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Quantitative rigidity of differential inclusions in two dimensions

Xavier Lamy* Andrew Lorent[†] Guanying Peng[‡]

Abstract

For any compact connected one-dimensional submanifold $K \subset \mathbb{R}^{2\times 2}$ without boundary which has no rank-one connection and is elliptic, we prove the quantitative rigidity estimate

$$\inf_{M \in K} \int_{B_{1/2}} |Du - M|^2 \, dx \le C \int_{B_1} \operatorname{dist}^2(Du, K) \, dx, \qquad \forall u \in H^1(B_1; \mathbb{R}^2).$$

This is an optimal generalization, for compact connected submanifolds of $\mathbb{R}^{2\times 2}$ without boundary, of the celebrated quantitative rigidity estimate of Friesecke, James and Müller for the approximate differential inclusion into SO(n). The proof relies on the special properties of elliptic subsets $K \subset \mathbb{R}^{2\times 2}$ with respect to conformal-anticonformal decomposition, which provide a quasilinear elliptic PDE satisfied by solutions of the exact differential inclusion $Du \in K$. We also give an example showing that no analogous result can hold true in $\mathbb{R}^{n\times n}$ for n > 3.

1 Introduction

In 1850, Liouville [23] proved that, given a domain $\Omega \subset \mathbb{R}^3$, any smooth map $u \colon \Omega \to \mathbb{R}^3$ satisfying the differential inclusion $Du(x) \in \mathbb{R}_+ O(n)$ for all $x \in \Omega$ must be either affine or a Möbius transform. A corollary to Liouvilles Theorem is that a C^3 function whose gradient belongs everywhere to SO(n) is an affine mapping. This phenomenon of being able to globally control a map satisfying a certain differential inclusion $Du \in K$ is known as "rigidity".

In [14] Friesecke, James and Müller solved a long standing open problem by proving an optimal quantitative rigidity estimate for K = SO(n). Specifically, they showed that for every bounded Lipschitz domain $\Omega \subset \mathbb{R}^n$, $n \geq 2$, there exists a constant $C(\Omega)$ such that, for K = SO(n),

$$\inf_{R \in K} \|Dv - R\|_{L^2(\Omega)} \le C(\Omega) \|\operatorname{dist}(Dv, K)\|_{L^2(\Omega)}, \qquad \forall v \in H^1(\Omega; \mathbb{R}^n).$$
 (1)

Here and below, $\operatorname{dist}(M,K)$ denotes the distance from a matrix $M \in \mathbb{R}^{n \times n}$ to a subset $K \subset \mathbb{R}^{n \times n}$ measured in the Euclidean norm. This result strengthened earlier work of a series

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of authors, including John [18], Rešetnjak [27], and Kohn [20], and it has had a number of important applications, in particular to thin film limits of elastic structures [14, 15]. One of the remarkable things about this result is that it is a striking fact about classical mathematical objects that could be understood by the mathematicians of hundreds of years ago.

A number of works have extended the above result (1) to cover various larger classes of matrices than K = SO(n). Chaudhuri and Müller [8] and later De Lellis and Székelyhidi [10] considered a set of the form $K = SO(n)A \cup SO(n)B$ where A and B are strongly incompatible in the sense of Matos [25]. Faraco and Zhong [13] proved an analogous quantitative rigidity result with $K = \mathfrak{m} \cdot SO(n)$ where $\mathfrak{m} \subset (0, +\infty)$ is compact. There the infimum in the left-hand side of (1) also needs to include the gradients of Möbius transforms, and the integral is over a smaller subset $\Omega' \subset\subset \Omega$. Similar results for maps defined on the sphere have been obtained recently by Luckhaus and Zemas [24]. Best constants for (1) are investigated by Lewicka and Müller in [22].

Our main result is an optimal generalization of the quantitative rigidity estimate of [14] in the context of compact connected submanifolds $K \subset \mathbb{R}^{2\times 2}$ without boundary.

Theorem 1.1. Let $K \subset \mathbb{R}^{2 \times 2}$ be a smooth, compact and connected 1-manifold without boundary. Assume that K has no rank-one connections, and is elliptic in the sense that there exists $C_* > 0$ such that

$$|M - M'|^2 \le C_* \det(M - M') \qquad \forall M, M' \in K.$$
(2)

Then for any $u \in H^1(B_1; \mathbb{R}^2)$ we have

$$\inf_{M \in K} \int_{B_{1/2}} |Du - M|^2 dx \le C \int_{B_1} \operatorname{dist}^2(Du, K) dx, \tag{3}$$

for some constant C = C(K) > 0.

Remark 1.2. A covering argument as in [13] shows that the estimate (3) in the balls $B_{1/2} \subset B_1$ automatically improves to

$$\inf_{M \in K} \int_{\Omega'} |Du - M|^2 dx \le C \int_{\Omega} \operatorname{dist}^2(Du, K) dx, \qquad \forall u \in H^1(\Omega; \mathbb{R}^2),$$

for any bounded domain $\Omega \subset \mathbb{R}^2$ and open subset $\Omega' \subset\subset \Omega \subset \mathbb{R}^2$, and a constant $C = C(K, \Omega', \Omega) > 0$.

This result is optimal among compact connected submanifolds $K \subset \mathbb{R}^{2\times 2}$ without boundary for the following reasons:

- First, it is classical that the no-rank-one-connections assumption is necessary for the rigidity of the exact differential inclusion (see e.g. [26, 19]).
- Second, ellipticity is necessary for the validity of the linearized version of (3) (see Remark 2.3), because non-ellipticity would imply that the tangent space $T_M K$ has a rank-one connection for some $M \in K$.
- Third, the two previous conditions (no rank-one connections and ellipticity) imply that the connected submanifold K must be of dimension 1 [32, Corollary 3.5 & 3.6].

Moreover, we provide in Section 5 an example showing that the two-dimensional setting is also optimal: there exists an elliptic 1-submanifold $K \subset \mathbb{R}^{3\times 3}$ without rank-one connection but which contains a so-called \mathcal{T}_4 configuration, a well-known obstruction to compactness of sequences $\{u_k\} \subset H^1$ satisfying $\operatorname{dist}(Du_k, K) \to 0$ in L^2 [7], and therefore to any type of quantitative rigidity estimate.

One of our motivations for studying differential inclusion into general submanifolds $K \subset \mathbb{R}^{2\times 2}$ is our previous work [21] where we obtained a rigidity result for a non-elliptic differential inclusion related to the so-called Aviles-Giga functional, and pointed out the nice consequences that a corresponding quantitative rigidity estimate would have. Theorem 1.1 is not valid for non-elliptic differential inclusions, but the ideas in the present work should be relevant to attain that goal.

While the statements of the quantitative rigidity results of [14, 8, 13] are elementary, their proofs are not. Their starting point, in addition to rigidity of the exact differential inclusion, is a linearized version of (1) for the differential inclusion $Du \in T_{M_0}K$ into a tangent space $T_{M_0}K$. For K = SO(n) and $M_0 = I$, this is Korn's inequality. A natural linearization procedure then provides a quantitative rigidity estimate, but in terms of the L^{∞} norm of $\operatorname{dist}(Du,K)$, rather than L^2 . Strengthening the L^2 bound on $\operatorname{dist}(Du,K)$ into an L^{∞} bound constitutes therefore the main difficulty. A key idea, introduced in [14], is to use the regularity of an elliptic PDE satisfied by solutions of the exact differential inclusion: for K = SO(n)the exact differential inclusion $Du \in SO(n)$ implies that the coordinate functions u_k are harmonic. For $K \subset \mathbb{R}_+SO(n)$ the coordinate functions satisfy the (n-2)-Laplace equation $\operatorname{div}(|\nabla u_k|^{n-2} \nabla u_k) = 0$. Such PDE follows from the universal identity $\operatorname{div} \operatorname{cof}(Du) = 0$ (where cof denotes the cofactor matrix), together with identities satisfied by matrices in the specific set K. It is satisfied by solutions of the exact differential inclusion, and for a general map u the error from solving that PDE can be controlled in terms of the right-hand side of (1). This allows to reduce the proof of (1) to maps solving that PDE. Elliptic regularity then provides, via a compactness argument, a uniform bound on dist(Du, K) and the linearization can be performed.

Following this scheme, the main ingredient to prove Theorem 1.1 is to embed K into the graph of a uniformly monotone vector field: this will be enough to turn the identity $\operatorname{div} \operatorname{cof}(Du) = 0$ into a quasilinear elliptic equation for the exact differential inclusion $Du \in K$.

Proposition 1.3. Let K be as in Theorem 1.1. There exist $G_1, G_2 : \mathbb{R}^2 \to \mathbb{R}^2$ smooth, globally Lipschitz vector fields such that

$$K \subset \left\{ \left(\begin{array}{c} A \\ iG_1(A) \end{array} \right) : A \in \mathbb{R}^2 \right\} \cap \left\{ \left(\begin{array}{c} -iG_2(B) \\ B \end{array} \right) : B \in \mathbb{R}^2 \right\},$$

and G_1, G_2 are uniformly monotone, that is

$$(G_j(X) - G_j(X')) \cdot (X - X') \ge \lambda |X - X'|^2 \qquad \forall X, X' \in \mathbb{R}^2,$$

for some constant $\lambda > 0$ depending only on K.

Proposition 1.3 relies on remarkable properties of elliptic subsets of $\mathbb{R}^{2\times 2}$ with respect to the decomposition into conformal and anticonformal parts, discovered in [32] and exploited in a striking manner in [5, 11, 12] (see also [19]). It is also related to the classical link between two-dimensional elliptic PDEs of second order and complex Beltrami equations, see e.g. the

introduction of [6], and [4, 2, 1, 3] for recent developments on nonlinear Beltrami equations. The proof of Proposition 1.3 is in fact extremely close to [3, Theorem 5]. Then the proof of Theorem 1.1 follows the scheme outlined above.

The article is organized as follows. In Section 2 we establish the two basic prerequisites to Theorem 1, rigidity for the exact differential inclusion and the linearized estimate. In Section 3 we give the proof of Proposition 1.3. In Section 4 we gather these ingredients to prove Theorem 1.1. In Section 5 we describe the counterexample in $\mathbb{R}^{3\times3}$.

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2 Basic ingredients: rigidity of the exact inclusion and linearized estimate

In this section, let $K \subset \mathbb{R}^{2\times 2}$ be as in Theorem 1.1. We prove the following two Lemmas.

Lemma 2.1. If $u \in H^1(B_1; \mathbb{R}^2)$ is such that $Du \in K$ a.e., then $Du \equiv M$ for some $M \in K$.

Lemma 2.2. For all $M \in K$ and $u \in H^1(B_1; \mathbb{R}^2)$ we have

$$\inf_{X \in T_M K} \int_{B_1} |Du - X|^2 \, dx \le C \int_{B_1} \operatorname{dist}^2(Du, T_M K) \, dx \tag{4}$$

for some constant C = C(K) > 0, where $T_M K$ denotes the linear tangent space to K at M.

Remark 2.3. The linearized estimate (4), or rather its weaker interior version

$$\inf_{X \in T_M K} \int_{B_{1/2}} |Du - X|^2 dx \le C \int_{B_1} \operatorname{dist}^2(Du, T_M K) dx, \tag{5}$$

is a necessary condition for (3) to be valid. Assume indeed that (3) is verified, fix $u \in C^1(\overline{B}_1; \mathbb{R}^2)$, and apply (3) to $v_{\epsilon}(x) = Mx + \epsilon u(x)$ for $\epsilon \ll 1$. There exists $M_{\epsilon} \in K$ such that

$$\int_{B_{1/2}} |M - M_{\epsilon} - \epsilon Du|^2 dx \le C \int_{B_1} \operatorname{dist}^2(M + \epsilon Du, K) dx$$
$$\le C \epsilon^2 \int_{B_1} \operatorname{dist}^2(Du, T_M K) dx + o(\epsilon^2).$$

Hence, letting $X_{\epsilon} = \epsilon^{-1}(M - M_{\epsilon})$, we have

$$\int_{B_{1/2}} |X_{\epsilon} - Du|^2 dx \le C \int_{B_1} \operatorname{dist}^2(Du, T_M K) dx + o(1).$$

In particular X_{ϵ} is bounded, and extracting a converging subsequence we obtain $X \in T_M K$ showing the validity of (5) for $u \in C^1(\overline{B}_1; \mathbb{R}^2)$, and then by density for $u \in H^1(\overline{B}_1; \mathbb{R}^2)$.

Proof of Lemma 2.1. Let $\ell = |K|$ and $\gamma \colon \mathbb{R}/\ell Z \to K$ be an arc-length parametrization of K. The ellipticity assumption (2) ensures that $\gamma'(t)$ is invertible for all $t \in \mathbb{R}$. Let $u \in H^1(B_1; \mathbb{R}^2)$ such that $Du \in K$ a.e., then u is smooth by [31], and since B_1 is simply connected there exists a smooth lifting $\theta \colon B_1 \to \mathbb{R}$ such that $Du = \gamma(\theta)$. Using that $\operatorname{div} \operatorname{cof}(Du) = 0$, where cof denotes the cofactor matrix, we find $\operatorname{cof}(\gamma'(\theta))\nabla\theta = 0$, hence $\nabla\theta = 0$ since $\operatorname{cof}(\gamma'(\theta))$ is invertible. Therefore Du is constant.

Proof of Lemma 2.2. For $M \in K$ we denote by $P_M : \mathbb{R}^{2 \times 2} \to \mathbb{R}^{2 \times 2}$ the orthogonal projection onto $(T_M K)^{\perp}$, so that

$$\operatorname{dist}^{2}(X, T_{M}K) = |P_{M}X|^{2}.$$

We denote by $P_M^{\alpha\beta jk} \in \mathbb{R}$ the coefficients of P_M , that is,

$$(P_M X)_{\alpha\beta} = \sum_{jk} P_M^{\alpha\beta jk} X_{jk} \qquad \forall X = (X_{jk}) \in \mathbb{R}^{2 \times 2}.$$

Define $\mathbb{P}_M(i\xi) \in \mathcal{L}(\mathbb{C}^2; \mathbb{C}^{2\times 2})$ by

$$(\mathbb{P}_M(i\xi)v)_{\alpha\beta} = \sum_{jk} P_M^{\alpha\beta jk} v_j \xi_k \qquad \forall \xi, v \in \mathbb{C}^2,$$

so the differential operator $u \mapsto P_M Du$ has symbol $\mathbb{P}_M(i\xi)$, i.e.,

$$(P_{M}(Df))_{\alpha\beta} = \frac{1}{2\pi} \int \left(\mathbb{P}_{M}(i\xi) \, \hat{f}(\xi) \right)_{\alpha\beta} e^{ix\cdot\xi} d\xi.$$

We claim that $\mathbb{P}_M(\xi)$ has trivial kernel for all non-zero $\xi \in \mathbb{C}^2$. Let indeed $v \in \mathbb{C}^2$ such that $\mathbb{P}_M(\xi)v = 0$. This implies that $P_M \operatorname{Re}(v \otimes i\xi) = P_M \operatorname{Im}(v \otimes i\xi) = 0$, because the coefficients $P_M^{\alpha\beta jk}$ are real-valued. In other words, the real and imaginary parts of $v \otimes i\xi$ both belong to $\ker P_M = T_M K$. Since $T_M K$ is a one-dimensional subspace of $\mathbb{R}^{2\times 2}$ which doesn't contain any rank-one matrix, we have $T_M K = \mathbb{R} X_0$ for some invertible matrix X_0 . Hence we deduce that $v \otimes i\xi = \lambda X_0$ for some $\lambda \in \mathbb{C}$ and an invertible matrix $X_0 \in T_M K$. But $v \otimes i\xi$ has zero determinant, so $\lambda = 0$ and we must have v = 0. This proves that $\mathbb{P}_M(\xi)$ has trivial kernel. Therefore we have the representation formula [29, Theorem 4.1] and the coercive inequality [29, Theorem 8.15] that follows from it,

$$\int_{B_1} |Du|^2 dx \le C \int_{B_1} |P_M Du|^2 dx + C \int_{B_1} |u|^2 dx, \tag{6}$$

for all $u \in H^1(B_1; \mathbb{R}^2)$. (In the notation of [29], N = 4, M = 2, and the index set $\{1, 2, 3, 4\}$ for j is in our case given by $\{1, 2\}^2$, and we can take $m_j = 1$ for $j \in \{1, 2\}^2$ and $l_i = 0$ for $i \in \{1, 2\}$.) The constant C > 0 in (6) depends a priori on the fixed matrix $M \in K$. Denote by C(M) the best possible constant in (6). Then for any $M, M' \in K$ we have

$$\int_{B_1} |Du|^2 dx \le 2C(M) \int_{B_1} |P_{M'}Du|^2 dx + C(M) \int_{B_1} |u|^2 dx + 2C(M) ||P_M - P_{M'}||^2 \int_{B_1} |Du|^2 dx.$$

For all $M \in K$, there exists $\delta(M) > 0$ sufficiently small such that for all $M' \in K \cap B_{\delta(M)}(M)$, we have $2C(M)\|P_M - P_{M'}\|^2 < 1/2$. It follows that

$$C(M') \le \frac{2C(M)}{1 - 2C(M)\|P_M - P_{M'}\|^2} < 4C(M) \quad \forall M' \in K \cap B_{\delta(M)}(M).$$

By compactness this implies that $M \mapsto C(M)$ is bounded on K, so we can take the constant C in (6) to depend only on K.

Moreover, if $u \in H^1(B_1; \mathbb{R}^2)$ satisfies $P_M Du = 0$ a.e., then $Du = \lambda X_0$ for some $\lambda \in L^2(B_1; \mathbb{R})$, and the distributional identity $0 = \operatorname{div} \operatorname{cof}(Du) = \operatorname{cof}(X_0) \nabla \lambda$ implies that λ is constant, hence $Du \equiv X$ for some $X \in T_M K$.

Therefore (4) follows from (6) via a compactness argument: assume by contradiction the existence of sequences $M^k \in K$, and $u^k \in H^1(B_1; \mathbb{R}^2)$ such that

$$\inf_{X \in T_{M^k} K} \int_{B_1} |Du^k - X|^2 \, dx = 1, \quad \int_{B_1} |P_{M^k} Du^k|^2 \, dx \longrightarrow 0.$$

Subtracting from u^k its average and $X^k x$ for the matrix X^k at which the infimum in the left-hand side is attained, we may in fact assume

$$\int_{B_1} u^k \, dx = \int_{B_1} Du^k \cdot X \, dx = 0 \qquad \forall X \in T_{M^k} K,$$

and

$$\int_{B_1} |Du^k|^2 \, dx = 1, \quad \int_{B_1} |P_{M^k} Du^k|^2 \, dx \longrightarrow 0.$$

Thus we may extract subsequences (not relabeled) $u^k \to u$ weakly in $H^1(B_1; \mathbb{R}^2)$ and strongly in $L^2(B_1; \mathbb{R}^2)$, and $M^k \to M \in K$. It follows that $P_{M^k}Du^k \rightharpoonup P_MDu$ in $L^2(B_1; \mathbb{R}^{2 \times 2})$, and thus by lower semicontinuity of the L^2 norm under weak convergence, we have $P_MDu = 0$ a.e., which implies $Du \equiv X$ for some $X \in T_MK$. Approximating X by a sequence $X^k \in T_{M^k}K$ and using $Du^k \rightharpoonup Du$ in $L^2(B_1; \mathbb{R}^{2 \times 2})$, we deduce that $0 = \int_{B_1} Du \cdot X \, dx = |B_1| |X|^2$, and thus $Du \equiv X = 0$. Further u satisfies $\int_{B_1} u \, dx = 0$, which implies $u \equiv 0$. Plugging u^k into (6) gives

$$1 = \int_{B_1} |Du^k|^2 dx \le C \int_{B_1} |P_{M^k} Du^k|^2 dx + C \int_{B_1} |u^k|^2 dx.$$

Passing to the limit as $k \to \infty$ and using the strong L^2 convergence, we have $1 \le C \int_{B_1} |u|^2 dx = 0$, which gives a contradiction.

3 Proof of Proposition 1.3

We only prove the existence of G_1 , the existence of G_2 is obtained by the same arguments. The proof relies on the properties of the conformal and anticonformal projections of K uncovered in [32, 12]. For any $z_+, z_- \in \mathbb{C}$, we denote by

$$[z_+, z_-] = \begin{pmatrix} \operatorname{Re} z_+ & -\operatorname{Im} z_+ \\ \operatorname{Im} z_+ & \operatorname{Re} z_+ \end{pmatrix} + \begin{pmatrix} \operatorname{Re} z_- & \operatorname{Im} z_- \\ \operatorname{Im} z_- & -\operatorname{Re} z_- \end{pmatrix} \in \mathbb{R}^{2 \times 2},$$

the 2×2 matrix whose conformal, respectively anticonformal, part is represented by z_+ , respectively z_- . For any $A \in \mathbb{R}^{2 \times 2}$, the decomposition $A = [z_+, z_-]$ is unique, and we have the identities

$$\det A = |z_{+}|^{2} - |z_{-}|^{2}, \quad |A|^{2} = 2|z_{+}|^{2} + 2|z_{-}|^{2}, \quad ||A|| = |z_{+}| + |z_{-}|,$$

where |A| and ||A|| denote the Hilbert-Schmidt and the operator norms of A, respectively. We denote by $p_+: [z_+, z_-] \mapsto z_+$ the projection onto the conformal part.

Using these notations, the ellipticity assumption (2) is equivalent to

$$2|z_{+}-z'_{+}|^{2}+2|z_{-}-z'_{-}|^{2} \le C_{*}(|z_{+}-z'_{+}|^{2}+|z_{-}-z'_{-}|^{2}),$$

for all $[z_+, z_-], [z'_+, z'_-] \in K$. This implies that the curve K is C_* -elliptic in the sense of [12, Def. 1], and it follows from the analysis in [32] (see also the proofs of Lemmas 1 and 2 in [12]) that $p_+(K)$ is a Jordan curve and

$$K = \{ [z, H(z)] : z \in p_+(K) \},$$

for some k-Lipschitz function $H: p_+(K) \to \mathbb{C}$ with $0 \le k = (C_* - 1)/(C_* + 1) < 1$. Further the explicit formula $p_+(A) = (a_{11} + a_{22})/2 - i(a_{12} - a_{21})/2$ for $A = (a_{ij})$ and the smoothness of K imply that $p_+(K) \subset \mathbb{C}$ is smooth.

Lemma 3.1. The function H admits a smooth extension $H: \mathbb{C} \to \mathbb{C}$ which is k-Lipschitz for some (possibly larger) $0 \le k < 1$.

Proof of Lemma 3.1. We first fix, thanks to Kirszbraun's theorem, a k-Lipschitz extension $\widehat{H}: \mathbb{C} \to \mathbb{C}$. In the rest of the proof we modify \widehat{H} to make it smooth while still agreeing with H on $p_+(K)$, at the cost of slightly increasing its Lipschitz constant.

Let $[z_+(t),z_-(t)]$, $t\in\mathbb{R}/\ell\mathbb{Z}$, denote a smooth arc-length parametrization of K, so that $2|\dot{z}_+|^2+2|\dot{z}_-|^2=1$, where $\dot{z}_\pm=\frac{d}{dt}z_\pm$. As $z_-=H(z_+)$ and H is k-Lipschitz, it follows that $|\dot{z}_-|\leq k|\dot{z}_+|$, so $|\dot{z}_+|^2\geq \frac{1}{2(1+k^2)}>0$. Therefore we may reparametrize and consider

$$K = \{ [z(s), H(z(s))] \colon s \in \mathbb{R}/\ell_+\mathbb{Z} \},\,$$

with z(s) an arc-length parametrization of $p_+(K)$ and ℓ_+ its length, and the map $s \mapsto H(z(s))$ is smooth by smoothness of K.

For small enough $\delta > 0$, the map

$$\varphi \colon \mathbb{R}/\ell_{+}\mathbb{Z} \times (-2\delta, 2\delta) \to \mathcal{U}_{2\delta} = \{ z \in \mathbb{C} \colon \operatorname{dist}(z, p_{+}(K)) < 2\delta \}$$
$$(s, r) \mapsto z(s) + ri\dot{z}(s),$$

is a smooth diffeomorphism. We first modify \widehat{H} by setting

$$\widetilde{H} = \widehat{H} \circ \Phi, \qquad \Phi(Z) = \begin{cases} z(s) + \lambda(r)i\dot{z}(s) & \text{if } Z = \varphi(s,r) \in \mathcal{U}_{2\delta}, \\ Z & \text{otherwise,} \end{cases}$$

where λ is the odd $(1 - \delta)^{-1}$ -Lipschitz function given for r > 0 by

$$\lambda(r) = \begin{cases} 0 & \text{for } 0 \le r \le \delta^2, \\ -\frac{\delta^2}{1-\delta} + \frac{r}{1-\delta} & \text{for } \delta^2 < r \le \delta, \\ r & \text{for } r > \delta. \end{cases}$$

In particular we have

$$\widetilde{H}(Z) = H(z(s)) \qquad \forall Z = \varphi(s, r) \in \mathcal{U}_{\delta^2},$$

so \widetilde{H} is smooth in \mathcal{U}_{δ^2} (by smoothness of $s\mapsto H(z(s))$ and φ^{-1}) and agrees with H on $p_+(K)$. Note that by definition of φ and λ we have $\Phi(Z)=Z$ in $\mathbb{C}\setminus U_\delta$ and therefore Φ is Lipschitz in \mathbb{C} . Since $D\varphi(s,0)\in SO(2)$ for all $s\in \mathbb{R}/\ell_+\mathbb{Z}$, we have $\|D\varphi\|\leq 1+C\delta$ on $\mathbb{R}/\ell_+\mathbb{Z}\times(-2\delta,2\delta)$. Further, we have $\|D(\varphi^{-1})\|=1$ on $p_+(K)$, and $\|D(\varphi^{-1})\|\leq 1+C\delta$ on $\mathcal{U}_{2\delta}$. Denoting by ψ the $(1-\delta)^{-1}$ -Lipschitz map $(s,r)\mapsto (s,\lambda(r))$, we write $\Phi(Z)=\varphi(\psi(\varphi^{-1}(Z)))$ and deduce that $\|D\Phi\|\leq 1+C\delta$ a.e. in $\mathcal{U}_{2\delta}$. This inequality is also true in the rest of \mathbb{C} by definition of Φ , so we conclude that \widetilde{H} is \widetilde{k} -Lipschitz in \mathbb{C} , with $\widetilde{k}=(1+C\delta)k<1$ for small enough $\delta>0$. Now δ is fixed and we define, for $\epsilon\in(0,\delta^2/4)$,

$$H_{\epsilon}(z) = \int_{\mathbb{C}} \widetilde{H}(z + \epsilon \chi(z)y) \rho(y) dy,$$

for a smooth kernel $\rho \geq 0$ with support in B_1 and $\int \rho(y) dy = 1$, and some smooth cut-off function χ with $\mathbf{1}_{\mathcal{U}_{\delta^2/4}} \leq 1 - \chi \leq \mathbf{1}_{\mathcal{U}_{\delta^2/2}}$. In $\mathcal{U}_{3\delta^2/4}$, the map H_{ϵ} is smooth thanks to the smoothness of \widetilde{H} in \mathcal{U}_{δ^2} . In $\mathbb{C} \setminus \overline{\mathcal{U}}_{\delta^2/2}$, we have $H_{\epsilon}(z) = \int_{\mathbb{C}} \widetilde{H}(z + \epsilon y) \rho(y) dy$ is also smooth. Therefore the map H_{ϵ} is smooth in \mathbb{C} . Further, $H_{\epsilon}(z) = \widetilde{H}(z)$ for $z \in \mathcal{U}_{\delta^2/4}$, and thus agrees with H on $p_+(K)$. Finally, denoting by L the Lipschitz constant of χ , we have

$$|H_{\epsilon}(z) - H_{\epsilon}(z')| \le \int_{B_1} |\widetilde{H}(z + \epsilon \chi(z)y) - \widetilde{H}(z' + \epsilon \chi(z')y)| \rho(y) \, dy$$

$$\le \int_{B_1} \widetilde{k}(1 + \epsilon L|y|) |z - z'| \rho(y) \, dy \le \widetilde{k}(1 + \epsilon L) |z - z'|.$$

So H_{ϵ} is k_{ϵ} -Lipschitz with $k_{\epsilon} \leq \hat{k}(1 + \epsilon L) < 1$ for small enough ϵ .

Lemma 3.2. The map $F: z \mapsto \bar{z} + H(z)$ is a smooth diffeomorphism from \mathbb{C} onto \mathbb{C} . Moreover F is (1+k)-Lipschitz and F^{-1} is $(1-k)^{-1}$ -Lipschitz.

Proof of Lemma 3.2. For any $w \in \mathbb{C}$ the equation

$$w = \bar{z} + H(z) \quad \Leftrightarrow \quad z = \bar{w} - \bar{H}(z),$$

admits a unique solution $z \in \mathbb{C}$ thanks to the fixed point theorem, since $z \mapsto \bar{w} - \bar{H}(z)$ is k-Lipschitz and $0 \le k < 1$. This shows that F is bijective. The inequalities

$$(1-k)|z-z'| \le |F(z) - F(z')| \le (1+k)|z-z'|,$$

follow directly from the fact that H is k-Lipschitz and imply the announced Lipschitz constants of F and F^{-1} . The inverse F^{-1} is smooth thanks to the Inverse Function Theorem, since $DF(z) \colon h \mapsto \bar{h} + DH(z)h$ is invertible for all $z \in \mathbb{C}$ because ||DH|| < 1.

Proof of Proposition 1.3 completed. Then, identifying \mathbb{C} with \mathbb{R}^2 , we define

$$G_1(A) = \overline{F^{-1}(A)} - H(F^{-1}(A)),$$

so that a short calculation shows that

$$[z, H(z)] = \begin{pmatrix} A \\ iG_1(A) \end{pmatrix}$$
 for $A = F(z)$.

The map G_1 is smooth and globally Lipschitz with Lipschitz constant $\Lambda = (1+k)/(1-k)$. Moreover, for all A = F(z), A' = F(z'), we have

$$(G_1(A) - G_1(A')) \cdot (A - A') = \det \begin{pmatrix} A - A' \\ i(G_1(A) - G_1(A')) \end{pmatrix}$$
$$= \det([z - z', H(z) - H(z')])$$
$$= |z - z'|^2 - |H(z) - H(z')|^2$$
$$\ge (1 - k^2)|z - z'|^2.$$

Since F is (1+k)-Lipschitz, this implies

$$(G_1(A) - G_1(A')) \cdot (A - A') \ge \lambda |A - A'|^2$$

for $\lambda = (1-k^2)/(1+k)^2 = (1-k)/(1+k) > 0$. This concludes the proof of Proposition 1.3. \square

4 Proof of Theorem 1.1

Step 1. We may assume that u is Lipschitz, thanks to the truncation argument of [14, Proposition A.1]. Let $K \subset B_R \subset \mathbb{R}^{2\times 2}$ for some R > 0. Then for all $X \in \mathbb{R}^{2\times 2}$ with |X| > 2R, we have $|X| \le 2 \operatorname{dist}(X, K)$. An application of [14, Proposition A.1] with $\lambda = 2R$ gives $v: B_1 \to \mathbb{R}^2$ satisfying

$$||Dv||_{L^{\infty}(B_1)} \le 2C_0 R,$$

$$\int_{B_1} |Du - Dv|^2 dx \le C_0 \int_{\{x \in B_1: |Du| > 2R\}} |Du|^2 dx$$

$$\le 4C_0 \int_{B_1} \operatorname{dist}^2(Du, K) dx,$$

for some constant C_0 (depending only on B_1). If there exists $M \in K$ such that

$$\int_{B_{1/2}} |Dv - M|^2 dx \le \widetilde{C}(K) \int_{B_1} \operatorname{dist}^2(Dv, K) dx,$$

then repeated applications of the triangle inequality give

$$\int_{B_{1/2}} |Du - M|^2 dx \le \left(16\widetilde{C}(K)C_0 + 4\widetilde{C}(K) + 8C_0\right) \int_{B_1} \operatorname{dist}^2(Du, K) dx.$$

Thus if Theorem 1.1 holds for all Lipschitz mappings v for some constant $\widetilde{C}(K)$, then it is also valid for all H^1 mappings with $C(K) = 16\widetilde{C}(K)C_0 + 4\widetilde{C}(K) + 8C_0$.

Step 2. We may assume in addition that $u \in C^2(B_1; \mathbb{R}^2)$ solves

$$\operatorname{div} G_1(Du_1) = \operatorname{div} G_2(Du_2) = 0$$
 in B_1 . (7)

Consider indeed $w \in C^2(B_1)$ such that w = u on ∂B_1 and

$$\operatorname{div} G_1(Dw_1) = \operatorname{div} G_2(Dw_2) = 0$$
 in B_1 .

The existence of such w is guaranteed by the ellipticity of the equation $0 = \operatorname{div} G_j(Dw_j) = \operatorname{tr}(DG_j(Dw_j)D^2w_j) = 0$, invoking e.g. [16, Theorem 12.5]: the inequality $\lambda |\xi|^2 \leq DG_j(A)\xi \cdot \xi \leq \Lambda |\xi|^2$, valid for all $A, \xi \in \mathbb{R}^2$ thanks to Proposition 1.3, ensures that the eigenvalues of the symmetric part $[DG_j(A)]_s = (DG_j(A) + DG_j(A)^T)/2$ of $DG_j(A)$ are bounded above and below (since $DG_j(A)\xi \cdot \xi = [DG_j(A)]_s\xi \cdot \xi$ for all $\xi \in \mathbb{R}^2$) and in particular condition (ii) in [16, Theorem 12.5] is satisfied. Letting v = u - w and using the uniform monotonicity of G_1 we find

$$\lambda \int_{B_1} |Dv_1|^2 dx \le \int_{B_1} (G_1(Du_1) - G_1(Dw_1)) \cdot Dv_1 dx.$$

Since div $G_1(Dw_1) = 0$ and div $(iDu_2) = 0$ we rewrite this as

$$\lambda \int_{B_1} |Dv_1|^2 dx \le \int_{B_1} (G_1(Du_1) + iDu_2) \cdot Dv_1 dx$$

$$\le \frac{1}{2\lambda} \int_{B_1} |G_1(Du_1) + iDu_2|^2 dx + \frac{\lambda}{2} \int_{B_1} |Dv_1|^2 dx,$$

and infer

$$\int_{B_1} |Dv_1|^2 dx \le \frac{1}{\lambda^2} \int_{B_1} |G_1(Du_1) + iDu_2|^2 dx.$$

According to Proposition 1.3 the function $M \mapsto G_1(A) + iB$, where A, B denote the first and second row of the matrix M, vanishes on K. Since that function is Lipschitz we deduce that $|G_1(Du_1) + iDu_2| \leq C \operatorname{dist}(Du, K)$, and therefore

$$\int_{B_1} |Dv_1|^2 dx \le C \int_{B_1} \operatorname{dist}^2(Du, K) dx.$$

Applying a similar argument to v_2 we obtain

$$\int_{B_1} |Dv|^2 dx \le C \int_{B_1} \operatorname{dist}^2(Du, K) dx.$$

Recalling that v = u - w and using the triangle inequality we deduce

$$\int_{B_1} \operatorname{dist}^2(Dw, K) \, dx \le C \int_{B_1} \operatorname{dist}^2(Du, K) \, dx,$$

$$\int_{B_{1/2}} |Du - M|^2 \, dx \le 2 \int_{B_{1/2}} |Dw - M|^2 \, dx + C \int_{B_1} \operatorname{dist}^2(Du, K) \, dx.$$

As a consequence, if Theorem 1.1 is valid for w then we obtain it for u. This proves Step 2.

Step 3. As $u_j \in C^2(B_1)$ satisfies (7), it is a weak solution of

$$\operatorname{div}\left(\partial_i(G_j(Du_j))\right) = \operatorname{div}\left(DG_j(Du_j)D(\partial_i u_j)\right) = 0 \qquad \text{ for } i = 1, 2.$$

Invoking the De Giorgi-Nash estimates [9] (see e.g. [17, Theorem 4.11] for the precise statement we use here) for $\partial_i u_j$, $i, j \in \{1, 2\}$, we obtain

$$||Du||_{C^{\alpha}(\overline{B}_{1/2})} \le C||Du||_{L^{2}(B_{1})},$$

for some $\alpha > 0$ and some constant C = C(K). Thanks to this estimate and the exact rigidity obtained in Lemma 2.1, we may argue exactly as in [13, Lemma 4.5] to deduce that

$$\inf_{M \in K} ||Du - M||_{L^{\infty}(B_{1/2})} \le \rho \left(\int_{B_1} \operatorname{dist}^2(Du, K) \, dx \right), \tag{8}$$

for some function ρ depending only on K and satisfying $\rho(\epsilon) \to 0$ as $\epsilon \to 0$.

Step 4. We finally combine Step 3 with the linearized estimate of Lemma 2.2, to obtain our main estimate (3). The basic idea, as in [14, 13], is to linearize $\operatorname{dist}^2(\cdot, K)$ around $M_0 \in K$ such that $|Du - M_0|$ is uniformly small. When doing so, (3) formally turns into the linearized estimate (4) of Lemma 2.2, and it remains to control the error terms. Due to the modification of u arising from the translation $X \in T_M K$ in the left-hand side of the linearized estimate (4), it is not directly obvious that the error terms are negligible. In [14] this problem is absent because their equivalent of (8) comes with an explicit $\rho(\epsilon) = C\epsilon^{\frac{1}{4}}$. In [13] it is taken care of via a topological degree argument [13, Proposition 4.7] (see also [28]) which allows to avoid the translation. It is however not obvious to us how to apply that topological degree argument in the absence of a group structure on K, and we present here an alternative method relying on elementary estimates.

We assume without loss of generality that

$$\int_{B_1} \operatorname{dist}^2(Du, K) \, dx = \epsilon \le \epsilon_0,\tag{9}$$

where $\epsilon_0 = \epsilon_0(K)$ is to be chosen in the course of the proof. If (9) is not valid then (3) is automatically satisfied for a large enough constant C because the left-hand side of (3) is bounded thanks to Step 1.

We fix $\delta_0 > 0$ depending only on K, such that the nearest-point projection Π_K onto K is uniquely defined and smooth in the neighborhood $\mathcal{N}_{2\delta_0}(K)$. We first choose ϵ_0 small enough that $\rho(\epsilon_0) \leq \delta_0$, so thanks to (8) the projection $\Pi_K(Du)$ is well-defined.

We claim that, for every $M \in K$, there exists $Y_M \in K$ such that

$$\int_{B_{1/2}} |Du - Y_M|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du, K) dx + C \int_{B_{1/2}} |Du - M|^4 dx,$$
and
$$|M - Y_M|^2 \le C \int_{B_{1/2}} |Du - M|^2 dx.$$
(10)

Here and in the rest of this proof we denote by C > 0 a generic constant depending only on K.

To prove (10), we first invoke Lemma 2.2, according to which we have

$$\inf_{X \in T_M K} \int_{B_{1/2}} |Du - M - X|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du - M, T_M K) dx.$$

Choosing $X = X_M \in T_M K$ attaining the infimum in the left-hand side, we obtain

$$\int_{B_{1/2}} |Du - M - X_M|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du - M, T_M K) dx. \tag{11}$$

Moreover the minimizing property of X_M implies that $\int_{B_{1/2}} (Du - M - X_M) dx$ is orthogonal to X_M . We deduce

$$|X_M|^2 = \int_{B_{1/2}} X_M \cdot (Du - M) \, dx \le \frac{1}{2} |X_M|^2 + C \int_{B_{1/2}} |Du - M|^2 \, dx,$$

and therefore

$$|X_M|^2 \le C \int_{B_{1/2}} |Du - M|^2 dx. \tag{12}$$

Recalling from the proof of Lemma 2.2 that P_M denotes the orthogonal projection onto $(T_M K)^{\perp}$, we estimate the integrand in the right-hand side of (11) using

$$dist(Du - M, T_M K) = |P_M(Du - M)|$$

$$\leq |Du - \Pi_K(Du)| + |\Pi_K(Du) - M - (I - P_M)(Du - M)|$$

$$\leq |Du - \Pi_K(Du)| + C|Du - M|^2.$$

The last inequality follows from the fact that $I - P_M = D\Pi_K(M)$. Since $|Du - \Pi_K(Du)| = \text{dist}(Du, K)$, plugging this into (11) we find

$$\int_{B_{1/2}} |Du - M - X_M|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du, K) dx + C \int_{B_{1/2}} |Du - M|^4 dx \tag{13}$$

Now we may choose $Y_M \in K$ such that

$$|M + X_M - Y_M| \le C|X_M|^2. (14)$$

Indeed, if $|X_M| \leq \delta_0$ then one can simply take $Y_M = \Pi_K(M + X_M)$ and use the fact that $D\Pi_K(M)X_M = (I - P_M)X_M = X_M$, and if $|X_M| \geq \delta_0$ one may take $Y_M = M$. From (13) and (14) we infer

$$\int_{B_{1/2}} |Du - Y_M|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du, K) dx + C \int_{B_{1/2}} |Du - M|^4 dx + C |X_M|^4.$$

Using (12) and Cauchy-Schwarz to estimate the last term, we deduce the first inequality in (10), and the estimate on $|M - Y_M|$ in (10) follows from (14) and (12) (taking into account that $|X_M| \leq C$ thanks to (12) and Step 1).

Our goal is to find $M \in K$ for which we can discard the last term in (10). To that end, we apply (10) to define by induction a sequence $(M_i) \subset K$ satisfying

$$M_{j+1} = Y_{M_j} \qquad \forall j \ge 0,$$

and M_0 to be chosen later. According to (10) we have

$$\int_{B_{1/2}} |Du - M_{j+1}|^2 dx \le C \int_{B_{1/2}} \operatorname{dist}^2(Du, K) dx + C \int_{B_{1/2}} |Du - M_j|^4 dx,$$
and
$$|M_j - M_{j+1}|^2 \le C \int_{B_{1/2}} |Du - M_j|^2 dx,$$

Therefore, setting

$$q_j = \int_{B_{1/2}} |Du - M_j|^2 \quad \text{and} \quad s_j = \sup_{B_{1/2}} |Du - M_j|,$$
 (15)

and recalling the definition of $\epsilon = \int_{B_1} \operatorname{dist}^2(Du, K) dx$ (9), we find

$$q_{j+1} \le C\epsilon + Cs_j^2 q_j$$
 and $s_{j+1} \le s_j + Cq_j^{\frac{1}{2}}$. (16)

We wish to find $j_0 \ge 0$ such that

$$q_{j_0} \le \Lambda \epsilon, \tag{17}$$

for some constant $\Lambda > 0$ depending only on K. Recalling the definitions of q_j (15) and ϵ (9), this concludes the proof of (3). We prove by induction on j_0 that, if $s_0 \leq 1/(2\Lambda)$ and

$$q_j > \Lambda \epsilon \quad \text{for } j = 0, \dots, j_0,$$
 (18)

then

$$q_j \le (\Lambda s_0)^{2j} s_0^2$$
 for $j = 0, \dots, j_0$. (19)

The constant $\Lambda > 0$ depends only on K and will be adjusted in the course of the proof. Initialization at $j_0 = 0$ is obvious since $|B_{1/2}| = \pi/4 \le 1$. Next we assume (18)-(19) and prove that either (17) or (19) is satisfied at $j_0 + 1$. Combining (19) with the second inequality in (16), we find

$$s_{j_0} \le s_0 + Cs_0 \sum_{j=0}^{j_0-1} (\Lambda s_0)^j \le (1+C)s_0,$$

provided $2\Lambda s_0 \leq 1$. Therefore the first inequality in (16) gives

$$q_{j_0+1} \le C\epsilon + C(1+C)^2 s_0^2 q_{j_0},$$

and using (19) we obtain

$$q_{j_0+1} \le C\epsilon + C(1+C)^2 s_0^4 (\Lambda s_0)^{2j_0}$$
$$= C\epsilon + \frac{C(1+C)^2}{\Lambda^2} (\Lambda s_0)^{2j_0+2} s_0^2.$$

Therefore we have

either
$$q_{j_0+1} \le 2C\epsilon$$
, or $q_{j_0+1} \le 2\frac{C(1+C)^2}{\Lambda^2}(\Lambda s_0)^{2j_0+2}s_0^2$.

In the first case (17) is satisfied for $j = j_0 + 1$ provided $\Lambda \geq 2C$, and in the second case (19) is satisfied at $j_0 + 1$, provided $\Lambda \geq (2C)^{1/2}(1+C)$. This concludes the proof of (19) by induction

As a consequence, assuming by contradiction that there is no $j_0 \ge 0$ satisfying (17), we would have

$$q_i > \Lambda \epsilon$$
 and $q_i \leq (\Lambda s_0)^{2j} s_0^2$ $\forall j \geq 0$.

As $s_0 \leq 1/(2\Lambda)$, this implies $q_j \to 0$ and, passing to the limit in $q_j > \Lambda \epsilon$ we deduce $\epsilon = 0$. But in that case by Lemma 2.1 we have $Du \equiv M_0$ and this choice of M_0 gives $q_0 = 0$.

We conclude that there exists $j_0 = j_0(\epsilon)$ satisfying (17), if we can choose M_0 in such a way that $s_0 \leq 1/(2\Lambda)$. But thanks to (8) we do have $s_0 \leq \rho(\epsilon) \leq 1/(2\Lambda)$ provided ϵ_0 is small enough.

5 A 3×3 counter-example

In this section we prove that the two-dimensional setting of Theorem 1.1 is optimal in the following sense: a connected 1-submanifold of $\mathbb{R}^{3\times3}$ which has no rank-one connection and is elliptic may not satisfy Šverák's compactness result [31], and even less a quantitative rigidity estimate.

We recall (see e.g. [30, § 2]) that an ordered set of $N \geq 4$ matrices $\{T_i\}_{i=1}^N \subset \mathbb{R}^{m \times n}$ without rank-one connections is said to form a \mathcal{T}_N configuration if there exist matrices $P_i, C_i \in \mathbb{R}^{m \times n}$ and numbers $\kappa_i > 1$ such that

$$T_1 = P + \kappa_1 C_1,$$

 $T_2 = P + C_1 + \kappa_2 C_2,$
...
 $T_N = P + C_1 + C_2 + \dots C_{N-1} + \kappa_N C_N,$

where C_i is rank-one for all i and $\sum_{i=1}^{N} C_i = 0$.

Proposition 5.1. There exists a smooth, compact and connected 1-submanifold $K \subset \mathbb{R}^{3\times 3}$ without boundary which is elliptic and has no rank-one connection, but contains a \mathcal{T}_4 configuration.

By a known construction, see e.g. [7, Theorem 3.1], Proposition 5.1 implies the existence of a sequence of maps $u_k \colon \mathbb{R}^3 \to \mathbb{R}^3$ such that

$$\int_{B_1} \operatorname{dist}^2(Du_k, K) \, dx \to 0 \quad \text{as } k \to 0,$$

but (Du_k) is not precompact in $L^2(B_{1/2})$. In particular, one certainly cannot hope for a quantitative estimate

$$\inf_{M \in K} \int_{B_{1/2}} |Du_k - M|^2 dx \le \rho \left(\int_{B_1} \operatorname{dist}^2(Du_k, K) dx \right),$$

for any function $\rho(\epsilon) \to 0$ as $\epsilon \to 0$.

Proposition 5.1 is a consequence of the construction below. Let a > 0 and define matrices T_1, T_2, T_3, T_4 by

$$T_1 = -T_3 = \begin{pmatrix} 1+a & 0 & 0 \\ 0 & 1 & 0 \\ 1+a & 0 & 0 \end{pmatrix}, \qquad T_2 = -T_4 = \begin{pmatrix} -1 & 0 & 0 \\ 0 & 1+a & 0 \\ -1 & 0 & 0 \end{pmatrix},$$

and C_1, C_2, C_3, C_4 by

$$C_1 = -C_3 = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix},$$
 $C_2 = -C_4 = \begin{pmatrix} -1 & 0 & 0 \\ 0 & 1 & 0 \\ -1 & 0 & 0 \end{pmatrix}.$

We have that $T_k - C_k$ is rank-one for k = 1, 2, 3, 4, and (with the convention that $C_5 = C_1$)

$$T_k - \frac{2+a}{a}(T_k - C_k) = C_{k+1}$$
 for $k = 1, 2, 3, 4$,

so $\{T_1, T_2, T_3, T_4\}$ forms a \mathcal{T}_4 configuration. Next we construct a curve $K = K_a$ as in Proposition 5.1, which contains this \mathcal{T}_4 configuration. Let $\theta_a = \arctan(1/(1+a))$, so

$$\cos \theta_a = \frac{1+a}{r_a}, \quad \sin \theta_a = \frac{1}{r_a}, \quad r_a = \sqrt{1+(1+a)^2}.$$

Let $\rho \colon \mathbb{R} \to \mathbb{R}$ be a smooth monotonically increasing function to be determined later that satisfies

- $\rho(\theta + 2\pi) = \rho(\theta) + 2\pi$ for $\theta \in \mathbb{R}$.
- $\rho(\hat{\theta}_k) = \hat{\theta}_k \text{ for } \hat{\theta}_k = \theta_a + (k-1)\frac{\pi}{2}, k = 1, 2, 3, 4.$

Define $K_a = \Gamma_a(\mathbb{R}/2\pi\mathbb{Z})$, with

$$\Gamma_{a}(\theta) := \begin{pmatrix} r_{a}\cos(\theta) & -\sin(8\theta - 8\theta_{a}) & \sin(6\rho(\theta) - 6\theta_{a}) \\ \sin(6\theta - 6\theta_{a}) & r_{a}\sin(\theta) & \sin(8\rho(\theta) - 8\theta_{a}) \\ r_{a}\cos(\theta) & \sin(8\theta - 8\theta_{a}) & \sin(6\rho(\theta) - 6\theta_{a}) \end{pmatrix}$$

$$(20)$$

Then we have

$$T_k = \Gamma_a(\hat{\theta}_k) \text{ for } k = 1, 2, 3, 4,$$

so K_a contains the \mathcal{T}_4 configuration $\{T_1, T_2, T_3, T_4\}$. Next we adjust the parameter a > 0 and the function ρ in order to ensure that K_a has no rank-one connection and is elliptic.

Notation. With the $M_{i_1i_2,j_1j_2}$ minor we mean the determinant of the 2×2 submatrix corresponding to the rows i_1, i_2 and columns j_1, j_2 .

Lemma 5.2. If a > 0 is such that $\theta_a \notin \frac{\pi}{48}\mathbb{Z}$, the curve K_a is elliptic, i.e. Rank $\Gamma'_a(\theta) > 1$ for all $\theta \in \mathbb{R}$.

Proof. The derivative Γ'_a is given by

$$\Gamma_a'(\theta) := \begin{pmatrix} -r_a \sin(\theta) & -8\cos(8\theta - 8\theta_a) & 6\rho'(\theta)\cos(6\rho(\theta) - 6\theta_a) \\ 6\cos(6\theta - 6\theta_a) & r_a\cos(\theta) & 8\rho'(\theta)\cos(8\rho(\theta) - 8\theta_a) \\ -r_a \sin(\theta) & 8\cos(8\theta - 8\theta_a) & 6\rho'(\theta)\cos(6\rho(\theta) - 6\theta_a) \end{pmatrix}$$

Assume Rank $\Gamma'_{a}(\theta) \leq 1$ for some $\theta \in \mathbb{R}$. Then calculating the $M_{12,12}$ minor we have

$$-r_a^2 \sin(\theta)\cos(\theta) + 48\cos(8\theta - 8\theta_a)\cos(6\theta - 6\theta_a) = 0$$

and calculating the $M_{23,12}$ minor we have

$$48\cos(8\theta - 8\theta_a)\cos(6\theta - 6\theta_a) + r_a^2\sin(\theta)\cos(\theta) = 0.$$

Adding and subtracting these two equations we obtain that

$$\sin(\theta)\cos(\theta) = 0$$
 and $\cos(8\theta - 8\theta_a)\cos(6\theta - 6\theta_a) = 0$.

The first equality implies $\theta \in \frac{\pi}{2}\mathbb{Z}$, and then the second equality becomes

$$\cos(8\theta_a)\cos(6\theta_a) = 0,$$

which is impossible for $\theta_a \notin \frac{\pi}{48}\mathbb{Z}$.

Lemma 5.3. For any a > 0 such that $\theta_a \notin \frac{\pi}{24}\mathbb{Z}$, for any $\epsilon > 0$, there exists a smooth monotonic function $\rho \colon \mathbb{R} \to \mathbb{R}$ with the following properties

- $\rho(\theta + 2\pi) = \rho(\theta) + 2\pi \text{ for } \theta \in \mathbb{R},$
- $\rho(\hat{\theta}_k) = \hat{\theta}_k \text{ for } \theta_k = \hat{\theta}_a + (k-1)\frac{\pi}{2}, \ k = 1, 2, 3, 4.$
- For any $\theta, \theta' \in [0, 2\pi) \cap \frac{\pi}{24}\mathbb{Z}$, we have $\rho(\theta) \rho(\theta') \notin \frac{\pi}{12}\mathbb{Z}$
- $\sup_{x \in \mathbb{R}} |\rho(x) x| < \epsilon$

Proof. Let $\delta = \frac{1}{2}\operatorname{dist}(\theta_a, \frac{\pi}{24}\mathbb{Z}) = \frac{1}{2}\operatorname{dist}(\{\hat{\theta}_k\}, \frac{\pi}{24}\mathbb{Z}) > 0$ and fix a smooth function φ such that $\sup \varphi \subset (-\delta, \delta), \ 0 \le \varphi \le \varphi(0) = 1$ and $|\varphi'| \le 2/\delta$. Define ρ on $[0, 2\pi)$ by setting

$$\rho(\theta) = \theta + \sum_{j=1}^{47} t_j \, \varphi\left(\theta - j\frac{\pi}{24}\right) \quad \text{for } \theta \in [0, 2\pi),$$

where $t_1, \ldots, t_{47} \in (-\eta, \eta)$ are to be fixed later and $\eta = \min(\epsilon, \delta/2)$. The choice of $\delta > 0$ ensures that $\rho(\hat{\theta}_k) = \hat{\theta}_k$, also since $|t_j| < \epsilon$ we have $|\rho - id| < \epsilon$, and finally since $|t_j| < \delta/2$ we have $\rho' > 0$ on $[0, 2\pi)$. Moreover the function $\rho - id$ is identically zero near 0 and 2π , so it can be extended to a smooth 2π periodic function, thus yielding a smooth monotonic extension $\rho \colon \mathbb{R} \to \mathbb{R}$ satisfying $\rho(\theta + 2\pi) = \rho(\theta) + 2\pi$ for $\theta \in \mathbb{R}$. It remains to argue that we can pick $t_1, \ldots, t_{47} \in [0, \eta)$ to ensure that the third condition in Lemma 5.3 is satisfied.

Denote $t_0 = 0$. By induction, we may for each j = 1, ..., 47 choose $t_j \in (-\eta, \eta)$ to ensure that

$$\rho(j\pi/24) - \rho(\ell\pi/24) = t_j - t_\ell + (j-\ell)\pi/24 \notin \frac{\pi}{12}\mathbb{Z}$$

for all $\ell \in \{0, ..., j-1\}$. This is possible because at each step there is only a discrete set of values of t_j to avoid.

Lemma 5.4. If a > 0 is such that $\theta_a \notin \frac{\pi}{48}\mathbb{Z}$, $\epsilon > 0$ is small enough and ρ is as in Lemma 5.3, then the curve $K_a \subset \mathbb{R}^{3\times 3}$ does not contain Rank-1 connections.

Proof. Note as in Lemma 5.2 that the assumption $\theta_a \notin \frac{\pi}{48}\mathbb{Z}$ implies

$$\cos(8\theta_a) \neq 0$$
 and $\cos(6\theta_a) \neq 0$.

We assume that there exist $\theta \neq \theta' \in \mathbb{R}/2\pi\mathbb{Z}$ such that

$$\operatorname{Rank}\left(\Gamma_a(\theta) - \Gamma_a(\theta')\right) = 1,\tag{21}$$

and we obtain a contradiction. We do this in several steps.

Step 1. We have

$$\theta + \theta' \in \pi \mathbb{Z}$$
 and $\theta - \theta' \in \frac{\pi}{3} \mathbb{Z} \cup \frac{\pi}{4} \mathbb{Z}$. (22)

Proof of Step 1. From (20) calculating the $M_{12,12}$ minor we have

$$0 = r_a^2 \left(\cos(\theta) - \cos(\theta')\right) \left(\sin(\theta) - \sin(\theta')\right)$$

$$+ \left(\sin(8\theta - 8\theta_a) - \sin(8\theta' - 8\theta_a)\right) \left(\sin(6\theta - 6\theta_a) - \sin(6\theta' - 6\theta_a)\right)$$

$$= -4r_a^2 \sin\left(\frac{\theta + \theta'}{2}\right) \cos\left(\frac{\theta + \theta'}{2}\right) \sin^2\left(\frac{\theta - \theta'}{2}\right)$$

$$+ 4\sin\left(3(\theta - \theta')\right) \sin\left(4(\theta - \theta')\right) \cos\left(3(\theta + \theta') - 6\theta_a\right) \cos\left(4(\theta + \theta') - 8\theta_a\right).$$
 (23)

And calculating the $M_{23,12}$ minor we have

$$0 = -r_a^2 \left(\cos(\theta) - \cos(\theta')\right) \left(\sin(\theta) - \sin(\theta')\right)$$

$$+ \left(\sin(8\theta - 8\theta_a) - \sin(8\theta' - 8\theta_a)\right) \left(\sin(6\theta - 6\theta_a) - \sin(6\theta' - 6\theta_a)\right)$$

$$= 4r_a^2 \sin\left(\frac{\theta + \theta'}{2}\right) \cos\left(\frac{\theta + \theta'}{2}\right) \sin^2\left(\frac{\theta - \theta'}{2}\right)$$

$$+ 4\sin\left(3\left(\theta - \theta'\right)\right) \sin\left(4\left(\theta - \theta'\right)\right) \cos\left(3\left(\theta + \theta'\right) - 6\theta_a\right) \cos\left(4\left(\theta + \theta'\right) - 8\theta_a\right).$$
 (24)

Adding and subtracting (23) and (24) we obtain the equations

$$0 = \sin(3(\theta - \theta'))\sin(4(\theta - \theta'))\cos(3(\theta + \theta') - 6\theta_a)\cos(4(\theta + \theta') - 8\theta_a)$$
(25)

and

$$0 = \sin\left(\frac{\theta + \theta'}{2}\right)\cos\left(\frac{\theta + \theta'}{2}\right)\sin^2\left(\frac{\theta - \theta'}{2}\right). \tag{26}$$

Since $\theta \neq \theta'$ in $\mathbb{R}/2\pi\mathbb{Z}$, the last factor of (26) is nonzero, so either the first or the second must be zero. This implies $\theta + \theta' \in \pi\mathbb{Z}$. As a consequence, the last two factors in (25) are equal to $\pm \cos(6\theta_a)$ and $\cos(8\theta_a)$ and are nonzero by our choice of a. So one of the first two factors of (25) must vanish, that is, $\theta - \theta' \in \frac{\pi}{3}\mathbb{Z} \cup \frac{\pi}{4}\mathbb{Z}$.

Step 2. We have
$$\theta - \theta' \in \frac{\pi}{4}\mathbb{Z}$$
.

Proof of Step 2. Considering the $M_{13,23}$ minor of $\Gamma_a(\theta) - \Gamma_a(\theta')$ we obtain the equation

$$0 = \left(\sin\left(8\theta - 8\theta_{a}\right) - \sin\left(8\theta' - 8\theta_{a}\right)\right) \times \left(\sin\left(6\rho\left(\theta\right) - 6\theta_{a}\right) - \sin\left(6\rho\left(\theta'\right) - 6\theta_{a}\right)\right)$$
$$= 4\sin\left(4\left(\theta - \theta'\right)\right)\cos\left(4\left(\theta + \theta'\right) - 8\theta_{a}\right) \times \sin\left(3\left(\rho\left(\theta\right) - \rho\left(\theta'\right)\right)\right)\cos\left(3\left(\rho\left(\theta\right) + \rho\left(\theta'\right)\right) - 6\theta_{a}\right).$$

From Step 1 we have $\theta + \theta' \in \pi \mathbb{Z}$ so the second factor is $\cos(8\theta_a) \neq 0$. The last factor is arbitrarily close to $\pm \cos(6\theta_a) \neq 0$ since $|\rho - id| \leq \epsilon$. The third factor is nonzero by construction of ρ (recall Lemma 5.3) because $\theta, \theta' \in \frac{\pi}{24}\mathbb{Z}$ by (22). So we must have $\sin(4(\theta - \theta')) = 0$ hence $\theta - \theta' \in \frac{\pi}{4}\mathbb{Z}$.

Step 3. We have $\theta - \theta' \in \pi \mathbb{Z}$.

Proof of Step 3. Considering the $M_{12,13}$ minor of $\Gamma_a(\theta) - \Gamma_a(\theta')$ we obtain the equation

$$0 = r_a \left(\cos(\theta) - \cos(\theta')\right) \left(\sin(8\rho(\theta) - 8\theta_a) - \sin(8\rho(\theta') - 8\theta_a)\right) - \left(\sin(6\theta - 6\theta_a) - \sin(6\theta' - 6\theta_a)\right) \times \left(\sin(6\rho(\theta) - 6\theta_a) - \sin(6\rho(\theta') - 6\theta_a)\right),$$

which we rewrite as

$$-4r_{a} \sin\left(\frac{\theta+\theta'}{2}\right) \sin\left(\frac{\theta-\theta'}{2}\right)$$

$$\times \sin\left(4\left(\rho\left(\theta\right)-\rho\left(\theta'\right)\right)\right) \cos\left(4\left(\rho\left(\theta\right)+\rho\left(\theta'\right)\right)-8\theta_{a}\right)$$

$$=4 \sin\left(3\left(\theta-\theta'\right)\right) \cos\left(3\left(\theta+\theta'\right)-6\theta_{a}\right)$$

$$\times \sin\left(3\left(\rho\left(\theta\right)-\rho\left(\theta'\right)\right)\right) \cos\left(3\left(\rho\left(\theta\right)+\rho\left(\theta'\right)-6\theta_{a}\right)\right). \tag{27}$$

Since $\theta - \theta' \in \frac{\pi}{4}\mathbb{Z}$ and $|\rho - id| \leq \epsilon$, the third factor in the left-hand side of (27) has absolute value $\leq 4\epsilon$. Since $\theta + \theta' \in \pi\mathbb{Z}$, the second factor in the right-hand side of (27) has absolute value equal to $|\cos(6\theta_a)| > 0$, and for small enough ϵ the absolute value of the last factor is $\geq |\cos(6\theta_a)|/2 > 0$. Taking also into account that the first and third factor in the right-hand side of (27) differ from each other by an error $\leq 3\epsilon$, we must have

$$\sin^2\left(3\left(\theta - \theta'\right)\right) \le c_a \epsilon,$$

for some $c_a > 0$ depending only on a. Because $\theta - \theta' \in \frac{\pi}{4}\mathbb{Z}$ by Step 2, provided ϵ is chosen small enough, this implies $\theta - \theta' \in \frac{\pi}{3}\mathbb{Z} \cap \frac{\pi}{4}\mathbb{Z} = \pi\mathbb{Z}$.

Step 4: Conclusion. From Step 1 and Step 3 we have $\theta + \theta', \theta - \theta' \in \pi \mathbb{Z}$, so $\theta, \theta' \in \pi \mathbb{Z}$. Since me may without loss of generality exchange the roles of θ and θ' and choose arbitrary representants in $\mathbb{R}/2\pi\mathbb{Z}$, this amounts to $\theta = 0$ and $\theta' = \pi$. Considering the $M_{12,13}$ minor of $\Gamma_a(0) - \Gamma_a(\pi)$ as in (27) we deduce

$$\sin(4(\rho(\pi) - \rho(0)))\cos(4(\rho(\pi) + \rho(0)) + 8\theta_a) = 0.$$

Using again that $|\rho - id| \le \epsilon$, the second factor has absolute value $\ge |\cos(8\theta_a)|/2$ provided ϵ is small enough, and the first factor is nonzero by construction of ρ , so we conclude that (21) is not possible for $\theta \ne \theta' \in \mathbb{R}/2\pi\mathbb{Z}$.

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