# ON THE EXISTENCE OF KAM TORI FOR PRESYMPLECTIC VECTOR FIELDS

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ABSTRACT. We prove the existence of a torus that is invariant with respect to the flow of a vector field that preserves the presymplectic form in an exact presymplectic manifold. The flow on this invariant torus is conjugate to a linear flow on a torus with a Diophantine velocity vector. The proof has an "a posteriori" format, the the invariant torus is constructed by using a Newton method in a space of functions, starting from a torus that is approximately invariant. The geometry of the problem plays a major role in the construction by allowing us to construct a special adapted basis in which the equations that need to be solved in each step of the iteration have a simple structure. In contrast to the classical methods of proof, this method does not assume that the system is close to integrable, and does not rely on using action-angle variables.

### 1. INTRODUCTION

The goal of this paper is to give a proof of the existence of a torus that is invariant with respect to the flow of a presymplectic vector field V in an exact presymplectic manifold  $(\mathcal{P}, \Omega)$ , i.e., in a manifold  $\mathcal{P}$  endowed with an exact constant-rank 2-form  $\Omega$  that is preserved under the flow of V.

Perhaps the most prominent occurrence of presymplectic manifolds in physics is 9 in the geometric theory of dynamical systems with constraints. These are systems 10 for which the transition from Lagrangian to Hamiltonian description is non-trivial 11 because some of the relations  $p_A := \frac{\partial L}{\partial \dot{q}^A}(q, \dot{q})$  expressing the generalized momenta  $p_A$  in terms of generalized velocities  $\dot{q}^A$  cannot be solved for  $\dot{q}^A$  since the matrix 12 13  $\left(\frac{\partial^2 L}{\partial \dot{q}^A \partial \dot{q}^B}\right)$  is degenerate; the relations that cannot be solved play the role of con-14 straints. The modern theory of constrained systems was initiated in the early 1950s 15 by Dirac [23, 24, 25] and developed by Bergmann and his collaborators for purposes 16 of quantization of field theories [4, 44, 10] (the book [52] offers an in-depth exposi-17 tion). Such situations occur in also classical electromagnetic theory [53, Ch. V], in 18 the description of relativistic particles [36], [53, Ch. VII], gauge fields [53, 38, 47]. 19 A geometric theory of constrained systems was proposed in the late 1970s by 20 Gotay and collaborators [35, 34, 33]. In their approach, the system is transformed 21 in stages, and the process ends up with a manifold that is typically presymplectic. 22 Presymplectic geometry is also related to equivalence between Lagrangian and 23 Hamiltonian formalisms for constrained systems [15, 8, 7], geometric approach to 24 maximum principles [6], geometric optics [18, 17, 27], etc. Other topics of inter-25 est are canonical transformations in presymplectic systems [14, 16], reduction of 26 presymplectic manifolds [22, 29, 28, 45, 2], etc. 27

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A major achievement in the theory of Hamiltonian systems in the second half of the XX century was the celebrated Kolmogorov-Arnold-Moser (KAM) Theorem. This theorem is the subject of multiple reviews and pedagogical expositions (see, e.g., [41, 46, 51, 50, 19, 20, 5, 11]), and its history is beautifully described in the recent book by Dumas [26].

The standard proofs of the KAM-type theorems perform an infinite sequence of 6 canonical transformations to convert a slightly perturbed integrable system on a symplectic manifold into action-angle variables. González, Jorba, de la Llave and 8 Villanueva developed a new method of proof in their seminal 2005 paper [21]. This 9 method relies heavily on the geometry of the system. One important ingredient in 10 their proof is the so-called *automatic reducibility*: if  $\mathcal{K}$  is a torus in the symplectic 11 manifolds  $\mathcal{P}$  that is invariant with respect to a map  $f: \mathcal{P} \to \mathcal{P}$ , then the tangent 12 bundle to  $\mathcal{K}$  is preserved under the derivative T f. This simplifies the structure of 13 the coefficient matrices in certain difference equations which, in turn, dramatically 14 simplifies the solution of the problem. 15

Methods similar to the ones developed in [21] have since been used in [32] to study the existence of non-twist tori in degenerate Hamiltonian systems, and in [30, 42] to prove the existence of lower dimensional invariant tori that are partially hyperbolic or elliptic. Since these methods are suitable for efficient numerical implementation, they have been used for this purpose in [12, 39, 13, 31]. Many aspects of these methods are considered in the recent book by Haro *et al* [37] (KAM theory is the subject of Chapter 4).

Alishah and de la Llave [3] used the ideas of [21] to prove a KAM theorem for presymplectic systems, for which the degeneracy of the presymplectic form complicates the matters. They considered a family  $\{f_{\lambda}\}$  of maps that preserve the presymplectic form, and found a value  $\bar{\lambda}$  of the parameter  $\lambda$  and an embedding K from a torus to the presymplectic manifold such that  $f_{\bar{\lambda}} \circ K = K \circ T_{\omega}$  where  $T_{\omega} : \theta \mapsto \theta + \omega$  is translation on the torus by a Diophantine vector  $\omega$ .

The main goal of this paper is to prove a KAM theorem for a family  $\{V_{\lambda}\}$  of 29 presymplectic vector fields on an exact presymplectic manifold  $(\mathcal{P}, \Omega)$ , with dim  $\mathcal{P} =$ 30 d + 2n, dim ker  $\Omega = d$ ,  $\Omega = d\tau$  for some  $\tau \in \Omega^1(\mathcal{P})$ . For most of the paper we consider  $\mathcal{P} \cong \mathbb{T}^d \times T^* \mathbb{T}^n \cong \mathbb{T}^{d+n} \times \mathbb{R}^n$ , where ker  $\Omega$  coincides with the first d31 32 dimensions. Our goal is to find a value  $\bar{\lambda}$  of the parameter  $\lambda \in \mathbb{R}^{d+2n}$  and an 33 embedding  $K: \mathbb{T}^{d+n} \to \mathcal{P}$  such that the submanifold  $\mathcal{K} := K(\mathbb{T}^{d+n})$  is invariant 34 with respect to the flow  $\Phi_{\bar{\lambda},t}$  of the vector field  $V_{\bar{\lambda}}$ , and K conjugates  $\Phi_{\bar{\lambda},t}$  to the 35 linear flow  $\phi_t : \mathbb{T}^{d+n} \to \mathbb{T}^{d+n} : \theta \mapsto \theta + t\omega$ , where  $\omega \in \mathbb{R}^d$  is a constant Diophantine 36 vector: 37

$$\Phi_{\bar{\lambda},t} \circ K = K \circ \phi_t , \qquad t \ge 0 . \tag{1.1}$$

The infinitesimal form of (1.1) is  $V_{\bar{\lambda},K(\theta)} = K_{*\theta} \omega_{\theta}$ , where  $K_{*\theta} : T_{\theta} \mathbb{T}^{d+n} \to T_{K(\theta)} \mathcal{P}$ is the derivative of K at  $\theta \in \mathbb{T}^{d+n}$  and we consider  $\omega \in \mathbb{R}^{d+n}$  as  $\omega_{\theta} \in T_{\theta} \mathbb{T}^{d+n} = \mathbb{R}^{d+n}$ .

Our proof of the theorem has an *a posteriori* format (as in [21]). In more detail, we assume the existence of  $\lambda_0$  and  $K_0 : \mathbb{T}^{d+n} \to \mathcal{P}$  that satisfy (1.1) only approximately, i.e.,  $\Phi_{\lambda_0,t} \circ K_0 \approx K_0 \circ \phi_t$ . Then we start a version of the Newton method to construct iteratively a sequence of better and better approximations

$$(\lambda_0, K_0) \mapsto (\lambda_1, K_1) \mapsto (\lambda_2, K_2) \mapsto (\lambda_3, K_3) \mapsto \cdots$$
(1.2)

whose limit  $(\lambda_{\infty}, K_{\infty})$  is the desired solution  $(\overline{\lambda}, K)$  satisfying (1.1). The loss of domain accompanying each step of the iteration is compensated by the fast convergence of the Newton method. This method is convenient to implement in numerical computations. Moreover, *a posteriori* theorems are suitable for validation of numerical results, i.e., they can be used to produce computer assisted proofs of existence of invariant manifolds.

In the paper we extensively use the geometry of the system to our advantage, 7 inspired by the ideas of [21, 3]. If  $\mathcal{K} = K(\mathbb{T}^{d+n})$  is the invariant torus, then at 8 each point  $k \in \mathcal{K}$  the kernel ker  $\Omega_k$  of the presymplectic form is a subspace of  $T_k \mathcal{K}$ , 9 so that we have ker  $\Omega_p \subseteq T_k \mathcal{K} \subseteq T_k \mathcal{P}$ . In fact, we have much more – a filtration of subbundles ker  $\Omega \subseteq T\mathcal{K} \subseteq T\mathcal{P}|_{\mathcal{K}}$  which is invariant with respect to the flow 10 11  $\Phi_{\bar{\lambda},t}$  of the vector field  $V_{\bar{\lambda}}$ . This and the invariance (1.1) allow us to construct a 12 special basis adapted to the filtration, in which the matrices of the operators have 13 zero blocks. Even if the torus is only approximately invariant (as in the case of 14  $(\lambda_i, K_i)$ ), these blocks, albeit non-zero, are small and we have good bounds on their 15 16 norms.

An important role is also played by the fact that  $\mathcal{K}$  is isotropic (i.e., that the the pull-back of  $\Omega$  to  $\mathcal{K}$  vanishes identically), and the approximately invariant tori are approximately isotropic. We found the following interesting quotations related to this fact. On page 45 of his classic 1973 monograph [43], Moser writes

Actually, more than asserted in Theorem 2.7 can be proven. It turns

out that the differential form  $\sum_{k=1}^{n} dy_k \wedge dx_k$  vanishes identically

on the tori (3.11), and one calls manifolds with this property and

<sup>24</sup> of maximal dimension Lagrange manifolds.

In this quotation, Theorem 2.7 is (as Moser calls it) the Kolmogorov-Arnold Theorem, and the tori (3.11) are the invariant tori whose existence is proved in the KAM theorem. On page 584 of their monumental book [1], Abraham and Marsden write

29 Moser [1973a] states that the invariant tori are Lagrangian sub-30 manifolds [...]. This fact can probably be exploited, although to

our knowledge it has not been.

The fact that  $\mathcal{K}$  is isotropic and of a maximum dimension (i.e., Lagrangian in the symplectic case) is crucial for the proofs in [21, 3] and in this paper.

The paper is organized as follows. In Sections 2.1-2.5 we introduce some basic 34 definitions and notations, discuss the integrability of the distribution ker  $\Omega$  and con-35 struct the symplectic manifold  $\mathcal{P}/\ker \Omega$ . The main theorem is stated in Section 2.6. 36 In Section 3 we study the geometric structures occurring when we know the true 37 38 solution  $(\lambda, K)$  of the problem – we prove that  $\mathcal{K}$  is isotropic (Section 3.1), give a detailed construction of the basis adapted to the filtration ker  $\Omega \subseteq T\mathcal{K} \subseteq T\mathcal{P}|_{\mathcal{K}}$ 39 (Section 3.2), and study the properties of the matrix of transition from a general 40 basis of  $T\mathcal{P}|_{\mathcal{K}}$  to the special adapted basis (Section 3.3). 41

Section 4 is devoted to the properties of approximate solutions. Approximately isotropic tori are studied in Section 4.1. In Section 4.2 we derive an equation for the corrections  $\varepsilon_j$  and  $\Delta_j$  needed to obtain a better approximation  $\lambda_{j+1} =$  $\lambda_j + \varepsilon_j$ ,  $K_{j+1} = K_j + \Delta_j$ . We solve this equation in Section 4.3, relying heavily on the machinery developed in Section 3. Finally, in Section 5, we collect several lemmata that justify the applicability of the Newton method for performing the iteration (1.2).

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#### 2. Preliminaries and general setup

<sup>2</sup> In this section we set up the problem and state our main result.

# <sup>3</sup> 2.1. Presymplectic manifolds and vector fields.

4 **Definition 2.1.** A presymplectic manifold is a pair  $(\mathcal{P}, \Omega)$ , where  $\mathcal{P}$  is a manifold 5 of any (finite) dimension and  $\Omega \in \Omega^2(\mathcal{P})$  is a closed 2-form with constant rank. If 6  $\Omega$  is exact, i.e., if  $\Omega = d\tau$  for some  $\tau \in \Omega^1(\mathcal{P})$ , then we say that  $(\mathcal{P}, \Omega)$  is an exact 7 presymplectic manifold.

8 Throughout this paper, we will always assume that

$$\dim \mathcal{P} = d + 2n , \qquad \text{rank } \Omega = 2n . \tag{2.1}$$

9 Most of the time we will consider the specific exact presymplectic manifold

$$\mathcal{P} := \mathbb{T}^d \times T^* \mathbb{T}^n \cong \mathbb{T}^d \times \mathbb{T}^n \times \mathbb{R}^n \tag{2.2}$$

- with an exact presymplectic form  $\Omega$  with ker  $\Omega = \mathbb{T}^d$ . We assume that  $\mathbb{T}^d \times T^* \mathbb{T}^n$  is endowed with an Euclidean structure, so that we can identify two-forms with linear operators and abstract tangent vectors with column vectors.
- In the definition below,  $\mathfrak{X}(\mathcal{P})$  stands for the vector fields on  $\mathcal{P}$ ,  $\mathcal{L}$  is the Lie derivative, and  $\iota$  is the interior product, i.e., the contraction with a vector field.
- **Definition 2.2.** Let  $V \in \mathfrak{X}(\mathcal{P})$  and  $\Phi_t : \mathcal{P} \to \mathcal{P}$  be the time-t flow of V. The vector field V is said to be presymplectic if  $\Phi_t^* \Omega = \Omega$  for all  $t \in \mathbb{R}$ .

**Lemma 2.3.** Let  $(\mathcal{P}, \Omega)$  be a presymplectic manifold, and  $V \in \mathfrak{X}(\mathcal{P})$ . Then the following conditions are equivalent:

- 19 (a) V is a presymplectic vector field;
- (b) the Lie derivative of the presymplectic form along V vanishes:  $\mathcal{L}_V \Omega = 0$ ;
- 21 (c) the 1-form  $\iota_V \Omega$  is closed.

22 Proof. The equivalence of (a) and (b) comes directly from the definition of a Lie 23 derivative, and the equivalence of (b) and (c) follows from Cartan's magic formula 24 and the closedness of  $\Omega: 0 = \mathcal{L}_V \Omega = \iota_V d\Omega + d(\iota_V \Omega) = d(\iota_V \Omega).$ 

25 2.2. Foliation induced by ker  $\Omega$ . In this section we discuss some results about 26 a general presymplectic manifold  $(\mathcal{P}, \Omega)$  (not necessarily (2.2)). For any  $p \in \mathcal{P}$ , 27 define

$$\ker \Omega_p := \{ W_p \in T_p \mathcal{P} \mid \iota_{W_p} \Omega_p = 0 \}$$
$$= \{ W_p \in T_p \mathcal{P} \mid \Omega_p(W_p, U_p) = 0 \ \forall U_p \in T_p \mathcal{P} \} \subseteq T_p \mathcal{P} .$$

<sup>28</sup> The subspaces ker  $\Omega_p$  form a differentiable distribution, ker  $\Omega$ , of rank d. Define

$$\mathfrak{X}^{\ker\Omega}(\mathcal{P}) := \{ W \in \mathfrak{X}(\mathcal{P}) \mid \iota_W\Omega = 0 \} = \{ W \in \mathfrak{X}(\mathcal{P}) \mid W_p \in \ker\Omega_p \; \forall \, p \in \mathcal{P} \}$$

<sup>29</sup> Using the classical Frobenius Theorem (see, e.g., [48, Section 3.5], or, especially,

- <sup>30</sup> [40, Appendix 3]) and the fact that  $\Omega$  is closed, one can easily obtain
- **Lemma 2.4.** If  $\Omega$  is presymplectic, the distribution ker  $\Omega$  is integrable.
- Lemma 2.4 implies that  $\mathcal{P}$  has a foliation with *d*-dimensional leaves such that the tangent space to the leaf through  $p \in \mathcal{P}$  at p is ker  $\Omega_p$ . We make the additional assumption that the collection of leaves of the foliation forms a smooth manifold  $\mathcal{Q}$ (for a discussion see, e.g., [40, Sec. 4.3.3 of Appendix 3]). Let  $\pi^{\mathcal{Q}} : \mathcal{P} \to \mathcal{Q}$  be the

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<sup>1</sup> canonical projection taking each point  $p \in \mathcal{P}$  to the leaf  $(\pi^{\mathcal{Q}})^{-1}(\pi^{\mathcal{Q}}(p))$  through it. <sup>2</sup> If  $\pi^{\mathcal{Q}}_*: T\mathcal{P} \to T\mathcal{Q}$  is the derivative of  $\pi^{\mathcal{Q}}$ , then ker  $\pi^{\mathcal{Q}}_* = \ker \Omega$ . The lemma below <sup>3</sup> (see [40, Section III.7]) states that  $\mathcal{Q}$  carries a natural symplectic structure.

4 Lemma 2.5. Let  $(\mathcal{P}, \Omega)$  be a presymplectic manifold such that ker  $\Omega$  determines an 5 integrable foliation of  $\mathcal{P}$  whose leaves form a smooth manifold  $\mathcal{Q}$ , and let  $\pi^{\mathcal{Q}} \colon \mathcal{P} \to$ 6  $\mathcal{Q}$  be the canonical projection. Then there exists a unique symplectic form  $\widetilde{\Omega} \in$ 7  $\Omega^2(\mathcal{Q})$  on the manifold  $\mathcal{Q}$  such that  $(\pi^{\mathcal{Q}})^*\widetilde{\Omega} = \Omega$ .

For the case (2.2) considered in this paper,  $\mathcal{Q} \cong T^* \mathbb{R}^n$ . It would be interesting to investigate the case when  $\mathcal{P}$  has a more complicated structure than (2.2).

<sup>10</sup> 2.3. Matrix representation of  $\Omega$  and  $\widetilde{\Omega}$ . We consider the case when  $\mathcal{P}$  has <sup>11</sup> product structure (2.2), so ker  $\Omega_p \cong \mathbb{T}^d$  for every  $p \in \mathcal{P}$  and the collection of leaves, <sup>12</sup>

$$\mathcal{Q} = \mathcal{P}/\ker \Omega = T^* \mathbb{T}^n , \qquad (2.3)$$

<sup>13</sup> is a symplectic manifold with symplectic form  $\Omega$ . Because of the assumed Euclidean <sup>14</sup> structure on  $\mathcal{P}$ , we can identify a 2-form on  $\mathcal{P}$  with a linear operator. For any  $p \in \mathcal{P}$ , <sup>15</sup> let  $J_p: T_p\mathcal{P} \to T_p\mathcal{P}$  be the linear operator corresponding to  $\Omega_p$ , defined by

$$\langle U_p, J_p W_p \rangle_{\mathbb{R}^{d+2n}} := \Omega_p(U_p, W_p) , \qquad U_p, W_p \in T_p \mathcal{P} \cong \mathbb{R}^{d+2n} , \qquad (2.4)$$

<sup>16</sup> where  $\langle \cdot, \cdot \rangle_{\mathbb{R}^{d+2n}}$  is the Euclidean inner product on  $\mathbb{R}^{d+2n}$ . Similarly, for any  $q \in \mathcal{Q}$ , <sup>17</sup> let  $\widetilde{J}_q : T_q \mathcal{Q} \to T_q \mathcal{Q}$  be the linear operator corresponding to the symplectic form <sup>18</sup>  $\widetilde{\Omega}_q$  on  $T_q \mathcal{Q}$  at  $q \in \mathcal{Q}$ : if  $\langle \cdot, \cdot \rangle_{\mathbb{R}^{2n}}$  is the Euclidean inner product on  $\mathbb{R}^{2n}$ , then

$$\langle \widetilde{U}_q, \widetilde{J}_q \widetilde{W}_q \rangle_{\mathbb{R}^{2n}} := \widetilde{\Omega}_q \left( \widetilde{U}_q, \widetilde{W}_q \right) , \qquad \widetilde{U}_q, \widetilde{W}_q \in T_q \mathcal{Q} \cong \mathbb{R}^{2n} .$$
 (2.5)

If we choose a basis for  $T_p \mathcal{P} \cong \mathbb{R}^d \times \mathbb{R}^{2n}$  such that the first d vectors form a basis of  $\mathbb{R}^d = T_p \ker \Omega$ , and the other 2n vectors form a basis of  $\mathbb{R}^{2n} \cong T_p T^* \mathbb{T}^n$ , then we can write  $J_p$  in a matrix form as

$$J_p = \begin{bmatrix} 0 & 0\\ 0 & \widetilde{J}_{\pi^{\mathcal{Q}}(p)} \end{bmatrix}, \qquad J_p^{\top} = -J_p , \quad \widetilde{J}_q^{\top} = -\widetilde{J}_q .$$
(2.6)

<sup>22</sup> Although  $J_p$  is not invertible, we define  $J_p^{-1}$  as the Moore-Penrose pseudoinverse:

$$J_{p}^{-1} := \begin{bmatrix} 0 & 0 \\ 0 & \tilde{J}_{\pi^{\mathcal{Q}}(p)}^{-1} \end{bmatrix}, \quad \text{so that} \quad J_{p} J_{p}^{-1} = \begin{bmatrix} 0 & 0 \\ 0 & \mathbb{I}_{2n} \end{bmatrix}. \quad (2.7)$$

# 23 2.4. Miscellaneous definitions and results.

**Definition 2.6.** For  $\gamma > 0$  and  $\sigma \ge d + n - 1$ , the set of all  $\omega \in \mathbb{R}^{d+n}$  satisfying

$$|\omega \cdot k| \ge \frac{\gamma}{|k|^{\sigma}} \qquad \forall \ k \in \mathbb{Z}^{d+n} \setminus \{0\}$$

$$(2.8)$$

- is called the set of Diophantine vectors and is denoted by  $\mathcal{D}(\gamma, \sigma)$ .
- For any  $\rho > 0$ , we define the torus "thickened" into the complex direction,

$$\mathbb{T}_{\rho}^{d+n} := \{ \theta \in \mathbb{C}^{d+n} / \mathbb{Z}^{d+n} \mid |\text{Im } \theta^{\alpha}| \le \rho, \ \alpha = 1, 2, \dots, d+n \} .$$

$$(2.9)$$

- 1 Let || stand for the supremum norm on  $\mathbb{R}^m$  or  $\mathbb{C}^m$  (for any m). Given  $\rho > 0$ , we 2 define the set of functions  $\mathcal{W}_{\rho}$  as follows:
  - $\mathcal{W}_{\rho} := \left\{ K : \mathbb{T}_{\rho}^{d+n} \to \mathcal{P} \mid \text{(a) } K \text{ is real analytic on the interior of } \mathbb{T}_{\rho}^{d+n} , \\ \text{(b) } K \text{ is continuous on the boundary of } \mathbb{T}_{\rho}^{d+n} , \quad (2.10) \right\}$

(c) K is periodic of period 1 in all arguments  $\}$ .

<sup>3</sup> Define  $||K||_{\rho} = \sup_{\theta \in \mathbb{T}_{\rho}^{d+n}} |K(\theta)|$ , so that  $(\mathcal{W}_{\rho}, || ||_{\rho})$  is a Banach space. For analytic <sup>4</sup> functions  $g : \mathcal{B} \to \mathbb{C}$  (where  $\mathcal{B} \subseteq \mathbb{C}$ ), and any  $\ell \in \mathbb{N}$ , define the norms

$$|g|_{C^{\ell},\mathcal{B}} := \sup_{0 \le |\mathbf{k}| \le \ell} \sup_{z \in \mathcal{B}} \left| D^{\mathbf{k}} g(z) \right| .$$

$$(2.11)$$

5 Lemma 2.7. (Cauchy bound) For  $K \in W_{\rho}$  and  $0 < \delta < \rho$ ,

$$\|DK\|_{\rho-\delta} \le C\delta^{-1} \|K\|_{\rho}$$
 . (2.12)

- 6 Lemma 2.8. (Rüssmann [49]) Let  $\omega = \begin{bmatrix} \omega^1 & \omega^2 & \cdots & \omega^{d+n} \end{bmatrix}^\top \in \mathcal{D}(\gamma, \sigma)$  and let 7 the function  $h: \mathbb{T}^{d+n} \to \mathcal{P}$  be analytic on  $\mathbb{T}_{\rho}^{d+n}$  and have zero average. Then for 8 any  $0 < \delta < \rho$ , the differential equation  $\partial_{\omega} v = h$ , where  $\partial_{\omega} := \sum_{\alpha=1}^{d+n} \omega^{\alpha} \frac{\partial}{\partial \theta^{\alpha}}$  is
- 9 the directional derivative in the direction of  $\omega$ , has a unique average-zero solution 10  $v: \mathbb{T}^{d+n} \to \mathcal{P}$  which is analytic in  $\mathbb{T}^{d+n}_{\rho-\delta}$ . The solution v satisfies the estimate

$$\|v\|_{\rho-\delta} < C\gamma^{-1}\delta^{-\sigma} \|h\|_{\rho} , \qquad (2.13)$$

<sup>11</sup> where C is a constant depending only on d, n, and  $\sigma$ .

<sup>12</sup> 2.5. General setup and matrix notations. Let  $\{V_{\lambda}\}$  be a (d + 2n)-parameter <sup>13</sup> family of presymplectic vector fields on the exact presymplectic manifold  $(\mathcal{P}, \Omega)$ .

14 Our goal is to construct a smooth embedding

$$K: \mathbb{T}^{d+n} \to \mathcal{P} \tag{2.14}$$

such that  $\mathcal{K} := K(\mathbb{T}^{d+n})$  be invariant with respect to the flow  $\Phi_{\bar{\lambda},t}$  of  $V_{\bar{\lambda}}$  for some value  $\bar{\lambda}$  and the flow  $\Phi_{\bar{\lambda},t}$  on  $\mathcal{K}$  be conjugate to the linear flow  $\phi_t$  on  $\mathbb{T}^{d+n}$ :

$$\Phi_{\bar{\lambda},t}(K(\theta)) = K(\phi_t(\theta)) \qquad \forall t \in \mathbb{R} , \quad \forall \theta \in \mathbb{T}^{d+n} ,$$
(2.15)

where  $\omega \in \mathbb{R}^{d+n}$  is a constant Diophantine vector and  $\phi_t : \mathbb{T}^{d+n} \to \mathbb{T}^{d+n} : \theta \mapsto \theta + t\omega$ . (For elements of  $\mathbb{T}^{d+n}$ , e.g.,  $\theta + t\omega$ , we assume that we take the fractional part of each component.) Differentiate (2.15) with respect to t and set t = 0 to obtain

$$V_{\bar{\lambda},K(\theta)} = K_{*\theta}\,\omega_{\theta} \qquad \forall \,\theta \in \mathbb{T}^{d+n} \ . \tag{2.16}$$

- <sup>21</sup> Here  $V_{\bar{\lambda},K(\theta)} \in T_{K(\theta)}\mathcal{K} \subseteq T_{K(\theta)}\mathcal{P}$  is the value of  $V_{\bar{\lambda}}$  at  $K(\theta) \in \mathcal{K}, \ \omega_{\theta} \in T_{\theta}\mathbb{T}^{d+n}$ <sup>22</sup> is the Diophantine vector  $\omega$  considered as an element of  $T_{\theta}\mathbb{T}^{d+n} = \mathbb{R}^{d+n}$ , and <sup>23</sup>  $K_{*\theta}: T_{\theta}\mathbb{T}^{d+n} \to T_{K(\theta)}\mathcal{K} \subseteq T_{K(\theta)}\mathcal{P}$  is the derivative of K at  $\theta$ .
- Instead of the differential-geometric notations used in (2.16), we will normally use matrix notations. In these notations (2.16) reads

$$V_{\bar{\lambda},K(\theta)} = DK_{\theta}\,\omega$$
 or  $V_{\bar{\lambda},K(\theta)} = \partial_{\omega}K_{\theta}$ , (2.17)

where  $\omega$  is considered as a constant column vector,  $\partial_{\omega} := \omega \cdot \nabla$ , and

$$DK_{\theta} := \left[ \left( DK_{\theta} \right)^{A}{}_{\alpha} \right] = \left[ \frac{\partial K^{A}}{\partial \theta^{\alpha}}(\theta) \right] \in \mathcal{M}_{d+2n,d+n}(\mathbb{R})$$

<sup>27</sup> Throughout the paper we will systematically use the notations collected in Table 1.

Coordinates	Range of indices	Description
$x = \left(x^A\right)$	$A, B = 1, 2, \dots, d + 2n$	Coordinates in $\mathcal{P} \cong \mathbb{T}^d \times T^* \mathbb{T}^n$
$\underline{x} = (\underline{x}^{\mu})$	$\mu,\nu=1,2,\ldots,d$	Coordinates in $\mathbb{T}^d$ (the first $d$ in $\mathcal{P}$ )
$\widetilde{x} = \left(\widetilde{x}^i\right)$	$i, j = 1, 2, \ldots, 2n$	Coordinates in $T^*\mathbb{T}^n$ (the last $2n$ in $\mathcal{P}$ )
$\theta = (\theta^\alpha)$	$\alpha, \beta = 1, 2, \dots, d+n$	Coordinates in $\mathbb{T}^{d+n}$

TABLE 1. Notations for indices and coordinates.

### <sup>1</sup> 2.6. Statement of the main theorem. Here we state our

### <sup>2</sup> Main Theorem. Assume that:

- 3 1)  $\omega \in \mathcal{D}(\gamma, \sigma)$  is a Diophantine vector;
- 4 2)  $\mathcal{P} = \mathbb{T}^d \times T^* \mathbb{T}^n;$
- 5 3) Ω is an exact presymplectic form on P of rank 2n such that the kernel of Ω
  6 coincides with the first d directions;
- 7 4)  $\{V_{\lambda}\}$  is a (d+2n)-parameter family of analytic presymplectic vector fields on  $\mathcal{P}$ ;
- 8 5)  $K_0: \mathbb{T}^{d+n} \to \mathcal{P}$  is an embedding belonging to the class  $\mathcal{W}_{\rho_0}$  (2.10);
- 9 6) the value  $\lambda_0$  of the parameter  $\lambda$  is such that the pair  $(\lambda_0, K_0)$  is non-degenerate 10 in the sense of Definition 4.7;
- 11 7) each vector field from the family  $\{V_{\lambda}\}$  can be holomorphically extended to some 12 complex neighborhood  $\mathcal{B}_r$  of  $K_0(\mathbb{T}_o^{d+n})$ , where

$$\mathcal{B}_r := \left\{ z \in \mathbb{C}^{d+2n} \mid \exists \theta \in \mathbb{T}_{\rho_0}^{d+n} \text{ such that } |z - K_0(\theta)| < r \right\}, \qquad (2.18)$$

13 for some r > 0 and such that  $|V_{\lambda}|_{C^2, \mathcal{B}_r}$  is finite.

14 Define the error function as  $e_{0,\theta} := V_{\lambda,K_0(\theta)} - \partial_{\omega}K_{0,\theta}$ . Then there exists a constant 15 c > 0 which depends on  $d, n, \sigma, \rho_0, \|DK_0\|_{\rho_0}, r, \|V_{\lambda}\|_{C^2,\mathcal{B}_r}, \|\frac{\partial V_{\lambda}}{\partial \lambda}\|_{\lambda=\lambda_0} \circ K_0\|_{\rho_0},$ 16 and  $|\{\operatorname{avg}(\Lambda_0)\}^{-1}|$  (see (4.32)), such that if  $0 < \delta_0 < \max\{1, \frac{\rho_0}{12}\}$  and the error  $e_0$ 17 satisfies

$$||e_0||_{\rho_0} \le \min\left\{\gamma^4 \delta_0^{4\sigma}, \, cr\gamma^2 \delta_0^{2\sigma} \, ||e_0||_{\rho_0}\right\} ,$$

then there exists a mapping  $K \in \mathcal{W}_{\rho_0-6\delta_0}$  and a vector  $\bar{\lambda} \in \mathbb{R}^{d+2n}$  such that (2.16) (or, equivalently, (2.17)) is satisfied. Moreover, the following inequalities hold:

$$\|K - K_0\|_{\rho_0 - 6\delta_0} < \frac{1}{c} \gamma^2 \delta_0^{-2\sigma} \|e_0\|_{\rho_0} , \qquad \left|\bar{\lambda} - \lambda_0\right| < \frac{1}{c} \gamma^2 \delta_0^{-2\sigma} \|e_0\|_{\rho_0}$$

20

### 3. True solutions

In this section we will assume that we know the exact solution of the problem and will use the geometry of the problem in order to construct bases of  $T\mathcal{K}$  with special properties that will be utilized in Section 4.

### <sup>24</sup> 3.1. Invariant tori are isotropic.

**Definition 3.1.** Let  $\{V_{\lambda}\}$  be a (d + 2n)-parameter family of presymplectic vector fields on the exact presymplectic manifold  $(\mathcal{P}, \Omega)$ , and  $\omega \in \mathbb{R}^{d+n}$  be a Diophantine vector. If for some value  $\overline{\lambda}$  of  $\lambda$  there exists an embedding  $K: \mathbb{T}^{d+n} \to \mathcal{P}$  such that (2.17) holds, we call K and K an invariant torus (KAM torus) or a true solution.

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**Remark 1.** This terminology is somewhat of a misnomer. We require more than  $\mathcal{K}$  being merely invariant under the flow of  $V_{\bar{\lambda}}$  – we want that the motion on  $\mathcal{K}$  be quasi-periodic, (i.e., we require that the dynamics on  $\mathcal{K}$  be conjugate to a linear flow on  $\mathbb{T}^{d+n}$  with frequency  $\omega$  independent over the rationals).

5 Notational convention. In Section 3 we assume that  $\lambda = \overline{\lambda}$ , and set  $V := V_{\overline{\lambda}}$ .

6 Definition 3.2. An invariant (in the sense of Definition 3.1) torus,  $\mathcal{K} = K(\mathbb{T}^{d+n})$ , 7 in the presymplectic manifold  $(\mathcal{P}, \Omega)$  is said to be isotropic if the pull-back,  $K^*\Omega \in$ 8  $\Omega^2(\mathbb{T}^{d+n})$ , of  $\Omega \in \Omega^2(\mathcal{P})$  to the torus  $\mathbb{T}^{d+n}$  vanishes identically.

In Lemma 3.3 below we will prove that an invariant torus is isotropic. Similar results for maps are well-known for the case of submanifolds invariant with respect to symplectic or presymplectic maps (see, e.g., [21, Section 4, Lemma 1] or [3, Lemma 2.5]). This fact is crucial in the proof of Lemma 4.2 which, in turn, is essential for the bounds needed to solve the linearized equation in Section 4.3.

We introduce the linear operator  $L_{\theta}: T_{\theta}\mathbb{T}^{d+n} \to T_{\theta}\mathbb{T}^{d+n}$  as the matrix representation of the pull-back  $(K^*\Omega)_{\theta}$ : for  $\mathsf{U}_{\theta}, \mathsf{W}_{\theta} \in T_{\theta}\mathbb{T}^{d+n}$ ,

$$\langle \mathsf{U}_{\theta}, L_{\theta} \, \mathsf{W}_{\theta} \rangle_{\mathbb{R}^{d+n}} := (K^* \Omega)_{\theta} \left( \mathsf{U}_{\theta}, \mathsf{W}_{\theta} \right) = \Omega_{K(\theta)} \left( K_{*\theta} \, \mathsf{U}_{\theta}, K_{*\theta} \, \mathsf{W}_{\theta} \right) \,.$$
 (3.1)

<sup>16</sup> The explicit expression for the matrix elements of  $L_{\theta}$  is

$$L_{\theta} = DK_{\theta}^{\dagger} J_{K(\theta)} DK_{\theta} \in \mathcal{M}_{d+n,d+n}(\mathbb{R}) .$$
(3.2)

17 **Lemma 3.3.** Let  $(\mathcal{P}, \Omega)$  be an exact presymplectic manifold,  $V \in \mathfrak{X}(\mathcal{P})$  be presym-18 plectic, and  $K : \mathbb{T}^{d+n} \to \mathcal{P}$  be a true solution. Then the invariant torus  $\mathcal{K} =$ 19  $K(\mathbb{T}^{d+n})$  is isotropic (i.e.,  $K^*\Omega$  and, hence,  $L_{\theta}$ , vanish identically).

<sup>20</sup> Proof. We prove the lemma in two steps: first we use the exactness of  $\Omega$  to show <sup>21</sup> that the average (over  $\mathbb{T}^{d+n}$ ) of each matrix element of L is zero, and then we <sup>22</sup> use the ergodicity of the flow  $\theta \mapsto \theta + t\omega$  on  $\mathbb{T}^{d+n}$  to demonstrate that  $K^*\Omega$  and, <sup>23</sup> therefore, L, are constant on  $\mathbb{T}^{d+n}$ .

Since the presymplectic form  $\Omega$  is exact, there exists a 1-form  $\tau \in \Omega^1(\mathcal{P})$  such that  $\Omega = d\tau$ , hence  $K^*\Omega = K^*(d\tau) = d(K^*\tau)$ . If

$$\tau_{K(\theta)} = \sum_{A=1}^{d+2n} \tau_A(K(\theta)) \,\mathrm{d} x^A \;,$$

then the pull-back  $K^* \tau \in \Omega^1(\mathbb{T}^{d+n})$  is given by

έ

$$(K^*\tau)_{\theta} = \sum_{\alpha=1}^{d+n} C_{\alpha}(\theta) \,\mathrm{d}\theta^{\alpha} , \qquad C_{\alpha}(\theta) := \sum_{A=1}^{d+2n} \tau_A(K(\theta)) \,\frac{\partial K^A}{\partial \theta^{\alpha}}(\theta) ,$$

and the matrix representation of  $(K^*\Omega)_{\theta}$  is  $(L_{\theta})^{\alpha}{}_{\beta} = \frac{\partial C_{\alpha}}{\partial \theta^{\beta}}(\theta) - \frac{\partial C_{\beta}}{\partial \theta^{\alpha}}(\theta)$ . Because of the periodicity of the functions  $C_{\alpha} : \mathbb{T}^{d+n} \to \mathbb{R}$ ,

$$\operatorname{avg}\left(\frac{\partial C_{\alpha}}{\partial \theta^{\beta}}\right) := \int_{\mathbb{T}^{d+n}} \frac{\partial C_{\alpha}}{\partial \theta^{\beta}}(\theta) \, \mathrm{d}\theta^{1} \, \mathrm{d}\theta^{2} \cdots \, \mathrm{d}\theta^{d+n}$$

$$= \int_{\mathbb{T}^{d+n-1}} \left(\int_{\mathbb{T}^{1}} \frac{\partial C_{\alpha}}{\partial \theta^{\beta}}(\theta) \, \mathrm{d}\theta^{\beta}\right) \, \mathrm{d}\theta^{1} \cdots \widehat{\mathrm{d}\theta^{\beta}} \cdots \, \mathrm{d}\theta^{d+n} = 0$$
(3.3)

(the term  $d\theta^{\hat{\beta}}$  is missing), so avg (L) = 0 and, hence, avg  $(K^*\Omega) = 0$ .

Now we will prove that L and  $K^*\Omega$  are constant on  $\mathbb{T}^{d+n}$ . Restrict the target space of the map K (2.14) from  $\mathcal{P}$  to the image  $\mathcal{K}$  of K, to obtain the diffeomorphism  $K_{\mathbf{r}} := K|_{\mathbb{T}^{d+n} \to \mathcal{K}} : \mathbb{T}^{d+n} \to \mathcal{K}$ . Since  $\mathcal{K}$  is invariant with respect to the flow of  $V \in \mathfrak{X}(\mathcal{P})$ , at each  $k \in \mathcal{K}, V_k \in T_k \mathcal{K}$ . Therefore, the restriction  $V|_{\mathcal{K}}$  of V to  $\mathcal{K}$  can be considered as a vector field on  $\mathcal{K}$ :  $V|_{\mathcal{K}} \in \mathfrak{X}(\mathcal{K})$ . For the same reason, the Lie derivative with respect to V has a natural restriction to sections of any 6 tensor power of the tangent and cotangent bundles of  $\mathcal{K}$ . The pull-back of  $V|_{\mathcal{K}}$ by the diffeomorphism  $K_r$  is  $K_r^* V := (K_r^{-1})_* (V|_{\mathcal{K}}) \in \mathfrak{X}(\mathbb{R}^{d+n})$ . If we consider the constant  $\omega \in \mathbb{R}^{d+n}$  as a tangent vector  $\omega_{\theta} \in T_{\theta}(\mathbb{T}^{d+n})$ , then the pull-back of 9  $V_{K(\theta)} = K_{*\theta} \, \omega_{\theta} \in T_{K(\theta)} \mathcal{K}$  is 10

$$(K_{\mathbf{r}}^*V)_{\theta} = \left(K_{\mathbf{r}}^{-1}\right)_{*K(\theta)} K_{*\theta} \,\omega_{\theta} = \left(K_{\mathbf{r}}^{-1} \circ K\right)_{*\theta} \omega_{\theta} = \omega_{\theta} \ .$$

By a well-known property of the Lie derivative,  $K^* \mathcal{L}_V \Omega = \mathcal{L}_{K^*_r V} K^* \Omega = \mathcal{L}_{\omega} K^* \Omega$ 11 (where all objects and operations are restricted to  $\mathcal{K}$ ). Since  $\dot{V}$  is presymplectic, 12  $\mathcal{L}_V \Omega = 0$ , which implies that its pull-back  $K^* \Omega$  to  $\mathbb{T}^{d+n}$  is constant on the orbits 13 of the flow  $\theta \mapsto \theta + t\omega$ ,  $t \in \mathbb{R}$ . But  $\omega$  is Diophantine, hence this flow is ergodic 14 on  $\mathbb{T}^{d+n}$ , therefore  $K^*\Omega = \text{const}$  and L = const. This together with the fact that 15  $\operatorname{avg}(L) = 0$  and  $\operatorname{avg}(K^*\Omega) = 0$  implies the desired result.  $\square$ 16

3.2. Construction and properties of an adapted basis of  $T_{K(\theta)}\mathcal{P}$ . In this 17 section we construct a basis of  $(T\mathcal{P})|_{\mathcal{K}}$  that is *adapted* to the invariant (with respect 18 to the flow of the presymplectic vector field V) filtration  $(\ker \Omega)|_{\mathcal{K}} \subseteq T\mathcal{K} \subseteq (T\mathcal{P})|_{\mathcal{K}}$ 19 of vector bundles over  $\mathcal{K}$ . 20

3.2.1. Adapted coordinates in  $\mathbb{T}^{d+n}$  and a basis of  $T_{K(\theta)}\mathcal{K}$ . We first construct a basis 21 of  $T_{K(\theta)}\mathcal{K}$  (at an arbitrary point  $K(\theta) \in \mathcal{K}$ ) as a push-forward  $\left\{K_{*\theta}\left(\frac{\partial}{\partial\theta^{\alpha}}\right)_{\theta}\right\}_{\alpha=1}^{d+n}$  of 22 the basis  $\left\{ \left(\frac{\partial}{\partial \theta^{\alpha}}\right)_{\theta} \right\}_{\alpha=1}^{d+n}$  of  $T_{\theta} \mathbb{T}^{d+n}$ . Every vector in  $T_{K(\theta)} \mathcal{K}$  has the form  $K_{*\theta} \cup_{\theta}$  for some  $\bigcup_{\theta} \in T_{\theta} \mathbb{T}^{d+n}$ . The com-23

24 ponents  $(K_{*\theta} \cup_{\theta})^A$  of  $K_{*\theta} \cup_{\theta}$  in the basis  $\left\{ \left( \frac{\partial}{\partial x^A} \right)_{K(\theta)} \right\}_{A=1}^{d+2n}$  are related to the components  $\bigcup_{\theta}^{\alpha}$  of  $\bigcup_{\theta}$  in the basis  $\left\{ \left( \frac{\partial}{\partial \theta^{\alpha}} \right)_{\theta} \right\}_{\alpha=1}^{d+n}$  by 25 26

$$(K_{*\theta} \mathsf{U}_{\theta})^{A} = \sum_{\alpha=1}^{d+n} \frac{\partial K^{A}}{\partial \theta^{\alpha}}(\theta) \mathsf{U}_{\theta}^{\alpha} .$$
(3.4)

If we think of  $(K_{*\theta} \cup_{\theta})^A$  as a column vector  $K_{*\theta} \cup_{\theta} = \left[ (K_{*\theta} \cup_{\theta})^1 \cdots (K_{*\theta} \cup_{\theta})^{d+2n} \right]^+$ 27 and of the (d+n) columns of the matrix 28

$$DK_{\theta} = \left[\frac{\partial K^{A}}{\partial \theta^{\alpha}}(\theta)\right] = \left[\frac{\partial K}{\partial \theta^{1}}(\theta) \quad \cdots \quad \frac{\partial K}{\partial \theta^{d+n}}(\theta)\right] \in \mathcal{M}_{d+2n,d+n}(\mathbb{R}) , \qquad (3.5)$$

as (d+2n)-dimensional vectors, then we can interpret (3.4) as expressing the ar-29 bitrary vector  $K_{*\theta} \cup_{\theta} \in T_{K(\theta)} \mathcal{K}$  as a linear combination of  $\frac{\partial K}{\partial \theta^1}(\theta), \ldots, \frac{\partial K}{\partial \theta^{d+n}}(\theta)$ . 30 31

Therefore, the column vectors  $\left\{\frac{\partial K}{\partial \theta^{\alpha}}(\theta)\right\}_{\alpha=1}^{d+n}$  form a basis of  $T_{K(\theta)}\mathcal{K}$ . To adapt the basis  $\left\{\frac{\partial K}{\partial \theta^{\alpha}}(\theta)\right\}_{\alpha=1}^{d+n}$  to the special subspace ker  $\Omega_{K(\theta)} \subseteq T_{K(\theta)}\mathcal{K}$ , recall that the first d coordinates in  $\mathcal{P} = \mathbb{T}^d \times T^*\mathbb{T}^n$  (2.2) correspond to ker  $\Omega$ . 32 33 Because of this, we reorder the coordinates  $\theta^{\alpha}$ ,  $\alpha = 1, \ldots, d+n$ , of  $\mathbb{T}^{d+n}$  in such a 34 way that the  $d \times d$  block in the upper left corner of  $DK_{\theta}$  (3.5) has full rank (then 35 the  $2n \times n$  block in the lower right corner of  $DK_{\theta}$  will also be of full rank). Then 36

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Define the matrices  $Z_{\theta} \in M_{d+2n,d}(\mathbb{R})$  and  $X_{\theta} \in M_{d+2n,n}(\mathbb{R})$  as the first d, respectively the last n, columns of  $DK_{\theta}$  (3.5), so that  $DK_{\theta} = [Z_{\theta} \ X_{\theta}]$ . Recall the notations from Table 1 for the coordinates  $x = (x^{A}) = (\underline{x}, \widetilde{x})$ , in  $\mathcal{P}$  (2.2):

$$\underline{x} = (\underline{x}^{\mu}) = (x^1, \dots, x^d) \quad , \quad \widetilde{x} = (\widetilde{x}^i) = (\widetilde{x}^1, \dots, \widetilde{x}^{2n}) = (x^{d+1}, \dots, x^{d+2n})$$

6 We will use these notations in matrices with d + 2n rows – the underscore for the 7 first d rows, and the tilde for the remaining 2n rows:

$$Z_{\theta} = \begin{bmatrix} \underline{Z}_{\theta} \\ \widetilde{Z}_{\theta} \end{bmatrix}, \quad X_{\theta} = \begin{bmatrix} \underline{X}_{\theta} \\ \widetilde{X}_{\theta} \end{bmatrix}, \quad DK_{\theta} = \begin{bmatrix} Z_{\theta} & X_{\theta} \end{bmatrix} = \begin{bmatrix} \underline{Z}_{\theta} & \underline{X}_{\theta} \\ \widetilde{Z}_{\theta} & \widetilde{X}_{\theta} \end{bmatrix}.$$
(3.6)

3.2.2. Adapted basis of  $T_{K(\theta)}\mathcal{P}$ . Having constructed a basis of  $T_{K(\theta)}\mathcal{K}$ , we need n more vectors that span the complement of  $T_{K(\theta)}\mathcal{K}$  in  $T_{K(\theta)}\mathcal{P}$ . We will construct them in such a way that, together with the columns  $\frac{\partial K}{\partial \theta^{d+1}}(\theta), \ldots, \frac{\partial K}{\partial \theta^{d+n}}(\theta)$  of  $X_{\theta}$ , they form a symplectic basis of  $T_{K(\theta)}\mathcal{Q}$  (2.3). To this end we will use the matrix representation  $\tilde{J}_{K(\theta)}$  (2.5) of the symplectic form  $\tilde{\Omega}$  on  $\mathcal{Q}$ , as well as the Gramian matrix of the vectors  $\frac{\partial K}{\partial \theta^{d+1}}(\theta), \ldots, \frac{\partial K}{\partial \theta^{d+n}}(\theta)$ , i.e., the matrix  $X_{\theta}^{\top}X_{\theta} \in M_{n,n}(\mathbb{R})$  of their inner products with respect to the Euclidean inner product on  $T^*\mathbb{T}^n$ . Define

$$R_{\theta} := \left( \widetilde{X}_{\theta}^{\top} \widetilde{X}_{\theta} \right)^{-1} \in \mathcal{M}_{n,n}(\mathbb{R}) , \qquad (3.7)$$

$$\widetilde{Y}_{\theta} := \widetilde{J}_{K(\theta)}^{-1} \widetilde{X}_{\theta} R_{\theta} \in \mathcal{M}_{2n,n}(\mathbb{R}) , \qquad (3.8)$$

$$Y_{\theta} := \begin{bmatrix} 0\\ \widetilde{Y}_{\theta} \end{bmatrix} = \begin{bmatrix} 0\\ \widetilde{J}_{K(\theta)}^{-1} \widetilde{X}_{\theta} R_{\theta} \end{bmatrix} \in \mathcal{M}_{d+2n,n}(\mathbb{R}) .$$
(3.9)

8 Since J<sub>K(θ)</sub><sup>-1</sup>, X<sub>θ</sub>, and R<sub>θ</sub> are of maximal rank, Y<sub>θ</sub> and Y<sub>θ</sub> are of maximal rank:
9 rank Y<sub>θ</sub> = rank Y<sub>θ</sub> = n. We think of the n columns of Y<sub>θ</sub> as of as vectors in T<sub>K(θ)</sub>P.
10 Let Z<sub>θ,μ</sub>, X<sub>θ,a</sub>, Y<sub>θ,a</sub>, with μ = 1,...,d, a = 1,...,n, stand for the columns of
11 Z<sub>θ</sub>, X<sub>θ</sub>, and Y<sub>θ</sub>. These d + 2n vectors are a basis of T<sub>K(θ)</sub>P with the properties

$$\operatorname{span} \{ Z_{\theta,\mu} \}_{\mu=1}^{d} = \ker \Omega_{K(\theta)}$$
$$\operatorname{span} \left\{ \{ Z_{\theta,\mu} \}_{\mu=1}^{d}, \{ X_{\theta,a} \}_{a=1}^{n} \right\} = T_{K(\theta)} \mathcal{K} .$$

<sup>12</sup> The first property implies that  $\Omega_{K(\theta)}(Z_{\theta,\mu}, \cdot) = 0$ . The construction of  $Y_{\theta,a}$  yields <sup>13</sup> (using (2.4), (2.5), (2.6), (3.7), (3.8), and (3.9))

$$\Omega_{K(\theta)}(X_{\theta,a}, Y_{\theta,b}) = \langle X_{\theta,a}, J_{K(\theta)}Y_{\theta,b} \rangle_{\mathbb{R}^{d+2n}} = X_{\theta,a}^{\top}J_{K(\theta)}Y_{\theta,b} = \left(X_{\theta}^{\top}J_{K(\theta)}Y_{\theta}\right)_{ab}$$
$$= \left(\left[\underline{X}_{\theta}^{\top} \ \widetilde{X}_{\theta}^{\top}\right] \left[\begin{array}{c}0\\\widetilde{X}_{\theta}R_{\theta}\end{array}\right]\right)_{ab} = \left(\widetilde{X}_{\theta}^{\top} \ \widetilde{X}_{\theta}R_{\theta}\right)_{ab} = (\mathbb{I}_{n})_{ab} = \delta_{ab}$$

<sup>14</sup> ( $\mathbb{I}_n = (\delta_{ab})$  is the unit  $n \times n$  matrix), hence the vectors  $X_{\theta,a}$  and  $Y_{\theta,b}$  form a <sup>15</sup> symplectic basis of  $T_{K(\theta)}\mathcal{Q} \cong (T_{K(\theta)}\mathcal{P})/\ker \Omega_{K(\theta)}$ . We write this symbolically as

$$\Omega_{K(\theta)}(X_{\theta}, Y_{\theta}) = X_{\theta}^{\top} J_{K(\theta)} Y_{\theta} = \widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Y}_{\theta} = \mathbb{I}_{n} .$$
(3.10)

Below we summarize the properties of the basis of  $T_{K(\theta)}\mathcal{P}$  constructed above, using matrix notations as in (3.10):

$$Z_{\theta}^{\top} J_{K(\theta)} Z_{\theta} = \widetilde{Z}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Z}_{\theta} = 0 ,$$

$$Z_{\theta}^{\top} J_{K(\theta)} X_{\theta} = \widetilde{Z}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{X}_{\theta} = 0 ,$$

$$Z_{\theta}^{\top} J_{K(\theta)} Y_{\theta} = \widetilde{Z}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Y}_{\theta} = \widetilde{Z}_{\theta}^{\top} \widetilde{X}_{\theta} R_{\theta} ,$$

$$X_{\theta}^{\top} J_{K(\theta)} X_{\theta} = \widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{X}_{\theta} = 0 ,$$

$$X_{\theta}^{\top} J_{K(\theta)} Y_{\theta} = \widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Y}_{\theta} = \mathbb{I}_{n} ,$$

$$Y_{\theta}^{\top} J_{K(\theta)} Y_{\theta} = \widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Y}_{\theta} = -R_{\theta} \widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{X}_{\theta} R_{\theta} .$$
(3.11)

- 3 3.2.3. Presymplecticity of V at K in adapted coordinates. If  $U, V, W \in \mathfrak{X}(\mathcal{P})$  with
- <sup>4</sup> V presymplectic, i.e.,  $\mathcal{L}_V \Omega = 0$  (recall Lemma 2.3), then

$$0 = (\mathcal{L}_V \Omega)(U, W) = \mathcal{L}_V (\Omega(U, W)) - \Omega (\mathcal{L}_V U, W) - \Omega (U, \mathcal{L}_V W)$$
  
$$= \sum_{A,B,C=1}^{d+2n} U^A \left( \frac{\partial \Omega_{AB}}{\partial x^C} V^C + \frac{\partial V^C}{\partial x^A} \Omega_{CB} + \Omega_{AC} \frac{\partial V^C}{\partial x^B} \right) W^B$$
  
$$= \sum_{A,B=1}^{d+2n} U_A \left( (DJ) V + (DV)^\top J + J DV \right)^A_B W^B ,$$

- <sup>5</sup> where we used the operator  $J = (J^A{}_B)$  (2.4), lowering an index of a vector signifies
- <sup>6</sup> transposition (i.e., contraction with the Euclidean metric tensor), and

$$(DV)^C{}_B = \frac{\partial V^C}{\partial x^B} , \qquad ((DJ)V)^A{}_B := \sum_{C=1}^{d+2n} \frac{\partial J^A{}_B}{\partial x^C} V^C$$

7 Therefore in matrix notations the presymplecticity condition reads

$$(DJ)V + (DV)^{\top}J + JDV = 0.$$
 (3.12)

8 Writing the derivative of V at  $K(\theta) \in \mathcal{K}$  as

$$DV_{K(\theta)} =: \begin{bmatrix} \frac{\partial V}{\partial \underline{x}} & \frac{\partial V}{\partial \overline{x}} \\ \frac{\partial \tilde{V}}{\partial \underline{x}} & \frac{\partial \tilde{V}}{\partial \overline{x}} \end{bmatrix}_{K(\theta)}, \qquad (3.13)$$

9 we can easily show that (3.12) is equivalent to the conditions

$$\left[\frac{\partial \widetilde{V}}{\partial \underline{x}}\right]_{K(\theta)} = 0 , \quad (D\widetilde{J})_{K(\theta)}V_{K(\theta)} + \left[\frac{\partial \widetilde{V}}{\partial \widetilde{x}}\right]_{K(\theta)}^{\top}\widetilde{J}_{K(\theta)} + \widetilde{J}_{K(\theta)}\left[\frac{\partial \widetilde{V}}{\partial \widetilde{x}}\right]_{K(\theta)} = 0 . \quad (3.14)$$

# 10 3.3. Change of basis matrix $M_{\theta}$ .

11 3.3.1. Definition of  $M_{\theta}$ . The adapted basis  $\{Z_{\theta,\mu}\}_{\mu=1}^{d}$ ,  $\{X_{\theta,a}\}_{a=1}^{n}$ ,  $\{Y_{\theta,a}\}_{a=1}^{n}$  of 12  $T_{K(\theta)}\mathcal{P}$  constructed in Section 3.2 has properties that are very useful for our anal-13 ysis. Given an arbitrary column vector  $U_{\theta}$ , considered as an element of  $T_{K(\theta)}\mathcal{P}$ , we 14 can find its components in the adapted basis as follows. Define the *change of basis* 15 *matrix*  $M_{\theta}$  of all vectors from the adapted basis, written as column vectors:

$$M_{\theta} := [DK_{\theta} \ Y_{\theta}] = [Z_{\theta} \ X_{\theta} \ Y_{\theta}] = \begin{bmatrix} \underline{Z}_{\theta} & \underline{X}_{\theta} & 0\\ \widetilde{Z}_{\theta} & \widetilde{X}_{\theta} & \widetilde{Y}_{\theta} \end{bmatrix} \in \mathcal{M}_{d+2n,d+2n}(\mathbb{R}) \ . \tag{3.15}$$

1 Then the vector  $U_{\theta}$  can be written as a superposition of the vectors from the adapted 2 basis as follows:

$$U_{\theta} = M_{\theta} \,\xi_{\theta} = \sum_{\mu=1}^{d} Z_{\theta,\mu} \,\xi_{\theta}^{\mu} + \sum_{a=1}^{n} X_{\theta,a} \,\xi_{\theta}^{d+a} + \sum_{a=1}^{n} Y_{\theta,a} \,\xi_{\theta}^{d+a+a} \,. \tag{3.16}$$

In the adapted basis, if we write the (d + 2n) components of the vector  $\xi_{\theta}$  as three blocks of length d, n, and n, as in (3.16), then the vectors from  $T_{K(\theta)}\mathcal{K}$  have the form  $\xi_{\theta} = [* * 0]^{\top}$ , where the stars represent numbers that are generally non-zero. 3.3.2. Computing  $(DV_{K(\theta)} - \partial_{\omega})M_{\theta}$ . In this section we will perform some computations related to the change of basis matrix  $M_{\theta}$  (3.15), which will be needed in Section 4. Differentiating the invariance condition (2.17), we obtain

$$DV_{K(\theta)} DK_{\theta} = \partial_{\omega} DK_{\theta}$$
 (3.17)

9 This, together with the definition (3.15) of  $M_{\theta}$ , gives us

$$(DV_{K(\theta)} - \partial_{\omega})M_{\theta} = (DV_{K(\theta)} - \partial_{\omega}) [DK_{\theta} Y_{\theta}] = [0 \ 0 \ (DV_{K(\theta)} - \partial_{\omega})Y_{\theta}] ,$$

<sup>10</sup> Our first goal is to find an explicit expression for  $(DV_{K(\theta)} - \partial_{\omega})Y_{\theta}$ . To this end we <sup>11</sup> have to compute

$$(DV_{K(\theta)} - \partial_{\omega})Y_{\theta} = \begin{bmatrix} \begin{bmatrix} \frac{\partial V}{\partial \widetilde{x}} \end{bmatrix}_{K(\theta)} \widetilde{Y}_{\theta} \\ \begin{bmatrix} \frac{\partial \widetilde{V}}{\partial \widetilde{x}} \end{bmatrix}_{K(\theta)} \widetilde{Y}_{\theta} - \partial_{\omega}\widetilde{Y}_{\theta} \end{bmatrix}$$

<sup>12</sup> By the Leibniz rule,

$$\partial_{\omega} \widetilde{Y}_{\theta} = \partial_{\omega} \big( \widetilde{J}_{K(\theta)}^{-1} \widetilde{X}_{\theta} R_{\theta} \big) = \partial_{\omega} \big( \widetilde{J}_{K(\theta)}^{-1} \big) \widetilde{X}_{\theta} R_{\theta} + \widetilde{J}_{K(\theta)}^{-1} \partial_{\omega} \big( \widetilde{X}_{\theta} \big) R_{\theta} + \widetilde{J}_{K(\theta)}^{-1} \widetilde{X}_{\theta} \partial_{\omega} R_{\theta} \; .$$

<sup>13</sup> The elementary identity  $0 = \partial_{\omega}(\mathbb{I}_{2n}) = \partial_{\omega}(\widetilde{J}_{K(\theta)}^{-1}\widetilde{J}_{K(\theta)})$ , the invariance (2.17), and <sup>14</sup> the presymplecticity condition (3.14) yield

$$\partial_{\omega} \left( \widetilde{J}_{K(\theta)}^{-1} \right) = -\widetilde{J}_{K(\theta)}^{-1} \partial_{\omega} \left( \widetilde{J}_{K(\theta)} \right) \widetilde{J}_{K(\theta)}^{-1} = -\widetilde{J}_{K(\theta)}^{-1} \left( D\widetilde{J} \right)_{K(\theta)} \left( \partial_{\omega} K_{\theta} \right) \widetilde{J}_{K(\theta)}^{-1} = - \left( \widetilde{J}^{-1} \left( D\widetilde{J} \right) V \widetilde{J}^{-1} \right)_{K(\theta)} = - \left( \widetilde{J} \left[ \frac{\partial \widetilde{V}}{\partial \widetilde{x}} \right] + \left[ \frac{\partial \widetilde{V}}{\partial \widetilde{x}} \right]^{\top} \widetilde{J} \right)_{K(\theta)}.$$

$$(3.18)$$

<sup>15</sup> The invariance (3.17) and the expressions (3.13) and (3.14) give us

$$\partial_{\omega} \widetilde{X}_{\theta} = \left[ \frac{\partial \widetilde{V}}{\partial \widetilde{x}} \right]_{K(\theta)} \widetilde{X}_{\theta} .$$
(3.19)

From the definition (3.7) of  $R_{\theta}$ , the expression (3.19) for  $\partial_{\omega} X_{\theta}$ , we easily obtain

$$\partial_{\omega}R_{\theta} = -2R_{\theta}\widetilde{X}_{\theta}^{\top} \left[\frac{\partial\widetilde{V}}{\partial\widetilde{x}}\right]_{K(\theta)}^{\text{sym}} \widetilde{X}_{\theta}R_{\theta} . \qquad (3.20)$$

17 where  $\left[\frac{\partial \tilde{V}}{\partial \tilde{x}}\right]_{K(\theta)}^{\text{sym}} := \frac{1}{2} \left( \left[\frac{\partial \tilde{V}}{\partial \tilde{x}}\right]_{K(\theta)} + \left[\frac{\partial \tilde{V}}{\partial \tilde{x}}\right]_{K(\theta)}^{\mathsf{T}} \right).$ 18 It will be convenient to introduce the operator

$$\widetilde{\Pi}_{\theta} := \mathbb{I}_{2n} - \widetilde{X}_{\theta} R_{\theta} \widetilde{X}_{\theta}^{\top} : \mathbb{R}^{2n} \to \mathbb{R}^{2n} .$$
(3.21)

<sup>19</sup> We collect some properties of  $\widetilde{\Pi}_{\theta}$  and  $\widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)} = \mathbb{I}_{2n} - \widetilde{Y}_{\theta}\widetilde{X}_{\theta}^{\top}\widetilde{J}_{K(\theta)}$  which <sup>20</sup> follow easily from (3.7), (3.8), and (3.11):

- $\Pi_{\theta}$  is symmetric;
- both  $\widetilde{\Pi}_{\theta}$  and  $\widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)}$  are idempotent;

- the *n* columns of the matrix  $\widetilde{X}_{\theta}$  are in the kernel of  $\widetilde{\Pi}_{\theta}$ ;
- the *n* columns of  $\widetilde{X}_{\theta}$  are eigenvectors of  $\widetilde{J}_{K(\theta)}^{-1} \widetilde{\Pi}_{\theta} \widetilde{J}_{K(\theta)}$  with eigenvalue 1,
- while the *n* columns of  $\widetilde{Y}_{\theta}$  are in the kernel of  $\widetilde{J}_{K(\theta)}^{-1} \widetilde{\Pi}_{\theta} \widetilde{J}_{K(\theta)}$ :

$$\widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)}\widetilde{X}_{\theta} = \widetilde{X}_{\theta} , \qquad \widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)}\widetilde{Y}_{\theta} = 0 , \qquad (3.22)$$

hence  $\widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)}$  and  $\left(\mathbb{I}_{2n} - \widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta}\widetilde{J}_{K(\theta)}\right)$  are projection operators corresponding to the splitting

$$T_{K(\theta)}\mathcal{Q} = \operatorname{span}\left\{\widetilde{X}_{\theta,a}\right\}_{a=1}^{n} \oplus \operatorname{span}\left\{\widetilde{Y}_{\theta,a}\right\}_{a=1}^{n}.$$

<sup>6</sup> Putting together (3.18), (3.19), and (3.20), and using the definition (3.21), we <sup>7</sup> obtain

$$\partial_{\omega} \widetilde{Y}_{\theta} = \left[\frac{\partial \widetilde{V}}{\partial \widetilde{x}}\right]_{K(\theta)} \widetilde{Y}_{\theta} + 2\widetilde{J}_{K(\theta)}^{-1} \widetilde{\Pi}_{\theta} \left[\frac{\partial \widetilde{V}}{\partial \widetilde{x}}\right]_{K(\theta)}^{\text{sym}} \widetilde{X}_{\theta} R_{\theta}$$
(3.23)

8 so, finally,

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$$(DV_{K(\theta)} - \partial_{\omega}) M_{\theta} = \begin{bmatrix} 0 & 0 & \left[\frac{\partial V}{\partial \tilde{x}}\right]_{K(\theta)} \tilde{Y}_{\theta} \\ 0 & 0 & -2\tilde{J}_{K(\theta)}^{-1} \tilde{\Pi}_{\theta} \left[\frac{\partial \tilde{V}}{\partial \tilde{x}}\right]_{K(\theta)}^{\text{sym}} \tilde{X}_{\theta} R_{\theta} \end{bmatrix} .$$
 (3.24)

9 3.3.3. Writing  $(DV_{K(\theta)} - \partial_{\omega})M_{\theta}$  as  $M_{\theta}C_{\theta}$ . Having computed  $(DV_{K(\theta)} - \partial_{\omega})M_{\theta}$ , 10 we will rewrite it in a form that plays a crucial role in Section 4:

$$(DV_{K(\theta)} - \partial_{\omega})M_{\theta} = M_{\theta}C_{\theta} := M_{\theta} \begin{bmatrix} 0 & 0 & T_{\theta} \\ 0 & 0 & S_{\theta} \\ 0 & 0 & U_{\theta} \end{bmatrix} .$$
(3.25)

Substitute (3.24) and (3.15) in (3.25) to obtain that  $T_{\theta}$ ,  $S_{\theta}$ , and  $U_{\theta}$  should satisfy

$$\underline{Z}_{\theta}T_{\theta} + \underline{X}_{\theta}S_{\theta} = \left[\frac{\partial V}{\partial \tilde{x}}\right]_{K(\theta)} \tilde{J}_{K(\theta)}^{-1} \tilde{X}_{\theta}R_{\theta} , \qquad (3.26)$$

$$\widetilde{Z}_{\theta}T_{\theta} + \widetilde{X}_{\theta}S_{\theta} + \widetilde{Y}_{\theta}U_{\theta} = -2\widetilde{J}_{K(\theta)}^{-1}\widetilde{\Pi}_{\theta} \left[\frac{\partial\widetilde{V}}{\partial\widetilde{x}}\right]_{K(\theta)}^{\text{sym}} \widetilde{X}_{\theta}R_{\theta} .$$
(3.27)

Multiplying (3.27) separately by  $\widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)}$  and  $\widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)}$  on the left and using (3.11) and the definition of  $R_{\theta}$  (3.7), we obtain

$$U_{\theta} = -2\widetilde{X}_{\theta}^{\top} \widetilde{\Pi}_{\theta} \left[ \frac{\partial \widetilde{V}}{\partial \widetilde{x}} \right]_{K(\theta)}^{\text{sym}} \widetilde{X}_{\theta} R_{\theta} , \quad (3.28)$$

$$-R_{\theta}\widetilde{X}_{\theta}^{\top}\widetilde{Z}_{\theta}T_{\theta} - S_{\theta} - R_{\theta}\widetilde{X}_{\theta}^{\top}\widetilde{J}_{K(\theta)}^{-1}\widetilde{X}_{\theta}R_{\theta}U_{\theta} = -2\widetilde{Y}_{\theta}^{\top}\widetilde{\Pi}_{\theta} \left[\frac{\partial\widetilde{V}}{\partial\widetilde{x}}\right]_{K(\theta)}^{\mathrm{sym}}\widetilde{X}_{\theta}R_{\theta} .$$
(3.29)

<sup>11</sup> Since  $\widetilde{X}_{\theta}^{\top} \widetilde{\Pi}_{\theta} = (\widetilde{\Pi}_{\theta} \widetilde{X}_{\theta})^{\top} = 0$ , (3.28) yields  $U_{\theta} = 0$ . From (3.26) and (3.29) we <sup>12</sup> obtain

$$T_{\theta} = \underline{\mathscr{Z}}_{\theta}^{-1} \left( \left[ \frac{\partial V}{\partial \tilde{x}} \right]_{K(\theta)} \tilde{Y}_{\theta} - 2\underline{X}_{\theta} \tilde{Y}_{\theta}^{\top} \widetilde{\Pi}_{\theta} \left[ \frac{\partial \tilde{V}}{\partial \tilde{x}} \right]_{K(\theta)}^{\operatorname{sym}} \tilde{X}_{\theta} R_{\theta} \right) ,$$

$$S_{\theta} = -\tilde{Y}_{\theta}^{\top} \tilde{J}_{K(\theta)} \tilde{Z}_{\theta} \underline{\mathscr{Z}}_{\theta}^{-1} \left[ \frac{\partial V}{\partial \tilde{x}} \right]_{K(\theta)} \tilde{Y}_{\theta} \qquad (3.30)$$

$$+ 2 \left( \mathbb{I}_{n} - \tilde{Y}_{\theta}^{\top} \tilde{J}_{K(\theta)} \tilde{Z}_{\theta} \underline{\mathscr{Z}}_{\theta}^{-1} \underline{X}_{\theta} \right) \tilde{Y}_{\theta}^{\top} \widetilde{\Pi}_{\theta} \left[ \frac{\partial \tilde{V}}{\partial \tilde{x}} \right]_{K(\theta)}^{\operatorname{sym}} \tilde{X}_{\theta} R_{\theta} ,$$

13 where we have set

$$\underline{\mathscr{Z}}_{\theta} := \underline{Z}_{\theta} - \underline{X}_{\theta} R_{\theta} \widetilde{X}_{\theta}^{\top} \widetilde{Z}_{\theta} = \underline{Z}_{\theta} + \underline{X}_{\theta} \widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Z}_{\theta} \in \mathcal{M}_{d,d}(\mathbb{R}) .$$
(3.31)

- The geometric meaning of  $\underline{\mathscr{Z}}_{\theta}$  is discussed in detail in [9, Section 3.5.4]. 1
- We summarize our findings in the following 2
- **Lemma 3.4.** Let  $(\mathcal{P}, \Omega)$  be an exact presymplectic manifold,  $V \in \mathfrak{X}(\mathcal{P})$  be a presym-3
- plectic vector field,  $K: \mathbb{T}^{d+n} \to \mathcal{P}$  be an invariant torus in the sense of Defini-
- tion 3.1, and the matrix  $M_{\theta}$  be defined by (3.15). Then the equality (3.25) holds,
- with  $U_{\theta} = 0$  and  $T_{\theta}$  and  $S_{\theta}$  given by (3.30).
- 3.3.4. Factorization of  $M_{\theta}$ . Later we will need the representation of  $M_{\theta}^{-1}$  that 7 follows from the lemma below.

**Lemma 3.5.** If the  $(d+2n) \times (d+2n)$  matrices  $Q_{\theta}$  and  $W_{\theta}$  are defined by

$$Q_{\theta} := \begin{bmatrix} \mathbb{I}_{d} & 0\\ 0 & \widetilde{X}_{\theta}^{\top}\\ 0 & \widetilde{Y}_{\theta}^{\top} \end{bmatrix} \begin{bmatrix} \mathbb{I}_{d} & 0\\ 0 & \widetilde{J}_{K(\theta)} \end{bmatrix} = \begin{bmatrix} \mathbb{I}_{d} & 0\\ 0 & \widetilde{X}_{\theta}^{\top} \widetilde{J}_{K(\theta)}\\ 0 & \widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \end{bmatrix} , \qquad (3.32)$$

$$W_{\theta} = \begin{bmatrix} \underline{Z}_{\theta} & \underline{X}_{\theta} & 0\\ 0 & 0 & \mathbb{I}_{n}\\ \widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Z}_{\theta} & -\mathbb{I}_{n} & \widetilde{Y}_{\theta}^{\top} \widetilde{J}_{K(\theta)} \widetilde{Y}_{\theta} \end{bmatrix}, \qquad (3.33)$$

then the following identity holds: 9

$$Q_{\theta}M_{\theta} = W_{\theta} \ . \tag{3.34}$$

This implies that  $M_{\theta}$  (3.15) is invertible if and only if  $W_{\theta}$  is invertible. 10

*Proof.* The columns of  $\widetilde{X}_{\theta}$  and  $\widetilde{Y}_{\theta}$  form a (symplectic) basis of  $\mathbb{R}^{2n}$ , which implies 11 that the rows of  $\widetilde{X}_{\theta}^{\top}$  and  $\widetilde{Y}_{\theta}^{\top}$  form a basis of  $\mathbb{R}^{2n}$ . Since  $\widetilde{J}_{K(\theta)}$  is an invertible matrix 12 (it corresponds to the symplectic form  $\Omega$  on Q, recall (2.3) and (2.5)), the rows of 13  $\widetilde{X}_{\theta}^{\top}\widetilde{J}_{K(\theta)}$  and  $\widetilde{Y}_{\theta}^{\top}\widetilde{J}_{K(\theta)}$  from a basis of  $\mathbb{R}^{2n}$ , so that the matrix  $Q_{\theta}$  given by (3.32) 14 is invertible. The identity (3.34) follows directly from (3.11). 15

### 4. Approximate solutions

In this section we will examine the case when K is merely an approximate solu-17 tion as defined below. We will build off of the results in Section 3 for true solutions 18 to derive similar results for approximate solutions. We start with the definition for 19 approximate solution. 20

**Definition 4.1.** Let  $(\mathcal{P}, \Omega)$  be an exact presymplectic manifold,  $\{V_{\lambda}\}$  be a (d+2n)-21 parameter family of presymplectic vector fields,  $\omega \in \mathcal{D}(\gamma, \sigma)$ , and  $K_0: \mathbb{T}^{d+n} \to \mathcal{P}$ 22 23 be an embedding. For a value  $\lambda_0$  of the parameter  $\lambda$ , define the error,

$$e_{0,\theta} := V_{\lambda_0, K_0(\theta)} - \partial_\omega K_{0,\theta} \in T_{K_0(\theta)} \mathcal{P} \cong \mathbb{R}^{d+2n} .$$

$$(4.1)$$

If some appropriately defined norm of  $e_0$  is sufficiently small, then we say that  $K_0$ 24 is an approximate solution. 25

- 4.1. Approximately isotropic tori. In Lemma 3.3 we showed that if  $K_{\theta}$  is a 26 true solution (i.e., if (2.17) is satisfied), then the invariant manifold  $\mathcal{K} = K(\mathbb{T}^{d+n})$ 27 is isotropic, i.e.,  $K^*\Omega = 0$ . The analogous result for this section will be that if  $K_{0,\theta}$
- 28
- is an approximate solution, then  $K_0$  is approximately isotropic, i.e.,  $K_0^*\Omega$  is small. 29

14

**Lemma 4.2.** Let  $(\mathcal{P}, \Omega)$  be an exact presymplectic manifold,  $\{V_{\lambda}\}$  be a (d + 2n)parameter family of analytic presymplectic vector fields, and  $K_0 \in \mathcal{W}_{\rho}$  (2.10) be an approximate solution with Diophantine frequency  $\omega \in \mathcal{D}(\gamma, \sigma)$ . Assume that  $V_{\lambda}$  extends holomorphically to some complex neighborhood  $\mathcal{B}_r$  (2.18) of the image of  $\mathbb{T}_{\rho}^{d+n}$  under  $K_0$ , for some r > 0. Let  $L_{0,\theta} : T_{\theta}\mathbb{T}^{d+n} \to \mathbb{T}_{\theta}\mathbb{T}^{d+n}$  be the matrix representation of the pull-back  $(K_0^*\Omega)_{\theta}$  as in (3.1) and (3.2):

$$L_{0,\theta} = DK_{0,\theta}^{\dagger} J_{K_0(\theta)} DK_{0,\theta} .$$
(4.2)

<sup>7</sup> Then there exists a constant C > 0 depending on  $d, n, \sigma, \rho, \|DK_0\|_{\rho}, \|V_{\lambda_0}\|_{C^1, \mathcal{B}_r},$ <sup>8</sup> and  $|J|_{C^1, \mathcal{B}_r}$ , such that for every  $\delta$  satisfying  $0 < \delta < \frac{\rho}{2}$ , the following bound holds: <sup>9</sup>

$$\|L_0\|_{\rho-2\delta} < C\gamma^{-1}\delta^{-(\sigma+1)} \|e_0\|_{\rho} \quad .$$
(4.3)

<sup>10</sup> Proof. For the directional derivative of  $L_{0,\theta}$  (4.2) we have, using (4.6), (4.1), and (3.12):

$$\begin{aligned} \partial_{\omega} L_{0,\theta} &= \partial_{\omega} \left( DK_{0,\theta}^{\top} J_{K_{0}(\theta)} DK_{0,\theta} + DK_{0,\theta}^{\top} \partial_{\omega} \left( J_{K_{0}(\theta)} \right) DK_{0,\theta} + DK_{0,\theta}^{\top} J_{K_{0}(\theta)} \partial_{\omega} \left( DK_{0,\theta} \right) \\ &= \left( DV_{\lambda_{0}, K_{0}(\theta)} DK_{0,\theta} - De_{0,\theta} \right)^{\top} J_{K_{0}(\theta)} DK_{0,\theta} \\ &+ DK_{0,\theta}^{\top} DJ_{K_{0}(\theta)} \left( V_{\lambda_{0}, K_{0}(\theta)} - e_{0,\theta} \right) DK_{0,\theta} \\ &+ DK_{0,\theta}^{\top} J_{K_{0}(\theta)} \left( DV_{\lambda_{0}, K_{0}(\theta)} DK_{0,\theta} - De_{0,\theta} \right) \\ &= DK_{0,\theta}^{\top} \left( DV_{\lambda_{0}, K_{0}(\theta)}^{\top} J_{K_{0}(\theta)} + DJ_{K_{0}(\theta)} V_{\lambda_{0}, K_{0}(\theta)} + J_{K_{0}(\theta)} DV_{\lambda_{0}, K_{0}(\theta)} \right) DK_{0,\theta} \\ &- \left( De_{0,\theta}^{\top} J_{K_{0}(\theta)} DK_{0,\theta} + DK_{0,\theta}^{\top} DJ_{K_{0}(\theta)} e_{0,\theta} DK_{0,\theta} + DK_{0,\theta}^{\top} J_{K_{0}(\theta)} De_{0,\theta} \right) \\ &= - \left( De_{0,\theta}^{\top} J_{K_{0}(\theta)} DK_{0,\theta} + DK_{0,\theta}^{\top} DJ_{K_{0}(\theta)} e_{0,\theta} DK_{0,\theta} + DK_{0,\theta}^{\top} J_{K_{0}(\theta)} De_{0,\theta} \right) . \end{aligned}$$

<sup>11</sup> From this and the Cauchy bound (2.12) we obtain

$$\|\partial_{\omega}L_0\|_{\rho-\delta} \le C_1 \|e_0\|_{\rho-\delta} + C_2 \|De_0\|_{\rho-\delta} \le C\delta^{-1} \|e_0\|_{\rho} .$$
(4.4)

Although  $K_0$  is only an approximate solution, the exactness of  $\Omega$  implies that avg  $(L_0) = 0$  (the proof of this repeats part of the proof of Lemma 3.3, with Kreplaced by  $K_0$ ). We apply (4.4) and the Rüssmann estimate (2.13) to obtain

$$\|L_0\|_{\rho-2\delta} \le C\gamma^{-1}\delta^{-\sigma} \|\partial_{\omega}L_0\|_{\rho-\delta} \le C\gamma^{-1}\delta^{-(\sigma+1)}\|e_0\|_{\rho} .$$

15

4.2. Linearized equation for the corrections. Given a family of presymplectic vector fields  $\{V_{\lambda}\}$ , the equation (2.17) can be difficult to solve for a value  $\bar{\lambda}$  of the parameter and an embedding  $K : \mathbb{T}^{d+n} \to \mathcal{P}$ . So instead of solving it directly for  $\lambda$  and K, we start with an approximate solution  $(\lambda_0, K_0)$  and construct an iterative process that produces better approximate solutions. As a result, we obtain a sequence  $(\lambda_j, K_j)$  (1.2) that converges to  $(\lambda_{\infty}, K_{\infty}) = (\bar{\lambda}, K)$ .

Let  $(\lambda_j, K_j)$  be an approximate pair. Define the error (cf. (4.1))

$$e_{j,\theta} := V_{\lambda_j, K_j(\theta)} - \partial_\omega K_{j,\theta} \in T_{K_j(\theta)} \mathcal{P} \cong \mathbb{R}^{d+2n} .$$

$$(4.5)$$

We will usually consider  $e_j$  as a mapping  $e_j : \mathbb{T}^{d+n} \to \mathbb{R}^{d+2n}$ , whose derivative, 24  $De_j : \mathbb{T}^{d+n} \to M_{d+2n,d+n}(\mathbb{R})$ , is given by

$$De_{j,\theta} = \left(DV_{\lambda_j, K_j(\theta)} - \partial_\omega\right) DK_{j,\theta} \in \mathcal{M}_{d+2n, d+n}(\mathbb{R}) .$$
(4.6)

Note that in (4.6),  $DV_{\lambda_j,K_j(\theta)}$  stands for the derivative of the vector field  $V_{\lambda_j}$  with respect to the spatial variables  $x \in \mathbb{R}^{d+2n}$ :

$$DV_{\lambda_j,K_j(\theta)} = \left[ \left( DV_{\lambda_j,K_j(\theta)} \right)^A {}_B \right] = \left[ \frac{\partial V_{\lambda_j,x}^A}{\partial x^B} \Big|_{x=K_j(\theta)} \right] \in \mathcal{M}_{d+2n,d+2n}(\mathbb{R})$$

- <sup>3</sup> Recall that the presymplecticity of  $V_{\lambda_j}$  imply that  $DV_{\lambda_j,K_j(\theta)}$  (3.13) satisfies (3.14).
- Let  $\varepsilon_j$  and  $\Delta_j$  be (j+1)st correction terms, i.e.,

$$\lambda_{j+1} := \lambda_j + \varepsilon_j , \qquad K_{j+1,\theta} := K_{j,\theta} + \Delta_{j,\theta} .$$
(4.7)

5 To derive an equation for the corrections  $\varepsilon_j$  and  $\Delta_{j,\theta}$ , we set  $\mathcal{F}[\lambda, K] := V_\lambda \circ K - \partial_\omega K$ , so that  $e_{j,\theta} = \mathcal{F}[\lambda_j, K_j](\theta)$ . Expanding  $\mathcal{F}[\lambda_{j+1}, K_{j+1}]$  about  $(\lambda_j, K_j)$ , we 7 obtain

$$e_{j+1,\theta} = e_{j,\theta} + DV_{\lambda_j, K_j(\theta)} \,\Delta_{j,\theta} - \partial_\omega \Delta_{j,\theta} + \left[\frac{\partial V_\lambda}{\partial \lambda}\right]_{\lambda_j, K_j(\theta)} \varepsilon_j + \mathcal{O}(|\Delta_j, \varepsilon_j|^2) \,. \tag{4.8}$$

<sup>8</sup> Therefore, if the corrections  $\varepsilon_j$  and  $\Delta_j$  satisfy the linear equation

$$\left(DV_{\lambda_j, K_j(\theta)} - \partial_{\omega}\right) \Delta_{j,\theta} = -e_{j,\theta} - \left[\frac{\partial V_{\lambda}}{\partial \lambda}\right]_{\lambda_j, K_j(\theta)} \varepsilon_j , \qquad (4.9)$$

the cancellations in the right-hand side of (4.8) guarantee quadratic convergence. 9 The system (4.9) of (d+2n) equations for the corrections  $\varepsilon_i$  and  $\Delta_i$  is a linear 10 algebraic equation with respect to  $\varepsilon_j$ , and a linear first-order partial differential 11 equation with respect to  $\Delta_j$ . Since  $DV_{\lambda_j, K_j(\theta)} \in M_{d+2n, d+2n}(\mathbb{R})$  is of a general 12 form, is not easy to solve (4.9) and to obtain estimates on its solution. In Section 4.3 13 we will use the matrix  $M_{\theta}$  (3.15) of change of basis from a general one to the 14 adapted basis constructed in Section 3.2; the calculations from Section 3.3 will be 15 very useful. 16

# 17 4.3. Solving the linearized equation.

18 4.3.1. Change of basis. We use an adapted basis in  $\mathbb{R}^{d+2n}$ , so that instead of the 19 unknown function  $\Delta_j$  we introduce the unknown function  $\xi_j : \mathbb{T}^{d+n} \to \mathbb{R}^{d+2n}$ 20 through the linear change of basis

$$\Delta_{j,\theta} =: M_{j,\theta} \,\xi_{j,\theta} \,. \tag{4.10}$$

<sup>21</sup> The change of basis matrix  $M_{j,\theta} \in M_{d+2n,d+2n}(\mathbb{R})$  is constructed similarly to the <sup>22</sup> matrix  $M_{\theta}$  in (3.15), but by using the approximate value  $\lambda_j$  and the approximate <sup>23</sup> embedding  $K_j$ . Namely, given an approximate invariant torus  $K_j$  (which we treat <sup>24</sup> as a map  $K_j: \mathbb{T}^{d+n} \to \mathbb{R}^{d+2n}$ ), we define

$$\begin{bmatrix} \underline{Z}_{j,\theta} & \underline{X}_{j,\theta} \\ \widetilde{Z}_{j,\theta} & \widetilde{X}_{j,\theta} \end{bmatrix} := [Z_{j,\theta} \ X_{j,\theta}] := DK_{j,\theta} \in \mathcal{M}_{d+2n,d+n}(\mathbb{R})$$

as in (3.6). If the matrix  $\widetilde{X}_{j,\theta}^{\top}\widetilde{X}_{j,\theta}$  is invertible (cf. Definition 4.3 below), define  $R_{j,\theta} := \left(\widetilde{X}_{j,\theta}^{\top}\widetilde{X}_{j,\theta}\right)^{-1} \in \mathcal{M}_{n,n}(\mathbb{R})$  as in (3.7),

$$\widetilde{Y}_{j,\theta} := \widetilde{J}_{K_j(\theta)}^{-1} \widetilde{X}_{j,\theta} R_{j,\theta} \in \mathcal{M}_{2n,n}(\mathbb{R}) , \qquad Y_{j,\theta} := \begin{bmatrix} 0\\ \widetilde{Y}_{j,\theta} \end{bmatrix} \in \mathcal{M}_{d+2n,n}(\mathbb{R})$$

as in (3.8) and (3.9), and the approximate change of basis matrix  $M_{j,\theta}$  as in (3.15):

$$M_{j,\theta} := [DK_{j,\theta} \ Y_{j,\theta}] = [Z_{j,\theta} \ X_{j,\theta} \ Y_{j,\theta}] = \begin{bmatrix} \underline{Z}_{j,\theta} & \underline{X}_{j,\theta} & 0\\ \widetilde{Z}_{j,\theta} & \widetilde{X}_{j,\theta} & \widetilde{Y}_{j,\theta} \end{bmatrix} \in \mathcal{M}_{d+2n,d+2n}(\mathbb{R})$$

Below we will need to invert the matrices  $M_{j,\theta}$  in order to solve the linearized equation, which motivates the following

**Definition 4.3.** The map  $K \in W_{\rho}$  is said to be a non-degenerate torus if the matrices  $\widetilde{X}_{\theta}^{\top}\widetilde{X}_{\theta} \in M_{n,n}(\mathbb{R})$  and  $M_{\theta} \in M_{d+2n,d+2n}(\mathbb{R})$  defined above (and therefore, the matrix  $W_{\theta}$  as in (3.33)) are invertible.

<sup>6</sup> We will always assume that  $K_j$  is a non-degenerate torus (this is a part of <sup>7</sup> Definition 4.7 below which is one of the assumptions in the Main Theorem).

As before, we think of the columns  $Z_{j,\theta,\mu}$ ,  $X_{j,\theta,a}$ , and  $Y_{j,\theta,a}$   $(\mu = 1, \ldots, d, a = 1, \ldots, n)$  of the matrices  $Z_{j,\theta}$ ,  $X_{j,\theta}$ , and  $Y_{j,\theta}$  as vectors in  $T_{K_j(\theta)}\mathcal{P}$ . If  $K_j$  is close to the true solution K, then these vectors still form a basis of  $T_{K_j(\theta)}\mathcal{P}$  as in the true case. By construction, the columns of  $Z_{j,\theta}$  and  $X_{j,\theta}$  span the tangent space to the approximately invariant torus  $\mathcal{K}_j := K_j(\mathbb{T}^{d+n})$ :

span 
$$\left\{ \left\{ Z_{j,\theta,\mu} \right\}_{\mu=1}^{d}, \left\{ X_{j,\theta,a} \right\}_{a=1}^{n} \right\} = T_{K_{j}(\theta)} \mathcal{K}_{j}$$
 (4.11)

<sup>13</sup> However, unlike the case of a true solution,  $\mathcal{K}_j$  is *not* invariant with respect to the <sup>14</sup> flow of  $V_{\lambda_j}$ , and ker  $\Omega_{K_j(\theta)}$  is generally not a subspace of  $T_{K_j(\theta)}\mathcal{K}_j$ .

Making the substitution (4.10) in (4.9) and assuming that  $M_{j,\theta}$  is invertible, we obtain the following equation for the new unknown function  $\xi_{j,\theta}$ :

$$M_{j,\theta}^{-1} \left( DV_{\lambda_j,K_j(\theta)} M_{j,\theta} - \partial_\omega M_{j,\theta} \right) \xi_{j,\theta} - \partial_\omega \xi_{j,\theta} = -M_{j,\theta}^{-1} \left( e_{j,\theta} + \left[ \frac{\partial V_\lambda}{\partial \lambda} \right]_{\lambda_j,K_j(\theta)} \varepsilon_j \right).$$

$$\tag{4.12}$$

<sup>17</sup> To rewrite the coefficient of  $\xi_{j,\theta}$  in (4.12) in a simple form, we want that

$$(DV_{\lambda_{j},K_{j}(\theta)} - \partial_{\omega})M_{j,\theta} = M_{j,\theta} (C_{j,\theta} + B_{j,\theta}) := M_{j,\theta} \left( \begin{bmatrix} 0 & 0 & T_{j,\theta} \\ 0 & 0 & S_{j,\theta} \\ 0 & 0 & 0 \end{bmatrix} + B_{j,\theta} \right)$$
(4.13)

where  $B_{j,\theta}$  is "small," i.e., vanishing if  $e_j$  becomes identically zero (cf. (3.25)). For the left-hand side, long and unenlightening calculations (cf. Section 3.3.2) yield

$$\begin{aligned} \left( DV_{\lambda_{j},K_{j}(\theta)} - \partial_{\omega} \right) M_{j,\theta} &= \begin{bmatrix} De_{j,\theta} & \left( DV_{\lambda_{j},K_{j}(\theta)} - \partial_{\omega} \right) Y_{j,\theta} \end{bmatrix} \\ &= \begin{bmatrix} \left( DV_{\lambda_{j},K_{j}(\theta)} - \partial_{\omega} \right) DK_{j,\theta} & \left( DV_{\lambda_{j},K_{j}(\theta)} - \partial_{\omega} \right) \begin{bmatrix} 0 \\ \widetilde{Y}_{j,\theta} \end{bmatrix} \end{bmatrix} \\ &= \begin{bmatrix} 0 & 0 & \left[ \frac{\partial V_{\lambda_{j}}}{\partial \widetilde{x}} \right]_{K_{j}(\theta)} \left[ \frac{\partial V_{\lambda_{j}}}{\partial \widetilde{x}} \right]_{K_{j}(\theta)} \widetilde{Y}_{j,\theta} \\ \\ &0 & 0 & -2\widetilde{J}_{K_{j}(\theta)}^{-1} \widetilde{\Pi}_{j,\theta} \left[ \frac{\partial \widetilde{V}_{\lambda_{j}}}{\partial \widetilde{x}} \right]_{K_{j}(\theta)}^{\operatorname{sym}} \widetilde{X}_{j,\theta} R_{j,\theta} \end{bmatrix} + \begin{bmatrix} De_{j,\theta} & \left[ \begin{array}{c} 0 \\ \mathscr{E}[e_{j}](\theta) \end{array} \right] \end{bmatrix} \,. \end{aligned}$$

<sup>20</sup> Here we have set  $\widetilde{\Pi}_{j,\theta} := \mathbb{I}_{2n} - \widetilde{X}_{j,\theta} R_{j,\theta} \widetilde{X}_{j,\theta}^{\top} : \mathbb{R}^{2n} \to \mathbb{R}^{2n}$  (cf. (3.21)),

$$\mathscr{E}[e_{j}](\theta) := \widetilde{J}_{K_{j}(\theta)}^{-1} D \widetilde{J}_{K_{j}(\theta)} e_{j,\theta} \widetilde{Y}_{j,\theta} - \widetilde{J}_{K_{j}(\theta)}^{-1} \left[\frac{\partial \widetilde{e}_{j,\theta}}{\partial \theta}\right] R_{j,\theta} + \widetilde{Y}_{j,\theta} \left( \left[\frac{\partial \widetilde{e}_{j,\theta}}{\partial \theta}\right]^{\top} \widetilde{X}_{j,\theta} + \widetilde{X}_{j,\theta}^{\top} \left[\frac{\partial \widetilde{e}_{j,\theta}}{\partial \theta}\right] \right) R_{j,\theta} \in \mathcal{M}_{2n,n}(\mathbb{R}) ,$$

$$(4.14)$$

1 where  $\left[\frac{\partial \tilde{e}_{j,\theta}}{\partial \theta}\right]$  is the  $2n \times n$  matrix with entries  $\left[\frac{\partial \tilde{e}_{j,\theta}}{\partial \theta}\right]^{i}_{a} := \frac{\partial e^{d+i}_{j,\theta}}{\partial \theta^{d+a}}, i = 1, \dots, 2n,$ 2  $a = 1, \dots, n$ , and  $D \widetilde{J}_{K_{j}(\theta)} e_{j,\theta}$  stands for the  $2n \times 2n$  matrix with entries

$$\left(D\widetilde{J}_{K_{j}(\theta)} e_{j,\theta}\right)^{i}{}_{k} := \left.\sum_{A=1}^{d+2n} \frac{\partial \widetilde{J}^{i}{}_{k}}{\partial x^{A}}\right|_{K_{j}(\theta)} e^{A}_{j,\theta} , \qquad i,k=1,\ldots,2n$$

For the entries of the matrix  $C_{j,\theta}$  (4.13), we obtain (similarly to Section 3.3.3)

$$T_{j,\theta} = \underline{\mathscr{Z}}_{j,\theta}^{-1} \left( \left[ \frac{\partial V_{\lambda_j}}{\partial \tilde{x}} \right]_{K_j(\theta)} \widetilde{Y}_{j,\theta} - 2\underline{X}_{j,\theta} \widetilde{Y}_{j,\theta}^{\top} \widetilde{\Pi}_{j,\theta} \left[ \frac{\partial \tilde{V}_{\lambda_j}}{\partial \tilde{x}} \right]_{K_j(\theta)}^{\mathrm{sym}} \widetilde{X}_{j,\theta} R_{j,\theta} \right) ,$$

$$S_{j,\theta} = -\widetilde{Y}_{j,\theta}^{\top} \widetilde{J}_{K_j(\theta)} \widetilde{Z}_{j,\theta} \underline{\mathscr{Z}}_{j,\theta}^{-1} \left[ \frac{\partial V_{\lambda_j}}{\partial \tilde{x}} \right]_{K_j(\theta)} \widetilde{Y}_{j,\theta} \qquad (4.15)$$

$$+ 2 \left( \mathbb{I}_n - \widetilde{Y}_{j,\theta}^{\top} \widetilde{J}_{K_j(\theta)} \widetilde{Z}_{j,\theta} \underline{\mathscr{Z}}_{j,\theta}^{-1} \underline{X}_{j,\theta} \right) \widetilde{Y}_{j,\theta}^{\top} \widetilde{\Pi}_{j,\theta} \left[ \frac{\partial \widetilde{V}_{\lambda_j}}{\partial \tilde{x}} \right]_{K_j(\theta)} \widetilde{X}_{j,\theta} R_{j,\theta} ,$$

4 where we have set  $\underline{\mathscr{Z}}_{j,\theta} := \underline{Z}_{j,\theta} + \underline{X}_{j,\theta} \widetilde{Y}_{j,\theta}^{\top} \widetilde{J}_{K_j(\theta)} \widetilde{Z}_{j,\theta} \in \mathrm{M}_{d,d}(\mathbb{R})$  (cf. (3.30), (3.31)). 5 Summarizing, with the help of (4.13), (4.14), and (4.15), we rewrote (4.12) as

$$(C_{j,\theta} + B_{j,\theta})\xi_{j,\theta} - \partial_{\omega}\xi_{j,\theta} = -M_{j,\theta}^{-1} \left( e_{j,\theta} + \left[ \frac{\partial V_{\lambda}}{\partial \lambda} \right]_{\lambda_j, K_j(\theta)} \varepsilon_j \right), \qquad (4.16)$$

<sup>6</sup> where  $C_{j,\theta}$  is defined in (4.13) and (4.15), and  $B_{j,\theta}$  is a "small" matrix, given by

$$B_{j,\theta} = M_{j,\theta}^{-1} \left[ De_{j,\theta} \quad \left[ \begin{array}{c} 0\\ \mathscr{E}[e_j](\theta) \end{array} \right] \right] \,. \tag{4.17}$$

7 4.3.2. Invertibility issues. We define the  $(d+2n) \times (d+2n)$  matrices  $Q_{j,\theta}$  and  $W_{j,\theta}$ 8 in exactly the same way as  $Q_{\theta}$  (3.32) and  $W_{\theta}$  (3.33) but with  $X_{\theta}$ ,  $Y_{\theta}$ ,  $Z_{\theta}$ , and  $J_{K(\theta)}$ 9 replaced by  $X_{j,\theta}$ ,  $Y_{j,\theta}$ ,  $Z_{j,\theta}$ , and  $J_{K_{j}(\theta)}$ , respectively. Since the rank of the  $(2n \times 2n)$ 10 matrix  $[\tilde{X}_{j,\theta}, \tilde{Y}_{j,\theta}]$  is maximal and  $\tilde{J}_{K_{j}(\theta)}$  is non-degenerate,  $Q_{j,\theta}$  is non-degenerate. 11 By a direct calculation we obtain (cf. (3.34))

$$Q_{j,\theta}M_{j,\theta} = W_{j,\theta} + P_{j,\theta} := W_{j,\theta} + \begin{bmatrix} 0 & 0 & 0 \\ \widetilde{X}_{j,\theta}^{\top} \widetilde{J}_{K_j(\theta)} \widetilde{Z}_{j,\theta} & \widetilde{X}_{j,\theta}^{\top} \widetilde{J}_{K_j(\theta)} \widetilde{X}_{j,\theta} & 0 \\ 0 & 0 & 0 \end{bmatrix},$$
(4.18)

where the matrix  $P_{j,\theta}$  is small (if  $K_j$  were a true solution,  $P_{j,\theta}$  would be zero).

**Lemma 4.4.** Assume that the hypotheses of Lemma 4.2 hold. Then there exists a constant C depending on d, n,  $\sigma$ ,  $\rho$ ,  $\|DK_j\|_{\rho}$ ,  $|V_{\lambda_j}|_{C^1,\mathcal{B}_r}$ , and  $|J|_{C^1,\mathcal{B}_r}$ , such that

15 for every  $\delta$  satisfying  $0 < \delta < \frac{\rho}{2}$ , the following bound holds:

$$\|W_j^{-1} P_j\|_{\rho-2\delta} \le C\gamma^{-1}\delta^{-(\sigma+1)} \|e_j\|_{\rho} .$$
(4.19)

<sup>16</sup> Proof. Recalling the bound (4.3) on the norm of the pull-back  $L_j$  (4.2) of the <sup>17</sup> presymplectic form  $\Omega$  to the torus  $\mathcal{K}_j = K_j(\mathbb{T}^{d+n})$ , we obtain

$$\begin{split} \|W_j^{-1}P_j\|_{\rho-2\delta} &\leq C \,\|P_j\|_{\rho-2\delta} \\ &\leq C_1 \,\|\widetilde{X}_j^\top \left(\widetilde{J} \circ K_j\right)\widetilde{Z}_j\|_{\rho-2\delta} + C_2 \,\|\widetilde{X}_j^\top \left(\widetilde{J} \circ K_j\right)\widetilde{X}_j\|_{\rho-2\delta} \\ &\leq C \,\|DK_j^\top \left(\widetilde{J} \circ K_j\right)DK_j\|_{\rho-2\delta} = C \,\|L_j\|_{\rho-2\delta} \\ &\leq C\gamma^{-1}\delta^{-(\sigma+1)} \,\|e_j\|_{\rho} \,\,. \end{split}$$

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- The approximate factorization (4.18) can be used to write the inverse matrix  $M_{i,\theta}^{-1}$  in a convenient form, and Lemma 4.4 yields some useful bounds:
- 3 **Lemma 4.5.** Assume that the hypotheses of Lemma 4.2 hold. Let  $0 < \delta < \frac{\rho}{2}$  and 4 the error  $e_i$  satisfy the bound

$$C\gamma^{-1}\delta^{-(\sigma+1)} \|e_j\|_{\rho} \le \frac{1}{2}$$
, (4.20)

<sup>5</sup> where C is the same constant as in (4.19). Then the matrix  $M_{j,\theta}$  is invertible, and <sup>6</sup>

$$M_{j,\theta}^{-1} = W_{j,\theta}^{-1} Q_{j,\theta} + M_{j,\theta}^{\rm E} , \qquad (4.21)$$

<sup>7</sup> where the error term,  $M_{i,\theta}^{\rm E}$ , is given by

$$M_{j,\theta}^{\rm E} = -\left(\mathbb{I}_{d+2n} + W_{j,\theta}^{-1} P_{j,\theta}\right)^{-1} W_{j,\theta}^{-1} P_{j,\theta} W_{j,\theta}^{-1} Q_{j,\theta} , \qquad (4.22)$$

<sup>8</sup> and satisfies the bound

$$\|M_{j}^{\mathrm{E}}\|_{\rho-2\delta} \leq C' \gamma^{-1} \delta^{-(\sigma+1)} \|e_{j}\|_{\rho} ; \qquad (4.23)$$

9 here C' is a constant that depends on the same parameters as the constant C in (4.19).

<sup>11</sup> Proof. The expression (4.22) comes from (4.18), and (4.23) follows from (4.19).

<sup>12</sup> 4.3.3. Bounds on the "small" parts. Recall that, in order to find an approximate <sup>13</sup> solution of the linearized equation (4.9), we changed the variable  $\Delta_{j,\theta}$  to  $\xi_{j,\theta}$  by <sup>14</sup> (4.10) to transform it to the form (4.12). Then we rewrote the coefficient of  $\xi_{j,\theta}$  in <sup>15</sup> (4.12) as a sum of a "big" part,  $C_{j,\theta}$  (given by (4.13) and (4.15)), and a "small" <sup>16</sup> part,  $B_{j,\theta}$  (4.17). In the Lemma below we give bounds on the "small" terms <sup>17</sup> in (4.12).

**Lemma 4.6.** Let  $K_j \in W_{\rho}$  and the error  $e_j$  be defined by (4.5). Let the pair ( $\lambda_j, K_j$ ) be non-degenerate (in the sense of Definition 4.7 below) for the family V<sub>λ</sub> of presymplectic analytic vector fields. If the error  $e_j$  satisfies (4.20), then the change of variables (4.10) transforms the linearized equation (4.9) to

$$(C_{j,\theta} + B_{j,\theta})\xi_{j,\theta} - \partial_{\omega}\xi_{j,\theta} = -M_{j,\theta}^{-1} \left( e_{j,\theta} + \partial_{\lambda}V_{\lambda_{j},K_{j}(\theta)} \varepsilon_{j} \right)$$
$$= -W_{j,\theta}^{-1} Q_{j,\theta} e_{j,\theta} - W_{j,\theta}^{-1} Q_{j,\theta} \left[ \frac{\partial V_{\lambda}}{\partial \lambda} \right]_{\lambda_{j},K_{j}(\theta)} \varepsilon_{j} \quad (4.24)$$
$$-M_{j,\theta}^{\mathrm{E}} e_{j,\theta} - M_{j,\theta}^{\mathrm{E}} \left[ \frac{\partial V_{\lambda}}{\partial \lambda} \right]_{\lambda_{j},K_{j}(\theta)} \varepsilon_{j} ,$$

<sup>22</sup> where  $C_{j,\theta}$  is defined by (4.13) and (4.15),  $B_{j,\theta}$  by (4.17),  $M_j^E$  by (4.22). Further-<sup>23</sup> more,

$$\|B_{j}\|_{\rho-2\delta} \leq C\delta^{-1} \|e_{j}\|_{\rho} ,$$

$$\|M_{j}^{\mathrm{E}} e_{j}\|_{\rho-2\delta} \leq C\gamma^{-1}\delta^{-(\sigma+1)} \|e_{j}\|_{\rho}^{2} ,$$

$$\|M_{j}^{\mathrm{E}} \left[\frac{\partial V_{\lambda}}{\partial \lambda}\Big|_{\lambda_{j}} \circ K_{j}\right] \varepsilon_{j}\Big\|_{\rho-2\delta} \leq C\gamma^{-1}\delta^{-(\sigma+1)} \left\|\frac{\partial V_{\lambda}}{\partial \lambda}\Big|_{\lambda_{j}} \circ K_{j}\right\| \|e_{j}\|_{\rho} |\varepsilon_{j}| .$$

$$(4.25)$$

*Proof.* Equation (4.24) follows directly from (4.16) and (4.21), so we only need to derive the bounds (4.25). Combining (4.17) and (4.21), we obtain

$$B_{j,\theta} = \left( W_{j,\theta}^{-1} Q_{j,\theta} + M_{j,\theta}^{\mathrm{E}} \right) \left[ De_{j,\theta} \quad \left[ \begin{array}{c} 0 \\ \mathscr{E}[e_j](\theta) \end{array} \right] \right]$$

<sup>1</sup> From the definition (4.14) of  $\mathscr{E}[e_j](\theta)$  and the Cauchy bound (2.12),

$$\|\mathscr{E}[e_j]\|_{\rho-2\delta} \le C_1 \|e_j\|_{\rho-2\delta} + C_2 \delta^{-1} \|e_j\|_{\rho-\delta} \le C \delta^{-1} \|e_j\|_{\rho-\delta}$$

<sup>2</sup> which, together with the bound (4.23) on  $M_j^E$ , yields the first bound in (4.25):

$$|B_{j}|_{\rho-2\delta} \leq \left( \left\| W_{j}^{-1}Q_{j} \right\|_{\rho-2\delta} + \left\| M_{j}^{\mathrm{E}} \right\|_{\rho-2\delta} \right) \left( \left\| De_{j} \right\|_{\rho-2\delta} + \left\| \mathscr{E}[e_{j}] \right\|_{\rho-2\delta} \right)$$
$$\leq \left( C_{1} + C_{2}\gamma^{-1}\delta^{-(\sigma+1)} \|e_{j}\|_{\rho} \right) \gamma^{-1} \|e_{j}\|_{\rho-\delta} \leq C\gamma^{-1} \|e_{j}\|_{\rho} .$$

<sup>3</sup> The remaining two bounds in (4.25) in are direct consequences of (4.23).

(**4.20**).

4 4.3.4. Solving the simplified equation. To use the Newton method for finding  $\xi_{j,\theta}$ , 5 it is enough to solve (4.24) retaining only the "big" terms, i.e., ignoring all terms of 6 higher order with respect to  $||e_j||$ . As we will show below (see (4.37)), the term  $\varepsilon_j$ 7 is of order of  $||e_j||$ . Lemma 4.6 allows us to keep only the leading terms in (4.24):

order of 
$$||c_j||$$
. Termina 4.0 arows us to keep only the reading terms in (4.24).

$$C_{j,\theta}\xi_{j,\theta} - \partial_{\omega}\xi_{j,\theta} = -W_{j,\theta}^{-1} Q_{j,\theta} e_{j,\theta} - \Lambda_{j,\theta} \varepsilon_j , \qquad (4.26)$$

8 where we have set

$$\Lambda_{j,\theta} := W_{j,\theta}^{-1} Q_{j,\theta} \left[ \frac{\partial V_{\lambda}}{\partial \lambda} \right]_{\lambda_j, K_j(\theta)} .$$
(4.27)

9 Let us denote the first *d* components of  $\xi_{j,\theta}$  by  $\xi_{j,\theta}^{z}$ , the next *n* components by  $\xi_{j,\theta}^{x}$ , 10 and the last *n* components by  $\xi_{j,\theta}^{y}$ . Using the specific form of  $C_{j,\theta}$  (4.13), we rewrite 11 (4.26) in the form

$$\partial_{\omega} \begin{bmatrix} \xi_{j,\theta}^{z} \\ \xi_{j,\theta}^{x} \\ \xi_{j,\theta}^{y} \end{bmatrix} = W_{j,\theta}^{-1} Q_{j,\theta} e_{j,\theta} + \Lambda_{j,\theta} \varepsilon_{j} + \begin{bmatrix} T_{j,\theta} \xi_{j,\theta}^{y} \\ S_{j,\theta} \xi_{j,\theta}^{y} \\ 0 \end{bmatrix} .$$
(4.28)

Integrating both sides of (4.28) over  $\mathbb{T}^{d+n}$ , we obtain that (4.28) has a solution if and only if the average over  $\mathbb{T}^{d+n}$  of its right-hand side is 0:

$$\operatorname{avg}\left(W_{j}^{-1}Q_{j}e_{j}\right) + \operatorname{avg}\left(\Lambda_{j}\right)\varepsilon_{j} + \begin{bmatrix}\operatorname{avg}\left(T_{j}\xi_{j}^{\mathrm{v}}\right)\\\operatorname{avg}\left(S_{j}\xi_{j}^{\mathrm{v}}\right)\\0\end{bmatrix} = 0.$$
(4.29)

Observe that the right-hand side of the last n equations of the system (4.28) to does not involve  $\xi_j^{\rm y}$ , so that the last n equations of (4.28) have the form

$$\partial_{\omega}\xi_{j,\theta}^{y} = \left(W_{j,\theta}^{-1} Q_{j,\theta} e_{j,\theta} + \Lambda_{j,\theta} \varepsilon_{j}\right)^{y}, \qquad (4.30)$$

where the superscript y stands for the last n components. The system (4.30) for  $\xi_{j,\theta}^{y}$  has a solution if and only if the last-hand side of (4.30) has zero average. To

satisfy this condition, we will require initially that  $\varepsilon_j$  satisfy the linear system

$$\operatorname{avg}\left(W_{j}^{-1} Q_{j} e_{j}\right) + \operatorname{avg}\left(\Lambda_{j}\right) \varepsilon_{j} = 0 , \qquad (4.31)$$

which in turn has a solution if and only if the matrix  $\operatorname{avg}(\Lambda_j)$  is invertible. In order to guarantee that we can solve (4.31) for  $\varepsilon_j$ , we introduce the following

**Definition 4.7.** The pair  $(\lambda, K)$  is said to be non-degenerate for a (d + 2n)parameter family of vector fields  $V_{\lambda}$  if  $K \in \mathcal{W}_{\rho}$  is a non-degenerate torus (recall

<sup>23</sup> Definition 4.3) and the matrix  $\Lambda$  (as in (4.27)) has a non-singular average:

rank 
$$\operatorname{avg}(\Lambda) = d + 2n$$
,  $\operatorname{avg}(\Lambda) := \int_{\mathbb{T}^{d+n}} \Lambda_{\theta} \, \mathrm{d}\theta^{1} \cdots \mathrm{d}\theta^{d+n}$ . (4.32)

We assume that the pair  $(\lambda_j, K_j)$  is non-degenerate, and set  $\varepsilon_j$  to be equal to the preliminary value

$$\varepsilon_j^{\text{prelim}} := -\left\{ \operatorname{avg}\left(\Lambda_j\right) \right\}^{-1} \operatorname{avg}\left(W_j^{-1} Q_j e_j\right) , \qquad (4.33)$$

so that (4.31) is satisfied;  $\varepsilon_i^{\text{prelim}}$  satisfies the bound

$$|\varepsilon_j^{\text{prelim}}| \le C \operatorname{avg}(e_j) \le C ||e_j||_{\rho} \quad . \tag{4.34}$$

- <sup>4</sup> This choice of  $\varepsilon_j$  guarantees the existence of a solution  $\xi_j^{\rm y}$  of (4.30) which (thanks
- 5 to (4.34) and the Rüssmann's inequality (2.13) satisfies the bound

$$\|\xi_{j}^{\mathsf{y}}\|_{\rho-\delta} \leq C\gamma^{-1}\delta^{-\sigma} \|W_{j}^{-1}Q_{j}e_{j} + \Lambda_{j}\varepsilon_{j}^{\operatorname{prelim}}\|_{\rho} \leq C\gamma^{-1}\delta^{-\sigma} \|e_{j}\|_{\rho} \quad (4.35)$$

<sup>6</sup> Having found  $\xi_j^{\rm y}$  from solving (4.30), we redefine  $\varepsilon_j$  as

$$\varepsilon_{j} := -\left\{ \operatorname{avg}\left(\Lambda_{j}\right)\right\}^{-1} \left( \operatorname{avg}\left(W_{j}^{-1} Q_{j} e_{j}\right) + \left[ \begin{array}{c} \operatorname{avg}\left(T_{j} \xi_{j}^{y}\right) \\ \operatorname{avg}\left(S_{j} \xi_{j}^{y}\right) \\ 0 \end{array} \right] \right)$$
(4.36)

<sup>7</sup> to satisfy the solvability condition (4.29) for (4.28). Note that the change from <sup>8</sup>  $\varepsilon_j^{\text{prelim}}$  (4.33) to  $\varepsilon_j$  (4.36) does not affect the component  $\xi_j^{\text{y}}$ . Thanks to (4.35),  $\varepsilon_j$ <sup>9</sup> satisfies

$$|\varepsilon_j| \le C \left( \|e_j\|_{\rho-\delta} + \left\| \xi_j^y \right\|_{\rho-\delta} \right) \le C \gamma^{-1} \delta^{-\sigma} \|e_j\|_{\rho} \quad .$$

$$(4.37)$$

With the new value of  $\varepsilon_j$  from (4.36), we solve (4.28) to find  $\xi_j$  which, according to the Rüssmann's inequality (2.13) and the bound (4.37), satisfies

$$\begin{aligned} \|\xi_{j}\|_{\rho-2\delta} &\leq C\gamma^{-1}\delta^{-\sigma} \left( \left\| W_{j}^{-1} Q_{j} e_{j} + \Lambda_{j}\varepsilon_{j} \right\|_{\rho-\delta} + C \left\| \xi_{j}^{\mathrm{v}} \right\|_{\rho-\delta} \right) \\ &\leq C\gamma^{-1}\delta^{-\sigma} \left( \|e_{j}\|_{\rho} + |\varepsilon_{j}| + \gamma^{-1}\delta^{-\sigma} \|e_{j}\|_{\rho} \right) \leq C\gamma^{-2}\delta^{-2\sigma} \|e_{j}\|_{\rho} . \end{aligned}$$

$$(4.38)$$

<sup>12</sup> We summarize our findings in the following lemma:

16

**Lemma 4.8.** Assume the hypotheses of Lemma 4.6. Then there exist a parameter  $\varepsilon_j$  and a function  $\xi_j$  that solve the reduced linear equation (4.26) and satisfy the bounds (4.37) and (4.38).

#### 5. The Newton Method

In this section we will collect the estimates for the *j*th step of the iterative scheme and show that the Newton method generates a Cauchy sequence of approximate solutions in a Banach space which converges to a true solution. We only give brief sketches, referring the reader to [21, 3] for details.

**Lemma 5.1.** If the assumptions of Lemma 4.8 are satisfied and  $r_j := ||K_j - K_0||_{\rho_j} < r$ , then there exist a function  $\Delta_j$  and a parameter  $\varepsilon_j \in \mathbb{R}^{d+2n}$  such that

$$\begin{split} \|\Delta_{j}\|_{\rho_{j}-2\delta_{j}} &\leq c_{j}\gamma^{-2}\delta_{j}^{-2\sigma} \|e_{j}\|_{\rho_{j}} \\ \|D\Delta_{j}\|_{\rho_{j}-3\delta_{j}} &\leq c_{j}\gamma^{-2}\delta_{j}^{-(2\sigma+1)} \|e_{j}\|_{\rho_{j}} \\ &|\varepsilon_{j}| \leq c_{j} \left|\operatorname{avg}(\Lambda_{j})^{-1}\right| \|e_{j}\|_{\rho_{j}} , \end{split}$$

$$(5.1)$$

where  $c_j$  is a constant that depends on  $n, d, r, \rho, |V_{\lambda_j}|_{\mathcal{C}^2, \mathcal{B}_r}, ||DK_j||_{\rho_j}, ||R_j||_{\rho_j}$ , and  $\left\|\frac{\partial V_{\lambda}}{\partial \lambda}\right\|_{\rho_{i}}$  Additionally, if 2

$$r_j + c_j \gamma^{-2} \delta_j^{-(2\sigma-1)} \|e_j\|_{\rho_j} < r , \qquad (5.2)$$

then3

$$\|e_{j+1}\|_{\rho_{j+1}} \le c_j \gamma^{-4} \delta_j^{-4\sigma} \|e_j\|_{\rho_j}^2 \quad .$$
(5.3)

*Proof.* The inequalities (5.1) follow from Lemmata 2.7 and 4.6, and (4.10).

To see that  $K_{j+1} \in \mathcal{B}_r$ , that is,  $K_{j+1}$  stays within the neighborhood where V is holomorphically extended, we use (5.1) and (5.2): 6

$$\begin{aligned} \|K_{j+1} - K_0\|_{\rho_{j+1} - 2\delta_{j+1}} &= \|K_j + \Delta_j - K_0\|_{\rho_{j+1} - 2\delta_{j+1}} \\ &\leq \|K_j - K_0\|_{\rho_j} + \|\Delta_j\|_{\rho_j - 2\delta_j} \\ &\leq r_j + c_j \gamma^{-2} \delta_j^{-(2\sigma+1)} \|e_j\|_{\rho_j} < r \end{aligned}$$

To prove (5.3), recall from (4.8) that  $\xi_j = M_j^{-1} \Delta_j$  was found by solving (4.26), 7 8 SC

$$DV_{\lambda_{j},K_{j}(\theta)} \Delta_{j,\theta} - \partial_{\omega} \Delta_{j,\theta} + \left[\frac{\partial V_{\lambda}}{\partial \lambda}\right]_{\lambda_{j},K_{j}(\theta)} \varepsilon_{j} + e_{j}$$
$$= M_{j,\theta} \left( B_{j,\theta}\xi_{j,\theta} + M_{j,\theta}^{\mathrm{E}}e_{j,\theta} + M_{j,\theta}^{\mathrm{E}}\left[\frac{\partial V_{\lambda}}{\partial \lambda}\Big|_{\lambda_{j}} \circ K_{j}\right] \right),$$

and each term on the right hand side is quadratically small from Lemma 4.6, hence 9

$$\left\| DV_{\lambda_j,K_j} \,\Delta_j - \partial_\omega \Delta_j + \left[ \frac{\partial V_\lambda}{\partial \lambda} \right]_{\lambda_j,K_j} \,\varepsilon_j + e_j \right\|_{\rho_j - 2\delta_j} \le C \gamma^{-3} \delta^{-(3\sigma+1)} \left\| e_j \right\|_{\rho_j}^2$$

Finally, recalling the Taylor expansion (4.8) of  $e_{j+1,\theta}$ , we see that the remainder 10 term is on the order of  $\|\Delta_j\|_{\rho_j-2\delta_j}^2$ . Thus we get the estimate (5.3). 11

The lemma below guarantees that, if the error is small and some invertibility 12 conditions are met at the *j*th step, then the invertibility holds at the (j+1)st step. 13

**Lemma 5.2.** Assume the setup of Lemma 5.1. If  $c_j \gamma^{-2} \delta_j^{-(\sigma+1)} \|e_j\|_{a_i} \leq \frac{1}{2}$ , then: 14

- 15
- 16
- (i) if  $\widetilde{X}_j^{\top} \widetilde{X}_j$  is invertible, then  $\widetilde{X}_{j+1}^{\top} \widetilde{X}_{j+1}$  is invertible; (ii) if  $W_j$  is invertible, then  $W_{j+1}$  is invertible; (iii) if  $\operatorname{avg}(\Lambda_j)$  is invertible, then  $\operatorname{avg}(\Lambda_{j+1})$  is invertible. 17

*Proof.* Recalling that  $\widetilde{X}_j$  is a part of the matrix  $DK_j$ , we obtain  $\widetilde{X}_{j+1}^{\top}\widetilde{X}_{j+1} =$ 18  $\widetilde{X}_{j}^{\top}\widetilde{X}_{j} + P_{j}$ , with  $P_{j} := \widetilde{X}_{j}^{\top}\widetilde{\Delta}_{j,\widetilde{x}} + \widetilde{\Delta}_{j,\widetilde{x}}^{\top}\widetilde{X}_{j} + \widetilde{\Delta}_{j,\widetilde{x}}^{\top}\widetilde{\Delta}_{j,\widetilde{x}}$ , where  $\widetilde{\Delta}_{j,\widetilde{x}} \in \mathcal{M}_{2n,2n}(\mathbb{R})$  has 19 entries  $(\widetilde{\Delta}_{j,\widetilde{x}})^i{}_k = \partial(\widetilde{\Delta}_j)^i/\partial\widetilde{x}^k$ . The bounds (5.1) give an upper bound on the size 20 of  $P_j$ . The matrix  $\widetilde{X}_j^{\top} \widetilde{X}_j$  is invertible by assumption,  $\mathbb{I}_n + (\widetilde{X}_j^{\top} \widetilde{X}_j)^{-1} P_j$  is invertible 21 by the Neumann series, so  $\widetilde{X}_{j+1}^{\top}\widetilde{X}_{j+1} = \widetilde{X}_{j}^{\top}\widetilde{X}_{j}(\mathbb{I}_{n} + (\widetilde{X}_{j}^{\top}\widetilde{X}_{j})^{-1}P_{j})$  is invertible, which completes the proof of (i). The proofs of (ii) and (iii) are similar.  $\Box$ 22 23

The lemma below shows how close the initial approximation has to be for the 24 Newton method to be iterated indefinitely and to converge to a true solution  $K_{\infty}$ , 25 and gives a bound on the difference between  $K_{\infty}$  and the initial approximation  $K_0$ . 26 The proof can be found in [21, Lemma 13]. 27

1 **Lemma 5.3.** Let  $\{c_j\}_{j\geq 0}$  be the sequence of constants from Lemmata 5.1 and 5.2. 2 For  $0 < \delta_0 < \min(\rho_0/12, 1)$  define

$$\delta_j := \delta_0 2^{-j}$$
,  $\rho_j := \rho_{j-1} - 6\delta_{j-1}$ ,  $r_j := \|K_j - K_0\|_{\rho_j}$ ,

 $\begin{array}{l} \rho_{\infty} := \lim_{j \to \infty} \rho_j, \ K_{\infty} := \lim_{j \to \infty} K_j. \ Then \ there \ exists \ a \ constant \ C > 0 \ depending \ on \ d, \ n, \ |V_{\lambda}|_{\mathcal{C}^2, \mathcal{B}_r}, \ |J_0|_{\mathcal{C}^1, \mathcal{B}_r}, \ \|DK_0\|_{\rho_0}, \ and \ |\{\operatorname{avg}(\Lambda_0)\}^{-1}| \ such \ that \ if \ \|e_0\|_{\rho_0} \ satisfies \ the \ conditions \end{array}$ 

$$C2^{4\sigma} \gamma^{-4} \delta_0^{-4\sigma} \|e_0\|_{\rho_0} \le \frac{1}{2} ,$$
  
$$C \left( 1 + \frac{2^{4\sigma}}{2^{2\sigma} - 1} \right) \gamma^{-2} \delta_0^{-2\sigma} \|e_0\|_{\rho_0} < r ,$$

3 then the Newton method converges to a true solution  $(\lambda_{\infty}, K_{\infty})$ . Furthermore,

$$\|K_{\infty} - K_0\|_{\rho_0 - 6\delta_0} \le \frac{2^{2\sigma}}{2^{2\sigma} - 1} c\gamma^{-2} \delta_0^{-2\sigma} \|e_0\|_{\rho_0}.$$

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