuno-type condition. Our result is inspired by these results.

1.6. Holomorphic normalization for the abelian CR singularity. In section 5, we obtain the convergent normalization for an *abelian* CR singularity which we now define. We first consider a product quadric Q which satisfies condition J. So the deck transformations of π_1 for the complexification of Q are generated by p involutions of which each fixes a hypersurface pointwise. We denote them by T_{11}, \dots, T_{1p} . Let $T_{2j} = \rho T_{1j} \rho$. It turns out that each T_{1j} commutes with all T_{ik} except one, T_{2k_j} for some $1 \leq k_j \leq p$. When we formulate $S_j = T_{1j}T_{2k_j}$ for $1 \leq j \leq p$, the S_1, \dots, S_p commute pairwise. Consider a general M that is a third-order perturbation of product quadric Q and satisfies condition J. We define $\sigma_j = \tau_{1j}\tau_{2k_j}$. In suitable coordinates, T_{ij} (resp. S_j) is the linear part of τ_{ij} (resp. σ_j) at the origin. We say that the complex tangent of a third order perturbation M of a product quadric at the origin is of *abelian type*, if $\sigma_1, \dots, \sigma_p$ commute pairwise. If each linear part S_j of σ_j has exactly two eigenvalues μ_j, μ_j^{-1} that are different from 1, then $\mathcal{S} := \{S_1, \dots, S_p\}$ is of Poincaré type if and only if $|\mu_j| \neq 1$ for all j . As mentioned previously, Moser and Webster actually dealt with n -dimensional real submanifolds in \mathbf{C}^n that have the minimal dimension of complex tangent subspace at a CR singular point. In their situation, there is only one possible composition, that is $\sigma = \tau_1\tau_2$. When the complex tangent has an

elliptic but non-vanishing Bishop invariant, σ has exactly two positive eigenvalues that are separated by 1, while the remaining eigenvalues are 1 with multiplicity $n - 2$.

As an application of Theorem 1.2, we will prove the following convergent normalization.

Theorem 1.3. *Let M be a germ of real analytic submanifold in \mathbf{C}^{2p} that is a third order perturbation of a product quadric given by (1.2)-(1.5) with an abelian CR singularity. Suppose that M has all eigenvalues of modulus different from one, i.e. it has no hyperbolic component ($h_* = 0$) while each γ_s in (1.4) satisfies $\operatorname{Re} \gamma_s < 1/2$ additionally. Then M is holomorphically equivalent to*

$$\widehat{M}: z_{p+j} = \Lambda_{1j}(\zeta)\zeta_j, \quad \Lambda_{1j}(0) = \lambda_j, \quad 1 \leq j \leq p,$$

where $\zeta = (\zeta_1, \dots, \zeta_p)$ are the solutions to

$$\begin{aligned} \zeta_e &= A_e(\zeta)z_e\bar{z}_e - B_e(\zeta)(z_e^2 + \bar{z}_e^2), \quad 1 \leq e \leq e_*, \\ \zeta_s &= A_s(\zeta)z_s\bar{z}_{s+s_*} - B_s(\zeta)(z_s^2 + \Lambda_{1s}^2(\zeta)\bar{z}_{s+s_*}^2), \quad e_* < s \leq e_* + s_*, \\ \zeta_{s+s_*} &= A_{s+s_*}(\zeta)\bar{z}_s z_{s+s_*} - B_{s+s_*}(\zeta)(z_{s+s_*}^2 + \Lambda_{1(s+s_*)}^2(\zeta)\bar{z}_s^2), \end{aligned}$$

while Λ_{1j} satisfies (5.6)-(5.7), and A_j, B_j are rational functions in Λ_{1j} defined by (5.16)-(5.17).

There are many non-product real submanifolds of abelian CR singularity.

Example 1.4. Let $0 < \gamma_i < \infty$. Let $R(z_1, \bar{z}_1) = |z_1|^2 + \gamma_1(z_1^2 + \bar{z}_1^2) + O(3)$ be a real-valued power series in z_1, \bar{z}_1 of real coefficients. Then the origin is an abelian CR singularity of

$$M: z_3 = R(z_1, \bar{z}_1), \quad z_4 = (z_2 + 2\gamma_2\bar{z}_2 + z_2z_3)^2.$$

We will also present a more direct proof of Theorem 1.3 by using a convergence theorem of Moser and Webster [MW83] and some formal results from section 4. The above $\Lambda_{11}, \dots, \Lambda_{1p}$ satisfy conditions $\Lambda_{1j}(0) = \lambda_j$ and (5.6)-(5.7) and are otherwise *arbitrary* convergent power series. The $\Lambda_{11}, \dots, \Lambda_{1p}$ may be subjected to further normalization. In [GS15], we find a unique holomorphic normal form by refining the above normalization for M satisfying a non resonance condition and a third order non-degeneracy condition (see Theorem 5.6 in [GS15]); in particular, it shows the existence of infinitely many formal invariants and non-product structures of the manifolds when $p > 1$.

As an application of Theorem 1.3, we will prove the following flattening result.

Corollary 1.5. *Let M be as in Theorem 1.3. In suitable holomorphic coordinates, M is contained in the linear subspace defined by $z_{p+e} = \bar{z}_{p+e}$ and $z_{p+s} = \bar{z}_{p+s+s_*}$ where $1 \leq e \leq e_*$ and $e_* < s \leq e_* + s_*$.*

1.7. Analytic hull of holomorphy. One of significances of the Bishop quadrics is that their higher order analytic perturbation at an elliptic complex tangent has a non-trivial hull of holomorphy. As another application of the above normal form, we will construct the local hull of holomorphy of M via higher dimensional non-linear analytic polydiscs.

Corollary 1.6. *Let M be as in Theorem 1.3. Suppose that M has only elliptic component of complex tangent. Then in suitable holomorphic coordinates, $\mathcal{H}_{loc}(M)$, the local hull of holomorphy of M , is filled by a real analytic family of analytic polydiscs of dimension p .*

For a precise statement of the corollary, see Theorem 5.5. The hulls of holomorphy for real submanifolds with a CR singularity have been studied extensively, starting with the work of Bishop. In the real analytic case with minimum complex tangent space at an elliptic complex tangent, we refer to Moser-Webster [MW83] for $\gamma > 0$, and Krantz-Huang [HK95] for $\gamma = 0$. For the smooth case, see Kenig-Webster [KW82, KW84], Huang [Hua98]. For global results on hull of holomorphy, we refer to [BG83, BK91].

1.8. Rigidity of quadrics. In Section 6, we prove the following theorem.

Theorem 1.7. *Let M be a germ of analytic submanifold that is an higher order perturbation of a product quadric Q in \mathbf{C}^{2p} given by (1.2)-(1.5). Assume that M is formally equivalent to Q . Suppose that each hyperbolic component has an eigenvalue μ_h which is either a root of unity or satisfies Brjuno condition, and each complex component has an eigenvalue μ_s is not a root of unity and satisfies the Brjuno condition. Then M is holomorphically equivalent to the product quadric.*

We emphasize that condition (1.5) ensure that M is diagonalizable (condition J). It is plausible that theorem remains valid when M satisfies condition J and μ_s satisfy the Brjuno condition or are roots of unity; however, the resonance condition requires some tedious changes of computation in section 6. The proof uses a theorem of linearization of holomorphic mappings in [Sto15]. Brjuno small divisors condition is defined by (7.37), with $\nu = \mu_h$ and $p = 1$. When $p = 1$, the result under the stronger Siegel condition is in [Gon94]. This last statement requires a small divisors condition to be true as shown in [Gon04]. When $p = 1$ with a vanishing Bishop invariant, such rigidity result was obtained by Moser [Mos85] and by Huang-Yin [HY09b] in a more general context.

1.9. Attached complex submanifolds. We now describe convergent results for attached complex submanifolds. The results are for a general M , including the one of which the complex tangent might not be of abelian type.

We say that a formal complex submanifold K is *attached* to M if $K \cap M$ contains at least two germs of totally real and formal submanifolds K_1, K_2 that intersect transversally at a given CR singularity. In [Kli85], Klingenberg showed that when M is non-resonant and $p = 1$, there is a unique formal holomorphic curve attached to M with a hyperbolic complex tangent. He also proved the convergence of the attached formal holomorphic curve under a Siegel small divisors condition. When $p > 1$, we will show that generically there is no formal complex submanifold that can be attached to M if the CR singularity has an elliptic component. When $p > 1$ and M is a higher order perturbation of a product quadric of $Q_{\gamma_h}, Q_{\gamma_s}$, we will encounter various interesting situations.

By adapting Klingenberg's proof for $p = 1$ and using a theorem of Pöschel [Pös86], we will prove the following.

Theorem 1.8. *Let M be a germ of analytic submanifold that is an higher order perturbation of a product quadric Q in \mathbf{C}^{2p} without elliptic components. Assume that the eigenvalues $\mu_1, \dots, \mu_p, \mu_1^{-1}, \dots, \mu_p^{-1}$ of $D\sigma(0)$ are distinct. Let $\epsilon_h^2, \epsilon_s^2 = 1$, $\nu_h := \mu_h^{\epsilon_h}$, $\nu_s := \bar{\mu}_s^{\epsilon_s}$ and $\nu_{s+s^*} := \bar{\nu}_s^{-1}$. Assume $\nu = (\nu_1, \dots, \nu_p)$ is weakly non resonant and Diophantine in the sense of Pöschel. Then M admits an attached complex submanifold M_ϵ .*

Weak non resonance is defined in (7.34), while Diophantine condition in the sense of Pöschel is defined in (7.37).

Finally, we prove the convergence of *all* attached formal submanifolds:

Theorem 1.9. *Let M be as in Theorem 1.8. Suppose that the $2p$ eigenvalues of σ are non-resonant. If the eigenvalues of σ satisfy a Bruno type condition, all attached formal submanifolds are convergent.*

The Brjuno-type condition, defined in (7.38), was introduced in [Sto15] for linearization on ideals.

1.10. Notation. We denote the identity map by I and by LF the linear part at the origin of a mapping $F: (\mathbf{C}^m, 0) \rightarrow (\mathbf{C}^n, 0)$. We also denote by $DF(z)$ or $F'(z)$, the Jacobian matrix of F at z . By an analytic (or holomorphic) function, we mean a *germ* of analytic function at a point (which will be defined by the context) otherwise stated. We denote by \mathcal{O}_n (resp. $\widehat{\mathcal{O}}_n$, \mathfrak{M}_n , $\widehat{\mathfrak{M}}_n$) the space of germs of holomorphic functions of \mathbf{C}^n at the origin (resp. of formal power series in \mathbf{C}^n , holomorphic germs, and formal germs vanishing at the origin). If $Q = (q_1, \dots, q_k) \in \mathbf{N}^k$, then $|Q| = q_1 + \dots + q_k$ and $x^Q = x_1^{q_1} \dots x_k^{q_k}$.

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2. CR SINGULARITIES AND DECK TRANSFORMATIONS

We consider a real submanifold M of \mathbf{C}^n . Let $T_{x_0}^{(1,0)}M$ be the space of tangent vectors of M at x_0 of the form $\sum_{j=1}^n a_j \frac{\partial}{\partial z_j}$. Let M have dimension n . In this paper, we assume that $T_{x_0}^{(1,0)}M$ has the largest possible dimension $p = n/2$ at a given point x_0 . In suitable holomorphic affine coordinates, we have $x_0 = 0$ and

$$(2.1) \quad M: z_{p+j} = E_j(z', \bar{z}'), \quad 1 \leq j \leq p.$$

Here we set $z' = (z_1, \dots, z_p)$ and we will denote $z'' = (z_{p+1}, \dots, z_{2p})$. Also, the E_j together with their first order derivatives vanish at 0. The tangent space T_0M is then the z' -subspace.

The main purpose of this section is to obtain some basic invariants and a relation between two families of involutions and the real analytic submanifolds which we want to normalize.

Note that M is totally real at $(z', z'') \in M$ if and only if $C(z', \bar{z}') \neq 0$, where $C(z', \bar{z}') := \det\left(\frac{\partial E_i}{\partial \bar{z}_j}\right)_{1 \leq i, j \leq p}$. We will assume that $C(z', \bar{z}')$ is not identically zero in any neighborhood of the origin. Then the zero set of C on M , denoted by M_{CRsing} , is called *CR singular set* of M , or the set of *complex tangents* of M . We assume that M is real analytic. Then M_{CRsing} is a possibly singular proper real analytic subset of M that contains the origin.

2.1. Existence of deck transformations and examples. We first derive some quadratic invariants. Applying a quadratic change of holomorphic coordinates, we obtain

$$(2.2) \quad E_j(z', \bar{z}') = h_j(z', \bar{z}') + q_j(\bar{z}') + O(|(z', \bar{z}')|^3).$$

Here we have used the convention that if $x = (x_1, \dots, x_n)$, then $O(|x|^k)$ denotes a formal power series in x without terms of order $< k$. A biholomorphic map f that preserves the form of the above submanifolds M and fixes the origin must preserve their complex tangent

spaces at the origin, i.e. $z'' = 0$. Thus if \tilde{z} denote the old coordinates and z denote the new coordinates then f has the form

$$\tilde{z}' = \mathbf{A}z' + \mathbf{B}z'' + O(|z|^2), \quad \tilde{z}'' = \mathbf{U}z'' + O(|z|^2).$$

Here \mathbf{A} and \mathbf{U} are non-singular $p \times p$ complex matrices. Now $f(M)$ is given by

$$\mathbf{U}z'' = h(\mathbf{A}z', \overline{\mathbf{A}z'}) + q(\overline{\mathbf{A}z'}) + O(|z|^3).$$

We multiply the both sides by \mathbf{U}^{-1} and solve for z'' ; the vectors of p quadratic forms $\{h(\tilde{z}', \overline{\tilde{z}'}), q(\tilde{z}')\}$ are transformed into

$$(2.3) \quad \{\hat{h}(z', z'), \hat{q}(\overline{z}')\} = \{\mathbf{U}^{-1}h(\mathbf{A}z', \overline{\mathbf{A}z'}), \mathbf{U}^{-1}q(\overline{\mathbf{A}z'})\}.$$

This shows that if M and \hat{M} are holomorphically equivalent, their corresponding quadratic terms are equivalent via (2.3). Therefore, we obtain a holomorphic invariant

$$q_* = \dim_{\mathbf{C}}\{z' : q_1(z') = \dots = q_p(z') = 0\}.$$

We remark that when M, \hat{M} are quadratic (i.e. when their corresponding E, \hat{E} are homogeneous quadratic polynomials), the equivalence relation (2.3) implies that M, \hat{M} are linearly equivalent, Therefore, the above transformation of h and q via \mathbf{A} and \mathbf{U} determines the classifications of the quadrics under local biholomorphisms as well as under linear biholomorphisms. We have shown that the two classifications for the quadrics are identical.

Recall that M is real analytic. Let us complexify such a real submanifold M by replacing \tilde{z}' by w' to obtain a complex n -submanifold of \mathbf{C}^{2n} , defined by

$$\mathcal{M}: \begin{cases} z_{p+i} = E_i(z', w'), \\ w_{p+i} = \bar{E}_i(w', z'), \quad i = 1, \dots, p. \end{cases}$$

We use (z', w') as holomorphic coordinates of \mathcal{M} and define the anti-holomorphic involution ρ on it by

$$(2.4) \quad \rho(z', w') = (\bar{w}', \bar{z}').$$

Occasionally we will also denote the above ρ by ρ_0 for clarity. We will identify M with a totally real and real analytic submanifold of \mathcal{M} via embedding $z \rightarrow (z, \bar{z})$. We have $M = \mathcal{M} \cap \text{Fix}(\rho)$ where $\text{Fix}(\rho)$ denotes the set of fixed points of ρ . Let $\pi_1: \mathcal{M} \rightarrow \mathbf{C}^n$ be the restriction of the projection $(z, w) \rightarrow z$ and let π_2 be the restriction of $(z, w) \rightarrow w$. It is clear that $\pi_2 = \overline{\pi_1} \rho$ on \mathcal{M} . Throughout the paper, π_1, π_2, ρ are restricted on \mathcal{M} unless stated otherwise.

Condition B that $q_* = 0$, introduced in section 1, ensures that π_1 is a branched covering. A necessary condition for $q_* = 0$ is that functions $q_1(z'), q_2(z'), \dots, q_p(z')$ are linearly independent, since the intersection of k germs of holomorphic hypersurfaces at 0 in \mathbf{C}^p has dimension at least $p - k$. (See [Chi89], p. 35; [Gun90][Corollary 8, p. 81].)

When $\pi_1: \mathcal{M} \rightarrow \mathbf{C}^{2p}$ is a branched covering, we define a *deck transformation* on \mathcal{M} for π_1 to be a germ of biholomorphic mapping F defined at $0 \in \mathcal{M}$ that satisfies $\pi_1 \circ F = \pi_1$. In other words, $F(z', w') = (z', f(z', w'))$ and

$$E_i(z', w') = E_i(z', f(z', w')), \quad i = 1, \dots, p.$$

Lemma 2.1. *Suppose that $q_* = 0$. Then M_{CRsing} is a proper real analytic subset of M and M is totally real away from M_{CRsing} , i.e. the CR dimension of M is zero. Furthermore, π_1 is a 2^p -to-1 branched covering. The group of deck transformations of π_1 consists of 2^ℓ commuting involutions with $0 \leq \ell \leq p$.*

Proof. Since $q^{-1}(0) = \{0\}$, then $z' \rightarrow q(z')$ is a finite holomorphic map; see [Chi89], p. 105. Hence its Jacobian determinant is not identically zero. In particular, $\det(\frac{\partial E_i}{\partial \bar{z}_j})_{1 \leq i, j \leq p}$ is not identically zero. This shows that M has CR dimension 0.

Since $w' \rightarrow q(w')$ is a homogeneous quadratic mapping of the same space which vanishes only at the origin, then

$$|q(w')| \geq c|w'|^2.$$

We want to verify that π_1 is a 2^p -to-1 branched covering. Let $\Delta_r = \{z \in \mathbf{C} : |z| < r\}$. We choose $C > 0$ such that $\pi_1(z, w) = (z', E(z', w'))$ defines a proper and onto mapping

$$(2.5) \quad \pi_1: \mathcal{M}_1 := \mathcal{M} \cap ((\Delta_\delta^p \times \Delta_{\delta^2}^p) \times (\Delta_{C\delta}^p \times \Delta_{C\delta^2}^p)) \mapsto \Delta_\delta^p \times \Delta_{\delta^2}^p.$$

By Sard's theorem, the regular values of π_1 have the full measure. For each regular value z , $\pi_1^{-1}(z)$ has exactly 2^p distinct points (see [Chi89], p. 105 and p. 112). It is obvious that \mathcal{M}_1 is smooth and connected. We fix a fiber F_z of 2^p points. Then the group of deck transformations of π_1 acts on F_z in such a way that if a deck transformation fixes a point in F_z , then it must be the identity. Therefore, the number of deck transformations divides 2^p and each deck transformation has period 2^ℓ with $0 \leq \ell \leq p$.

We first show that each deck transformation f of π_1 is an involution. We know that f is periodic and has the form

$$z' \rightarrow z', \quad w' \rightarrow \mathbf{A}w' + \mathbf{B}z' + O(2),$$

where \mathbf{A}, \mathbf{B} are matrices. Assume that f has period m . Then $\hat{f}(z', w') = (z', \mathbf{A}w' + \mathbf{B}z')$ satisfies $\hat{f}^m = I$ and f is locally equivalent to \hat{f} ; indeed $\hat{f}g f^{-1} = g$ for

$$g = \sum_{i=1}^m (\hat{f}^i)^{-1} \circ f^i.$$

Therefore, it suffices to show that \hat{f} is an involution. We have

$$\hat{f}^m(z', w') = (z', \mathbf{A}^m w' + (\mathbf{A}^{m-1} + \dots + \mathbf{A} + \mathbf{I})\mathbf{B}z').$$

Since f is a deck transformation, then $E(z', w')$ is invariant under f . Recall from (2.2) that $E(z', \bar{z}')$ starts with quadratic terms of the form $h(z', \bar{z}') + q(\bar{z}')$. Comparing quadratic terms in $E(z', w') = E \circ \hat{f}(z', w')$, we see that the linear map \hat{f} has invariant functions

$$z'' = h(z', w') + q(w').$$

We know that $\mathbf{A}^m = \mathbf{I}$. By the Jordan normal form, we choose a linear transformation $\tilde{w}' = \mathbf{S}w'$ such that \mathbf{SAS}^{-1} is the diagonal matrix $\text{diag } \mathbf{a}$ with $\mathbf{a} = (a_1, \dots, a_p)$. In (z', \tilde{w}') coordinates, the mapping \hat{f} has the form $(z', \tilde{w}') \rightarrow (z', (\text{diag } \mathbf{a})\tilde{w}' + \mathbf{SB}z')$. Now

$$\tilde{h}_j(z', \tilde{w}') + \tilde{q}_j(\tilde{w}') := h_j(z', \mathbf{S}^{-1}\tilde{w}') + q_j(\mathbf{S}^{-1}\tilde{w}')$$

are invariant under \hat{f} . Hence $\tilde{q}_j(\tilde{w}')$ are invariant under $\tilde{w}' \mapsto (\text{diag } \mathbf{a})\tilde{w}'$. Since the common zero set of $q_1(w'), \dots, q_p(w')$ is the origin, then

$$V = \{\tilde{w}' \in \mathbf{C}^p: \tilde{q}(\tilde{w}') = 0\} = \{0\}.$$

We conclude that $\tilde{q}(\tilde{w}_1, 0, \dots, 0)$ is not identically zero; otherwise V would contain the \tilde{w}_1 -axis. Now $\tilde{q}((\text{diag } \mathbf{a})\tilde{w}') = \tilde{q}(\tilde{w}')$, restricted to $\tilde{w}' = (\tilde{w}_1, 0, \dots, 0)$, implies that $a_1 = \pm 1$. By the same argument, we get $a_j = \pm 1$ for all j . This shows that $\mathbf{A}^2 = \mathbf{I}$. Let us combine it with

$$\mathbf{A}^m = \mathbf{I}, \quad (\mathbf{A}^{m-1} + \dots + \mathbf{A} + \mathbf{I})\mathbf{B} = \mathbf{0}.$$

If $m = 1$, it is obvious that $\hat{f} = I$. If $m = 2^\ell > 1$, then $(\mathbf{A} + \mathbf{I})\mathbf{B} = \mathbf{0}$. Thus $\hat{f}^2(z', w') = (z', \mathbf{A}^2 w' + (\mathbf{A} + \mathbf{I})\mathbf{B}z') = (z', w')$. This shows that every deck transformation of π_1 is an involution.

For any two deck transformations f and g , fg is still a deck transformation. Hence $(fg)^2 = I$ implies that $fg = gf$. \square

Before we proceed to discussing the deck transformations, we give some examples. The first example turns out to be a holomorphic equivalent form of a real submanifold that admits the maximum number of deck transformations and satisfies other mild conditions.

Example 2.2. Let $\mathbf{B} = (b_{jk})$ be a non-singular $p \times p$ matrix. Let M be defined by

$$(2.6) \quad z_{p+j} = \left(\sum_k b_{jk} \bar{z}_k + R_j(z', \bar{z}') \right)^2, \quad 1 \leq j \leq p,$$

where each $R_j(0, \bar{z}')$ starts with terms of order at least 2. Then M admits 2^p deck transformations for π_1 . Indeed, let $\mathbf{E}_1, \dots, \mathbf{E}_{2^p}$ be the set of diagonal $p \times p$ matrices with $\mathbf{E}_j^2 = \mathbf{I}$, and let \mathbf{R} be the column vector $(R_1, \dots, R_p)^t$. For each \mathbf{E}_j let us show that there is a deck transformation $(z', w') \rightarrow (z', \tilde{w}')$ satisfying

$$(2.7) \quad \mathbf{B}\tilde{w}' + \mathbf{R}(z', \tilde{w}') = \mathbf{E}_j(\mathbf{B}w' + \mathbf{R}(z', w')).$$

Since \mathbf{B} is invertible, it has a unique solution

$$\tilde{w}' = \mathbf{B}^{-1}\mathbf{E}_j\mathbf{B}w' + O(|z'|) + O(|w'|^2).$$

Finally, $(z', w') \rightarrow (z', \tilde{w}')$ is an involution, as if $(z', w', \tilde{w}') = (z', w', f(z', w'))$ satisfy (2.7) if and only if $(z', f(z', w'), w')$, substituting for (z', w', \tilde{w}') in (2.7), satisfy (2.7).

We now present an example to show that the deck transformations can be destroyed by perturbations when $p > 1$. This is the major difference between real submanifolds with $p > 1$ and the ones with $p = 1$. The example shows that the number of deck transformations can be reduced to any number 2^ℓ by a higher order perturbation.

Example 2.3. Let $N_{\gamma, \epsilon}$ be a perturbation of Q_γ defined by

$$z_{p+j} = z_j \bar{z}_j + \gamma_j \bar{z}_j^2 + \epsilon_{j-1} \bar{z}_{j-1}^3, \quad 1 \leq j \leq p.$$

Here $\epsilon_j \neq 0$ for all j , $\epsilon_0 = \epsilon_p$ and $z_0 = z_p$. Let τ be a deck transformation of $N_{\gamma, \epsilon}$ for π_1 . We know that τ has the form

$$z'_j = z_j, \quad w'_j = A_j(z', w') + B_j(z', w') + O(|(z', w')|^3).$$

Here A_j are linear and B_j are homogeneous quadratic polynomials. We then have

$$(2.8) \quad z_j A_j(z', w') + \gamma_j A_j^2(z', w') = z_j w_j + \gamma_j w_j^2,$$

$$(2.9) \quad z_j B_j(z', w') + 2\gamma_j (A_j B_j)(z', w') + \epsilon_{j-1} A_{j-1}^3(z', w') = \epsilon_{j-1} w_{j-1}^3.$$

We know that $L\tau$ is a deck transformation for Q_γ . Thus $a_j(w') := A_j(0, w') = \pm w_j$. Set $z_j = 0$ in (2.9) to get $a_j(w') | \epsilon_{j-1} (w_{j-1}^3 - a_{j-1}^3(w'))$. Thus $a_{j-1}(w') = w_{j-1}$. Hence, the matrix of $L\tau$ is triangular and its diagonal entries are 1. Since $L\tau$ is periodic then $L\tau = I$. Since τ is periodic, then $\tau = I$.

Based the above example, we impose the basic condition D that the branched covering π_1 of \mathcal{M} admits the maximum 2^p deck transformations.

We first derive some properties of real submanifolds under condition D.

2.2. Real submanifolds and Moser-Webster involutions. The main result of this subsection is to show the equivalence of classification of the real submanifolds satisfying condition D with that of families of involutions $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$. More precisely, condition J is not imposed. The relation between two classifications plays an important role in the Moser-Webster theory for $p = 1$. This will be the base of our approach to the normal form problems.

Let \mathcal{F} be a family of holomorphic maps in \mathbf{C}^n with coordinates z . Let $L\mathcal{F}$ denote the set of linear maps $z \rightarrow f'(0)z$ with $f \in \mathcal{F}$. Let $\mathcal{O}_n^{\mathcal{F}}$ denote the set of germs of holomorphic functions h at $0 \in \mathbf{C}^n$ so that $h \circ f = h$ for each $f \in \mathcal{F}$. Let $[\mathfrak{M}_n]_1^{L\mathcal{F}}$ be the subset of linear functions of $\mathfrak{M}_n^{L\mathcal{F}}$.

Lemma 2.4. *Let G be an abelian group of holomorphic (resp. formal) involutions fixing $0 \in \mathbf{C}^n$. Then G has 2^ℓ elements and they are simultaneously diagonalizable by a holomorphic (resp. formal) transformation. If $k = \dim_{\mathbf{C}}[\mathfrak{M}_n]_1^{LG}$ then $\ell \leq n - k$. Assume furthermore that $\ell = n - k$. In suitable holomorphic (z_1, \dots, z_n) coordinates, the group G is generated by Z_{k+1}, \dots, Z_n with*

$$(2.10) \quad Z_j: z'_j = -z_j, \quad z'_i = z_i, \quad i \neq j, \quad 1 \leq i \leq n.$$

In the z coordinates, the set of convergent (resp. formal) power series in $z_1, \dots, z_k, z_{k+1}^2, \dots, z_n^2$ is equal to \mathcal{O}_n^G (resp. $\widehat{\mathcal{O}}_n^G$), and with $Z = Z_{n-k} \cdots Z_n$,

$$(2.11) \quad [\mathfrak{M}_n]_1^G = [\mathfrak{M}_n]_1^Z, \quad \text{Fix}(Z) = \bigcap_{j=k+1}^n \text{Fix}(Z_j).$$

Proof. We first want to show that G has 2^ℓ elements. Suppose that it has more than one element and we have already found a subgroup of G that has 2^i elements f_1, \dots, f_{2^i} . Let g be an element in G that is different from the 2^i elements. Since g is an involution and commutes with each f_j , then

$$f_1, \dots, f_{2^i}, \quad g f_1, \dots, g f_{2^i}$$

form a group of 2^{i+1} elements. We have proved that every finite subgroup of G has exactly 2^ℓ elements. Moreover, if G is infinite then it contains a subgroup of 2^ℓ elements for every $\ell \geq 0$. Let $\{f_1, \dots, f_{2^\ell}\}$ be such a subgroup of G . It suffices to show that $\ell \leq n - k$. We first

linearize all f_j simultaneously. We know that Lf_1, \dots, Lf_{2^ℓ} commute pairwise. Note that $I + f_1'(0)^{-1}f_1$ linearizes f_1 . Assume that f_1 is linear. Then $f_1 = Lf_1$ and Lf_2 commute, and $I + f_2'(0)^{-1}f_2$ commutes with f_1 and linearizes f_2 . Thus f_j can be simultaneously linearized by a holomorphic (resp. formal) change of coordinates. Without loss of generality, we may assume that each f_j is linear. We want to diagonalize all f_j simultaneously. Let E_i^1 and E_i^{-1} be the eigenspaces of f_i with eigenvalues 1 and -1 , respectively. Since $f_i = f_j^{-1}f_i f_j$, each eigenspace of f_i is invariant under f_j . Then we can decompose

$$(2.12) \quad \mathbf{C}^n = \bigoplus_{(i_1, \dots, i_s)} E_1^{i_1} \cap \dots \cap E_s^{i_s}, \quad s = 2^\ell.$$

Here (i_1, \dots, i_s) runs over $\{-1, 1\}^s$ with subspaces $E^{(i_1, \dots, i_s)} := E_1^{i_1} \cap \dots \cap E_s^{i_s}$. On each of these subspaces, $f_j = I$ or $-I$. We are ready to choose a new basis for \mathbf{C}^n whose elements are in the subspaces. Under the new basis, all f_j are diagonal.

Let us rewrite (2.12) as $\mathbf{C}^n = V_1 \oplus V_2 \oplus \dots \oplus V_d$. Here $V_j = E^{I_j}$ and $I_1 = (1, \dots, 1)$. Also, $I_j \neq (1, \dots, 1)$ and $\dim V_j > 0$ for $j > 1$. We have $\dim_{\mathbf{C}} \text{Fix}(G) = \dim_{\mathbf{C}} V_1 = \dim_{\mathbf{C}} [\mathfrak{M}_n]_1^{LG} = k$. Therefore, $d - 1 \leq n - \dim_{\mathbf{C}} V_1 \leq n - k$. We have proved that in suitable coordinates G is contained in the group generated by Z_{k+1}, \dots, Z_n . The remaining assertions follow easily. \square

We will need an elementary result about invariant functions.

Lemma 2.5. *Let Z_{k+1}, \dots, Z_n be defined by (2.10). Let $F = \{f_{k+1}, \dots, f_n\}$ be a family of germs of holomorphic mappings at the origin $0 \in \mathbf{C}^n$. Suppose that the family F is holomorphically equivalent to $\{Z_{k+1}, \dots, Z_n\}$. Let $b_1(z), \dots, b_n(z)$ be germs of holomorphic functions that are invariant under F . Suppose that for $1 \leq j \leq k$, $b_j(0) = 0$ and the linear part of b_j at the origin is \tilde{b}_j . Suppose that for $i > k$, $b_i(z) = O(|z|^2)$ and the quadratic part of b_i at the origin is b_i^* . Suppose that $\tilde{b}_1, \dots, \tilde{b}_k$ are linear independent, and that b_{k+1}^*, \dots, b_n^* are linearly independent modulo $\tilde{b}_1, \dots, \tilde{b}_k$, i.e.*

$$\sum c_i b_i^*(z) = \sum d_j(z) \tilde{b}_j(z) + O(|z|^3)$$

holds for some constants c_i and formal power series d_j , if and only if all c_i are zero. Then invariant functions of F are power series in b_1, \dots, b_n . Furthermore, F is uniquely determined by b_1, \dots, b_n . The same conclusion holds if F and b_j are given by formal power series.

Proof. Without loss of generality, we assume that F is $\{Z_{k+1}, \dots, Z_n\}$. Hence, for all j , there is a formal power series a_j such that $b_j(z) = a_j(z_1, \dots, z_k, z_{k+1}^2, \dots, z_n^2)$. Let us show that the map $w \rightarrow a(w) = (a_1(w), \dots, a_n(w))$ is invertible.

By Lemma 2.4, $\tilde{b}_1(z), \dots, \tilde{b}_k(z)$ are linear combinations of z_1, \dots, z_k , and vice versa. By Lemma 2.4 again, b_{k+1}^*, \dots, b_n^* are linear combinations of z_{k+1}^2, \dots, z_n^2 modulo z_1, \dots, z_k . This shows that

$$b_i^*(z) = \sum_{j>k} c_{ij} z_j^2 + \sum_{\ell \leq k} d_{i\ell}(z) \tilde{b}_\ell(z), \quad i > k.$$

Since b_{k+1}^*, \dots, b_n^* are linearly independent modulo $\tilde{b}_1, \dots, \tilde{b}_k$. Then (c_{ij}) is invertible; so is the linear part of a .

To show that F is uniquely determined by its invariant functions, let \tilde{F} be another such family that is equivalent to $\{Z_{k+1}, \dots, Z_n\}$. Assume that F and \tilde{F} have the same invariant functions. Without loss of generality, assume that \tilde{F} is $\{Z_{k+1}, \dots, Z_n\}$. Then z_1, \dots, z_k are invariant by each F_j , i.e. the i th component of $F_j(z)$ is z_i for $i \leq k$. Also $F_{j,\ell}^2(z) = z_\ell^2$ for $\ell > k$. We get $F_{j,\ell} = \pm z_\ell$. Since z_ℓ is not invariant by \tilde{F} , then it is not invariant by F either. Then $F_{j_\ell,\ell}(z) = -z_\ell$ for some $\ell_j > k$. Since F_{j_ℓ} is equivalent to some Z_i , the set of fixed points of F_{j_ℓ} is a hypersurface. This shows that $F_{j_\ell} = Z_\ell$. So the family F is $\{Z_{k+1}, \dots, Z_n\}$. \square

We now want to find a special set of generators for the deck transformations and its basic properties, which will be important to our study of the normal form problems.

Lemma 2.6. *Let M be defined by (2.1) and (2.2) with $q_* = 0$. Suppose that the group \mathcal{T}_i of deck transformations of $\pi_i: \mathcal{M} \rightarrow \mathbf{C}^p$ has exactly 2^p elements. Then the followings hold.*

- (i) \mathcal{T}_1 is generated by p distinct involutions τ_{1j} such that $\text{Fix}(\tau_{11}), \dots, \text{Fix}(\tau_{1p})$ are hypersurfaces intersecting transversally at 0. And $\tau_1 = \tau_{11} \cdots \tau_{1p}$ is the unique deck transformation of which the set of fixed points has dimension p . Moreover, $\text{Fix}(\tau_1) = \bigcap \text{Fix}(\tau_{1j})$.
- (ii) $\mathcal{O}_n^{\mathcal{T}_1}$ (resp. $\widehat{\mathcal{O}}_n^{\mathcal{T}_1}$) is precisely the set of convergent (resp. formal) power series in z' and $E(z', w')$. $\mathcal{O}_n^{\mathcal{T}_2}$ (resp. $\widehat{\mathcal{O}}_n^{\mathcal{T}_2}$) is the set of convergent (resp. formal) power series in w' and $\overline{E}(w', z')$. In particular, in (z', w') coordinates of \mathcal{M} , \mathcal{T}_1 and \mathcal{T}_2 satisfy

$$(2.13) \quad \begin{aligned} [\mathfrak{M}_n]_1^{L\mathcal{T}_1} \cap [\mathfrak{M}_n]_1^{L\mathcal{T}_2} &= \{0\}, \\ \dim \text{Fix}(\tau_i) &= p, \quad \text{Fix}(\tau_1) \cap \text{Fix}(\tau_2) = \{0\}. \end{aligned}$$

Here $[\mathfrak{M}_n]_1$ is the set of linear functions in z', w' without constant terms.

Proof. (i). Since z_1, \dots, z_p are invariant under deck transformations of π_1 , we have $p' = \dim_{\mathbf{C}}[\mathcal{O}_n]_1^{L\mathcal{T}_1} \geq p$. By Lemma 2.4, π_1 has at most $2^{2p-p'}$ deck transformations. Therefore, $p' = p$. By Lemma 2.4 again, we may assume that in suitable (ξ, η) coordinates, the deck transformations are generated by Z_{p+1}, \dots, Z_{2p} defined by (2.10) in which $z = (\xi, \eta)$. It follows that $Z = Z_{p+1} \cdots Z_{2p}$ is the unique deck transformation of π_1 , of which the set of fixed points has dimension p .

(ii). We have proved that in (ξ, η) coordinates the deck transformations are generated by the above Z_{p+1}, \dots, Z_{2p} . Thus, the invariant holomorphic functions of Z_{p+1}, \dots, Z_{2p} are precisely the holomorphic functions in $\xi_1, \dots, \xi_p, \eta_1^2, \dots, \eta_p^2$. Since z_1, \dots, z_p and $E_i(z', w')$ are invariant under deck transformations, then on \mathcal{M}

$$(2.14) \quad z' = f(\xi, \eta_1^2, \dots, \eta_p^2), \quad E(z', w') = g(\xi, \eta_1^2, \dots, \eta_p^2).$$

Since (z', w') are local coordinates of \mathcal{M} , the differentials of z_1, \dots, z_p under any coordinate system of \mathcal{M} are linearly independent. Computing the differentials of z' in variables ξ, η by using (2.14), we see that the mapping $\xi \rightarrow f(\xi, 0)$ is a local biholomorphism. Expressing both sides of the second identity in (2.14) as power series in ξ, η , we obtain

$$E(f(\xi, 0), w') = g(\xi, \eta_1^2, \dots, \eta_p^2) + O(|(\xi, \eta)|^3).$$

We set $\xi = 0$, compute the left-hand side, and rewrite the identity as

$$(2.15) \quad g(0, \eta_1^2, \dots, \eta_p^2) = q(w') + O(|(\xi, \eta)|^3).$$

As coordinate systems, (z', w') and (ξ, η) vanish at $0 \in \mathcal{M}$. We now use $(z', w') = O(|(\xi, \eta)|)$. By (2.14), $f(0) = g(0) = 0$ and $g(\xi, 0) = O(|\xi|^2)$. Let us verify that the linear parts of $g_1(0, \eta), \dots, g_p(0, \eta)$ are linearly independent. Suppose that $\sum_{j=1}^p c_j g_j(0, \eta) = O(|\eta|^2)$. Replacing ξ, η by $O(|(z', w')|)$ in (2.15) and setting $z' = 0$, we obtain

$$\sum_{j=1}^p c_j q_j(w') = O(|w'|^3), \quad i.e. \quad \sum_{j=1}^p c_j q_j(w') = 0.$$

As remarked after condition B was introduced, $q_* = 0$ implies that $q_1(w'), \dots, q_p(w')$ are linearly independent. Thus all c_j are 0. We have verified that $\xi \rightarrow f(\xi, 0)$ is biholomorphic near $\xi = 0$. Also $\eta \rightarrow g(0, \eta)$ is biholomorphic near $\eta = 0$ and $g(\xi, 0) = O(|\xi|^2)$. Therefore, $(\xi, \eta) \rightarrow (f, g)(\xi, \eta)$ is invertible near 0. By solving (2.14), the functions $\xi, \eta_1^2, \dots, \eta_p^2$ are expressed as power series in z' and $E(z', w')$.

It is clear that z_1, \dots, z_p are invariant under τ_{1j} . From linearization of \mathcal{T}_1 , we know that the space of invariant linear functions of $L\mathcal{T}_1$ is the same as the space of linear invariant functions of $L\tau_1$, which has dimension p . This shows that z_1, \dots, z_p span the space of linear invariant functions of $L\tau_1$. Also w_1, \dots, w_p span the space of linear invariant functions of $L\tau_2$. We obtain $[\mathfrak{M}_n]_1^{L\mathcal{T}_1} \cap [\mathfrak{M}_n]_1^{L\mathcal{T}_2} = \{0\}$. We have verified (2.13).

In view of the linearization of \mathcal{T}_1 in (i), we obtain $\dim \text{Fix}(\tau_1) = \dim \text{Fix}(\mathcal{T}_1) = p$. Moreover, $\text{Fix}(\tau_i)$ is a smooth submanifold of which the tangent space at the origin is $\text{Fix}(L\tau_i)$. We choose a basis u_1, \dots, u_p for $\text{Fix}(L\tau_1)$. Let v_1, \dots, v_p be any p vectors such that $u_1, \dots, u_p, v_1, \dots, v_p$ form a basis of \mathbf{C}^n . In new coordinates defined by $\sum \xi_i u_i + \eta_i v_i$, we know that linear invariant functions of $L\tau_1$ are spanned by ξ_1, \dots, ξ_p . The linear invariant functions in (ξ, η) that are invariant by $L\tau_2$ are spanned by $f_j(\xi, \eta) = \sum_k (a_{jk} \xi_k + b_{jk} \eta_k)$ for $1 \leq j \leq p$. Since $[\mathfrak{M}_n]^{L\tau_1} \cap [\mathfrak{M}_n]^{L\tau_2} = \{0\}$, then $\xi_1, \dots, \xi_p, f_1, \dots, f_p$ are linearly independent. Equivalently, (b_{jk}) is non-singular. Now $\text{Fix}(L\tau_2)$ is spanned by vectors $\sum_k (a_{jk} u_k + b_{jk} v_k)$. This shows that $\text{Fix}(L\tau_1) \cap \text{Fix}(L\tau_2) = \{0\}$. Therefore, $\text{Fix}(\tau_1)$ intersects $\text{Fix}(\tau_2)$ transversally at the origin and the intersection must be the origin. \square

Note that the proof of the above lemma actually gives us a more general result.

Corollary 2.7. *Let \mathfrak{J} be a group of commuting holomorphic (formal) involutions on \mathbf{C}^n .*

(i) *$\text{Fix}(L\mathfrak{J}) = \{0\}$ if and only if $[\mathfrak{M}_n]_1^{L\mathfrak{J}}$ has dimension 0.*

(ii) *Let $\tilde{\mathfrak{J}}$ be another family of commuting holomorphic (resp. formal) involutions such that $[\mathfrak{M}_n]^{L\mathfrak{J}} \cap [\mathfrak{M}_n]^{L\tilde{\mathfrak{J}}} = \{0\}$. Then $\text{Fix}(L\mathfrak{J}) \cap \text{Fix}(L\tilde{\mathfrak{J}}) = \{0\}$. Moreover, $\text{Fix}(\mathfrak{J}) \cap \text{Fix}(\tilde{\mathfrak{J}}) = \{0\}$ if \mathfrak{J} and $\tilde{\mathfrak{J}}$ consist of convergent involutions.*

In view of Lemma 2.6, we will refer to the family

$$\mathcal{T}_1 := \{\tau_{11}, \dots, \tau_{1p}\}, \quad \mathcal{T}_2 := \{\tau_{21}, \dots, \tau_{2p}\}, \quad \rho$$

as the Moser-Webster involutions, where $\tau_{2j} = \rho \tau_{1j} \rho$. The significance of the two sets of involutions is the following proposition that transforms the normalization of the real manifolds into that of two families of commuting involutions.

For clarity, recall the anti-holomorphic involution $\rho_0: (z', w') \rightarrow (\overline{w'}, \overline{z'})$.

Proposition 2.8. *Let M and \widetilde{M} be two real analytic submanifolds of the form (2.1) and (2.2) that admit Moser-Webster involutions $\{\mathcal{T}_1, \rho_0\}$ and $\{\widetilde{\mathcal{T}}_1, \rho_0\}$, respectively. Then M and \widetilde{M} are holomorphically equivalent if and only if $\{\mathcal{T}_1, \rho_0\}$ and $\{\widetilde{\mathcal{T}}_1, \rho_0\}$ are holomorphically equivalent, i.e. if there is a biholomorphic map f commuting with ρ_0 such that $f\tau_{1j}f^{-1} = \widetilde{\tau}_{1i_j}$ for $1 \leq j \leq p$. Here $\{i_1, \dots, i_p\} = \{1, \dots, p\}$.*

Let $\mathcal{T}_1 = \{\tau_{11}, \dots, \tau_{1p}\}$ be a family of p distinct commuting holomorphic involutions. Suppose that $\text{Fix}(\tau_{11}), \dots, \text{Fix}(\tau_{1p})$ are hypersurfaces intersecting transversely at the origin. Let ρ be an anti-holomorphic involutions and let \mathcal{T}_2 be the family of involutions $\tau_{2j} = \rho\tau_{1j}\rho$ with $1 \leq j \leq p$. Suppose that

$$(2.16) \quad [\mathfrak{M}_n]_1^{L\mathcal{T}_1} \cap [\mathfrak{M}_n]_1^{L\mathcal{T}_2} = \{0\}.$$

There exists a real analytic real n -submanifold

$$(2.17) \quad M \subset \mathbf{C}^{2p}: z_{p+j} = A_j^2(z', \bar{z}'), \quad 1 \leq j \leq p$$

such that the set of Moser-Webster involutions $\{\widetilde{\mathcal{T}}_1, \rho_0\}$ of M is holomorphically equivalent to $\{\mathcal{T}_1, \rho\}$.

Proof. We recall from (2.5) the branched covering

$$\pi_1: \mathcal{M}_1 := \mathcal{M} \cap ((\Delta_\delta^p \times \Delta_{\delta^2}^p) \times (\Delta_{C\delta}^p \times \Delta_{C\delta^2}^p)) \longrightarrow \Delta_\delta^p \times \Delta_{\delta^2}^p.$$

Here $C \geq 1$. Let π_1 be restricted to \mathcal{M}_1 . Then $\pi_2 = \overline{\pi_1} \circ \rho$ is defined on $\rho(\mathcal{M}_1)$. Note that

$$\pi_2: \rho(\mathcal{M}_1) \longrightarrow \Delta_\delta^p \times \Delta_{\delta^2}^p.$$

We have $\pi_1^{-1}(z) \cap \text{Fix}(\rho) = \{(z, \bar{z})\}$ for $z \in M$ and $\pi_1(\text{Fix}(\rho)) = M$. Let $\mathcal{B}_0 \subset \Delta_\delta^p \times \Delta_{\delta^2}^p$ be the branched locus. Take $\mathcal{B} = \pi_1^{-1}(\mathcal{B}_0)$. We will denote by $\widetilde{\mathcal{M}}_1, \widetilde{\mathcal{B}}$ and $\widetilde{\mathcal{B}}_0$ the corresponding data for \widetilde{M} . Here $\widetilde{\mathcal{M}}_1$ is an analogous branched covering over $\pi_1(\widetilde{\mathcal{M}}_1)$. We assume that the latter contains $f(\pi_1(\mathcal{M}_1))$ if \widetilde{M} is equivalent to M via f .

Assume that f is a biholomorphic map sending M into \widetilde{M} . Let f^c be the restriction of biholomorphic map $f^c(z, w) = (f(z), \bar{f}(w))$ to \mathcal{M} . Let M be defined by $z'' = E(z', \bar{z}')$ and \widetilde{M} be defined by $z'' = \widetilde{E}(z', \bar{z}')$. By $f(M) \subset \widetilde{M}$, $f = (f', f'')$ satisfies

$$f''(z', E(z', \bar{z}')) = \widetilde{E}(f'(z', E(z', \bar{z}')), \overline{f'}(\bar{z}', \overline{E}(\bar{z}', z'))).$$

Using the defining equations for \mathcal{M} , we get $f^c(\mathcal{M}) \subset \widetilde{\mathcal{M}}$ and $\rho f^c = f^c \rho$ on $\mathcal{M} \cap \rho(\mathcal{M})$. We will also assume that $f^c(\mathcal{M}_1)$ is contained in $\widetilde{\mathcal{M}}_1$. It is clear that f^c sends a fiber $\pi_1^{-1}(z)$ onto the fiber $\pi_1^{-1}(f(z))$ for $z \in \Omega = \pi_1(\mathcal{M}_1) \setminus (\mathcal{B}_0 \cup f^{-1}(\widetilde{\mathcal{B}}_0))$, since the two fibers have the same number of points and f is injective. Thus $f^c\tau_{1j} = \widetilde{\tau}_{1i_j}f^c$ on $\pi_1^{-1}(\Omega)$. Here i_j is of course locally determined on $\pi_1^{-1}(\Omega)$. Since \mathcal{B} has positive codimension in \mathcal{M}_1 then $\mathcal{M}_1 \setminus \mathcal{B}$ is connected. Hence i_j is well-defined on $\pi_1^{-1}(\Omega)$. Then $f^c\tau_{1j} = \widetilde{\tau}_{1i_j}f^c$ on $\mathcal{M}_1 \setminus \mathcal{B}$. This shows that f^c conjugates simultaneously the deck transformations of \mathcal{M} to the deck transformations of $\widetilde{\mathcal{M}}$ for π_1 . The same conclusion holds for π_2 .

Conversely, assume that there is a biholomorphic map $g: \mathcal{M} \rightarrow \widetilde{\mathcal{M}}$ such that $\rho g = g \rho$ and $g\tau_{1i} = \widetilde{\tau}_{1j_i}g$. Since $\tau_{11}, \dots, \tau_{12p}$ are distinct and $\mathcal{M}_1 \setminus \mathcal{B}$ is connected, then $\bigcup_{j \neq i} \{x \in \mathcal{M}_1 \setminus \mathcal{B}: \tau_{1i}(x) = \tau_{1j}(x)\}$ is a complex subvariety of positive codimension in $\mathcal{M}_1 \setminus \mathcal{B}$. Its image under the proper projection π_1 is a subvariety of positive codimension in $\Delta_\delta^p \times \Delta_{\delta^2}^p \setminus \mathcal{B}_0$. This

shows that the latter contains a non-empty open subset ω such that $\{\tau_{11}(x), \dots, \tau_{12^p}(x)\} = \pi_1^{-1}\pi_1(x)$ has 2^p distinct points for each $\pi_1(x) \in \omega$. Therefore, $\tau_{11}, \dots, \tau_{12^p}$ are all deck transformations of π_1 over ω . Hence they are all deck transformations of $\pi_1: \mathcal{M}_1 \setminus \mathcal{B} \rightarrow \Delta_\delta^p \times \Delta_{\delta^2}^p \setminus \mathcal{B}_0$, too. This shows that $\pi_1^{-1}(\pi_1(x)) = \{\tau_{1j}(x): 1 \leq j \leq 2^p\}$ for $x \in \mathcal{M}_1 \setminus \mathcal{B}$. Now, g sends $\tau_{1j}(x)$ to $\tilde{\tau}_{1i_j}(g(x))$ for each j . Hence $f(z) = \pi_1 g \pi_1^{-1}(z)$ is well-defined and holomorphic for $z \in \Delta_\delta^p \times \Delta_{\delta^2}^p \setminus \mathcal{B}_0$. By the Riemann extension for bounded holomorphic functions, f extends to a holomorphic mapping, still denoted by f , which is defined near the origin. We know that f is invertible and in fact the inverse can be obtained by extending the mapping $z \rightarrow \pi_1 g^{-1} \pi_1(z)$. If $z = (z', E(z', w')) \in M$, then $w' = \bar{z}'$ and $f(z) = \pi_1 g \pi_1^{-1}(z) = \pi_1 g(z, \bar{z})$ with $(z, \bar{z}) \in \text{Fix}(\rho)$. Since $\rho g = g \rho$, then $g(z, \bar{z}) \in \text{Fix}(\rho)$. Thus $f(z) = \pi_1 g(z, \bar{z}) \in \widetilde{M}$.

Assume that $\{\tau_{1j}\}$ and ρ are germs of involutions defined at the origin of \mathbf{C}^n . Assume that they satisfy the conditions in the proposition. From Lemma 2.4 it follows that $\tau_{11}, \dots, \tau_{1p}$ generate a group of 2^p involutions, while the p generators are the only elements of which each fixes a hypersurface pointwise. To realize them as deck transformations of the complexification of a real analytic submanifold, we apply Lemma 2.4 to find a coordinate map $(\xi, \eta) \rightarrow \phi(\xi, \eta) = (A, B)(\xi, \eta)$ such that invariant holomorphic functions of $\{\tau_{1j}\}$ are precisely holomorphic functions in

$$z' = (A_1(\xi, \eta), \dots, A_p(\xi, \eta)), \quad z'' = (B_1^2(\xi, \eta), \dots, B_p^2(\xi, \eta)).$$

Note that B_j is skew-invariant under τ_{1j} and is invariant under τ_{1i} for $i \neq j$ and A is invariant under all τ_{1j} . Set

$$w'_j = \overline{A_j \circ \rho(\xi, \eta)}, \quad w''_j = \overline{B_j^2 \circ \rho(\xi, \eta)}.$$

Since $\tau_{2j} = \rho \tau_{1j} \rho$, the holomorphic functions invariant under all τ_{2j} are precisely the holomorphic functions in the above w', w'' . We now draw conclusions for the linear parts of invariant functions and involutions. Since ϕ is biholomorphic, then LA_1, \dots, LA_p are linearly independent. They are also invariant under $L\tau_{1j}$. Since $\tau_{2j} = \rho \tau_{1j} \rho$, the p functions $\overline{LA_i \circ \rho}$ are linearly independent and invariant under $L\tau_{2j}$. Thus

$$LA_1, \dots, LA_p, \overline{LA_1 \circ \rho}, \dots, \overline{LA_p \circ \rho}$$

are linearly independent, since $[\mathfrak{M}_n]_1^{L\tau_1} \cap [\mathfrak{M}_n]_1^{L\tau_2} = \{0\}$. This shows that the map $(\xi, \eta) \rightarrow (z', w') = (A(\xi, \eta), \overline{A \circ \rho(\xi, \eta)})$ has an inverse $(\xi, \eta) = \psi(z', w')$. Define

$$M: z'' = (B_1^2, \dots, B_p^2) \circ \psi(z', \bar{z}').$$

The complexification of M is given by

$$\mathcal{M}: z'' = (B_1^2, \dots, B_p^2) \circ \psi(z', w'), \quad w'' = (\overline{B_1^2}, \dots, \overline{B_p^2}) \circ \bar{\psi}(w', z').$$

Note that $\phi \circ \psi(z', w') = (z', B \circ \psi(z', w'))$ is biholomorphic. In particular, we can write

$$B_j^2 \circ \psi(z', \bar{z}') = h_j(z', \bar{z}') + q_j(\bar{z}') + b_j(z') + O(|(z', \bar{z}')|^3).$$

Here $q_j(\bar{z}') = \tilde{q}_j^2(\bar{z}')$, and $\tilde{q}(w')$ is the linear part of $w' \rightarrow B \circ \psi(0, w')$. Therefore, $|q(w')| \geq c|w'|^2$ and $q_* = 0$. By Lemma 2.1, $\pi_1: \mathcal{M} \rightarrow \mathbf{C}^p$ is a 2^p -to-1 branched covering defined near $0 \in \mathcal{M}$. Since B^2 is invariant by τ_{1j} , then $z'' = B^2 \circ \psi(z', w')$ is invariant by $\psi^{-1} \tau_{1j} \psi(z', w')$. Also A is invariant under τ_{1j} . Then $z' = A \circ \psi(z', w')$ is invariant by $\psi^{-1} \tau_{1j} \psi(z', w')$. This

show that $\{\psi^{-1}\tau_{1j}\psi\}$ has the same invariant functions as of the deck transformations of π_1 . By Lemma 2.5, $\{\psi^{-1}\tau_{1j}\psi\}$ agrees with the set of deck transformations of π_1 . For $\rho_0(z', w') = (\overline{w'}, \overline{z'})$ we have $\rho_0\psi^{-1} = \psi^{-1}\rho$. This shows that M is a realization for $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$. \square

Remark 2.9. We choose the realization in such a way that z_{p+j} are square functions. This particular holomorphic equivalent form of M will be crucial to study the asymptotic manifolds in section 7.

Next we want to compute the deck transformations for a product quadric. We will first recall the Moser-Webster involutions for elliptic and hyperbolic complex tangents. We will then compute the deck transformations for complex tangents of complex type.

The Moser-Webster theory deals with the case $p = 1$ for a real analytic surface

$$z_2 = |z_1|^2 + \gamma(z_1^2 + \overline{z}_1^2) + O(|z_1|^3), \quad \text{or} \quad z_2 = z_1^2 + \overline{z}_1^2 + O(|z_1|^3).$$

Here $\gamma \geq 0$ is the Bishop invariant of M . One of most important properties of the Moser-Webster theory is the existence of the above mentioned deck transformations. When $\gamma \neq 0$, there is a pair of Moser-Webster involutions τ_1, τ_2 with $\tau_2 = \rho\tau_1\rho$ such that τ_1 generates the deck transformations of π_1 . In fact, τ_1 is the only possible non-trivial deck transformation of π_1 . When $\gamma \neq 1/2$, in suitable coordinates their composition $\sigma = \tau_1\tau_2$ is of the form

$$\tau: \xi' = \mu\xi + O(|(\xi, \eta)|^2), \quad \eta' = \mu^{-1}\eta + O(|(\xi, \eta)|^2).$$

Here $\rho(\xi, \eta) = (\overline{\eta}, \overline{\xi})$ when $0 < \gamma < 1/2$, and $\rho(\xi, \eta) = (\overline{\xi}, \overline{\eta})$ when $\gamma > 1/2$. When the complex tangent is elliptic and $0 < \gamma < 1/2$, σ is *hyperbolic* with $|\mu| > 1$; when the complex tangent is hyperbolic, i.e. $1/2 < \gamma \leq \infty$, then σ is *elliptic* with $|\mu| = 1$. When the complex tangent is parabolic, the linear part of σ is not diagonalizable and 1 is the eigenvalue.

We will see later that with $p \geq 2$, there is yet another simple model that is not in the product. This is the quadric in \mathbf{C}^4 defined by

$$(2.18) \quad Q_{\gamma_s}: z_3 = z_1\overline{z}_2 + \gamma_s\overline{z}_2^2 + (1 - \gamma_s)z_1^2, \quad z_4 = \overline{z}_3.$$

Here γ_s is a complex number. We will, however, exclude $\gamma_s = 0$ or equivalently $\gamma_s = 1$ by condition B. We also exclude $\gamma_s = 1/2$ by condition J. Note that $\gamma_s = 1/2$ does not correspond to a product Bishop quadrics either, by examining the CR singular sets. Under these mild non degeneracy conditions, we will show that γ_s is an invariant when it is normalized to the range

$$(2.19) \quad \text{Re } \gamma_s \leq 1/2, \quad \text{Im } \gamma_s \geq 0, \quad \gamma_s \neq 0.$$

In this case, the complex tangent is said of *complex* type.

We have introduced the types of the complex tangent at the origin. Of course a product of quadrics, or a product quadric, can exhibit a combination of the above basic 4 types. We will see soon that quadrics have other invariants when $p > 1$. Nevertheless, in our results, the above invariants that describe the types of the complex tangent will play a major role in the convergence or divergence of normalizations.

Let us first recall involutions in [MW83] where the complex tangents are elliptic (with non-vanishing Bishop invariant) or hyperbolic. When $\gamma > 0$, the non-trivial deck transformations of

$$Q_\gamma: z_2 = |z_1|^2 + \gamma(z_1^2 + \overline{z}_1^2)$$

for π_1, π_2 are τ_1, τ_2 , respectively. They are

$$(2.20) \quad \tau_1: z'_1 = z_1, \quad w'_1 = -w_1 - \gamma^{-1}z_1; \quad \tau_2 = \rho\tau_1\rho$$

with ρ being defined by (2.4). Here the formula is valid for $\gamma = \infty$ (i.e. $\gamma^{-1} = 0$). Note that τ_1 and τ_2 do not commute and $\sigma = \tau_1\tau_2$ satisfies

$$\sigma^{-1} = \tau_i\sigma\tau_i = \rho\sigma\rho, \quad \tau_i^2 = I, \quad \rho^2 = I.$$

When the complex tangent is not parabolic, the eigenvalues of σ are μ, μ^{-1} with $\mu = \lambda^2$ and $\gamma\lambda^2 - \lambda + \gamma = 0$. For the elliptic complex tangent, we can choose a solution $\lambda > 1$, and in suitable coordinates we obtain

$$(2.21) \quad \begin{aligned} \tau_1: \xi' &= \lambda\eta + O(|(\xi, \eta)|^2), & \eta' &= \lambda^{-1}\xi + O(|(\xi, \eta)|^2), \\ \tau_2 &= \rho\tau_1\rho, & \rho(\xi, \eta) &= (\bar{\eta}, \bar{\xi}), \\ \sigma: \xi' &= \mu\xi + O(|(\xi, \eta)|^2), & \eta' &= \mu^{-1}\eta + O(|(\xi, \eta)|^2), & \mu &= \lambda^2. \end{aligned}$$

When the complex tangent is hyperbolic, i.e. $1/2 < \gamma \leq \infty$, τ_i and σ still have the above form, while $|\mu| = 1 = |\lambda|$ and

$$\rho(\xi, \eta) = (\bar{\xi}, \bar{\eta}).$$

When the complex tangent is parabolic, i.e. $\gamma = 1/2$, the pair of involutions still exists. However, $L\sigma$ is not diagonalizable and 1 is its only eigenvalue. We recall from [MW83] that

$$(2.22) \quad \gamma = \frac{1}{\lambda + \lambda^{-1}}.$$

For the complex type, new situations arise. Such a quadric has the form

$$(2.23) \quad Q_{\gamma_s}: z_3 = z_1\bar{z}_2 + \gamma_s\bar{z}_2^2 + (1 - \gamma_s)z_1^2, \quad z_4 = \bar{z}_3.$$

Here γ_s is a complex number. Let us first check that such a quadric is not the product of two Bishop quadrics : Its CR singular set is defined by

$$(z_1 + 2\gamma_s\bar{z}_2)(z_2 + 2(1 - \bar{\gamma}_s)\bar{z}_1) = 0,$$

which is the union of a complex line and a real surface when $\gamma_s = 0, 1$, or a union of two totally real surfaces. The CR singular set of a quadric defined by $z_3 = |z_1|^2 + \gamma_1(z_1^2 + \bar{z}_1^2)$ and $z_4 = |z_2|^2 + \gamma_2(z_2^2 + \bar{z}_2^2)$ is given by

$$(z_1 + 2\gamma_1\bar{z}_1)(z_2 + 2\gamma_2\bar{z}_2) = 0.$$

It is the union of two Bishop surfaces when $\gamma_1 \neq 1/2$ and $\gamma_2 \neq 1/2$, and it contains a submanifold of dimension 3 otherwise.

By condition B, we know that $\gamma_s \neq 0, 1$. Let us compute the deck transformations of the complexification of (2.23). According to Lemma 2.6 (i), the deck transformations for π_1 are generated by two involutions

$$\tau_{11}: \begin{cases} z'_1 = z_1, \\ z'_2 = z_2, \\ w'_1 = -w_1 - (1 - \bar{\gamma}_s)^{-1}z_2, \\ w'_2 = w_2; \end{cases} \quad \tau_{12}: \begin{cases} z'_1 = z_1, \\ z'_2 = z_2, \\ w'_1 = w_1, \\ w'_2 = -w_2 - \gamma_s^{-1}z_1. \end{cases}$$

We still have ρ defined by (2.4). Let $\tau_{2j} = \rho\tau_{1j}\rho$. Then τ_{21}, τ_{22} generate the deck transformations of π_2 . Note that

$$\tau_{21}: \begin{cases} z'_1 = -z_1 - (1 - \gamma_s)^{-1}w_2, \\ z'_2 = z_2, \\ w'_1 = w_1, \\ w'_2 = w_2; \end{cases} \quad \tau_{22}: \begin{cases} z'_1 = z_1, \\ z'_2 = -z_2 - \bar{\gamma}_s^{-1}w_1, \\ w'_1 = w_1, \\ w'_2 = w_2. \end{cases}$$

Recall that $\tau_i = \tau_{i1}\tau_{i2}$ is the unique deck transformation of π_i that has the smallest dimension of the set of fixed-points among all deck transformations. They are

$$(2.24) \quad \tau_1: \begin{cases} z'_1 = z_1, \\ z'_2 = z_2, \\ w'_1 = -w_1 - (1 - \bar{\gamma}_s)^{-1}z_2, \\ w'_2 = -w_2 - \gamma_s^{-1}z_1; \end{cases} \quad \tau_2: \begin{cases} z'_1 = -z_1 - (1 - \gamma_s)^{-1}w_2, \\ z'_2 = -z_2 - \bar{\gamma}_s^{-1}w_1, \\ w'_1 = w_1, \\ w'_2 = w_2. \end{cases}$$

And $\tau_1\tau_2$ is given by

$$\sigma_s: \begin{cases} z'_1 = -z_1 - (1 - \gamma_s)^{-1}w_2, \\ z'_2 = -z_2 - \bar{\gamma}_s^{-1}w_1, \\ w'_1 = (1 - \bar{\gamma}_s)^{-1}z_2 + ((\bar{\gamma}_s - \bar{\gamma}_s^2)^{-1} - 1)w_1, \\ w'_2 = \gamma_s^{-1}z_1 + ((\gamma_s - \gamma_s^2)^{-1} - 1)w_2. \end{cases}$$

In contrast to the elliptic and hyperbolic cases, τ_{11} and $\rho\tau_{11}\rho$ commute; in other words, $\tau_{11}\rho\tau_{11}\rho$ is actually an involution. And τ_{12} and $\rho\tau_{12}\rho$ commute, too. However, τ_{11} and τ_{22} do not commute, and τ_{12}, τ_{21} do not commute either. Thus, we form compositions

$$\sigma_{s1} = \tau_{11}\tau_{22}, \quad \sigma_{s2} = \tau_{12}\tau_{21}, \quad \sigma_{s2}^{-1} = \rho\sigma_{s1}\rho.$$

By a simple computation, we have

$$\sigma_{s1}: \begin{cases} z'_1 = z_1, \\ z'_2 = -z_2 - \bar{\gamma}_s^{-1}w_1, \\ w'_1 = (1 - \bar{\gamma}_s)^{-1}z_2 + ((\bar{\gamma}_s - \bar{\gamma}_s^2)^{-1} - 1)w_1, \\ w'_2 = w_2; \end{cases}$$

$$\sigma_{s2}: \begin{cases} z'_1 = -z_1 - (1 - \gamma_s)^{-1}w_2, \\ z'_2 = z_2, \\ w'_1 = w_1, \\ w'_2 = \gamma_s^{-1}z_1 + ((\gamma_s - \gamma_s^2)^{-1} - 1)w_2. \end{cases}$$

We verify that $\sigma_{s1}\sigma_{s2} = \sigma_s = \tau_1\tau_2$. This allows us to compute the eigenvalues of $\sigma_{s1}\sigma_{s2}$ easily:

$$(2.25) \quad \mu_s, \quad \mu_s^{-1}, \quad \bar{\mu}_s^{-1}, \quad \bar{\mu}_s,$$

$$(2.26) \quad \mu_s = \bar{\gamma}_s^{-1} - 1.$$

In fact we compute them by observing that the first two in (2.25) and 1 with multiplicity are eigenvalues of σ_{s1} , while the last two in (2.25) and 1 with multiplicity are eigenvalues of

σ_{s2} . Note that σ_s is diagonalizable if and only if $\gamma_s \neq 1/2$. When $\gamma_s < 1/2$ and $\gamma_s \neq 0$, μ_s and μ_s^{-1} are eigenvalues of multiplicity two while condition J holds. Therefore, for $\gamma_s \neq 1/2$, i.e. $\mu_s \neq 1$, we can find a linear transformation of the form

$$\psi: (z_1, w_2) \rightarrow (\xi_2, \eta_2) = \bar{\phi}(z_1, w_2), \quad (z_2, w_1) \rightarrow (\xi_1, \eta_1) = \phi(w_1, z_2)$$

such that $\sigma_{s1}, \sigma_{s2}, \sigma_s = \sigma_{s1}\sigma_{s2}$ are simultaneously diagonalized as

$$(2.27) \quad \begin{array}{l} \sigma_{s1}: \quad \xi'_1 = \mu_s \xi_1, \quad \eta'_1 = \mu_s^{-1} \eta_1, \quad \xi'_2 = \xi_2, \quad \eta'_2 = \eta_2, \\ \sigma_{s2}: \quad \xi'_1 = \xi_1, \quad \eta'_1 = \eta_1, \quad \xi'_2 = \bar{\mu}_s^{-1} \xi_2, \quad \eta'_2 = \bar{\mu}_s \eta_2, \\ \sigma_s: \quad \xi'_1 = \mu_s \xi_1, \quad \eta'_1 = \mu_s^{-1} \eta_1, \quad \xi'_2 = \bar{\mu}_s^{-1} \xi_2, \quad \eta'_2 = \bar{\mu}_s \eta_2. \end{array}$$

Under the transformation ψ , the involution ρ , defined by (2.4), takes the form

$$(2.28) \quad \rho(\xi_1, \xi_2, \eta_1, \eta_2) = (\bar{\xi}_2, \bar{\xi}_1, \bar{\eta}_2, \bar{\eta}_1).$$

Moreover, for $j = 1, 2$, we have $\tau_{2j} = \rho\tau_{1j}\rho$ and

$$(2.29) \quad \tau_{1j}: \quad \xi'_j = \lambda_j \eta_j, \quad \eta'_j = \lambda_j^{-1} \xi_j; \quad \xi'_i = \xi_i, \quad \eta'_i = \eta_i, \quad i \neq j;$$

$$(2.30) \quad \lambda_1 = \lambda_s, \quad \lambda_2 = \bar{\lambda}_s^{-1}, \quad \mu_s = \lambda_s^2.$$

Condition J on Q_{γ_s} is equivalent to $\gamma_s \neq 1/2$. By a permutation of coordinates that preserves ρ , we obtain a unique holomorphic invariant μ_s satisfying

$$(2.31) \quad |\mu_s| \geq 1, \quad \text{Im } \mu_s \geq 0, \quad \mu_s \neq 0, -1, \quad 0 \leq \arg \lambda_s \leq \pi/2, \quad \lambda_s \neq 0, i.$$

Note that $\text{Re } \gamma_s < 1/2$ implies that $|\mu_s| \neq 1$.

For later purpose, we summarize some facts for complex type in the following.

Proposition 2.10. *Let $Q_{\gamma_s} \subset \mathbf{C}^4$ be the quadric defined by (2.18) and (2.19) with $\gamma_s \neq 0, 1$. Then π_1 admits two deck transformations τ_{11}, τ_{12} such that the set of fixed points of each τ_{1j} has dimension 3. Also, $\tau_{2j} = \rho\tau_{1j}\rho$ are the deck transformations of π_2 and*

$$\tau_{11}\tau_{21} = \tau_{21}\tau_{11}, \quad \tau_{12}\tau_{22} = \tau_{22}\tau_{12}.$$

Let $\sigma_{s1} = \tau_{11}\tau_{22}$, $\sigma_{s2} = \tau_{12}\tau_{21}$, $\tau_i = \tau_{i1}\tau_{i2}$, and $\sigma_s = \tau_1\tau_2$. Then

$$\sigma = \sigma_{s1}\sigma_{s2} = \sigma_{s2}\sigma_{s1}, \quad \sigma_{s2}^{-1} = \rho\sigma_{s1}\rho, \quad \sigma_s^{-1} = \rho\sigma_s\rho.$$

Assume further that $\gamma_s \neq 1/2$. In suitable coordinates $\sigma_{s1}, \sigma_{s2}, \sigma, \rho_s$ are given by (2.27)-(2.28), while σ_s has 4 (possibly repeated) eigenvalues given by (2.25)-(2.26). The σ has four distinct eigenvalues if and only if Q_{γ_s} can be holomorphically normalized so that $\text{Im } \gamma_s > 0$ and $\text{Re } \gamma_s < 1/2$.

In summary, we have found linear coordinates for the product quadrics such that the normal forms of S, T_{ij}, ρ of the corresponding $\sigma, \sigma_j, \tau_{ij}, \rho_0$ are given by

$$(2.32) \quad S: \quad \xi'_j = \mu_j \xi_j, \quad \eta'_j = \mu_j^{-1} \eta_j;$$

$$(2.33) \quad S_j: \quad \xi'_j = \mu_j \xi_j, \quad \eta'_j = \mu_j^{-1} \eta_j, \quad \xi'_k = \xi_k, \quad \eta'_k = \eta_k, \quad k \neq j;$$

$$(2.34) \quad T_{ij}: \quad \xi'_j = \lambda_{ij} \eta_j, \quad \eta'_j = \lambda_{ij}^{-1} \xi_j, \quad \xi'_k = \xi_k, \quad \eta'_k = \eta_k, \quad k \neq j;$$

$$(2.35) \quad \rho: \quad \begin{cases} (\xi'_e, \eta'_e, \xi'_h, \eta'_h) = (\bar{\eta}_e, \bar{\xi}_e, \bar{\xi}_h, \bar{\eta}_h), \\ (\xi'_s, \xi'_{s+s^*}, \eta'_s, \eta'_{s+s^*}) = (\bar{\xi}_{s+s^*}, \bar{\xi}_s, \bar{\eta}_{s+s^*}, \bar{\eta}_s). \end{cases}$$

Throughout the paper, the indices h, e, s have the ranges $1 \leq e \leq e_*$, $e_* < h \leq e_* + h_*$, and $e_* + h_* < s \leq p - s_*$. Notice that we can always normalize ρ_0 into the above normal form ρ .

The commutativity of $\sigma_h, \sigma_e, \sigma_{s1}, \sigma_{s2}$ will be important to understand the convergence of normalization for the abelian CR singularity.

3. FORMAL DECK TRANSFORMATIONS AND CENTRALIZERS

In section 2 we show the equivalence of the classification of real analytic submanifolds M that admit the maximum number of deck transformations and the classification of the families of involutions $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$ that satisfy some mild conditions (see Proposition 2.8). As a consequence we show that a real submanifold is formally equivalent to a quadric if and only if its family of Moser-Webster involutions is formally linearizable.

3.1. Formal submanifolds. We first need some notation. Let I be an ideal of the ring $\mathbf{R}[[x]]$ of formal power series in $x = (x_1, \dots, x_N)$. Since $\mathbf{R}[[x]]$ is noetherian, then I and its radical \sqrt{I} are finitely generated. We say that I defines a formal submanifold M of dimension $N - k$ if \sqrt{I} is generated by r_1, \dots, r_k such that at the origin all r_j vanish and dr_1, \dots, dr_k are linearly independent. For such an M , let $I(M)$ denote \sqrt{I} and let T_0M be defined by $dr_1(0) = \dots = dr_k(0) = 0$. If $F = (f_1, \dots, f_N)$ is a formal mapping with $f_j \in \mathbf{R}[[x]]$, we say that its set of (formal) fixed points is a submanifold if the ideal generated by $f_1(x) - x_1, \dots, f_N(x) - x_N$ defines a submanifold. Let I, \tilde{I} be ideals of $\mathbf{R}[[x]], \mathbf{R}[[y]]$ and let $\sqrt{I}, \sqrt{\tilde{I}}$ define two formal submanifolds M, \tilde{M} , respectively. We say that a formal map $y = G(x)$ maps M into \tilde{M} if $\tilde{I} \circ G \subset \sqrt{I}$. If M, \tilde{M} are in the same space, we write $M \subset \tilde{M}$ if $\tilde{I} \subset \sqrt{I}$. We say that a formal map F fixes M pointwise if $I(M)$ contain each component of the mapping $F - \text{I}$.

3.2. Formal deck transformations. Consider a formal $(2p)$ -submanifold in \mathbf{C}^{2p} defined by

$$(3.1) \quad M: z_{p+j} = E_j(z', \bar{z}'), \quad 1 \leq j \leq p.$$

Here E_j are formal power series in z', \bar{z}' . We assume that

$$(3.2) \quad E_j(z', \bar{z}') = h_j(z', \bar{z}') + q_j(\bar{z}') + O(|(z', \bar{z}')|^3)$$

and h_j, q_j are homogeneous quadratic polynomials. The formal complexification of M is defined by

$$\begin{cases} z_{p+i} = E_i(z', w'), & i = 1, \dots, p, \\ w_{p+i} = \bar{E}_i(w', z'), & i = 1, \dots, p. \end{cases}$$

We define a *formal deck transformation* of π_1 to be a formal biholomorphic map

$$\tau: (z', w') \rightarrow (z', f(z', w')), \quad \tau(0) = 0$$

such that $\pi_1 \tau = \pi_1$, i.e. $E \circ \tau = E$. Recall that condition B says that $q_* = \dim\{z' \in \mathbf{C}^n: q(z') = 0\}$ is zero, i.e. q vanishes only at the origin in \mathbf{C}^p .

Lemma 3.1. *Let M be a formal submanifold defined by (3.1)-(3.2). Suppose that M satisfies condition B. Then formal deck transformations of π_1 are commutative involutions. Each formal deck transformation τ of $\pi_1: \mathcal{M} \rightarrow \mathbf{C}^p$ is uniquely determined by its linear part $L\tau$ in the (z', w') coordinates, while $L\tau$ is a deck transformation for the complexification for $\pi_1: \mathcal{Q} \rightarrow \mathbf{C}^p$, where \mathcal{Q} is the complexification of the quadratic part Q of M . If M is real analytic, all formal deck transformations of π_1 are convergent.*

Proof. Let us recall some results about the quadric Q . We already know that $q_* = 0$ implies that π_1 for the complexification of Q is a branched covering. As used in the proof of Lemma 2.1, π_1 is an open mapping near the origin and its regular values are dense. In particular, we have

$$(3.3) \quad \det \partial_{w'} \{h(z', w') + q(w')\} \neq 0.$$

Let τ be a formal deck transformation for M . To show that τ is an involution, we note that its linear part at the origin, $L\tau$, is a deck transformation of Q . Hence $L\tau$ is an involution. Replacing τ by the deck transformation τ^2 , we may assume that τ is tangent to the identity. Write

$$\tau(z', w') = (z', w' + u(z', w')).$$

We want to show that $u = 0$. Assume that $u(z', w') = O(|(z', w')|^k)$ and let u_k be homogeneous and of degree k such that $u(z', w') = u_k(z', w') + O(|(z', w')|^{k+1})$. We have

$$E(z', w' + u(z', w')) = E(z', w').$$

Comparing terms of order $k + 1$, we get

$$\partial_{w'} \{h(z', w') + q(w')\} u_k(z', w') = 0.$$

By (3.3), $u_k = 0$. This shows that each formal deck transformation τ of π_1 for M is an involution. As mentioned above, $L\tau$ is a deck transformation of π_1 for Q . Also if $\tau, \tilde{\tau}$ are commuting formal involutions then $\tau^{-1}\tilde{\tau}$ is an involution and $\tau = \tilde{\tau}$ if and only if $L\tau = L\tilde{\tau}$.

Assume now that M is real analytic. We want to show that each formal deck transformation τ is convergent. By a theorem of Artin [Art68] applied to the solution u , there is a convergent $\tilde{\tau}(z', w') = \tau(z', w') + O(|(z', w')|^2)$ such that $E \circ \tilde{\tau} = E$, i.e. $\tilde{\tau}$ is a deck transformation. Then $\tilde{\tau}^{-1}\tau$ is a deck transformation tangent to the identity. Since it is a formal involution by the above argument, then it must be identity. Therefore, $\tau = \tilde{\tau}$ converges. \square

Analogous to real analytic submanifolds, we say that a formal manifold defined by (3.1)-(3.2) satisfies condition D if its formal branched covering π_1 admits 2^p formal deck transformations.

Recall from section 2 that it is crucial to distinguish a special set of generators for the deck transformations in order to relate the classification of real analytic manifolds to the classification of certain $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$. The set of generators is uniquely determined by the dimension of fixed-point sets. We want to extend these results at the formal level.

Proposition 3.2. *Let M, \tilde{M} be formal p -submanifolds in \mathbf{C}^n of the form (3.1)-(3.2). Suppose that M, \tilde{M} satisfy condition D. Then the following hold :*

- (i) *M and \tilde{M} are formally equivalent if and only if their associated families of involutions $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$ and $\{\tilde{\tau}_{11}, \dots, \tilde{\tau}_{1p}, \rho\}$ are formally equivalent.*

(ii) Let $\mathcal{T}_1 = \{\tau_{11}, \dots, \tau_{1p}\}$ be a family of formal holomorphic involutions which commute pairwise. Suppose that the tangent spaces of $\text{Fix}(\tau_{11}), \dots, \text{Fix}(\tau_{1p})$ are hyperplanes intersecting transversally at the origin. Let ρ be an anti-holomorphic formal involution and let $\mathcal{T}_2 = \{\tau_{21}, \dots, \tau_{2p}\}$ with $\tau_{2j} = \rho\tau_{1j}\rho$. Suppose that $\sigma = \tau_1\tau_2$ has distinct eigenvalues for $\tau_i = \tau_{i1} \cdots \tau_{ip}$, and

$$[\mathfrak{M}_n]_1^{L\mathcal{T}_1} \cap [\mathfrak{M}_n]_1^{L\mathcal{T}_2} = \{0\}.$$

There exists a formal submanifold defined by

$$(3.4) \quad z'' = (B_1^2, \dots, B_p^2)(z', \bar{z}')$$

for some formal power series B_1, \dots, B_p such that M satisfies condition D. The set of involutions $\{\tilde{\tau}_{11}, \dots, \tilde{\tau}_{1p}, \rho_0\}$ of M is formally equivalent to $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$.

Proof. (i) Let M and \tilde{M} be given by $z'' = E(z', \bar{z}')$ and $\tilde{z}'' = \tilde{E}(\tilde{z}', \bar{\tilde{z}}')$, respectively. Suppose that f is a formal holomorphic transformation sending M into \tilde{M} . We have

$$(3.5) \quad f''(z', E(z', w')) = \tilde{E}(f'(z', E(z', w')), \bar{f}'(w', \bar{E}(w', z'))).$$

Here $f = (f', f'')$. Recall that $\rho_0(z', w') = (\bar{w}', \bar{z}')$. Define a formal mapping $(z', w') \rightarrow (\tilde{z}', \tilde{w}') = F(z', w')$ by

$$(3.6) \quad F(z', w') := (f'(z', E(z', w')), \bar{f}'(w', \bar{E}(w', z'))).$$

It is clear that $F\rho_0 = \rho_0F$. By Lemma 2.5, we know that \tilde{z}' and $\tilde{z}'' = \tilde{E}(\tilde{z}', \bar{\tilde{z}}')$ generate invariant formal power series of $\{\tilde{\tau}_{1j}\}$. Thus, $\tilde{z}' \circ F(z', w') = f'(z', E(z', w'))$ and $\tilde{E} \circ F(z', w')$ are invariant by $F^{-1} \circ \tilde{\tau}_{1j} \circ F$. By (3.5) and the definition of F ,

$$\tilde{E} \circ F(z', w') = f''(z', E(z', w')).$$

This shows that $f(z', E(z', w'))$ is invariant under $F^{-1} \circ \tilde{\tau}_{1j} \circ F$. Since f is invertible, then z' and $E(z', w')$ are invariant under $F^{-1} \circ \tilde{\tau}_{1j} \circ F$. Therefore, $\{\tau_{1j}\}$ and $\{F^{-1} \circ \tilde{\tau}_{1j} \circ F\}$ are the same by Lemma 2.5 as they have the same invariant functions.

Assume now that $\{\tau_{1j}\} = \{F^{-1} \circ \tilde{\tau}_{1j} \circ F\}$ for some formal biholomorphic map F commuting with ρ_0 . Recall that \tilde{z}', \tilde{z}'' are invariant by $\tilde{\tau}_{1j}$. Then $\tilde{z}' \circ F$ and $\tilde{E} \circ F$ are invariant by $\{\tau_{1j}\}$. By Lemma 2.5, invariant power series of τ_{1j} are generated by $z', E(z', w')$. Thus we can write

$$(3.7) \quad \begin{aligned} \tilde{z}' \circ F(z', w') &= f'(z', E(z', w')), \\ \tilde{E} \circ F(z', w') &= f''(z', E(z', w')) \end{aligned}$$

for some formal power series map $f = (f', f'')$. Since $\rho_0F = F\rho_0$, then by (3.6)

$$F(z', w') = (f'(z', 0), \bar{f}'(w', 0)) + O(|(z', w')|^2).$$

Since F is (formally) biholomorphic then $z' \rightarrow f'(z', 0)$ is biholomorphic. Then

$$f''(0, E(0, w')) = \tilde{E}(0, \bar{f}'(w', 0)) + O(|w'|^3).$$

We have $E(0, w') = q(w') + O(|w'|^3)$ and $\tilde{E}(0, w') = \tilde{q}(w') + O(|w'|^3)$. Here $q(w'), \tilde{q}(w')$ are quadratic. By condition $q_* = 0$, we know that $\tilde{q}_1, \dots, \tilde{q}_p$ and hence $\tilde{q}_1 \circ L, \dots, \tilde{q}_p \circ L$ are linearly independent. Here L is the linear part of the mapping $w' \rightarrow \bar{f}'(w', 0)$, which is invertible. This shows that the linear part of $w' \rightarrow f''(0, w')$ is biholomorphic. By (3.7),

$f''(z', 0) = O(|z'|^2)$. Hence $f = (f', f'')$ is biholomorphic. By a simple computation, we have $f(M) = \tilde{M}$, i.e.

$$\tilde{E}(f'(z), \overline{f'(z)}) = f''(z)$$

for $z'' = E(z', \bar{z}')$.

(ii) Assume that $\{\tau_{1j}\}$ and ρ are given in the (ξ, η) space. We want to show that a formal holomorphic equivalence class of $\{\tau_{1j}, \rho\}$ can be realized by a formal submanifold satisfying condition D. The proof is almost identical to the realization proof of Proposition 2.8 and we will be brief. Using a formal, instead of convergent, change of coordinates, we know that invariant formal power series of $\{\tau_{1j}\}$ are generated by

$$z' = (A_1(\xi, \eta), \dots, A_p(\xi, \eta)), \quad z'' = (B_1^2(\xi, \eta), \dots, B_p^2(\xi, \eta)),$$

where B_j is skew-invariant by τ_{1j} , and A, B_i are invariant under τ_{1j} for $i \neq j$. Moreover, $\phi(\xi, \eta) = (A, B)(\xi, \eta)$ is formally biholomorphic. Set

$$w'_j = \overline{A_j \circ \rho(\xi, \eta)}, \quad w''_j = \overline{B_j^2 \circ \rho(\xi, \eta)}.$$

Then $(\xi, \eta) \rightarrow (A(\xi, \eta), \overline{A \circ \rho(\xi, \eta)})$ has an inverse ψ . Define

$$M: z'' = (B_1^2, \dots, B_p^2) \circ \psi(z', \bar{z}').$$

The complexification of M is given by

$$\mathcal{M}: z'' = (B_1^2, \dots, B_p^2) \circ \psi(z', w'), \quad w'' = (\overline{B_1^2}, \dots, \overline{B_p^2}) \circ \bar{\psi}(w', z').$$

Note that $\phi \circ \psi(z', w') = (z', B \circ \psi(z', w'))$. Since $\phi\psi$ is invertible, the linear part D of $B \circ \psi$ satisfies $|D(0, w')| \geq |w'|/C$. This shows that $q_* = 0$. As in the proof of Proposition 2.8, we can verify that M is a realization for $\{\tau_{1j}, \rho\}$. \square

4. NORMAL FORMS OF COMMUTING BIHOLOMORPHISMS

In this section, we shall consider a family of commuting germs of holomorphic diffeomorphisms at a common fixed point, say $0 \in \mathbf{C}^n$. We shall give conditions that ensure that the family can be transformed simultaneously and holomorphically to a normal form. This means that there exists a germ of biholomorphism at the origin which conjugates each germ of biholomorphism in the family to a mapping that commutes with the linear part of every mapping of the family.

4.1. Centralizers and Decomposition.

Definition 4.1. Let \mathcal{F} be a family of formal mappings on \mathbf{C}^n fixing the origin. Let $\mathcal{C}(\mathcal{F})$ be the *centralizer* of \mathcal{F} , i.e. the set of formal holomorphic mappings g that fix the origin and commute with each element f of \mathcal{F} , i.e., $f \circ g = g \circ f$.

Let $\mathcal{C}_2(\mathcal{F})$ be the ‘‘higher order formal centralizer’’ of \mathcal{F} , that is

$$\mathcal{C}_2(\mathcal{F}) = \{H \in (\widehat{\mathfrak{M}}_n^2)^n \mid H \circ F = F \circ H, F \in \mathcal{F}\}.$$

We now deal with the following decomposition problem: Let \mathcal{C} be a set of analytic mappings. We want to decompose an arbitrary invertible mapping into the composition of an element of a centralizer of \mathcal{C} and an element which is normalized with respect to \mathcal{C} . We first prove a general convergence decomposition, which will be used several times. Let e_j denote the standard j th unit vector of \mathbf{C}^n .

Definition 4.2. Let \mathcal{A} be a group of permutations of $\{1, \dots, n\}$. Then \mathcal{A} acts on the right (resp. on the left) on $\widehat{\mathcal{O}}_n^n$ by permutation of variables $z = (z_1, \dots, z_n)$ as follows: Let $F(z) = \sum_{|Q|>0} F_Q z^Q$ be a formal mapping from \mathbf{C}^n to \mathbf{C}^n , and let $\nu, \mu \in \mathcal{A}$; set

$$\nu \circ F \circ \mu(z) := \sum_i \sum_{Q \in \mathbf{N}^n} F_{\nu(i), \mu^{-1}(Q)} z^Q e_i.$$

Define the components $(\mathcal{A}F)_i, (F\mathcal{A})_i$, and consequently $(\mathcal{A}F\mathcal{A})_i$ by

$$\begin{aligned} (\mathcal{A}F)_i(z) &:= \sum_{Q \in \mathbf{N}^n} \max_{\nu \in \mathcal{A}} |F_{\nu(i), Q}| z^Q, & (F\mathcal{A})_i(z) &:= \sum_{Q \in \mathbf{N}^n} \max_{\mu \in \mathcal{A}} |F_{i, \mu^{-1}(Q)}| z^Q, \\ (\mathcal{A}F\mathcal{A})_i(z) &= \sum_{Q \in \mathbf{N}^n} \max_{(\nu, \mu) \in \mathcal{A}^2} |F_{\nu(i), \mu^{-1}(Q)}| z^Q. \end{aligned}$$

We see that $F\mathcal{A}$ is the smallest (w.r.t. \prec) power series mapping that majorizes F and is right-invariant under \mathcal{A} , while $\mathcal{A}F$ is the smallest power series mapping that majorizes F and is left-invariant under \mathcal{A} . In particular, if F, G are mappings without constant or linear terms, then

$$(4.1) \quad \mathcal{A}(F \circ (I + G))\mathcal{A} \prec (\mathcal{A}F\mathcal{A})(\mathcal{A}I\mathcal{A} + \mathcal{A}G\mathcal{A}),$$

where the last relation holds if the composition is well-defined.

To simplify our notation, we will take \mathcal{A} to be the full permutation group of $\{1, \dots, n\}$. We will denote

$$F_{sym} = \mathcal{A}F\mathcal{A}.$$

Lemma 4.3. Let $\widehat{\mathcal{H}}$ be a real subspace of $(\widehat{\mathfrak{M}}_n^2)^n$. Let $\pi : (\widehat{\mathfrak{M}}_n^2)^n \rightarrow \widehat{\mathcal{H}}$ be a \mathbf{R} linear projection (i.e. $\pi^2 = \pi$) that preserves the degrees of the mappings and let $\widehat{\mathcal{G}} := (I - \pi)(\widehat{\mathfrak{M}}_n^2)^n$. Suppose that there is a positive constant C such that

$$(4.2) \quad \pi(E) \prec CE_{sym}$$

for any $E \in (\widehat{\mathfrak{M}}_n^2)^n$. Let F be a formal map tangent to the identity. There exists a unique decomposition

$$(4.3) \quad F = HG^{-1}$$

with $G - I \in \widehat{\mathcal{G}}$ and $H - I \in \widehat{\mathcal{H}}$. If F is convergent, then G and H are also convergent.

Proof. If f is a formal mapping, we define the k -jet:

$$J^k f(z) = \sum_{|Q| \leq k} f_Q z^Q.$$

Write $F = I + f$, $G = I + g$ and $H = I + h$. We need to solve $FG = H$, i.e to solve

$$h - g = f(I + g).$$

Since $f'(0) = 0$, then for any $k \geq 2$, the k -jet of $f(I + g)$ depends only on the $(k - 1)$ -jet of g . Since π is linear and preserves degrees, (4.2) implies that J^k commutes with π . Hence we can define, for all $k \geq 2$,

$$-J^k(g) := \pi(J^k(f(I + g))), \quad J^k(h) := (I - \pi)(J^k(f(I + g))).$$

This solves the formal decomposition uniquely. Assume that F is a germ of holomorphic mapping. Hence, we have

$$g \prec C(f(I + g))_{sym} \prec C f_{sym}(I_{sym} + g_{sym}).$$

Since g_{sym} is the smallest left and right \mathcal{A} invariant power series that dominates g , we have

$$g_{sym} \prec C f_{sym}(I_{sym} + g_{sym}).$$

Therefore, g_{sym} is dominated by the solution u to

$$u = C f_{sym}(I_{sym} + u), \quad u(0) = 0.$$

Notice that u is real analytic near the origin by the implicit function theorem. So, g_{sym} is convergent, and $g, h = g + f(I + g)$ are convergent. \square

Remark 4.4. Let \mathcal{A}, \mathcal{B} be two subgroups of permutations. Instead of using the full permutations group, we could have used $G_{sym} := \mathcal{A}\mathcal{G}\mathcal{B}$. We have

$$G \prec \mathcal{A}\mathcal{G}\mathcal{B} \prec \mathcal{C}\mathcal{A}(F \circ (I + G))\mathcal{B} \prec (\mathcal{A}F\mathcal{A})(\mathcal{A}I\mathcal{B} + \mathcal{A}\mathcal{G}\mathcal{B}).$$

4.2. Abelian family of biholomorphisms. Let $\mathbf{D}_i := \text{diag}(\mu_{i1}, \dots, \mu_{in})$ with $1 \leq i \leq \ell$ be diagonal invertible matrices of \mathbf{C}^n . Let $D_i: x \rightarrow \mathbf{D}_i x$ be the linear mappings. Let D denote the family $\{D_i\}_{i=1, \dots, \ell}$ of linear mappings.

Definition 4.5. We say that $F = (f_1, \dots, f_n)$ is *normalized* with respect to D if it is tangent to the identity and it does not have components along the centralizer of D , i.e. for each Q with $|Q| \geq 2$,

$$f_{j,Q} = 0, \quad \text{if } \mu_i^Q = \mu_{ij} \text{ for all } i.$$

Let $\mathcal{C}^c(D)$ denote the set of formal mappings normalized with respect to D . Let $\mathcal{C}_2^c(D)$ be the set of all $H \in (\widehat{\mathfrak{M}}_n^2)^n$ satisfying $I + H \in \mathcal{C}^c(D)$.

Let us consider a family $F := \{F_i\}_{i=1}^\ell$ of germs of holomorphic diffeomorphisms of $(\mathbf{C}^n, 0)$ of which the linear of $F_i(x)$ at the origin is D_i . Thus

$$F_i(x) = \mathbf{D}_i x + f_i(x), \quad f_i(0) = 0, \quad D f_i(0) = 0.$$

The group of germs of (resp. formal) biholomorphisms tangent to identity acts on the family F by $\Phi_* F := \{\Phi^{-1} \circ F_i \circ \Phi: 1 \leq i \leq \ell\}$.

Let $\{F_i\}_{i=1, \dots, \ell}$ be a family of **commuting germs of biholomorphisms** with $\ell < \infty$. Let us recall a result by M. Chaperon (see theorem 4 in [Cha86], page 132):

Proposition 4.6. *If the family of diffeomorphisms is abelian then there exists a formal diffeomorphism Φ , which is tangent to the identity, such that*

$$\widehat{F}_i(\mathbf{D}_j x) = \mathbf{D}_j \widehat{F}_i(x), \quad 1 \leq i, j \leq \ell$$

where $\widehat{F}_i := \Phi_* F_i$, for $1 \leq i \leq \ell$. We call the family $\{\widehat{F}_i\}$ a *formal normal form of the family F (or a normalized family) with respect to the family D of linear maps.*

For convenience, we restrict changes of holomorphic coordinates to the ones that are tangent to the identity. Also $\Phi_* \{F_i\}_{i=1}^\ell = \{\widehat{F}_i\}_{i=1}^\ell$ means that

$$\Phi_* F_i = \widehat{F}_i, \quad 1 \leq i \leq \ell.$$

These restrictions will be removed by mild changes. For instance, if Φ transforms a family F into a family \hat{F} that commutes with LF , the family of the linear part of the F , then $(L\Phi)^{-1}(LF_i)L\Phi = L\hat{F}_i$. Therefore, $\Phi(L\Phi)^{-1}$ is tangent to the identity and transforms F into $(L\Phi)\hat{F}(L\Phi)^{-1}$ which commutes with LF .

Let $\hat{\mathcal{O}}_n^D$ be the ring of formal invariants of the family D , that is

$$\hat{\mathcal{O}}_n^D := \{f \in \hat{\mathcal{O}}_n \mid f(\mathbf{D}_i x) = f(x), i = 1, \dots, \ell\}.$$

If $Q \in \mathbf{N}^n$, $Q \neq 0$, then $x^Q \in \hat{\mathcal{O}}_n^D$ if and only if

$$\mu_i^Q := \mu_{i1}^{q_1} \cdots \mu_{in}^{q_n} = 1, \quad \forall 1 \leq i \leq \ell.$$

If $|Q| > 1$, then $x^Q e_j \in \mathcal{C}_2(D)$ if and only if $\mu_i^Q = \mu_{ij}, \forall 1 \leq i \leq \ell$. With notation of Definition 4.1, we have

Lemma 4.7. *Any formal diffeomorphism Φ of $(\mathbf{C}^n, 0)$, tangent to identity, can be written uniquely as $\Phi = \Phi_1 \circ \Phi_0^{-1}$ with $\Phi_1 \in \mathcal{C}^c(D)$ and $\Phi_0 \in \mathcal{C}(D)$. Furthermore, Φ_0, Φ_1 are convergent when Φ is convergent.*

Proof. This follows from Lemma 4.3, where $\hat{\mathcal{H}}$ is replaced by $\mathcal{C}_2(D)$ and π is defined by

$$\pi \left(\sum f_{j,Q} x^Q e_j \right) = \sum_j \sum_{x^Q e_j \notin \mathcal{C}_2(D)} f_{j,Q} x^Q e_j. \quad \square$$

Lemma 4.8. *Let $\hat{F} := \{\hat{F}_i\}$ and $\tilde{F} := \{\tilde{F}_i\}$ be two formal normal forms of the abelian family of diffeomorphisms $F := \{F_i\}$. There exists a formal diffeomorphism Φ , tangent to identity at the origin, such that $\Phi \in \mathcal{C}(D)$ and $\Phi \circ \tilde{F}_i = \hat{F}_i \circ \Phi$. Furthermore, there is a unique $\Phi \in \mathcal{C}^c(D)$ that transforms the family F into a normal form.*

Proof. Since both \hat{F} and \tilde{F} are normal forms of F , there exists a formal diffeomorphism Φ , tangent to identity at the origin, such that $\tilde{F}_i \circ \Phi = \Phi \circ \hat{F}_i$. According to Lemma 4.7, we can decompose $\Phi = \Phi_1 \circ \Phi_0^{-1}$ where $\Phi_0 \in \mathcal{C}(D)$ and $\Phi_1 \in \mathcal{C}^c(D)$. Hence, we have $\Phi_1^{-1} \circ \tilde{F}_i \circ \Phi_1 = \Phi_0^{-1} \circ \hat{F}_i \circ \Phi_0$. Let us set $G_i := \Phi_0^{-1} \circ \hat{F}_i \circ \Phi_0$. Then G_i is a formal diffeomorphism satisfying $G_i(x) - \mathbf{D}_i x \in \mathcal{C}(D)$. Let us show by induction on $N \geq 2$ that if $\Phi_1 = I + \Phi_1^N + O(N+1)$ with Φ_1^N being homogeneous of degree N , then $\Phi_1^N = 0$. Indeed, a computation shows that

$$\{G_i\}_N = \{\hat{F}_i\}_N + D_i \circ \Phi_1^N - \Phi_1^N \circ D_i.$$

Express Φ_1^N as sum of monomial mappings. The monomial mappings are not in $\mathcal{C}(D)$, while those of F_i and G_i are. We obtain $\Phi_1^N = 0$.

To verify the last assertion, assume that $\Psi_* F = \hat{F}$ and $\tilde{\Psi}_* F = \tilde{F}$ are in the normal form. Suppose that $\Psi, \tilde{\Psi}$ are normalized. Then $(\Psi^{-1}\tilde{\Psi})_* \hat{F} = \tilde{\Psi}_*(\Psi^{-1})_* \hat{F}$ is in the normal form. Write $\Psi^{-1}\tilde{\Psi} = \psi_1 \psi_0^{-1}$ with $\psi_1 \in \mathcal{C}^c(D)$ and $\psi_0 \in \mathcal{C}(D)$. Then $(\psi_1)_* \hat{F}$ is in a normal form. From the above proof, we know that $\psi_1 = I$. Now $\Psi = \tilde{\Psi} \psi_0$, which implies that $\Psi = \tilde{\Psi}$. \square

Lemma 4.9. *If a formal normal form of F is completely integrable so are all other normal forms of F ; in particular, the unique Φ in Lemma 4.8 transforms F into a completely integrable normal form.*

Proof. By Lemma 4.8, we transform a normal form $\{\hat{F}_i\}$ into another one $\{\tilde{F}_i\}$ by applying a transformation Φ that commutes with each D_j . Hence, we have $\hat{F}_i := \Phi^{-1} \circ \tilde{F}_i \circ \Phi$, for all $i = 1, \dots, \ell$. Let us write $\Phi(x) = \sum_{Q \in \mathbf{N}^n, 1 \leq j \leq n} \phi_{j,Q} x^Q e_j$. Suppose that $\{\hat{F}_i\}$ is completely integrable, then

$$\Phi \circ \hat{F}_i(x) = \sum_{Q \in \mathbf{N}^n} \phi_{j,Q} \mu_i(x)^Q x^Q e_j.$$

Since Φ commutes with each D_j , then

$$\Phi \circ \hat{F}_i(x) = \text{diag}(\mu_{i1}(x), \dots, \mu_{in}(x)) \cdot \Phi(x).$$

The conjugacy equation leads to

$$\mu_{ij}(x) \cdot \Phi_j(x) = \tilde{F}_{ij}(\Phi(x)), \quad 1 \leq j \leq n.$$

As a consequence, we have $\tilde{F}_i(x) = \text{diag}((\tilde{\mu}_{i1}(x), \dots, \tilde{\mu}_{in}(x))) \cdot x$ with $(\tilde{\mu}_{ij} \circ \Phi(x)) \cdot \Phi_j(x) = \mu_{ij}(x) \cdot \Phi_j(x)$, i.e. $\tilde{\mu}_{ij} = \mu_{ij} \circ \Phi^{-1}$.

Each function $\tilde{\mu}_{ij}$ is an invariant function of D since

$$\tilde{\mu}_{ij}(\mathbf{D}_k x) = \mu_{ij} \circ \Phi^{-1}(\mathbf{D}_k x) = \mu_{ij} \circ D_k(\Phi^{-1}(x)) = \mu_{ij} \circ \Phi^{-1}(x).$$

The second and third conditions of the definition of the complete integrability is obviously satisfied by $\{\tilde{F}_i\}$ since $\tilde{\mu}_{ij} = \mu_{ij} \circ \Phi^{-1}$. \square

Lemma 4.10. *If a formal normal form of F is linear so are all other normal forms of F .*

Proof. According to Lemma 4.8, we transform a linear normal form $\{\hat{F}_i\}$ into another one $\{\tilde{F}_i\}$ by applying a transformation Φ that commutes with each D_j . Since $\hat{F}_i(x) = \mathbf{D}_i x$, we have $\tilde{F}_i = \Phi^{-1}(D_i \Phi(x))$, for all $i = 1, \dots, \ell$. Since Φ commutes with each map $x \mapsto \mathbf{D}_i x$, then

$$\tilde{F}_i = \Phi^{-1}(D_i \Phi(x)) = \Phi^{-1}(\Phi(\mathbf{D}_i x)) = \mathbf{D}_i x. \quad \square$$

Definition 4.11. We say that *the family D is of Poincaré type* if there exist constants $d > 1$ and $c > 0$ such that, for each $(j, Q) \in \{1, \dots, n\} \times \mathbf{N}^n$ that satisfies $\mu_m^Q - \mu_{mj} \neq 0$ for some m , there exists $(i, Q') \in \{1, \dots, \ell\} \times \mathbf{N}^n$ such that $\mu_k^{Q'} = \mu_k^Q$ for all $1 \leq k \leq \ell$, $\mu_i^{Q'} - \mu_{ij} \neq 0$, and

$$\max(|\mu_i^{Q'}|, |\mu_i^{-Q'}|) > c^{-1} d^{|Q'|}, \quad Q' - Q \in \mathbf{N}^n \cup (-\mathbf{N}^n).$$

Such a condition has appeared in the definition of the good set in [BHV10].

Definition 4.12. Let $f = \sum_{Q \in \mathbf{N}^n} f_Q x^Q$ and $g = \sum_{Q \in \mathbf{N}^n} g_Q x^Q$ be two formal power series. We say that g *majorizes* f , written as $f \prec g$, if $g_Q \geq 0$ and $|f_Q| \leq g_Q$ for all $Q \in \mathbf{N}^n$. Set

$$\bar{f} := \sum_{Q \in \mathbf{N}^n} |f_Q| x^Q.$$

Theorem 4.13. *Let F be an abelian family of germs of holomorphic diffeomorphisms at the origin of \mathbf{C}^n . Assume that it is formally completely integrable and that its linear part at the origin is of Poincaré type. Then F is holomorphically conjugated to a normal form $\hat{F} = \{\hat{F}_1, \dots, \hat{F}_\ell\}$ so that each \hat{F}_i is defined by*

$$x'_j = \mu_{ij}(x) x_j, \quad j = 1, \dots, n$$

where $\mu_{ij}(x)$ are germs of holomorphic functions invariant under D and $\mu_{ij}(0) = \mu_{ij}$. In fact, the unique normalized mapping Φ in Lemma 4.8 is convergent.

The primary example is the following Moser-Webster normal form of reversible mappings:

$$\hat{\sigma} : \xi' = M_1(\xi\eta, \zeta)\xi \quad \eta' = M_1^{-1}(\xi\eta, \zeta)\eta, \quad \zeta' = \zeta.$$

where $(\xi, \eta) \in \mathbf{C}^2$, $\zeta \in \mathbf{C}^{n-2}$, and $|M_1(0)| > 1$. Our convergence proof is inspired by the proof in Moser-Webster [MW83].

Proof. The last assertion follows from Lemma 4.7 and Lemma 4.9. Let us conjugate, simultaneously, each $F_i = \mathbf{D}_i x + f_i$ to $\hat{F}_i := \hat{\mathbf{D}}_i(x)x$ by the action of $\Phi(x) = x + \phi(x)$ where $\phi(0) = 0$ and $\phi'(0) = 0$. Here, $\hat{\mathbf{D}}_i(x)$ denotes the matrix $\text{diag}(\hat{\mu}_{i1}(x), \dots, \hat{\mu}_{in}(x))$ and each $\hat{\mu}_{ij}(x)$ is a formal power series invariant under D , i.e. $\hat{\mu}_{ij}(x) \in \hat{\mathcal{O}}_n^D$. We can assume that Φ does not have a non-zero component along the centralizer of D ; indeed, by Lemma 4.9, we can assume that Φ is normalized w.r.t D . Then, for each $i = 1, \dots, \ell$, we have

$$F_i \circ \Phi(x) = \mathbf{D}_i x + f_i(\Phi)(x) + \mathbf{D}_i \phi(x), \quad \Phi \circ \hat{F}_i(x) = \hat{\mathbf{D}}_i(x)x + \phi(\hat{F}_i)(x).$$

Equation $F_i \circ \Phi = \Phi \circ \hat{F}_i$ reads

$$(4.4) \quad \left(\phi(\hat{\mathbf{D}}_i(x)x) - \mathbf{D}_i \phi(x) \right) + \left(\hat{\mathbf{D}}_i(x) - \mathbf{D}_i \right) x = f_i(\Phi)(x) \quad i = 1, \dots, \ell.$$

Our convergence proof is based on two conditions: the existence of a formal $\phi \in \mathcal{C}^c(D)$ that satisfies the above equation, and the Poincaré type condition on the linear part D . We already know that ϕ is unique. We shall project equation (4.4) along the “non-resonant” space (i.e. the space $\mathcal{C}^c(D)$ of normalized mappings w.r.t. \mathbf{D}). The mapping ϕ also solves this last equation and we shall majorize it using that projected equation.

Let us first decompose these equations along the “resonant” and “non-resonant” parts, i.e. $\mathcal{C}_2(D)$ and $\mathcal{C}_2^\xi(D)$. Since $\phi = \sum_{Q \in \mathbf{N}^n, |Q| \geq 2} \phi_{j,Q} x^Q e_j$ is normalized then $\phi_{j,Q} = 0$ for some $Q \in \mathbf{N}^n$, $|Q| \geq 2$ and $1 \leq j \leq n$, if we have $\mu_m^Q = \mu_{mj}$ for all m . We recall that, since each \mathbf{D}_i is a diagonal matrix, then a map belongs to the centralizer of D if and only if each monomial map of its Taylor expansion at the origin belongs to this centralizer as well. Since the $\hat{\mu}_{ij}$ is a formal invariant function then

$$\phi(\hat{\mathbf{D}}_i(x)x) = \sum_{Q \in \mathbf{N}^n, |Q| \geq 2} \phi_{j,Q} \hat{\mu}_i^Q(x) x^Q e_j =: \sum_{Q' \in \mathbf{N}^n, |Q'| \geq 2} \psi_{j,Q'} x^{Q'} e_j.$$

The latter contains only non-resonant terms, that is that if $\mu_i^{Q'} = \mu_{ij}$ for all i , then $\psi_{j,Q'} = 0$. Indeed, $\hat{\mu}_i^Q(x)$ contains monomials of the form x^P with $\mu_i^P = 1$ for all $1 \leq i \leq \ell$. Hence, $\psi_{j,Q'}$ is a linear combination of $\phi_{j,Q}$ such that $Q' = Q + P$ with $\mu_i^P = 1$ for all i . Therefore, if $\mu_i^{Q'} = \mu_{ij}$ for all i , then for all these Q 's, we have $\mu_i^Q = \mu_i^{Q'} = \mu_{ij}$ for all i so that $\phi_{j,Q} = 0$; that is $\psi_{j,Q'} = 0$.

Hence, the projection on the resonant mappings in $\mathcal{C}_2(D)$ leads to

$$(4.5) \quad \left(\hat{\mathbf{D}}_i(x) - \mathbf{D}_i \right) x = \{f_i(\Phi)(x)\}_{res}, \quad i = 1, \dots, \ell.$$

Here for any formal mapping $g(x) = O(|x|^2)$ on \mathbf{C}^n , we define the projection on $\mathcal{C}_2(D)$ by

$$(4.6) \quad (g(x))_{res} = \sum_j \sum_{\forall i, \mu_i^Q = \mu_{ij}} g_{j,Q} x^Q e_j.$$

The projection g onto $\mathcal{C}_2^c(D)$ is defined as $g(x) - (g(x))_{res}$, i.e. it is the projection of g on the non-resonant mappings.

Let us consider the projection on the non-resonant mappings. We first need to decompose power series according to a non-homogeneous equivalence relation on their coefficients. Let us define the equivalence relation on $\{1, \dots, n\} \times \mathbf{N}^n$ by

$$(j, Q) \sim (\tilde{j}, \tilde{Q}), \text{ if } \mu_{ij} - \hat{\mu}_i^Q(x) = \mu_{i\tilde{j}} - \hat{\mu}_i^{\tilde{Q}}(x) \text{ for all } 1 \leq i \leq \ell.$$

Here the identities hold as formal power series. Let Δ be the set of the equivalence classes on the non-resonant multiindex set

$$\left\{ (j, Q) \in \{1, \dots, n\} \times \mathbf{N}^n : (\mu_1^Q - \mu_{1j}, \dots, \mu_\ell^Q - \mu_{\ell j}) \neq 0, |Q| > 1 \right\}.$$

If $\mu_k^Q - \mu_{kj} \neq 0$ for some k , clearly $\hat{\mu}_k^Q - \mu_{kj}$ is not identically zero. We can decompose any formal power series map f along these equivalent classes and the resonant part of the mapping. Let $\delta \in \Delta$ and $f = \sum_{Q \in \mathbf{N}^n, 1 \leq j \leq n} f_{j,Q} x^Q e_j$ with $f = O(2)$. We can write

$$(4.7) \quad f_\delta(x) := \sum_{(j,Q) \in \delta} f_{j,Q} x^Q e_j, \quad \sum_{\delta \in \Delta} f_\delta(x) \prec \bar{f}(x).$$

We denote by $\widehat{\mathfrak{M}}_{n,\delta}^n$ the vector space of such maps. To a given equivalence class δ , we associate a representative (j_δ, Q_δ) , and we shall identify an equation among n equations in (4.4) for estimation. Since ϕ contains no resonant mappings, then

$$(4.8) \quad \phi = \sum_{\delta \in \Delta} \phi_\delta.$$

Let us decompose the projection onto non-resonant mappings in $\mathcal{C}_2^c(D)$ of equation (4.4) along each equivalence class δ as follows. Using the definition of the equivalence class Δ , we obtain

$$(4.9) \quad \left[\hat{\mu}_i^{Q_\delta}(x) - \mu_{ij_\delta} \right] \phi_\delta(x) = \{f_i(\Phi)\}_\delta(x), \quad \forall i = 1, \dots, \ell$$

where $\{f\}_\delta$ denotes the projection of f on $\widehat{\mathfrak{M}}_{n,\delta}^n$, defined by (4.7).

For each $(j_\delta, Q_\delta) \in \delta$, we know that $\mu_k^{Q_\delta} - \mu_{kj_\delta} \neq 0$ for some k . By the Poincaré type condition, there exist i and $Q'_\delta \in \mathbf{N}^n$ such that

$$(4.10) \quad \mu_i^{Q'_\delta} - \mu_{ij_\delta} \neq 0; \quad \mu_m^{Q'_\delta} = \mu_m^{Q_\delta}, \quad \forall 1 \leq m \leq \ell; \quad Q'_\delta - Q_\delta \in \mathbf{N}^n \cup (-\mathbf{N}^n)$$

and, furthermore, one of the following holds:

$$(4.11) \quad |\mu_i^{Q'_\delta}| \leq cd^{-|Q'_\delta|}, \quad \text{or} \quad |\mu_i^{-Q'_\delta}| \leq cd^{-|Q'_\delta|}.$$

Here, $d > 1$ does not depend on Q_δ . So, let us use the i th equation of (4.9) to estimate ϕ_δ . We have, for that i ,

$$(4.12) \quad \phi_\delta = \left[\hat{\mu}_i^{Q_\delta} - \mu_{ij_\delta} \right]^{-1} \{f_i(\Phi)\}_\delta.$$

Therefore, we have established the uniqueness of ϕ under (4.8) and (4.12), and under the condition that ϕ satisfies the equation when (4.4) is projected onto $\mathcal{C}^c(D)$. The existence of ϕ is ensured by assumption. We now consider the convergence of ϕ . By (4.10) we obtain $\hat{\mu}_i^{Q'_\delta - Q_\delta} = 1$. This allows us to rewrite (4.12) as

$$(4.13) \quad \phi_\delta = \left[\hat{\mu}_i^{Q'_\delta} - \mu_{ij_\delta} \right]^{-1} \{f_i(\Phi)\}_\delta.$$

We majorize this power series.

Recall that $\hat{\mu}_{ij}(0) = \mu_{ij}$. Let us set $M_{ij}(x) := \mu_{ij}^{-1} \hat{\mu}_{ij}(x)$. We have $M_{ij}(0) = 1$ and we decompose $M_{ij}(x) = \sum_{Q \in \mathbf{N}^n} M_{ij,Q} x^Q$. Let us set $\mu^* := \max_{ij} \{|\mu_{ij}|, |\mu_{ij}^{-1}|\}$, and

$$m_i = \sum_{Q \in \mathbf{N}^n} \max_{1 \leq j \leq n} |M_{ij,Q}| x^Q, \quad m = \sum_{Q \in \mathbf{N}^n} \max_{1 \leq i \leq \ell, 1 \leq j \leq n} |M_{ij,Q}| x^Q.$$

Note that $m(0) = 1$. Then $M_{ij} \prec m$ and

$$M_{ij}^{-1} = \frac{1}{1 + (M_{ij} - 1)} \prec \frac{1}{1 - (m - 1)} = \frac{1}{2 - m}.$$

Here and in what follows, if $f(x)$ is a formal power series with $f(0) = 0$, then for any number $a \neq 0$, $\frac{1}{a - f(x)}$ stands for the formal power series in x for

$$\frac{1}{a} \left\{ 1 + \sum_{n=1}^{\infty} (a^{-1} f(x))^n \right\}.$$

To simplify notation in (4.13), let us write Q for Q'_δ and j for j_δ . We want to show that,

$$(4.14) \quad (\hat{\mu}_i^Q - \mu_{ij})^{-1} \prec S(m - 1).$$

Here $S(t)$ is a convergent power series independent of all (j, Q) 's of the form (j_δ, Q'_δ) 's. Fix d_1 with $1 < d_1 < d$. We consider the first case that $\mu^* c d^{-|Q|} > d_1^{-|Q|}$. Since $d > d_1$, we have only finitely many such Q 's (recall that each Q has the form Q'_δ). The function $M_i \mapsto \mu_{ij} - \mu_i^Q M_i^Q$ is holomorphic in $M_i \in \mathbf{C}^n$ at $M_i = (1, \dots, 1)$ and does not vanish at this point. Hence, the function

$$(\mu_{ij} - \hat{\mu}_i^Q)^{-1} = (\mu_{ij} - \mu_i^Q M_i^Q)^{-1}$$

is also holomorphic at $M_i = (1, \dots, 1)$. For all Q 's in the first case, we have

$$(\mu_{ij} - \hat{\mu}_i^Q)^{-1} \prec \frac{C}{1 - C(\bar{M}_{i1} - 1 + \dots + \bar{M}_{in} - 1)} \prec \frac{C}{1 - nC(m - 1)}.$$

Consider the second case that $\mu^* c d^{-|Q|} \leq d_1^{-|Q|}$. For the first case in (4.11), we obtain

$$\begin{aligned} (\hat{\mu}_i^Q - \mu_{ij})^{-1} &= -\mu_{ij}^{-1} (1 - \mu_{ij}^{-1} \mu_i^Q M_i^Q)^{-1} \prec \mu^* [1 - \mu^* c d^{-|Q|} m^{|Q|}]^{-1} \\ &\prec \mu^* [1 - d_1^{-|Q|} m^{|Q|}]^{-1} \prec \mu^* [1 - d_1^{-1} m]^{-1}. \end{aligned}$$

For the second case in (4.11), we have

$$\begin{aligned}
(\hat{\mu}_i^Q - \mu_{ij})^{-1} &= -\mu_i^{-Q} M_i^{-Q} \left[1 - \mu_{ij} \hat{\mu}_i^{-Q} M_i^{-Q} \right]^{-1} \\
&\prec cd^{-|Q|} (2-m)^{-|Q|} \left[1 - \mu^* cd^{-|Q|} (2-m)^{-|Q|} \right]^{-1} \\
&\prec (\mu^*)^{-1} d_1^{-|Q|} (2-m)^{-|Q|} \left[1 - d_1^{-|Q|} (2-m)^{-|Q|} \right]^{-1} \\
&\prec (\mu^*)^{-1} \left[1 - d_1^{-1} (2-m)^{-1} \right]^{-1}.
\end{aligned}$$

We have obtained the estimates for the second case. Therefore, we have shown that for any $Q = Q'_\delta$ and $j = j_\delta$

$$(4.15) \quad (\hat{\mu}_i^Q - \mu_{ij})^{-1} \prec S(m-1).$$

Here $S(t)$ is a convergent power series independent of all (j, Q) 's of the form (j_δ, Q'_δ) 's.

Let us set

$$f^* := \sum_{Q \in \mathbb{N}^n} \max_{1 \leq i \leq \ell, 1 \leq j \leq n} |f_{ij, Q}| x^Q e_j.$$

By the definition of the equivalence relation on multiindices, we have

$$(4.16) \quad \sum_{\delta \in \Delta} f_\delta^* \prec f^*.$$

According to (4.13) and (4.15), we have $\phi_\delta \prec S(m-1) \{f^*(\bar{\Phi})\}_\delta$. Now (4.7) and (4.16) imply

$$(4.17) \quad \phi \prec S(m-1) f^*(\bar{\Phi}).$$

Let us project (4.5) onto the k th components of $\mathcal{C}_2(D)$ as follows. For a power series map g , we define

$$g_{res, k}(x) = \sum_{\mu^Q = \mu_k} g_{k, Q} x^Q.$$

By the definition of g_{res} in (4.6), $g_{res} = (g_{res, 1}, \dots, g_{res, n})$. We have

$$\mu_{ik} (M_{i, k}(x) - 1) x_k = (\hat{\mu}_{ik}(x) - \mu_{ik}) x_k = \{f_{ik}(\Phi)\}_{res, k}(x).$$

Therefore, for all $1 \leq k \leq n$,

$$(4.18) \quad (m-1)x_k \prec \frac{1}{\min_{i, j} |\mu_{i, j}|} f^*(\bar{\Phi}).$$

Let us set $\mu_* := \frac{1}{\min_{i, j} |\mu_{i, j}|}$. We set $x_1 = t, \dots, x_n = t$ in $\bar{\Phi}(x)$ and $m(x)$. Let $\phi(t)$, $\bar{\Phi}(t)$, and $m(t)$ still denote $\phi(t, \dots, t)$, $\bar{\Phi}(t, \dots, t)$, and $m(t, \dots, t)$, respectively. Let

$$tW(t) := \phi(t) + (m(t) - 1)t.$$

We have $W(0) = 0$, $\phi(t) \prec tW(t)$, and $(m(t) - 1) \prec W(t)$. From estimates (4.17) and (4.18), we obtain

$$(4.19) \quad tW(t) \prec \mu_* f^*(\bar{\Phi}(t)) + S(m(t) - 1) f^*(\bar{\Phi}(t)).$$

Since $f_{ij}(x) = O(|x|^2)$, there exists a constant c_1 such that

$$f^*(x) \prec \frac{c_1(\sum_j x_j)^2}{1 - c_1(\sum_j x_j)}.$$

Hence, estimate (4.19) reads

$$\begin{aligned} (4.20) \quad tW(t) &\prec (\mu_* + S(m(t) - 1)) \frac{c_1(n(t + \phi))^2}{1 - c_1n(t + \phi)} \\ &\prec (\mu_* + S(W(t))) \frac{c_1t^2(n(1 + W(t)))^2}{1 - c_1nt(1 + W(t))}. \end{aligned}$$

Let us consider the equation in the unknown U with $U(0) = 0$:

$$(4.21) \quad U(t)(1 - c_1nt(1 + U(t))) = (\mu_* + S(W(t))) c_1t(n(1 + U(t)))^2.$$

According to the implicit function theorem, there exists a unique germ of holomorphic function $U(t)$, solution to (4.21) with $U(0) = 0$. According to inequality (4.20), the function W is dominated by U : $W(t) \prec U(t)$. This can be seen by induction on the degree of the Taylor polynomials at the origin. Therefore, W converges at the origin. The theorem is proved. \square

5. REAL MANIFOLDS WITH AN ABELIAN CR-SINGULARITY

Let us consider a real analytic manifold M with a CR-singularity at the origin, which is an higher order perturbation of a product quadric. We assume that for its complexification \mathcal{M} , π_1 has 2^p deck transformations generated by $\{\tau_{11}, \dots, \tau_{1p}\}$. Let $\tau_{2j} = \rho \circ \tau_{1j} \circ \rho$.

Let us consider the following germs of holomorphic diffeomorphisms :

$$(5.1) \quad \sigma_i := \tau_{1i} \circ \tau_{2i}, \quad 1 \leq i \leq e_* + h_*,$$

$$(5.2) \quad \sigma_s := \tau_{1s} \circ \tau_{2(s_*+s)}, \quad \sigma_{s+s_*} = \tau_{1(s+s_*)} \circ \tau_{2s}, \quad e_* + h_* < s \leq p - s_*.$$

Notice that the above property holds for quadrics of the complex case by Proposition 2.10. We assume that the linear parts T_{ij}, S_j, S of $\tau_{ij}, \sigma_j, \sigma$ and ρ are given by (2.32)-(2.35). The family $\{\sigma_i\}$ is reversible with respect to ρ . More precisely, we have the following relations

$$\sigma_i^{-1} = \rho \sigma_i \rho, \quad 1 \leq i \leq e_* + h_*; \quad \sigma_{s+s_*}^{-1} = \rho \sigma_s \rho, \quad e_* + h_* < s \leq p - s_*.$$

Definition 5.1. We say that the manifold M has an **abelian CR-singularity** at the origin if its complexification \mathcal{M} has the maximum number of deck transformation and if the family $\{\sigma_1, \dots, \sigma_p\}$ of germs of biholomorphisms at the origin of \mathbf{C}^{2p} is abelian, i.e.

$$\sigma_i \sigma_j = \sigma_j \sigma_i.$$

The aim of this section is to show that such an analytic perturbation with an abelian CR-singularity and no hyperbolic component is holomorphically conjugate to a normal form. We shall give two proofs of this result. The first one rests on Moser-Webster result [MW83][theorem 4.1] applied successively to each σ_i . It is to be emphasized that it is fortunate that we can apply such a result to our situation including the new type of CR singularity of *complex* type. The other one is based on the fact that the family $\{\sigma_i\}$ is formally completely integrable and their linear part is of Poincaré type. We then apply Theorem 4.13.

Before we apply the above theorem, let us first exhibit an example of real manifolds with an abelian CR singularity. We start with a Bishop surface $M_0 \subset \mathbf{C}^2$ defined by $z_2 = R(z_1, \bar{z}_1)$ with

$$(5.3) \quad R(z_1, \bar{z}_1) = z_1 \bar{z}_1 + \gamma_1(z_1^2 + \bar{z}_1^2) + O(3), \quad R(z_1, \bar{z}_1) = \overline{R(\bar{z}_1, z_1)}.$$

Let τ_1^0, τ_2^0 be the Moser-Webster involutions of M_0 . On the complexification \mathcal{M}_0 of M_0 , the z_1 and $z_2 := R(z_1, w_1)$ are invariant by τ_1^0 . Condition (5.3) implies that $w_2 := \overline{R(w_1, z_1)} = z_2$ is also invariant by τ_1^0 . Analogously, w_1, w_2 and z_2 are invariant by τ_2^0 . We are ready to verify the following

Proposition 5.2. *Let R be given by (5.3) and let*

$$(5.4) \quad M: z_3 = R(z_1, \bar{z}_1), \quad z_4 = (z_2 + 2\gamma_2 \bar{z}_2 + C(z_2, z_3, \bar{z}_2))^2,$$

where $C(z_2, R(z_1, \bar{z}_1), \bar{z}_2) = O(3)$. Then M has an abelian CR singularity at 0.

Proof. Let τ_{11} be defined by $(z'_1, w'_1) = \tau_1^0(z_1, w_1)$ and $(z'_2, w'_2) = (z_2, w_2)$. Let us verify that τ_{11} is a deck transformation of π_1 , i.e. all z_j are invariant by it. Obviously, z_2, w_2 are invariant by τ_{11} . We know that $z_1, z_3 = R(z_1, w_1)$ are invariant by τ_1^0 and hence by τ_{11} . On \mathcal{M} , $z_4 = (z_2 + 2\gamma_2 w_2 + C(z_2, z_3, w_2))^2$ is then invariant by τ_{11} . Therefore, τ_{11} is an involution and it fixes a hypersurface in \mathcal{M} . Let τ_{12} be defined by $(z'_1, z'_2, w'_1) = (z_1, z_2, w_1)$ and

$$z_2 + 2\gamma_2 w'_2 + C(z_2, R(z_1, w_1), w'_2) = -z_2 - 2\gamma_2 w_2 - C(z_2, R(z_1, w_1), w_2).$$

By the implicit function theorem, $w'_2 = -\gamma_2^{-1} z_2 - w_2 + f(z_2, R(z_1, w_1), w_2)$ with f being convergent. It is clear that τ_{12} leaves z_j invariant and it fixes a hypersurface.

For the abelian property, we note that z_2, w_2 and $R(z_1, w_1)$ are invariant by τ_{11}, τ_{21} . By a straightforward computation, we verify that τ_{11} and τ_{21} commute with τ_{12} . Now τ_{21} and τ_{11} commute with $\tau_{22} = \rho\tau_{12}\rho$. This shows that $\tau_{11}\tau_{21}$ commutes with $\tau_{12}\tau_{22}$. \square

5.1. Normal forms for abelian CR singularity.

Theorem 5.3. *Let M be a germ of real analytic submanifold in \mathbf{C}^n at an abelian CR-singularity at the origin. Suppose that M is a higher order perturbation of a product quadric of which $\gamma_1, \dots, \gamma_p$ satisfy (1.5). Suppose that M does not have a hyperbolic component (i.e. $e_* \geq 0, s_* \geq 0, h_* = 0$) and $\operatorname{Re} \gamma_s < 1/2$. Then there exists a germ of biholomorphism ψ that commutes with ρ and such that, for $1 \leq i \leq p$ and $k = 1, 2$*

$$(5.5) \quad \psi^{-1} \circ \sigma_i \circ \psi : \begin{cases} \xi'_i = M_i(\xi\eta)\xi_i \\ \eta'_i = M_i^{-1}(\xi\eta)\eta_i \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq i, \end{cases} \quad \psi^{-1} \circ \tau_{ki} \circ \psi : \begin{cases} \xi'_i = \Lambda_{ki}(\xi\eta)\eta_i \\ \eta'_i = \Lambda_{ki}^{-1}(\xi\eta)\xi_i \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq i. \end{cases}$$

Moreover, $\Lambda_{2j} = \Lambda_{1j}^{-1}$ and

$$(5.6) \quad \Lambda_{1e} = \overline{\Lambda_{1e} \circ \rho_z}, \quad 1 \leq e \leq e_*; \quad \Lambda_{1s}^{-1} = \overline{\Lambda_{1(s+s_*)} \circ \rho_z}, \quad e_* < s \leq p - s_*,$$

$$(5.7) \quad \rho_z : \zeta_e \rightarrow \bar{\zeta}_e, \quad \zeta_s \rightarrow \bar{\zeta}_{s+s_*}, \quad \zeta_{s+s_*} \rightarrow \bar{\zeta}_s.$$

Proof. We will present two convergence proofs: one is based on a convergent theorem of Moser and Webster and another is based on Theorem 4.13. We first use some formal results obtained by Moser and Webster [MW83] and some results in section 4. The conditions of the theorem imply that, for all i , $|\mu_i| \neq 1$.

Since M is a higher order perturbation of a product quadric, there are linear coordinates such that, for $1 \leq i \leq p$ and $k = 1, 2$, τ_{ki} and σ_i are higher order perturbations of

$$S_i : \begin{cases} \xi'_i = \mu_i \xi_i \\ \eta'_i = \mu_i^{-1} \eta_i \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq i, \end{cases} \quad T_{ki} : \begin{cases} \xi'_i = \lambda_{ki} \eta_i \\ \eta'_i = \lambda_{ki}^{-1} \xi_i \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq i. \end{cases}$$

For elliptic coordinates, this was computed in [MW83] and recalled in (2.21). For complex coordinates, this is computed in (2.27) and (2.29). Recall that $\sigma_1, \dots, \sigma_p$ are defined by (5.1)-(5.2) Since $|\mu_1| \neq 1$, then by theorem 4.1 of Moser-Webster ([MW83]), there is a unique convergent transformation ψ_1 normalized w.r.t. S_1 (see Definition 4.5) such that $\sigma_1^* := \psi_1^{-1} \circ \sigma_1 \circ \psi_1$ and $\tau_{i1}^* := \psi_1^{-1} \circ \tau_{i1} \circ \psi_1$ are given by

$$(5.8) \quad \sigma_1^* : \begin{cases} x'_1 = M_1(\xi, \eta) \xi_1 \\ \eta'_1 = M_1^{-1}(\xi, \eta) \eta_1 \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq 1, \end{cases} \quad \tau_{k1}^* : \begin{cases} \xi'_1 = \Lambda_{k1}(\xi, \eta) \eta_1 \\ \eta'_1 = \Lambda_{k1}^{-1}(\xi, \eta) \xi_1 \\ \xi'_j = \xi_j \\ \eta'_j = \eta_j, \quad j \neq 1. \end{cases}$$

Here $k = 1, 2$. It is a simple fact (e.g. see Lemma 4.9, $D = \{S_1\}$) that there is a unique $\phi_1 \in \mathcal{C}(S_1)$ such that $\phi_1^{-1} \sigma_1 \phi_1$ is in the centralizer of S_1 . Therefore, $\phi_1 = \psi_1$ is also convergent.

Furthermore, we have $M_1(\xi, \eta) = \Lambda_{11}(\xi, \eta) \Lambda_{21}^{-1}(\xi, \eta)$; and $\Lambda_{11}, \Lambda_{21}, M_1$ are invariant by S_1 . In the new coordinates, let us denote τ_{im}, σ_m by the same symbols for $m > 1$. However, $\sigma_1 = \sigma_1^*$ and $\tau_{k1} = \tau_{k1}^*$. Since each σ_m commutes with σ_1 , then σ_m is in the centralizer of S_1 . Indeed, according to [MW83][Lemma 3.1](or Lemma 4.7 with $D = \{S_1\}$), we can decompose $\sigma_m = \sigma_m^1 \sigma_m^0$ where σ_m^1 is normalized w.r.t S_1 and σ_m^0 is in the centralizer of S_1 . Write $\sigma_1 \sigma_m = \sigma_m \sigma_1$ as

$$(\sigma_m^1)^{-1} \sigma_1 \sigma_m^1 = \sigma_m^0 \sigma_1 (\sigma_m^0)^{-1}.$$

Since $\sigma_m^0 \sigma_1 (\sigma_m^0)^{-1}$ belongs to $\mathcal{C}(S_1)$, so does $(\sigma_m^1)^{-1} \sigma_1 \sigma_m^1$. Then applying the uniqueness of ψ_1 stated earlier to σ_m^1 , we conclude that $\sigma_m^1 = I$ and $\sigma_m = \sigma_m^0$ is in the centralizer of S_1 .

Let us verify that σ_m^0 or in general each (formal) transformation φ in $\mathcal{C}(S_1)$ preserves the form of σ_1^* and τ_{i1}^* . Indeed, φ^{-1} commutes with S_1 too. Thus $\varphi^{-1} \sigma_1^* \varphi$ commutes with S_1 and its linear part is S_1 . The linear part of φ must preserve the eigenspaces of S_1 and hence it is given by

$$\xi_1 \rightarrow a \xi_1, \quad \eta_1 \rightarrow b \eta_1, \quad (\xi_*, \eta_*) \rightarrow \phi(\xi_*, \eta_*)$$

for $\xi_* = (\xi_2, \dots, \xi_n)$ and $\eta_* = (\eta_2, \dots, \eta_m)$. The linear part of τ_{k1}^* is given by

$$\xi_1 \rightarrow \lambda_{k1} \eta_1, \quad \eta_1 \rightarrow \lambda_{k1}^{-1} \xi_1, \quad (\xi_*, \eta_*) \rightarrow (\xi_*, \eta_*).$$

By a simple computation, the linear part of $\varphi^{-1} \tau_{k1}^* \varphi$ still has this form with $\lambda_{k1} \frac{b}{a}$ instead of λ_{k1} . According to [MW83][lemma 3.2], there a unique normalized mapping Ψ that

normalizes $\varphi^{-1}\sigma_1^*\varphi$ and the $\varphi^{-1}\tau_{k1}^*\varphi$'s. According to the uniqueness property of Lemma 4.8, $\Psi = Id$. Therefore, φ preserves the forms of τ_{i1}^* and σ_1^* .

Let ψ_2 be the unique biholomorphic map normalized w.r.t. S_2 such that $\psi_2^{-1}\sigma_2\psi_2 = \sigma_2^*$ and $\psi_2^{-1}\tau_{k2}\psi_2 = \tau_{k2}^*$ are in the normal form :

$$(5.9) \quad \sigma_2^* : \begin{cases} \xi_2' = M_2(\xi, \eta)\xi_2 \\ \eta_2' = M_2^{-1}(\xi, \eta)\eta_2 \\ \xi_j' = \xi_j \\ \eta_j' = \eta_j, \quad j \neq 2, \end{cases} \quad \tau_{k2}^* : \begin{cases} \xi_2' = \Lambda_{k2}(\xi, \eta)\eta_2 \\ \eta_2' = \Lambda_{k2}^{-1}(\xi, \eta)\xi_2 \\ \xi_j' = \xi_j \\ \eta_j' = \eta_j, \quad j \neq 2. \end{cases}$$

Here $k = 1, 2$, and M_2 and Λ_{k2} are invariant by S_2 . Since σ_2 commutes with S_1 , we have

$$(S_1^{-1}\psi_2S_1)^{-1} \circ \sigma_2 \circ (S_1^{-1}\psi_2S_1) = S_1^{-1}\sigma_2^*S_1.$$

Note that $S_1^{-1}\sigma_2^*S_1$ (resp. $S_1^{-1}\tau_{k2}^*S_1$) has the form (5.9) in which M_2 (resp. Λ_{k2}) is replaced by $M_2 \circ S_1$ (resp. $\Lambda_{k2} \circ S_1$). In other words $S_1^{-1}\sigma_2^*S_1$ and $S_1^{-1}\tau_{k2}^*S_1$ are still of the form (5.9). Since S_1 is diagonal, then $S_1^{-1}\psi_2S_1$ remains normalized w.r.t. S_2 . Applying the above uniqueness on ψ_2 for σ_2 , we conclude that $\psi_2 = S_1^{-1}\psi_2S_1$. This shows that ψ_2 preserves the forms of τ_{k1}^* and σ_1^* . By the same argument as above, we have $\sigma_m^* \in \mathcal{C}(S_1, S_2)$.

In summary, we have found holomorphic coordinates so that $\tau_{ij} = \tau_{im}^*$ and $\sigma_m = \sigma_m^*$ for $m = 1, 2$. As mentioned previously, we know that $\sigma_1^*, \sigma_2^*, \sigma_3, \dots, \sigma_m$ commute with S_1 and S_2 . In particular, M_1, M_2 are invariant by S_1, S_2 . Repeating this procedure, we find a holomorphic map ϕ so that all $\phi^{-1}\sigma_j\phi = \sigma_j^*$ and $\phi^{-1}\tau_{kj}\phi = \tau_{kj}^*$ are in the normal forms. Furthermore, M_i and $\Lambda_{k,i}$ are invariant by $\mathcal{S} = \{S_1, \dots, S_p\}$.

By Lemma 4.7, we decompose $\phi = \phi_1\phi_0^{-1}$ where ϕ_1 is normalized w.r.t. \mathcal{S} and ϕ_0 is in the centralizer of \mathcal{S} . Then $\phi_1^{-1}\sigma_j\phi_1 = \sigma_j^*$ and $\phi_1^{-1}\tau_{ij}\phi_1 = \tau_{ij}^*$ are in the normal forms, since ϕ_0 commutes with S_j . We want to show that ϕ_1 commutes with ρ .

Note that $\sigma_e^{-1} = \rho\sigma_e\rho$ and $\sigma_{s+s_*}^{-1} = \rho\sigma_{s_*}\rho$. Thus $(\rho\phi_1\rho)^{-1}\sigma_j(\rho\phi_1\rho) = \tilde{\sigma}_j^*$ where $\tilde{\sigma}_e^* := \rho(\sigma_e^*)^{-1}\rho$ and $\tilde{\sigma}_{s_*}^* := \rho(\sigma_{s+s_*}^*)^{-1}\rho$. It is easy to see that $\rho\phi_1\rho$ is still normalized w.r.t. \mathcal{S} (see also Definition 6.4). By Lemma 4.8, we know that there is a unique normalized formal mapping ϕ_1 such that $\phi_1^{-1}\sigma_j\phi_1$ are in the centralizer of \mathcal{S} . Since $\tilde{\sigma}_j^*$ belongs to the centralizer of \mathcal{S} , then we have $\rho\phi_1\rho = \phi_1$.

Now, $\tau_{2j}^* = \rho\tau_{1j}^*\rho$ follows from $\tau_{2j} = \rho\tau_{1j}\rho$. This shows that

$$\Lambda_{2e} = \overline{\Lambda_{1e}^{-1} \circ \rho_z}, \quad \Lambda_{2s} = \overline{\Lambda_{1(s+s_*)} \circ \rho_z}, \quad \Lambda_{2(s+s_*)} = \overline{\Lambda_{1s} \circ \rho_z},$$

where $1 \leq e \leq e_*$ and $e_* < s \leq p - s_*$. Let ϕ_2 be defined by

$$\xi_j' = (\Lambda_{1j}^{1/2} M_j^{1/4})(\xi\eta)\xi_j, \quad \eta_j' = (\Lambda_{1j}^{-1/2} M_j^{-1/4})(\xi\eta)\eta_j, \quad 1 \leq j \leq p.$$

For a suitable choice of the roots, we have $\phi_2\rho = \rho\phi_2$. Furthermore, ϕ_2 preserves all invariant functions of \mathcal{S} . Hence, each $\phi_2^{-1} \circ \phi_1^{-1} \circ \tau_{ki} \circ \phi_1 \circ \phi_2$ has the form τ_{kj}^* stated in Theorem 5.3.

We now present another proof by using the more general Theorem 4.13.

Note that the above proof is valid at the formal level without using the convergence result of Moser and Webster. More specifically, if τ_{ij} are given by formal power series with $\sigma_1, \dots, \sigma_p$ commuting pairwise, there exists a formal map ψ that is tangent to the identity

and commutes with ρ such that (5.5) holds. Since each μ_j is not a root of unity, then (5.5) implies that the conjugate family $\{\sigma_m^*\}$ is a completely integrable normal form.

Let σ_i be defined as above. Let S_i be its linear part at the origin of \mathbf{C}^n . The eigenvalues $\{\mu_{ij}\}_{1 \leq j \leq n}$ of S_i are either μ_i , μ_i^{-1} or 1. More precisely, if $Q \in \mathbf{N}^n$, $|Q| \geq 2$ then

$$(5.10) \quad \mu_m^Q - \mu_{mj} = \mu_m^{q_m - q_{m+p}} - \begin{cases} \mu_m & \text{if } j = m \\ \mu_m^{-1} & \text{if } j = m + p \\ 1 & \text{otherwise.} \end{cases}$$

We need to verify the condition that the family of linear part $\{S_1, \dots, S_p\}$ is of the Poincaré type. So we can apply Theorem 4.13.

Suppose that $(j, Q) \in \{1, \dots, 2p\} \times \mathbf{N}^{2p}$ satisfies $\mu_l^Q - \mu_{lj} \neq 0$ for some $1 \leq l \leq 2p$. Set $d = \{\min_i \max(|\mu_i|, |\mu_i^{-1}|)\}^{1/(2p)}$. We define

$$Q' = Q - \sum_{i=1}^p \min(q_i, q_{i+p})(e_i + e_{i+p}) := (q'_1, \dots, q'_{2p}).$$

Then $\mu_i^Q = \mu_i^{Q'}$ for all i . Take $i = l$ if $|Q'| \leq 2p$. In this case, we easily get

$$(5.11) \quad \mu_i^{Q'} - \mu_{ij} \neq 0, \quad |\mu_i^{Q'}| > c^{-1}d^{|Q'|}$$

by choosing a sufficiently large c . Assume that $|Q'| > 2p$. Take i such that

$$q'_i + q'_{i+p} = \max_k (q'_k + q'_{k+p}).$$

Then $q'_i + q'_{i+p} \geq |Q'|/p > 2$. By (5.10), we get the first equality in (5.11). We note that $(q'_i, q'_{i+p}) = (q_i, 0)$ or $(0, q_{i+p})$. Thus

$$\max(|\mu_i^{Q'}|, |\mu_i^{-Q'}|) = (\max(|\mu_i|, |\mu_i^{-1}|))^{q'_i + q'_{i+p}} \geq d^{|Q'|}.$$

This shows that $\{D\sigma_1(0), \dots, D\sigma_p(0)\}$ is of the Poincaré type.

We now apply Theorem 4.13 as follows. We decompose $\psi = \psi_1 \psi_0^{-1}$ such that $\psi_1 \in \mathcal{C}^c(S_1, \dots, S_p)$ and $\psi_0 \in \mathcal{C}(S_1, \dots, S_p)$. Then each $\sigma_i^* = \psi_1^{-1} \sigma_i \psi_1$ still has the form in (5.5); in particular, $\{\sigma_1^*, \dots, \sigma_p^*\}$ is a completely integrable formal normal form. By Theorem 4.13, ψ_1 is convergent. Now, $\psi_1^{-1} \tau_{kj} \psi_1 = \psi_0^{-1} (\psi^{-1} \tau_{kj} \psi) \psi_0$ are still of the form (5.5); however (5.6) and $\Lambda_{2j} \Lambda_{1j} = 1$ might not hold. As in the first proof, we can verify that $\psi_1 \rho = \rho \psi_1$. Applying another change of coordinates that commutes with ρ and each S_j as before, we achieve (5.5)-(5.6) and $\Lambda_{2j} \Lambda_{1j} = 1$. The proof of the theorem is complete. \square

As a corollary of Theorem 5.3, we have the following normal form for real submanifolds. In order to study the holomorphic flatness and hull of holomorphy, we choose a realization similar to the case of Moser-Webster for $p = 1$.

Theorem 5.4. *Let M be as in Theorem 5.3. Then M is holomorphically equivalent to*

$$(5.12) \quad \widehat{M}: z_{p+j} = \Lambda_{1j}(\zeta) \zeta_j, \quad 1 \leq j \leq p,$$

where $\zeta = (\zeta_1, \dots, \zeta_p)$ are the convergent solutions to

$$(5.13) \quad \zeta_e = A_e(\zeta)|z_e|^2 - B_e(\zeta)(z_e^2 + \bar{z}_e^2),$$

$$(5.14) \quad \zeta_s = A_s(\zeta)z_s\bar{z}_{s+s_*} - B_s(\zeta)(z_s^2 + \Lambda_{1s}^2(\zeta)\bar{z}_{s+s_*}^2),$$

$$(5.15) \quad \zeta_{s+s_*} = A_{s+s_*}(\zeta)\bar{z}_s z_{s+s_*} - B_{s+s_*}(\zeta)(z_{s+s_*}^2 + \Lambda_{1(s+s_*)}^2(\zeta)\bar{z}_s^2).$$

Here $\Lambda_{1j}(\zeta) = \lambda_j + O(\zeta)$ ($1 \leq j \leq p$) satisfy (5.6) and

$$(5.16) \quad A_e = \frac{1 + \Lambda_{1e}^2}{(1 - \Lambda_{1e}^2)^2}, \quad A_j = \frac{\Lambda_{1j} + \Lambda_{1j}^3}{(1 - \Lambda_{1j}^2)^2}, \quad j = s, s + s_*,$$

$$(5.17) \quad B_j = \frac{\Lambda_{1j}}{(1 - \Lambda_{1j}^2)^2}, \quad j = e, s, s + s_*.$$

In particular, \widehat{M} is contained in $z_{p+e} = \bar{z}_{p+e}$ and $z_{p+s}\tilde{\Lambda}_{1s}^2(z'') = \bar{z}_{p+s+s_*}$, where $\zeta_j = z_{p+j}\tilde{\Lambda}_{1j}(z'')$, $1 \leq j \leq p$, is the inverse mapping of $z_{p+j} = \zeta_j\Lambda_{1j}(\zeta)$, $1 \leq j \leq p$.

Proof. We use a realization which is different from (2.17). We assume that M already has the normal form as in Theorem 5.3. Thus for $j = 1, \dots, p$, we have

$$(5.18) \quad \tau_{1j}: \xi'_j = \Lambda_{1j}(\xi\eta)\eta_j, \quad \eta'_j = \Lambda_{1j}^{-1}(\xi\eta)\xi_j, \quad (\xi'_k, \eta'_k) = (\xi_k, \eta_k), \quad k \neq j.$$

Let us define

$$f_j(\xi, \eta) = \xi_j + \xi_j \circ \tau_{1j}, \quad g_j = \overline{f_j \circ \rho}, \quad 1 \leq j \leq p.$$

The latter implies that the biholomorphic mapping $\varphi(\xi, \eta) = (f(\xi, \eta), g(\xi, \eta))$ transforms ρ into the standard complex conjugation $(z', w') \rightarrow (\bar{w}', \bar{z}')$. Define

$$F_j(\xi, \eta) = \xi_j \circ \tau_{1j}(\xi, \eta)\xi_j, \quad 1 \leq j \leq p.$$

Using the expressions of τ_{1j} given by (5.18), we verify that f_j and F_j are invariant by τ_{1k} . Note that the linear part of $f_j(\xi, \eta)$ is $\xi_j + \lambda_j\eta_j$ for $1 \leq j \leq p$, and the quadratic part of $F_j(\xi, \eta)$ is $\lambda_j\xi_j^2$. By Lemma 2.5, f_1, \dots, f_p and F_1, \dots, F_p generate all invariant functions of $\{\tau_{11}, \dots, \tau_{1p}\}$.

Using $\Lambda_{1e} \circ \rho_z = \Lambda_{1e}$ and $\overline{\Lambda_{1s} \circ \rho_z} = \Lambda_{1(s+s_*)}^{-1}$, rewrite $z_j = f_j(\xi, \eta)$, $w_j = g_j(\xi, \eta)$ as

$$\begin{aligned} \xi_e &= \frac{z_e - \Lambda_{1e}(\xi\eta)w_e}{1 - \Lambda_{1e}^2}, & \eta_e &= \frac{w_e - \Lambda_{1e}(\xi\eta)z_e}{1 - \Lambda_{1e}^2}, \\ \xi_s &= \frac{z_s - \Lambda_{1s}^2(\xi\eta)w_{s+s_*}}{1 - \Lambda_{1s}^2(\xi\eta)}, & \eta_s &= \frac{\Lambda_{1s}(\xi\eta)(w_{s+s_*} - z_s)}{1 - \Lambda_{1s}^2(\xi\eta)}, \\ \xi_{s+s_*} &= \frac{z_{s+s_*} - \Lambda_{1(s+s_*)}^2(\xi\eta)w_s}{1 - \Lambda_{1(s+s_*)}^2(\xi\eta)}, & \eta_{s+s_*} &= \frac{\Lambda_{1(s+s_*)}(\xi\eta)(w_s - z_{s+s_*})}{1 - \Lambda_{1(s+s_*)}^2(\xi\eta)}. \end{aligned}$$

Using the above formulas and $w_j = \bar{z}_j$, we compute $\zeta_j = \xi_j\eta_j$ to obtain (5.13)-(5.15).

Note that $F_j(\xi, \eta) = \zeta_j\Lambda_{1j}(\zeta)$. This shows that $z_{p+j} = F_j \circ \varphi^{-1}(z', \bar{z}')$ have the form (5.12). Again, we use the formula of τ_{1k} to verify that $z = (z', z'')$ are invariant by all $\varphi\tau_{1k}\varphi^{-1}$. On the other hand, $z = (z', z'')$ generate invariant functions of the deck transformations of π_1 for the complexification of \widehat{M} given by (5.12). This shows that $\{\varphi\tau_{11}\varphi^{-1}, \dots, \varphi\tau_{1p}\varphi^{-1}\}$ and the deck transformations of π_1 , of which each family consists of commuting involutions,

have the same invariant functions. By Lemma 2.5, we know that the two families must be identical. This shows that (5.12) is a realization for $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$.

To verify the last assertion of the theorem, we first note that by (5.6) the solutions ζ_1, \dots, ζ_p to (5.13)-(5.15) satisfy $\zeta_e = \bar{\zeta}_e$ and $\zeta_{s+s_*} = \bar{\zeta}_s$. By (5.6), on M we have $\bar{z}_{p+e} = \Lambda_{1e}(\zeta)\bar{\zeta}_e = z_{p+e}$. Also $\bar{z}_{p+s+s_*} = \Lambda_{1(s+s_*)}(\zeta)\bar{\zeta}_{s+s_*} = \Lambda_{1s}^{-1}(\zeta)\zeta_s = \Lambda_{1s}^{-2}(\zeta)z_{p+s}$. From (5.12), solve $\zeta_j = \tilde{\Lambda}_{1j}(z'')z_{p+j}$ ($1 \leq j \leq p$) to obtain $\Lambda_{1j}^{-1}(\zeta) = \tilde{\Lambda}_{1j}(z'')$. The proof is complete. \square

5.2. Hull of holomorphy for the abelian CR singularity. Let X be a subset of \mathbf{C}^n . We define the hull of holomorphy of X , denoted by $\mathcal{H}(X)$, to be the intersection of domains of holomorphy in \mathbf{C}^n that contain X . Let B_r^n be the ball in \mathbf{C}^n of radius r and centered at the origin.

By Theorem 5.4, we assume that M has pure elliptic type and it is equivalent to

$$M: z_{p+j} = \Lambda_{1j}(\zeta)\zeta_j, \quad 1 \leq j \leq p,$$

where $\zeta_j = \zeta_j(z')$ ($j = 1, \dots, p$) are the convergent real-valued solutions to (5.13). For $\zeta \in \mathbf{R}^p$ with small $|\zeta|$, we know that $\Lambda_{1j}(\zeta) > 1$.

Near the origin in \mathbf{R}^p , we define a real analytic diffeomorphism:

$$R: \zeta \rightarrow (\Lambda_{11}(\zeta)\zeta_1, \dots, \Lambda_{1p}(\zeta)\zeta_p).$$

If ϵ is small enough, for each $x'' \in [0, \epsilon]^p$, we can define $\zeta = R^{-1}(x'')$. Note that R sends $\zeta_j = 0$ into $x_{p+j} = 0$ for each j . We can write

$$R^{-1}(x'') = (x_{p+1}S_1(x''), \dots, x_{2p}S_p(x''))$$

with $S_j(0) > 0$. Then $M \cap \{z'' = x''\}$ is given by (5.12)-(5.15). For $x'' \in [0, \epsilon]^p$ let $D_j(x'')$ be the compact set in the complex plane whose boundary is defined by the j th equation in (5.12)-(5.15) where $\zeta = R^{-1}(x'')$. When $x_{p+j} > 0$, the boundary of $D_j(x'')$ is an ellipse with

$$(5.19) \quad D_j(x'') \subset B_{C_1\sqrt{x_{p+j}}}^1.$$

Here and in what follows constants will depend only on $\lambda_1, \dots, \lambda_p$. Thus

$$D(x'') := D_1(x'') \times \dots \times D_p(x'') \times \{x''\} \subset \mathbf{C}^p \times \mathbf{R}^p$$

is a product of ellipses and its dimension equals the number of positive numbers among x_{p+1}, \dots, x_{2p} . We will call $D(x'')$ an *analytic polydisc* and

$$\partial^* D(x'') := \partial D_1(x'') \times \dots \times \partial D_p(x'') \times \{x''\}$$

its distinguished boundary which is contained in M . Set $D(0'') = \partial^* D(0'') = \{0\}$. Thus, M is foliated by $\partial^* D(x'')$ as x'' vary in $[0, \epsilon]^p$ and ϵ is sufficiently small. Specifying the ϵ later, we will use this foliation and Hartogs' figures in analytic polydiscs to find the local hull of holomorphy of M at the origin.

As x'' vary in $[0, \epsilon]^p$, let M_ϵ be the union of $\partial^* D(x'')$, and \mathcal{H}_ϵ the union of $D(x'')$. Both \mathcal{H}_ϵ and M_ϵ are compact subsets in \mathbf{C}^{2p} . Note that

$$B_{\epsilon_*}^{2p} + M_\epsilon := \{a + b: a \in B_{\epsilon_*}^{2p}, b \in M_\epsilon\}$$

is contained in a given neighborhood of M_ϵ , if ϵ_* is sufficiently small. Analogously, $B_{\epsilon_*}^{2p} + \mathcal{H}_\epsilon$ is a connected open neighborhood of \mathcal{H}_ϵ . Let us first verify that a function that is holomorphic in a connected neighborhood of M_ϵ in \mathbf{C}^{2p} extends holomorphically to a neighborhood of \mathcal{H}_ϵ such that the extension agrees with the original function on a possibly smaller neighborhood of M_ϵ . Assume that f is holomorphic in a neighborhood \mathcal{U} of $\partial_\epsilon^* D := \cup_{x'' \in [0, \epsilon]^p} \partial^* D(x'')$. Note that \mathcal{H}_ϵ is defined by

$$(5.20) \quad A_j(x'')|z_j|^2 - B_j(x'')(z_j^2 + \bar{z}_j^2) \leq x_{p+j}, \quad 1 \leq j \leq p;$$

$$(5.21) \quad y'' = 0, \quad x'' \in [0, \epsilon]^p$$

with

$$A_j(x'') = \frac{1 + \Lambda_{1j}^2(R^{-1}(x''))}{S_j(x'')(1 - \Lambda_{1j}^2(R^{-1}(x''))^2)}, \quad B_j(x'') = \frac{\Lambda_{1j}(R^{-1}(x''))}{S_j(x'')(1 - \Lambda_{1j}^2(R^{-1}(x''))^2)}.$$

Let δ be a small positive number. For $x'' \in [-\delta, \epsilon]^p$, let $D_j^\delta(x'') \subset \mathbf{C}$ be defined by

$$A_j(x'')|z_j|^2 - B_j(x'')(z_j^2 + \bar{z}_j^2) \leq x_{p+j} + \delta.$$

Let P_ϵ^δ (resp. $\partial^* P_\epsilon^\delta$) be the set of $z = (z', z'')$ such that $y'' \in [-\delta, \delta]^p$, $x'' \in [-\delta, \epsilon]^p$, and $z_j \in D_j^\delta(x'')$ (resp. $z_j \in \partial D_j^\delta(x'')$) for $1 \leq j \leq p$. Let $\mathcal{U}_\epsilon^\delta$ (resp. $\mathcal{U}_\epsilon^{\delta_1}$) be a small neighborhood of P_ϵ^δ (resp. $P_\epsilon^{\delta_1}$). Assume that $0 < \delta_1 < \delta$ and δ_1 is sufficiently small. We may also assume that $\mathcal{U}_\epsilon^{\delta_1}$ is contained in $\mathcal{U}_\epsilon^\delta$ and $\partial^* P_\epsilon^\delta \subset \mathcal{U}$. Thus, for $(z', z'') \in \mathcal{U}_\epsilon^{\delta_1}$, we can define

$$(5.22) \quad F(z', z'') = \frac{1}{(2\pi i)^p} \int_{\zeta_1 \in \partial D_1^\delta(x'')} \cdots \int_{\zeta_p \in \partial D_p^\delta(x'')} \frac{f(\zeta, z'') d\zeta_1 \cdots d\zeta_p}{(\zeta_1 - z_1) \cdots (\zeta_p - z_p)}.$$

When z is sufficiently small, $F(z) = f(z)$ as f is holomorphic near the origin. Fix $z_0 \in \mathcal{U}_\epsilon^{\delta_1}$. We want to show that F is holomorphic at z_0 . So F is a desired extension of f . By continuity, when $z = (z_1, \dots, z_{2p})$ tends to z_0 , x'' tends to x_0'' and $\partial D_j^\delta(x'')$ tends to $\partial D_j^\delta(x_0'')$, while $z_j \in D_j^\delta(x_0'')$ when z is sufficiently close to z_0 . By Cauchy theorem, for z sufficiently close to z_0 we change the repeated integral for $\zeta_j \in \partial D_j^\delta(x_0'')$, $1 \leq j \leq p$. The domain of integration is thus fixed. The integrand is holomorphic in z . Hence F is holomorphic at $z = z_0$.

Next we want to show that \mathcal{H}_ϵ is the hull of holomorphy of M_ϵ in $B_{\epsilon_0}^{2p}$ for suitable ϵ, ϵ_0 that can be arbitrarily small.

Let us first show that \mathcal{H}_ϵ is the intersection of domains of holomorphy in \mathbf{C}^{2p} . Recall that \mathcal{H}_ϵ is defined by (5.20)-(5.21). Define for $\delta' := (\delta_1, \dots, \delta_p)$ with $\delta_j > 0$

$$\begin{aligned} \rho_j^{\delta'} &= A_j(x'')|z_j|^2 - B_j(x'')(z_j^2 + \bar{z}_j^2) - x_{p+j} + (\delta_1^{-1} + \cdots + \delta_p^{-1}) \sum_{i=1}^p y_{p+i}^2 \\ &+ \sum_{i \neq j} \delta_i^{-1} \{A_i(x'')|z_i|^2 - B_i(x'')(z_i^2 + \bar{z}_i^2) - x_{p+i}\}. \end{aligned}$$

When $p = 1$, the last summation is 0. The complex Hessian of $\rho_j^{\delta'}$ is

$$\begin{aligned} \sum_{\alpha, \beta=1}^{2p} \frac{\partial^2 \rho_j^{\delta'}}{\partial z_\alpha \bar{z}_\beta} t_\alpha \bar{t}_\beta &= A_j(x'') |t_j|^2 + \frac{\delta_1^{-1} + \cdots + \delta_p^{-1}}{2} \sum_i |t_{p+i}|^2 + \sum_{i \neq j} \frac{1}{\delta_i} A_i(x'') |t_i|^2 \\ &+ \operatorname{Re} \sum_k a_{jk}(x''; z_j) t_j \bar{t}_{p+k} + \sum_{k, \ell} b_{j, k\ell}(x''; z_j) t_{p+k} \bar{t}_{p+\ell} \\ &+ \operatorname{Re} \sum_{i \neq j} \sum_k \frac{1}{\delta_i} c_{j, ik}(x''; z_i) t_i \bar{t}_{p+k} + \sum_{i \neq j} \sum_{k, \ell} \frac{1}{\delta_i} d_{j, k\ell}(x''; z_i) t_{p+k} \bar{t}_{p+\ell}. \end{aligned}$$

Here $a_{jk}(x''; 0) = b_{j, k\ell}(x''; 0) = c_{j, ik}(x''; 0) = d_{j, k\ell}(x''; 0) = 0$, and i, j, k, ℓ are in $\{1, \dots, p\}$. From the Cauchy-Schwarz inequality, it follows that for $z \in B_{\epsilon_0}^{2p}$ with $\epsilon_0 > 0$ sufficiently small and $0 < \delta_j < 1$,

$$2 \sum_{\alpha, \beta=1}^{2p} \frac{\partial^2 \rho_j^{\delta'}}{\partial z_\alpha \bar{z}_\beta} t_\alpha \bar{t}_\beta \geq A_j(x'') |t_j|^2 + \frac{\delta_1^{-1} + \cdots + \delta_p^{-1}}{2} \sum_j |t_{p+j}|^2 + \sum_{i \neq j} \delta_i^{-1} A_i(x'') |t_i|^2.$$

Therefore, each $\rho_j^{\delta'}$ is strictly plurisubharmonic on $|z| < \epsilon_0$ for all $0 < \delta_i < 1$. Hence for $\delta_* = (\delta_0, \dots, \delta_p) = (\delta_0, \delta') \in (0, 1)^{p+1}$,

$$\rho_\epsilon^{\delta_*}(z) = \max_j \{ \rho_j^{\delta'}(z), |y''|^2 - \delta_0^2, x_{p+j}^2 - \epsilon^2 \}$$

is plurisubharmonic on $B_{\epsilon_0}^{2p}$. By (5.19), $D(x'')$ is contained in $B_{C_2 \epsilon^{1/2}}^{2p}$ for $x'' \in [0, \epsilon]^p$. We now fix $\epsilon < (\epsilon_0/C_2)^2$ to ensure

$$(5.23) \quad D(x'') \subset B_{\epsilon_0}^{2p}, \quad \forall x'' \in [0, \epsilon]^p.$$

This shows that $\mathcal{H}_\epsilon^{\delta_*} := \{z \in B_{\epsilon_0}^{2p} : \rho_\epsilon^{\delta_*}(z) < 0\}$ is a domain of holomorphy.

Let us verify that $\mathcal{H}_\epsilon = \bigcap_{\epsilon' > \epsilon, \delta_0 > 0, \dots, \delta_p > 0} \mathcal{H}_{\epsilon'}^{\delta_*}$. Fix $z \in \mathcal{H}_\epsilon$. From (5.23) we get $z \in B_{\epsilon_0}^{2p}$. We have $y'' = 0$. Hence (5.20) hold and $x_{p+j}^2 \leq \epsilon^2$. Clearly, $\rho_j^{\delta'}(z) < 0$ for each j and $\delta \in (0, 1)^p$. This shows that $z \in \mathcal{H}_\epsilon$ is in the intersection. For the other inclusion, let us assume that z is in the intersection. Then $y'' = 0$. With $\rho_j^{\delta'}(z) < 0$, we let δ_i tend to 0 for $i \neq j$. We conclude

$$A_i(x'') |z_i|^2 - B_i(x'')(z_i^2 + \bar{z}_i^2) \leq x_{p+i}$$

for all $i \neq j$, and hence for all i as $p > 1$. When $p = 1$ the above inequality can be obtained directly from $\rho_1^{\delta_*}$. We also see that $0 \leq x_{p+j} \leq \epsilon$. We have verified (5.20) and (5.21). This shows that $z \in \mathcal{H}_\epsilon$.

In view of (5.20)-(5.21), the boundary of \mathcal{H}_ϵ is the union $\bigcup_{j=1}^p \mathcal{H}_j^\epsilon$ with \mathcal{H}_j^ϵ being defined by

$$\begin{aligned} A_j(x'') |z_j|^2 - B_j(x'')(z_j^2 + \bar{z}_j^2) &= x_{p+j}, \\ A_i(x'') |z_i|^2 - B_i(x'')(z_i^2 + \bar{z}_i^2) &\leq x_{p+i}, \quad 1 \leq i \leq p, \quad i \neq j; \\ y'' = 0, \quad x_{p+i} &\leq \epsilon, \quad 1 \leq i \leq p. \end{aligned}$$

Therefore, we have proved the following theorem.

Theorem 5.5. *Let M be a germ of real analytic submanifold at an abelian CR singularity. Assume that the complex tangent of M is purely elliptic at the origin. There is a base of neighborhoods $\{U_j\}$ of the origin in \mathbf{C}^n which satisfies the following: For each U_j , the hull of holomorphy $H(M \cap U_j)$ of $M \cap U_j$ is foliated by embedded complex submanifolds with boundaries. Furthermore, near the origin $H(M \cap U_j)$ is the transversal intersection of p real analytic submanifolds of dimension $3p$ with boundary. The boundary of $H(M \cap U_j)$ contains $M \cap U_j$; and two sets are the same if and only if $p = 1$.*

Remark 5.6. The proof shows that the hull $H(M \cap U_j)$ is foliated by analytic polydiscs, i.e. holomorphic embeddings of closed unit polydisc in \mathbf{C}^p .

6. RIGIDITY OF PRODUCT QUADRICS

The aim of this section is to prove the following rigidity theorem: Let us consider a higher order analytic perturbation of a product quadric. If this manifold is formally equivalent to the product quadric, then under a small divisors condition, it is also holomorphically equivalent to it. Notice that when $p > 1$, there are real submanifold M with a linearizable σ such that M is not formally equivalent to the quadric, or equivalently, the $\{\tau_{1j}, \rho\}$ is not formally linearizable; see [GS15].

The proof goes as follows : Since the manifold is formally equivalent to the quadric, the associated involutions $\{\tau_{1i}\}$ and $\{\tau_{2i}\}$ are simultaneously linearizable by a formal biholomorphism that commutes with ρ . In particular, $\sigma_1, \dots, \sigma_p$, as defined by (5.1) and (5.2), are formally linearizable and they commute pairwise. These are germs of biholomorphisms with a diagonal linear part. According to [Sto15][theorem 2.1], this abelian family can be holomorphically linearized under a collective Brjuno type condition (7.38). Furthermore, the transformation commutes with ρ . Then, we linearize simultaneously and holomorphically both $\tau_1 := \tau_{11} \circ \dots \circ \tau_{1p}$ and $\tau_2 := \tau_{21} \circ \dots \circ \tau_{2p}$ by a transformation that commutes with both ρ and \mathcal{S} , the family of linear parts of the $\sigma_1, \dots, \sigma_p$. Finally, we linearize simultaneously and holomorphically both families $\{\tau_{1i}\}$ and $\{\tau_{2i}\}$ by a transformation that commutes with ρ , \mathcal{S} , T_1 and T_2 .

These last two steps will be obtained through a majorant method and the application of a holomorphic implicit function theorem. This is obtained in Proposition 6.6. They first require a complete description of the various centralizers and their associated normalized mappings, i.e. suitable complements. This is a goal of Proposition 6.3.

Throughout this section, we do not assume that μ_1, \dots, μ_p are non resonant in the sense that $\mu^Q \neq 1$ if $Q \in \mathbf{Z}^p$ and $Q \neq 0$. In fact, we will apply our results to M which might be resonant. However, we will retain the assumption that σ has distinct eigenvalues when we apply the results to the manifolds.

6.1. Centralizers. We recall from (5.1) and (5.2), the definition and property of germs of holomorphic diffeomorphisms : $\sigma_i := \tau_{1i} \circ \tau_{2i}$, $\sigma_i^{-1} = \rho \sigma_i \rho$, $1 \leq i \leq e_* + h_*$; $\sigma_s := \tau_{1s} \circ \tau_{2(s_* + s)}$, $\sigma_{s+s_*} := \tau_{1(s+s_*)} \circ \tau_{2s}$, $\sigma_{s+s_*}^{-1} = \rho \sigma_s \rho$, $e_* + h_* < s \leq p - s_*$. Recall the linear maps

$$(6.1) \quad \begin{aligned} S: \xi'_j &= \mu_j \xi_j, & \eta'_j &= \mu_j^{-1} \eta_j; \\ S_j: \xi'_j &= \mu_j \xi_j, & \eta'_j &= \mu_j^{-1} \eta_j, & \xi'_k &= \xi_k, & \eta'_k &= \eta_k, & k &\neq j; \end{aligned}$$

$$(6.2) \quad T_{ij}: \xi'_j = \lambda_{ij} \eta_j, \quad \eta'_j = \lambda_{ij}^{-1} \xi_j, \quad \xi'_k = \xi_k, \quad \eta'_k = \eta_k, \quad k \neq j;$$

while ρ is given by

$$(6.3) \quad \rho: \begin{cases} (\xi'_e, \eta'_e, \xi'_h, \eta'_h) = (\bar{\eta}_e, \bar{\xi}_e, \bar{\xi}_h, \bar{\eta}_h), \\ (\xi'_s, \xi'_{s+s_*}, \eta'_s, \eta'_{s+s_*}) = (\bar{\xi}_{s+s_*}, \bar{\xi}_s, \bar{\eta}_{s+s_*}, \bar{\eta}_s). \end{cases}$$

Recall that indices e, h, s have the ranges $1 \leq e \leq e_*$, $e_* < h \leq e_* + h_*$, and $e_* + h_* < s \leq p - s_*$. The basic conditions on $\mu_j = \lambda_j^2$ are the following:

$$(6.4) \quad |\mu_h| = 1, \quad \mu_{s+s_*} = \bar{\mu}_s^{-1}, \quad \mu_e > 1, \quad |\mu_s| \geq 1, \quad \mu_s^k \neq 1, k = 1, 2, \dots$$

In particular, a μ_j may be repeated and μ_h can be 1.

We need to introduce notation for the indices to describe various centralizers regarding T_{1j}, S_j and ρ . We first introduce index sets for the centralizer of $\mathcal{S} := \{S_1, \dots, S_p\}$, $T_1 := T_{i1} \circ \dots \circ T_{ip}$, ρ . We recall that $\mathcal{T}_i := \{T_{i1}, \dots, T_{ip}\}$.

Let $(P, Q) \in \mathbf{N}^p \times \mathbf{N}^p$ and $1 \leq j \leq p$. By definition, $\xi^P \eta^Q e_j$ belongs to the centralizer of \mathcal{S} if and only if it commutes with each S_i . In other words, $\xi^P \eta^Q e_j \in \mathcal{C}(\mathcal{S})$ if and only if

$$(6.5) \quad \mu_k^{p_k - q_k} = 1, \quad \forall k \neq j; \quad \mu_j^{p_j - q_j} = \mu_j.$$

Note that the same condition holds for $\xi^Q \eta^P e_{p+j}$ to belong to $\mathcal{C}(\mathcal{S})$. This leads us to define the set of multiindices

$$\mathcal{R}_j := \{(P, Q) \in \mathbf{N}^{2p} : \mu_j^{p_j - q_j} = \mu_j, \mu_i^{p_i - q_i} = 1, \forall i \neq j\}, \quad 1 \leq j \leq p.$$

We observe that if $(P, Q) \in \mathcal{R}_j$, then (6.4) implies that

$$(6.6) \quad \begin{aligned} p_j &= q_j + 1, \quad j \neq h; \quad p_i = q_i, \quad \forall i \neq j, h; \\ \lambda_h^{p_h - q_h} &= \pm 1, \quad h \neq j; \quad \lambda_j^{p_j - q_j - 1} = \pm 1, \quad j = h. \end{aligned}$$

Here we have used the assumption that μ_s are not root of unity, which simplifies greatly the results and computation in this section.

For convenience, we define for $P = (p_e, p_h, p_s, p_{s+s_*})$ and $Q = (q_e, q_h, q_s, q_{s+s_*})$

$$(6.7) \quad \begin{aligned} \rho(PQ) &:= (q_e, p_h, p_{s+s_*}, p_s, p_e, q_h, q_{s+s_*}, q_s), \\ \rho_a(PQ) &:= (q_e, p_h, p_{s+s_*}, p_s), \quad \rho_b(PQ) := (p_e, q_h, q_{s+s_*}, q_s), \end{aligned}$$

$$(6.8) \quad \bar{f}_{\rho(PQ)} := (\bar{f} \circ \rho)_{PQ}.$$

Here $p_e = (p_1, \dots, p_{e_*})$ denotes the ‘‘elliptic coordinates’’ of P . Hence,

$$(6.9) \quad \rho(PQ) = (\rho_a(PQ), \rho_b(PQ)) = (\rho_b(QP), \rho_a(QP)).$$

According to (6.5) and equation (6.3) of ρ , the restriction of ρ to \mathcal{R}_h is an involution, which will be denoted by ρ_h . Moreover, ρ is a bijection ρ_s from \mathcal{R}_s onto \mathcal{R}_{s+s_*} . We define an involution on \mathcal{R}_e by

$$(6.10) \quad \rho_e(PQ) := (\rho_b(PQ), \rho_a(PQ)) = (p_e, q_h, q_{s+s_*}, q_s, q_e, p_h, p_h, p_{s+s_*}, p_s).$$

Note that ρ_e is not a restriction of ρ , and ρ_s is not an involution either.

Next, we introduce sets of indices to be used to compute the centralizers on $\mathcal{T}_1, \mathcal{T}_2, \rho$. Set

$$\mathcal{N}_j := \mathcal{R}_j \cap \{(P, Q) : p_i \geq q_i, \quad \forall i \neq j\}, \quad 1 \leq j \leq p.$$

Since there is no restriction for $p = 1$, we have $\mathcal{N}_j = \mathcal{R}_j$ for $j = e$ or h . Let us set

$$\begin{aligned} A_{jk}(P, Q) &:= \max\{p_k, q_k\}, & k \neq j, & \quad A_{jj}(P, Q) = p_j; \\ B_{jk}(P, Q) &:= \min\{p_k, q_k\}, & k \neq j, & \quad B_{jj}(P, Q) = q_j. \end{aligned}$$

We define a mapping

$$(A_j, B_j): \mathcal{R}_j \rightarrow \mathcal{N}_j$$

with $A_j := (A_{j1}, \dots, A_{jp})$ and $B_j := (B_{j1}, \dots, B_{jp})$. For $(P, Q) \in \mathcal{N}_j$ with $j = e, h$, we have $A_j \circ \rho_j(P, Q) = (p_e, p_h, p_{s+s_*}, p_s)$ and $B_j \circ \rho_j(P, Q) = (q_e, q_h, q_{s+s_*}, q_s)$. In other words, on \mathcal{N}_j for $j = e$ or h , $A_j \circ \rho_j$ just interchanges the s th and the $(s + s_*)$ th coordinates for each s , so does $B_j \circ \rho_j$, while $A_{s+s_*}\rho_s$ and $B_{s+s_*}\rho_s$ have the same property on \mathcal{N}_s . Furthermore,

$$(6.11) \quad (A_h, B_h)\rho = \rho(A_h, B_h) \quad \text{on } \mathcal{R}_h,$$

$$(6.12) \quad (A_{s+s_*}, B_{s+s_*})\rho = \rho(A_s, B_s) \quad \text{on } \mathcal{R}_s.$$

Finally, with the convention that the product over an empty set is 1, we define, for $(P, Q) \in \mathcal{R}_j$:

$$(6.13) \quad \nu_{PQ} := \begin{cases} \prod_{h'} \lambda_{h'}^{q_{h'} - p_{h'}}, & j \neq h, \\ \lambda_h^{p_h - q_h - 1} \prod_{h' \neq h} \lambda_{h'}^{q_{h'} - p_{h'}}, & j = h; \end{cases}$$

$$(6.14) \quad \nu_{PQ}^+ := \begin{cases} \prod_{h' | q_{h'} > p_{h'}} \lambda_{h'}^{q_{h'} - p_{h'}}, & j \neq h; \\ \prod_{h' \neq h, q_{h'} > p_{h'}} \lambda_{h'}^{q_{h'} - p_{h'}}, & j = h. \end{cases}$$

Here $e_* < h', h \leq e_* + h_*$. For convenience, we however define

$$\nu_{QP} := \nu_{PQ}, \quad (P, Q) \in \mathcal{R}_j.$$

If $p = 1$ we set $\nu_{PQ}^+ = 1$.

Lemma 6.1. *Let $(P, Q) \in \mathcal{R}_j$. Then $\lambda_j^{-1} \lambda^{P-Q} = \nu_{PQ}$, and*

$$(6.15) \quad \nu_{PQ} = \pm 1; \quad \nu_{PQ}^+ = \pm 1; \quad \nu_{PQ}^+ = 1, \quad (P, Q) \in \mathcal{N}_j;$$

$$(6.16) \quad \nu_{\rho_e(PQ)} = \nu_{PQ}, \quad j = e; \quad \nu_{\rho(PQ)} = \nu_{PQ};$$

$$(6.17) \quad \nu_{\rho_e(PQ)}^+ = \nu_{PQ}^+ \nu_{PQ}, \quad j = e; \quad \nu_{\rho(PQ)}^+ = \nu_{PQ}^+.$$

Proof. The first identity follows from the definition of ν_{PQ} and \mathcal{R}_j . From the definition of \mathcal{R}_j , we have $(\lambda_i^{p_i - q_i})^2 = \mu_i^{p_i - q_i} = 1$ for $i = h'$ in (6.13)-(6.14). We also have $\mu_h^{p_h - q_h - 1} = 1$ for terms in (6.13)-(6.14). Thus

$$\lambda_{h'}^{p_{h'} - q_{h'}} = \pm 1, \quad \lambda_h^{p_h - q_h - 1} = \pm 1.$$

Thus we obtain (6.15); the rest identities follow from the definition of ρ_e, ρ , and the above identities. \square

Lemma 6.2. For all multiindices (P, Q) , $\xi^P \eta^Q \circ \rho = \bar{\xi}^{\rho_a(PQ)} \bar{\eta}^{\rho_b(PQ)}$. For all multiindices $(P, Q) \in \mathcal{R}_e \cup \mathcal{R}_h$, we have

$$(6.18) \quad \overline{\lambda^{\rho_a(P,Q) - \rho_b(P,Q)}} = \lambda^{Q-P}, \quad \overline{\mu^{\rho_b - \rho_a}} = \mu^{P-Q},$$

$$(6.19) \quad \xi^P \eta^Q \circ \rho \circ T_1 = \lambda^{Q-P} \bar{\xi}^{\rho_b(PQ)} \bar{\eta}^{\rho_a(PQ)},$$

$$(6.20) \quad \xi^P \eta^Q \circ \rho \circ S^{-1} = \mu^{P-Q} \bar{\xi}^{\rho_a(PQ)} \bar{\eta}^{\rho_b(PQ)}.$$

Proof. Identity (6.18) follows from (6.7) and the fact that λ_e and μ_e are reals, $\lambda_h^{-1} = \bar{\lambda}_h$, $p_s = q_s$, and $p_{s+s_*} = q_{s+s_*}$. A direct computation shows that

$$\xi^P \eta^Q \circ \rho \circ T_1 = \bar{\lambda}^{\rho_a - \rho_b} \bar{\xi}^{\rho_b(PQ)} \bar{\eta}^{\rho_a(PQ)}, \quad \xi^P \eta^Q \circ \rho \circ S^{-1} = \bar{\mu}^{\rho_b - \rho_a} \bar{\xi}^{\rho_a(PQ)} \bar{\eta}^{\rho_b(PQ)}.$$

The result follows from (6.18). \square

Finally we note that

$$(6.21) \quad \iota_e: (P, Q) \rightarrow (A_e, B_e) \circ \rho_e(PQ) = (A_e, B_e)(\rho_b(P, Q), \rho_a(P, Q))$$

defines an involution on \mathcal{N}_e . We now can describe the centralizers.

Proposition 6.3. Let $\mathcal{S} = \{S_1, \dots, S_p\}$, $\mathcal{T}_i = \{T_{i1}, \dots, T_{ip}\}$ and ρ be given by (6.1)-(6.3). Let $\varphi = I + (U, V)$ be a formal biholomorphic map that is tangent to the identity.

(i) $\varphi \in \mathcal{C}(\mathcal{S})$ if and only if

$$(6.22) \quad U_{j,PQ} = 0 = V_{j,QP}, \quad \forall (P, Q) \notin \mathcal{R}_j.$$

Also, $\varphi \in \mathcal{C}(\mathcal{S}, \rho)$ if and only if (6.22) holds and

$$(6.23) \quad U_{h,PQ} = \bar{U}_{h,\rho(PQ)}, \quad (P, Q) \in \mathcal{R}_h; \quad U_{s+s_*,PQ} = \bar{U}_{s,\rho(PQ)}, \quad (P, Q) \in \mathcal{R}_{s+s_*};$$

$$(6.24) \quad V_{e,QP} = \bar{V}_{e,\rho_e(PQ)}, \quad (P, Q) \in \mathcal{R}_e;$$

$$(6.25) \quad V_{h,QP} = \bar{V}_{h,\rho(QP)}, \quad (P, Q) \in \mathcal{R}_h; \quad V_{s+s_*,QP} = \bar{V}_{s,\rho(QP)}, \quad (P, Q) \in \mathcal{R}_{s+s_*}.$$

(ii) $\varphi \in \mathcal{C}(\mathcal{S}, T_1)$ if and only if (6.22) holds and

$$(6.26) \quad V_{j,QP} = \nu_{PQ} U_{j,PQ}, \quad \forall (P, Q) \in \mathcal{R}_j.$$

(iii) $\varphi \in \mathcal{C}(\mathcal{S}, T_1, \rho)$ if and only if (6.22), (6.23) and (6.26) hold, and

$$(6.27) \quad U_{e,PQ} = \nu_{PQ} \bar{U}_{e,\rho_e(PQ)}, \quad (P, Q) \in \mathcal{R}_e.$$

(iv) Let $p > 1$. $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2)$ if and only if (6.22) and (6.26) hold, and

$$(6.28) \quad U_{j,PQ} = \nu_{PQ}^+ U_{j,(A_j, B_j)(P, Q)}, \quad (P, Q) \in \mathcal{R}_j \setminus \mathcal{N}_j.$$

Also, $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2, \rho)$ if and only if (6.22), (6.26) and (6.28) hold, and

$$(6.29) \quad U_{e,PQ} = \bar{U}_{e,(A_e, B_e)\rho_e(PQ)}, \quad (P, Q) \in \mathcal{N}_e,$$

$$(6.30) \quad U_{h,PQ} = \bar{U}_{h,\rho(PQ)}, \quad (P, Q) \in \mathcal{N}_h,$$

$$(6.31) \quad U_{s+s_*,PQ} = \bar{U}_{s,\rho(PQ)}, \quad (P, Q) \in \mathcal{N}_{s+s_*}.$$

We remark that condition (6.28) holds trivially when $(P, Q) \in \mathcal{N}_j$, in which case it becomes $U_{j,PQ} = U_{j,PQ}$.

Proof. To simplify notation, we abbreviate

$$\rho_a = \rho_a(PQ), \quad \rho_b = \rho_b(PQ), \quad A_j = A_j(P, Q), \quad B_j = B_j(P, Q).$$

Recall that $\lambda_e = \bar{\lambda}_e$, $\lambda_h = \bar{\lambda}_h^{-1}$ and $\lambda_{s+s_*} = \bar{\lambda}_s^{-1}$. By definition,

$$S_e = T_{1e}T_{2e}, \quad S_h = T_{1h}T_{2h}, \quad S_s = T_{1s}T_{2(s+s_*)}, \quad S_{s+s_*} = T_{1(s+s_*)}T_{2s}.$$

In the proof, we will use the fact that S_j is reversible by both involutions in the composition for S_j . In particular,

$$(6.32) \quad T_{1j}S_jT_{1j} = S_j^{-1}, \quad \forall j.$$

However, we have $T_{2(s+s_*)}S_sT_{2(s+s_*)} = S_s^{-1}$ and $T_{2s}S_{s+s_*}T_{2s} = S_{s+s_*}^{-1}$. For simplicity, we will derive identities by using (6.32) and

$$(6.33) \quad S_e^{-1} = \rho S_e \rho, \quad S_h^{-1} = \rho S_h \rho, \quad S_{s+s_*}^{-1} = \rho S_{s+s_*} \rho.$$

Finally, we need one more identity. Recall that

$$\begin{aligned} T_{1e}T_{2j} &= T_{2j}T_{1e}, \quad j \neq e; & T_{1h}T_{2j} &= T_{2j}T_{1h}, \quad j \neq h; \\ T_{1s}T_{2j} &= T_{2j}T_{1s}, \quad j \neq s + s_*; & T_{1(s+s_*)}T_{2j} &= T_{2j}T_{1(s+s_*)}, \quad j \neq s. \end{aligned}$$

Therefore, for any j we have the identity

$$(6.34) \quad T_1S_jT_1 = S_j^{-1}.$$

In what follows, we will derive all identities by using (6.32), (6.33) and (6.34), as well as $S_iS_j = S_jS_i$, $T_{1i}T_{1j} = T_{1j}T_{1i}$ and $T_2 = \rho T_1 \rho$.

(i) The centralizer of \mathcal{S} is easy to describe. Namely, $\varphi \in \mathcal{C}(\mathcal{S})$ if and only if

$$\begin{aligned} U_j \circ S_j &= \mu_j U_j, \quad U_j \circ S_k = U_j, \quad k \neq j, \\ V_j \circ S_j &= \mu_j^{-1} V_j, \quad V_j \circ S_k = V_j, \quad k \neq j. \end{aligned}$$

For $\varphi \rho = \rho \varphi$, we need

$$(6.35) \quad U_h = \overline{U_h \circ \rho}, \quad U_{s+s_*} = \overline{U_s \circ \rho},$$

$$(6.36) \quad V_e = \overline{V_e \circ \rho}, \quad V_h = \overline{V_h \circ \rho}, \quad V_{s+s_*} = \overline{V_s \circ \rho}.$$

Hence, using (6.8)-(6.10), we have $\overline{U}_{e,PQ} = V_{e,\rho(PQ)} = V_{e,\rho_e(QP)}$. The other equalities are obtained in the same way.

(ii) If $\varphi \in \mathcal{C}(\mathcal{S}, T_1) \subset \mathcal{C}(S, T_1)$, then it satisfies

$$(6.37) \quad V_j = \lambda_j^{-1} U_j \circ T_1.$$

This implies (6.26).

(iii) Assume furthermore that $\varphi \in \mathcal{C}(\mathcal{S}, T_1, \rho)$. Eliminating V_e from (6.37) and (6.36), we obtain

$$U_e = \lambda_e \overline{U_e \circ \rho \circ T_1}.$$

According to (6.19), we obtain (6.27) by

$$U_{e,\rho_b \rho_a} = \lambda_e \bar{\lambda}^{Q-P} \overline{U}_{e,PQ} = \nu_{PQ} \overline{U}_{e,PQ}.$$

(iv) Let $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2)$. Then, in particular, we have

$$U_j = U_j(T_{1k}), \quad k \neq j; \quad V_j = \lambda_j^{-1} U_j \circ T_1.$$

Let $(P, Q) \in \mathcal{R}_j \setminus \mathcal{N}_j$. We compose U_j successively by each T_{1k} if $q_k > p_k$. We emphasize that such a k is a hyperbolic index. The previous identity yields

$$(6.38) \quad U_{j,PQ} = L_{j,PQ} U_{j,A_j B_j}, \quad L_{j,PQ} := \prod_{k \neq j, p_k < q_k} \lambda_k^{q_k - p_k}.$$

By the definition of ν_{PQ}^+ , we conclude

$$(6.39) \quad L_{j,PQ} = \nu_{PQ}^+, \quad (P, Q) \in \mathcal{R}_j.$$

If $(P, Q) \in \mathcal{N}_j$, then $(A_j, B_j) = (P, Q)$ and we have $L_{j,PQ} = \nu_{PQ}^+ = 1$, so that the relation (6.38) just becomes the identity $U_{j,PQ} = U_{j,PQ}$.

Assume now that $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2, \rho)$. In addition to the previous conditions, we have (6.35) and (6.36). Hence, (6.23), (6.27) and (6.38) lead to:

$$\begin{aligned} \nu_{PQ} \overline{U}_{e, \rho_e(PQ)} &= U_{e, PQ} = L_{e, PQ} U_{e, A_e B_e}, & (P, Q) \in \mathcal{R}_e; \\ \overline{U}_{h, \rho_h(PQ)} &= U_{h, PQ} = L_{h, PQ} U_{h, A_h B_h}, & (P, Q) \in \mathcal{R}_h; \\ \overline{U}_{s+s^*, \rho(PQ)} &= U_{s, PQ} = L_{s, PQ} U_{s, A_s B_s}, & (P, Q) \in \mathcal{R}_s. \end{aligned}$$

Since ρ_e, ρ_h are involutions on \mathcal{R}_e and \mathcal{R}_h , respectively, and since ρ is a bijection from \mathcal{R}_s onto \mathcal{R}_{s+s^*} , we obtain

$$\begin{aligned} \nu_{\rho_e(PQ)} \overline{U}_{e, PQ} &= L_{e, \rho_e(PQ)} U_{e, (A_e, B_e) \circ \rho_e(PQ)}, & (P, Q) \in \mathcal{R}_e; \\ \overline{U}_{h, PQ} &= L_{h, \rho_h(AB)} U_{h, (A_h, B_h) \circ \rho_h(PQ)}, & (P, Q) \in \mathcal{R}_h; \\ \overline{U}_{s+s^*, PQ} &= L_{s, \rho(AB)} U_{s, (A_s, B_s) \circ \rho(PQ)}, & (P, Q) \in \mathcal{R}_{s+s^*}. \end{aligned}$$

By (6.39), we copy the values $L_{j, \rho(PQ)} = \nu_{\rho(PQ)}^+$ from (6.17). We have

$$\begin{aligned} \nu_{\rho_j(PQ)}^+ &= \nu_{PQ}^+, & \text{if } j \neq e, \text{ and } (P, Q) \in \mathcal{R}_j; \\ \nu_{\rho_e(PQ)}^+ &= \nu_{PQ} \nu_{PQ}^+, & \text{if } (P, Q) \in \mathcal{R}_e; \\ \nu_{\rho_e(PQ)} &= \nu_{PQ}, & \text{if } (P, Q) \in \mathcal{R}_e. \end{aligned}$$

Finally, we obtain

$$\begin{aligned} U_{j, PQ} &= \nu_{PQ}^+ \overline{U}_{j, (A_j, B_j) \circ \rho_j(PQ)}, & (P, Q) \in \mathcal{R}_j, \quad j = e, h; \\ U_{s+s^*, PQ} &= \nu_{PQ}^+ \overline{U}_{s, (A_s, B_s) \circ \rho(PQ)}, & (P, Q) \in \mathcal{R}_{s+s^*}. \end{aligned}$$

Therefore, we have derived necessary conditions for the centralizers. Let us verify that the conditions are also sufficient. Of course, the verification for (i) is straightforward. Furthermore, that $\varphi = I + (U, V)$ commutes with S_1, \dots, S_p is equivalent to $U_{j, PQ} = V_{j, QP} = 0$ for all $(P, Q) \in \mathcal{R}_j$, which is also trivial in cases (ii) and (iii).

For (ii), (6.22) and (6.26) imply that φ commutes with T_1 . We verify that φ commutes with ρ . In other words, (6.24) and (6.25) hold. The latter follows immediately from (6.23) and (6.26). For the former, take $(P, Q) \in \mathcal{R}_e$. By (6.26) and (6.27), we get $V_{e, QP} = \nu_{PQ} U_{e, PQ} = \overline{U}_{e, \rho_e(PQ)}$, which is (6.24).

For (iii), let us verify that (6.28), (6.22), and (6.26) are sufficient conditions for $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2)$. By (6.26), we get $\varphi T_1 = T_1 \varphi$. Also, for $\varphi \in \mathcal{C}(\mathcal{T}_1)$ it remains to show that for

$(P, Q) \in \mathcal{R}_j$

$$(6.40) \quad (U_j \circ T_{1k})_{PQ} = U_{j,PQ}, \quad k \neq j; \quad (U_j \circ T_{1j})_{QP} = \lambda_j V_{j,QP}.$$

We introduce (P_j, Q_j) via $\xi^P \eta^Q \circ T_{1j} = \lambda_j^{p_j - q_j} \xi^{P_j} \eta^{Q_j}$ and also denote (P_j, Q_j) by $(P, Q)_j$. We remark that (6.28) also holds for $(P, Q) \in \mathcal{N}_j$. Therefore, we will use (6.28) for all $(P, Q) \in \mathcal{R}_j$.

For $k \neq j, h$, we have $(P_k, Q_k) = (P, Q)$. Thus in this case we immediately get the first identity in (6.40). Using (6.28) twice, we obtain for $j \neq h$

$$\begin{aligned} (U_j \circ T_{1h})_{PQ} &= \lambda_h^{p_h - q_h} U_{j,(PQ)_h} = \lambda_h^{p_h - q_h} \nu_{(PQ)_h}^+ U_{j,(A_j, B_j)(P, Q)} \\ &= \lambda_h^{p_h - q_h} \nu_{(PQ)_h}^+ \bar{\nu}_{PQ}^+ U_{j,PQ} = U_{j,PQ}. \end{aligned}$$

Combining with the identities which we have proved, we get $(U_j \circ T_{1j})_{QP} = (U_j \circ T_1)_{QP} = (\lambda_j V_j)_{QP}$ for $j \neq h$. This gives us all the identities in (6.40) for $(P, Q) \in \mathcal{R}_j$. These identities are trivial when (P, Q) is not in \mathcal{R}_j . Therefore, we have shown that these conditions are sufficient for $\varphi \in \mathcal{C}(\mathcal{T}_1, \mathcal{T}_2)$.

Finally, we need to verify that (6.22), (6.26), and (6.28)-(6.31) imply that φ and ρ commute. In other words, we need to verify (6.23) and (6.27), by (iii). We have

$$(6.41) \quad (A_e B_e) \circ \rho_e \circ (A_e B_e) = (A_e B_e) \circ \rho_e \quad \text{on } \mathcal{R}_e.$$

Let $(P, Q) \in \mathcal{R}_e$. By (6.28), (6.29), (6.41) and (6.28), we get

$$U_{e,PQ} = \nu_{PQ}^+ U_{e,(A_e, B_e)(P, Q)} = \nu_{PQ}^+ \bar{U}_{e,(A_e, B_e)\rho_e(P, Q)} = \nu_{PQ}^+ \nu_{\rho_e(P, Q)}^+ \bar{U}_{e, \rho_e(P, Q)}.$$

By (6.17), (6.15), $\nu_{PQ}^+ \nu_{\rho_e(P, Q)}^+ = \nu_{PQ}$. We obtain (6.27). Let us prove (6.23) with $PQ \in \mathcal{R}_{s+s^*}$. Using (6.28), (6.31) with $PQ = (A_{s+s^*} B_{s+s^*})(PQ)$, (6.12), (6.28) with $PQ = \rho(PQ)$ successively, we get

$$\begin{aligned} U_{s+s^*, PQ} &= \nu_{PQ}^+ U_{s+s^*, A_{s+s^*} B_{s+s^*}(PQ)} = \nu_{PQ}^+ \bar{U}_{s, \rho(A_{s+s^*} B_{s+s^*}(PQ))} \\ &= \nu_{PQ}^+ \bar{U}_{s, A_s B_s(\rho(PQ))} = \nu_{PQ}^+ \nu_{\rho(PQ)}^+ \bar{U}_{s, \rho(PQ)}. \end{aligned}$$

which gives us (6.23) by (6.17). To prove (6.23) for hyperbolic index, apply successively from left to right, (6.28), (6.30) with $PQ = (A_h B_h)(PQ)$, (6.11) and (6.28) with $PQ = \rho_h(PQ)$:

$$\begin{aligned} \nu_{PQ}^+ U_{h, PQ} &= U_{h, (A_h B_h)(PQ)} = \bar{U}_{h, \rho_h(A_h B_h)(PQ)} \\ &= \bar{U}_{h, (A_h B_h)\rho_h(PQ)} = \nu_{\rho_h(PQ)}^+ \bar{U}_{h, \rho_h(PQ)}. \end{aligned}$$

By (6.17) again, we obtain (6.23). The proof is complete. \square

6.2. Normalized mappings. We have described the conditions on centralizers. We now determine complements of these conditions to define normalized mappings.

Definition 6.4. Let $\varphi = I + (U, V)$ be a formal mapping tangent to the identity.

(i) We say that φ is *normalized* with respect to S_1, \dots, S_p if

$$U_{j, PQ} = 0 = V_{j, QP}, \quad \text{if } (P, Q) \in \mathcal{R}_j, \quad \forall j.$$

Furthermore, $\rho\varphi\rho$ is normalized w.r.t. S_1, \dots, S_p if and only if φ is.

(ii) We say that φ is *normalized* w.r.t. $\{\mathcal{S}, T_1\}$, if

$$(6.42) \quad V_{j,PQ} = -\nu_{PQ}U_{j,PQ}, \quad (P, Q) \in \mathcal{R}_j.$$

(iii) We say that φ is *normalized* w.r.t. $\{\mathcal{S}, T_1, \rho\}$ if

$$(6.43) \quad U_{h,PQ} = -\bar{U}_{h,\rho(PQ)}, \quad \forall (P, Q) \in \mathcal{R}_h;$$

$$(6.44) \quad U_{s+s_*,PQ} = -\bar{U}_{s,\rho(PQ)}, \quad \forall (P, Q) \in \mathcal{R}_{s+s_*};$$

$$(6.45) \quad U_{e,PQ} = -\nu_{PQ}\bar{U}_{e,\rho_e(PQ)}, \quad \forall (P, Q) \in \mathcal{R}_e.$$

(iv) Let $p > 1$. We say that φ is *normalized* w.r.t. $\{\mathcal{T}_1, \mathcal{T}_2\}$ if

$$(6.46) \quad U_{j,PQ} = 0, \quad (P, Q) \in \mathcal{N}_j.$$

We say that φ is *normalized* w.r.t. $\{\mathcal{T}_1, \mathcal{T}_2, \rho\}$ if

$$(6.47) \quad U_{e,PQ} = -\bar{U}_{e,(A_e, B_e) \circ \rho_e(P, Q)}, \quad \forall (P, Q) \in \mathcal{N}_e;$$

$$(6.48) \quad U_{h,PQ} = -\bar{U}_{h,\rho(P, Q)}, \quad \forall (P, Q) \in \mathcal{N}_h;$$

$$(6.49) \quad U_{s+s_*,PQ} = -\bar{U}_{s,\rho(P, Q)}, \quad \forall (P, Q) \in \mathcal{N}_{s+s_*}.$$

The set of normalized mapping w.r.t. to a family \mathcal{F} is denoted $\mathcal{C}^c(\mathcal{F})$.

Lemma 6.5. *Let F be a formal map which is tangent to the identity. There exists a unique formal decomposition $F = HG^{-1}$ with $G \in \mathcal{C}(\mathcal{S}, T_1, \rho)$ (resp. $\mathcal{C}(\mathcal{T}_1, \mathcal{T}_2, \rho)$) and $H \in \mathcal{C}^c(\mathcal{S}, T_1, \rho)$ (resp. $\mathcal{C}^c(\mathcal{T}_1, \mathcal{T}_2, \rho)$). If F is convergent, then G and H are also convergent.*

Proof. We will apply Lemma 4.3 as follows. Let \hat{H} be the set of mappings in $\mathcal{C}_2^c(\mathcal{S}, \mathcal{T}_1, \rho)$. Note that \hat{H} is a \mathbf{R} -linear subspace of $(\widehat{\mathfrak{M}}_n^2)^n$. We will define a \mathbf{R} -linear projection π from $(\widehat{\mathfrak{M}}_n^2)^n$ onto \hat{H} such that π preserves the degree of F if F is homogeneous. We will show that $\hat{G} = (I - \pi)\hat{H}$ agrees with $\mathcal{C}_2(\mathcal{S}, \mathcal{T}_1, \rho)$. We will derive estimates on π stated in Lemma 4.3, from which we conclude the convergence of H, G . The same argument will be applied to the second case of $\mathcal{C}(\mathcal{T}_1, \rho)$ and $\mathcal{C}^c(\mathcal{T}_1, \rho)$.

For the first case, let us define a projection $\pi: (\widehat{\mathfrak{M}}_n^2)^n \rightarrow \hat{H}$. We decompose

$$(U, V) = (U' + U'', V' + V''), \quad \pi(U, V) = (U', V').$$

We first define

$$(6.50) \quad U'_{j,PQ} = U_{j,PQ}, \quad V'_{j,PQ} = V_{j,PQ}, \quad U''_{j,PQ} = 0, \quad V''_{j,PQ} = 0,$$

for $(P, Q) \notin \mathcal{R}_j$. Suppose that $(P, Q) \in \mathcal{R}_e$. We have

$$U_{e,PQ} = U'_{e,PQ} + U''_{e,PQ}, \quad U_{e,\rho_e(PQ)} = U'_{e,\rho_e(PQ)} + U''_{e,\rho_e(PQ)}.$$

According to (6.45) and (6.27), we need to seek solutions that satisfy

$$(6.51) \quad U'_{e,PQ} + \nu_{PQ}\bar{U}'_{e,\rho_e(PQ)} = 0, \quad U''_{e,PQ} - \nu_{PQ}\bar{U}''_{e,\rho_e(PQ)} = 0.$$

Hence, for $(P, Q) \in \mathcal{R}_e$ we choose

$$U'_{e,PQ} = \frac{1}{2}(U_{e,PQ} - \nu_{PQ}\bar{U}_{e,\rho_e(PQ)}), \quad U''_{e,PQ} = \frac{1}{2}(U_{e,PQ} + \nu_{PQ}\bar{U}_{e,\rho_e(PQ)}).$$

We verify directly that the solutions satisfy (6.51) as follows:

$$\begin{aligned} U'_{e,PQ} + \nu_{PQ} \overline{U}'_{e,\rho_e(PQ)} &= \frac{1}{2}(U_{e,PQ} - \nu_{PQ} \overline{U}_{e,\rho_e(PQ)}) \\ &\quad + \frac{1}{2}(\nu_{PQ} \overline{U}_{e,\rho_e(PQ)} - \nu_{PQ} \nu_{\rho_e(PQ)} \overline{U}_{e,PQ}) = 0. \end{aligned}$$

Here we have used that ρ_e is an involution on \mathcal{R}_e and $\nu_{\rho_e(PQ)} \nu_{PQ} = 1$ from (6.16).

For $(P, Q) \in \mathcal{R}_h$, we achieve (6.43) and the first identity in (6.23) by taking

$$U'_{h,PQ} = \frac{1}{2}(U_{h,PQ} - \overline{U}_{h,\rho(PQ)}), \quad U''_{h,PQ} = \frac{1}{2}(U_{h,PQ} + \overline{U}_{h,\rho(PQ)}).$$

For $(P, Q) \in \mathcal{R}_{s+s^*}$, we achieve the second identity in (6.23) and (6.44) by taking

$$U'_{s+s^*,PQ} = \frac{1}{2}(U_{s+s^*,PQ} - \overline{U}_{s,\rho(PQ)}), \quad U''_{s+s^*,PQ} = \frac{1}{2}(U_{s+s^*,PQ} + \overline{U}_{s,\rho(PQ)}).$$

We have determined coefficients for $U'_{j,PQ}, U''_{j,PQ}$ with $(P, Q) \in \mathcal{R}_j$. Let us set for $(P, Q) \in \mathcal{R}_j$,

$$(6.52) \quad V''_{j,QP} = \lambda_j^{-1} \lambda^{P-Q} U''_{j,PQ},$$

$$(6.53) \quad V'_{j,QP} = V_{j,QP} - V''_{j,QP}.$$

This fulfills the conditions on V'_j and V''_j easily. Note that the first identity means that (U'', V'') commutes with T_1 . We have obtained the required formal decomposition.

To prove the convergence, we start with

$$(6.54) \quad \lambda_j^{-1} \lambda^{P-Q} = \nu_{PQ} = \pm 1$$

for $(P, Q) \in \mathcal{R}_j$. So π is indeed an \mathbf{R} -linear projection which preserves degrees. Since $|\nu_{PQ}| = 1$, we have that

$$|U'_{PQ}| \leq \max_{(P',Q')} |U_{P'Q'}|.$$

Here (P', Q') runs over all permutations of (P, Q) in $2p$ coordinates. The same holds for V' . Hence, with the notation of Lemma 4.3, we have

$$\{\pi(U, V)\}_{sym} \prec (U, V)_{sym}.$$

The existence and uniqueness as well as the convergence also follow from Lemma 4.3.

We now consider the second case of $\mathcal{C}(\mathcal{T}_1, \mathcal{T}_2, \rho)$ by minor changes. Let us define a projection $\pi: ((\widehat{\mathfrak{M}}_n^2)^n \rightarrow \hat{H}$. Here \hat{H} is the space associated with the mappings satisfying the normalized conditions (6.46)-(6.49). Let $\hat{G} = (I - \pi)\hat{H}$. We decompose as above

$$(U, V) = (U' + U'', V' + V''), \quad \pi(U, V) = (U', V').$$

Recalling that $\iota_e = (A_e, B_e) \circ \rho_e$ is an involution on \mathcal{N}_e , we choose :

$$(6.55) \quad U''_{j,PQ} = \frac{1}{2}(U_{j,PQ} + \bar{U}_{j,\rho_j(PQ)}), \quad (P, Q) \in \mathcal{N}_h,$$

$$(6.56) \quad U'_{j,PQ} = \frac{1}{2}(U_{j,PQ} - \bar{U}_{j,\rho_j(PQ)}), \quad (P, Q) \in \mathcal{N}_h,$$

$$(6.57) \quad U''_{e,PQ} = \frac{1}{2}(U_{e,PQ} + \bar{U}_{e,\iota_e(PQ)}), \quad (P, Q) \in \mathcal{N}_e,$$

$$(6.58) \quad U'_{e,PQ} = \frac{1}{2}(U_{e,PQ} - \bar{U}_{e,\iota_e(PQ)}), \quad (P, Q) \in \mathcal{N}_e,$$

$$(6.59) \quad U''_{s+s^*,PQ} = \frac{1}{2}(U_{s+s^*,PQ} + \bar{U}_{s,\rho(PQ)}), \quad (P, Q) \in \mathcal{N}_{s+s^*},$$

$$(6.60) \quad U'_{s+s^*,PQ} = \frac{1}{2}(U_{s+s^*,PQ} - \bar{U}_{s,\rho(PQ)}), \quad (P, Q) \in \mathcal{N}_{s+s^*}.$$

We still use (6.50) for $(P, Q) \notin \mathcal{R}_j$. For $(P, Q) \in \mathcal{R}_j$, define $V''_{j,QP}$ by (6.52) and $V'_{j,QP} = V_{j,QP} - V''_{j,QP}$, after we set

$$(6.61) \quad U''_{j,PQ} = \nu_{PQ}^+ U''_{j,(A_j,B_j)(P,Q)}, \quad U'_{j,PQ} = U_{j,PQ} - U''_{j,PQ}, \quad (P, Q) \in \mathcal{R}_j \setminus \mathcal{N}_j.$$

Let us verify that $\pi(U, V) = (U', V')$ is in \hat{H} . To verify (6.47) for $j = e$, via (6.56) we compute

$$U'_{e,PQ} + \bar{U}'_{e,(A_e,B_e) \circ \rho_e(PQ)} = \frac{1}{2}(U_{e,PQ} - \bar{U}_{e,\iota_e(PQ)}) + \frac{1}{2}(U_{e,\iota_e(PQ)} - \bar{U}_{e,PQ}) = 0.$$

We also know that ρ is an involution on \mathcal{N}_h and it is a bijection from \mathcal{N}_{s+s^*} onto \mathcal{N}_s . Analogously, we verify (6.48) and (6.49) via (6.56) and (6.60). Note that $(P, Q) \rightarrow (A_j, B_j)(P, Q)$ is a projection on \mathcal{N}_j . Analogously, we verify (6.46) via (6.61). This shows that $\pi(U, V)$ is in \hat{H} . We can also verify that $(U'', V'') = (I - \pi)(U, V)$ satisfies the conditions on the centralizer, i.e. it is in \hat{G} .

As before, we have

$$|U'_{j,PQ}|, |U''_{j,PQ}| \leq C \max_i \max_{(P',Q') \text{ permutation of } (P,Q)} |U_{i,P'Q'}|.$$

Equations (6.52) lead to the same inequality for V'' and hence for $V' = V - V''$. Hence, again the result follows from Lemma 4.3. \square

6.3. Convergence of linearizations.

Proposition 6.6. *Assume that the family of involutions $\{\mathcal{T}_1, \mathcal{T}_2, \rho\}$ is formally linearizable. Assume further that $\sigma_1, \dots, \sigma_p$ defined by (5.1)-(5.2), are linear.*

- (i) *There is a biholomorphic mapping in the centralizer of $\{\mathcal{S}, \rho\}$ which linearizes τ_1 and τ_2 .*
- (ii) *Assume further that $\tau_1 = T_1$ and $\tau_2 = T_2$. Then $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$ is holomorphically linearizable.*

Proof. (i) Suppose that Ψ is a formal mapping satisfying

$$\Psi^{-1} \tau_{1j} \Psi = T_{1j}, \quad \Psi \rho = \rho \Psi.$$

Then $T_{1j} = (L\Psi) \circ T_{1i_j} \circ (L\Psi)^{-1}$, and $L\Psi$ commutes with ρ . Replacing Ψ by $\Psi \circ L\Psi^{-1}$, we may assume that Ψ is tangent to the identity and $i_j = j$. We decompose $\Psi = \Psi_1\Psi_0^{-1}$, where Ψ_1 is normalized w.r.t. \mathcal{S}, T_1, ρ and Ψ_0 is in the centralizer of \mathcal{S}, T_1, ρ . Since Ψ, Ψ_0 commute with the S_j 's and ρ , then Ψ_1 commutes with the S_j 's and ρ too. We now let Ψ denote Ψ_1 .

To be more specific, let us write

$$\tau_1: \begin{cases} \xi'_i = \lambda_i \eta_i + f_i(\xi, \eta) & i = 1, \dots, p, \\ \eta'_i = \lambda_i^{-1} \xi_i + g_i(\xi, \eta) & i = 1, \dots, p, \end{cases}$$

and

$$\Psi: \begin{cases} \xi'_i = \xi_i + U_i(\xi, \eta) & i = 1, \dots, p, \\ \eta'_i = \eta_i + V_i(\xi, \eta) & i = 1, \dots, p. \end{cases}$$

Let us write that Ψ conjugates τ_1 to

$$T_1: \xi'_i = \lambda_i \eta_i, \quad \eta'_i = \lambda_i^{-1} \xi_i, \quad i = 1, \dots, p.$$

We have $\Psi \circ T_1 = \tau_1 \circ \Psi$; that is

$$(6.62) \quad \lambda_i V_i - U_i \circ T_1 = -f_i \circ \Psi(\xi, \eta) \quad i = 1, \dots, p,$$

$$(6.63) \quad \lambda_i^{-1} U_i - V_i \circ T_1 = -g_i \circ \Psi(\xi, \eta) \quad i = 1, \dots, p.$$

Since Ψ is normalized with respect $\{\mathcal{S}, T_1, \rho\}$, it satisfies Definition 6.4 (iii). Since Ψ commutes with each S_j , then $U_{j,PQ} = V_{j,QP} = 0$ for $(P, Q) \notin \mathcal{R}_j$. Since it also commutes with ρ , then by (6.23) and (6.43)-(6.44) we obtain $U_{j,PQ} = 0$ for $(P, Q) \in \mathcal{R}_j$ and $j = h, s, s + s_*$.

We need to majorize $U_{e,PQ}, V_{e,QP}$ for $(P, Q) \in \mathcal{R}_e$. By (6.63) and (6.54), we obtain

$$U_{e,PQ} - \nu_{PQ}^{-1} V_{e,QP} = -\lambda_e \{g_e \circ \Psi\}_{PQ}.$$

Using (6.24) and (6.45), we obtain $V_{e,QP} = \bar{U}_{e, \rho_e(PQ)} = -\nu_{PQ}^{-1} U_{e,PQ}$, and hence

$$U_{e,PQ} = -\frac{1}{2} \lambda_e \{g_e \circ \Psi\}_{PQ}, \quad V_{e,QP} = \frac{1}{2} \nu_{PQ} \lambda_e \{g_e \circ \Psi\}_{PQ}.$$

Therefore, we have

$$|V_{e,QP}|, |U_{e,PQ}| \leq C |\{g_e \circ \Psi\}_{PQ}|.$$

The above holds for $(P, Q) \in \mathcal{R}_e$. It holds trivially for $(P, Q) \notin \mathcal{R}_e$. In view of (4.1), we then have

$$\psi_{sym} \prec C g_{sym} \circ \Psi_{sym} = g_{sym} \circ (I_{sym} + \psi_{sym}).$$

Therefore, ψ_{sym} is convergent at the origin and so is Ψ .

(ii) Assume now that $\sigma = S, \tau_1 = T_1, \tau_2 = T_2$ are linear. Suppose that Ψ linearizes the $\{\tau_{ij}\}$ and commutes with ρ . We decompose $\Psi = \Psi_1\Psi_0^{-1}$ with Ψ_1 being normalized w.r.t. $\mathcal{T}_1, \mathcal{T}_2, \rho$ and with Ψ_0 being in the centralizer of $\mathcal{T}_1, \mathcal{T}_2, \rho$. Since $\Psi^{-1}\tau_{ij}\Psi = T_{ij}$, we have $\Psi_1^{-1}\tau_{ij}\Psi_1 = \Psi_0^{-1}T_{ij}\Psi_0 = T_{ij}$. Hence, Ψ_1 linearizes the τ_{ij} and is normalized w.r.t. $\mathcal{T}_1, \mathcal{T}_2, \rho$. Since Ψ, Ψ_0 commute with \mathcal{S}, T_1 and ρ , so does Ψ_1 .

We recall

$$T_{1j}: \begin{cases} \xi'_j = \lambda_j \eta_j \\ \eta'_j = \lambda_j^{-1} \xi_j \\ \xi'_k = \xi_k, & k \neq j \\ \eta'_k = \eta_k, & k \neq j, \end{cases} \quad \tau_{1j}: \begin{cases} \xi'_j = \lambda_j \eta_j + f_{jj}(\xi, \eta) \\ \eta'_j = \lambda_j^{-1} \xi_j + g_{jj}(\xi, \eta) \\ \xi'_k = \xi_k + f_{jk}(\xi, \eta), & k \neq j \\ \eta'_k = \eta_k + g_{jk}(\xi, \eta), & k \neq j. \end{cases}$$

Since we have $\Psi \circ T_{1j} = \tau_{1j} \circ \Psi$, we obtain the following relations

$$(6.64) \quad \begin{cases} \lambda_j V_j - U_j \circ T_{1j} = -f_{jj} \circ \Psi \\ \lambda_j^{-1} U_j - V_j \circ T_{1j} = -g_{jj} \circ \Psi \\ U_k - U_k \circ T_{1j} = -f_{jk} \circ \Psi, & k \neq j \\ V_k - V_k \circ T_{1j} = -g_{jk} \circ \Psi, & k \neq j. \end{cases}$$

Since $\Psi \in \mathcal{C}(\mathcal{S}, T_1, \rho)$, combining (6.23), (6.26) with the normalizing conditions (6.48), (6.49), we find that $U_{j,PQ} = 0 = V_{j,QP}$ for $(P, Q) \in \mathcal{N}_j$ and $j = h, s, s + s_*$. Using $\Psi\rho = \rho\Psi$, we get $V_e = \overline{U_e} \circ \rho$. By the first equation above, we get

$$\lambda_e \overline{U_e} \circ \rho \circ T_{1e} - U_e = -f_{ee} \circ \tau_{1e} \circ \Psi.$$

For $(P, Q) \in \mathcal{N}_e$, we have $(\lambda_e U_e \circ \rho \circ T_{1e})_{PQ} = \overline{U_{e, A_e B_e(\rho_e(PQ))}}$. By (6.47), we get

$$(6.65) \quad U_{e,PQ} = \frac{1}{2} \{f_{e,e} \circ \tau_{1e} \circ \Psi\}_{PQ}, \quad V_{e,QP} = \nu_{PQ} U_{e,PQ}, \quad (PQ) \in \mathcal{N}_e.$$

We now majorize $U_{j,PQ}, V_{j,QP}$ for $(P, Q) \in \mathcal{R}_j \setminus \mathcal{N}_j$. Fix $(P, Q) \in \mathcal{R}_k \setminus \mathcal{N}_k$. Start with some j such that $p_j < q_j$. In the second last identity in (6.64), let us compose on the right by $T_{1j'}$ with $j' \neq k, j$ to get

$$U_k \circ T_{1j'} - U_k \circ T_{1j} \circ T_{1j'} = -f_{j,k} \circ \Psi \circ T_{1j'} = f_{j,k} \circ \tau_{1j'} \circ \Psi.$$

Let $\{\ell_1, \dots, \ell_d\}$ be the set of $i \neq k$ such that $p_i < q_i$. Composing successively with the $T_{1\ell_j}$'s and adding, we get

$$U_k - U_k \circ T_{1\ell_1} \circ \dots \circ T_{1\ell_d} = - \sum_{i=1}^d f_{\ell_i, k} \circ \tau_{1\ell_1} \circ \dots \circ \tau_{1\ell_{i-1}} \circ \Psi.$$

Hence, if $PQ \in \mathcal{R}_k \setminus \mathcal{N}_k$, then

$$U_{k,PQ} = \{U_k \circ T_{1\ell_1} \circ \dots \circ T_{1\ell_d}\}_{PQ} - \left\{ \sum_{i=1}^d f_{\ell_i, k} \circ \tau_{1\ell_1} \circ \dots \circ \tau_{1\ell_{i-1}} \circ \Psi \right\}_{PQ}.$$

The first term on the right-hand side, $U_{k, (A_k, B_k)(PQ)}$, is either zero or majorized by (6.65). The summations have finitely many combinations. This shows that $U_k \prec a_k \circ \overline{\Psi}$. By (6.26), we obtain $V_j \prec (U_j)_{sym}$. This shows that $(U, V) \prec b \circ (I_{sym} + (U, V)_{sym})$ for some analytic mapping $b = O(2)$. Using Lemma 4.3, we obtain the convergence of U_k, V_k . \square

Theorem 6.7. *Let M be a germ of analytic submanifold that is an higher order perturbation of a product quadric Q in \mathbf{C}^{2p} . Assume that M satisfies condition J and it is formally equivalent to Q . Suppose that each hyperbolic component has an eigenvalue μ_h which is either a root of unity or satisfies the Brjuno condition (7.38) in which $\mathcal{I} = 0$, $\ell = 1, n =$*

1, $\mu_i = \mu_{i,j} = \mu_h$, and each μ_s is not a root of unity and satisfies the Brjuno condition. Then M is holomorphically equivalent to the product quadric.

Proof. We first apply a theorem (with $\mathcal{I} = 0$) in [Sto15] that linearize simultaneously and holomorphically the $\sigma_1, \dots, \sigma_p$. Note that the small divisor condition in this special case is equivalent that each μ_h is either a root of unity or a Brjuno number. Then, we apply successively the two assertions of Proposition 6.6. Hence, in good holomorphic coordinates, $\{\tau_{11}, \dots, \tau_{1p}, \rho\}$ are linear. Then, by Proposition 2.8, the manifold is holomorphically equivalent to the quadric. \square

As in the case of Theorem 5.3, we can also prove the first part of the above proof by applying Rüssmann's theorem [Rüs02] successively to each σ_i . This is due to the commutativity property and the special type of the linear parts that lead to the relatively simpler relations on $\mathcal{C}(S_i)$ and $\mathcal{C}^c(S_i)$ for each fixed i .

7. EXISTENCE OF ATTACHED COMPLEX MANIFOLDS

We are interested in complex submanifolds K in \mathbf{C}^{2p} that intersect the real submanifold M at the origin. Recall that M has real dimension $2p$. Generically, the origin is an isolated intersection point if $\dim K = p$. Let us consider the situation when the intersection has dimension p . Without further restrictions, there are many such complex submanifolds; for instance, we can take a p -dimensional totally real and real analytic submanifold K_1 of M . We then let K be the complexification of K_1 . To ensure the uniqueness or finiteness of the complex submanifolds K , we therefore introduce the following.

Definition 7.1. Let M be a formal real submanifold of dimension $2p$ in \mathbf{C}^{2p} . We say that a formal complex submanifold K is *attached* to M if $K \cap M$ contains at least two germs of totally real and formal submanifolds K_1, K_2 of dimension p that intersect transversally at the origin. Such a pair $\{K_1, K_2\}$ are called a pair of *asymptotic* formal submanifolds of M .

Before we present the details, let us describe the main steps to derive the results. We first derive the results at the formal level. We then apply the results of [Pös86] and [Sto15]. The proof of the co-existence of convergent and divergent attached submanifolds will rely on a theorem of Pöschel on stable invariant submanifolds and Siegel's small divisor technique.

We now describe the formal results. When $p = 1$, a non-resonant hyperbolic M admits a unique attached formal holomorphic curve [Kli85]. When $p > 1$, new situations arise. First, we show that there are obstructions to attach formal submanifolds. However, the formal obstructions disappear when M admits the maximum number of deck transformations and M is non-resonant. These two conditions allow us to express M in an equivalent form (3.4). This equivalent form for M , which has not been used so far, will play an essential role in our proof for $p > 1$.

We will consider a real submanifold M which is a higher order perturbation of a non-resonant product quadrics. By adapting the proof of Klingenberg [Kli85] to the manifold M (3.4), we will show the existence of a unique attached formal submanifold for a prescribed non-resonance condition. As in [Kli85], we also show that the complexification of K in \mathcal{M} is a pair of invariant formal submanifolds $\mathcal{K}_1, \mathcal{K}_2$ of σ . Furthermore, K is convergent if and only if \mathcal{K}_1 is convergent.

Let us first recall the values of the Bishop invariants. The types of the invariants play an important role for the existence and the convergence of attached formal complex submanifolds. From (2.22) and (2.26), we recall that

$$(7.1) \quad \gamma_e^{-1} = \lambda_e + \lambda_e^{-1}, \quad \gamma_h^{-1} = \lambda_h + \bar{\lambda}_h, \quad \gamma_s^{-1} = 1 + \bar{\lambda}_s^2,$$

$$(7.2) \quad 0 < \gamma_e < 1/2, \quad \gamma_h > 1/2, \quad \gamma_s \in (-\infty, 1/2) + i(0, \infty), \quad \gamma_{s+s_*} = 1 - \bar{\gamma}_s.$$

Here we exclude the case that $\operatorname{Re} \gamma_s = 1/2$ or $\gamma_s < 1/2$ as we will assume that σ has distinct eigenvalues. We normalize

$$(7.3) \quad \lambda_e > 1, \quad |\lambda_h| = 1, \quad |\lambda_s| > 1, \quad \lambda_{s+s_*} = \bar{\lambda}_s^{-1};$$

$$(7.4) \quad \arg \lambda_h \in (0, \pi/2), \quad \arg \lambda_s \in (0, \pi/2).$$

Recall that $\mu_j = \lambda_j^2$. By (7.1), we have

$$(7.5) \quad \gamma_j^2 = \frac{\mu_j}{(1 + \mu_j)^2}, \quad j = e, h; \quad \gamma_s \bar{\gamma}_{s+s_*} = \frac{\bar{\mu}_s}{(1 + \bar{\mu}_s)^2}.$$

We first verify the following.

Lemma 7.2. *Let γ_j, λ_j be given by (7.1)-(7.4). Let $\mu_j = \lambda_j^2$. Assume that $\mu_1, \mu_1^{-1}, \dots, \mu_p, \mu_p^{-1}$ are distinct. Then $\gamma_e^2, \gamma_h^2, \bar{\gamma}_s \gamma_{s+s_*}, \gamma_s \bar{\gamma}_{s+s_*}$ are distinct p numbers. The latter is equivalent to $\gamma_1, \dots, \gamma_p$ being distinct.*

Proof. Note that $x^{-1} + x$ and x^{-1} decrease strictly on $(0, 1)$. So γ_e^2, γ_h^2 are distinct. We also have $\gamma_s \bar{\gamma}_{s+s_*} = \gamma_s - \gamma_s^2$. If a, b are complex numbers, then $a - a^2 = b - b^2$ if and only if $a = b$ or $a + b = 1$. Since γ_s is not real, then $\gamma_s \bar{\gamma}_{s+s_*}$ are different from γ_e^2 and γ_h^2 . For any distinct complex numbers a_1, a_2 in $(-\infty, 1/2) + i(0, \infty)$. We have $1 - a_2 \neq 1 - a_1, a_1, a_2$. The lemma is proved. \square

Let us first investigate the numbers of pairs of formal asymptotic submanifolds and attached formal submanifolds.

Lemma 7.3. *Let M be a formal submanifold that is a third order perturbation of a product quadric Q in \mathbf{C}^{2p} . Assume that M has distinct eigenvalues*

$$\mu_1, \dots, \mu_p, \quad \mu_1^{-1}, \dots, \mu_p^{-1}.$$

- (i) *If M admits an attached formal submanifold, its CR singularity has no elliptic component.*
- (ii) *If Q has no elliptic components, then Q has at least $2^{h_*+s_*-1}$ pairs of asymptotic totally real and real analytic submanifolds that are contained in a single attached complex submanifold.*
- (iii) *There is no formal submanifold attached to*

$$M: z_3 = (z_1 + 2\gamma_1 \bar{z}_1)^2 + (z_2 + 2\gamma_2 \bar{z}_2)^3, \quad z_4 = (z_2 + 2\gamma_2 \bar{z}_2)^2.$$

Here M has a hyperbolic complex tangent at the origin.

- (iv) *Assume that M has no elliptic component and it admits the maximum number of formal deck transformations. Given $\epsilon_h, \epsilon_s = \pm 1$, let $\nu = \nu_\epsilon := (\nu_1, \dots, \nu_p)$ with*

$$(7.6) \quad \nu_h := \mu_h^{\epsilon_h}, \quad \nu_s := \bar{\mu}_s^{\epsilon_s}, \quad \nu_{s_*+s} := \mu_s^{-\epsilon_s}$$

Suppose that

$$(7.7) \quad \nu^Q \neq \nu_j^{-1}, \quad \forall Q \in \mathbf{N}^p, \quad |Q| > 0, \quad 1 \leq j \leq p.$$

Then M admits a unique pair of asymptotic formal submanifolds K_1, K_2 such that each K_i is defined by $z' = \rho_i(z')$ for a formal anti-holomorphic involution ρ_i and the linear part of $\overline{\rho_2}^{-1}\overline{\rho_1}$ has eigenvalues ν_1, \dots, ν_p . In particular, if (7.7) holds for each ν of the form (7.6) then M admits exactly $2^{h_*+s_*-1}$ pairs of asymptotic formal submanifolds.

Proof. (i) Let M be defined by

$$(7.8) \quad z_{p+j} = Q_j(z', \overline{z}') + H_j(z', \overline{z}'), \quad 1 \leq j \leq p$$

where $H_j(z', \overline{z}') = O(|z'|^3)$ and each Q_j is quadratic. Let $\{K_1, K_2\}$ be a pair of asymptotic formal submanifolds of M intersecting a formal complex submanifold K . We know that the totally real spaces T_0K_1, T_0K_2 are contained in T_0M , the z' -subspace. Let K'_i be the projection of K_i onto the z' -subspace, then K'_1, K'_2 are still totally real. Let K'_1 be defined by

$$K'_1: \overline{z}' = \mathbf{A}z' + R(z'), \quad \overline{\mathbf{A}}\mathbf{A} = \mathbf{I}, \quad R(z') = O(2)$$

such that $\rho_1(z') := \overline{\mathbf{A}}\overline{z}' + \overline{R}(\overline{z}')$ defines an anti-holomorphic formal involution. Let K_2 be the (formal) fixed-point set of the anti-holomorphic involution $\rho_2(z') = \overline{\tilde{\mathbf{A}}}\overline{z}' + \overline{\tilde{R}}(\overline{z}')$ with $\tilde{R}(z') = O(2)$. Then K_1, K_2 intersect transversally at the origin if and only if $\det(\tilde{\mathbf{A}} - \mathbf{A}) \neq 0$. Let us define holomorphic mappings

$$(7.9) \quad \overline{\rho}_i(z') := \overline{\rho_i(z')}, \quad i = 1, 2.$$

Then K is given by

$$(7.10) \quad z''_{p+j} = Q_j(z', \overline{\rho}_i(z')) + H_j(z', \overline{\rho}_i(z')), \quad i = 1, 2, \quad j = 1, \dots, p.$$

The two equations agree, if and only if

$$(7.11) \quad Q_j(z', \overline{\rho}_1(z')) + H_j(z', \overline{\rho}_1(z')) = Q_j(z', \overline{\rho}_2(z')) + H_j(z', \overline{\rho}_2(z')), \quad 1 \leq j \leq p.$$

Then the asymptotic totally real submanifolds $\{K_1, K_2\}$ are defined by

$$(7.12) \quad K_i: z_{p+j} = Q_j(z', \overline{z}') + H_j(z', \overline{z}'), \quad 1 \leq j \leq p, \quad \rho_i(z') = z'.$$

Recall that

$$\begin{aligned} Q_j(z', \overline{z}) &= (z_j + 2\gamma_j \overline{z}_j)^2, \quad j = e, h; \\ Q_s(z', \overline{z}') &= (z_{s+s_*} + 2\gamma_{s+s_*} \overline{z}_s)^2, \\ Q_{s+s_*}(z', \overline{z}') &= (z_s + 2\gamma_s \overline{z}_{s+s_*})^2. \end{aligned}$$

Let us first find necessary conditions on the linear parts of ρ_i for (7.11) to be solvable. Let $w' = \mathbf{A}z'$ and $\tilde{w}' = \tilde{\mathbf{A}}z'$. Comparing the quadratic terms in (7.11) for $i = 1, 2$, we see that

$$\begin{aligned} (z_j + 2\gamma_j w_j)^2 &= (z_j + 2\gamma_j \tilde{w}_j)^2, \\ (z_{s+s_*} + 2\gamma_{s+s_*} w_s)^2 &= (z_{s+s_*} + 2\gamma_{s+s_*} \tilde{w}_s)^2, \\ (z_s + 2\gamma_s w_{s+s_*})^2 &= (z_s + 2\gamma_s \tilde{w}_{s+s_*})^2. \end{aligned}$$

Here $\gamma_{s+s_*} = 1 - \bar{\gamma}_s$, by (7.2). For each j , $w_j \neq \tilde{w}_j$. Otherwise, the fixed points of ρ_1 and ρ_2 do not intersect transversally. Therefore, the above 3 identities can be written as

$$\begin{aligned} z_j + 2\gamma_j w_j &= -(z_j + 2\gamma_j \tilde{w}_j), \\ z_{s+s_*} + 2\gamma_{s+s_*} w_s &= -(z_{s+s_*} + 2\gamma_{s+s_*} \tilde{w}_s), \\ z_s + 2\gamma_s w_{s+s_*} &= -(z_s + 2\gamma_s \tilde{w}_{s+s_*}). \end{aligned}$$

In the matrix form, we get

$$(7.13) \quad \tilde{\mathbf{A}} = -\gamma^{-1} - \mathbf{A}, \quad \text{with} \quad \gamma := \begin{pmatrix} \gamma_{e_*} & \mathbf{0} & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \gamma_{h_*} & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & \mathbf{0} & \gamma_{s_*} \\ \mathbf{0} & \mathbf{0} & \tilde{\gamma}_{s_*} & \mathbf{0} \end{pmatrix}.$$

Here in matrices $\tilde{\gamma}_{s_*} = \mathbf{I}_{s_*} - \bar{\gamma}_{s_*}$. Let us express in block matrices

$$\mathbf{A} = \begin{pmatrix} \mathbf{A}_{e_* e_*} & \mathbf{A}_{e_* h_*} & \mathbf{A}_{e_* s_*} & \mathbf{A}_{e_* (2s_*)} \\ \mathbf{A}_{h_* e_*} & \mathbf{A}_{h_* h_*} & \mathbf{A}_{h_* s_*} & \mathbf{A}_{h_* (2s_*)} \\ \mathbf{A}_{s_* e_*} & \mathbf{A}_{s_* h_*} & \mathbf{A}_{s_* s_*} & \mathbf{A}_{s_* (2s_*)} \\ \mathbf{A}_{(2s_*) e_*} & \mathbf{A}_{(2s_*) h_*} & \mathbf{A}_{(2s_*) s_*} & \mathbf{A}_{(2s_*) (2s_*)} \end{pmatrix}$$

where the diagonal block matrices are of sizes $e_* \times e_*$, $h_* \times h_*$, $s_* \times s_*$, and $s_* \times s_*$, respectively. By $\mathbf{A}\bar{\mathbf{A}} = \mathbf{I}$, $\tilde{\mathbf{A}}\bar{\tilde{\mathbf{A}}} = \mathbf{I}$ and (7.13) we get $(\tilde{\mathbf{A}}\bar{\tilde{\mathbf{A}}} - \mathbf{A}\bar{\mathbf{A}})\bar{\gamma} = \mathbf{0}$, i.e. $\gamma^{-1} + \mathbf{A} + \gamma^{-1}\bar{\mathbf{A}}\bar{\gamma} = \mathbf{0}$. Recall that $\gamma_1^2, \dots, \gamma_{e_*+h_*}^2$ are real and distinct. It is easy to see that $\mathbf{A}_{e_* h_*} = \mathbf{0}$, $\mathbf{A}_{h_* e_*} = \mathbf{0}$, and $\mathbf{A}_{e_* e_*}$, $\mathbf{A}_{h_* h_*}$ are diagonal. Also,

$$(7.14) \quad \mathbf{A}_{e_* e_*} + \bar{\mathbf{A}}_{e_* e_*} = -\gamma_{e_*}^{-1}, \quad \mathbf{A}_{h_* h_*} + \bar{\mathbf{A}}_{h_* h_*} = -\gamma_{h_*}^{-1}.$$

In block matrices, we obtain

$$(7.15) \quad \gamma_j^{-1} \bar{\mathbf{A}}_{j(2s_*)} \bar{\gamma}_{s_*} = -\mathbf{A}_{j s_*}, \quad \tilde{\gamma}_{s_*}^{-1} \bar{\mathbf{A}}_{(2s_*)j} \bar{\gamma}_j = -\mathbf{A}_{s_* j};$$

$$(7.16) \quad \gamma_j^{-1} \bar{\mathbf{A}}_{j s_*} \bar{\gamma}_{s_*} = -\mathbf{A}_{j(2s_*)}, \quad \gamma_{s_*}^{-1} \bar{\mathbf{A}}_{s_* j} \bar{\gamma}_j = -\mathbf{A}_{(2s_*)j};$$

$$(7.17) \quad \tilde{\gamma}_{s_*}^{-1} \bar{\mathbf{A}}_{(2s_*)(2s_*)} \bar{\gamma}_{s_*} = -\mathbf{A}_{s_* s_*}, \quad \tilde{\gamma}_{s_*}^{-1} \bar{\mathbf{A}}_{(2s_*)s_*} \bar{\gamma}_{s_*} = -\mathbf{A}_{s_* (2s_*)} - \tilde{\gamma}_{s_*}^{-1},$$

$$(7.18) \quad \gamma_{s_*}^{-1} \bar{\mathbf{A}}_{s_* (2s_*)} \bar{\gamma}_{s_*} = -\mathbf{A}_{(2s_*)s_*} - \gamma_{s_*}^{-1}, \quad \gamma_{s_*}^{-1} \bar{\mathbf{A}}_{s_* s_*} \bar{\gamma}_{s_*} = -\mathbf{A}_{(2s_*)(2s_*)}.$$

In the first 4 equations, we have $j = e_*, h_*$.

By Lemma 7.2, we know that γ_e^2, γ_h^2 , and $\gamma_s \bar{\gamma}_{s+s_*}$ are distinct. Thus, $\mathbf{A}_{j s_*} = \mathbf{A}_{j(2s_*)} = \mathbf{0}$ and $\mathbf{A}_{s_* j} = \mathbf{A}_{(2s_*)j} = \mathbf{0}$ for $j = e_*, h_*$. Since $\gamma_s \bar{\gamma}_{s+s_*}$ is different from all $\gamma_{s+s_*} \bar{\gamma}_s$, then $\mathbf{A}_{s_* s_*} = \mathbf{A}_{(2s_*)(2s_*)} = \mathbf{0}$ while $\mathbf{A}_{s_* (2s_*)}$, $\mathbf{A}_{(2s_*) s_*}$ are diagonal. Now $\mathbf{A}\bar{\mathbf{A}} = \mathbf{I}$ implies that

$$(7.19) \quad \mathbf{A}_{e_* e_*} \bar{\mathbf{A}}_{e_* e_*} = \mathbf{I}, \quad \mathbf{A}_{h_* h_*} \bar{\mathbf{A}}_{h_* h_*} = \mathbf{I}, \quad \mathbf{A}_{s_* (2s_*)} \bar{\mathbf{A}}_{(2s_*) s_*} = \mathbf{I}.$$

Combining the first identities in (7.14) and (7.19), we know that the diagonal e th element a_e of $\mathbf{A}_{e_* e_*}$ must satisfy $a_e + \bar{a}_e = -\gamma_e^{-1}$, $a_e \bar{a}_e = 1$. Since $0 < \gamma_e < 1/2$, there is no such solution a_e if $e_* > 0$. We have verified (i).

(ii) For the hyperbolic components, by (7.1) we have $\gamma_h^{-1} = \lambda_h + \bar{\lambda}_h$ with $|\lambda_h| = 1$. By the second identities in (7.14), (7.19), and by (7.13), we obtain $(a_h, \tilde{a}_h) = (-\lambda_h, -\bar{\lambda}_h)$ or $(-\bar{\lambda}_h, -\lambda_h)$. For the complex components, we use $\mathbf{A}_{s_* (2s_*)} \bar{\mathbf{A}}_{(2s_*) s_*} = \mathbf{I}$ and multiply the

second identity in (7.17) by the diagonal matrix $\mathbf{A}_{s_*(2s_*)}$. Thus the s th diagonal element a_s of $\mathbf{A}_{s_*(2s_*)}$ satisfies $a_s(a_s + \tilde{\gamma}_s^{-1}) + \tilde{\gamma}_s^{-1}\bar{\gamma}_s = 0$. By the last identity in (7.5), we get

$$a_s^2 + (1 + \mu_s^{-1})a_s + \mu_s^{-1} = 0.$$

Obviously $a_s = -1, -\mu_s^{-1}$ are solutions. By (7.13), we get $(a_s, \tilde{a}_s) = (-1, 1 - \tilde{\gamma}_s^{-1}) = (-1, -\mu_s^{-1})$ or $(-\mu_s^{-1}, \mu_s^{-1} - \tilde{\gamma}_s^{-1}) = (-\mu_s^{-1}, -1)$. Each tuple determines a tuple (b_s, \tilde{b}_s) by (7.19), with b_s being the diagonal entries of $\mathbf{A}_{(2s_*)_{s_*}}$. We verify that $(b_s, \tilde{b}_s) = (\bar{a}_s^{-1}, \tilde{\bar{a}}_s^{-1})$. There are exactly $2^{h_*+s_*-1}$ solutions for \mathbf{A} and $\tilde{\mathbf{A}}$ since we can only determine the pairs $\{\mathbf{A}_{h_*h_*}, \tilde{\mathbf{A}}_{h_*h_*}\}, \{\mathbf{A}_{s_*(2s_*)}, \tilde{\mathbf{A}}_{s_*(2s_*)}\}$. Indeed, we have

$$(7.20) \quad \mathbf{A} = \begin{pmatrix} \text{diag}(a_h) & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & \text{diag}(a_s) \\ \mathbf{0} & \text{diag}(b_s) & \mathbf{0} \end{pmatrix}, \quad \tilde{\mathbf{A}} = \begin{pmatrix} \text{diag}(\tilde{a}_h) & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & \text{diag}(\tilde{a}_s) \\ \mathbf{0} & \text{diag}(\tilde{b}_s) & \mathbf{0} \end{pmatrix},$$

$$(7.21) \quad \text{diag } \nu := \tilde{\mathbf{A}}^{-1}\mathbf{A} = \begin{pmatrix} \text{diag}(\tilde{a}_h^{-1}a_h) & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \text{diag}(\tilde{b}_s^{-1}b_s) & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & \text{diag}(\tilde{a}_s^{-1}a_s) \end{pmatrix},$$

$$(7.21) \quad \nu = \mu_\epsilon = (\mu_h^{\epsilon_h}, \bar{\mu}_s^{\epsilon_s}, \mu_s^{-\epsilon_s}), \quad \epsilon_h^2, \epsilon_s^2 = 1, \quad \nu_{s+s_*} = \bar{\nu}_s^{-1},$$

where there are $2^{h_*+s_*}$ distinct combinations. Thus, we get exactly $2^{h_*+s_*-1}$ pairs $\{K_\epsilon^1, K_\epsilon^2\}$ of asymptotic linear submanifolds indexed by $\epsilon = (\epsilon_1, \dots, \epsilon_{h_*+s_*})$ with $\epsilon_j^2 = 1$ for the product quadric. The attached formal submanifolds associated to these linear asymptotic submanifolds are unique and restricting to $\epsilon_i = 1$ for all i , it is given by

$$z_{p+h} = (1 - 4\gamma_h^2)z_h^2, \quad z_{p+s} = (1 - 2\gamma_{s+s_*})^2 z_{s+s_*}^2, \quad z_{p+s+s_*} = (1 - 2\gamma_s)^2 z_s^2.$$

Here we have used $(1 - 4\gamma_h^2) = (1 - 2\gamma_h\lambda_h)^2$.

In summary, we have shown that there are exactly $2^{h_*+s_*-1}$ pairs of linear anti-holomorphic involutions $\{\rho_1, \rho_2\}$. In (iv) we show that under the non-resonant conditions on μ_1, \dots, μ_p , they are the only pairs of anti-holomorphic involutions. This finishes the proof of (ii).

(iii). Let us continue the computation for the perturbations. We have determined linear parts of antiholomorphic involutions ρ_i . We expand components of $R(z')$ as $R_j(z') = \sum_{k=2}^{\infty} R_{j;k}(z')$. Here $R_{j;k}$ are homogeneous terms of degree k . We expand \tilde{R}_j analogously. Suppose that terms of order up to $k-1$ in R_j, \tilde{R}_j have been determined. For the hyperbolic components, we need to solve the equations

$$(7.22) \quad 4\sqrt{1 - 4\gamma_h^2}z_h(R_{h;k}(z') + \tilde{R}_{h;k}(z')) = \dots,$$

where the right-hand side has been determined. Indeed, let us compute the terms of degree k in (7.11) to obtain

$$(1 - 2\gamma_j\lambda_j)^2 z_j^2 + 2(1 - 2\gamma_j\lambda_j)z_j R_{j;k} = (1 - 2\gamma_j\lambda_j^{-1})^2 z_j^2 + 2(1 - 2\gamma_j\lambda_j^{-1})z_j \tilde{R}_{j;k} + \mathcal{R}$$

where \mathcal{R} is a polynomial that depends on $\tilde{R}_{j;l}, R_{j;l}, l < k$. Since $(1 - 2\gamma_j\lambda_j) = -(1 - 2\gamma_j\lambda_j^{-1})$, we obtain (7.22).

When $p > 1$, the system of equations (7.22) cannot be solved even formally, unless the right-hand side is divisible by z_h . When $p = 1$, the equation (7.22) is clearly solvable.

In fact, under the non-resonant condition on μ_1 , the formal anti-holomorphic involutions $\{\rho_1, \rho_2\}$ can be uniquely determined.

Let us keep the above notation and compute for the example stated in (iii). We need to solve

$$\begin{aligned} (z_1 + 2\gamma_1\tilde{w}_1)^2 + (z_2 + 2\gamma_2\tilde{w}_2)^3 &= (z_1 + 2\gamma_1w_1)^2 + (z_2 + 2\gamma_2w_2)^3, \\ (z_2 + 2\gamma_2\tilde{w}_2)^2 &= (z_2 + 2\gamma_2w_2)^2. \end{aligned}$$

Again $\tilde{w}_2 - w_2$ cannot be identically zero. Thus $\tilde{w}_2 = -w_2 - \gamma_2^{-1}z_2$. Then we need to solve

$$(z_1 + 2\gamma_1\tilde{w}_1)^2 = (z_1 + 2\gamma_1w_1)^2 + 2(z_2 + 2\gamma_2w_2)^3.$$

By (ii), we know that $w_1 = -\lambda_1z_1 + R_1(z')$ and $w_2 = -\lambda_2z_2 + R_2(z')$ with $R_i(z') = O(2)$. Also $\tilde{w}_1 = -\bar{\lambda}_1z_1 + \tilde{R}_1(z')$ and $\tilde{w}_2 = -\bar{\lambda}_2z_2 + \tilde{R}_2(z')$. Comparing the cubic terms implies that z_1 must divide $2(1 - 2\gamma_2\lambda_2)^3z_2^3$, which is a contradiction.

(iv) For a general M , following Klingenberg [Kli85] we reformulate the problem by considering the following equations

$$\begin{aligned} h(z') &= Q(z', \bar{\rho}_i(z')) + H(z', \bar{\rho}_i(z')), \quad i = 1, 2, \\ h^*(\bar{\rho}_i(z')) &= \bar{Q}(\bar{\rho}_i(z'), z') + \bar{H}(\bar{\rho}_i(z'), z'), \quad i = 1, 2. \end{aligned}$$

Here $h, h^*, \bar{\rho}_i$ are unknowns. Initially, we require that $\bar{\rho}_1, \bar{\rho}_2$ be arbitrary biholomorphic maps, except their linear parts match with $z' \rightarrow Az'$ and $z' \rightarrow \tilde{A}z'$. This will ensure that the solutions $\bar{\rho}_i$ are unique and they are involutions.

As demonstrated in (iii), in general there is no formal submanifold attached to M . Thus we assume that M is a higher order perturbation of product quadric without elliptic component and it admits the maximum number of deck transformation.

We may assume that

$$(7.23) \quad z_{p+h} = (z_h + 2\gamma_h\bar{z}_h + E_h(z', \bar{z}'))^2,$$

$$(7.24) \quad z_{p+s} = (z_s + 2\gamma_s\bar{z}_{s+s_*} + E_s(z', \bar{z}'))^2,$$

$$(7.25) \quad z_{p+s+s_*} = (z_{s+s_*} + 2\gamma_{s+s_*}\bar{z}_s + E_{s+s_*}(z', \bar{z}'))^2.$$

For late references, we express the above in an abbreviated form:

$$(7.26) \quad M \subset \mathbf{C}^{2p}: z_{p+j} = (L_j(z', \bar{z}') + E_j(z', \bar{z}'))^2, \quad 1 \leq j \leq p.$$

We fix linear parts of ρ_i such that

$$\rho_1(z') = \bar{A}\bar{z}' + \bar{R}(\bar{z}'), \quad \rho_2(z') = \bar{\tilde{A}}\bar{z}' + \bar{\tilde{R}}(\bar{z}').$$

For $i = 1, 2$ we then need to solve $\bar{\rho}_i$ from

$$(7.27) \quad z_h + 2\gamma_h\bar{\rho}_{ih} + E_h(z', \bar{\rho}_i) = (-1)^i f_h(z'),$$

$$(7.28) \quad z_s + 2\gamma_s\bar{\rho}_{is+s_*} + E_s(z', \bar{\rho}_i) = (-1)^i f_s(z'),$$

$$(7.29) \quad z_{s+s_*} + 2\gamma_{s+s_*}\bar{\rho}_{is} + E_{s+s_*}(z', \bar{\rho}_i) = (-1)^i f_{s+s_*}(z'),$$

$$(7.30) \quad 2\gamma_h z_h + \bar{\rho}_{ih} + \bar{E}_h(\bar{\rho}_i, z') = (-1)^i f_h^*(\bar{\rho}_i),$$

$$(7.31) \quad 2\gamma_s z_{s+s_*} + \bar{\rho}_{is} + \bar{E}_s(\bar{\rho}_i, z') = (-1)^i f_s^*(\bar{\rho}_i),$$

$$(7.32) \quad 2\gamma_{s+s_*} z_s + \bar{\rho}_{is+s_*} + \bar{E}_{s+s_*}(\bar{\rho}_i, z') = (-1)^i f_{s+s_*}^*(\bar{\rho}_i).$$

Suppose that we have determined terms of $R_j, \tilde{R}_j, f_j, f_j^*$ of order $< k$. We have

$$\bar{\rho}_1(z') = \mathbf{A}z' + R(z'), \quad \bar{\rho}_1^{-1}(z') = \mathbf{A}^{-1}z' - \mathbf{A}^{-1}R'(\mathbf{A}^{-1}z'),$$

where the terms in $R' - R$ of order k depend only on terms of R of order $< k$. For terms of order k , by eliminating f_j , we need to solve

$$(7.33) \quad R_{jQ} + \tilde{R}_{jQ} = \dots$$

where the dots denote terms which have been determined. We compose from right in the last 3 identities for $i = 1$ (resp. $i = 2$) by $\bar{\rho}_1^{-1}$ (resp. $\bar{\rho}_2^{-1}$). Eliminating f^* from the new identities, we obtain

$$\mathbf{A}^{-1}R(\mathbf{A}^{-1}z') + \tilde{\mathbf{A}}^{-1}\tilde{R}(\tilde{\mathbf{A}}^{-1}z') = \dots$$

Recall that $\tilde{\mathbf{A}}^{-1}\mathbf{A} = \text{diag } \nu$ with $\nu := \nu_\epsilon$. Multiplying on the left by $\tilde{\mathbf{A}}$, using $\tilde{\mathbf{A}}\mathbf{A}^{-1} = (\text{diag } \nu)^{-1}$, and evaluating at $z' = \mathbf{A}\tilde{z}'$, we thus need to solve (7.33) and

$$\nu_j^{-1}R_{j,Q} + \nu^Q\tilde{R}_{j,Q} = \dots$$

This shows that R_j, \tilde{R}_j are uniquely determined as

$$(7.34) \quad \nu^Q \neq \nu_j^{-1}, \quad Q \in \mathbf{N}^p, \quad |Q| > 1, \quad 1 \leq j \leq p.$$

To verify that ρ_i are involutions, we compose by $\bar{\rho}_i^{-1}$ from right in (7.27)-(7.29), and we apply complex conjugate to the coefficients of the new identities. This results in (7.30)-(7.32) in which $(\bar{\rho}_i, f_j^*)$ are replaced by $(\overline{(\bar{\rho}_i)^{-1}}, \bar{f}_i)$. We can also start with (7.30)-(7.32) and apply the same procedure to get (7.27)-(7.29), in which $(\bar{\rho}_i, f_i)$ are replaced by $(\overline{(\bar{\rho}_i)^{-1}}, \bar{f}_i^*)$. By the uniqueness of the solutions, we conclude that $\overline{(\bar{\rho}_i)^{-1}} = \bar{\rho}_i$ as both sides have the same linear part. We now have $\overline{(\bar{\rho}_i)^{-1}}(\bar{z}') = \bar{\rho}_i(z')$. Hence, by (7.9), $\bar{z}' = \bar{\rho}_i(\rho_i(z')) = \bar{\rho}_i^2(z')$. This shows that each ρ_i is an involution. \square

We now can prove the following theorem.

Theorem 7.4. *Let M be a real analytic submanifold in \mathbf{C}^{2p} defined by (7.26) without elliptic components. Assume that in (ξ, η) coordinates, $D\sigma(0)$ is diagonal and has distinct eigenvalues $\mu_1, \dots, \mu_p, \mu_1^{-1}, \dots, \mu_p^{-1}$. Let $\nu = \nu_\epsilon$ be of the form (7.21) and satisfy (7.34). Then M admits a unique pair of formal asymptotic submanifold $\{K_1^\epsilon, K_2^\epsilon\}$ such that the complexification of K_1^ϵ in \mathcal{M} is an invariant formal submanifold \mathcal{H}_ϵ of σ that is tangent to*

$$(7.35) \quad \cap_{\epsilon_j=1}\{\eta_j = 0\} \cap \cap_{\epsilon_j=-1}\{\xi_j = 0\}.$$

Furthermore, the complexification of K_2^ϵ equals $\tau_1\mathcal{H}_\epsilon$.

Proof. Let $K_i = K_i^\epsilon$. We will follow Klingenberg's approach for $p = 1$, by using the deck transformations. Suppose that K is an attached formal complex submanifold which intersects with M at two totally real formal submanifolds K_1, K_2 . We first embed $K_1 \cup K_2$ into \mathcal{M} as M is embedded into \mathcal{M} . Let \mathcal{K}_i be the complexification of K_i in \mathcal{M} . Since ρ fixes K_i pointwise, then $\rho\mathcal{K}_i = \mathcal{K}_i$.

We want to show that $\tau_1(\mathcal{K}_1) = \mathcal{K}_2$; thus \mathcal{K}_i is invariant under σ . We can see that \mathcal{K}_i is defined by

$$(7.36) \quad \bar{\rho}_i(z') = w'.$$

On \mathcal{K}_1 , by (7.27) and (7.29) we have $\tilde{L}(z', w') + E(z', w') = -f(z')$. The latter defines a complex submanifold of dimension p . Thus it must be \mathcal{K}_1 . On \mathcal{M} ,

$$(\tilde{L}_j(z', w') + E_j(z', w'))^2 = z_{p+j}$$

are invariant by τ_1 . Thus each $\tilde{L}_j(z', w') + E_j(z', w')$ is either invariant or skew-invariant by τ_1 . Computing the linear part, we conclude that they are all skew-invariant by τ_1 . Hence $\tau_1(\mathcal{K}_1)$ is defined by $\tilde{L}(z', w') + E(z', w') = f(z')$, which is the defining equations for \mathcal{K}_2 . We must identify the tangent space of \mathcal{K}_1 at the origin. Let us verify (7.35) for $\epsilon_j = 1$ for all j , while the general case is analogous. Let A, S be the linear parts of $\bar{\rho}_1$ and $\sigma = \tau_1\tau_2$. Define $e(z', w') = w' - A(z')$. The tangent space to \mathcal{K}_1 at the origin is $\{e(z', w') = 0\}$. From the proof of Lemma 7.3 (ii), the matrix of A is

$$\mathbf{A} = \begin{pmatrix} \text{diag}(a_h) & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & \text{diag}(a_s) \\ \mathbf{0} & \text{diag}(b_s) & \mathbf{0} \end{pmatrix} = \begin{pmatrix} -\text{diag}(\lambda_h) & \mathbf{0} & \mathbf{0} \\ \mathbf{0} & \mathbf{0} & -\mathbf{I} \\ \mathbf{0} & -\mathbf{I} & \mathbf{0} \end{pmatrix}.$$

Thus $e_h = w_h + \lambda_h z_h$, $e_s = w_s + z_{s+s^*}$, and $e_{s+s^*} = w_{s+s^*} + z_s$. Using the formulas (2.20) and (2.24) of τ_1, τ_2 when M is the product quadric, we can verify that $e_{s+s^*} \circ S = \bar{\mu}_s e_{s+s^*}$, $e_s \circ S = \mu_s^{-1} e_s$, and $e_h \circ S = \bar{\mu}_h e_h$. Therefore, $e_j(z', w') = c_j \eta_j$.

Finally, if \mathcal{K}_1 is convergent, then (7.36) implies that $\bar{\rho}_1$ is convergent. Hence K_1 , the fixed point set of ρ_1 , is convergent. \square

We now study the convergence of attached formal submanifolds. Let us first recall a theorem of Pöschel [Pös86]. Let ν and ϵ be as in (7.21). Define

$$\omega_\nu(k) = \min_{1 < |P| \leq 2^k, P \in \mathbf{N}^p} \min_{1 \leq i \leq p} \{ |\nu^P - \nu_i|, |\nu^P - \nu_i^{-1}| \}.$$

Suppose that

$$(7.37) \quad - \sum \frac{\log \omega_\nu(k)}{2^k} < \infty.$$

Then the unique invariant formal submanifold of σ that is tangent to the \mathcal{H}_ϵ defined by (7.35) is convergent.

We now obtain a consequence of Theorem 7.4 and Pöschel's theorem.

Theorem 7.5. *Let M be as in Theorem 7.4. Let $\nu = \mu_\epsilon$ be given by (7.21). Assume that $\nu = (\mu_1^{\epsilon_1}, \dots, \mu_p^{\epsilon_p})$ satisfy (7.37). Then M admits an attached complex submanifold.*

To study the convergence of all attached formal manifolds, we use a theorem in [Sto15] to conclude simultaneous convergence of all attached formal submanifolds. In fact the conclusion is much more stronger. It is based on the simultaneous linearization of the σ_j 's on the resonant ideal, i.e. the ideal \mathcal{I} generated by $\xi_1 \eta_1, \dots, \xi_p \eta_p$. Define with $D\sigma_i(0) := \text{diag}(\mu_{i,1}, \dots, \mu_{i,n})$,

$$(7.38) \quad \omega_{\mathcal{S}, \mathcal{I}}(k) = \inf \left\{ \max_{1 \leq i \leq l} |\mu_i^Q - \mu_{i,j}| \neq 0 : 2 \leq |Q| \leq 2^k, 1 \leq j \leq n, Q \in \mathbf{N}^n, x^Q \notin \mathcal{I} \right\}$$

where $\mu_i^Q := \mu_{i,1}^{q_1} \cdots \mu_{i,n}^{q_n}$. As in [Sto15], we say that the family \mathcal{S} is *Diophantine* on \mathcal{I} , if the sequence of numbers (7.38) satisfies (7.37).

Theorem 7.6. *Let M be as in Theorem 7.5, given by (7.26). Assume furthermore that M is non resonant. Suppose that \mathcal{S} is Diophantine on \mathcal{I} or that M has an abelian CR singularity of pure complex type. Then all attached formal submanifolds are convergent. Moreover, the complex submanifolds K_j attached to pairs of antiholomorphic involutions $\{\rho_{j1}, \rho_{j2}\}$ have the form*

$$(7.39) \quad K_j: z_{p+i} = (L_i(z', \bar{\rho}_{j1}(z')) + E_i(z', \bar{\rho}_{j1}(z')))^2, \quad 1 \leq i \leq p.$$

Proof. When M has an abelian CR singularity of pure complex type, from the normal form of τ_{ij} in Theorem 5.3 we know that all invariant submanifolds of σ that are tangent to (7.35) are convergent. The non-abelian CR singularity case is a consequence of the theorem of simultaneous linearization of the σ_j 's along the resonant ideal \mathcal{I} [Sto15][theorem 2.1] and Theorem 7.4. Since, in good holomorphic coordinates, σ is linear on the zero set of the resonant ideal, the solutions $\{\bar{\rho}_1, \bar{\rho}_2\}$ to (7.27)-(7.32) are linear and there are $2^{h_*+s_*-1}$ pairs $\{\bar{\rho}_{j1}, \bar{\rho}_{j2}\}$ of solutions. The equation (7.39) is derived in (7.10) for a general situation. \square

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