# Extremal problems in Euclidean combinatorial geometry 

Danila Cherkashin

March 12, 2024

## Contents

1 Introduction ..... 4
2 Steiner trees ..... 6
2.1 Basics ..... 6
2.1.1 Topology and embedding class of a tree ..... 8
2.1.2 Connectedness in $\mathbb{P}_{d}$ ..... 12
2.2 A universal Steiner tree ..... 13
2.3 Connectivity of the subset of $\mathbb{P}_{d}$ with a unique Steiner tree ..... 18
2.3.1 Canonical realization of an embedding class ..... 18
2.3.2 Proof of Theorem [2.1.4 ..... 18
2.4 Proof of Theorem|2.1.1 ..... 22
2.4.1 Subanalytic subsets of a real analytic manifold ..... 22
2.4.2 Proof of Theorem 2 2.1.1 ..... 24
2.5 Steiner trees in real analytic Riemannian manifolds ..... 29
2.5.1 $\quad$ Realization space $\mathcal{R}_{\text {geoemb }}$ for an arbitrary metric space ..... 29
2.5.2 $\quad$ Manifold structure on $\mathcal{R}_{\text {geoemb }}$ ..... 30
2.5.3 $\quad$ The subvariety $\mathcal{R}_{\text {locmin }}$ of $\mathcal{R}_{\text {geoemb }}$ ..... 31
2.5.4 $\quad$ Example of $\mathcal{A}_{T_{1}, T_{2}}$ with non-empty interior ..... 32
3 Gilbert-Steiner problem ..... 34
3.1 Basics ..... 34
3.2 Preliminaries ..... 35
3.3 Main result ..... 37
3.4 Examples of branching points of degree 4 ..... 38
3.5 Open questions ..... 40
4 Maximal distance minimizers ..... 41
4.1 Introduction ..... 41
4.1.1 Class of problems ..... 41
4.1.2 Dual problem ..... 41
4.1.3 The first parallels with average distance minimization problem ..... 42
4.1.4 Notation ..... 42
4.1.5 $\quad$ Existence. Absence of loops. Ahlfors regularity and other simple properties ..... 42
4.1.6 Local maximal distance minimizers ..... 43
4.2 Regularity ..... 43
4.2.1 Tangent rays. Blow up limits in $\mathbb{R}^{d}$ ..... 43
4.2.2 Properties of branching points in $\mathbb{R}^{2}$ ..... 44
4.2.3 Continuity of one-sided tangents in $\mathbb{R}^{2}$ ..... 45
4.2.4 Planar example of infinite number of corner points ..... 46
4.2.5 Every $C^{1,1}$-smooth simple curve is a minimizer ..... 46
4.3 Explicit examples for maximal distance minimizers ..... 47
4.3.1 Simple examples. Finite number of points and $r$-neighbourhood. Inverse mini- mizers ..... 47
4.3.2 Circle. Curves with big radius of curvature ..... 48
4.3.3 Rectangle ..... 49
4.4 Tools ..... 49
4.4.1 Energetic points ..... 49
4.4.2 Convexity argument ..... 51
4.4.3 Lower bounds on the length of a minimizer ..... 51
4.5 More properties of minimizers ..... 52
$4.5 .1 \quad \Gamma$-convergence ..... 52
4.5.2 Approximation by Steiner trees ..... 52
4.5.3 NP-hardness ..... 53
4.5.4 Penalized form ..... 54
4.5.5 Uniqueness ..... 55
4.6 On minimizers for a planar convex closed smooth curve ..... 55
4.6.1 The class of $M$ considered in the section ..... 55
4.6.2 Pseudo-networks ..... 56
4.6.3 $\quad$ Structural properties of minimizers in the annulus $N \backslash N_{r}$ ..... 58
4.6.4 Derivation in the picture ..... 71
4.7 Horseshoe theorem ..... 76
4.7.1 Sketch of the proof ..... 76
4.7.2 Lemmas for the first step ..... 77
4.7.3 Central lemma ..... 82
4.7.4 Finishing the proof ..... 90
4.8 Open questions ..... 92
4.8.1 Regularity ..... 92
4.8.2 Explicit solutions ..... 93
4.8.3 Uniqueness ..... 94
5 Johnson-type graphs ..... 95
5.1 Basics ..... 95
5.1.1 Independence and chromatic numbers of $J(d, k, t)$ and $K(d, k, t)$ ..... 96
5.1.2 Known facts about the graphs $J_{ \pm}(d, k, t)$ and $K_{ \pm}(d, k, t)$ ..... 97
5.1.3 Results ..... 98
5.2 Tools ..... 100
5.2.1 Trivial bounds on the chromatic numbers ..... 100
5.2.2 Katona averaging method ..... 101
5.2.3 Nontrivial intersecting families ..... 101
5.2.4 An isodiametric inequality ..... 102
5.2.5 Simple hypergraphs and Reed-Solomon codes ..... 102
5.2.6 Steiner systems ..... 103
5.2.7 Families with even or odd intersections ..... 103
5.3 Examples ..... 104
5.4 Proofs ..... 105
5.4.1 Proof of Theorem 5 .1.7 ..... 105
5.4.2 Proof of Theorem 5 5.1.8 ..... 106
5.4.3 Proof of Theorem|5.1.9 ..... 108
5.4.4 Proof of Corollary 5.1.1 ..... 108
5.4.5 Proof of Theorem 5.1.10 ..... 108
5.4.6 Proof of Theorem 5.1.11 ..... 109
5.4.7 Proof of Theorem 5.1 .12 ..... 110
5.5 Independent numbers in the case $k \leq 3$ ..... 111
5.5.1 The case $k=2$ ..... 111
5.5.2 The case $k=3, t=-1$ ..... 112
5.5.3 The case $k=3, t=0$ ..... 112
5.5.4 The case $k=3, t=-2$ ..... 113
5.6 Open questions ..... 114
6 Chromatic numbers of 2-dimensional spheres ..... 116
6.1 Introduction ..... 116
6.2 Proof of Theorem|6.1.2 ..... 118
6.2.1 Step 1. Each color is a dense set ..... 118
6.2.2 Step 2. Stability of embedding ..... 121
6.3 Open questions ..... 125
7 Chromatic numbers of 3-dimensional slices ..... 127
7.1 Introduction ..... 127
7.1.1 Nelson-Hadwiger problem and its planar generalizations ..... 127
7.1.2 The chromatic numbers of real 3-dimensional slices ..... 128
7.1.3 The chromatic numbers of 2-dimensional rational slices ..... 128
7.2 Notation and auxiliary lemmas ..... 129
7.3 Proof of Theorem 17.1.3 ..... 132
7.4 Proof of Theorem 7.1 .2 ..... 135
7.5 Proof of Proposition 7.1.1 ..... 138
7.6 Further questions ..... 139

## Chapter 1

## Introduction

The dissertation is devoted to extremal problems in the intersection of Euclidean geometry and combinatorics. Consider a distance graph $G\left(\mathbb{R}^{d}\right)$ of a Euclidean space, which is a complete weighted graph with vertex set $\mathbb{R}^{d}$ and the weights from Euclidean metrics. A typical framework is $G$ or its "subgraph" $G(V, \rho)=\left(V, E_{\rho}\right)$, where $V$ is a subset of $\mathbb{R}^{d}$ and $E_{d}$ consists of pairs of vertices at a distance of $\rho$. We consider both finite and infinite $V$. We focus on several classical combinatorial problems: Steiner tree problem, finding a maximal independent set and finding the chromatic number. Note that these three problems belong to the initial Karp's list of 21 NP-complete problems.

Chapter 2 is devoted to the Euclidean Steiner tree problem. Theorem 2.1.1 states that for $d=2$ a random $n$-point input leads to a unique solution. Then Theorem 2.1.4 shows the connectedness of the set of $d$-dimensional $n$-point configurations having a unique Steiner tree (as a subset of $\left.\left(\mathbb{R}^{d}\right)^{n}\right)$. Also, Theorem 2.2.1 provides an example of a Steiner tree for an input $\mathcal{A}$ of a positive Hausdorff dimension, which cannot be considered as a union of the solutions for $\mathcal{A}_{1}, \mathcal{A}_{2} \subset \mathcal{A}$.

The Gilbert-Steiner problem is a generalization of the Steiner tree problem on a specific optimal mass transportation. The cost function for the transportation of a mass $m$ on distance $l$ is chosen to be $m^{p} \cdot l$, where $p \in(0,1)$ is a parameter. The difference with the optimal transportation problem is that the extra geometric points may be of use; such points are called branching. Chapter 3 proves that every branching point in a solution of the planar Gilbert-Steiner problem has degree 3, see Theorem 3.1.2 and more general Theorem 3.3.1.

In Chapter 4 we consider the problem of minimizing the maximal distance to a given compact set $M$ among the sets of a given length $\ell$. First, we give a survey on the results in this problem. Then we find maximal distance minimizers for a closed planar curve of a small enough curvature (see Theorem 4.3.2 and finish with open questions.

Chapter 5 has a deal with independence and chromatic numbers of Johnson-type graphs. The Johnson-type graph $J_{ \pm}(d, k, t)$ in defined the following way: the vertex set consists of all vectors from $\{-1,0,1\}^{d}$ with exactly $k$ nonzero coordinates; edges connect the pairs of vertices with scalar product $t$. Theorems 5.1.7 and 5.1.8 determine the independence number of $J_{ \pm}(d, k, t)$ for an odd negative $t$ and $d>d_{0}(k, t)$. Theorem 5.1 .12 shows that the asymptotic of the chromatic numbers for $k=3$, $t=-2$ is doubly logarithmic in $d$.

In Chapter 6 we show that for a positive $\rho<2$ the chromatic number of $G\left(\mathbb{S}^{2}, \rho\right)$ is at least four, where $\mathbb{S}^{2}$ is a 2-dimensional sphere with unit radius, see Theorem 6.1.2. Note that for $\rho=2$ the corresponding graph is a matching, so its chromatic number is two.

Finally, Theorem 7.1 .2 in Chapter 7 establishes that for every positive $\varepsilon$ the chromatic number of $\mathbb{R}^{3} \times[0, \varepsilon]^{6}$ is at least 10 . For a small enough $\varepsilon$ the upper bound is 15 and it comes from a well-known permutohedron tiling.

The thesis is based on the following papers and preprints:

1. "On the horseshoe conjecture for maximal distance minimizers", D. Cherkashin, Y. Teplitskaya, ESAIM: Control, Optimisation and Calculus of Variations 24 (3), 1015-1041, 2018 (IF2018 $1.295, \mathrm{Q} 2)$
2. "A self-similar infinite binary tree is a solution to the Steiner problem", D. Cherkashin, Y. Teplitskaya, Fractal and Fractional 7 (5), 414, 2023 (IF2022 5.4, Q1)
3. "Independence numbers of Johnson-type graphs", D. Cherkashin, S. Kiselev, Bulletin of the Brazilian Mathematical Society, New Series 54 (3), 30, 2023 (IF2022 0.7, Q3)
4. "On the chromatic number of 2-dimensional spheres", D. Cherkashin, V. Voronov, Discrete \& Computational Geometry, 71, 467-479, 2024 (IF2022 0.8, Q3)
5. "On the chromatic numbers of 3-dimensional slices", V.A. Voronov, A.Y. Kanel-Belov, G.A. Strukov, D.D. Cherkashin, Zapiski Nauchnykh Seminarov POMI 518, 94-113, 2022 (in Russian, English translation to appear in Journal of Mathematical Sciences, SJR2022 0.314, Q3)
6. "On the chromatic numbers of Johnson-type graphs", D.D. Cherkashin, Zapiski Nauchnykh Seminarov POMI 518, 192-200, 2022 (in Russian, English translation to appear in Journal of Mathematical Sciences, SJR2022 0.314, Q3)
7. "On uniqueness in Steiner problem", M. Basok, D. Cherkashin, N. Rastegaev, Y. Teplitskaya, preprint arXiv:1809.01463, to appear in IMRN, 2024 (IF2022 1.0, Q2)
8. "On minimizers of the maximal distance functional for a planar convex closed smooth curve" D.D. Cherkashin, A.S. Gordeev, G.A. Strukov, Y.I. Teplitskaya, preprint arXiv:2011.10463, submitted to St. Petersburg Mathematical Journal
9. "Branching points in the planar Gilbert-Steiner problem have degree 3", D. Cherkashin, F. Petrov, preprint arXiv:2309.04202, to appear in Pure and Applied Functional Analysis, the volume is dedicated to the memory of Anatoly Moiseevich Vershik
10. "An overview of maximal distance minimizers problem", D. Cherkashin, Y. Teplitskaya, preprint arXiv:2212.05607, submitted to Serdica Mathematical Journal

Acknowledgements. I would like to thank my coauthors, namely Mikhail Basok, Alexey Gordeev, Alexei Kanel-Belov, Sergei Kiselev, Fedor Petrov, Emanuele Paolini, Nikita Rastegaev, Georgii Strukov, Yana Teplitskaya and Vsevolod Voronov.

I am appreciated to Peter Boyvalenkov for the organization of the procedure.

## Chapter 2

## Steiner trees

This chapter is based of papers [5, 20]. We consider both finite and infinite forms of the Euclidean Steiner tree problem:

Problem 2.0.1. For a given finite set $P=\left\{x_{1}, \ldots, x_{n}\right\} \subset \mathbb{R}^{d}$ find a connected set $\mathcal{S}$ t with minimal length (one-dimensional Hausdorff measure $\mathcal{H}^{1}$ ) containing $P$.

Problem 2.0.2. For a given compact set $\mathcal{A} \subset \mathbb{R}^{d}$ find a set $\mathcal{S t}$ with minimal length (one-dimensional Hausdorff measure $\mathcal{H}^{1}$ ) such that $\mathcal{S} t \cup \mathcal{A}$ is connected.

Results on Problem 2.0.1. We prove that the set of $n$-point configurations for which the solution to the planar Steiner problem is not unique has the Hausdorff dimension at most $2 n-1$ (as a subset of $\left.\mathbb{R}^{2 n}\right)$. Moreover, we show that the Hausdorff dimension of the set of $n$-point configurations for which at least two locally minimal trees have the same length is also at most $2 n-1$. The methods we use essentially rely upon the theory of subanalytic sets developed in [9]. Motivated by this approach we develop a general setup for the similar problem of uniqueness of the Steiner tree where the Euclidean plane is replaced by an arbitrary analytic Riemannian manifold $M$. In this setup we argue that the set of configurations possessing two locally-minimal trees of the same length either has dimension equal to $n \operatorname{dim} M-1$ or has a non-empty interior. We provide an example of a two-dimensional surface for which the last alternative holds.

We study the set of $n$-point configurations for which there is a unique solution to the Steiner problem in $\mathbb{R}^{d}$. We show that this set is path-connected.

Results on Problem 2.0.2. A solution to Problem 2.0.2 for $\mathcal{A}$ is called indecomposable if it cannot be represented as a union of solutions for $\mathcal{A}_{1}, \mathcal{A}_{2} \subset \mathcal{A}$. We construct several self-similar indecomposable solutions, in particular for $\mathcal{A}$ having a positive Hausdorff dimension.

### 2.1 Basics

Throughout this chapter $n \geq 4, d \geq 2$ are natural numbers. All the solutions of Problem 2.0.1 for $n \leq 3$ are known in the explicit form since 17-th century. The finite Euclidean Steiner problem has an intricate history, which is studied in the paper [10]. Brazil, Graham, Thomas and Zachariasen done a detailed research and discovered that a statement and basic results were rediscovered (at least) three times. Up to a modern knowledge, it was first stated by Gergonne in 1811, then by Gauss in 1836. The first known publication is dated from 1934 and belongs to Jarník and Kössler [66]. The
problem has become well-known as "Steiner problem" after the great success of the book "What is Mathematics?" by Courant and Robbins [29].

A solution to Problem 2.0 .1 is called Steiner tree. It is known that such a set $\mathcal{S} t$ always exists (but is not necessarily unique, see Fig. 2.1) and that every such $\mathcal{S} t$ is a finite union of segments. Thus, $\mathcal{S} t$ can be represented as a graph, embedded into the Euclidean space, such that its set of vertices contains $P$ and all its edges are straight line segments. This graph is connected and does not contain cycles, i.e. is a tree, which explains the naming of $\mathcal{S}$. It is known that the maximal degree of the vertices of $\mathcal{S} t$ is at most 3 . Moreover, only vertices $x_{i}$ can have degree 1 or 2 , all the other vertices have degree 3 and are called Steiner points while the vertices $x_{i}$ are called terminals. Vertices of degree 3 are called branching points. The angle between any two adjacent edges of $\mathcal{S} t$ is at least $2 \pi / 3$. That means that for a branching point the angle between any two segments incident to it is exactly $2 \pi / 3$, and these three segments belong to the same 2-dimensional plane.

The number of Steiner points in $\mathcal{S} t$ does not exceed $n-2$. A Steiner tree with exactly $2 n-2$ vertices is called full. Every terminal point of a full Steiner tree has degree one.

For a given finite set $P \subset \mathbb{R}^{d}$ consider a connected acyclic set $S$ containing $P$. Then $S$ is called a locally minimal tree if $\overline{S \cap B_{\varepsilon}(x)}$ is a Steiner tree for $(\{x\} \cap P) \cup\left(S \cap \partial B_{\varepsilon}(x)\right)$ for every point $x \in S$ and small enough $\varepsilon>0$. Clearly every Steiner tree is locally minimal and not vice versa. Locally minimal trees have all the mentioned properties of Steiner trees except the minimal length condition. So locally minimal trees inherit the definitions of terminals, Steiner points, branching points and fullness. Proof of the listed properties of Steiner and locally minimal trees together with an additional information on them can be found in the book [63] and in article [55].

Garey, Graham and Johnson [53] proved that the Steiner problem is NP-hard, then Rubinstein, Thomas and Wormald [112] proved that the hardness property remains even in the case of terminals, belonging to a pair of parallel lines as well as in the case of terminals on the sides of the angle which smaller than $2 \pi / 3$.

Similar problems could also be considered in abstract metric spaces. In the most general form the problem would be to connect a set (not necessarily finite or countable) of subsets of an arbitrary metric space in a minimal way with respect to the metric [99], see Section 2.5.1.

Problem 2.0.1 may have several solutions starting with $n=4$ (see Fig. 2.1). Theorem 2.1.1implies the uniqueness of a solution for a general input.

Let us proceed to basic properties of Problem 2.0.2. A general setting for the problem was given in [99]: the ambient space $X$ can be any connected complete metric space with the Heine-Borel property (closed bounded sets are compact) and the given set of points can be any compact subset of the ambient space. In this setting there always exists a set $\mathcal{S}$ t, with minimal 1-dimensional Hausdorff measure $\mathcal{H}^{1}$, such that $\mathcal{S} t \cup \mathcal{A}$ is connected.

As shown in [99] every solution $\mathcal{S} t$ having a finite length has the following properties:

- $\mathcal{S} t \cup \mathcal{A}$ is compact,
- $\mathcal{S} t \backslash \mathcal{A}$ has at most countably many connected components, each of which has positive length,
- $\overline{\mathcal{S} t}$ contains no loops (homeomorphic images of the circle $\mathbb{S}^{1}$ ),
- the closure of every connected component of $\mathcal{S} t$ is a topological tree (a connected, locally connected compact set without loops) with endpoints on $\mathcal{A}$ (so that in particular it has at most a countable number of branching points), with at most one endpoint on each connected component of $\mathcal{A}$ and all the branching points having finite order (i.e. finite number of branches leaving them),
- if $\mathcal{A}$ has a finite number of connected components, then $\mathcal{S} t \backslash \mathcal{A}$ also has finitely many connected components, the closure of each of which is a finite geodesic embedded graph with endpoints on $\mathcal{A}$, and with at most one endpoint on each connected component of $\mathcal{A}$,
- the set $\mathcal{S t} \backslash \mathcal{A}$ is a locally finite geodesic embedded minimal graph.

We also call a solution to Problem 2.0.2 Steiner tree; the above properties explain such naming in the case of $\mathcal{A}$ being a totally disconnected set. It has similar properties with a solution of Problem 2.0.1 and we specify it below. Denote by $\mathbb{M}(\mathcal{A})$ the set of Steiner trees for $\mathcal{A}$. The points of $\mathcal{A}$ touched by the Steiner tree $\mathcal{S} t$ will be called terminals. All points of $\mathcal{A}$ are terminal points if $\mathcal{A}$ is totally disconnected. From now we focus on $X=\mathbb{R}^{2}$ (but the claims of this paragraph also hold for $\mathbb{R}^{d}$ ). A combination of the last enlisted property from [99] with well-known facts on Euclidean Steiner trees (see [55, 63]) gives the following properties. The edges of the locally finite graph $\mathcal{S t} \backslash \mathcal{A}$ are straight line segments. The maximal degree (in graph-theoretic sense) of a vertex is at most 3. Moreover, only terminals can have degree 1 or 2 , all the other vertices have degree 3 and are called Steiner points. Vertices of the degree 3 are called branching points. The angle between any two adjacent edges of a Steiner tree is at least $2 \pi / 3$, in particular a Steiner tree in a neighbourhood of a branching point is a regular tripod: all three angles are equal to $2 \pi / 3$.

In 1980-s and 1990-s explicit solutions of the Steiner problem attracted the attention of several notable mathematicians, in particular Graham. It is worth noting that Du, Hwang and Weng [33] completely solved the Steiner problem when $\mathcal{A}$ is the set of vertices of a regular polygon. Rubinstein and Thomas [111] generalizes the result for when the points of $\mathcal{A}$ are uniformly enough distributed on a circle.

Let us also mention that Chung and Graham [25] and Burkard, Dudás and Maier [11] determined the set of Steiner trees for ladders. A ladder is a collection of $2 n$ lattice points of $\mathbb{Z}^{2}$ which forms a rectangle $1 \times(n-1)$; the structure of a solution is absolutely different for odd and even $n$.

More recently, following the paper [99, the question of finding examples of non trivial infinite Steiner trees was raised. Of course it is easy to find infinite trees by merging together an infinite number of finite trees. Much more difficult is to find an infinite Steiner tree which is indecomposable.

The first example of an infinite, indecomposable Steiner tree was given in [101].
Let us denote by $\mathbb{P}_{d}:=\left(\mathbb{R}^{d}\right)^{n} \backslash$ diag the space of labelled $n$-point configurations $x_{1}, \ldots, x_{n} \in \mathbb{R}^{d}$ of distinct points in the Euclidean space, where diag is the union of $(d n-d)$-dimensional subspaces $x_{i}=x_{j}, i \neq j$. Note that every point of $\mathbb{P}_{d}$ corresponds to some labelled non-degenerate configuration; so let us consider $\mathbb{P}_{d}$ as a configuration space.

A configuration $P \in \mathbb{P}_{d}$ is ambiguous if there are several Steiner trees for $P$. Ivanov and Tuzhilin proved [65] that the complement to the set of ambiguous configurations contains an open dense subset of $\mathbb{P}_{2}$. Edelsbrunner and Strelkova [36] asked whether the measure of ambiguous configurations is zero or not. We provide a positive answer by proving the following stronger statement.

Theorem 2.1.1 (Basok-Cherkashin-Rastegaev-Teplitskaya [5]). Assume that $n \geq 4$. Then the set of planar ambiguous configurations in $\mathbb{P}_{2}$ has Hausdorff dimension $2 n-1$.

### 2.1.1 Topology and embedding class of a tree

For the sake of convenience and completeness, we would like to begin our discussion with a careful introduction of the concept of "topology" of a tree often used in the context of the Steiner problem. As it is usually done in the literature (see, for instance [55, 63]), we define the topology of a Steiner tree


Figure 2.1: An example of non-unique solution. Terminals form a square
to be the corresponding abstract topological graph with labelled terminals and unlabelled Steiner points. Thus, a topology $T$ is a topological space with a tree structure, and some vertices of $T$, including all its leaves and vertices of degree 2, are labelled. Moreover, we assume that all vertices of $T$ have degrees at most 3 , as it naturally holds for any Steiner tree.

Note that two trees embedded in a different (non-homotopic) way into the plane may have the same topology (on the other hand, for $d \geq 3$ homotopic equivalence means exactly that the topologies coincide). To distinguish non-homotopic embeddings Edelsbrunner and Strelkova [35] introduced another invariant way to describe the topological type of the tree which we call the "embedding class". Below we introduce several ways to define the embedding class of a tree commonly used. We include the proof of their equivalence for the sake of completeness.

Let $T$ be a combinatorial tree, id est $T=(V, E)$ as a graph. Let $\vec{E}(T)$ denote the set of oriented edges of $T$ (in particular, $|\vec{E}(T)|=2|E(T)|)$. Given an edge $\vec{e} \in \vec{E}$ denote by $o(\vec{e})$ the origin and by $t(\vec{e})$ the tail. Let us say that a bijection $\sigma: \vec{E}(T) \rightarrow \vec{E}(T)$ determines a cyclic order around each vertex of $T$ if $o(\sigma(\vec{e}))=o(\vec{e})$ for any $\vec{e}$ and for any pair $(v, \vec{e})$ such that $o(\vec{e})=v$ the set $\vec{e}, \sigma(\vec{e}), \sigma^{2}(\vec{e}), \ldots$ is exactly the set of oriented edges emanating from $v$.

The following classical lemma defines the embedding class:
Lemma 2.1.1. Let a positive integer $n$ be fixed. The following three sets are in natural bijection:

1. The set $P M_{1}$ of pairs $(T, \sigma)$, where $T$ is a combinatorial tree with all vertices of degree at most 3, with $n$ labelled vertices, including all leaves of $T$ and vertices of degree 2, and $\sigma: \vec{E}(T) \rightarrow \vec{E}(T)$ is a bijection determining a cyclic order around each vertex.
2. The set $P M_{2}$ of pairs $(T,[f])$, where $T$ is a combinatorial tree with all vertices of degree at most 3, with $n$ labelled vertices, including all leaves of $T$ and vertices of degree 2, $f$ is a bijection between $\vec{E}(T)$ and the set $\partial D$ of edges of the regular $|\vec{E}(T)|$-gon $D$ (0-gon is assumed to be empty) oriented clockwise such that if $o\left(f\left(\vec{e}_{1}\right)\right)=t\left(f\left(\vec{e}_{2}\right)\right)$, then $o\left(\vec{e}_{1}\right)=t\left(\vec{e}_{2}\right)$, and $[f]$ is the equivalence class of $f$ with respect to the cyclic shift on $\partial D$.
3. The set $P M_{3}$ of pairs $(T,[\iota])$, where $T$ is a topology with $n$ labelled vertices, ८ is some embedding of $T$ into the plane and [८] is the homotopy class of $\iota$ in the space of embeddings.

Remark 2.1.1. Each of these three sets can be considered as the set of plane maps or ribbon graphs with tree-like skeletons and some labelled vertices (see [81]).

Proof. We construct a map $F_{i}: P M_{i} \rightarrow P M_{i+1}$ for $i=1,2,3($ where we compute indices mod 3$)$, and then show that the composition of these maps is identity.

Let $(T, \sigma) \in P M_{1}$ be given, let $N=|\vec{E}(T)|$. If $N=0$, then we just take an empty $D$. Assume that $N>0$. Define $\alpha: \vec{E}(T) \rightarrow \vec{E}(T)$ to be the involution reversing the orientation and set $\varphi=\alpha \circ \sigma$. It is easy to verify that $\varphi$ is a cyclic permutation of $\vec{E}(T)$. Set $f: \vec{E}(T) \rightarrow \partial D$ to be any bijection which respects the cyclic order imposed by $\varphi$. It is easy to check that ( $T,[f]$ ) belongs to $P M_{2}$, thus we get the map $F_{1}$.

Let $(T,[f]) \in P M_{2}$ be given. If $T$ is one point, then we define $\iota$ arbitrary. Assume that $T$ has at least two vertices, let $N=|\vec{E}(T)|$ and $D$ be the regular $N$-gon. As above, let $\alpha: \vec{E}(T) \rightarrow \vec{E}(T)$ to be the involution reversing the orientation. Now, glue each edge $\vec{e} \in \partial D$ with $f \circ \alpha \circ f^{-1}(\vec{e})$ in opposite direction. It is straightforward to see that in this way we get an oriented surface $S$ out of $D$, and a natural embedding $\iota$ of $T$ into $S$. Computing the Euler characteristic we find out that $S$ is a sphere and hence $\iota$ corresponds to a planar embedding of $T$. Set $F_{2}(T,[f])=(T,[\iota])$ (where the topology on $T$ comes from $\iota$ naturally).

Let $(T,[\iota]) \in P M_{3}$ be given. If $T$ is one point, then we take $\sigma$ to be the only map between empty sets. Assume that $T$ has at least two vertices. Let $v$ be a vertex and $\vec{E}_{v}(T)$ be the set of oriented edges emanating from $v$. Then given $\vec{e} \in \vec{E}_{v}(T)$ define $\sigma(\vec{e})$ to be first edge in $\vec{E}_{v}(T)$ coming after $\vec{e}$ when going around $\iota(v)$ in counterclockwise direction. Set $F_{3}(T,[\iota])=(T, \sigma)$.

The fact that $F_{3} \circ F_{2} \circ F_{1}=\mathrm{id}$ is a simple exercise which we leave to the reader. Note that given a labelled tree $T$ the amount of all possible $\sigma$ such that $(T, \sigma) \in P M_{1}$, or $f$ such that $(T,[f]) \in P M_{2}$ is finite; this shows that $F_{1}$ is a bijection. On the other hand, the fact that the number of homotopy classes of embeddings $\iota$ for a given topology is finite is not obvious. Hence, at the moment we have only a right inverse for $F_{3}$. Let us sketch the construction the inverse map for $F_{2}$ to overcome this difficulty. Choose an embedding $\iota: T \rightarrow \mathbb{C}$ and consider the simply-connected surface $\widehat{\mathbb{C}} \backslash \iota(T)$, where $\widehat{\mathbb{C}}$ is the Riemann sphere. Let $\mathbb{D}$ be the unit disc and $\psi: \mathbb{D} \rightarrow \widehat{\mathbb{C}} \backslash \iota(T)$ be the uniformization map. One can show that $\psi$ extends to the boundary of $\mathbb{D}$ in a unique way such that $\psi: \overline{\mathbb{D}} \rightarrow \widehat{\mathbb{C}}$ is continuous. Moreover, each point of $\iota(T)$ corresponds to several prime ends of the domain $\widehat{\mathbb{C}} \backslash \iota(T)$ (see [104, Chapter 2]); there are two prime ends for each inner point of an edge, and $\operatorname{deg} v$ prime ends for each vertex $v$. Let $v_{1}, \ldots, v_{N}$ be all the preimages of vertices of $T$ on $\partial \mathbb{D}$, the count of the prime ends implies that $N=|\vec{E}(T)|$. Then $\overline{\mathbb{D}}$ together with these points has the combinatorics of the regular $N$-gon, whence we get the morphism $f$ such that $(T, f) \in P M_{2}$. The fact that this construction inverses $F_{2}$ is straightforward.

Note that the regular polygon $D$ from the set $P M_{2}$ naturally corresponds to the outer face of the planar graph $\iota(T)$ for $\iota$ coming from $P M_{3}$. Using Lemma 2.1.1, we identify $P M_{1}, P M_{2}$ and $P M_{3}$, so that given, say $(T, \sigma) \in P M_{1}$ we will always assume that we are also given the corresponding $(T,[f]) \in P M_{2}$ and $(T,[\iota]) \in P M_{3}$ and will use the corresponding notation if it does not lead to a confusion.

Now, we introduce another (fourth) way to encode embeddings of a topological tree, which was originally used by Edelsbrunner and Strelkova.

Let $(T, \sigma) \in P M_{1}$ and $\varphi=\alpha \circ \sigma$, where $\alpha: \vec{E}(T) \rightarrow \vec{E}(T)$ is the involution reversing the orientation. Assume that $T$ has $n$ labels. Let $A=\{1,2,3, \ldots, n, b\}$ be the alphabet on $n+1$ letters,
$n$ of them are numbers from 1 to $n$, and $(n+1)$-th is the special letter $b$. Let

$$
\mathbf{C}=\left(\bigcup_{k \geq 0} A^{k}\right) / \text { cyclic shift }
$$

be the set of all words build from this alphabet considered up to the cyclic shift. Let $\vec{e} \in \vec{E}(T)$ be arbitrary and $N=|\vec{E}(T)|$. Then, given $(T, \sigma)$, define the word $C(T, \sigma) \in \mathbf{C}$ by the following rule: fix a vector $\vec{e}_{0} \in \vec{E}(T)$ and set $\vec{e}_{i}=\varphi\left(\vec{e}_{i-1}\right)$, then define $C(T, \sigma)=a_{0} a_{1} \ldots a_{N-1}$, where $a_{i}$ is the label of $o\left(\varphi\left(\vec{e}_{i}\right)\right)$ if $o\left(\varphi\left(\vec{e}_{i}\right)\right)$ is labelled, and $a_{i}=b$ else; if $T$ consists of one vertex, then the word $C(T, \sigma)$ is the empty word.

Let $(T,[f]) \in P M_{2}$ correspond to $(T, \sigma)$ and $D$ be the regular $N$-gon. As $D$ can be seen as the outer face of the planar graph $\iota(T)$, there is a many-to-one correspondence between the vertices of $D$ and the vertices of $T$. Then the word $C(T, \sigma)$ is nothing but the list of vertices obtained by going along the boundary of $D$; each time we met a vertex walking along $\partial D$, we add its label to $C(T, \sigma)$, or the letter $b$ if the vertex does not have a label. For example, we have $C=1 b b 2 b 3 b b 4 b$ and $C=1 b 2 b b 3 b 4 b b$ (and we could also write $C=b 2 b b 3 b 4 b b 1$ in the latter case as we factorized by a cyclic shift) for the left and the right trees on the Fig. 2.1 respectively.

Lemma 2.1.2. The morphism $(T, \sigma) \mapsto C(T, \sigma)$ is injective from the set of pairs $(T, \sigma)$ to $\mathbf{C}$.
Proof of Lemma 2.1.2. Let $(T, \sigma)$ be given and $N=|\vec{E}(T)|$. Then the length of the word $C(T, \sigma)$ is $N$, hence $C$ distinguishes pairs $(T, \sigma)$ with different cardinality $N$ of the set of edges of the tree. We will show that $C$ distinguishes different pairs with the same $N$ by induction. If $N=0,2,4$, then there is nothing to prove, assume that $N>4$ and $W=C(T, \sigma)$. We need to show that if $W=C\left(T_{1}, \sigma_{1}\right)$, then $(T, \sigma)=\left(T_{1}, \sigma_{1}\right)$. Define

$$
\begin{array}{l|l}
I_{1}=\{i & \left.a_{i} \text { occurs } 1 \text { time in } W\right\} \\
I_{2}=\{i & \left.a_{i} \text { occurs } 2 \text { times in } W\right\}
\end{array}
$$

We clearly have a bijection between the labels $\left\{a_{i} \mid i \in I_{1}\right\}$ and $\left\{a_{i} \mid i \in I_{2}\right\}$ and the vertices of degree 1 and 2 in $T$ respectively, and the same for $T_{1}$. Assume that we can find $i \in I_{1}$ such that $i+1 \in I_{2}$ (here $N+1=1$ ). Then consider the word $W^{\prime}$ obtained from $W$ by removing $a_{i}$ and $a_{i+1}$. Then $W^{\prime}=C\left(T^{\prime}, \sigma^{\prime}\right)$, where $T^{\prime}$ is obtained from $T$ by removing the edge $a_{i} a_{i+1}$ and keeping all labels, and $\sigma^{\prime}$ is computed from $\sigma$ in the natural way (note that $T$ has at least one edge since we assume that $N>4)$. In the same time, $W^{\prime}=C\left(T_{1}^{\prime}, \sigma_{1}^{\prime}\right)$, where $\left(T_{1}^{\prime}, \sigma_{1}^{\prime}\right)$ is obtained from $\left(T_{1}, \sigma_{1}\right)$ in the same procedure. By the induction hypothesis $\left(T^{\prime}, \sigma^{\prime}\right)=\left(T_{1}^{\prime}, \sigma_{1}^{\prime}\right)$. From here, it is easy to see that $(T, \sigma)=\left(T_{1}, \sigma_{1}\right)$.

Assume now that for any $i \in I_{1}$ we have $i+1 \notin I_{2}$. It follows that one can find $i \in I_{1}$ such that $i+2 \in I_{1}$ also. Consider the word $W^{\prime}$ obtained from $W$ by removing $a_{i-1}, a_{i}, a_{i+1}$ and $a_{i+3}$. This word corresponds to $C\left(T^{\prime}, \sigma^{\prime}\right)$, where $T^{\prime}$ is obtained from $T$ by removing two vertices $a_{i}$ and $a_{i+2}$ and labelling their common parent by $a_{i+2}$ (note that their parent must have degree 3). Note that $T^{\prime}$ has at least one edge since $N>4$. Doing the same with $T_{1}$ we again get two pairs ( $T^{\prime}, \sigma^{\prime}$ ) and ( $T_{1}^{\prime}, \sigma_{1}^{\prime}$ ) such that $W^{\prime}=C\left(T^{\prime}, \sigma^{\prime}\right)=C\left(T_{1}^{\prime}, \sigma_{1}^{\prime}\right)$, which implies that $\left(T^{\prime}, \sigma^{\prime}\right)=\left(T_{1}^{\prime}, \sigma_{1}^{\prime}\right)$ by the induction and, eventually, $(T, \sigma)=\left(T^{\prime}, \sigma^{\prime}\right)$.

Starting from now we will call the word $C(T, \sigma)$ an embedding class. Using Lemma 2.1.1 and Lemma 2.1.2 we will feel free to identify the embedding class with the homotopy class of embeddings defined in several ways presented in aforementioned lemmas.

### 2.1.2 Connectedness in $\mathbb{P}_{d}$

Let us return to our analysis of Steiner trees. We say that a topology $T$ of a tree $S$ is full if the corresponding tree is full. Further, let us call a topology $T$ realizable for a configuration $P \in \mathbb{P}_{d}$ if there exists such a locally minimal tree $S(P)$ with topology $T$; we will denote this tree by $S_{T}(P)$.

Proposition 2.1.1 (Melzak, [91]). If a topology $T$ is realizable for $P \in \mathbb{P}_{2}$ then the realization $S_{T}(P)$ is unique.

Proposition 2.1.1 shows that $S_{T}(P)$ is uniquely defined. Moreover one can construct (or show that it is impossible) $S_{T}(P)$ in a linear time $O(n)$, see [62]. However, we rarely know a priori, which topology gives a Steiner tree. Although the number of possible topologies for an $n$ points configuration is finite, checking all of them may consume a lot of time, since this number of topologies grows very fast with $n$, see [55, 63]. Indeed, the Steiner tree problem is NP-complete [53].

We need the following generalization of Proposition 2.1.1. For a full topology $T$ define $D(T)$ as the set of topologies that can be obtained from $T$ by shrinking some edges connecting a terminal with a Steiner point (these edges should have pairwise different ends).

Proposition 2.1.2 (Gilbert-Pollak [55], Hwang-Weng [64]). Let $T$ be a full topology and $P \in \mathbb{P}_{2}$. Consider the function $L\left(y_{1}, \ldots, y_{n-2}\right):\left(\mathbb{R}^{2}\right)^{n-2} \rightarrow \mathbb{R}$ which is the length of a tree on the vertex set $P \cup\left\{y_{1}, \ldots, y_{n-2}\right\}$ with straight edges and topology $T$ (we allow $y_{i}$ coincide with terminals). Then $L$ has a unique local minimum and so there is at most one realization with a topology from $D(T)$.

A generic topology is a topology without terminals of degree 3 .
Observation 2.1.1. (i) Every generic topology $R$ belongs to exactly one set $D(T)$, because the reverse procedure (replacing every vertex $A$ of degree 2 in $R$ on a Steiner point $b$ and add edge bA) leads to a full topology $T$.
(ii) Suppose that $\mathcal{S}$ t is the unique Steiner tree for some $P \in \mathbb{P}_{2}$ and has a generic topology $R \in D(T)$ for some full topology $T$. Then for some positive $\eta>0$ and any other full topology $T^{\prime}$ the length of the realization from $D\left(T^{\prime}\right)$ exceeds $\mathcal{H}^{1}(\mathcal{S} t)$ by at least $\eta$. If one perturbs every terminal by at most $\eta /(2 n)$, then by triangle inequality a perturbed configuration $P^{\prime}$ has a unique Steiner tree $\mathcal{S t}\left(P^{\prime}\right)$ and the topology of $\mathcal{S} t\left(P^{\prime}\right)$ belongs to $D(T)$.
(iii) Configurations $P \in \mathbb{P}_{2}$ for which there exists a locally minimal tree with non-generic topology have the Hausdorff dimension $2 n-2$.

In this section we study the way realizations and minimal realizations of different embedding classes divide the configuration space.

A similar research topic appears in [36, 35], where the connectedness of some sets related to an embedding class $E C$, is studied. Let $\Omega(E C)$ be the subset of $\mathbb{P}_{d}$ consisting of all $P \in \mathbb{P}_{d}$ for which $E C$ is realizable. Note that for every embedding class $E C$ the set $\Omega(E C)$ is path-connected.

Theorem 2.1.2 (Edelsbrunner-Strelkova, [36, 35]). Let EC be an embedding class. Then the subset of $\mathbb{P}_{d}$ for which the Steiner tree is unique and has the embedding class EC is path-connected.

In the planar case they also obtained the following result.
Theorem 2.1.3 (Edelsbrunner-Strelkova, [36, 35]). Let EC be a full embedding class. Then the subset of $\mathbb{P}_{2}$ for which the Steiner tree has the embedding class EC is path-connected.

The second result of this chapter is the following.
Theorem 2.1.4 (Basok-Cherkashin-Rastegaev-Teplitskaya [5]). The subset of $\mathbb{P}_{d}$ for which there is a unique Steiner tree is path-connected.

The proof of Theorem 2.1.4 is constructive (modulo Theorem 2.1.2) and planar (again modulo Theorem 2.1.2). Moreover the embedding class of the Steiner tree is known at every point of a constructed path.

### 2.2 A universal Steiner tree

In this subsection we provide a construction of a unique Steiner tree with an infinite number of Steiner points. It appeared in [101] and then was simplified and improved in [20, 100].

Let $S_{\infty}$ be an infinite tree with vertices $y_{0}, y_{1}, y_{2}, \ldots$ and edges given by $y_{0} y_{1}$ and $y_{k} y_{2 k}, y_{k} y_{2 k+1}$, $k \geq 1$. Thus, $S_{\infty}$ is an infinite binary tree with an additional vertex $y_{0}$ attached to the common parent $y_{1}$ of all other vertices $y_{k}, k \geq 2$. The goal of the mentioned papers is to embed $S_{\infty}$ in the plane in such a way that the image of each finite subtree of $S_{\infty}$ will be the unique Steiner tree for the set of its vertices having degree 1 or 2 . We define the embedding below by specifying the positions of $y_{0}, y_{1}, y_{2}, \ldots$ on the plane.


Figure 2.2: The first three levels in the construction of $\Sigma_{\infty}$. The set $\Sigma_{3}$ is thick blue.
Let $\Lambda=\left\{\lambda_{i}\right\}_{i=0}^{\infty}$ be a sequence of positive real numbers. Define an embedding $\Sigma(\Lambda)$ of $S_{\infty}$ as a rooted binary tree with the root $y_{0}=(0,0)$ the first descendant $y_{1}=(1,0)$ and the ratio between edges of $(i+1)$-th and $i$-th levels being $\lambda_{i}$. For small enough $\left\{\lambda_{i}\right\}$ the set $\Sigma(\Lambda)$ is a tree, see Fig. 2.2.

Let $A_{\infty}(\Lambda)$ be the union of the set of all leaves (limit points) of $\Sigma(\Lambda)$ and $\left\{y_{0}\right\}$.

Theorem 2.2.1 (Cherkashin-Teplitskaya, [20]). A binary tree $\Sigma(\Lambda)$ is a Steiner tree for $A_{\infty}(\Lambda)$ provided that $\lambda_{i}=\lambda<1 / 300$.

Very recently Theorem 2.2.1 was significantly improved.
Theorem 2.2.2 (Paolini-Stepanov, [100]). A binary tree $\Sigma(\Lambda)$ is a unique Steiner tree for $A_{\infty}(\Lambda)$ provided that $\lambda_{i}=\lambda<1 / 25$.
Proof of Theorem 2.2.1. Let the values of $\lambda<\frac{1}{300}, \varepsilon=\frac{\lambda^{2}}{1-\lambda}$ be fixed during the proof. The following auxiliary constructions are drawn in Figure 2.3. Let $Y_{1} B_{1} C_{1}$ be an isosceles triangle with the Fermat-Torricelli point $T_{1}$, such that $\left|Y_{1} T_{1}\right|=1,\left|T_{1} B_{1}\right|=\left|T_{1} C_{1}\right|=\lambda$; then, by the cosine rule $\left|Y_{1} B_{1}\right|=\left|Y_{1} C_{1}\right|=\sqrt{1+\lambda+\lambda^{2}}$ and $\left|B_{1} C_{1}\right|=\sqrt{3} \lambda$. Analogously, let $Y_{2} B_{2} C_{2}$ be an isosceles triangle with the Fermat-Torricelli point $T_{2}$, such that $\left|Y_{2} T_{2}\right|=1 / 4,\left|T_{2} B_{2}\right|=\left|T_{2} C_{2}\right|=\lambda$; then $\left|Y_{2} B_{2}\right|=\left|Y_{2} C_{2}\right|=\sqrt{1 / 16+\lambda / 4+\lambda^{2}}$ and $\left|B_{2} C_{2}\right|=\sqrt{3} \lambda$.


Figure 2.3: The construction of triangles in lemmas.
Let $b_{i} \subset B_{\varepsilon}\left(B_{i}\right)$ and $c_{i} \subset B_{\varepsilon}\left(C_{i}\right)$ be symmetric sets with respect to the axis of symmetry $l_{i}$ of $Y_{i} B_{i} C_{i}$, where $i=1,2$. Finally, let $Y_{\text {up }}, Y_{\text {down }}$ be such points that $Y_{\text {up }} Y_{\text {down }} \| B_{2} C_{2}, Y_{2} \in\left[Y_{\text {up }} Y_{\text {down }}\right]$ and $\left|Y_{2} Y_{\text {up }}\right|=\left|Y_{2} Y_{\text {down }}\right|=1 / 2$.

The following proposition is more-or-less known (see, for instance, Lemma A. 6 in [101]), but we prove is for the sake of completeness. Recall that a regular tripod is a union of three segments with a common end and pairwise angles equal to $2 \pi / 3$.

Proposition 2.2.1. (i) For every $\mathcal{S} \in \mathbb{M}\left(\left\{Y_{1}\right\} \cup b_{1} \cup c_{1}\right)$, the set $\mathcal{S} \backslash B_{10 \varepsilon}\left(B_{1}\right) \backslash B_{10 \varepsilon}\left(C_{1}\right)$ is a regular tripod.
(ii) Every $\mathcal{S} \in \mathbb{M}\left(\left[Y_{\text {up }} Y_{\text {down }}\right] \cup b_{2} \cup c_{2}\right)$ is a regular tripod outside of $B_{10 \varepsilon}\left(B_{2}\right) \cup B_{10 \varepsilon}\left(C_{2}\right)$.

Proof. In this proof, $i \in\{1,2\}$. Suppose that $\mathcal{S}$ intersects with every circle $\partial B_{\rho}\left(B_{i}\right)$ in at least 2 points for $\varepsilon \leq \rho \leq 10 \varepsilon$ (see Fig. 2.4). Then, we may replace $\mathcal{S}$ with a shorter competitor, as follows. Put $\mathcal{S}_{b}=\mathcal{S} \cap B_{\varepsilon}\left(B_{i}\right)$. By the definition and the co-area inequality,

$$
\mathcal{H}^{1}(\mathcal{S}) \geq \mathcal{H}^{1}\left(\mathcal{S}_{b}\right)+2 \cdot 9 \varepsilon+\mathcal{H}^{1}\left(\mathcal{S}_{i}\right)
$$

where $\mathcal{S}_{1} \in \mathbb{M}\left(\left\{Y_{1}\right\} \cup \partial B_{10 \varepsilon}\left(B_{1}\right) \cup c_{1}\right), \mathcal{S}_{2} \in \mathbb{M}\left(\left[Y_{\text {up }} Y_{\text {down }}\right] \cup \partial B_{10 \varepsilon}\left(B_{2}\right) \cup c_{2}\right)$. Now, take $\mathcal{S}_{i} \cup \mathcal{S}_{b} \cup$ $\partial B_{\varepsilon}\left(B_{i}\right) \cup \mathcal{R}_{B}$, where $\mathcal{R}_{B}$ is the radius connecting $\mathcal{S}_{i}$ with $\partial B_{\varepsilon}\left(B_{i}\right)$. The length of this competitor is

$$
\mathcal{H}^{1}\left(\mathcal{S}_{b}\right)+2 \pi \varepsilon+9 \varepsilon+\mathcal{H}^{1}\left(\mathcal{S}_{i}\right),
$$

which gives a contradiction since $2 \pi<9$. The symmetric construction shows that the situation where $\mathcal{S}$ intersects with every circle $\partial B_{\rho}\left(C_{i}\right)$ in at least 2 points for $\varepsilon \leq \rho \leq 10 \varepsilon$ is also impossible.


Figure 2.4: Picture of the proof of Proposition 1.
Thus, there are $\rho_{b}, \rho_{c} \in[\varepsilon, 10 \varepsilon]$, such that $\mathcal{S} \cap \partial B_{\rho_{b}}\left(B_{i}\right)$ is a point $B_{i}^{\prime}$ and $\mathcal{S} \cap \partial B_{\rho_{c}}\left(C_{i}\right)$ is a point $C_{i}^{\prime}$. Clearly, $\mathcal{S}=\mathcal{S}_{i} \cup \mathcal{S}_{b} \cup \mathcal{S}_{c}$, where $\mathcal{S}_{b}=\mathcal{S} \cap B_{\rho_{b}}\left(B_{i}\right), \mathcal{S}_{c}=\mathcal{S} \cap B_{\rho_{c}}\left(C_{i}\right)$ and $\mathcal{S}_{1} \in \mathbb{M}\left(\left\{Y_{1}\right\} \cup\left\{B_{i}^{\prime}\right\} \cup\left\{C_{i}^{\prime}\right\}\right)$, $\mathcal{S}_{2} \in \mathbb{M}\left(\left[Y_{\text {up }} Y_{\text {down }}\right] \cup\left\{B_{i}^{\prime}\right\} \cup\left\{C_{i}^{\prime}\right\}\right)$. Clearly $\mathcal{S}_{i}$ is a tripod or the union of two segments. We claim that $\mathcal{S}_{i}$ is a tripod. By the triangle inequality:

$$
\begin{equation*}
\left|\mathcal{H}^{1}\left(\left[T_{i} B_{i}\right] \cup\left[T_{i} C_{i}\right]\right)-\mathcal{H}^{1}\left(\left[T_{i} B_{i}^{\prime}\right] \cup\left[T_{i} C_{i}^{\prime}\right]\right)\right|<20 \varepsilon . \tag{2.1}
\end{equation*}
$$

Now, let us prove item (i). By (2.1), the length of the (non-regular) tripod $\left[T_{1} Y_{1}\right] \cup\left[T_{1} C_{1}^{\prime}\right] \cup\left[T_{1} B_{1}^{\prime}\right]$ connecting $Y_{1}, B_{1}^{\prime}$ and $C_{1}^{\prime}$ is, at most, $1+2 \lambda+20 \varepsilon$. For the same reason, the length of two segments is at least

$$
\sqrt{1+\lambda+\lambda^{2}}+\sqrt{3} \lambda-30 \varepsilon>1+\left(\frac{1}{2}+\sqrt{3}\right) \lambda-30 \varepsilon
$$

Recall that $\varepsilon=\frac{\lambda^{2}}{1-\lambda} ;$ it is straightforward to check that

$$
1+\left(\frac{1}{2}+\sqrt{3}\right) \lambda-30 \varepsilon>1+2 \lambda+20 \varepsilon
$$

for $\lambda<1 / 300$. Thus, we show that $\mathcal{S}_{1}$ contains a tripod connecting $Y_{1}, B_{1}^{\prime}$ and $C_{1}^{\prime}$; by the minimality argument, it is regular.

Let us deal with item (ii). By (2.1), the length of the (non-regular) tripod $\left[T_{2} Y_{2}\right] \cup\left[T_{2} C_{2}^{\prime}\right] \cup\left[T_{2} B_{2}^{\prime}\right]$ connecting $Y_{2}, B_{2}^{\prime}$ and $C_{2}^{\prime}$ is, at most, $1 / 4+2 \lambda+20 \varepsilon$. Again, the two-segment construction has a length of at least

$$
1 / 4+\lambda / 2+\sqrt{3} \lambda-30 \varepsilon
$$

The rest of the calculations coincide with the first item.
Lemma 2.2.1. There exists $\mathcal{S} \in \mathbb{M}\left(\left\{Y_{1}\right\} \cup b_{1} \cup c_{1}\right)$, which is symmetric with respect to $l_{1}$.
Proof. Let $F$ be a point at the ray $\left[Y_{1} T_{1}\right.$ ), such that $\left|Y_{1} F\right|=1+\frac{3}{2} \lambda$ (see Fig. 2.5), and denote by $D E F$ the equilateral triangle, such that $Y_{1}$ is the middle of the segment $[D E]$ and $B_{1} C_{1}$ is parallel to $D E$. Consider segments $\left[Z_{l} Z_{r}\right] \subset[D F]$ and $\left[V_{l} V_{r}\right] \subset[E F]$, such that $\left|Z_{l} Z_{r}\right|=\left[V_{l} V_{r}\right]=\lambda$ and $Z:=D F \cap T_{1} B_{1}, V:=E F \cap T_{1} C_{1}$ are centers of the segments. Note that $l_{1}$ is a symmetry axis of $D E F$, and $\left[Z_{l} Z_{r}\right]$ and $\left[V_{l} V_{r}\right]$ are also symmetric with respect to $l_{1}$.

By Proposition 2.2.1 (i), every minimal set $\mathcal{S}$ is a regular tripod $Y_{1} B_{1}^{\prime} C_{1}^{\prime}$ out of $B_{10 \varepsilon}\left(B_{1}\right) \cup B_{10 \varepsilon}\left(C_{1}\right)$. We claim that the tripod $Y_{1} B_{1}^{\prime} C_{1}^{\prime}$ intersects segments $\left[Z_{l} Z_{r}\right]$ and $\left[V_{l} V_{r}\right]$. Indeed, consider Cartesian coordinates in which $Y_{1}=(0,0), B_{1}=(1+\lambda / 2, \sqrt{3} \lambda / 2)$ and $C_{1}=(1+\lambda / 2,-\sqrt{3} \lambda / 2)$. Then $Z=(1+$ $3 \lambda / 8,3 \sqrt{3} \lambda / 8), Z_{l}=(1+3 \lambda / 8-\sqrt{3} \lambda / 4,3 \sqrt{3} \lambda / 8+\lambda / 4)$, and $Z_{r}=(1+3 \lambda / 8+\sqrt{3} \lambda / 4,3 \sqrt{3} \lambda / 8-\lambda / 4)$. Since the center $T_{1}^{\prime}$ of $Y_{1} B_{1}^{\prime} C_{1}^{\prime}$ lies inside triangle $Y_{1} B_{1}^{\prime} C_{1}^{\prime}$, it has an $x$-coordinate smaller than the $x$-coordinate of $B_{1}^{\prime}$ and a $y$-coordinate smaller than the $y$-coordinate of $B_{1}^{\prime}$.

We consider the following auxiliary data for the Steiner problem: $A_{\text {mid }}=\left[Z_{l} Z_{r}\right] \cup\left[V_{l} V_{r}\right] \cup\left\{Y_{1}\right\}$, $A_{u p}=\left[Z_{l} Z_{r}\right] \cup b_{1}, A_{\text {down }}=\left[V_{l} V_{r}\right] \cup c_{1}$. By the results from [99], as mentioned in the introduction, every $\mathbb{M}\left(A_{i}\right)$ is not empty. Segments $\left[Z_{l} Z_{r}\right]$ and $\left[V_{l} V_{r}\right]$ split every $\mathcal{S} \in \mathbb{M}\left(\left\{Y_{1}\right\} \cup b_{1} \cup c_{1}\right)$ into three parts, connecting $A_{\text {mid }}, A_{\text {up }}$, and $A_{\text {down }}$, so

$$
\begin{equation*}
\mathcal{H}^{1}(\mathcal{S}) \geq \mathcal{H}^{1}\left(\mathcal{S}_{\text {mid }}\right)+\mathcal{H}^{1}\left(\mathcal{S}_{\text {up }}\right)+\mathcal{H}^{1}\left(\mathcal{S}_{\text {down }}\right) \tag{2.2}
\end{equation*}
$$

where $\mathcal{S}_{i} \in \mathbb{M}\left(A_{i}\right)$. We claim that the equality in (2.2) holds.
It is known (see the barycentric coordinate system) that the sum of distances from a point inside a closed equilateral triangle to the sides does not depend on a point. Thus, $\mathbb{M}\left(A_{\text {mid }}\right)$ is a set of regular tripods, and each tripod is symmetric with respect to $l_{1}$. Moreover, for every point $x \in\left[Z_{l} Z_{r}\right]$, there is a unique regular tripod $\mathcal{S}_{x} \in \mathbb{M}\left(A_{\text {mid }}\right)$, and $\mathcal{S}_{x}$ is orthogonal to $\left[Z_{l} Z_{r}\right]$ at $x$.

Now consider any $\mathcal{S}_{\text {down }} \in \mathbb{M}\left(A_{\text {down }}\right)$. Let $\mathcal{S}_{\text {up }}$ be a set that is symmetric to $\mathcal{S}_{\text {down }}$ with respect to $l_{1}$; clearly, $\mathcal{S}_{u p} \in \mathbb{M}\left(A_{\text {up }}\right)$. For $x \in \mathcal{S}_{\text {down }} \cap\left[V_{l} V_{r}\right]$, the set $\mathcal{S}_{x} \cup \mathcal{S}_{u p} \cup \mathcal{S}_{\text {down }}$ connects $\left\{Y_{1}\right\} \cup b_{1} \cup c_{1}$, and reaches the equality in (2.2), so $\mathcal{S}_{x} \cup \mathcal{S}_{u p} \cup \mathcal{S}_{\text {down }}$ is a Steiner tree for $\left\{Y_{1}\right\} \cup b_{1} \cup c_{1}$. By the construction, it is symmetric with respect to $l_{1}$.


Figure 2.5: Picture of the proof of Lemma 1.
Lemma 2.2.2. There exists $\mathcal{S} \in \mathbb{M}\left(\left[Y_{\text {up }} Y_{\text {down }}\right] \cup b_{2} \cup c_{2}\right)$, which is symmetric with respect to $l_{2}$.
Proof. By Proposition 2.2.1(ii), every Steiner tree $\mathcal{S}$ coincides with a regular tripod outside of $B_{10 \varepsilon}\left(B_{2}\right) \cup B_{10 \varepsilon}\left(C_{2}\right)$. Clearly, its longest segment is perpendicular to $Y_{u p} Y_{\text {down }}$ (see Fig. 2.6). We
want to show that it touches $Y_{\text {up }} Y_{\text {down }}$ in $Y_{2}$, i.e., one of the three segments is a subset of $l_{2}$. Assuming the contrary, suppose that $l_{2} \cap \mathcal{S}$ is a point, denote it by $L$, and let $n \| B_{2} C_{2}$ be the line containing $L$. Then $n$ divides $\mathcal{S}$ into three connected components; denote them by $\mathcal{S}_{Y}, \mathcal{S}_{b}$, and $\mathcal{S}_{c}$, respectively.


Figure 2.6: Picture of the proof of Lemma 2.
Without the loss of generality, $L$ belongs to $\mathcal{S}_{c}$.
Let us construct competitors $\mathcal{S}_{1}$ and $\mathcal{S}_{2}$, connecting $\left[Y_{\text {up }} Y_{\text {down }}\right]$, $b_{2}$, and $c_{2}$. Let $\mathcal{S}_{1}=\left[Y_{2} L\right] \cup \mathcal{S}_{c} \cup \mathcal{S}_{c}^{\prime}$, where $\mathcal{S}_{c}^{\prime}$ is a reflection of $\mathcal{S}_{c}$ with respect to $l_{2}$. Put $h:=\operatorname{dist}\left(Y_{2}, \mathcal{S}_{Y} \cap\left[Y_{u p} Y_{\text {down }}\right]\right)$. Thus

$$
\mathcal{H}^{1}\left(\mathcal{S}_{1}\right)=\mathcal{H}^{1}\left(\mathcal{S}_{Y}\right)-\sqrt{3} h+2 \mathcal{H}^{1}\left(\mathcal{S}_{c}\right) .
$$

Let $\mathcal{S}_{2}:=\mathcal{T} \cup \mathcal{S}_{b} \cup \mathcal{S}_{b}^{\prime}$, where $\mathcal{S}_{b}^{\prime}$ is a reflection of $\mathcal{S}_{b}$, with respect to $l_{2}$, and $\mathcal{T}$ is a regular tripod connecting $Y_{2}$ with $n \cap \mathcal{S}_{b}$ and $n \cap \mathcal{S}_{b^{\prime}}$. Thus,

$$
\mathcal{H}^{1}\left(\mathcal{S}_{2}\right)=\mathcal{H}^{1}\left(\mathcal{S}_{Y}\right)+\sqrt{3} h+2 \mathcal{H}^{1}\left(\mathcal{S}_{b}\right) .
$$

Since $\mathcal{S}$ is a Steiner tree, one has $\mathcal{H}^{1}(\mathcal{S}) \leq \mathcal{H}^{1}\left(\mathcal{S}_{1}\right), \mathcal{H}^{1}(\mathcal{S}) \leq \mathcal{H}^{1}\left(\mathcal{S}_{2}\right)$ and clearly $\mathcal{H}^{1}(\mathcal{S})=$ $\frac{\mathcal{H}^{1}\left(\mathcal{S}_{1}\right)+\mathcal{H}^{1}\left(\mathcal{S}_{2}\right)}{2}$. Then $\mathcal{H}^{1}(\mathcal{S})=\mathcal{H}^{1}\left(\mathcal{S}_{1}\right)=\mathcal{H}^{1}\left(\mathcal{S}_{2}\right)$, and so $\mathcal{S}_{1}, \mathcal{S}_{2}$ belong to $\mathbb{M}\left(\left[Y_{\text {up }} Y_{\text {down }}\right] \cup b_{2} \cup c_{2}\right)$. As $\mathcal{S}_{1}$ and $\mathcal{S}_{2}$ are symmetric with respect to $l_{2}$, the statement is proven.

Now we are ready to prove Theorem 2.2.1, i.e., to show that for $\lambda_{j}=\lambda<\frac{1}{300}$, the set $\Sigma(\Lambda)$ is a Steiner tree for the set of terminals $A_{\infty}$.

Let $b_{1}$ and $c_{1}$ be the subsets of terminals that are descendants of $y_{2}$ and $y_{3}$, respectively. Since $\varepsilon=\lambda^{2}+\lambda^{3}+\cdots+\lambda^{k}+\ldots$, we have $b_{i} \subset B_{\varepsilon}\left(B_{i}\right), c_{i} \subset B_{\varepsilon}\left(C_{i}\right)$. Applying Lemma 2.2.1 to $Y_{1}=y_{0}$, $B_{1}=y_{2}, C_{1}=y_{3}, b_{1}$, and $c_{1}$, we show that there is a Steiner tree for $A_{\infty}$, which is symmetric with respect to the line ( $y_{0} y_{1}$ ).

Let $\left[Z_{l} Z_{r}\right]$ and $\left[V_{l} V_{r}\right]$ be the segments from the previous application of Lemma 2.2.1. Now define $b_{2}$ and $c_{2}$ as descendants of $y_{4}$ and $y_{5}$, respectively. Then, applying Lemma 2.2.2 to $\left[Y_{\text {up }} Y_{\text {down }}\right]=\left[Z_{l} Z_{r}\right]$, $B_{2}=y_{4}, C_{2}=y_{5}, b_{2}$, and $c_{2}$ (these data are similar to those required with the scale factor $\lambda$ ), we show that there is a Steiner tree containing $\left[y_{0} y_{1}\right]$ and branching at $y_{1}$ (because $y_{1}$ belongs to the axis of the symmetries of $b$ and $c$ ).

Since $\lambda_{i}$ is constant, the upper and lower components of $\Sigma(\Lambda) \backslash\left[y_{0} y_{1}\right]$ are similar (with the scale factor $\lambda$ ) to $\Sigma(\Lambda)$. Thus, the second application of Lemmas 2.2.1 and 2.2.2 shows that there is a Steiner tree containing $\left[y_{0} y_{1}\right] \cup\left[y_{1} y_{2}\right] \cup\left[y_{1} y_{3}\right]$. This procedure recovers $\Sigma(\Lambda)$ step by step; so after the $k$-th step, we know that the length of every Steiner tree for $A_{\infty}$ is at least

$$
\sum_{i=0}^{k-1}(2 \lambda)^{i}
$$

Thus, the length of every Steiner tree for $A_{\infty}$ is at least the length of $\Sigma(\Lambda)$, which implies $\Sigma(\Lambda) \in$ $\mathbb{M}\left(A_{\infty}\right)$.

Now fix some $\lambda<\frac{1}{25}$, put $\Sigma_{\infty}=\Sigma(\Lambda)$ and define $\Sigma_{k}$ as the union of the edges from the first $k$ levels of $\Sigma_{\infty}$. We will use the following corollary of Theorem 2.2 .2 , which explains why a full binary Steiner tree is universal, i.e. it contains a subtree with a given combinatorial structure.

Corollary 2.2.1. In the conditions of Theorem 2.2.2 each connected closed subset $S$ of $\Sigma_{\infty}$ contained in $\Sigma_{k}$ for some $k$ has a natural tree structure. Moreover, every such an $S$ is the unique Steiner tree for any set $P$ containing the set of the vertices with the degree 1 and 2 of $S$.

Proof. Let $S \subset \Sigma_{\infty}$ and $P \subset S$ satisfy the conditions of the corollary. The fact that $S$ is a tree is straightforward. Let $S^{\prime} \neq S$ be any Steiner tree for $S$ and assume that $\mathcal{H}^{1}\left(S^{\prime}\right) \leq \mathcal{H}^{1}(S)$. Then it is clear that $\mathcal{H}^{1}\left(\left(\Sigma_{k} \backslash S\right) \cup S^{\prime}\right) \leq \mathcal{H}^{1}\left(\Sigma_{k}\right)$, but on the other hand $\left\{y_{0}\right\} \cup A_{k} \subset\left(\Sigma_{k} \backslash S\right) \cup S^{\prime}$, which contradicts to Theorem 2.2.2.

### 2.3 Connectivity of the subset of $\mathbb{P}_{d}$ with a unique Steiner tree

### 2.3.1 Canonical realization of an embedding class

Using the construction of the tree $\Sigma_{\infty}$ from the previous section we define the canonical realization tree $\mathcal{S} t_{E C}$ for any embedding class $E C$. Note that a canonical realization is planar.

Fix a topological tree $S$ with the embedding class $E C$ and pick some vertex $v$ of $S$ of degree one. Then identify $S$ with the subtree of $\Sigma_{\infty}$ by mapping $v$ to the root $y_{0}$ of $\Sigma_{\infty}$ and mapping all the other vertices following the steps of the breadth-first search algorithm started from $v$, where at every vertex of degree 2 of $S$ we choose the left direction in $\Sigma_{\infty}$ (i.e. map the only child to $y_{2 k}$ if the parent was mapped to $y_{k}$ ).

### 2.3.2 Proof of Theorem 2.1.4

In this section we prove Theorem 2.1.4. Note that there are no ambiguous configurations on at most 3 points, so Theorem 2.1.4 clearly holds for $n \leq 3$. Thus we have to deal with $n \geq 4$ to prove the theorem. First we deal with the planar case.

Let us denote by $\mathbb{P}_{2}^{u} \subset \mathbb{P}_{2}$ the subset of configurations having a unique Steiner tree. Observe that, due to Theorem 2.1.2, Theorem 2.1.4 will follow from the following

Theorem 2.3.1. Let $T_{1}, T_{2}$ be two embedding classes and $P_{1}, P_{2} \in \mathbb{P}_{2}^{u}$ be the two configurations of terminal points of the corresponding canonical realizations $\mathcal{S} t_{E C_{1}}, \mathcal{S} t_{E C_{2}}$. Then there is a path in $\mathbb{P}_{2}^{u}$ connecting $P_{1}$ and $P_{2}$.

Define the special (non-labelled) all-left linear tree ALT to be the path on $n$ vertices starting at $y_{0}$ and turning left at every branching point of $\Sigma_{\infty}$, i.e. $A L T$ is the subgraph of $\Sigma_{\infty}$ with the vertices $y_{0}, y_{1}, \ldots, y_{2^{k}}, \ldots, y_{2^{n-2}}$. To establish Theorem 2.3.1 we will show that any $\mathcal{S} t_{E C}$ corresponding to an embedding class with $n$ terminal vertices can be continuously deformed to $A L T$ inside the space of unique Steiner trees with some deformation preserving the labeling of terminal vertices. We construct such a deformation in several steps described below. In each step we continuously deform the set $P$ of the terminal points of $\mathcal{S} t_{E C}$ to the set $P^{\prime}$ of the terminals of $\mathcal{S} t_{E C^{\prime}}$ inside $\mathbb{P}_{2}^{u}$ by moving several points from $P$ one by one.

We need a preliminary lemma.
Lemma 2.3.1. Let EC be an embedding class with generic topology $R$. Suppose that $\mathcal{S} t_{E C}$ contains a leaf $B=y_{2 k}$ adjacent to the terminal $A=y_{k}$ of degree 2. Then one can continuously move $B$ to $y_{2 k+1}$ along a path $\gamma$ in such a way that the whole configuration will remain in $\mathbb{P}_{2}^{u}$ at any point of $\gamma$ and has the embedding class EC.

Proof. We construct a desired part of $\gamma$ explicitly (see Fig. 2.7). First move $B$ into $\mu$-neighborhood of $A$ inside the segment $\left[y_{2 k} y_{k}\right]$, where a small enough $\mu$ will be defined in the next paragraph; by Corollary 2.2.1 the Steiner tree is unique and has the embedding class $E C$ at any configuration from this part of $\gamma$.


Figure 2.7: The construction of $\gamma$ in Lemma 2.3.1.
Let $\bar{R}$ be the topology of $\mathcal{S} t_{E C} \backslash[A B]$; obviously $\bar{R}$ is also generic. Observation 2.1.1 (i) states that $\bar{R}$ lies in exactly one set $D(\bar{O})$, where $\bar{O}$ is a full topology. By Corollary 2.2.1 $\mathcal{S} t_{E C} \backslash[A B]$ is a unique Steiner tree with $n-1$ terminals, so by Observation 2.1.1 (ii) there is $\eta>0$ such that for any other full topology $R^{\prime}$ the length of the realization from $D\left(R^{\prime}\right)$ exceeds $\mathcal{H}^{1}(\mathcal{S} t \backslash[A B])$ by at least $\eta$. Put $\mu=\eta / 2$.

Now rotate $B$ around $A$ : let $B(\alpha), \alpha \in[0,2 \pi / 3]$ be a such point that $|B A|=|B(\alpha) A|$ and the clockwise-oriented angle $y_{[k / 2]} A B(\alpha)$ is equal to $2 \pi / 3+\alpha$. In particular $B(0) \in\left[A y_{2 k}\right), B(2 \pi / 3) \in$ [Ay $y_{2 k+1}$ ).

Let $\mathcal{S} t(\alpha)$ be a Steiner tree for the terminals of $\mathcal{S} t \backslash[A B]$ and $B(\alpha)$. Then

$$
\mathcal{H}^{1}(\mathcal{S t}(\alpha)) \leq \mathcal{H}^{1}(\mathcal{S} t \backslash[A B] \cup[A B(\alpha)])=\mathcal{H}^{1}(\mathcal{S} t \backslash[A B])+\mu<\mathcal{H}^{1}(\mathcal{S} t \backslash[A B])+\eta
$$

Let $O$ be the full topology such that $R \in D(O)$. Then the topology of $\mathcal{S} t(\alpha)$ belongs to $D(O)$. By Proposition 2.1.2 $\mathcal{S} t(\alpha)$ is uniquely defined. The set $\mathcal{S} t \backslash[A B] \cup[A B(\alpha)]$ is locally minimal, and has the topology from $D(R)$, so it coincides with $\mathcal{S t}(\alpha)$. Thus not only the topology but the embedding class is preserved during this part of the path.

Finally, move $B$ from $B(2 \pi / 3)$ to $y_{2 k}$ inside the segment $\left[y_{2 k} y_{k}\right]$.

Let us now fix an embedding class $E C$ and construct the desired deformation of $\mathcal{S} t_{E C}$ to $A L T$ inside the space of unique Steiner trees.

Step 1. Transform $\mathcal{S} t_{E C}$ into a full Steiner tree $\mathcal{S} t_{E C^{\prime}}$ inside $\mathbb{P}_{2}^{u}$. To make such a transformation we need to move all terminal vertices of $\mathcal{S} t_{E C}$ of degree 2 or 3 to make them leaves.

Suppose first that $\mathcal{S} t_{E C}$ contains a terminal $A=y_{j}$ of degree 2. By the construction of $\mathcal{S} t_{E C}$, the vertex $A$ is adjacent to vertices $B=y_{2 j}$ and $C=y_{\lfloor j / 2\rfloor}$, which may be terminals or Steiner points. Move $A$ towards $y_{2 j+1}$ along the edge $y_{j} y_{2 j+1}$ of $\Sigma_{\infty}$ until it hits $y_{2 j+1}$ (see. Fig. 2.8). Corollary 2.2.1


Figure 2.8: Elimination of points with degree 2 in $\mathcal{S} t_{E C}$
ensures that the obtained deformation of the set of terminal points lies inside $\mathbb{P}_{2}^{u}$. Applying this deformation to each terminal vertex of degree 2 one by one we eventually get rid of those.

Assume now that $\mathcal{S} t_{E C}$ has a terminal point of degree 3. Since $\mathcal{S} t_{E C}$ has no terminal of degree two and the number of Steiner points is at most the number of leafs minus two, one may move terminals of degree three one by one in the neighborhood of different leaves by a path in $\mathbb{P}_{2}^{u}$. From now on, the topology of a tree is generic.

Now consider any point $A$ in an $\varepsilon$-neighborhood of a leaf $B$ for some small $\varepsilon$. Then continue moving $A$ while moving $B$ simultaneously in the same direction until $A$ reaches $y_{k}$ and $B$ reaches $B(\pi / 3)$ (see Fig. 2.9). Now stop moving $A$, but rotate $B$ around $A$ until it hits the ray $A y_{2 k+1}$, then extend $B$ to $y_{2 k+1}$ and $A$ to $y_{2 k}$. Now all our terminal points again belong to the set $\left\{y_{0}, y_{1}, \ldots\right\}$ and the unique Steiner tree is given by the canonical realization $\mathcal{S} t_{E C^{\prime}}$ for some new embedding class $E C^{\prime}$. The fact that the set of terminal points was staying inside $\mathbb{P}_{2}^{u}$ while we were moving them follows from Corollary 2.2.1 and the proof of Lemma 2.3.1.


Figure 2.9: Elimination of points with degree 3 in $\mathcal{S} t_{E C}$

Note that $\mathcal{S t}_{E C^{\prime}}$ still has no terminal points of degree 2 and has one less terminal point of degree 3 than $\mathcal{S} t_{E C}$. Hence we can do this procedure until we obtain a canonically realized full tree.

Step 2. Permute the labels of terminal points of $\mathcal{S} t_{E C}$ if necessary. Now we can assume that $\mathcal{S} t_{E C}$ is a full tree. By Theorem 2.1 .2 one may put label 1 into the root by a path in $\mathbb{P}_{2}^{u}$ and do not touch the root of the tree later on. Let $A$ and $B$ be the two terminal points which we want to swap.

Since $\mathcal{S} t_{E C}$ is a full tree, it has exactly $n-2$ Steiner points. In particular, we can choose two Steiner points of $\mathcal{S} t_{E C}$ that are connected with two terminal points of $\mathcal{S} t_{E C}$; let $y_{k}$ be the one of them which is not adjacent to the root of $\Sigma_{\infty}$. Denote the terminals adjacent to $y_{k}$ by $B=y_{2 k}$ and $A=y_{2 k+1}$.

We may swap $A$ with any label. First swap $A$ and $B$ as shown at Fig. 2.10 move $B$ into $y_{k}$ and $A$ in a small neighborhood of $y_{k}$, then turn and finally make a reverse procedure. By Lemma 2.3.1 the Steiner tree is unique during the middle part of this procedure; by Corollary 2.2.1 the Steiner tree is unique during other parts.


Figure 2.10: Swapping the labels of terminals connecting with a common branching point

Then swap $A$ with any terminal $C \neq B$ of $\mathcal{S} t_{i}^{\prime}$ (see Fig. 2.11). Start with the previous procedure and stop it at the point $B=B(\pi / 3)$ (in the notation of Lemma 2.3.1). Then $A$ moves inside the tree into a neighborhood of $C=y_{l}$ and $B$ comes to $y_{k}$. We are going to apply Lemma 2.3.1 to $A$ and $C$ : move $C$ to $C(\pi / 3)$ and $A$ to $y_{l}$. Then $C$ rotates to $y_{2 l+1}$, after that $A$ moves to $y_{2 l}$. Now the positions of $A$ and $C$ are symmetric so we may do the reverse procedure after swapping $A$ and $C$.


Figure 2.11: Swapping the labels of arbitrarily terminals

Finally to swap labels of arbitrary terminals $C$ and $D$ we swap $A=y_{2 k+1}$ with $C, C=y_{2 k+1}$ with $D$ and $D=y_{2 k+1}$ with $A$. Since the set of all transpositions spans the symmetric group we may construct a path in $\mathbb{P}_{2}^{u}$ connecting $\mathcal{S} t_{E C}$ and the same tree with an arbitrary permutation of its labels. Till the end of the section all trees are not labelled.

Step 3. Connect $\mathcal{S} t_{i}^{\prime}$ with the all-left linear tree $A L T$ by a path in $\mathbb{P}_{2}^{u}$. While there is a terminal point $A=y_{j}$ of $\mathcal{S} t_{i}^{\prime}$ not belonging to $A L T$, consider such a vertex with the largest $j$. It implies that the degree of $A$ is 1 . Our aim is to move $A$ inside $\Sigma_{\infty}$ to the first vertex $y_{w}$ of $A L T$ which does not belong to $\mathcal{S} t_{i}^{\prime}$.

Consider the case when $A$ is adjacent to a branching point $y_{l}$ then $j=2 l+1$ and $B=y_{2 l}$ is also a terminal of $\mathcal{S} t_{i}^{\prime}$ because of the maximality of $j$. Move $A$ into $y_{l}$ and rotate $B$ into $B(\pi / 3)$ (in the notation of Lemma 2.3.1). Then $A$ moves into the tree and $B$ moves into $y_{l}$.

Now $A$ is either inside the tree or $A$ is a terminal connected with a vertex of degree 2 . Move $A$ into $y_{w}$, the only problem is that $A$ cannot coincide with the terminal of degree 2 . Movement through a terminal of degree 2 is depicted in Fig. 2.12.


Figure 2.12: Movement through a terminal of degree 2

Finally all the vertices of $\mathcal{S} t_{i}^{\prime}$ belong to $A L T$, so we are done. Since we connect $\mathcal{S} t_{E C_{1}}$ with $A L T$ and $\mathcal{S} t_{E C_{2}}$ with $A L T$, the desired $\gamma$ is constructed.

Proof of Theorem 2.1.4. Let $P_{1}, P_{2} \in \mathbb{P}_{d}$ be configurations with unique Steiner trees $\mathcal{S} t\left(P_{1}\right)$ and $\mathcal{S t}\left(P_{2}\right)$ having embedding classes $T_{1}$ and $T_{2}$, respectively. By Theorem 2.1 .2 there is a path $\gamma_{i}$ between $\mathcal{S} t\left(P_{i}\right)$ and $\mathcal{S} t_{T_{i}}$ in $\mathbb{P}_{d}$ such that the Steiner tree is unique during $\gamma_{i}$. We have constructed the path $\gamma$ between $\mathcal{S} t_{T_{1}}$ and $\mathcal{S} t_{T_{2}}$ in $\mathbb{P}_{2} \subset \mathbb{P}_{d}$; the Steiner tree is also unique during $\gamma$. The gluing of $\gamma_{1}, \gamma$ and $\gamma_{2}^{-1}$ finishes the proof.

### 2.4 Proof of Theorem 2.1.1

This section is devoted to the proof of Theorem 2.1.1. The proof is using the theory of subanalytic sets, and for the sake of completeness we begin our exposition with a brief reminder of some definitions and facts from this theory.

### 2.4.1 Subanalytic subsets of a real analytic manifold

All the facts expounded in this section are well-known and may be skipped by an advanced reader. During our exposition we mostly follow Sections 2 and 3 from the paper [9].

Let $M$ be a real analytic manifold and $\mathcal{O}_{M}$ denote the sheaf of real analytic functions on $M$, that is, for any open $U \subset M$ the set $\mathcal{O}_{M}(U)$ is the space of real analytic functions defined on $U$. We introduce the following definitions:

1. A subset $A \subset M$ is called an analytic submanifold if for each $p \in A$ there exists a neighborhood $U \subset M$ of it such that either $A \cap U=U$, or there exist a finite collection $f_{1}, f_{2}, \ldots, f_{k} \in \mathcal{O}_{M}(U)$
such that $A \cap U$ is the set of common zeros of $f_{1}, \ldots, f_{k}$ and for any $x \in A \cap U$ the gradients $\nabla f_{1}(x), \ldots, \nabla f_{k}(x)$ are linearly independent.
2. A subset $A \subset M$ is called analytic if for each $p \in M$ there exists a neighborhood $U$ of $p$ such that either $A \cap U=U$, or there exists a finite set of functions $f_{1}, \ldots, f_{k} \in \mathcal{O}_{M}(U)$ such that $A \cap U$ is the set of common zeros of $f_{1}, \ldots, f_{k}$. Note that we require this property for all $p \in M$, not only for $p \in A$.
3. A subset $A \subset M$ is called semianalytic if for each point $p \in M$ there exists a neighborhood $U \subset M$ and a finite number of subsets $A_{i, j} \subset U$ such that $A \cap U=\cup_{i} \cap_{j} A_{i, j}$ and each $A_{i, j}$ is of the form $\{f>0\}$ or $\{f=0\}$ for some $f \in \mathcal{O}_{M}(U)$. A semianalytic subset $A$ is called smooth if it is an analytic submanifold.

The following lemma follows from [9, Proposition 2.10]:
Lemma 2.4.1. Let $M$ be a real analytic manifold, $A \subset M$ be a semianalytic subset and $p \in M$ be an arbitrary point. Then there exists a neighborhood $U \subset M$ of $p$ and a finite collection of disjoint subsets $A_{1}, A_{2}, \ldots, A_{k} \subset U$ such that

1. each of $A_{1}, \ldots, A_{k}$ is a semianalytic subset of $U$ and an analytic submanifold of $M$, and
2. $A \cap U$ is a disjoint union of $A_{1}, \ldots, A_{k}$.

Semianalytic sets admit many properties similar to those of semialgebraic sets (i.e. those given by polynomial inequalities), but the theories are not identical. An important difference is that projections of semianalytic sets are not necessarily semianalytic (see [9, Example 2.14]), while projections of semialgebraic sets are always semialgebraic. This motivates the following definition:

Definition 2.4.1. A subset $X \subset M$ is called subanalytic if for any $p \in M$ there is a neighborhood $U$ of $p$, an analytic manifold $N$ and a relatively compact semianalytic subset $A \subset M \times N$ such that $X \cap U=\pi(A)$, where $\pi$ is the projection on $M$.

The following lemma follows immediately from this definition:
Lemma 2.4.2. Let $N, M$ be two analytic manifolds and $f: N \rightarrow M$ be an analytic map. Assume that $A \subset N$ is semianalytic and for each $p \in M$ there is a neighborhood $U$ of it such that $f^{-1}(U) \cap A$ is relatively compact in $N$. Then $f(A)$ is a subanalytic subset of $M$.

Proof. Let $\Gamma_{f} \subset M \times N$ be the graph of the mapping $f$. The graph is an analytic subset of $M \times N$ since $f$ is analytic. Let $p \in M$ and $U \subset M$ be a semianalytic relatively compact neighborhood such that $f^{-1}(U) \cap A$ is relatively compact in $N$. Define $B=(U \times A) \cap \Gamma_{f} \subset M \times N$, then $B$ is semianalytic and relatively compact. We have $f(A) \cap U=\pi(B)$, where $\pi$ is the projection on $M$. Since $p$ were arbitrary, we conclude that $f(A)$ is subanalytic.

Subanalytic sets, although not being semianalytic in general, still have a lot of nice properties. Direct products, finite intersections and unions, closures, complements and, thus, interiors of subanalytic sets are still subanalytic (see [9, Chapter 3]). The following lemma describes a local structure of subanalytic sets, see [9, Lemma 3.4]):

Lemma 2.4.3. Let $N, M$ be analytic manifolds and $A \subset N \times M$ be a relatively compact semianalytic subset. Then there exists a finite collection of smooth connected semianalytic subsets $A_{1}, \ldots, A_{k} \subset$ $N \times M$ such that

1. $A=\sqcup_{j=1}^{k} A_{j}$,
2. for any $j$ the rank of $d \pi$ on $T_{x} A_{j}$ does not depend on $x \in A_{j}$.

From this lemma we get an immediate corollary:
Corollary 2.4.1. Let $M$ be an analytic manifold and $X \subset M$ be a subanalytic subset. Then there exists a countable collection of connected analytic submanifolds $X_{1}, X_{2}, X_{3}, \ldots$ of $M$ such that $X=$ $X_{1} \cup X_{2} \cup X_{3} \cup \ldots$

Proof. Since the topology of $M$ has a countable base, it is enough to prove the statement of the corollary locally. Passing to a neighborhood of some point if necessary we can assume that there is a relatively compact semianalytic subset $A \subset M \times N$ for some real analytic manifold $N$ such that $X=\pi(A)$. Let $A_{1}, \ldots, A_{k} \subset M \times N$ be such as in Lemma 2.4.3. For each $j=1, \ldots, k$ there is a countable collection of open subsets $U_{j 1}, U_{j 2}, \ldots$ of $M \times N$ covering $A_{j}$ and such that $\pi\left(U_{j i} \cap A_{j}\right)$ is a connected analytic submanifold of $M$. Then we have

$$
X=\pi(A)=\bigcup_{j=1}^{k} \bigcup_{i \geq 1} \pi\left(U_{j i} \cap A_{j}\right)
$$

For future needs we now recall the notion of a fiber product. Let $X, Y, U$ be some sets and $f: X \rightarrow U, g: Y \rightarrow U$ be some maps between these sets. The fiber product of $X$ and $Y$ over the base $U$ is defined as

$$
\begin{equation*}
X \times_{f=g} Y=\{(x, y) \in X \times Y \quad \mid f(x)=g(y)\} \tag{2.3}
\end{equation*}
$$

Note that we have a natural projection $X \times_{f=g} Y \rightarrow U$ which sends $(x, y)$ to $f(x)$.
Lemma 2.4.4. Assume that $M, N, U$ are real analytic manifolds and $f: M \rightarrow U$ and $g: N \rightarrow U$ are real analytic maps. Let $X \subset M$ and $Y \subset N$ be sub- or semianalytic subsets. Then $X \times_{f=g} Y$ is a sub- or semianalytic subset of $M \times_{f=g} N$ respectively.

Proof. Follows immediately from definitions. Indeed, $X \times_{f=g} Y$ is the intersection of the sub- or semianalytic set $X \times Y$ and the subset $\{(x, y) \in M \times N \mid f(x)=g(y)\}$ inside $M \times N$, hence is subor semianalytic respectively.

### 2.4.2 Proof of Theorem 2.1.1

In this section we prove Theorem 2.1.1. Recall that a number $n \geq 4$ of terminals is fixed and $\mathbb{P}_{2}$ is equal to $\left(\mathbb{R}^{2}\right)^{n}$ with diagonals removed. Denote the subset of ambiguous configurations by $\mathcal{A}$. Let also $\mathcal{A}_{\text {non-generic }}$ denote the set of all configurations admitting non-generic Steiner tree. Recall that $\operatorname{dim} \mathcal{A}_{\text {non-generic }}=2 n-2$ by Observation 2.1.1.

We begin with the following
Lemma 2.4.5. The Hausdorff dimension of $\mathcal{A}$ is at least $2 n-1$.
Proof. Let $T_{1}, \ldots, T_{N}$ be all possible full topologies on $n$ points and $V\left(T_{j}\right)$ denotes the set of vertices of $T_{j}$. Given a map $f: V\left(T_{j}\right) \rightarrow \mathbb{R}^{2}$, let $L(f)$ the total length of the segments connecting $f\left(V\left(T_{j}\right)\right)$
accordingly to topology $T_{j}$ (note that we do not claim any restrictions, in particular absence of cycles or local minimality), i.e.

$$
L(f)=\sum_{v w \text { is an edge of } T_{j}}|f(v)-f(w)| .
$$

As it follows from Proposition 2.1.2, for any $P \in \mathbb{P}_{2}$ and $j=1, \ldots, N$ there exists precisely one map $f: V\left(T_{j}\right) \rightarrow \mathbb{R}^{2}$ which maps the terminals of $T_{j}$ to the points from $P$ keeping the enumeration and which minimizes $L(f)$ among all such maps. Set $L_{j}(P)=L(f)$ in this case. Note that $L_{j}$ is a continuous function on $\mathbb{P}_{2}$. Define

$$
B_{j}=\left\{P \in \mathbb{P}_{2} \mid \quad L_{j}(P)<L_{i}(P) \forall i \neq j\right\} .
$$

It follows that $B_{j}$ 's are open and disjoint sets. We also have $B_{j} \neq \varnothing$; indeed, by Corollary 2.2 .1 each $T_{j}$ is the topology of some Steiner tree which is unique. Note that $A=\mathbb{P}_{2} \backslash\left(\cup_{j=1}^{N} B_{j}\right) \subset \mathcal{A} \cup \mathcal{A}_{\text {non-generic }}$ by Observation 2.1.1. The lemma now follows from $\operatorname{dim} \mathcal{A}_{\text {non-generic }}=\operatorname{dim} \mathbb{R}^{2 n} \backslash \mathbb{P}_{2}=2 n-2$ and Lemma 2.4.6.

Lemma 2.4.6. Assume that $N \geq 2, m \geq 1$ and $B_{1}, \ldots, B_{N} \subset \mathbb{R}^{m}$ are non-empty disjoint open sets. Put $A=\mathbb{R}^{m} \backslash\left(\cup_{j=1}^{N} B_{j}\right)$. Then $\operatorname{dim} A \geq m-1$, where $\operatorname{dim}$ is the Hausdorff dimension.

Proof. Let $P_{1}, P_{2} \in \mathbb{R}^{m}$ be such that $P_{1} \in B_{1}$ and $P_{2} \in B_{2}$. Let $l$ be the line passing through these points and $V \subset \mathbb{R}^{m}$ be the subspace of codimension 1 orthogonal to $l$. Let $\pi: \mathbb{R}^{m} \rightarrow V$ be the orthogonal projection. Note that $v \in \pi(A)$ if and only if the line $\pi^{-1}(v)$ intersects $A$. In particular, $v_{0}=\pi(l) \in \pi(A)$ and moreover there exists $\varepsilon>0$ such that $v \in \pi(A)$ if $\left|v-v_{0}\right| \leq \varepsilon$ since $B_{1}, B_{2}$ are open. It follows that $\pi(A)$ has a non-empty interior as a subset of $V$ and $\operatorname{dim} \pi(A)=m-1$. Since $\pi$ is 1 -Lipschitz, it implies that $\operatorname{dim} A \geq m-1$.

The converse estimate $\operatorname{dim} \mathcal{A} \leq 2 n-1$ is more involved and requires some additional constructions. Let $T$ be some (not necessarily full or generic) topology with $n$ terminals. Enumerate the Steiner points of $T$ arbitrarily, let $k$ be the total amount of them. Given $(P, q)=\left(p_{1}, \ldots, p_{n}, q_{1}, \ldots, q_{k}\right) \in$ $\mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k}$, we can identify the corresponding points on the plane with the vertices of $T$ following the enumeration and connect a pair of corresponding points by a straight segment for each edge of $T$. Thus, any such $(P, q)$ defines a map from $T$ to the plane. Let $\mathcal{R}_{\text {geoemb }}(T) \subset \mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k}$ be the following:

$$
\mathcal{R}_{\text {geoemb }}(T)=\left\{(P, q) \in \mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k} \mid(P, q) \text { defines an embedding of } T\right\}
$$

Obviously, $\mathcal{R}_{\text {geoemb }}(T)$ is an open subset of $\mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k}$. Let $L_{T}: \mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k} \rightarrow \mathbb{R}$ be the function that computes the length of the image of $T$; note that $L_{T}$ is real analytic on $\mathcal{R}_{\text {geoemb }}(T)$ and continuous everywhere. Define

$$
\mathcal{R}_{\text {lencrit }}(T)=\left\{(P, q) \in \mathcal{R}_{\text {geoemb }}(T) \mid \nabla_{q} L_{T}(P, q)=0\right\}
$$

where $\nabla_{q}$ is the gradient with respect to the variable $q$. Since $L_{T}$ is real-analytic, $\mathcal{R}_{\text {lencrit }}(T)$ is an analytic subset of $\mathcal{R}_{\text {geoemb }}(T)$. We have the following

Lemma 2.4.7. The following statements hold:

1. Let $(P, q) \in \mathcal{R}_{\text {lencrit }}(T)$. Consider the function $L_{T}(P, \cdot)$ as a continuous function from $\left(\mathbb{R}^{2}\right)^{k}$ to $\mathbb{R}$. Then $q$ is the unique point of the global minimum of $L_{T}(P, \cdot)$.
2. Let $P_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathbb{P}_{2}$ be the projection. Then $P_{T}$ restricted to $\mathcal{R}_{\text {lencrit }}(T)$ is injective and the set $P_{T}\left(\mathcal{R}_{\text {lencrit }}(T)\right)$ is open in $\mathbb{P}_{2}$.

Proof. Note that for any fixed $P$ the function $L_{T}(P, q)$ tends to infinity as the distance between $q$ and the origin tends to infinity. It follows that for any $P$ there is a point $q(P) \in\left(\mathbb{R}^{2}\right)^{k}$ where $L_{T}(P, \cdot)$ attains its global minimum. By Proposition 2.1.2 such a point is always unique and there are no other local minima of $L_{T}(P, \cdot)$. In particular, the mapping $P \mapsto q(P)$ is a well-defined mapping $\mathbb{P}_{2} \rightarrow\left(\mathbb{R}^{2}\right)^{k}$. From the continuity of $L_{T}$ it follows that this mapping is continuous.

To prove the first item, it is enough to show that whenever $(P, q) \in \mathcal{R}_{\text {lencrit }}(T)$, the value $L_{T}(P, q)$ is a local minimum of $L_{T}(P, \cdot)$. Let $v_{1}, \ldots, v_{n}$ be the terminals of $T$, let $\tilde{T}_{1}, \ldots, \tilde{T}_{l}$ be the connected components of $T \backslash\left\{v_{1}, \ldots, v_{n}\right\}$ containing Steiner points and let $T_{i}$ be the closure of $\tilde{T}_{i}$ in $T$. Then each $T_{i}$ is a full topology (with terminals being a subset of terminals of $T$ ). Let $f: T \rightarrow \mathbb{R}^{2}$ be the embedding corresponding to $(P, q)$, then it is easy to see that the differential condition $\nabla_{q} L_{T}(p, q)=0$ is equivalent to the fact that $f\left(T_{i}\right)$ is a locally minimal tree, which implies the statement.

For the second item, consider the mapping $M: \mathbb{P}_{2} \rightarrow \mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k}$ which sends $P$ to $(P, q(P))$. Then $M$ is continuous. But since $P_{T}\left(\mathcal{R}_{\text {lencrit }}(T)\right)=M^{-1}\left(\mathcal{R}_{\text {geoemb }}(T)\right)$ and $\mathcal{R}_{\text {geoemb }}(T)$ is open in $\mathbb{P}_{2} \times\left(\mathbb{R}^{2}\right)^{k}$, we conclude that $\mathcal{R}_{\text {lencrit }}(T)$ is open.

Recall that $P_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathbb{P}_{2}$ is the projection. Given two topologies $T_{1}, T_{2}$ with $n$ labelled vertices, define

$$
\begin{aligned}
& \mathcal{A}_{T_{1}, T_{2}}=\left\{P \in P_{T_{1}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{1}\right)\right) \cap P_{T_{2}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{2}\right)\right) \mid\right. \\
& \left.L_{T_{1}}\left(P, q_{1}\right)=L_{T_{2}}\left(P, q_{2}\right), \text { where }\left(P, q_{i}\right) \in \mathcal{R}_{\text {lencrit }}\left(T_{i}\right)\right\} .
\end{aligned}
$$

In the next lemma we will use the notion of a subanalytic subset introduced in Section 2.4.1.
Lemma 2.4.8. Let $T_{1} \neq T_{2}$ be two generic topologies with $n$ terminals. Then there exists an open set $U \subset \mathbb{P}_{2}$ such that $\mathcal{A}_{T_{1}, T_{2}}$ is a subanalytic subset of $U$. In particular, $\mathcal{A}_{T_{1}, T_{2}}$ is a union of a countable collection of connected analytic submanifolds of $U$.

Proof. Recall that for any $T$ the map $P_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathbb{P}_{2}$ is the projection. Define $U=$ $P_{T_{1}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{1}\right)\right) \cap P_{T_{2}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{2}\right)\right)$. By Lemma 2.4.7 we have that $U$ is an open subset of $\mathbb{P}_{2}$, and we have $\mathcal{A}_{T_{1}, T_{2}} \subset U$ by the definition of $\mathcal{A}_{T_{1}, T_{2}}$.

Introduce the temporary notation $\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)_{U}=\mathcal{R}_{\text {lencrit }}\left(T_{i}\right) \cap P_{T_{i}}^{-1}(U)$ and

$$
\mathcal{R}_{\text {geoemb }}\left(T_{i}\right)_{U}=\mathcal{R}_{\text {geoemb }}\left(T_{i}\right) \cap P_{T_{i}}^{-1}(U)
$$

for simplicity. Recall the definition of a fiber product was introduced in (2.3). As we can see from the definition,

$$
\mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U} \subset U \times\left(\mathbb{R}^{2}\right)^{k_{1}} \times\left(\mathbb{R}^{2}\right)^{k_{2}}
$$

is an open subset, where $k_{i}$ is the number of Steiner points of $T_{i}$ and we identify $U \times\left(\mathbb{R}^{2}\right)^{k_{1}} \times\left(\mathbb{R}^{2}\right)^{k_{2}}$ with $\left(U \times\left(\mathbb{R}^{2}\right)^{k_{1}}\right) \times_{P_{T_{1}}=P_{T_{2}}}\left(U \times\left(\mathbb{R}^{2}\right)^{k_{2}}\right)$. Therefore $\mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U}$ is a real analytic submanifold of $\mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U}$. The set $\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)_{U}$ is an analytic subset of $\mathcal{R}_{\text {geoemb }}\left(T_{i}\right)_{U}$, hence by Lemma 2.4.4 $\mathcal{R}_{\text {lencrit }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {lencrit }}\left(T_{2}\right)_{U}$ is an analytic subset of $\mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U}$.

Define now

$$
\mathcal{R}=\left\{\left(P, q_{1}, q_{2}\right) \in \mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U} \mid L_{T_{1}}\left(P, q_{1}\right)=L_{T_{2}}\left(P, q_{2}\right)\right\} .
$$

Then $\mathcal{R}$ is an analytic subset of $\mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U}$. Denote by $\pi$ the natural projection $\pi: \mathcal{R}_{\text {geoemb }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {geoemb }}\left(T_{2}\right)_{U} \rightarrow U$. Then we have

$$
\mathcal{A}_{T_{1}, T_{2}}=\pi\left(\mathcal{R} \cap\left(\mathcal{R}_{\text {lencrit }}\left(T_{1}\right)_{U} \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {lencrit }}\left(T_{2}\right)_{U}\right)\right)
$$

It follows from Lemma 2.4 .2 that $\mathcal{A}_{T_{1}, T_{2}}$ is a subanalytic subset of $U$.
The last assertion of the lemma follows from Corollary 2.4.1.
Let us say that trees $S_{1}$ and $S_{2}$ are codirected at a terminal $v$ is for a small enough neighborhood $U \ni v$ one has $S_{1} \cap U=S_{2} \cap U$.

Lemma 2.4.9. Let $T_{1} \neq T_{2}$ be two generic topologies with $n$ terminals, and assume that $P \in$ $\operatorname{Int}\left(\mathcal{A}_{T_{1}, T_{2}}\right)$ (here Int stands for the interior in $\left.\mathbb{P}_{2}\right)$ and $S_{1}(P), S_{2}(P)$ are the images of $T_{1}, T_{2}$ on the plane. Then for each terminal $v$ we have $\operatorname{deg}_{T_{1}} v=\operatorname{deg}_{T_{2}} v$ and the trees $S_{1}(P), S_{2}(P)$ are codirected at $v$.

Proof. Given $P \in \operatorname{Int} \mathcal{A}_{T_{1}, T_{2}}$, denote by $q_{1}(P)$ and $q_{2}(P)$ the configurations of Steiner points such that we have $\left(P, q_{i}(P)\right) \in \mathcal{R}_{\text {lencrit }}\left(T_{i}\right), i=1,2$. Note that it is enough to prove the lemma for the points $P$ from a dense subset of the interior, since $q_{i}(P)$ depends continuously on $P$ (cf. the proof of Lemma 2.4.7). Recall that we denote by $P_{T_{i}}: \mathcal{R}_{\text {geoemb }}\left(T_{i}\right) \rightarrow \mathbb{P}_{2}$ the projection and $P_{T_{i}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)\right)$ is open in $\mathbb{P}_{2}$; recall also that $P_{T_{i}}$ restricted to $\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)$ is one-to-one by Lemma 2.4.7 and $q_{i}$ is the inverse mapping. From Lemma 2.4 .3 and the fact that $P_{T_{i}}$ is one-to-one restricted to $\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)$ we deduce that $q_{i}$ is differentiable on $P_{T_{i}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)\right)$ outside a subset which is nowhere dense in $P_{T_{i}}\left(\mathcal{R}_{\text {lencrit }}\left(T_{i}\right)\right)$. Thus, deforming $P$ inside $\operatorname{Int}\left(\mathcal{A}_{T_{1}, T_{2}}\right)$ a little bit we can achieve that that both $q_{1}$ and $q_{2}$ are differentiable in a neighborhood of $P$.

Further, we claim that deforming $P$ a little bit more we can assume that for any terminal vertex of degree 2 in $S_{i}(P)$ the angle between the corresponding edges in $S_{i}(P)$ is not equal to $\pi$ or $2 \pi / 3$. Indeed, let $v$ be such a vertex in, say, $S_{1}(P)$. Then $v$ divides $S_{1}(P)$ into two subtrees $S_{1}^{+}(P)$ and $S_{1}^{-}(P)$. Rotating $S_{1}^{-}(P)$ around $v$ a little bit we can assure that the angle at $v$ in $S_{1}(P)$ is not equal to $\pi$ or $2 \pi / 3$, and $q_{1}, q_{2}$ are still differentiable in a neighborhood of $P$. Repeating this for all terminal vertices of degree 2 in $S_{1}(P)$ and $S_{2}(P)$ we get the result; note that directions of edges in $S_{1}(P), S_{2}(P)$ depend on $P$ continuously, because $q_{1}(P), q_{2}(P)$ are continuous functions of $P$.

Given a point $v$ from $P$ and an oriented edge $\vec{e}$ of $S_{i}(P)$ emanating from $v$ in $S_{i}(P)$ denote by $\eta_{i}(v, \vec{e}) \in \mathbb{R}^{2}$ the unit vector in $\mathbb{R}^{2}$ codirected with $\vec{e}$. Set

$$
\eta_{i}(v)=\sum_{\vec{e} \in \vec{E}\left(S_{i}(P)\right): o(\vec{e})=v} \eta_{i}(v, \vec{e}) ;
$$

note that the sum consists of one or two elements for each $v$ since all terminal vertices have degrees 1 or 2 . A direct computation using angle condition at Steiner points of $S_{i}(P)$ shows that for any $\mu \in\left(\mathbb{R}^{2}\right)^{n}$ the derivative in the direction $\mu$ of $L_{T_{i}}\left(P, q_{i}(P)\right)$ is given by

$$
\frac{\partial}{\partial \mu} L_{T_{i}}\left(P, q_{i}(P)\right)=-\sum_{j=1}^{n} \eta_{i}\left(v_{j}\right) \cdot \mu_{j}
$$

where $v_{j}$ is the $j$-th terminal. Since $L_{T_{1}}\left(\cdot, q_{1}(\cdot)\right)$ and $L_{T_{2}}\left(\cdot, q_{2}(\cdot)\right)$ are equal in a neighborhood of $P$, we conclude that $\eta_{1}\left(v_{j}\right)=\eta_{2}\left(v_{j}\right)$ for any $j=1, \ldots, n$. If $\operatorname{deg}_{S_{1}(P)} v_{j}=\operatorname{deg}_{S_{2}(P)} v_{j}=1$, then this means that $S_{1}(P)$ and $S_{2}(P)$ are codirected at $v_{j}$. Assume that $\operatorname{deg}_{S_{1}(P)} v_{j}=2$; then $\left|\eta_{1}\left(v_{j}\right)\right| \neq 1$
since the angle between the two edges emanating from $v_{j}$ is not equal to $2 \pi / 3$ by our assumption. From this inequality and the equality $\eta_{1}\left(v_{j}\right)=\eta_{2}\left(v_{j}\right)$ we find out that $\operatorname{deg}_{S_{2}(P)} v_{j}=2$ also. Further, we have $\eta_{1}\left(v_{j}\right) \neq 0$ since the angle between the two edges emanating from $v_{j}$ is not equal to $\pi$ by our assumption. As a consequence, there is only one unordered pair of unit vectors ( $\mu_{+}, \mu_{-}$) such that $\eta_{1}\left(v_{j}\right)=\mu_{+}+\mu_{-}$, hence the pair of edges emanating from $v_{j}$ in both $S_{1}(P)$ and $S_{2}(P)$ must have these directions, so that $S_{1}(P)$ and $S_{2}(P)$ are codirected at $v_{j}$. We conclude that $S_{1}(P)$ and $S_{2}(P)$ are codirected at all terminals.

Let us now formulate the following theorem by Oblakov:
Theorem 2.4.1 (Oblakov [96]). Assume that $S_{1}$ and $S_{2}$ are two locally minimal trees connecting the same set of terminals $P \in \mathbb{P}_{2}$ and codirected at this set. Then $S_{1}$ and $S_{2}$ coincide.

Given two generic topologies $T_{1}$ and $T_{2}$ with $n$ terminals and $P \in \mathcal{A}_{T_{1}, T_{2}}$, denote by $S_{1}(P)$ and $S_{2}(P)$ the embeddings of $T_{1}$ and $T_{2}$ as in the lemma above. Define

$$
\mathcal{A}_{T_{1}, T_{2}}^{\min }=\left\{P \in \mathcal{A}_{T_{1}, T_{2}} \mid S_{1}(P) \text { and } S_{2}(P) \text { are both locally minimal trees }\right\} .
$$

We have the following
Corollary 2.4.2. Let $T_{1} \neq T_{2}$ be two generic topologies with $n$ terminals. Then we have

$$
\operatorname{Int}\left(\mathcal{A}_{T_{1}, T_{2}}\right) \cap \mathcal{A}_{T_{1}, T_{2}}^{\min }=\varnothing,
$$

where Int stands for the interior in $\mathbb{P}_{2}$.
Proof. Follows from Lemma 2.4 .9 and Theorem 2.4.1.
Lemma 2.4.10. We have

$$
\mathcal{A} \subset \bigcup_{\substack{T_{1} \neq T_{2} \\ T_{1}, T_{2} \text { are generic }}} \mathcal{A}_{T_{1}, T_{2}}^{\min } \cup\left\{P \in \mathbb{P}_{2} \mid \text { there is a Steiner tree for } P \text { with a non-generic topology }\right\} .
$$

Proof. Let $P \in \mathcal{A}$, then $P$ is connected by two Steiner trees $S_{1}$ and $S_{2}$. If both $S_{1}$ and $S_{2}$ have the same topology $T$, then we get the contradiction with Lemma 2.4.7. Thus the topologies of $S_{1}$ and $S_{2}$ are different and the lemma follows.

We can now prove Theorem 2.1.1
Proof of Theorem 2.1.1. By Lemma 2.4.5 the dimension of the set $\mathcal{A}$ of ambiguous configurations is at least $2 n-1$. Therefore, by Lemma 2.4.10 and Observation 2.1.1, (iv), we have

$$
\operatorname{dim} \mathcal{A} \leq \max \left\{\operatorname{dim} \mathcal{A}_{T_{1}, T_{2}}^{\min } \mid T_{1} \neq T_{2}, T_{1}, T_{2} \text { are generic }\right\} .
$$

Let two generic topologies $T_{1} \neq T_{2}$ be fixed. Let $P \in \mathcal{A}_{T_{1}, T_{2}}^{\min }$. By Lemma 2.4 .8 we have $\mathcal{A}_{T_{1}, T_{2}}=$ $X_{1} \cup X_{2} \cup \ldots$, where $X_{1}, X_{2}, \ldots$ are connected analytic submanifolds of an open subset $U \subset \mathbb{P}_{2}$. Therefore, by Corollary 2.4 .2

$$
\mathcal{A}_{T_{1}, T_{2}}^{\min } \subset \bigcup_{i: \operatorname{dim} X_{i} \leq 2 n-1} X_{i} .
$$

It follows that $\operatorname{dim} \mathcal{A}_{T_{1}, T_{2}}^{\min } \leq 2 n-1$. Since $T_{1}, T_{2}$ were arbitrary, we conclude that $\operatorname{dim} \mathcal{A} \leq 2 n-1$.

Remark 2.4.1. A reasonable question would be if the dimension of the whole $\mathcal{A}_{T_{1}, T_{2}}$ (not only $\mathcal{A}_{T_{1}, T_{2}}^{\min }$ ) is less or equal to $2 n-1$. We claim that it can be proven in a similar way we prove it for $\mathcal{A}_{T_{1}, T_{2}}^{\min }$. Indeed, due to Lemma 2.4 .8 it is enough to prove that the interior of $\mathcal{A}_{T_{1}, T_{2}}$ is empty. If $P \in \operatorname{Int} \mathcal{A}_{T_{1}, T_{2}}$ and $S_{1}(P)$ and $S_{2}(P)$ are the corresponding embeddings of $T_{1}, T_{2}$, then by Lemma 2.4.9 trees $S_{1}(P)$ and $S_{2}(P)$ are codirected. Our claim is now that Theorem 2.4.1 can still be applied in this case to say that $S_{1}(P)$ and $S_{2}(P)$ coincide. Note that $S_{1}(P)$ and $S_{2}(P)$ are not necessarily locally minimal networks as we allow them to have angles less than $\frac{2 \pi}{3}$ at terminals of degree 2. Nevertheless, the proof of Theorem 2.4.1 given by Oblakov [96] still applies to this case.

### 2.5 Steiner trees in real analytic Riemannian manifolds

A question on the uniqueness of Steiner trees in a Riemannian manifold was raised in [34].
We should not expect that Theorem 2.1.1 can be directly generalized to the case when $\mathbb{R}^{2}$ is replaced with an arbitrary manifold $M$ (cf. Section 2.5.4). Nevertheless, if $M$ is a real analytic manifold, then we still can expect that the set of ambiguous configurations of terminals either has a non-empty interior, or dimension strictly less than the set of all configurations of terminals.

The aim of this section is to build a similar framework to the one used in the proof of Theorem 2.1.1 in Section 2.4 .2 in the case of arbitrary real analytic manifold $M$. Using this framework we reduce the alternative stated above to Conjecture 2.5.1 about analytic sets.

### 2.5.1 Realization space $\mathcal{R}_{\text {geoemb }}$ for an arbitrary metric space

Let us begin with rephrasing Problem 2.0.1 with an arbitrary proper metric space $M$ instead of $\mathbb{R}^{d}$.
Problem 2.5.1. Let $M$ be a metric space. For a given finite set $Q=\left\{p_{1}, \ldots, p_{n}\right\} \subset M$ find a connected set $\mathcal{S}$ with minimal length (one-dimensional Hausdorff measure) containing $Q$.

Due to the following theorem, solutions to Problem 2.5.1 still lie among geodesically embedded trees:

Theorem 2.5.1 (Paolini-Stepanov, [99]). Assume that $M$ is proper (i.e. all closed balls in $M$ are compact) and pathwise connected. Then a solution to Problem 2.5.1 exists. Moreover, for any solution $\mathcal{S t}(Q)$ the following statements hold:
(i) $\mathcal{S t}$ is compact;
(ii) St contains no embedded loops (homeomorphic images of $\mathbb{S}^{1}$ );
(iii) $\mathcal{S} t \backslash Q$ has a finite number of connected components, each component has strictly positive length, and the closure of each component is a finite geodesic embedded graph with endpoints in $Q$;
(iv) the closure of every connected component of $\mathcal{S} t \backslash Q$ is a topological tree with endpoints in $Q$, and all the branching points having finite degree.

This theorem motivates the following definition. Given a positive integer $n$ and a topology $T$, define

$$
\begin{equation*}
\mathcal{R}_{\text {geoemb }}(T)=\{f: T \rightarrow M \mid f \text { is an embedding which maps all edges to (shortest) geodesics }\} / \sim, \tag{2.4}
\end{equation*}
$$

where $f_{1} \sim f_{2}$ if and only if $f_{1}(v)=f_{2}(v)$ for any labelled vertex $v$ and $f_{1}(T)=f_{2}(T)$ as subsets of $M$. Note that this $\mathcal{R}_{\text {geoemb }}(T)$ is a straightforward generalization of $\mathcal{R}_{\text {geoemb }}(T)$ introduces in Section 2.4.2. Let $\mathcal{K}(M)$ be the set of all compact subsets of $M$ endowed with the Hausdorff metric. Introduce two maps: first, $\mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathcal{K}(M)$ which maps $f$ to $f(T)$, second, $P_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow M^{n}$ which sends $f$ to the collection $\left(f\left(v_{1}\right), \ldots, f\left(v_{n}\right)\right)$ of images of $n$ terminals. The topology on $\mathcal{R}_{\text {geoemb }}(T)$ is defined to be the pullback of the product topology under the map $\mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathcal{K}(M) \times M^{n}$ given by the two maps above (the map $P_{T}$ is added to keep track of the enumeration of terminals).

### 2.5.2 Manifold structure on $\mathcal{R}_{\text {geoemb }}$

Let us now assume that $M$ is a connected real analytic manifold with a Riemannian metric $d$ which depends analytically on the point of $M$. We define the intrinsic metric $d_{i n}$ on $M$ as usual; note that $\left(M, d_{i n}\right)$ is a proper metric space. Given a point $p$, denote by $\exp _{p}$ the exponential map defined with respect to $d$; since we have not required $\left(M, d_{i n}\right)$ to be complete, $\exp _{p}$ is defined only for an open subset of the tangent space $T_{p} M$. Set

$$
\widetilde{T_{p} M}=\left\{w \in T_{p} M \mid w \neq 0 \text { and } \exp _{p}(w) \text { is defined }\right\}
$$

Denote by $T M$ the tangent bundle of the manifold $M$. Elements of $T M$ are parameterized by pairs $(p, w)$ where $p \in M$ and $w \in T_{p} M$. Let $\widetilde{T M} \subset T M$ be the union of $\widetilde{T_{p} M}$ over all $p \in M$. Then $\widetilde{T M}$ is an open subset of $T M$ and $\exp : \widetilde{T M} \rightarrow T M$, given by $(p, w) \mapsto\left(\exp _{p}(w), \operatorname{dexp}_{p}(w)\right)$ on each fiber, is a real analytic map mapping $\widetilde{T M}$ onto its image in $T M$ diffeomorphically (see [26, Section 8]). Moreover, this diffeomorphism is analytic at any point.

Let us show that for any topology $T$ the set $\mathcal{R}_{\text {geoemb }}(T)$ has a natural structure of an analytic manifold. Let $e_{1}, \ldots, e_{m}$ be the edges of $T$ and $\vartheta$ be an arbitrary orientation on edges of $T$. Define the $\operatorname{map} \varphi_{\vartheta}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow(\widetilde{T M})^{m}$ by

$$
\varphi_{\vartheta}(f)=\left(\left(f\left(o\left(e_{1}\right)\right), w_{1}\right), \ldots,\left(f\left(o\left(e_{m}\right)\right), w_{m}\right)\right) \in(\widetilde{T M})^{m}
$$

where $o\left(e_{k}\right)$ is the origin of $e_{k}$ oriented according to $\vartheta$ and $w_{k}$ is the tangent vector at $o\left(e_{k}\right)$ such that the geodesic $f\left(e_{k}\right)$ is given by $\left\{\exp _{o\left(e_{k}\right)}\left(t w_{k}\right)\right\}_{0 \leq t \leq 1}$.
Lemma 2.5.1. The following statements hold true:
(i) For any orientation $\vartheta$, the map $\varphi_{\vartheta}$ is a homeomorphism between $\mathcal{R}_{\text {geoemb }}(T)$ and a smooth real analytic submanifold $\mathcal{U}_{\vartheta} \subset(\widetilde{T M})^{m}$ of dimension $(m+1) \cdot \operatorname{dim} M$.
(ii) The morphism $P_{T} \circ \varphi_{\vartheta}^{-1}: \mathcal{U}_{\vartheta} \rightarrow M^{n}$ is analytic.
(iii) Given two orientations $\vartheta_{1}, \vartheta_{2}$, the map $\varphi_{\vartheta_{2}} \circ \varphi_{\vartheta_{1}}^{-1}: \mathcal{U}_{\vartheta_{1}} \rightarrow \mathcal{U}_{\vartheta_{2}}$ is analytic.

Proof. We first prove (i). Note that $\varphi_{\vartheta}$ is injective. Let us show now that $\mathcal{U}_{\vartheta}=f\left(\mathcal{R}_{\text {geoemb }}(T)\right)$ is a real analytic submanifold. In fact we have

$$
\begin{equation*}
\mathcal{U}_{\vartheta}=\left\{\left(\left(x_{1}, w_{1}\right), \ldots,\left(x_{m}, w_{m}\right)\right) \in(\widetilde{T M})^{m} \quad \mid \exp _{x_{i}}\left(w_{i}\right)=x_{j} \text { if } x_{j}=\operatorname{tail}\left(e_{i}\right), x_{i}=x_{j} \text { if } o\left(e_{i}\right)=o\left(e_{j}\right)\right\} \tag{2.5}
\end{equation*}
$$

where tail $\left(e_{i}\right)$ is the tail of $e_{i}$ oriented according to $\vartheta$. It follows that $\mathcal{U}_{\vartheta}$ is an analytic subset of $(\widetilde{T M})^{m}$. The smoothness easily follows from the fact that exp is a diffeomorphism.

Finally, we have that the inverse mapping $\varphi_{\vartheta}^{-1}: \mathcal{U}_{\vartheta} \rightarrow \mathcal{R}_{\text {geoemb }}(T)$ is continuous because $\exp$ is continuous. We conclude that $\varphi_{\vartheta}$ is continuous because $\mathcal{U}_{\vartheta}$ is locally compact.

The items (ii) and (iii) follow easily from the fact that $\exp$ is analytic.

From Lemma 2.5.1 we see that $\mathcal{R}_{\text {geoemb }}(T)$ has a natural structure of a real analytic manifold.
Lemma 2.5.2. The map $P_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow M^{n}$ is analytic and $P_{T}\left(\mathcal{R}_{\text {geoemb }}(T)\right)$ is open in $M^{n}$. The differential $\mathrm{d} P$ has the maximal rank at any point. In particular, the fiber of $P$ is smooth and has the dimension $\operatorname{dim} \mathcal{R}_{\text {geoemb }}(T)-n \operatorname{dim} M$.

Proof. Choose an arbitrary orientation $\vartheta$ on the edges of $T$ having the property that any terminal is an origin of some edge. The lemma now follows from Lemma 2.5.1 and the description (2.5) of $\mathcal{U}_{\vartheta}$.

### 2.5.3 The subvariety $\mathcal{R}_{\text {locmin }}$ of $\mathcal{R}_{\text {geoemb }}$

Define the length function $L_{T}: \mathcal{R}_{\text {geoemb }}(T) \rightarrow \mathbb{R}$ by

$$
L_{T}(f)=\mathcal{H}^{1}(f(T)) .
$$

Let $\mathcal{R}_{\text {locmin }}(T) \subset \mathcal{R}_{\text {geoemb }}(T)$ be the subset given by

$$
\mathcal{R}_{\text {locmin }}(T)=\left\{f \in \mathcal{R}_{\text {geoemb }}(T) \mid f \text { is a local minimum of } L_{T} \text { on the set } P_{T}^{-1}\left(P_{T}(f)\right)\right\} .
$$

Note that, given the orientation $\vartheta$ on the edges of $T$ we have

$$
L \circ \varphi_{\vartheta}^{-1}\left(\left(x_{1}, w_{1}\right), \ldots,\left(x_{m}, w_{m}\right)\right)=\left|w_{1}\right|+\cdots+\left|w_{m}\right|
$$

(cf. Lemma 2.5.1), hence $L_{T}$ is an analytic function on $\mathcal{R}_{\text {geoemb }}(T)$. Recall that due to Lemma 2.5.2, the fibers of $P_{T}$ are smooth. Define the vertical gradient of the function $L_{T}$ at a point $f \in \mathcal{R}_{\text {geoemb }}$ to be the restriction of $\nabla L_{T}$ to the tangent space to the fiber $P_{T}^{-1}\left(P_{T}(f)\right)$ of $P_{T}$ over $f$. Define

$$
\mathcal{R}_{\text {lencrit }}(T)=\left\{f \in \mathcal{R}_{\text {geoemb }} \mid \text { vertical gradient of } L_{T} \text { at } f \text { is zero }\right\} .
$$

Clearly, $\mathcal{R}_{\text {locmin }}(T) \subset \mathcal{R}_{\text {lencrit }}(T)$; note that if $M=\mathbb{R}^{2}$ with the Euclidean metric, then the converse inclusion is also true, but in general we do not have equality between these two sets. We expect nevertheless that $\mathcal{R}_{\text {locmin }}(T)$ is a subanalytic subset of $\mathcal{R}_{\text {geoemb }}(T)$. By Lemma 2.5.2, this statement would immediately follow from the following assertion:

Conjecture 2.5.1. Assume that $U \subset \mathbb{R}^{m}$ and $V \subset \mathbb{R}^{k}$ are open subsets, $f: U \times V \rightarrow \mathbb{R}$ is real analytic. Put

$$
A=\{(x, y) \in U \times V \mid x \text { is a local minimum for } f(\cdot, y) \text { on } U \text {, if } y \text { is fixed }\} .
$$

Then $A$ is a semianalytic subset of $U \times V$.
Assume that for any $T$ the set $\mathcal{R}_{\text {locmin }}(T)$ is subanalytic for a moment. In this case we immediately have the following

Proposition 2.5.1. Let $T_{1}, T_{2}$ be two topologies with $n$ terminals and assume that

$$
\begin{gathered}
\mathcal{A}_{T_{1}, T_{2}}=\left\{P \in M^{n} \quad \mid \exists f_{1} \in \mathcal{R}_{\text {locmin }}\left(T_{1}\right), f_{2} \in \mathcal{R}_{\text {locmin }}\left(T_{2}\right):\right. \\
\left.P_{T_{1}}\left(f_{1}\right)=P_{T_{2}}\left(f_{2}\right)=P \text { and } L_{T_{1}}\left(f_{1}\right)=L_{T_{2}}\left(f_{2}\right)\right\},
\end{gathered}
$$

Then either $\mathcal{A}_{T_{1}, T_{2}}$ has a non-empty interior, or the Hausdorff dimension of $\mathcal{A}_{T_{1}, T_{2}}$ is strictly less than $n \operatorname{dim} M$.

Note that both alternatives in Proposition 2.5.1 can occur, see Section 2.5.4.
Proof. The proof essentially follows proof of Lemma 2.4.8. We consider the fiber product

$$
\mathcal{R}_{\text {locmin }}\left(T_{1}\right) \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {locmin }}\left(T_{2}\right)
$$

which is subanalytic due to Lemma 2.4.4. Inside $\mathcal{R}_{\text {locmin }}\left(T_{1}\right) \times_{P} \mathcal{R}_{\text {locmin }}\left(T_{2}\right)$ we consider the set $\tilde{\mathcal{A}}_{T_{1}, T_{2}}$ cut out by the equation $L_{T_{1}}\left(f_{1}\right)=L_{T_{2}}\left(f_{2}\right)$. Then $\tilde{\mathcal{A}}_{T_{1}, T_{2}}$ is subanalytic and $\mathcal{A}_{T_{1}, T_{2}}=P\left(\tilde{\mathcal{A}}_{T_{1}, T_{2}}\right)$, where $P=P_{T_{1}}: \mathcal{R}_{\text {locmin }}\left(T_{1}\right) \times_{P_{T_{1}}=P_{T_{2}}} \mathcal{R}_{\text {locmin }}\left(T_{2}\right) \rightarrow M^{n}$ is the projection. The proposition now follows from Corollary 2.4.1.

Let us finalize this section with a short discussion on Conjecture 2.5.1. Note that if $k=0$ so that $V$ is just a point, then the statement in the conjecture is straightforward (see also [41 for a much more general statement). Indeed, define first

$$
Z=\{x \in U \mid \nabla f(x)=0\},
$$

Clearly, $A \subset Z$. Using [9, Proposition 2.10] we find that each $(x, y) \in U$ has a neighborhood $U^{\prime}$ such that $Z \cap U^{\prime}=Z_{1} \cup \cdots \cup Z_{k}$ where each $Z_{j}$ is a connected semianalytic subset of $U^{\prime}$ which is also an analytic submanifold of $U$. In particular, $f$ restricted to each $Z_{j}$ is a constant, say, $C_{j}$. Define

$$
B_{j}=\overline{\left\{x \in U \mid f(x)<C_{j}\right\}} .
$$

Then $B_{j}$ is a semianalytic subset (as the closure of a semianalytic subset is semianalytic [9, Proposition 2.10]). We have

$$
A \cap U^{\prime}=\bigcup_{j=1}^{k}\left(Z_{j} \backslash B_{j}\right)
$$

Since $x$ was chosen arbitrarily, this proves that $A$ is semianalytic.
With some more involved arguments using Weierstrass preparation theorem we can prove the conjecture when $\operatorname{dim} U=1$ and $V$ is arbitrary, but the general case remains unclear to us.

### 2.5.4 Example of $\mathcal{A}_{T_{1}, T_{2}}$ with non-empty interior

In this subsection we construct an example of a Riemann surface $M$ with a locally flat metric such that there exist two topologies $T_{1}, T_{2}$ on 8 terminals for which $\mathcal{A}_{T_{1}, T_{2}}$ has non-empty interior (see Proposition 2.5 .1 for the definition of $\mathcal{A}_{T_{1}, T_{2}}$ ).

Let $T_{1}$ and $T_{2}$ be the two topologies introduced by Ivanov and Tuzhilin in [65, Fig. 1], see Figure 2.13. Let $T_{1}$ be the topology of the tree drawn by solid lines for certainty. We fix a set of terminals $x=\left(x_{1}, \ldots, x_{8}\right) \in \mathbb{P}_{2}$ and fix immersions $f_{i}: T_{i} \rightarrow \mathbb{R}^{2}$ into plane as on Figure 2.13. Following [65] we have the following

Lemma 2.5.3. For any configuration $\left(\tilde{x}_{1}, \ldots, \tilde{x}_{8}\right)$ sufficiently close to $\left(x_{1}, \ldots, x_{8}\right)$ the corresponding immersions $\tilde{f}_{1}, \tilde{f}_{2}$ of the topologies $T_{1}, T_{2}$ realizing them as locally minimal trees are codirected and have the same length.

Proof. Follows immediately from the Melzak algorithm. Note that being codirected at some point implies being codirected in a neighborhood, which in turn is equivalent to the equality of length (cf. Lemma 2.4.9.

Let

$$
T=T_{1} \sqcup T_{2} / \text { glued along terminal half-edges },
$$

i.e. we glue $T_{1}$ and $T_{2}$ along the portion of edges emanating from $x_{1}, \ldots, x_{8}$ which coincide on the picture (other intersection points on the picture are not glued). Let $f: T \rightarrow \mathbb{R}^{2}$ be the corresponding immersion.

Note that $T$ inherits a metric from $f(T)$ : use Euclidean metric on $f(T)$ to measure distances along an edge of $T$, and use the inner metric to measure the distance between two arbitrary points.


Figure 2.13: Two locally minimal trees with self-intersections
Fix an $\varepsilon>0$ and for each $t \in T$ define the surface $M_{t}$ to be a copy the $\varepsilon$-neighborhood of $f(t)$ in $\mathbb{R}^{2}$. The map $f$ extends to the embedding $F_{t}: M_{t} \rightarrow \mathbb{R}^{2}$. Given $t_{1}, t_{2} \in T$ on the distance at most $10 \varepsilon$ from each other, glue $M_{t_{1}}$ and $M_{t_{2}}$ such that $F_{t_{1}}=F_{t_{2}}$. As a result we obtain Riemann surface $M$ containing $T$, and an immersion $F: M \rightarrow \mathbb{R}^{2}$ such that the image $F(M)$ is the $\varepsilon$-neighborhood of $f(T)$. Note that $M$ is endowed with a locally flat metric for which $F$ is a local isometry.

Let $g_{i}: T_{i} \rightarrow M, i=1,2$, be the natural maps. Note that $f_{i}=F \circ g_{i}$ and $g_{i}$ are injective. Denote by $X_{1}, \ldots, X_{8} \in M$ the points such that $F\left(X_{i}\right)=x_{i}$. The following proposition is straightforward.

Proposition 2.5.2. Trees $g_{1}\left(T_{1}\right), g_{2}\left(T_{2}\right)$ are locally minimal on $M$ and $\left(g_{1}, g_{2}\right) \in \operatorname{Int} \mathcal{A}_{T_{1}, T_{2}}$.
We are not able to extend this example on minimal trees.

## Chapter 3

## Gilbert-Steiner problem

Here we follow paper [17].

### 3.1 Basics

One of the first models for branched transport was introduced by Gilbert [54]. The difference with the optimal transportation problem is that the extra geometric points may be of use; this explains the naming in honor of Steiner. Sometimes it is also referred to as optimal branched transport; a large part of book [8] is devoted to this problem. Let us proceed with the formal definition.
Definition 3.1.1. Let $\mu^{+}$, $\mu^{-}$be two finite measures on a metric space $(X, \rho(\cdot, \cdot))$ with finite supports such that total masses $\mu^{+}(X)=\mu^{-}(X)$ are equal. Let $V \subset X$ be a finite set containing the support of the signed measure $\mu^{+}-\mu^{-}$, the elements of $V$ are called vertices. Further, let $E$ be a finite collection of unordered pairs $\{x, y\} \subset V$ which we call edges. So, $(V, E)$ is a simple undirected finite graph. Assume that for every $\{x, y\} \in E$ two non-zero real numbers $m(x, y)$ and $m(y, x)$ are defined so that $m(x, y)+m(y, x)=0$. This data set is called $a\left(\mu^{+}, \mu^{-}\right)$-flow if

$$
\mu^{+}-\mu^{-}=\sum_{\{x, y\} \in E} m(x, y) \cdot\left(\delta_{y}-\delta_{x}\right)
$$

where $\delta_{x}$ denotes a delta-measure at $x$ (note that the summand $m(x, y) \cdot\left(\delta_{y}-\delta_{x}\right)$ is well-defined in the sense that it does not depend on the order of $x$ and $y$ ).

Let $C:[0, \infty) \rightarrow[0, \infty)$ be a cost function. The expression

$$
\sum_{\{x, y\} \in E} C(|m(x, y)|) \cdot \rho(x, y)
$$

is called the Gilbert functional of the $\left(\mu^{+}, \mu^{-}\right)$-flow.
The Gilbert-Steiner problem is to find the flow which minimizes the Gilbert functional with cost function $C(x)=x^{p}$, for a fixed $p \in(0,1)$; we call a solution minimal flow.

Vertices from supp $\left(\mu^{+}\right) \backslash \operatorname{supp}\left(\mu^{-}\right)$are called terminals. A vertex from $V \backslash \operatorname{supp}\left(\mu^{+}\right) \backslash \operatorname{supp}\left(\mu^{-}\right)$ is called a branching point. Formally, we allow a branching point to have degree 2, but such points may be easily eliminated.

Local structure in the Gilbert-Steiner problem was discussed in [8, and the paper [85] deals with planar case. A local picture around a branching point $b$ of degree 3 is clear due to the initial paper of Gilbert. Similarly to the finding of the Fermat-Torricelli point in the celebrated Steiner problem one can determine the angles around $b$ in terms of masses (see Lemma 3.2.1).

Theorem 3.1.1 (Lippmann-Sanmartín-Hamprecht [85], 2022). A solution of the planar GilbertSteiner problem has no branching point of degree at least 5 .

A similar setup was independently considered in Minkowski spaces by Volz, Brazil, Ras, Swanepoel and Thomas 121. In a Euclidean space (of an arbitrary dimension) they obtained that a degree of a branching is at most 3 provided that $c$ is a concave monotone function, the function $\left(c^{2}\right)^{\prime}$ is convex and $c(0)>0$. Note that for $p \in(1 / 2,1)$ the convexity condition fails.

The goal of this chapter is to give some conditions on a cost function under which all branching points in a planar solution have degree 3. They are slightly stronger than the Schoenberg [115] conditions of the embedding of the metric of the form $\rho(x, y):=f(x-y)$ to a Hilbert space. In particular, this covers the case of the standard cost function $x^{p}, 0<p<1$. The following main theorem is the part of a more general Theorem 3.3.1.

Theorem 3.1.2. A solution of the planar Gilbert-Steiner problem has no branching point of degree at least 4.

### 3.2 Preliminaries

We need the following lemmas.
Lemma 3.2.1 (Folklore). Let $P Q R$ be a triangle and $w_{1}, w_{2}, w_{3}$ be non-negative reals. For every point $X \in \mathbb{R}^{2}$ consider the value

$$
L(X):=w_{1} \cdot|P X|+w_{2} \cdot|Q X|+w_{3} \cdot|R X| .
$$

Then
(i) a minimum of $L(X)$ is achieved at a unique point $X_{\text {min }}$;
(ii) if $X_{\text {min }}=P$ then $w_{1} \geq w_{2}+w_{3}$ or there is a triangle $\Delta$ with sides $w_{1}, w_{2}, w_{3}$ and $\angle P$ is at least the outer angle between $w_{2}$ and $w_{3}$ in $\Delta$.

Hereafter the metric space is the Euclidean plane $\mathbb{R}^{2}$.
The following concept only slightly changes from that of Schoenberg [115], introduced for describing which metrics of the form $\rho(x, y)=f(x-y)$ on the real line can be embedded to a Hilbert space.

Definition 3.2.1. Let $\lambda$ be a Borel measure on $\mathbb{R}$ for which

$$
\begin{equation*}
\int \min \left(x^{2}, 1\right) d \lambda(x)<\infty \tag{3.1}
\end{equation*}
$$

Assume additionally that the support of $\lambda$ is uncountable. A function $f: \mathbb{R}_{\geqslant 0} \rightarrow \mathbb{R}_{\geqslant 0}$ of the form

$$
\begin{equation*}
f(t)=\sqrt{\int \sin ^{2}(t x) d \lambda(x)}=\frac{1}{2}\left\|e^{2 i t x}-1\right\|_{L^{2}(\lambda)} \tag{3.2}
\end{equation*}
$$

is called admissible.

The only difference with [115] is that we require that the support of the measure $\lambda$ is uncountable which guarantees that the corresponding embedding has full dimension (see below).

Remark 3.2.1. As $\lambda$ is a Borel measure, a continuous function $\sin ^{2}(t x)$ is $\lambda$-measurable. Under conditions (3.1), the integral in (3.2) is finite, so $f(t)<\infty$ for all $t \geqslant 0$.

Further we are going to consider only admissible cost functions. Note that admissibility implies some properties one may expect from a cost function. In particular, $f(0)=0$ and $f$ is subadditive: for non-negative $t, s$ we have

$$
\begin{aligned}
f(t)+f(s) & =\frac{1}{2}\left\|e^{2 i t x}-1\right\|_{L^{2}(\lambda)}+\frac{1}{2}\left\|e^{2 i s x}-1\right\|_{L^{2}(\lambda)} \\
& =\frac{1}{2}\left\|e^{2 i t x}-1\right\|_{L^{2}(\lambda)}+\frac{1}{2}\left\|e^{2 i(s+t) x}-e^{2 i t x}\right\|_{L^{2}(\lambda)} \geq \frac{1}{2}\left\|e^{2 i(t+s) x}-1\right\|_{L^{2}(\lambda)}=f(t+s) .
\end{aligned}
$$

On the other hand it does not imply monotonicity (for instance, if $\operatorname{supp} \lambda \subset[0.9,1.1]$ then $f(\pi)<$ $f(\pi / 2)$ ).

Hereafter $L^{2}(\lambda)$ for a measure $\lambda$ on $\mathbb{R}$ is understood as a real Hilbert space of complex-valued square summable w.r.t. $\lambda$ functions (strictly speaking, of classes of equivalences of such functions modulo coincidence $\lambda$-almost everywhere).

Proposition 3.2.1. If $\lambda$ is a Borel measure on $\mathbb{R}$ with uncountable support such that

$$
\int \min \left(x^{2}, 1\right) d \lambda(x)<\infty
$$

then any finite collection of functions of the form $e^{i a x}-1, a \in \mathbb{R}$, is affinely independent in $L^{2}(\lambda)$.
Proof. Assume the contrary. Then there exist distinct real numbers $a_{1}, \ldots, a_{n}$ and non-zero real coefficients $t_{1}, \ldots, t_{n}$ such that $\sum t_{j}=0$ and $\sum t_{j}\left(e^{i a_{j} x}-1\right)=0 \lambda$-almost everywhere. But the analytic function $\sum t_{j}\left(e^{i a_{j} x}-1\right)$ is either identically zero, or has at most countably many (and isolated) zeroes. In the latter case, it is not zero $\lambda$-almost everywhere, since the support of $\lambda$ is uncountable. The former case is not possible: indeed, if $\sum t_{j} e^{i a_{j} x} \equiv 0$, then taking the Taylor expansion at 0 we get $\sum t_{j} a_{j}^{k}=0$ for all $k=0,1,2, \ldots$. Therefore $\sum t_{j} W\left(a_{j}\right)=0$ for any polynomial $W$. Choosing $W(t)=\prod_{j=2}^{n}\left(t-a_{j}\right)$ we get $t_{1}=0$, a contradiction.

One can see from the proof that the condition on uncountability of the support may be weakened.
Lemma 3.2.2. Let $C$ be an admissible cost function. For real numbers $m_{1}$ and $m_{2}$ we define $h\left(m_{1}, m_{2}\right)$ as the value of the outer angle between $C\left(\left|m_{1}\right|\right)$ and $C\left(\left|m_{2}\right|\right)$ in the triangle with sides $C\left(\left|m_{1}\right|\right), C\left(\left|m_{2}\right|\right), C\left(\left|m_{1}+m_{2}\right|\right)$ (it exists by the discussion before Proposition 3.2.1). Suppose that $O V_{1}, O V_{2}$ are edges in a minimal flow with masses $m_{1}$ and $m_{2}$. Then the angle between $O V_{1}$ and $O V_{2}$ is at least $h\left(m_{1}, m_{2}\right)$.

Proof. Assume the contrary, then by Lemma 3.2.1 with $P=O, Q=V_{1}, R=V_{2}, w_{1}=C\left(\left|m_{1}\right|\right), w_{2}=$ $C\left(\left|m_{2}\right|\right), w_{3}=C\left(\left|m_{1}+m_{2}\right|\right)$ we have $X_{\min } \neq O$. Then we can replace $\left[O V_{1}\right] \cup\left[O V_{2}\right]$ with $\left[X_{\min } O\right] \cup$ $\left[X_{\text {min }} V_{1}\right] \cup\left[X_{\text {min }} V_{2}\right]$ with the corresponding masses in our flow; this contradicts the minimality of the flow.

Lemma 3.2.3. For $0<p<1$, the function $f(x)=x^{p}$ is admissible.

Proof. Consider the measure $d \lambda=x^{-2 p-1} d x$ on $[0, \infty)$. Then $\int_{0}^{\infty} \min \left(x^{2}, 1\right) d \lambda<\infty$ and for $t>0$ we have

$$
\int_{0}^{\infty} \sin ^{2}(t x) d \lambda(x)=\int_{0}^{\infty} \sin ^{2}(t x) x^{-2 p-1} d x=t^{2 p} \int_{0}^{\infty} \sin ^{2} y y^{-2 p-1} d y
$$

thus the measure $\lambda$ multiplied by an appropriate positive constant proves the result.
Example 3.2.1. For another natural choice $d \lambda=4 c e^{-2 c x} d x, c>0$, we get an admissible function $f(t)=t / \sqrt{t^{2}+c^{2}}$.

The following lemma is essentially well-known, but for the sake of completeness and for covering degeneracies and the equality cases we provide a proof.

Lemma 3.2.4. Let $X$ be a finite-dimensional Euclidean space, let the points $A_{0}, A_{1}, A_{2}, \ldots, A_{n-1}$, $A_{n}=A_{0}, A_{n+1}=A_{1}$ in $X$ be chosen so that $A_{i} \neq A_{i+1}$ for all $i=1,2, \ldots, n$. Denote $\varphi_{i}:=$ $\pi-\angle A_{i-1} A_{i} A_{i+1}$ for $i=1,2, \ldots, n$. Then $\sum \varphi_{i} \geqslant 2 \pi$, and if the equality holds then the points $A_{1}, \ldots, A_{n}$ belong to the same two-dimensional affine plane.

Proof. Let $u$ be a randomly chosen unit vector in $X$ (with respect to a uniform distribution on the sphere). For $j=1,2, \ldots, n$ denote by $U(j)$ the following event: $\left\langle u, A_{j}\right\rangle=\max _{1 \leqslant i \leqslant n}\left\langle u, A_{i}\right\rangle$, where $\langle\cdot, \cdot\rangle$ denotes the inner product in $X$; and by $V(j)$ the event $\left\langle u, A_{j}\right\rangle=\max _{j-1 \leqslant i \leqslant j+1}\left\langle u, A_{i}\right\rangle$. Obviously, $\operatorname{prob} U(j) \leqslant \operatorname{prob} V(j)$. Also, prob $V(j)=\frac{\varphi_{j}}{2 \pi}$, since the set of directions of $u$ for which $V(j)$ holds is the dihedral angle of measure $\varphi_{j}$. Thus, since always at least one event $U(j)$ holds, we get

$$
1 \leqslant \sum_{j=1}^{n} \operatorname{prob} U(j) \leqslant \sum_{j=1}^{n} \operatorname{prob} V(j)=\frac{1}{2 \pi} \sum_{j=1}^{n} \varphi_{j}
$$

This proves the inequality. It remains to prove that it is strict assuming that not all the points belong to a two-dimensional plane. Note that if every three consecutive points $A_{j-1}, A_{j}, A_{j+1}$ are collinear, then all the points $A_{1}, \ldots, A_{n}$ are collinear that contradicts to our assumption. If $A_{j-1}, A_{j}, A_{j+1}$ are not collinear, denote by $\alpha$ the two-dimensional plane they belong to. There exists $i$ for which $A_{i} \notin \alpha$. Then $\operatorname{prob} U(j)<\operatorname{prob} V(j)$, since there exist planes passing through $A_{j}$ which separate the triangle $A_{j-1} A_{j} A_{j+1}$ and the point $A_{i}$, and the measure of directions of such planes is strictly positive. Therefore, our inequality is strict.

### 3.3 Main result

Theorem 3.3.1. Let $\mu^{+}, \mu^{-}$be two measures with finite support on the Euclidean plane $\mathbb{R}^{2}$, and assume that the cost function $C$ is admissible. Then if $a\left(\mu^{+}, \mu^{-}\right)$-flow has a branching point of degree at least 4, then there exists a $\left(\mu^{+}, \mu^{-}\right)$-flow with strictly smaller value of Gilbert functional.

Proof. Assume the contrary. Let $O$ be a branching point, $O V_{1}, O V_{2}, \ldots, O V_{k}, k \geqslant 4$, be the edges incident to $O$, enumerated counterclockwise. Further the indices of $V_{i}$ 's are taken modulo $k$, so that $V_{1}=V_{k+1}$ etc. Denote $m_{i}=m\left(O V_{i}\right)$, then by the definition of flow we get $\sum m_{i}=0$. By Lemma 3.2.2, $\angle V_{i} O V_{i+1} \geqslant h\left(m_{i}, m_{i+1}\right)$.

Consider the functions $A_{j}(x):=e^{i\left(m_{1}+\ldots+m_{j}\right) x}-1$ for $j=1,2, \ldots$ (here $i$ is the imaginary unit). Then $\sum m_{j}=0$ yields that $A_{j+k} \equiv A_{j}$ for all $j>0$.

Since the cost function $C(t)$ is admissible, there exists a Borel measure $\lambda$ on $\mathbb{R}$ with uncountable support such that $\int \min \left(x^{2}, 1\right) d \lambda(x)<\infty$ and

$$
C(t)=\sqrt{\int 4 \sin ^{2} \frac{t x}{2} d \lambda(x)}
$$

Using the identity $\left|e^{i a}-e^{i b}\right|^{2}=4 \sin ^{2} \frac{a-b}{2}$ for real $a, b$ we note that for $j, s>0$ in the Hilbert space $L^{2}(\lambda)$ we have

$$
\left\|A_{j+s}-A_{j}\right\|^{2}=C\left(\left|m_{j+1}+\ldots+m_{j+s}\right|\right)^{2} .
$$

In particular, the lengths of the sides of the triangle $A_{j-1} A_{j} A_{j+1}$ are equal to $C\left(\left|m_{j}\right|\right), C\left(\left|m_{j+1}\right|\right)$ and $C\left(\left|m_{j}+m_{j+1}\right|\right)$. Therefore $\varphi_{j}:=\pi-\angle A_{j-1} A_{j} A_{j+1}=h\left(m_{j}, m_{j+1}\right)$. By Lemma 3.2.4 we get $\sum \varphi_{j} \geqslant 2 \pi$.

By Lemma 3.2.1. this yields $2 \pi=\sum_{j=1}^{k} \angle V_{j} O V_{j+1} \geqslant \sum \varphi_{j} \geqslant 2 \pi$. Therefore, the equality must take place. Again by Lemma 3.2 .4 it follows that the points $A_{j}$ belong to the same 2-dimensional subspace. But by Proposition 3.2.1, distinct points between $A_{j}$ 's are affinely independent. Therefore, there exist at most three distinct $A_{j}$ 's, and if exactly three, they are not collinear. It is easy to see that the equality $\sum \varphi_{j}=2 \pi$ under these conditions does not hold when $k>3$. A contradiction.

### 3.4 Examples of branching points of degree 4

Let us start with an example in three dimensions. Consider four masses $m_{1}, m_{2}, m_{3}, m_{4}$ of zero sum, such that no two of them give zero sum. Repeat the beginning of the proof of Theorem 3.3.1 to get the simplex $A_{1} A_{2} A_{3} A_{4}$ in 3-dimensional space. Now consider unit edges $O B_{i}$ in $\mathbb{R}^{3}$ with directions $A_{i-1} A_{i}$, $1 \leq i \leq 4$. By the construction the angles between vectors $O B_{i}$ and $O B_{i+1}$ are exactly $h\left(m_{i}, m_{i+1}\right)$. Suppose that angles $\angle B_{1} O B_{3}$ and $\angle B_{2} O B_{4}$ are at least $h\left(m_{1}, m_{3}\right)$ and $h\left(m_{2}, m_{4}\right)$, respectively (for instance, it happens for $\left.C(x)=x^{p}, 1 / 2 \leq p<1, m_{1}=m_{2}=m_{3}=1, m_{4}=-3\right)$.

Then we claim that for a concave monotone admissible cost function $C$ the flow

$$
\sum_{i=1}^{4} m_{i} \cdot\left(\delta_{O}-\delta_{B_{i}}\right)
$$

is a solution of the corresponding Gilbert-Steiner problem.
First, if we fix the graph structure then the position of $O$ is optimal by the following lemma, because the closeness of the polychain $A_{1} A_{2} A_{3} A_{4} A_{1}$ gives exactly (3.3).

Lemma 3.4.1 (Weighted geometric median, [125). Consider different non-collinear points $A, B, C$, $D \in \mathbb{R}^{3}$ and let $w_{1}, w_{2}, w_{3}, w_{4}$ be non-negative reals. Then

$$
L(X):=w_{1} \cdot|A X|+w_{2} \cdot|B X|+w_{3} \cdot|C X|+w_{4} \cdot|D X|
$$

has unique local (and global) minimum satisfying

$$
\begin{equation*}
w_{1}-\overline{e_{A}}+w_{2} \overline{e_{B}^{-}}+w_{3} \overline{e_{C}^{-}}+w_{4} \overline{e_{D}}=0 \tag{3.3}
\end{equation*}
$$

where $\overline{e_{A}}, \overline{e_{B}}, \overline{e_{C}}, \overline{e_{D}}$ are unit vectors codirected with $X A, X B, X C, X D$, respectively.

For a concave monotone cost function $c$ there is an optimal flow without cycles (see Proposition 7.8 in [8]).

Since any two masses have nonzero sum, every flow is connected. Thus every possible competitor has 2 branching points of degree 3. Consider the case in which branching points $U$ and $V$ are connected with $B_{1}, B_{2}$ and $B_{3}, B_{4}$, respectively. By the convexity of length, the Gilbert functional $L(U, V)$ considered on the set of all possible $U$ and $V\left(\mathbb{R}^{3} \times \mathbb{R}^{3}\right)$ is a convex function. Let us show that $U=V=O$ is a local minimum. Indeed, consider $U_{\varepsilon}=O+\varepsilon u$ and $V_{\delta}=O+\delta v$ for arbitrary unit vectors $u, v$ and small positive $\varepsilon, \delta$. By the convexity of length

$$
\begin{gathered}
L\left(U_{\varepsilon}, V_{\delta}\right)-L(O, O) \geq \\
w(U V) \cdot\|\varepsilon u-\delta v\|-\varepsilon\left\langle w_{1} e_{1}, u\right\rangle-\varepsilon\left\langle w_{2} e_{2}, u\right\rangle-\delta\left\langle w_{3} e_{3}, v\right\rangle-\delta\left\langle w_{4} e_{4}, v\right\rangle= \\
w(U V) \cdot\|\varepsilon u-\delta v\|-\varepsilon\left\langle w_{12} e_{12}, u\right\rangle-\delta\left\langle w_{34} e_{34}, v\right\rangle
\end{gathered}
$$

where $w_{12} e_{12}=w_{1} e_{1}+w_{2} e_{2}$ and $w_{34} e_{34}=w_{3} e_{3}+w_{4} e_{4}$ for unit $e_{12}$ and $e_{34}$. By the construction one has $w_{12}=w_{34}=w(U V)$ and $e_{12}+e_{34}=0$, so

$$
L\left(U_{\varepsilon}, V_{\delta}\right)-L(O, O) \geq w(U V) \cdot\left(\|\varepsilon u-\delta v\|-\left\langle e_{12}, \varepsilon u-\delta v\right\rangle\right)
$$

Since $e_{12}$ is unit, the derivative is non-negative for every $u, v$.
The case in which $U$ and $V$ are connected with $B_{2}, B_{3}$ and $B_{4}, B_{1}$, respectively, is completely analogous. In the remaining case ( $U$ is connected with $B_{1}, B_{3}$ and $V$ is connected with $B_{2}, B_{4}$ ) we have $w_{12}=w_{34} \leq w(U V)$ due to $\angle B_{1} O B_{3} \geq h\left(m_{1}, m_{3}\right)$ and $\angle B_{2} O B_{4} \geq h\left(m_{2}, m_{4}\right)$. Thus $U=V=O$ is also a local minimum.

It is known [54] that $L$ has a unique local and global minimum, which finishes the example.
Now proceed with planar examples of 4 -branching for some non-admissible cost-function $C$. Then we may repeat the 3 -dimensional argument starting with planar $A_{1} A_{2} A_{3} A_{4}$.

The simplest way to produce an example is to consider an isosceles trapezoid $A_{1} A_{2} A_{3} A_{4}$ and apply Ptolemy's theorem. This case corresponds to $m_{1}=m_{3}$ and $m_{1}+m_{2}+m_{3}+m_{4}=0$. Then $\left|A_{1} A_{2}\right|=$ $\left|A_{3} A_{4}\right|=C\left(\left|m_{1}\right|\right),\left|A_{2} A_{3}\right|=C\left(\left|m_{2}\right|\right),\left|A_{4} A_{1}\right|=C\left(\left|m_{4}\right|\right)$ and $\left|A_{1} A_{3}\right|=\left|A_{2} A_{4}\right|=C\left(\left|m_{1}+m_{2}\right|\right)$. The existence of such trapezoid means

$$
\begin{equation*}
C\left(\left|m_{1}+m_{2}\right|\right)^{2}=C\left(\left|m_{1}\right|\right)^{2}+C\left(\left|m_{2}\right|\right) \cdot C\left(\left|2 m_{1}+m_{2}\right|\right) \tag{3.4}
\end{equation*}
$$

If we assume that $C$ is monotone and subadditive then (3.4) means that the isosceles trapezoid exists; note that we need values of $C$ only at 4 points.

Now we give an example of a monotone, subadditive and concave cost function with 4 -branching. For this purpose put $m_{1}=m_{2}=m_{3}=1$ and $m_{4}=-3, C(1)=1, C(2)=1.9, C(3)=2.61$; clearly (3.4) holds. Now one can easily interpolate a desired $C$, for instance

$$
C(t)= \begin{cases}t, & t \leq 1 \\ 0.1+0.9 t, & 1<t \leq 2 \\ 0.48+0.71 t, & 2<t \leq 3 \\ 1.11+0.5 t & 3<t\end{cases}
$$

Finally, the inequalities $\pi=\angle B_{1} O B_{3}>h\left(m_{1}, m_{3}\right)$ and $\angle B_{2} O B_{4}>\angle B_{2} O B_{3}=h\left(m_{2}, m_{3}\right)=$ $h\left(m_{2}, m_{4}\right)$ hold.

### 3.5 Open questions

It would be interesting to describe all cost functions for which the conclusion of Theorem 3.3.1 holds.
Now let us focus on the cost function $C(x)=x^{p}$. Having a knowledge that every branching point has degree 3 one can adapt Melzak algorithm [91] from Steiner trees to Gilbert-Steiner problem. The idea of the algorithm is that after fixing the combinatorial structure one can find two terminals $t_{1}, t_{2}$ connected with the same branching point $b$. Then one may reconstruct the solution for $V$ from the solution for $V \backslash\left\{t_{1}, t_{2}\right\} \cup\left\{t^{\prime}\right\}$ for a proper $t^{\prime}$ which depend only on $t_{1}, t_{2}$ (in fact one has to check 2 such $t^{\prime}$ ). When the underlying graph is a matching we finish in an obvious way. Application of this procedure for all possible combinatorial structures gives a slow but mathematically exhaustive algorithm in the planar case.

However there is no known algorithm in $\mathbb{R}^{d}$ for $d>2$ (see Problem 15.12 in [8]). Recall that we have to consider a high-degree branching.

A naturally related problem is to evaluate the maximal possible degree of a branching point in the $d$-dimensional Euclidean space for every $d$. Note that the dependence on the cost function may be very complicated. In particular, an upper bound on the degree which does not depend on $p$ is of interest. It is worth noting that $p<1 / 2$ implies that $h\left(m_{1}, m_{2}\right)>\pi / 2$ for every $m_{1}, m_{2} \neq 0$ and thus the degree is at most $d+1$.

Some other questions are collected in Section 15 of 8] (some of them are solved, in particular Problem 15.1 is solved in [27]).

## Chapter 4

## Maximal distance minimizers

### 4.1 Introduction

This chapter is based on papers [19, 18, 24]. We are interested in the solutions of the following maximal distance minimizer problem.
Problem 4.1.1. For a given compact set $M \subset \mathbb{R}^{d}$ and $l>0$ to find a connected compact set $\Sigma$ of length (one-dimensional Hausdorff measure $\mathcal{H}^{1}$ ) at most $l$ that minimizes

$$
\max _{y \in M} \operatorname{dist}(y, \Sigma)
$$

where dist stands for the Euclidean distance.
It appeared in a very general form by Buttazzo, Oudet and Stepanov in 13 and then in was specified by Miranda, Paolini and Stepanov in [92, 98].

A maximal distance minimizer is a solution of Problem 4.1.1. Such sets can be considered as networks of radiating Wi-Fi cables with a bounded length arriving to each customer (for the set $M$ of customers) at the distance $r$, where such $r$ is the smallest possible.

### 4.1.1 Class of problems

Maximal distance minimization problem could be considered as a particular example of shape optimization problem. A shape optimization problem is a minimization problem where the unknown variable runs over a class of domains; then every shape optimization problem can be written in the form $\min F(\Sigma): \Sigma \in A$ where $A$ is the class of admissible domains and $F()$ is the cost function that one has to minimize over $A$.

So for a given compact set $M$ and positive number $l \geq 0$ let the admissible set $A$ be a set of all closed connected set $\Sigma^{\prime}$ with length constraint $\mathcal{H}^{1}\left(\Sigma^{\prime}\right) \leq l$; and let cost function be the energy $F_{M}(\Sigma)=\max _{y \in M} \operatorname{dist}(y, \Sigma)$. Also $F_{M}(\emptyset):=\infty$.

### 4.1.2 Dual problem

Define the dual problem to Problem 4.1.1 as follows.
Problem 4.1.2. For a given compact set $M \subset \mathbb{R}^{d}$ and $r>0$ to find a connected compact set $\Sigma$ of the minimal length (one-dimensional Hausdorff measure $\mathcal{H}^{1}$ ) such that

$$
\max _{y \in M} \operatorname{dist}(y, \Sigma) \leq r
$$

In a nondegenerate case (i.e. for $F_{M}(\Sigma)>0$ ) the primal and dual problems have the same sets of solutions for the corresponding $r$ and $l$ (see [98]) and hence an equality $F_{M}(\Sigma)=r$ is reached for a minimizer $\Sigma$.

### 4.1.3 The first parallels with average distance minimization problem

Maximal distance minimization problem is somehow similar to another shape optimization problem: average distance minimization problem (see the survey of Lemenant [84]) and it seems interesting to compare the known results and open questions concerning these two problems. In the average distance minimization problem's statement the admissible set $A$ is the same as in Maximal distance minimization problem, but the function $F\left(\Sigma_{a}\right)$ is defined as $\int_{M} A\left(\operatorname{dist}\left(y, \Sigma_{a}\right)\right) d \phi(x)$ where $A: \mathbb{R}^{+} \rightarrow$ $\mathbb{R}^{+}$is a nondecreasing function and $\phi()$ is a finite nonnegative measure with compact nonempty support in $\mathbb{R}^{d}$.

Minimization problems for average distance and maximum distance functionals are used in economics and urban planning with similar interpretations. If it is required to find minimizers under the cardinality constraint $\sharp \Sigma \leq k$, instead of the length and the connectedness constraints, where $k \in \mathbb{N}$ is given and $\sharp$ denotes the cardinality, then the corresponding problems are referred to as optimal facility location problems.

### 4.1.4 Notation

For a given set $X \subset \mathbb{R}^{d}$ we denote by $\bar{X}$ its closure, by $\operatorname{Int}(X)$ its interior and by $\partial X$ its topological boundary.

Let $B_{\rho}(x)$ stand for the open ball of radius $\rho$ centered at a point $x$, and let $B_{\rho}(T)$ be the open $\rho$-neighborhood of a set $T$ i.e.

$$
B_{\rho}(T):=\bigcup_{x \in T} B_{\rho}(x)
$$

(in other words $B_{\rho}(T)$ is Minkowski sum of a ball $B_{\rho}$ centered in the origin and $T$ ). Note that the condition

$$
\max _{y \in M}^{\operatorname{dist}}(y, \Sigma) \leq r
$$

is equivalent to $M \subset \overline{B_{r}(\Sigma)}$.
For given points $b, c$ we use the notation $[b c],[b c)$ and (bc) for the corresponding closed line segment, ray and line respectively.

### 4.1.5 Existence. Absence of loops. Ahlfors regularity and other simple properties

For the both problems existence of solutions is proved easily: according to the classical Blaschke and Gołąb Theorems, the class of admissible sets is compact for the Hausdorff distance and both of the functions (maximal distance and also the average distance) is continuous for this convergence because of the uniform convergence of $x \rightarrow \operatorname{dist}(x, \Sigma)$.

Definition 4.1.1. A closed set $\Sigma$ is said to be Ahlfors regular if there exists some constants $C_{1}$, $C_{2}>0$ and a radius $\varepsilon_{0}>0$ such that $C_{1} \varepsilon \leq \mathcal{H}^{1}\left(\Sigma \cap B_{\varepsilon}(x)\right) \leq C_{2} \varepsilon$ for every $x \in \Sigma$ and $\varepsilon<\varepsilon_{0}$.

In the work [98] Paolini and Stepanov proved

- the absence of closed loops for maximum distance minimizers and, under general conditions on $\phi$, the absence of closed loops for average distance minimizers;
- the Ahlfors regularity of maximum distance minimizers and, under the additional summability condition on $\phi$, the Ahlfors regularity of average distance minimizers. Gordeev and Teplitskaya [58] refine Ahlfors constants of maximum distance minimizers to the best possible, i.e. show that $\mathcal{H}^{1}\left(\Sigma \cap B_{\varepsilon}(x)\right)=\operatorname{ord}_{x} \Sigma \cdot \varepsilon+o(\varepsilon)$, where $\operatorname{ord}_{x} \Sigma \in\{1,2,3\}$.
- Recall that maximal distance minimization problem and the dual problem have the same sets of solutions (the planar case was proved before by Miranda, Paolini, Stepanov in [92]). It particularly implies that maximal distance minimizers must have maximum available length $l$. Paolini and Stepanov also proved that average distance minimizers (with additional assumptions on $\phi$ ) have maximum available length.

In the work [6] the following basic results were shown.
(i) Let $\Sigma$ be an $r$-minimizer for some $M$. Then $\Sigma$ is an $r$-minimizer for $\overline{B_{r}(\Sigma)}$.
(ii) Let $\Sigma$ be an $r$-minimizer for $\overline{B_{r}(\Sigma)}$. Then $\Sigma$ is an $r^{\prime}$-minimizer for $\overline{B_{r^{\prime}}(\Sigma)}$, where $0<r^{\prime}<r$.

### 4.1.6 Local maximal distance minimizers

Definition 4.1.2. Let $M \subset \mathbb{R}^{d}$ be a compact set and let $r>0$. A closed connected set $\Sigma \subset \mathbb{R}^{d}$ with $\mathcal{H}^{1}(\Sigma)<\infty$ is called a local minimizer if $\mathcal{F}_{M}(\Sigma) \leq r$ and there exists $\varepsilon>0$ such that for any connected set $\Sigma^{\prime}$ satisfying $\mathcal{F}_{M}\left(\Sigma^{\prime}\right) \leq r$ and $\operatorname{diam}\left(\Sigma \triangle \Sigma^{\prime}\right) \leq \varepsilon$ the inequality $\mathcal{H}^{1}(\Sigma) \leq \mathcal{H}^{1}\left(\Sigma^{\prime}\right)$ holds, where $\triangle$ is the symmetric difference.

Any maximal distance minimizer is also a local minimizer. Usually the properties of maximal distance minimizers are also true for local maximal distance minimizers (see [58]).

### 4.2 Regularity

### 4.2.1 Tangent rays. Blow up limits in $\mathbb{R}^{d}$

Definition 4.2.1. We say that a ray ( $a x]$ is a tangent ray of a set $\Gamma \subset \mathbb{R}^{d}$ at the point $x \in \Gamma$ if there exists a sequence of points $x_{k} \in \Gamma \backslash\{x\}$ such that $x_{k} \rightarrow x$ and $\angle x_{k} x a \rightarrow 0$.

Then we have the following regularity theorem.
Theorem 4.2.1 (Gordeev-Teplitskaya [58]). Let $\Sigma$ be a minimizer for a compact set $M \subset \mathbb{R}^{d}$ and $r>0$. Then there are at most three tangent rays at any point of $\Sigma$, and the pairwise angles between the tangent rays are at least $2 \pi / 3$. Furthermore, tangent rays coincide with one-sided tangents, particularly the angles between one-sided tangents cannot be equal to 0, i.e. there is one to one correspondence between tangent rays at an arbitrary point $x \in \Sigma$ and connected components of $\Sigma \backslash\{x\}$. Moreover, if $d=2$, then $\Sigma$ is a finite union of simple curves with one-sided tangents continuous from the corresponding side.

In works concerning average distance minimizers the notion of blow up limits is used. Santambrogio and Tilli in 113 proved that for any average distance minimizer blow up sequence $\Sigma_{\varepsilon}:=\varepsilon^{-1}\left(\Sigma_{a} \cap\right.$ $\left.B_{\varepsilon}(x)-x\right)$ with $x \in \Sigma_{a}$, converges in $B_{1}(0)$ (for the Hausdorff distance) to some limit $\Sigma_{0}(x)$ when $\varepsilon \rightarrow 0$, and the limit is one of the following below (see Fig. 4.1 which is analogues to a picture from [84]), up to a rotation.

a general $x$ has tangent line $\psi(x)=0$

$x$ is a leaf
$\psi(x)>0$

$x$ is a corner point
$\psi(x)>0$

$x$ is a branching point

$$
\psi(x)=0
$$

Figure 4.1: All possible variants of tangent rays at any point of a maximal distance minimizer (or blow up limits of an average distance minimizer)

It is clear that for maximal distance minimizers blow up limits also exists and are more or less the same: $\Sigma_{0}$ can be a radius, a diameter, a corner points with the angle between the segments greater or equal to $2 \pi / 3$ or a center of a regular tripod. Herewith at the second and third case (id est when $\psi(x)>0)$ the point $x$ has to be energetic; see the following definition.
Definition 4.2.2. A point $x \in \Sigma$ is called energetic, if for all $\rho>0$ one has

$$
F_{M}\left(\Sigma \backslash B_{\rho}(x)\right)>F_{M}(\Sigma)
$$

Herewith if a point $x$ of a maximal distance minimizer $\Sigma$ is energetic then there exists such a point $y \in M$ (may be not unique) such that dist $(x, y)=r$ and $B_{r}(y) \cap \Sigma=\emptyset$; such $y$ is called corresponding to $x$.

If a point $x \in \Sigma$ is not energetic then in a sufficiently small neighbourhood it is a center of a regular tripod or a segment (and coincides there with its one-sided tangents).

A key object in all the study of the average distance problem is the pull-back measure of $\mu$ with respect to the projection onto $\Sigma_{a}$, where $\Sigma_{a}$ is a solution of the average distance minimizer problem. More precisely, if $\mu$ does not charge the Ridge set (which is defined as the set of all $x \in \mathbb{R}^{d}$ for which the minimum distance to $\Sigma_{a}$ is attained at more than one point) of $\Sigma_{a}$ (this is the case for instance when $\mu$ is absolutely continuous with respect to the Lebesgue measure), then it is possible to choose a measurable selection of the projection multimap onto $\Sigma$, i.e. a map $\pi_{\Sigma}: M \rightarrow \Sigma$ such that $d(x, \Sigma)=d\left(x, \pi_{\Sigma_{a}}\right)$ (this map is uniquely defined everywhere except the Ridge set). Then one can define the measure $\psi$ as being $\psi(A):=\mu\left(\pi_{\Sigma_{a}}^{-1}(A)\right)$, for any Borel set $A \subset M$. In other words $\psi=\pi_{\Sigma_{a}} \sharp \mu$.

For the maximal distance minimizers in $\mathbb{R}^{d}$ we can define measure $\psi$ the similar way, but replace $M$ by $\partial B_{r}(\Sigma)$ and with $\left(n-1\right.$ )-dimensional Hausdorff measure as $\mu$ (or accordingly $\overline{B_{r}(\Sigma)}$ and $n$-dimensional Hausdorff measure). Thus Fig. 4.1 is true both for maximal and average distance minimizers.

### 4.2.2 Properties of branching points in $\mathbb{R}^{2}$

Recall that by Theorem4.2.1 that for every planar compact set $M$ and a positive number $r$ a maximal distance minimizer can have only a finite number of points with 3 tangent rays.

In the plane it is also known (see [12]) that every average distance minimizer is topologically a tree composed of a finite union of simple curves joining with a number of 3 .

Every branching point of a planar maximal distance minimizer should be the center of a regular tripod. If $x \in \Sigma \subset \mathbb{R}^{2}$ has 3 tangent rays then there exists such a neighbourhood of $x$ in which the minimizer coincides with its tangent rays. Id est, there exists such $\varepsilon>0$ that $\Sigma \cap \overline{B_{\varepsilon}(x)}=$ $[a x] \cup[b x] \cup[c x]$ where $\{a, b, c\}=\Sigma \cap \partial B_{\varepsilon}(x)$ and $\angle a x b=\angle b x c=\angle c x a=2 \pi / 3$. For planar average distance minimizers it is proved that any branching point admits such a neighbourhood in which three pieces of $\Sigma$ are $C^{1,1}$.

### 4.2.3 Continuity of one-sided tangents in $\mathbb{R}^{2}$

Definition 4.2.3. We will say that the ray (ax] is a one-sided tangent of a set $\Gamma \subset \mathbb{R}^{d}$ at a point $x \in \Gamma$ if there exists a connected component $\Gamma_{1}$ of $\Gamma \backslash\{x\}$ such that $x \in \overline{\Gamma_{1}}$ and that any sequence of points $x_{k} \in \Gamma_{1}$ with the property $x_{k} \rightarrow x$ satisfies $\angle x_{k} x a \rightarrow 0$. In this case we will also say that (ax] is tangent to the connected component $\Gamma_{1}$.

In the plane the continuity of one-sided tangents from the corresponding side holds (see [58):
Lemma 4.2.1. Let $\Sigma \subset \mathbb{R}^{2}$ be a (local) maximal distance minimizer and let $x \in \Sigma$. Let $\Sigma_{1}$ be a connected component of $\Sigma \backslash\{x\}$ with one-sided tangent (ax] (it has to exist) and let $\bar{x} \in \Sigma_{1}$.

1. For any one-sided tangent $(\bar{a} \bar{x}]$ of $\Sigma$ at $\bar{x}$ the equality $\angle((\bar{a} \bar{x}),(a x))=o_{|\bar{x} x|}(1)$ holds.
2. Let $(\bar{a} \bar{x}]$ be a one-sided tangent at $\bar{x}$ of any connected component of $\Sigma \backslash\{\bar{x}\}$ not containing $x$. Then $\angle((\bar{a} \bar{x}],(a x])=o_{|\bar{x} x|}(1)$.

For planar average distance minimizers it is proved (see 84]) that away from branching points an average distance minimizer $\Sigma_{a}$ is locally at least as regular as the graph of a convex function, namely that the Right and Left tangent maps admit some Right and Left limits at every point and are semicontinuous. More precisely, for a given parametrization $\gamma$ of an injective Lipschitz arc $\Gamma \subset \Sigma_{a}$, by existence of blow up limits one can define the Left and Right tangent half-lines at every point $x \in \Gamma$ by

$$
T_{R}(x):=x+\mathbb{R}^{+} . \lim _{h \rightarrow 0} \frac{\gamma\left(t_{0}+h\right)-\gamma\left(t_{0}\right)}{h}
$$

and

$$
T_{L}(x):=x+\mathbb{R}^{+} . \lim _{h \rightarrow 0} \frac{\gamma\left(t_{0}-h\right)-\gamma\left(t_{0}\right)}{h} .
$$

Then the following planar theorem for average distance minimizers holds.
Theorem 4.2.2 (Lemenant, 2011 [83]). Let $\Gamma \subset \Sigma_{a}$ be an open injective Lipschitz arc. Then the Right and Left tangent maps $x \rightarrow T_{R}(x)$ and $x \rightarrow T_{L}(x)$ are semicontinuous, id est for every $y_{0} \in \Gamma$ there holds $\lim _{y \rightarrow y_{0} ; y<\gamma y_{0}} T_{L}(y)=T_{L}\left(y_{0}\right)$ and $\lim _{y \rightarrow y_{0} ; y>\gamma y_{0}} T_{R}(y)=T_{R}\left(y_{0}\right)$. In addition the limit from the other side exists and we have $\lim _{y \rightarrow y_{0} ; y>\gamma y_{0}} T_{L}(y)=T_{R}\left(y_{0}\right)$ and $\lim _{y \rightarrow y_{0} ; y<\gamma y_{0}} T_{R}(y)=T_{L}\left(y_{0}\right)$.

An immediate consequence of the theorem is the following corollary:
Corollary 4.2.1. Assume that $\Gamma \subset \Sigma$ is a relatively open subset of $\Sigma$ that contains no corner points nor branching points. Then $\Gamma$ is locally a $C^{1}$-regular curve.

### 4.2.4 Planar example of infinite number of corner points

Recall that each maximal distance minimizer in the plane is a finite union of simple curves. These curves should have continuous one-sided tangents but do not have to be $C^{1}$ : there exists a minimizer with an infinite number of points without tangent lines. The following example is provided in [6].

Fix positive reals $r, R$ and let $N$ be a large enough integer. Consider a sequence of points $\left\{a_{i}\right\}_{i=1}^{\infty}$ chosen from the circumference $\partial B_{R}(o)$ such that $N \cdot\left|a_{2} a_{1}\right|=r$,

$$
\left|a_{i+1} a_{i+2}\right|=\frac{1}{2}\left|a_{i} a_{i+1}\right|
$$

and $\angle a_{i} a_{i+1} a_{i+2}>\frac{\pi}{2}$ for every $i \in \mathbb{N}$ (see Fig. 4.2). Let $a_{\infty}$ be the limit point of $\left\{a_{i}\right\}$. Finally, let $a_{\infty+1}$ be the point in the tangent line to $B_{r}(o)$ at $a_{\infty}$, such that

$$
\left|a_{\infty} a_{\infty+1}\right|=r / N
$$

We claim that polyline

$$
\Sigma=\bigcup_{i=1}^{\infty}\left[a_{i} a_{i+1}\right]
$$

is a unique maximal distance minimizer for the following $M$.
Let $v_{1} \in\left(a_{1} a_{2}\right]$ be such point that $\left|v_{1} a_{1}\right|=r$. For $i \in \mathbb{N} \cup\{\infty\} \backslash\{1\}$ define $v_{i}$ as the point satisfying $\left|v_{i} a_{i}\right|=r$ and $\angle a_{i-1} a_{i} v_{i}=\angle a_{i+1} a_{i} v_{i}>\pi / 2$. Define $v_{\infty}$ as the limit point of $\left\{v_{i}\right\}$. Finally, let $v_{\infty+1}$ be such point that $v_{\infty+1} a_{\infty} \perp v_{\infty} a_{\infty}$ and $\left|v_{\infty+1} a_{\infty}\right|=r$. Clearly $M:=\left\{v_{i}\right\}_{i=1}^{\infty+1}$ is a compact set.


Figure 4.2: The example of a minimizer with infinite number of corner points

Theorem 4.2.3 (Basok-Cherkashin-Teplitskaya, 2022 [6]). Let $\Sigma$ and $M$ be defined above. Then $\Sigma$ is a unique maximal distance minimizer for $M$.

### 4.2.5 Every $C^{1,1}$-smooth simple curve is a minimizer

For planar average distance minimizers Tilli proved in [119] that any simple $C^{1,1}$-curve is a minimizer for some given data. Paper [6] generalizes Tilli's result on $d$-dimensional space. The same statement with a similar but much simpler explanation is true for maximal distance minimizers.

Theorem 4.2.4 (Basok-Cherkashin-Teplitskaya, 2022 [6]). Let $\gamma \subset \mathbb{R}^{d}$ be a simple $C^{1,1}$-curve. Then $\gamma$ is a maximal distance minimizer for a small enough $r$ and $M=\overline{B_{r}(\gamma)}$.

### 4.3 Explicit examples for maximal distance minimizers

Recall that Theorems 4.2.3 and 4.2.4 provide explicit examples; however they are obtained by "reverse engineering": the input $M$ is constructed in a way to give the minimizer property to a desired $\Sigma$. This section is devoted to known explicit results.

### 4.3.1 Simple examples. Finite number of points and $r$-neighbourhood. Inverse minimizers

Here we consider Problem 4.1.2 in a case when $M$ is a finite set. Then it is closely related with the Steiner problem (Problem 2.0.1).

Any maximal distance minimizer for any finite set $M \subset \mathbb{R}^{d}$ is a finite union of at most $2 \sharp M-$ 1 segments. In this case maximal distance minimization problem comes down to connecting $r$ neighborhoods of all the points from $M$. If $\overline{B_{r}(a)}$ are disjoint for every $a \in M$ then a maximal distance minimizer is a Steiner tree connecting some points from $\partial B_{r}(a), a \in M$.

The following observations and statements of this paragraph are from the paper [6].
Remark 4.3.1. (i) Let $\Sigma$ be a maximal distance minimizer for some $M$ and $r>0$. Then $\Sigma$ is a maximal distance minimizer for $\overline{B_{r}(\Sigma)}$ and $r$.
(ii) Let $\Sigma$ be a minimizer for $\overline{B_{r}(\Sigma)}$ and $r>0$. Then $\Sigma$ is a minimizer for $\overline{B_{r^{\prime}}(\Sigma)}$ and $r^{\prime}$, where $0<r^{\prime}<r$.


Figure 4.3: A maximal distance minimizer for a certain 3-point set $M=\{a, b, c\}$
A topology $T$ of a labelled Steiner tree (or a labelled locally minimal tree) $\mathcal{S} t$ is the corresponding abstract graph with labelled terminals and unlabelled Steiner points.

Theorem 4.3.1 (Basok-Cherkashin-Teplitskaya, 2022 [6]). Let St be a Steiner tree for a labelled set of terminals $A=\left(a_{1}, \ldots, a_{n}\right), a_{i} \in \mathbb{R}^{d}$ such that every Steiner tree for an n-tuple in the closed $2 r$-neighbourhood of $A$ (with respect to $\rho$ ) has the same topology as $\mathcal{S}$ for some positive $r$. Then $\mathcal{S} t$ is an r-minimizer for an n-tuple $M$.

In the plane a Steiner tree for a random input is unique with unit probability, see [5]. Also in the plane we have a general inverse statement to Theorem 4.3.1.

Proposition 4.3.1 (Basok-Cherkashin-Teplitskaya, 2022 [6]). Suppose that St is a full Steiner tree for terminals $a_{1}, \ldots, a_{n} \in \mathbb{R}^{2}$, which is not unique. Then $\mathcal{S} t$ can not be a minimizer for $M$ being an $n$-tuple of points.

To illustrate Proposition 4.3.1 consider a square $a_{1} a_{2} a_{3} a_{4}$. There are two Steiner trees for $a_{1}, a_{2}, a_{3}, a_{4}$ (see the left-hand side of Fig. 4.4), let us pick the solid one. The right-hand side of Fig. 4.4 shows that an $r$-minimizer for every positive $r$ has the topology of the dotted Steiner tree.

In all known examples a $\mathcal{S} t$ with $n$ terminals is an $r$-minimizer for a set $M$ of $n$ points and a small enough positive $r$ if and only if $\mathcal{S} t$ in the unique Steiner tree for its terminals. So the planar case of several non-full solutions is open, and also it is interesting to derive any analogue of Proposition 4.3.1 for $d>2$.


Figure 4.4: An example to Proposition 4.3.1

### 4.3.2 Circle. Curves with big radius of curvature

Theorem 4.3.2 (Cherkashin-Teplitskaya, 2018 [18]). Let $r$ be a positive real, $M$ be a convex closed curve with the radius of curvature at least $5 r$ at every point, $\Sigma$ be an arbitrary minimizer for $M$. Then $\Sigma$ is a union of an arc of $M_{r}$ and two segments that are tangent to $M_{r}$ at the ends of the arc (so-called horseshoe, see Fig. 4.5). In the case when $M$ is a circumference with the radius $R$, the condition $R>4.98 r$ is enough.

Also Theorem 4.3.2 admits a corollary on local minimizers in the sense of Definition 4.1.2.
Corollary 4.3.1 (Cherkashin-Teplitskaya, 2018 [18]). Let $\hat{\Sigma}$ be a local minimizer for some closed convex curve $M$ with minimal radius of curvature $R>5 r$. Then if $\hat{\Sigma}$ is not a horseshoe, one has $\mathcal{H}^{1}(\hat{\Sigma})-\mathcal{H}^{1}(\Sigma) \geq(R-5 r) / 2$, where $\Sigma$ is an arbitrary (global) minimizer.

Miranda, Paolini and Stepanov [92] conjectured that all the minimizers for a circumference of radius $R>r$ are horseshoes. Theorem 4.3.2 solves this conjecture with the assumption $R>4.98 r$; for $4.98 r \geq R>r$ the conjecture remains open.


Figure 4.5: A minimizer for a convex closed planar curve $M$ with the radius of curvature at least $5 r$ at every point, so-called horseshoe (left). A minimizer for $M=\partial B_{R}(x)$, where $R>4.98 r$ (right)


Figure 4.6: $M$ is $r$-neighbourhood for a sufficiently smooth curve $\Sigma$ and small enough $r>0$

### 4.3.3 Rectangle

Theorem 4.3.3 (Cherkashin-Gordeev-Strukov-Teplitskaya, 2021 [14]). Let $M=a_{1} a_{2} a_{3} a_{4}$ be a rectangle. Then there is a positive number $r_{0}(M)$ such that for any positive $r<r_{0}(M)$ every minimizer of the maximum distance functional has a topology of 21 segments, shown on the leftmost side of Fig. 4.7. The middle part of the figure shows an enlarged fragment of the minimizer in the vicinity of point $a_{1}$; the marked angles are equal to $\frac{2 \pi}{3}$. The rightmost side of the figure shows an even more enlarged fragment of the minimizer in the vicinity of $a_{1}$.

Any minimizer of the maximum distance functional has length $\operatorname{Per}(M)-c r+o(r)$, where $\operatorname{Per}(M)$ is the perimeter of the rectangle $M$, and $c$ is a constant approximately equal to 8.473981.

In fact, every maximal distance minimizer is very close (in the sense of Hausdorff distance) to the one depicted in the picture.

### 4.4 Tools

### 4.4.1 Energetic points

For the planar problem the notion of energetic points (which is also correct in $\mathbb{R}^{d}$ ) is very useful.
Recall that a point $x \in \Sigma$ is called energetic, if for all $\rho>0$ one has $F_{M}\left(\Sigma \backslash B_{\rho}(x)\right)>F_{M}(\Sigma)$. The set of all energetic points of $\Sigma$ is denoted by $G_{\Sigma}$. Each minimizer $\Sigma$ can be split into three disjoint


Figure 4.7: The minimizer for a rectangle $M$ with $r<r_{0}(M)$.
subsets:

$$
\Sigma=E_{\Sigma} \sqcup \mathcal{X}_{\Sigma} \sqcup \mathcal{S}_{\Sigma},
$$

where $X_{\Sigma} \subset G_{\Sigma}$ is the set of isolated energetic points (i.e. every $x \in X_{\Sigma}$ is energetic and there is a $\rho>0$ possibly depending on $x$ such that $\left.B_{\rho}(x) \cap G_{\Sigma}=\{x\}\right), E_{\Sigma}:=G_{\Sigma} \backslash X_{\Sigma}$ is the set of non isolated energetic points and $S_{\Sigma}:=\Sigma \backslash G_{\Sigma}$ is the set of non energetic points also called the Steiner part.

Note that it is possible for a (local) minimizer in $\mathbb{R}^{d}, d>2$ to have no non-energetic points at all. Moreover, in some sense, any (local) minimizer does not have non-energetic points in a larger dimension:

Example 4.4.1. Let $\Sigma$ be a (local) minimizer for a compact set $M \subset \mathbb{R}^{d}$ and $r>0$. Then $\bar{\Sigma}:=$ $\Sigma \times\{0\} \subset \mathbb{R}^{d+1}$ is a (local) minimizer for $\bar{M}=(M \times\{0\}) \cup(\Sigma \times\{r\}) \subset \mathbb{R}^{d+1}$ and $\mathcal{E}_{\bar{\Sigma}}=\bar{\Sigma}$.

Recall that for every point $x \in G_{\Sigma}$ there exists a point $y \in M$ (may be not unique) such that dist $(x, y)=r$ and $B_{r}(y) \cap \Sigma=\emptyset$. Thus all points of $\Sigma \backslash \overline{B_{r}(M)}$ can not be energetic and thus $\overline{\Sigma \backslash \overline{B_{r}(M)}}$ is so-called Steiner forest id est each connected component of it is a Steiner tree with terminal points on the $\partial B_{r}(M)$.

In the plane it makes sense to define energetic rays.
Definition 4.4.1. We say that a ray ( $a x$ ] is the energetic ray of the set $\Sigma$ with a vertex at the point $x \in \Sigma$ if there exists non stabilized sequence of energetic points $x_{k} \in G_{\Sigma}$ such that $x_{k} \rightarrow x$ and $\angle x_{k} x a \rightarrow 0$.

Remark 4.4.1. Let $\left\{x_{k}\right\} \subset G_{\Sigma}$ and let $x \in E_{\Sigma}$ be the limit point of $\left\{x_{k}\right\}: x_{k} \rightarrow x$. By basic property of energetic points for every point $x_{k} \in G_{\Sigma}$ there exists a point $y_{k} \in M$ (may be not unique) such that $\operatorname{dist}\left(x_{k}, y_{k}\right)=r$ and $B_{r}\left(y_{k}\right) \cap \Sigma=\emptyset$. In this case we will say that $y_{k}$ corresponds to $x_{k}$.

Let $y$ be an arbitrary limit point of the set $\left\{y_{k}\right\}$. Then the set $\Sigma$ does not intersect r-neighbourhood of $y: B_{r}(y) \cap \Sigma=\emptyset$ and the point $y$ belongs to $M$ and corresponds to $x$.

Let $[s x] \subset \Sigma$ be a simple curve. Let us define $\operatorname{turn}([s x])$ as the upper limit (supremum) over all sequences of points of the curve:

$$
\operatorname{turn}([\breve{s x}])=\sup _{n \in \mathbb{N}, s \preceq t^{1} \prec \cdots \prec t^{n} \prec x} \sum_{i=2}^{n} \widehat{t^{i}, t^{i-1}}
$$

where $t^{i}$ denotes the ray of the one-sided tangent to the curve $\left[s t_{i}\right] \subset\left[s x\left[\right.\right.$ at point $t_{i}$, and $t_{1}, \ldots, t_{n}$ is the partition of the curve [sx[ in the order corresponding to the parameterization, for which $s$ is the
beginning of the curve and $x$ is the end. In this case, the angle $\left(\widehat{t^{i}, t^{i+1}}\right) \in[-\pi, \pi[$ between two rays is counted from ray $t^{i}$ to ray $t^{i+1}$; positive direction is counterclockwise.

Let $\breve{s x}$ lay in the sufficiently small neighbourhood of $x$. Then if $B_{r}(y(x)) \cap[\breve{s x}]=\emptyset$, it is true that

$$
|\operatorname{turn}([s x])|<2 \pi .
$$

This property is the first one which is true for the plane and false in $\mathbb{R}^{d}$ with $d>2$, so this is the main difference between planar and non-planar cases. In the plane the turn is a very useful tool, see for example the proof of Theorem 4.3.2 [18].

The second main differ between plane and other Euclidean spaces is also concerning angles: in the plane if you know the angles $\widehat{t^{i}, t^{i-1}}$ for $i=2, \ldots k$ then you know the angle $\widehat{t^{1}, t^{k}}$ which is not true for $\mathbb{R}^{d}$ with $d>2$.

### 4.4.2 Convexity argument

Suppose that we fix some $M_{0} \subset M$ and consider a (possibly infinite) tree $T$ whose vertices are encoded by points of $M_{0}$. Let us pick an arbitrary point from $\overline{B_{r}(m)}$ for every $m \in M_{0}$ and connect such points by segments with respect to $T$. Consider the length $L$ of such a representation of $T$; note that we allow the representation to contain cycles or edges of zero length.

Then $L$ is a convex function from $\left(\mathbb{R}^{d}\right)^{M_{0}}$ to $\mathbb{R}$. Also if $v, u \in \overline{B_{r}(m)}$, then $\alpha v+(1-\alpha) u$ also lies in $\overline{B_{r}(m)}$. It implies that the sets of local and global minimums of $L$ coincide and form a convex set. It usually means that $L$ is a unique local minimum.

This approach allows us to show that if one fixes a topology of a solution, then the corresponding Steiner-type problem has a unique solution. The proofs of Theorems 4.2.3 and 4.3.1 heavily use it.

### 4.4.3 Lower bounds on the length of a minimizer

The proof of the following folklore inequality can be found, for instance in [93].
Lemma 4.4.1. Let $\gamma$ be a compact connected subset of $\mathbb{R}^{d}$ with $\mathcal{H}^{1}(\gamma)<\infty$. Then

$$
\mathcal{H}^{d}\left(\left\{x \in \mathbb{R}^{d}: \operatorname{dist}(x, \gamma) \leq t\right\}\right) \leq \mathcal{H}^{1}(\gamma) \omega_{d-1} t^{d-1}+\omega_{d} t^{d}
$$

where $\omega_{k}$ denotes the volume of the unit ball in $\mathbb{R}^{k}$.
The following corollary is very close to a theorem of Tilli on average distance minimizers [119].
Corollary 4.4.1. Let $V$ and $r$ be positive numbers. Then for every set $M$ with $\mathcal{H}^{d}(M)=V a$ maximal distance $r$-minimizer has the length at least

$$
\max \left(0, \frac{V-\omega_{d} r^{d}}{\omega_{d-1} r^{d-1}}\right)
$$

Theorem 4.2 .4 follows from the fact that for a $C^{1,1}$-curve and small enough $r$ the inequality in Corollary 4.4.1 is sharp. Let us provide a lower bound from [6] on the length of a minimizer in the planar case.

Proposition 4.4.1. Let $M$ be a planar convex set and $\Sigma$ is an r-minimizer for $M$. Then

$$
\left.\mathcal{H}^{1} \Sigma\right) \geq \frac{\left.\mathcal{H}^{1} \partial M\right)-2 \pi r}{2}
$$

### 4.5 More properties of minimizers

### 4.5.1 $\quad$-convergence

$\Gamma$-convergence is an important tool in studying minimizers based on approximation of energy. For Euclidean space the following definition of $\Gamma$-convergence can be used. Let $X$ be a first-countable space and $F_{n}: X \rightarrow \overline{\mathbb{R}}$ a sequence of functionals on $X$. Then $F_{n}$ are said to $\Gamma$-converge to a $\Gamma$-limit $F: X \rightarrow \overline{\mathbb{R}}$ if the following two conditions hold.

- Lower bound inequality. For every sequence $x_{n} \in X$ such that $x_{n} \rightarrow x$ as $n \rightarrow+\infty$,

$$
F(x) \leq \liminf _{n \rightarrow \infty} F_{n}\left(x_{n}\right)
$$

- Upper bound inequality. For every $x \in X$, there is a sequence $x_{n}$ converging to $x$ such that

$$
F(x) \geq \limsup _{n \rightarrow \infty} F_{n}\left(x_{n}\right)
$$

In the case of maximal distance minimizers for a given compact set $M$ and a number $l>0$ we can consider the space $X$ of connected compact sets with one-dimensional Hausdorff measure at most $l$, equipped with the Hausdorff distance (the distance $d_{H}$ between $A, C \in X$ is the smallest $\rho$ such that $A \subset \overline{B_{\rho}(C)}$ and $C \subset \overline{\left.B_{\rho}(A)\right)}$.
Proposition 4.5.1. If a sequence of compacts $M_{i}$ converges to $M$ then $F_{M_{i}} \Gamma$-converges to $F_{M}$.
Proof. By the definition of $F_{M}$ and triangle inequalities we have

$$
\begin{equation*}
\left|F_{M_{i}}\left(S_{i}\right)-F_{M}(S)\right| \leq\left|F_{M_{i}}\left(S_{i}\right)-F_{M}\left(S_{i}\right)\right|+\left|F_{M}\left(S_{i}\right)-F_{M}(S)\right| \leq d_{H}\left(M_{i}, M\right)+d_{H}\left(S_{i}, S\right) \tag{4.1}
\end{equation*}
$$

for every connected $S_{i}$ and $S$. So by (4.1) every sequence of $S_{i}$ with limit $S$ we have the first condition of $\Gamma$-convergence. For the second condition consider $S_{i}$ being a Steiner tree for a finite $1 / i$-network $N_{i} \subset S$. By the definition $\mathcal{H}^{1}\left(S_{i}\right) \leq \mathcal{H}^{1}(S) \leq l$. Again, by 4.1) $F_{M_{i}}\left(S_{i}\right)$ converges to $F_{M}(S)$.

### 4.5.2 Approximation by Steiner trees

A crucial property of $\Gamma$-convergence is that in the notation of Proposition 4.5.1 every cluster point of the sequence of minimizers of $F_{M_{i}}$ is a minimizer of $F_{M}$. Now let $M_{n}$ be a finite $1 / n$-network for $M$, so that every minimizer for $M_{n}$ is a finite Steiner tree.

Unfortunately, in the case of several minimizers for $M$ we cannot be sure that every minimizer is approximated. On the other hand it can be approximated a posteriori. Let $\Sigma$ be a minimizer for $M$ and let $\mathcal{E}_{k} \subset \Sigma$ be a finite $1 / k$-network and $\Sigma_{k}$ be an arbitrary solution of the Steiner problem for $\mathcal{E}_{k}$. By the definition we have

$$
\mathcal{H}^{1}\left(\Sigma_{k}\right) \leq \mathcal{H}^{1}(\Sigma)
$$

On the other hand, for any subsequential limit (with respect to the Hausdorff distance) $\Sigma^{\prime}$ of the sequence $\Sigma_{k}$ we have $\Sigma \subset \Sigma^{\prime}$ and so

$$
\mathcal{H}^{1}(\Sigma) \leq \mathcal{H}^{1}\left(\Sigma^{\prime}\right) \leq \liminf _{k \rightarrow \infty} \mathcal{H}^{1}\left(\Sigma_{k}\right)
$$

by Gołąb's theorem. It follows that $\Sigma_{k}$ converges to $\Sigma$ and $\mathcal{H}^{1}\left(\Sigma_{k}\right)$ converges to $\mathcal{H}^{1}(\Sigma)$.
Summing up, every maximal distance minimizer is a limit of finite Steiner trees. Similar results are also proved in [2]. More detailed and structural relations of finite Steiner trees and maximal distance minimizer are considered in Section 4.3.1.

### 4.5.3 NP-hardness

It is well-known that Euclidean Steiner problem is NP-hard [53] even is we restrict the terminals to two lines in the plane [112]. The first source of hardness is that if we fix a topology then one can write the length in the explicit form. However the expression may have $\Omega(n)$ square roots.

To avoid it Garey, Graham and Johnson [53] introduce a discrete version of the Steiner problem. Of course a minimizer of a new problem does not inherit any geometric properties, in particular we have no $2 \pi / 3$-condition at a branching point. Such a discretization appears to be NP-complete (and so the initial one is NP-hard), namely, Garey, Graham and Johnson used a reduction of the X3C problem to this version of the Steiner problem. The X3C problem is to decide whether a family of 3 -sets $\mathcal{F} \subset 2^{[3 n]}$ has a subfamily of $n$ sets which cover [3n]. It is well-known that X3C is NP-complete.

First we need the following reduction to the classical Steiner problem.
Theorem 4.5.1 (Garey-Graham-Johnson [53]). For a given $\mathcal{F} \subset 2^{[3 n]}$ one can construct in a polynomial time in $n$ an input $X(\mathcal{F}) \subset \mathbb{R}^{2}$ whose size is also polynomial in $n$ such that
(i) if $\mathcal{F}$ has an n-set covering then a solution of the Steiner problem for $X(\mathcal{F})$ has the length at most L;
(ii) if $\mathcal{F}$ does not have an n-set covering then a solution of the Steiner problem for $X(\mathcal{F})$ has the length at least $L+12|X(\mathcal{F})|$.

Moreover $L=L(\mathcal{F})$ can be extracted from the construction of $X(\mathcal{F})$ in an explicit form.
Now let us repeat Garey-Graham-Johnson rounding in the case of maximal distance minimizers. The following problem is a discrete approximation of Problem 4.1.2 analogous to the discrete version of Steiner problem used in [53]. Following [53] we replace the length function with is ceiling because it is not known if the problem of determining whether $\sum \sqrt{n_{i}}<L$ is NP or not ( $n_{i}, L$ are integer).

Problem 4.5.1. Let $M$ be a finite set of points in the plane with integer coordinates and $r, \ell \in \mathbb{N}$. Decide whether exists a connected graph whose vertices have integer coordinates and edges are segments with the sum of the ceiling function of the length over edges at most $\ell$ such that every point of $M$ lies at a distance at most $r$ from some vertex of the graph.

Now we are ready to obtain the following corollary of Garey-Graham-Johnson results and the approximation.

Proposition 4.5.2. Problem 4.5.1 is NP-complete.
Proof. Let $\mathcal{F} \subset 2^{[3 n]}$ be an arbitrary family. Consider the set $X(\mathcal{F})$ from Theorem 4.5.1. Fix any $r \in \mathbb{N}$ and let $k>10 r|X(\mathcal{F})|$ be a large integer number. Define $k X(\mathcal{F})$ as a set homothetic to $X(\mathcal{F})$ with the scale factor $k$. Let $M$ be the set of closest points of $k X(\mathcal{F})$ in the integer grid $\mathbb{Z}^{2}$. Put also $\ell=k L(\mathcal{F})+k$.

Then if $F$ has an $n$-set covering, then a solution $\mathcal{S}$ t of the Steiner problem for $k X(\mathcal{F})$ has the length at most $k L(\mathcal{F})$. Now we replace in $\mathcal{S} t$ every vertex with the closest point from $\mathbb{Z}^{2} ;$ denote the resulting set by $\mathcal{S} t^{D}$. By the definition $\mathcal{S} t^{D}$ is a graph whose vertices have integer coordinates and it connects $M$; also it has at most $2|X(\mathcal{F})|-3$ segments. After the rounding the length of every segment of $\mathcal{S t}$ grows by at most $\sqrt{2}$. Hence the ceiling function of the length of an edge in $\mathcal{S} t^{D}$ is at most the length of the corresponding edge of $\mathcal{S} t$ plus 3 . Thus sum of the ceiling function of the length over edges $\mathcal{S} t^{D}$ is at most

$$
k L+6|X(\mathcal{F})|<\ell
$$

and the answer to Problem 4.5.1 is positive.
On the other hand let us show that in the case when $F$ has no $n$-set covering, the answer to Problem 4.5.1 is negative. Consider a solution $\Sigma^{D}$ of Problem 4.5.1 for $X(\mathcal{F})$. Assume the contrary, so that the sum of the ceiling function of length over the edges of $\Sigma^{D}$ is at most $\ell$. It implies $\mathcal{H}^{1}\left(\Sigma^{D}\right) \leq \ell$. Consider the homothety $\Sigma_{1 / k}^{D}$ of $\Sigma^{D}$ with the scale factor $1 / k$, one has

$$
\mathcal{H}^{1}\left(\Sigma_{1 / k}^{D}\right)=\frac{\left.\mathcal{H}^{1} \Sigma^{D}\right)}{k} \leq \frac{\ell}{k} .
$$

By definition for every $x \in M$ ) there is a point $\sigma \in \Sigma^{D}$ at a distance at most $r$ from $x$. Hence for every $x \in X(\mathcal{F})$ there is a point $\sigma \in \Sigma_{1 / k}^{D}$ at a distance at most $(r+1) / k<1$ from $x$. Thus the length of a Steiner tree for $X(\mathcal{F})$ is at most

$$
\frac{\mathcal{H}^{1}\left(\Sigma^{D}\right)}{k}+|X(\mathcal{F})| \cdot \frac{r+1}{k} \leq \frac{\ell+(r+1) \cdot|X(\mathcal{F})|}{k} \leq L+|X(\mathcal{F})| .
$$

We got a contradiction with Theorem 4.5.1 and thus finished the reduction of the X3C problem to Problem 4.5.1 with the input $M, r, \ell$.

Finally one can easily compute the sum of the ceiling function of length over edges of a competitor for Problem 4.5.1 in polynomial time.

### 4.5.4 Penalized form

Let $M$ be a given compact set. Let us consider a problem of minimization $F_{M}(S)+\lambda \mathcal{H}^{1}(S)$ for some $\lambda>0$, where $F_{M}(S)=\max _{y \in M} \operatorname{dist}(y, S)$ among all connected compact sets $S$. We will call this problem $\lambda$-penalized.

Clearly every set $T$ which minimizes $\lambda$-penalized problem for some $\lambda$ is a maximal distance minimizer for a given data $M$ and the restriction of energy $r:=F_{M}(T)$. Hence the solutions of this problem inherit all regularity properties of maximal distance minimizers.

As usual in variational calculus on a restricted class, it may happen for a small variation $\Phi_{\varepsilon}(\Sigma)$ of $\Sigma$, that the length constraint $\mathcal{H}^{1}\left(\Phi_{\varepsilon}(\Sigma)\right) \leq l$ is violated. Hence to compute Euler-Lagrange equation associated to the maximal distance minimization problem a possible way is to consider first the penalized functional $F_{M}(S)+\lambda \mathcal{H}^{1}(S)$ for some constant $\lambda$, for which any competitor $\Sigma$ is admissible without length constraint.

Hence it is also make sense to consider local penalisation problem: the problem of searching a connected compact set $S$ of a finite length, such that $\mathcal{H}^{1}(S)+\lambda F_{M}(S) \leq \mathcal{H}^{1}(T)+\lambda F_{M}(T)$ for every connected compact $T$ with $\operatorname{diam}(S \triangle T)<\varepsilon$ for sufficiently small $\varepsilon>0$. The solutions of these problems also inherit properties of local maximal distance minimizers.

## Proposition 4.5.3. Consider

$$
\left.\min _{\Sigma \text { compact and connected }} F_{M}(\Sigma)+\lambda\left(\mathcal{H}^{1} \Sigma\right)-l\right)^{+}
$$

for any constant $\lambda>1$. Then this problem is equivalent to the maximal distance minimization problem.

Proof. The same as for average distance minimizers (see Proposition 23 in [84]). We use the fact that for a connected set $S \backslash T_{\varepsilon}$ if $S$ is a maximal distance minimizer and $\mathcal{H}^{1}\left(T_{\varepsilon}\right)=\varepsilon$ there holds $r-F_{M}\left(S \backslash T_{\varepsilon}\right) \leq \varepsilon$.

### 4.5.5 Uniqueness

Let us start with the following simple observation. The set of minimizers for $M$ being a circle $B_{R}(o)$ is uncountable for $r<R$. Indeed, any minimizer has no loops and does not reduce to a point, so its rotations rarely coincide.

Note that for every compact $M \subset \mathbb{R}^{d}$ and $r$ equal to the radius of the smallest ball containing $M$, there is a unique point $o$ such that $M \subset \overline{B_{r}(o)}$, id est the solution of Problem 4.1.2 is unique. For a larger $r$ one has an uncountable number of solutions. This motivate us to consider only small enough $r$. Let us call a finite point configuration $M$ ambiguous if Problem 4.1.2 has several solutions for $M$ and $r<r_{0}(M)$. The following statement is a straightforward corollary from the Theorem 2.1.1.

Proposition 4.5.4. For $n \geq 4$ the set of planar n-point ambiguous configurations $M$ has Hausdorff dimension $2 n-1$ (as a subset of $\mathbb{R}^{2 n}$ ).

Proof. Fix $n \geq 4$. Theorem 2.1.1 states that the Hausdorff dimension of planar $n$-point configurations with multiply Steiner trees is $2 n-1$.

Recall that a topology $T$ of a labelled Steiner tree is the corresponding abstract graph; a topology is full if every its terminal has degree 1 . We call a topology generic if it has no terminals of degree 3. For a generic topology $R$ one can replace vertex $A$ of degree 2 with a Steiner point $b$ and add edge $b A$; the resulting topology $T(R)$ is full.

By Proposition 2.1.1 the set of all configurations for which there is a realization of any degenerate topology has Hausdorff dimension $2 n-2$.

Let us show that if a Steiner tree for a finite $M$ is unique then $M$ is not ambiguous. Consider any $n$-point planar configuration $M$ with unique Steiner tree $\mathcal{S} t$ whose topology is generic. Let the length of the second locally minimal tree be $\mathcal{H}^{1}(\mathcal{S} t)+a$ and choose $r<a /(2 n)$.

Then a maximal distance minimizer for a given $M$ and $r$ is obtained by a convexity argument for a topology $T$. Thus the Hausdorff dimension of planar $n$-point ambiguous configuration is at most $2 n-1$.

To show that the Hausdorff dimension $2 n-1$ of the set of planar $n$-point ambiguous configurations is at least $M$ we word-by-word repeat the argument of Lemma 2.4.5.

Note that we need $n \geq 4$ is the proof since there is only one full topology for each $n \leq 3$.

### 4.6 On minimizers for a planar convex closed smooth curve

### 4.6.1 The class of $M$ considered in the section

Fix a positive real $r$ and a closed convex curve $M$ with the minimal radius of curvature $R>r$ (this implies $C^{1,1}$-smoothness of $M$ ). Introduce the notation: $N:=\operatorname{conv}(M)$; let $M_{r}$ be the inner part of the boundary of $B_{r}(M)$, and finally put $N_{r}=\operatorname{conv}\left(M_{r}\right)$. In the literature, $M_{r}$ is often called parallel or equidistant curve. Note that $M_{r}$ also is a closed convex curve $M$ with the minimal radius of curvature $R-r$ and inherits $C^{1,1}$-smoothness. For each point $y \in M$, the circle $\partial B_{r}(y)$ touches $M_{r}$; let us call the tangent point $n_{y}$. The segment $\left[y n_{y}\right]$ is orthogonal to the curves $M$ and $M_{r}$.

Further $\Sigma$ denotes an arbitrary minimizer for $M$. We need the following simple observation. The definitions of $\partial T$ are given in Section 4.6.2.

Lemma 4.6.1. Under the assumptions of this section


Figure 4.8: Definitions of $N, M_{r}, N$, and $N_{r}$
(i) a connected component $T$ of the set $\overline{\Sigma \cap \operatorname{Int} N_{r}}$ is a Steiner tree, moreover, $\partial T \subset M_{r}$;
(ii) the set $\Sigma \cap \operatorname{Int} N_{r}$ is a subset of the Steiner part of $\mathcal{S}_{\Sigma}$.

Proof. Note that, by the basic property of minimizers (b), an energetic point $x \in \Sigma$ is located at a distance $r$ apart the point $y(x) \in M$, therefore $\operatorname{dist}(x, M) \leq r$, which entails item (ii).

Now we prove item (i). Consider a connected component $T$ of the set $\overline{\Sigma \cap \operatorname{Int} N_{r}}$; by already proven item (ii), the set $\partial T$ belongs to $M_{r}$. By Theorem 4.2.1 (i) the set $\partial T$ is finite. If $T$ is not a Steiner tree for $T \cap M_{r}$, then there is a shorter tree $T^{\prime}$ for the same points. Then replacing $\Sigma$ with $\Sigma \backslash T \cup T^{\prime}$ preserves connectivity, does not increase energy, and reduces the length of the minimizer. Contradiction.

### 4.6.2 Pseudo-networks

We will also often use the surprising fact that there is a unique simple path between two points in a path-connected acyclic set, in particular in a Steiner tree or a locally minimal tree.

For a given tree $T$ we denote the set of its vertices of degree 1 or 2 as $\partial T$.
Remark 4.6.1. Let $T$ be an arbitrary (full) Steiner tree, and $L$ an arbitrary line. Then the closure of an arbitrary connected component $T \backslash L$ is also a (full) Steiner tree for a finite subset of $\partial T \cup(T \cap L)$.

Definition 4.6.1. Define a "wind rose" as a set of six rays starting at the origin point with angle $\pi / 3$ between any two adjacent rays; each ray is given a weight (a real number), which satisfy the following property: the weight of a ray is the sum of weights of two rays adjacent to it.

It follows, in particular, that the sum of the weights of two opposite rays (the ones forming a line) is zero.

By full Steiner pseudo-network let us call a connected set $S$ which contains $C$, if for any wind rose $\mathcal{R}$ such that
(i) $S$ consists of finite number of segments which are parallel to $\mathcal{R}$
the following holds:
(ii) no point of $C$ is incident to exactly two segments,
(iii) for any $x \in S \backslash C$ and small enough $\varepsilon>0$, sum of weights of rays of $\mathcal{R}$ which are parallel to rays of the form $[x y), y \in \partial B_{\varepsilon}(x) \cap S$, is zero.

It is clear that a full Steiner tree or a full locally minimal tree is a full Steiner pseudo-network. Note that a full Steiner pseudo-networks may contain cycles, and a point may be incident to 1,2 , 3 , 4 or 6 segments; in what follows this number is called the degree of the point. For a given full Steiner pseudo-network $T$ we denote by $\partial T$ the set of its points of degree 1. The diagram in Fig. 4.9 illustrates the inclusions of certain classes of sets.


Figure 4.9: Inclusion relations between the used notions

For a given pseudo-network $T$ let us denote by $\partial T$ set of vertices of degree 1 .
Remark 4.6.2. Suppose that $T$ is a full Steiner pseudo-network, and $\mathcal{R}$ is an arbitrary wind rose satisfying (i). Let us assign to an each vertex $x \in \partial T$ a weight of a ray of $\mathcal{R}$, which is parallel to a directed segment of $T$ entering $x$ (such segment is unique by definition of $\partial T$ ). Then sum of assigned numbers over all $x \in \partial T$ is zero.

Lemma 4.6.2. Let $T$ be a full Steiner pseudo-network, $L$ an arbitrary line intersecting $T$ at a finite number of points. Then

$$
\sharp(\partial T \cap L) \leq 2 \sharp(\partial T \backslash L) .
$$

Proof of Lemma 4.6.2: Let $L^{+}, L^{-}$be the two open half-planes bounded by $L$. Note that it is sufficient to prove the inequality for a closure of an arbitrary component of $\overline{T \cap L^{+}}$and $\overline{T \cap L^{-}}$, denote such closure as $S$. By definition, $S$ is also a full Steiner pseudo-network.

Without loss of generality, we can assume that the center of $\mathcal{R}$ lies in the same half-plane as $S$. Let us choose the weights on the rays of the rose $\mathcal{R}$ in such a way that all those and only those rays that intersect $L$ have positive weights: obviously, such a choice exists since $L$ intersects 3 or 2 neighboring rays $\mathcal{R}$. In the first case we provide them with weights $1,2,1$, in the second case 1,1 . Then the remaining rays will have weights $-1,-2,-1$ or $0,-1,-1,0$ (we list all weights counterclockwise). For each vertex $x \in \partial S$ there is a unique ray from $\mathcal{R}$ codirected with the segment $S$ included in $x$; let us assign the weight of this ray to vertex $x$. Then the sum over endpoints $S$ belonging to $L$ is at least $\sharp(S \cap L)$. On the other hand, the sum over all other endpoints of $S$ is at least $-2 \sharp(\partial S \backslash L)$. Due to Remark 4.6.2 the sum of the weights over all endpoints of $S$ is equal to zero, that is

$$
\sharp(\partial S \cap L)-2 \sharp(\partial S \backslash L) \leq 0,
$$

which completes the proof of the lemma.

Remark 4.6.3. Let $T$ be a full Steiner pseudo-network fully lying on one side of line L, such that equality in Lemma 4.6.2 is achieved. Then all leaf vertices in $\partial T \backslash L$ have weight -2 , therefore all segments of $T$ incident to vertices from $\partial T \backslash L$ are pairwise collinear.

Lemma 4.6.3. Consider a regular tripod with ends $a, b$ and $c$ and $a$ branch point $f$. Let $g$ be the intersection point of the lines (af) and (bc). Then

$$
\frac{\pi}{3}<\angle a g b, \angle a g c<\frac{2 \pi}{3}
$$

Proof. The angle $\angle a g b$ is an exterior angle of the triangle formed by the bisector of the angle $\angle b f c$ and the lines $(b c)$ and $(f c)$. Since the angle at $f$ in this triangle is equal to $\pi / 3$, the angle $\angle a g b$ is strictly greater than $\pi / 3$. Similar reasoning for $\angle a g c$ and the identity $\angle a g b+\angle a g c=\pi$ imply the required inequalities.

Observation 4.6.1. A full Steiner pseudo-network with three endpoints is a regular tripod.

### 4.6.3 Structural properties of minimizers in the annulus $N \backslash N_{r}$

Recall that we work in the setting from Subsection 4.6.1. The proofs of the next few lemmas are essentially contained in the paper [18], but we will rewrite them almost verbatim for our more general case.

Lemma 4.6.4. The inclusion $\Sigma \subset N$ holds .
Proof. Assume the contrary and consider the projection of the closure of a connected component of the set $\Sigma \backslash N$ onto the (closed convex) set $N$. It is well known (see Chapter 1.2 in [114]) that the projection $\operatorname{pr}(x)$ of a point $x$ onto a convex body is defined, and also that

$$
|\operatorname{pr}(x)-\operatorname{pr}(y)| \leq|x-y| .
$$

Thus, the length does not increase after the projection of a connected closed set onto a convex set. Obviously, equality is achieved only if the set and its image are parallel segments.

If there are at least two connected components of the set $\Sigma \backslash N$, then none of them is a segment parallel to the corresponding segment $N$. If the only component of $\Sigma \backslash N$ is a segment parallel to the segment $N$, then $\Sigma \backslash N=\Sigma$. Hence $M$ does not lie in $B_{r}(\Sigma)$, since $N=\operatorname{conv}(M)$ contains a ball of radius $R>r$, where $R$ is defined in Section 4.6.1, and the width of $B_{r}(\Sigma)$ in the direction orthogonal to $\Sigma$ is equal to $2 r$.

On the other hand, for any $x \in \Sigma \backslash N, y \in M \backslash\{\operatorname{pr}(x)\}$ angle $\operatorname{pr}(x)$ in triangle $x \operatorname{pr}(x) y$ is at least $\pi / 2$, therefore the distance between any pair of points from $\Sigma$ and $M$ does not increase during projection, which means that the energy does not increase. Preservation of connectivity follows from the fact that each connected component of the set $\Sigma \backslash N$ remains connected under projection, and the set $\Sigma$ inside $N$ does not change. The resulting contradiction with the optimality of $\Sigma$ finishes the proof.

Consider the closure of an arbitrary connected component $\Sigma \backslash N_{r}$; denote it by $S$ and reserve this designation in the current section. Let us call points from $S \cap M_{r}$ entering points. We will call the continuous image of an (open, half-open, closed) segment (respectively open, half-open, closed) arc. From the connectedness and closedness of $S \subset\left(N \backslash \operatorname{Int} N_{r}\right)$ it follows that $\overline{B_{r}(S)} \cap M$ is a closed
arc; let us denote it $Q(S)$. In what follows we will show that $Q(S) \neq M$. Since $M$ at each point has a strictly positive radius of curvature, it is $C^{1,1}$ smooth and has a tangent at each point. Then by angular measure of the arc $\gamma \subset M$ we mean the directed angle between the tangent rays to $M$ at the ends of $\gamma$. This value is also equal to the limit of the sum of external angles in polychains approximating $\gamma$, with nodes consequentially lying on $M$. The angular measure of an arc $\gamma$ is a non-negative quantity not exceeding $2 \pi$.

We say that a subset $m \subset M$ is covered by a subset $\sigma \subset \Sigma$ if $m \subset \overline{B_{r}(\sigma)}$. Recall that point $n_{q}$ is defined in Section 4.6.1.

Proposition 4.6.1. Let $y \in M$ correspond to two energetic points $x_{1}, x_{2} \in \Sigma \backslash N_{r}$. Then the points $x_{1}$ and $x_{2}$ lie on opposite sides of the line $\left(y n_{y}\right)$.

Proof. Since $M$ is a smooth curve with a radius of curvature greater than $r$ at each point, for each of the points $x_{1}$ and $x_{2}$ the set covered by the point is an arc; let us call these $\operatorname{arcs} Q\left(x_{1}\right)$ and $Q\left(x_{2}\right)$ respectively. Note that the $\operatorname{arcs} Q\left(x_{1}\right)$ and $Q\left(x_{2}\right)$ have a common endpoint $y$. Suppose that $x_{1}$ and $x_{2}$ lie on the same side of the line $\left(y n_{y}\right)$, then one of the arcs is a subset of the other, without loss of generality, $Q\left(x_{2}\right) \subset Q\left(x_{1}\right)$. If $Q\left(x_{2}\right)$ is a proper subset of $Q\left(x_{1}\right)$, then for a sufficiently small $\rho>0$ the arc covered by the set $B_{\rho}\left(x_{2}\right) \cap \Sigma$ is contained in $Q\left(x_{1}\right)$, that is, $x_{2}$ is not an energetic point. So $Q\left(x_{1}\right)=Q\left(x_{2}\right)=: Q$.

Let for any $\rho>0$ set $B_{\rho}\left(x_{1}\right) \cap \Sigma$ covers a larger arc than $Q$. Then, similarly to the previous argument, $x_{2}$ is not energetic. Finally, if there is $\rho_{1}>0$ such that $B_{\rho_{1}}\left(x_{1}\right) \cap \Sigma$ covers exactly $Q$, then $x_{1}$ is not an energetic point, because the entire $\operatorname{arc} Q$ is covered by point $x_{2}$.

Lemma 4.6.5. Under the assumptions of this section, $Q(S) \neq M$.
Proof. If $Q(S)=M$, then $S$ is connected and covers $M$, hence $\Sigma=S$. Then $S=\overline{S \backslash N_{r}}$, and since $S$ is the closure of the connected component $\Sigma \backslash N_{r}=S \backslash N_{r}$, the set $S \backslash N_{r}$ is connected.

Let us assume that there are two different corresponding points $y_{1}, y_{2} \in M$ (they can correspond to either the same energetic point or different ones). Then the ring $N \backslash N_{r}$ is divided by the balls $B_{r}\left(y_{1}\right)$ and $B_{r}\left(y_{2}\right)$ into two connected components $C_{1}$ and $C_{2}$. Since $\Sigma \cap B_{r}\left(y_{i}\right)=\emptyset, i=1,2$, and $S \backslash N_{r}$ is connected, the set $S$ belongs to the closure of one of the components, say $\overline{C_{1}}$. Then $\Sigma$ does not cover one of the two $\operatorname{arcs}\left[y_{1} y_{2}\right]$ in $M$. Therefore, there is at most one corresponding point.

So, all energetic points in $S$ correspond to a single point $y$. Then, according to Proposition 4.6.1, any two energetic points in $S \backslash N_{r}$ lie on opposite sides of the line ( $y n_{y}$ ), which means there are no more than two energetic points in $S \backslash N_{r}$. Also, the point $n_{y}$ can belong to $S$ and be energetic. Thus, there are no more than three energetic points in $S$, and all of them lie on the circle $\partial B_{r}(y)$. Then $\Sigma=S$ is a Steiner tree on the set of its energetic points, and if there are at least two of these points, $\Sigma$ intersects $B_{r}(y)$, which contradicts to the definition of the corresponding point. This means that $\Sigma$ contains unique energetic point, so $\Sigma$ is itself a point. But the set $M$ has a radius of curvature greater than $r$ at each point, which means it cannot be covered by one point, a contradiction.

Having proved Lemma 4.6.5, we can consider the ends of $Q(S)$; let us call them $e_{1}$ and $e_{2}$. Directly from the definitions $B_{r}\left(e_{i}\right) \cap S=\emptyset$, where $i=1,2$. Consider the arc $Z_{Q(S)} \subset M_{r}$, consisting of all possible points $n_{q}$, where $q \in Q(S)$. By definition, $S$ belongs to the (closed) domain $D=D(S)$ bounded by $\partial B_{r}\left(e_{1}\right), \partial B_{r}\left(e_{2}\right)$, the arc $Z_{Q(S)}$ and sometimes arc $Q(S)$, see Fig. 4.10. It is easy to see that the angular measures of the arcs $\partial B_{r}\left(e_{1}\right) \cap D, \partial B_{r}\left(e_{2}\right) \cap D$ do not exceed $\pi / 2$.

The following lemma essentially appeared in [18] (the proof of these statements does not use the additional requirement $R>5 r$, which is inherited from the main theorem of the paper [18]).


Figure 4.10: Two possibilities for a domain $D \supset S$

Lemma 4.6.6. The following statements are true.
(i) Let $x \in S$ be an energetic point. Then an arbitrary corresponding point $y(x)$ belongs to $\left\{e_{1}, e_{2}\right\}$.
(ii) Each of the ends $e_{1}, e_{2}$ corresponds to at most one energetic point from $S \backslash N_{r}$.
(iii) The set of non-isolated energetic points $\mathcal{E}_{\Sigma}$ is a subset of $M_{r}$.
(iv) The set $S$ is a locally minimal tree for its entering points and energetic points.
(v) The set $S \backslash N_{r}$ contains one or two energetic points.
(vi) Let the energetic point $x \in S \backslash N_{r}$ have a unique corresponding point $y(x)$. Then
(a) if $x$ has degree 1 (that is, by item (iii) it is the end of some segment $[z x] \subset \Sigma$ ), then $z, x$ and $y(x)$ lie on the same line;
(b) if $x$ has degree 2 (that is, by item (iii) it is the end of the segments $\left[z_{1} x\right],\left[x z_{2}\right] \subset \Sigma$ ), then the ray $[y(x) x)$ contains the bisector of angle $z_{1} x z_{2}$.

Proof. Assume the contrary to item (i). Then $B_{r}(y) \cup\left\{n_{y}\right\}$ divides $D$ into two non-empty regions, with $\partial B_{r}\left(e_{1}\right) \backslash B_{r}(y)$ and $\partial B_{r}\left(e_{2}\right) \backslash B_{r}(y)$ lie in different regions. But the set $S$ should intersect both $\operatorname{arcs} \partial B_{r}\left(e_{i}\right)$, since $e_{1}, e_{2} \in Q(S)$, and therefore $S$ contains $n_{y}$. But then $S \backslash\left\{n_{y}\right\}$ consists of two connected components, which contradicts the definition of $S$.

Assume the contrary to item (ii). Then $\partial B_{r}\left(e_{i}\right) \cap D$ contains two energetic points $x_{1}, x_{2} \in S \backslash N_{r}$, which contradicts the Proposition 4.6.1.

Let the point $x \in \Sigma \backslash N_{r}$ belong to $\mathcal{E}_{\Sigma}$. Since items (i) and (ii) imply that each connected component $\Sigma \backslash N_{r}$ has at most two energetic points, there is an infinite set of components containing a point from an arbitrarily small neighborhood $x$. This contradicts the finiteness of the length of $\Sigma$. Lemma 4.6.1 (ii) completes the proof of item (iii).

Note that the neighborhood of each point of the Steiner part of the set $S \backslash M_{r}$ is a segment or a regular tripod. Theorem 4.2.1(i) and the connection of $S$ complete the proof of item (iv).

It immediately follows from items (i) and (ii) that there are no more than two energetic points in the set $S \backslash N_{r}$. Let us assume that $S \backslash N_{r}$ does not contain energetic points. Then $S$ is a locally minimal tree for vertices on $M_{r}$, hence $S \subset N_{r}$; the resulting contradiction shows the validity of item (v).

Suppose the contrary to item (vi)(a), then the point $t:=\partial B_{r}(y(x)) \cap[y(x) z]$ is different from $x$. But then the set $(\Sigma \backslash[x z]) \cup[z t]$ is connected, covers $M$ and has a length strictly less than that of $\Sigma$, a contradiction.

Suppose the contrary to item (vi)(b), that is, that $\angle z_{1} x y(x) \neq \angle y(x) x z_{2}$. Let $L$ be the tangent line to $B_{r}(y(x))$ at point $x$. Then

$$
\mathcal{H}^{1}\left(\Sigma \cap B_{\varepsilon}(x)\right)-\mathcal{H}^{1}\left(\left[z_{1}^{\prime} t\right] \cup\left[t z_{2}^{\prime}\right]\right)=\Omega(\varepsilon),
$$

where $z_{1}^{\prime}$ and $z_{2}^{\prime}$ are the intersections of the segments $\left[z_{1} x\right]$ and $\left[z_{2} x\right]$ with $\partial B_{\varepsilon}(x)$, and $t$ is such a point on the line $L$ such that $\angle z_{1}^{\prime} t x+\angle x t z_{2}^{\prime}=\pi$ (see the left side of Fig. 4.11). On the other hand, $\operatorname{dist}\left(t, \partial B_{r}(y(x))\right)=O\left(\varepsilon^{2}\right)$; let us define $v \in \partial B_{r}(y(x))$ as the point for which $\operatorname{dist}\left(t, \partial B_{r}(y(x))\right)=$ dist $(t, v)$ holds. Then the set $\left(\Sigma \backslash B_{\varepsilon}(x)\right) \cup\left[z_{1}^{\prime} t\right] \cup\left[z_{2}^{\prime} t\right] \cup[t v]$ is connected, covers $M$ and has length strictly less than $\Sigma$. The resulting contradiction completes the proof of the lemma.

Remark 4.6.4. Note that if the entering point $x$ is energetic, then by Lemma 4.6.6(i) one has $x=n_{e_{1}}$ or $x=n_{e_{2}}$. Moreover, such a point $x$, in contrast to energetic points that are not also entering points, can belong to a set of non-isolated energetic points.

Since $S$ covers exactly $Q(S)$, the sets $S \cap B_{r}\left(e_{1}\right)$ and $S \cap B_{r}\left(e_{2}\right)$ are empty, and the sets $S \cap \partial B_{r}\left(e_{1}\right)$ and $S \cap \partial B_{r}\left(e_{2}\right)$ are non-empty. For each pair of points $x_{1} \in S \cap \partial B_{r}\left(e_{1}\right), x_{2} \in S \cap \partial B_{r}\left(e_{2}\right)$ consider the area bounded by the arc $Q(S)$, with radii $\left[e_{1} x_{1}\right]$ and $\left[e_{2} x_{2}\right]$ and the path from $x_{1}$ to $x_{2}$ to $S$. Note that the intersection of any two regions of the described type is also a region of the described type. Consequently, there is a minimal such region, let us call it $E(S)$ (one example of the $E(S)$ region is shown on the right side of Fig. 4.11). Note that in view of items (iv) and (i) of the Lemma 4.6.6 and the minimality of $E(S)$, one has $S \cap \operatorname{Int} E(S)=\emptyset$.


Figure 4.11: The left side shows the situation near the point $x$ in part (vi)(b) of the Lemma 4.6.6. On the right is one of the possible options $E(S)$.

Lemma 4.6.7. The region $E(S)$ is convex.
Proof. The boundary $E(S)$ consists of the arc $Q(S)$ and a polychain. Recall that $M$ is convex, so it suffices to show that all interior (with respect to $E(S)$ ) angles of the polychain are at most $\pi$, including the angles between the polychain and the tangents to $Q(S)$ at points $e_{1}$ and $e_{2}$.

The angle between the polychain and the tangent to $Q(S)$ at point $e_{i}$ does not exceed $\pi / 2$, since the corresponding radius $\left[e_{i} x_{i}\right]$ lies inside the right angle between the radius $\left[e_{i} n_{e_{i}}\right.$ ] and tangent to $Q(S)$ at point $e_{i}$.

If the polychain consists of one point, then $x_{1}=x_{2}$ and the interior angle at this point is not greater than $\pi$, since otherwise $\left[e_{1} e_{2}\right] \subset B_{r}\left(e_{1}\right) \cup B_{r}\left(e_{2}\right)$, then there is $\left[e_{1} e_{2}\right] \cap S=\emptyset$, but [ $e_{1} e_{2}$ ] separates $x_{1}$ from $M_{r}$, which contradicts the connectedness of $S$.

Consider the angle of the polychain at point $x_{i}$. If $x_{i}$ is not energetic, then its neighborhood is a segment tangent to the circle $\partial B_{r}\left(e_{i}\right)$, and the angle of interest to us is equal to $\pi / 2$. If $x_{i}$ is an
energetic point of degree 1 , then by Lemma 4.6.6(vi)(a) the angle is equal to $\pi$. If $x_{i}$ is an energetic point of degree 2, then by Lemma 4.6.6(vi)(b) and in view of the condition $S \cap \operatorname{Int} E(S)=\emptyset$ the angle is not greater than $2 \pi / 3$.

It remains to deal with the angles of the polychain at the internal vertices. By Lemma 4.6.6(i) the polychain does not contain energetic points, except, possibly, the ends. This means that all internal (relative to $E(S)$ ) angles of the polychain at internal vertices are equal to $2 \pi / 3$ or $4 \pi / 3$. But in the latter case $S \cap \operatorname{Int} E(S) \neq \emptyset$, which is impossible.

The following theorem is the main statement of this subsection.
Theorem 4.6.1 (Cherkashin-Gordeev-Strukov-Teplitskaya [24]). Let $S$ be the closure of the connected component $\Sigma \backslash N_{r}$. Then
(i) convex hull of $S$ is a segment, a triangle or a quadrilateral, the vertices of $\operatorname{conv}(S)$ are energetic points of $S \backslash N_{r}$ and entering points of $S$, with at most two points of each type;
(ii) $S$ contains at most three entering points.

Proof of item (i). If $S$ has less than three entering points, then the assertions of item (i) immediately follow from Lemma 4.6 .6 (iv) and (v): $S$ is a local Steiner tree whose only entering points and energetic points are not a convex combination of points from the neighborhood, and there are at most two energetic points $S \backslash N_{r}$.

By Lemma 4.6.5 the set $A_{0}=\overline{B_{r}(Q(S))} \cap M_{r}$ is an arc. Let us consider the largest subarc $A \subset A_{0}$ by inclusion, both ends of which are entering points $S$ (we will call them extreme). Let us consider an arbitrary entering point $x$ contained in the interior of this arc (there is one, because we assumed that there are at least three entering points), and draw a tangent $L$ to $M_{r}$ in it. By Remark 4.6.4 point $x$ is not energetic, therefore its neighborhood in $\Sigma$ is a segment or tripod. By definition of $S$, a neighborhood $x$ in $S$ is a segment not belonging to $L$. Due to the convexity of $N_{r}$, the tangent $L$ does not intersect Int $N_{r}$. This means that the set $S$ intersects $L$ on both sides of $x$, say at points $t_{1}$ and $t_{2}$. Then the entering point $x$ is a convex combination of points $t_{1}$ and $t_{2}$. Thus, we have proven that all entering points, except the two extreme ones, lie inside the convex hull $S$.

Proof of item (ii). Let us assume that there are at least three entering points and define the arc $A \subset M_{r}$ similarly to item (i); let us denote its ends $w_{1}$ and $w_{2}$. By Remark 4.6.4 if the entering point is energetic, then it is inevitably extreme.

Let $w$ be an arbitrary entering point from the interior of $A$. Let us denote by $R=R(w)$ the ray starting at the point $w$ and extending beyond the point $w$ the segment $S$ containing it.

Lemma 4.6.8. The ray $R$ intersects the segment $\left[w_{1} w_{2}\right]$.
Proof. If the ray $R$ does not intersect the segment $\left[w_{1} w_{2}\right]$, then it either touches $M_{r}$ at point $w$ or intersects $M_{r}$ at some point $u \in A$. Let us recall that the point $w$ cannot be energetic according to the Remark 4.6.4, hence the tangent case contradicts the connectedness of $S \backslash N_{r}$.

Thus, the ray $R$ intersects $M_{r}$ at some point $u \in A$. Since $w$ belongs to the Steiner part of the minimizer, its neighborhood in $\Sigma$ is a segment or tripod. In the first case, there is a single connected component $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$ containing $w$. In view of the Lemma 4.6.1, the closure $T$ of this component is a full Steiner tree, and therefore, in view of the Remark 4.6.1, contains an end vertex different from $w$ in each of the closed half-planes separated by the line $(w u)$. Moreover, all end vertices $T$ belong


Figure 4.12: Illustration to the proof of Proposition 4.6.1 (ii)
to the set $M_{r}$, which means there is a vertex $w^{\prime}$ belonging to the half-open arc $[\breve{w u}] \backslash\{w\} \subset A \subset M_{r}$; Moreover, $w^{\prime}$ does not belong to the set $S$ due to the absence of cycles in $\Sigma$. The point $w^{\prime}$ is not energetic, since in this case $B_{r}\left(y\left(w^{\prime}\right)\right) \cup\left\{w^{\prime}\right\}$ divides the set $D$, and hence $S$, into two non-empty parts; on the other hand, $B_{r}\left(y\left(w^{\prime}\right)\right) \cap \Sigma=\emptyset=B_{r}\left(y\left(w^{\prime}\right)\right) \cap S$ and $w^{\prime} \notin S$. Therefore $w^{\prime}$ is the entering point of the closure of some connected component of the set $\Sigma \backslash N_{r}$; let us call this closure $S^{\prime}$. Since $\Sigma$ does not contain cycles, $D\left(S^{\prime}\right) \subset D(S)$ and $Q\left(S^{\prime}\right) \subset Q(S)$. The last inclusion implies the absence of energetic points in $S^{\prime}$, which contradicts Lemma 4.6.6(v).

The case of a tripod is analyzed in exactly the same way.
For each entering point $w \notin\left\{w_{1}, w_{2}\right\}$ we denote by $i(w)$ the point of intersection of the ray $R(w)$ with the segment $\left[w_{1} w_{2}\right]$. We also put $i\left(w_{1}\right)=w_{1}, i\left(w_{2}\right)=w_{2}$. Let $\mathcal{S t}(S)$ denote the union of $S$ and all segments $[w i(w)]$. By Lemma 4.6.6(iv) the set $\mathcal{S} t(S)$ is the union of several full Steiner pseudo-networks. Let us consider two cases.
(a) Let the set $S \backslash N_{r}$ have one energetic point $x$. By Lemma 4.6.6(iii) $x \in \mathcal{X}_{\Sigma}$.

- Let $x$ have degree 1. In this case, $\mathcal{S} t(S)$ is a full Steiner pseudo-network. By Lemma 4.6.2 applied to $\mathcal{S} t(S)$ and the line $\left(w_{1} w_{2}\right)$, the pseudo-network $\mathcal{S t}(S)$ intersects $\left(w_{1} w_{2}\right)$ at most twice, since $x$ is the only endpoint of $\mathcal{S} t(S)$ outside $\left(w_{1} w_{2}\right)$. Then, by Lemma 4.6.8, the set $S$ has at most two entering points.
- Let $x$ have degree 2 . Then the sets $S \backslash\{x\}$ and $\mathcal{S} t(S) \backslash\{x\}$ have exactly two connected components; let us denote by $S_{1}, S_{2}$ and $\mathcal{S} t_{1}, \mathcal{S} t_{2}$, respectively, their closures. Obviously, $\mathcal{S} t_{1}$ and $\mathcal{S} t_{2}$ are full Steiner pseudo-networks, with $\mathcal{S} t_{1} \cup \mathcal{S} t_{2}=\mathcal{S} t(S)$.
Let us show that each of the sets $S_{1}$ and $S_{2}$ has at most two entering points. Let us assume the opposite: without loss of generality, let the set $S_{1}$ have at least three entering points. Then the full Steiner pseudo-network $\mathcal{S} t_{1}$ has at least three endpoints on the line ( $w_{1} w_{2}$ ) and only one outside this line. Contradiction with Lemma 4.6.2.
Thus, it is necessary to exclude only the situation in which both sets have two entering points, that is, according to Observation 4.6.1 they are tripods.
For each tripod $\mathcal{S} t_{i}$ (where $i \in\{1 ; 2\}$ ) we denote its branch point by $v_{i}$. Let $g_{i}$ be the intersection point of the lines $\left(x v_{i}\right)$ and $\left(w_{1} w_{2}\right)$. Then by Lemma 4.6.3

$$
\frac{\pi}{3}<\angle x g_{i} w_{j}<\frac{2 \pi}{3}
$$

where $i, j \in 1,2$. But then $\angle v_{1} x v_{2}=\angle g_{1} x g_{2}=\pi-\angle x g_{1} w_{2}-\angle x g_{2} w_{1}<\pi / 3<2 \pi / 3$, which contradicts the local minimality of $S$, namely the Lemma 4.6.6(iv).
Thus, in this case $S$ cannot have more than three entering points.
(b) Let the set $S \backslash N_{r}$ have two energetic points $x_{1}$ and $x_{2}$. Let us show that $\partial E(S) \cap S$ is the path between $x_{1}$ and $x_{2}$ in $S$; let us call this path $P$. Assuming the contrary, then, without loss of generality, there exists a point $x \in S \cap \partial B_{r}\left(e_{1}\right)$ such that $\angle x e_{1} n_{e_{1}}>\angle x_{1} e_{1} n_{e_{1}}$. By Lemma 4.6.7 and the fact that $S \cap \operatorname{Int} E(S)=\emptyset$ the segment $\left[e_{1} e_{2}\right]$ does not intersect $S$, that is, from the inequality $\angle x e_{1} \underline{n_{e_{1}}>\angle x_{1} e_{1} n_{e_{1}} \text { it follows that } Q\left(x_{1}\right) \text { is a proper subset of } Q(x) \text {, where } Q\left(x^{*}\right) ~}$ denotes the arc $\overline{B_{r}\left(x^{*}\right)} \cap M$. Then $x_{1}$ is not energetic, a contradiction.

By Lemma 4.6.7 the domain $E(S)$ is convex, therefore the sum of the exterior angles of the broken part $\partial E(S)$ and the angular measure of the $\operatorname{arc} Q(S)$ is equal to $2 \pi$.
By Lemma 4.6.6(iv), the vertices of $\partial E(S)$ from the set $S$ are $x_{1}, x_{2}$, and branch points. The external angle at the branch points is $\pi / 3$. By Lemma 4.6.6(vi) at energetic points of degree 1 the external angle is equal to 0 , and at energetic points of degree 2 it is at least $\pi / 3$. Recall that the sum of the external angles $\partial E(S)$ at the vertices $e_{1}$ and $e_{2}$ is at least $\pi$ in total. If the sum of these angles is $\pi$, then $x_{1}=n_{e_{1}}$ and $x_{2}=n_{e_{2}}$, which contradicts the fact that $x_{1}, x_{2} \in S \backslash N_{r}$. Also, the angular measure $Q(S)$ is non-negative. Therefore, the exterior angle is nonzero in at most two vertices of $S$.
Let $x_{i}$ (where $i \in\{1 ; 2\}$ ) have degree 1 and be connected by the segment $\Sigma$ to $x_{3-i}$. Then we can remove the segment $\left[x_{1} x_{2}\right]$ from $\mathcal{S t}(S)$ and, by Lemma 4.6.2, the remaining full Steiner pseudo-network $\overline{\mathcal{S} t(S) \backslash\left[x_{1} x_{2}\right]}$ has no more than two intersection points with ( $w_{1} w_{2}$ ), which means $S$ has at most two entering points.
In the remaining situation, if the point $x_{i}$ is a point of degree 1 , then the segment $\left[x_{i} x_{3-i}\right]$ does not belong to $\Sigma$. In this case, we denote by $v_{i}$ the branch point to which point $x_{i}$ is connected by the segment $\Sigma$. If $x_{i}$ has degree 2 , we define $v_{i}=x_{i}$ (see Fig. 4.12).
Regardless of the degrees of the points $x_{i}$, the points $v_{1}$ and $v_{2}$ are vertices of $P$ at which $\partial E(S)$ has a nonzero exterior angle. Since $E(S)$ contains at most 2 such vertices from $S$, either $v_{1}=v_{2}$, or $\Sigma$ contains a non-degenerate segment $\left[v_{1} v_{2}\right]$.
If $v_{1}=v_{2}$, then removing the segments $\left[v_{1} x_{1}\right]$ and $\left[v_{2} x_{2}\right]$ from $\mathcal{S} t(S)$ allows us to apply the Lemma 4.6.2 to the line $\left(w_{1} w_{2}\right)$ and the full pseudo-network $\mathcal{S} t(S) \backslash\left(\left[v_{1} x_{1}\right] \cup\left[v_{2} x_{2}\right]\right)$, and find that in this case $S$ has at most two entering points.
If $\Sigma$ contains a non-degenerate segment $\left[v_{1} v_{2}\right]$, then three cases are possible.

- Let both points $x_{1}$ and $x_{2}$ be of degree 1 , then $\mathcal{S} t(S)$ is a full Steiner pseudo-network. Application of Lemma 4.6.2 to the line $\left(w_{1} w_{2}\right)$ shows that $\mathcal{S} t(S)$ intersects $w_{1} w_{2}$ in no more than four points, which means $S$ contains no more than 4 entering points, and if there are 4 entering points, then equality is achieved in the lemma, and by Remark 4.6.3 the rays $\left[v_{1} x_{1}\right)$ and $\left[v_{2} x_{2}\right)$ are codirected. This contradicts the fact that $P$ contains exactly 2 branch points $v_{1}$ and $v_{2}$.
- Let both points $x_{1}$ and $x_{2}$ have degree 2 and, accordingly, coincide with points $v_{1}$ and $v_{2}$. Then $P=\left[x_{1} x_{2}\right]=\left[v_{1} v_{2}\right]$. Then removing $\left[v_{1} v_{2}\right] \backslash\left\{v_{1}, v_{2}\right\}$ splits $\mathcal{S} t(S)$ into two full pseudo-networks $\mathcal{S} t_{1} \ni v_{1}$ and $\mathcal{S} t_{2} \ni v_{2}$, in neither of which, In view of Lemma 4.6.2, there cannot be more than two entering points. Let each of them have two entering points, then by Observation 4.6.1 $\mathcal{S} t_{1}$ and $\mathcal{S} t_{2}$ are tripods. Let us denote by $u_{i}$ the branching point of the corresponding tripod, and by $g_{i}$ the intersection of the lines $\left(v_{i} u_{i}\right)$ and $\left(w_{1} w_{2}\right)$. By Lemma 4.6.3 the angles $g_{1}$ and $g_{2}$ of the quadrilateral $g_{1} g_{2} v_{2} v_{1}$ are strictly greater than $\pi / 3$, and by the property of the local Steiner tree the angles $v_{2}$ and $v_{1}$ are at least $2 \pi / 3$, which when summed gives an immediate contradiction.
- Let one of the points (without loss of generality, $x_{1}$ ) have degree 1, and the other have degree 2 (see Fig. 4.12). Removing the set $\left[v_{1} v_{2}\right] \cup\left[v_{1} x_{1}\right] \backslash\left\{v_{1}, v_{2}\right\}$ from $\mathcal{S t}(S)$ splits $\mathcal{S} t(S)$ into two pseudo-networks $\mathcal{S} t_{1} \ni v_{1}$ and $S t_{2} \ni v_{2}$. Then, by Lemma 4.6.2, each of the pseudo-networks $\mathcal{S} t_{1}$ or $\mathcal{S} t_{2}$ has at most two entering points, and if they both have two entering points, by Observation 4.6.1 they are tripods. Let us denote by $u_{i}$ the branching point of the corresponding tripod, and by $g_{i}$ the intersection of the lines $\left(v_{i} u_{i}\right)$ and $\left(w_{1} w_{2}\right)$. By Lemma 4.6.3 the angles $g_{1}$ and $g_{2}$ of the quadrilateral $g_{1} g_{2} v_{2} v_{1}$ are strictly greater than $\pi / 3$, by the property of the local Steiner tree the value of the angle $v_{2}$ is at least $2 \pi / 3$, and finally the angle $v_{1}$ is equal to $2 \pi / 3$. The resulting contradiction completes the proof of the theorem.

Now we are going to focus on the relations between the components. Denote the set of closures of a connected components of $\Sigma \backslash N_{r}$ by $V_{C}(G)$ and the set of maximal (with respect to the inclusion) arcs of $\Sigma \cap M_{r}$ of positive length by $V_{A}(G)$ (further $V_{A}(G) \cap V_{C}(G)$ will be associated with a subset of the vertex set of a graph).

Lemma 4.6.9. If $S \in V(G):=V_{C}(G) \sqcup V_{A}(G)$ does not reduce to a point, then

$$
Q_{S} \not \subset \bigcup_{S^{\prime} \in V(G) \backslash\{S\}} Q_{S^{\prime}}
$$

Proof of Lemma 4.6.9. The fact that $S$ has an energetic point immediately implies that $Q_{S}$ does not belong to the union of $Q_{S^{\prime}}$ over $S^{\prime} \in V(G) \backslash\{S\}$. Suppose the contrary, i.e. that $S$ has no energetic point.

If $S$ is the closure of a connected component of $\Sigma \backslash N_{r}$, then by Lemma 4.7.3 $S$ is a locally minimal tree for its entering points, but $m(S) \leq 2$, hence $S$ is a segment with endpoints on $M_{r}$, which is impossible for a connected component of $\Sigma \backslash N_{r}$.

If $S$ is a non degenerate arc $[\breve{b c}]$, then $[\breve{b c}] \subset S_{\Sigma}$, which is impossible by the definition of $V_{A}(G)$.
Lemma 4.6.10. The set $V(G)=V_{C}(G) \sqcup V_{A}(G)$ is finite.
Proof of Lemma 4.6.10. Suppose the contrary. Consider an arbitrary $\varepsilon>0$ (which later will be chosen sufficiently small). First, note that Lemma 4.6 .9 implies that every point of $M$ belongs to at most two different $\operatorname{arcs} q_{S}$, where $S \in V(G)$ (otherwise, there are three arcs of $M$ containing a point $x \in M$, so one of them is contained in the union if others, which is impossible by Lemma 4.6.9. . Thus the sum of $\mathcal{H}^{1}\left(q_{S}\right)$ over $V(G)$ is at most $2 \mathcal{H}^{1}(M)$, and therefore there is only a finite number of connected components and $\operatorname{arcs}$ with $\mathcal{H}^{1}\left(q_{S}\right) \geq \varepsilon$. Denote by $V_{\varepsilon}(G)$ the infinite set of such $S \in V(G)$ that $\mathcal{H}^{1}\left(q_{S}\right)<\varepsilon$.

Obviously, if $V(G)$ is an infinite set, then $V_{C}(G)$ is an infinite set. Let us show that there are infinitely many chords of $M_{r}$ in $\Sigma$ intersecting $\operatorname{Int}\left(N_{r}\right)$ (if $N$, and hence $N_{r}$, is strictly convex then in fact every chord of $M_{r}$ intersects $\operatorname{Int}\left(N_{r}\right)$ ). Suppose the contrary. Then $\Sigma \backslash \operatorname{Int}\left(N_{r}\right)$ has a finite number of connected components; but $V_{C}(G)$ is infinite, hence there are components containing infinitely many elements of $V_{C}(G)$; let $K$ be one of these components containing at least five different elements of $V_{C}(G)$. Obviously, $q_{K}:=\overline{B_{r}(K)} \cap \Sigma$ is connected. By Lemma 4.6.9 $K \backslash M_{r}$ contains 5 energetic points, such that they belong to different elements of $V_{C}(G)$. Call them $W_{1}, W_{2}, W_{3}, W_{4}$, $W_{5}$ such that $Q_{W_{1}}, Q_{W_{2}}, Q_{W_{3}}, Q_{W_{4}}, Q_{W_{5}} \in q_{K}$ belong to $M_{r}$ in the natural (clockwise) order. Then $B_{r}\left(Q_{W_{i}}\right) \cap \Sigma=\emptyset, i=1, \ldots, 5$ and therefore $K$ should contain the points $I_{2}, I_{3}, I_{4} \in M_{r}$ such that

$$
\operatorname{dist}\left(Q_{W_{2}}, I_{2}\right)=\operatorname{dist}\left(Q_{W_{3}}, I_{3}\right)=\operatorname{dist}\left(Q_{W_{4}}, I_{4}\right)=r
$$

(because $K \backslash I_{j}$ must be disconnected, $j=2,3,4$ ). Consider the path between $I_{2}$ and $I_{4}$ in $K$. It should coincide with $\left[I_{2} I_{4}\right] \subset M_{r}$, otherwise we reduce the length of $\Sigma$, projecting the path on $M_{r}$. So $W_{3}$ should belong to $M_{r}$ which is impossible by the choice of $W_{i}, i=1, \ldots, 5$ and gives the desired contradiction. Thus the set $C h$ of chords of $M_{r}$ in $\Sigma$ intersecting $\operatorname{Int}\left(N_{r}\right)$ is infinite.

There is at most a finite number of chords of length at least $\varepsilon$ because $\mathcal{H}^{1}(\Sigma)$ is finite. Let us exclude from the infinite set $C h$ a finite set of chords of length at least $\varepsilon$ and a finite set of chords adjacent to a component not in $V_{\varepsilon}(G)$; denote the resulting set by $C h^{\prime}$ : chords in $C h^{\prime}$ are adjacent only to the elements of $V_{\varepsilon}(G)$ and have length strictly less than $\varepsilon$. Let us show that any of the chords in $C h^{\prime}$ connects components without Steiner points. Suppose the contrary. The following three cases have to be considered:
(i) A chord in $C h^{\prime}$ is adjacent to a connected component $S \in V_{\varepsilon}(G)$ with $m(S)=2$ containing a Steiner point. Then the angle between the entering segments of the component is at most $2 \pi / 3$ (in fact, it must be between $\pi / 3$ and $2 \pi / 3$ ). Recall that $\mathcal{H}^{1}\left(q_{S}\right)<\varepsilon$, hence by the triangle inequality $S$ is a subset of an $\varepsilon$-neighbourhood of $M_{r}$ (otherwise dist $(x, y) \leq r-\varepsilon$ for some $x \in S, y \in M$, so $B_{\varepsilon}(y) \cap M \subset q_{S}$ which contradicts $\mathcal{H}^{1}\left(q_{S}\right)<\varepsilon$ ). So, when $\varepsilon$ is sufficiently small, recalling smoothness of $M_{r}$ one has that one of the entering segments has angle with $M_{r}$ at least $\pi / 12$. It implies that the entering point $I$ of this segment is not energetic, so by Lemma 4.7.2 its neighbourhood is a segment and it is an end of a chord $[I J] \subset \Sigma$ of $M_{r}$. So by the constraint on the radius of curvature of $M$ chord $[I J]$ has length more than $\varepsilon$, which gives a contradiction with the assumption that our chord is in $C h^{\prime}$.
(ii) A chord in $C h^{\prime}$ is adjacent to a connected component $S \in V_{\varepsilon}(G)$ with $m(S)=1$ containing a Steiner point. Then it has the combinatorial type (b) on Fig. 4.33. Let us consider the triangle $\triangle Q C I$, where $Q$ is an end of $q_{S}, C$ is the branching point of $S, I$ is the entering point of $S$. Since $\angle Q C I=2 \pi / 3$, we have $\angle Q I C \leq \pi / 3$, so the angle between the entering segment [ $C I$ ] and $M_{r}$ is at least $\pi / 6$. Then again the chord $[I J]$ has length more than $\varepsilon$, that contradicts the choice of the chord.
(iii) Finally, a chord in $C h^{\prime}$ is adjacent to an $\operatorname{arc} S \in V_{\varepsilon}(G)$ containing a Steiner point $x$. Then $x \in M_{r}$, and $x$ is an end of a chord of $M_{r}$ in $\Sigma$ which forms angle $\pi / 3$ with $M_{r}$. Again by the condition on the radius of curvature of $M_{r}$ and with the choice of $\varepsilon$ sufficiently small, this chord has length more than $\varepsilon$ which is impossible.

Let us consider any chord $\left[I_{1} I_{2}\right] \in C h^{\prime}$, such that it connects some components from $V_{\varepsilon}(G)$ (which do not have Steiner points as proven). Note that the set $\left.] I_{2} I_{1}\right) \cap \Sigma($ resp. $\left.] I_{1} I_{2}\right) \cap \Sigma$ ) contains an energetic point (it may coincide with $I_{1}\left(I_{2}\right)$; if $I_{1}\left(I_{2}\right)$ is not energetic, an energetic point on $\left.] I_{2} I_{1}\right) \cap \Sigma$ (resp. ] $\left.I_{1} I_{2}\right) \cap \Sigma$ ) exists by Lemma 4.7.2 and the absence of Steiner points in the considered connected components and arcs); denote the nearest to $I_{1}$ (resp. $I_{2}$ ) energetic point of $\left.] I_{2} I_{1}\right) \cap \Sigma($ resp. $\left.] I_{1} I_{2}\right) \cap \Sigma$ ) by $W_{1}$ (resp. $W_{2}$ ).

Consider the region $P$ bounded by the segments [ $W_{1} Q_{W_{1}}$ ], [ $W_{2} Q_{W_{2}}$ ], [ $W_{1} W_{2}$ ] and the lesser arc $\left[Q_{W_{1}} Q_{W_{2}}\right]$ of $M$. Let us show that the intersection of $\operatorname{Int}(P)$ with $\Sigma$ is nonempty. There are two tangent lines to $M_{r}$ parallel to [ $W_{1} W_{2}$ ]; let $l$ be the nearest line to $\left[W_{1} W_{2}\right.$ ]. Note that $\left[I_{1} I_{2}\right] \in C h^{\prime} \subset$ $C h$, so $\left[I_{1} I_{2}\right] \cap \operatorname{Int}\left(N_{r}\right) \neq \emptyset$ and $l \cap\left[W_{1} W_{2}\right]=\emptyset$. Consider a point $w \in l \cap M_{r}$ and note that $Q_{w}$ is not covered by $\Sigma$, because $\operatorname{dist}\left(Q_{W}, \Sigma\right)=\operatorname{dist}\left(Q_{w},\left[W_{1}, W_{2}\right]\right)>\operatorname{dist}\left(Q_{w}, l\right)=r$. We got a contradiction, so $\operatorname{Int}(P) \cap \Sigma \neq \emptyset$.

Let us pick a point $x \in \operatorname{Int}(P) \cap \Sigma$ and consider the path in $\Sigma$ connecting $x$ with the segment [ $W_{1} W_{2}$ ]. The existence of this path gives that for some $i \in\{1,2\}$ (say, without loss of generality,
$i=1$ ) one has $W_{i}=I_{i}$ (in fact, $] W_{1} W_{2}\left[\subset S_{\Sigma}\right.$, which means that this path connects $x$ with $W_{1}$ without touching $] W_{1} W_{2}$ ], but a neighbourhood in $\Sigma$ of an energetic point of $\Sigma \backslash N_{r}$ is either a single line segment or two line segments with angle at least $2 \pi / 3$, see Fig. 4.33 and 4.34 , and thus $W_{1} \in M_{r}$ ) and $B_{\delta}\left(I_{1}\right) \cap \operatorname{Int}(P) \cap \Sigma \neq \emptyset$ for sufficiently small $\delta>0$. Let $k$ be the tangent line to $M_{r}$ at $I_{1}=W_{1}$. Since $\left|I_{1} I_{2}\right| \leq \varepsilon$, the angle between $k$ and $\left[I_{1} I_{2}\right]$ is $O(\varepsilon)$. Consider an arbitrary point $y \in \partial B_{\delta}\left(I_{1}\right) \cap \operatorname{Int}(P) \cap \Sigma$. Since $B_{r}\left(Q_{I_{1}}\right) \cap \Sigma=\emptyset$ and $\left|y I_{1}\right|=\delta$ the angle between $k$ and $\left[y I_{1}\right]$ is $O(\delta)$. Let $z$ be a projection of $y$ on $\left[I_{1} I_{2}\right]$. Then $\angle y I_{1} z=O(\varepsilon+\delta)$ is the smallest angle (for sufficiently small $\varepsilon, \delta$ ) in right-angled triangle $\Delta y I_{1} z$. Hence one can replace $] I_{1} z[$ by $[z y]$ in $\Sigma$. The new set is still connected, covers $M$ and has strictly lower length than $\Sigma$. We got in this way a contradiction with the optimality of $\Sigma$, concluding the proof.

Note that a singleton of $\Sigma \cap M_{r}$ (a maximal arc $\xi \subset \Sigma \cap M_{r}$ of zero length not contained in the closure of a connected component of $\Sigma \backslash N_{r}$ ) cannot be energetic (by the previous Lemma the union of $Q_{S}$ over $S \in V(G) \backslash \xi$ is closed as a finite union of closed sets, hence it coincides with $M$ because $q_{\xi}=\left\{q_{\xi}\right\}$ ), so a neighbourhood of $\xi$ is a segment or a tripod (the latter is impossible by Lemma 4.7.2). Summing up, every point of $\Sigma \cap M_{r}$ is contained in a maximal arc of $M_{r}$ of positive length or in the closure of a connected component of $\Sigma \backslash N_{r}$. Also by Lemma 4.7 .3 every connected component of $\Sigma \backslash N_{r}$ contains at most 5 segments, thus $\Sigma$ consists of a finite number of segments and arcs of $M_{r}$.

Lemma 4.6.11. Let $[b i] \subset \Sigma$ be a chord of $M_{r}$. Then $i \in S_{\Sigma}$ and moreover there exists such an $\varepsilon>0$ that $\overline{B_{\varepsilon}(i)} \cap \Sigma=\left[i_{1} i_{2}\right]$, for some $i_{1}, i_{2} \in \partial B_{\varepsilon}(i)$.

Proof of Lemma 4.6.11. Note that in $\Sigma$ there are at most two chords of $M_{r}$ ending at $i$. It is true because of the properties of a locally minimal tree: the angle between two segments ending at the same point is greater or equal to $2 \pi / 3$.

Let us show that $i \in S_{\Sigma}$. Assume the contrary: let $i \in G_{\Sigma}$. Then $B_{r}\left(q_{i}\right) \cap \Sigma=\emptyset$. There are two possibilities:
(1) $i \in S$, where $S \in V_{C}(G)$;
(2) $i \in S$, where $S \in V_{A}(G)$ (as mentioned after Lemma 4.6.10 $S$ is non degenerate i.e. does not reduce to a single point $i$ ).

Recall that $\Sigma$ consists of a finite number of segments and a finite number of arcs of $M_{r}$. In the case (1) the smoothness of $M_{r}$, Lemma 4.7.3 and the fact $B_{r}\left(q_{i}\right) \cap \Sigma=\emptyset$ imply that the intersection of a small neighbourhood of $i$ with $S \backslash N_{r}$ is a subset of the tangent line to $M_{r}$ at $i$.

Thus the set $\Sigma \cap B_{\varepsilon}(i) \backslash \operatorname{Int}\left(N_{r}\right)$ is contained in the union of the tangent line $\tau$ to $M_{r}$ at $i$ and the arc $M_{r} \cap \partial B_{\varepsilon}(i)$. Both $\tau \cap B_{\varepsilon}(i)$ and $M_{r} \cap B_{\varepsilon}(i)$ are split by $i$ into 2 segments [ $\left.i e_{1}^{\prime}\right]$, $\left[i e_{2}^{\prime}\right]$ and $2 \operatorname{arcs}\left[i e_{1}\right]$, $\left[i e_{2}\right]$ of $M_{r}$, respectively, where $e_{1}, e_{2}, e_{1}^{\prime}, e_{2}^{\prime} \in \partial B_{\varepsilon}(i)$. we may assume $e_{1}$ in the same halfplane with $e_{1}^{\prime}$ bounded by the normal to $M_{r}$ passing throw $i$. At least one arc and one segment (say, $\left[i e_{1}\right]$ and $\left[i e_{1}^{\prime}\right]$ ) have angle at most $\pi / 2$ with the chord $[i b]$. The cases (i) and (ii) below deal with the situation with nonempty set $\Sigma \cap\left(\left[i e_{1}\right] \cup\left[i e_{1}^{\prime}\right]\right)$. In the remaining cases $\Sigma \cap B_{\varepsilon}(i) \backslash \operatorname{Int}\left(N_{r}\right)$ is a subset of $\left[i e_{2}^{\prime}\right] \cup\left[i e_{2}\right]$ and therefore in (iii)-(vi) we deal with all the possible cases of $B_{\varepsilon}(i) \cap\left[i e_{2}^{\prime}\right]$ and $B_{\varepsilon}(i) \cap\left[i e_{2}\right]$ empty/nonempty:
(i) There is such a segment $[i e] \subset \Sigma$, that $(i e)$ is the tangent line to $M_{r},|i e|=\varepsilon$ and $\angle b i e \leq \pi / 2$;
(ii) There is such an $\varepsilon>0$ and an arc $[\dot{i e}] \subset \Sigma \cap M_{r}$ that $|i e|=\varepsilon$ and $\angle b i e \leq \pi / 2$;


Figure 4.13: The case (iv) from Lemma 4.6.11; (a) the (impossible) part of the minimizer; (b) a better competitor.
(iii) There is such a small $\varepsilon>0$ that $\overline{B_{\varepsilon}(i)} \cap \Sigma$ is equal to $[f i] \cup[i e]$ where $f, e \in \partial B_{\varepsilon}(i),[f i] \subset[b i]$ and $[i e]$ is a subset of the tangent line to $M_{r}$ at point $i$;
(iv) There is such a small $\varepsilon>0$ that $\overline{B_{\varepsilon}(i)} \cap \Sigma$ is equal to $[f i] \cup[\breve{i e}]$, where $f, e \in \partial B_{\varepsilon}(i),[f i] \subset[i b]$ and $[i e] \subset M_{r}$;
(v) There is such a small $\varepsilon>0$ that $\overline{B_{\varepsilon}(i)} \cap \Sigma$ contains $[f i] \cup[i c] \cup[\breve{i d}]$ where $[i c]$ is a subset of the tangent line to $M_{r}$ at point $i,[f i] \subset[b i],[i d] \subset M_{r}$ and $\angle c i d<\pi / 6$;
(vi) There is such an $\varepsilon>0$ that $\overline{B_{\varepsilon}(i)}$ is a subset of chord [ib].
we will show that all these cases are impossible. Let $\xi$ stand for the segment [ie] in the cases (i) and (iii), and for $[\breve{i e}]$ in the cases (ii) and (iv).

CASES (i), (ii): Let $f:=[b i] \cap B_{\varepsilon}(i)$ and $l^{\varepsilon}$ be the lesser arc of $\partial B_{\varepsilon}(i)$ limited by intersections with $\partial B_{r}\left(q_{i}\right)$ and $M_{r}$. It is easy to see that $\mathcal{H}^{1}\left(l^{\varepsilon}\right)=O\left(\varepsilon^{2}\right)$ and $|f i|+\mathcal{H}^{1}(\xi)-\mathcal{H}^{1}(\mathcal{S} t(f, i, e))=$ $c \varepsilon+o(\varepsilon)$ with $c>0$, where $\mathcal{S} t(f, i, e)$ is a Steiner tree connecting points $f, i, e$. Then the length of $\Sigma^{\prime}:=\Sigma \backslash([f i] \cup \xi) \cup l^{\varepsilon} \cup \mathcal{S} t(f, i, e)$ is less than $\mathcal{H}^{1}(\Sigma)$ for sufficiently small $\varepsilon$. Moreover $\Sigma^{\prime}$ is still connected and $F_{M}\left(\Sigma^{\prime}\right) \leq F_{M}(\Sigma)$. This gives us a contradiction with optimality of $\Sigma$.

CASES (iii), (iv): Note that $|f i|=|i e|=\varepsilon$ (see Fig. 4.13(a)), so $\mathcal{H}^{1}(\xi)=\varepsilon+o(\varepsilon)$ when $\varepsilon \rightarrow 0^{+}$, because $M_{r}$ is smooth. Let $h$ be the point of intersection of $\left[e q_{i}\right]$ and $\partial B_{r}\left(q_{i}\right)$ (see Fig. 4.13(b)). Note that $\left(i q_{i}\right)$ is perpendicular to the tangent line to $M_{r}$ at the point $i$. Thus

$$
\begin{aligned}
|e h| & =\left|e q_{i}\right|-\left|q_{i} h\right|=\sqrt{|e i|^{2}+r^{2}}-r=\sqrt{\varepsilon^{2}+r^{2}}-r \\
& =r \sqrt{1+o(\varepsilon)}-r=o(\varepsilon) .
\end{aligned}
$$

Now, since the angle between $\xi$ and the segment $[f i]$ is less than $\pi$, we get

$$
|e f|=\sqrt{2 \varepsilon^{2}-2 \varepsilon^{2} \cos \angle e i f}=\sqrt{2} \varepsilon \sqrt{1-\cos \angle e i f}<2 \varepsilon-c \varepsilon, \text { for some } c>0
$$



Figure 4.14: The case of an outer segment in the proof of Lemma 4.6.12; (a) the (impossible) part of the minimizer; (b) the better competitor.
and therefore

$$
|e h|+|e f|<\mathcal{H}^{1}(\xi)+|i f|=2 \varepsilon+o(\varepsilon)
$$

for sufficiently small $\varepsilon>0$. So we have a contradiction with the optimality of $\Sigma$, because we show that $\left(\Sigma \backslash B_{\varepsilon}(i)\right) \cup[e h] \cup[e f]$ is the better competitor.

CASE (v): Let $h \in[i c)$ be such a point that $(d h) \perp(i c)$. Then the set

$$
\left.\left.\Sigma^{\prime}=\Sigma \backslash\right] \stackrel{i d}{ }\right] \cup[h d]
$$

is still connected, has energy $F_{M}$ not greater than $\Sigma$ and strictly smaller length, since $|h d|<|i d| / 2 \leq$ $\mathcal{H}^{1}([i d]) / 2$. It means $\Sigma^{\prime}$ is the better competitor than $\Sigma$, again a contradiction.

CASE (vi): In this case $S \in V_{A}(G)$ and $S=\{i\}$, which is impossible.
So all cases are impossible and we have a contradiction which implies $i \in S_{\Sigma}$. Because of Lemma 4.7.2 $i$ can not be a Steiner point. Then there exists an $\varepsilon>0$ such that $S_{\Sigma} \cap B_{\varepsilon}(i)$ is a segment.

Lemma 4.6.12. Every (maximal with respect to inclusion) arc $[\breve{b} c] \in V_{A}(G)$ is continued by segments lying on tangent lines to $M_{r}$ in the sense that there exists such an open $U \supset[\breve{b c}]$ that $\Sigma \cap \bar{U}=$ $\left[b^{\prime} b\right] \cup[\breve{b c}] \cup\left[c c^{\prime}\right]$, where $\left[b^{\prime} b\right]$ and $\left[c c^{\prime}\right]$ are subsets of tangent lines to $M_{r}$ at points $b$, c respectively.

Proof of Lemma 4.6.12. Let $\breve{b c}$ be as in the statement being proven.
Suppose that there is a segment $[i j] \subset \Sigma$ such that $i=] b c\left[\cap[i j]\right.$. we claim that $B_{\varepsilon}(i) \cap \Sigma \subset[\breve{b c}]$. In fact, by Lemma $4.6 .11[i j]$ cannot be a part of a chord of $M_{r}$, so $\left.] i j\right] \subset \Sigma \backslash \operatorname{Int}\left(N_{r}\right)$. Note that in this case $i$ is energetic (because $B_{\varepsilon}(i)$ is not a segment or a tripod for every $\varepsilon>0$ ). Hence $B_{r}\left(q_{i}\right) \cup \Sigma=\emptyset$, so $[i j]$ is a part of the tangent line to $M_{r}$ at $i$. Let us choose an $\varepsilon>0$ and set $\left\{d_{1}, d_{2}\right\}:=[\breve{b c}] \cap \partial B_{\varepsilon}(i)$, $e:=[i j] \cap \partial B_{\varepsilon}(i)$. If $\varepsilon>0$ is sufficiently small one of the angles $\angle d_{1} i e, \angle d_{2} i e$ is less than $\pi / 6$ (say $\left.\angle d_{1} i e\right)$. Let $h \in[i j)$ be such a point that $\left(d_{1} h\right) \perp(i j)$. Then the set

$$
\left.\left.\Sigma^{\prime}:=\Sigma \backslash\right] \breve{d}_{1}\right] \cup\left[h d_{1}\right]
$$

is still connected, has energy $F_{M}$ not greater than $F_{M}(\Sigma)$ and strictly smaller length, since $\left|h d_{1}\right|<$ $\left|i d_{1}\right| / 2 \leq \mathcal{H}^{1}\left(\left[i \breve{d}_{1}\right]\right) / 2$. It means that $\Sigma^{\prime}$ is better competitor than $\Sigma$. we got a contradiction, showing thus $B_{\varepsilon}(i) \cap \Sigma \subset[\breve{b c}]$ for $\left.i \in\right] \breve{b} c[$.


Figure 4.15: Picture to Lemma 4.6.12. An end of an arc of $M_{r} \cap \Sigma$ cannot be an endpoint of $\Sigma$.

Let us prove now that $B_{\varepsilon}(b) \backslash[\breve{b c}]$ is a subset of the tangent line to $M_{r}$ at $b$ (the analogous statement for the point $c$ is completely symmetric). By Lemma 4.6.11 there is no chord of $M_{r}$ in $\Sigma$ with endpoint $b$. So the set $B_{\varepsilon}(b) \backslash[b c]$ is a subset of $\Sigma \backslash N_{r}$.
we claim first that $b$ is not an endpoint of $\Sigma$ i.e. $B_{\varepsilon}(b) \backslash[\breve{b c}] \neq \emptyset$. Assume the contrary and recall that $q_{b}, q_{c} \in M$ are such points that $\operatorname{dist}\left(b, q_{b}\right)=\operatorname{dist}\left(c, q_{c}\right)=r$. Then one can set $b_{1}:=\partial B_{\varepsilon}(b) \cap[\breve{b c}]$ and replace $\left[b_{1} b\right]$ by the segment $\left[b_{1} i\right]:=\left[b_{1} q_{b}\right] \backslash B_{r}\left(q_{b}\right)$, producing the competitor of strictly lower length because $[\breve{b c}] \backslash\left[\breve{b_{1} b}\right] \cup\left[b_{1} i\right]=\left[\breve{b_{1} c}\right] \cup\left[b_{1} i\right]$ still covers the $\operatorname{arc}\left[\breve{q_{b}} q_{c}\right]$ of $M$ (when $\varepsilon$ is sufficiently small), see Fig. 4.15 .

Therefore we have proven that for sufficiently small $\varepsilon>0$ the set $B_{\varepsilon}(b) \backslash[\breve{b c}]$ is a nonempty subset of $\Sigma \backslash N_{r}$. If $b$ is energetic then $B_{r}\left(q_{b}\right) \cap \Sigma=\emptyset$, hence $B_{\varepsilon}(b) \backslash[b c]$ is a subset of the tangent line to $M_{r}$ at point $b$ showing the claim. So $b \in S_{\Sigma}$, hence $B_{\varepsilon}(b)$ is a segment or a tripod for sufficiently small $\varepsilon>0$. But the case of a tripod is impossible by Lemma 4.7.2, while the case of a segment is only possible recalling smoothness of $M_{r}$ (and part of $M_{r}$ in a neighbourhood of $b$ is in fact flat).

Summing up, the only segments intersecting $[\breve{b c}]$ are segments tangent to $M_{r}$ at points $b$ and $c$. As a consequence of Lemma 4.6.10 $\Sigma$ consists of a finitely many segments and maximal arcs of $M_{r}$, so when $\varepsilon$ is small, $B_{\varepsilon}([b c])$ contains only 2 segments which is proven to be tangent to $M_{r}$ at points $b$ and $c$, respectively. The statement is proven.

Lemma 4.6.13. Let $c \in M_{r} \cap \Sigma$. Then $\Sigma$ has the tangent line at $c$, in particular for every $\varepsilon>0$ there is a $\delta>0$ such that for every couple of points $b, d \in \Sigma \cap b_{\delta}(c) \backslash c$, holds $\min (|\angle b c d-\pi|,|\angle b c d|)<\varepsilon$.

Proof of Lemma 4.6.13. Consider a point $c \in M_{r} \cap \Sigma$. By Lemma 4.6.12 if $c$ belongs to some non degenerate arc of $\Sigma \cap M_{r}$ with an energetic point in its interior (i.e. an element of $V_{A}(G)$ ) the statement is true. Note that if there is a chord $[i c] \subset \Sigma$ of $M_{r}$ then Lemma 4.6.11 implies the claim. Thus $B_{\varepsilon}(c) \cap \operatorname{Int}\left(N_{r}\right)=\emptyset$. If $c \in S_{\Sigma}$ then by Lemma 4.7 .2 its neighbourhood cannot be a tripod, so it is a segment and the statement of Lemma is obvious. It remains to consider the case when $B_{\varepsilon}(c) \cap \operatorname{Int}\left(N_{r}\right)=\emptyset$ and $c$ is energetic, which implies $B_{r}\left(q_{c}\right) \cap \Sigma=\emptyset$ so the set $B_{\varepsilon}(c) \cap \Sigma$ is just a segment (because $\Sigma$ consists of a finite number of arcs of $M_{r}$ and segments by Lemma 4.6.10) which must be a subset of the tangent line to $M_{r}$ at $c$, the claim follows.

### 4.6.4 Derivation in the picture

Motivation. Let the point $y \in M$ correspond to two energetic points $x_{1}, x_{2} \in \Sigma$ from the ring $N \backslash N_{r}$. Then by Lemma 4.6.6(i) and (ii) they belong to the closures of different connected components $\Sigma \backslash N_{r}$; let us call them $S_{1}$ and $S_{2}$. The purpose of this section is to describe the structure of the $\Sigma$ minimizer in small neighborhoods of the points $x_{1}$ and $x_{2}$.

Informally speaking, we can try to "move" $y$ along $M$, that is, change the sets $S_{1}$ and $S_{2}$ in small neighborhoods of $x_{1}$ and $x_{2}$ in such a way that the resulting sets together are still covered all the same points $M$ as before; but so that the boundary between the $\operatorname{arcs} M$ they cover would no longer pass at $y$, but at point $M$ at a small distance from it. Since $\Sigma$ is a minimizer, such an operation cannot reduce the length of the entire set, that is, the sum of changes in lengths is non-negative.

Below we formally define the described operation for any of the sets $S_{1}$ and $S_{2}$ when $y$ is shifted by any sufficiently small distance. By directing the shift value $y$ to zero, we obtain "the derivative of the length $\Sigma$ in the vicinity of the point $x_{1}$ (or $x_{2}$ )" when $y$ moves along $M$, we calculate this derivative explicitly and describe the structure $\Sigma$ in small neighborhoods of $x_{1}$ and $x_{2}$ in terms of conditions on "derivatives of the lengths of $\Sigma$ in neighborhoods of $x_{1}$ and $x_{2}$ " when $y$ moves along $M$.

Definition of derivative. Let $S$ be the closure of the connected component $\Sigma \backslash N_{r}$. We denote one of the ends of the arc $Q(S)$ by $y_{1} \in M$. Let $x \in\left(\partial B_{r}\left(y_{1}\right) \backslash N_{r}\right) \cap S$ be the energetic point for which $y_{1}$ is corresponding. By Lemma 4.6.6(iv) the set $S$ is a local Steiner tree for its entering points and energetic points; in addition, $x \in \partial S$, since the energetic point $x$ cannot have a degree 3. By Lemma 4.6.6(i) $x$ cannot have more than two corresponding points. If there are two corresponding points, we denote the second by $y_{2}$. Also, $x$ can have degree 1 or 2 . Let us denote the degree of $x$ by $d \in\{1,2\}$, and the number of points corresponding to $x$ by $k \in\{1,2\}$. Thus, there are 4 possible cases, each of which we will consider in detail below.

Let us fix a $l>0$ such that $\overline{B_{l}(x)} \cap \Sigma$ is the union of $d$ segments of the form $\left[z_{i} x\right], z_{i} \in \partial B_{l}(x)$, 1 leqi $\leq d$. For a sufficiently small $0 \leq \varepsilon<\varepsilon_{0}\left(l, r,\left\{y_{j}\right\}_{j=1}^{k},\left\{z_{i}\right\}_{i=1}^{d}, x\right)$ for $y_{1}^{\varepsilon}$ denote the point obtained by shifting the point $y_{1}$ along $M$ by $\varepsilon$ (that is, such that the arc $M$ with ends at $y_{1}$ and $y_{1}^{\varepsilon}$ has length $\varepsilon$ ) in such a direction that $y_{1}^{\varepsilon} \notin Q(S)$. For a sufficiently small modulo $0>\varepsilon>-\varepsilon_{0}\left(l, r,\left\{y_{j}\right\}_{j=1}^{k}\left\{z_{i}\right\}_{i=1}^{d}, x\right)$ we denote $y_{1}^{\varepsilon}$ is the point obtained by shifting $y_{1}$ along $M$ by $-\varepsilon$ in the opposite direction (that is, in such a way that $y_{1}^{\varepsilon} \in Q(S)$ ). Let us denote $y_{1}^{0}=y_{1}$. In the case of $k=2$, we denote $y_{2}^{\varepsilon}=y_{2}$ for any $\varepsilon$. Let

$$
\Gamma(\varepsilon)=\min _{x^{\prime}} \sum_{i=1}^{d}\left|z_{i} x^{\prime}\right|
$$

where the minimum is taken over all points $x^{\prime}$ such that $\left|y_{j}^{\varepsilon} x^{\prime}\right|=r$ for $1 \leq j \leq k$. Let us denote by $x_{\varepsilon}$ the point at which the value $\Gamma(\varepsilon)$ is reached.

Note that $x_{0}=x$, since $\Sigma$ is a minimizer. The derivative $\Gamma(\varepsilon)$ at the origin $\Gamma^{\prime}(0)$ will be called the derivative of the length $\Sigma$ in the neighborhood of the point $x$ as $y_{1}(x)$ moves along $M$. We will show that this derivative exists by calculating it explicitly in each of the four cases. From the explicit form of the derivative, in particular, it will be clear that it does not depend on the auxiliary parameter $l$ used in the definition.

In further discussions on the angle between a curve and a ray with a vertex on the curve, we always mean the smaller of the angles between the ray and the tangent to the curve drawn at the intersection point of the curve and the ray.

Case 1. Let $d=1, k=1$ (see the left part of Fig. 4.16). By Lemma 4.6.6(vi)(a) the points $z_{1}, x$ and $y_{1}$ lie on the same line. Since $x_{\varepsilon}$ is the closest to $z_{1}$ among the points $\partial B_{r}\left(y_{1}^{\varepsilon}\right), x_{\varepsilon}=\left[z_{1} y_{1}^{\varepsilon}\right] \cap \partial B_{r}\left(y_{1}^{\varepsilon}\right)$. Since $M$ is smooth, the distance from the point $y_{1}^{\varepsilon}$ to the tangent to $M$ at the point $y_{1}$ is $o(\varepsilon)$. Let the angle between $\left(z_{1} y_{1}\right]$ and $M$ be equal to $\alpha$. Note that $\alpha<\pi / 2$ since $x \notin N_{r}$. Moreover, $\angle y_{1}^{\varepsilon} y_{1} x=\pi-\alpha-o(\varepsilon)$ for $\varepsilon>0$, since in this case $y_{1}^{\varepsilon} \notin Q(S)$, and $\angle y_{1}^{\varepsilon} y_{1} x=\alpha+o(\varepsilon)$ for $\varepsilon<0$, since in this case $y_{1}^{\varepsilon} \in Q(S)$. By the cosine theorem for the triangle $z_{1} y_{1} y_{1}^{\varepsilon}$

$$
\left|z_{1} y_{1}^{\varepsilon}\right|=\sqrt{\left|z_{1} y_{1}\right|^{2}+2\left|z_{1} y_{1}\right| \varepsilon \cos \alpha+\varepsilon^{2}}+o(\varepsilon)=\left|z_{1} y_{1}\right|+\varepsilon \cos \alpha+o(\varepsilon) .
$$

Then, since

$$
\begin{gathered}
\Gamma(\varepsilon)-\Gamma(0)=\left(\left|z_{1} y_{1}^{\varepsilon}\right|-r\right)-\left(\left|z_{1} y_{1}\right|-r\right)=\left|z_{1} y_{1}^{\varepsilon}\right|-\left|z_{1} y_{1}\right|=\varepsilon \cos \alpha+o(\varepsilon) \\
\Gamma^{\prime}(0)=\cos \alpha .
\end{gathered}
$$

So, in this case, the derivative of the length $\Sigma$ in the neighborhood of the point $x$ when the point $y_{1}$ moves along $M$ is equal to $\cos \alpha$.


Figure 4.16: The first and second cases

Case 2. Let $d=2, k=1$ (see the right side of Fig. 4.16). By Lemma 4.6.6(vi)(b) the ray $\left[y_{1} x\right)$ contains the bisector of the angle $z_{1} x z_{2}$. Let $\beta=\frac{1}{2} \angle z_{1} x z_{2}$, the angle between $\left[y_{1} x\right)$ and $M$ is equal to $\alpha$. Since at point $x_{\varepsilon}$ the minimum value $\left|z_{1} x_{\varepsilon}\right|+\left|z_{2} x_{\varepsilon}\right|$ among the points $\partial B_{r}\left(y_{1}^{\varepsilon}\right)$, the ray $\left[y_{1}^{\varepsilon} x_{\varepsilon}\right)$ contains the bisector of the angle $z_{1} x_{\varepsilon} z_{2}$, which can be understood by repeating the proof verbatim Lemmas 4.6.6(vi)(b).

If we write these two statements about bisectors algebraically, then $x, x_{\varepsilon} \in N$ are defined as solutions to the following system

$$
f_{1}\left(x^{*}\right)=\left(x^{*}-y, x^{*}-y\right)=r^{2}, \quad f_{2}\left(x^{*}\right)=\frac{\left(x^{*}-z_{1}, x^{*}-y\right)}{\sqrt{\left(x^{*}-z_{1}, x^{*}-z_{1}\right)}}-\frac{\left(x^{*}-z_{2}, x^{*}-y\right)}{\sqrt{\left(x^{*}-z_{2}, x^{*}-z_{2}\right)}}=0
$$

for $y=y_{1}$ and $y=y_{1}^{\varepsilon}$ respectively. Outside the points $z_{1}, z_{2}$, the system is smooth, and the gradient $f_{1}$ is equal to $2\left(x^{*}-y\right)$, that is, parallel to the bisector of the angle $z_{1} x^{*} z_{2}$, which is the level line $f_{2}$. Therefore the implicit function theorem implies $\left|x x_{\varepsilon}\right|=O(\varepsilon)$.

Let us draw a tangent to $B_{r}\left(y_{1}\right)$ through $x$, and parallel to it draw a straight line through $x_{\varepsilon}$. Let the last straight line intersect the ray $\left[y_{1} x\right)$ at point $x_{\text {new }}$. Since $\left|z_{1} x\right|=\left|z_{2} x\right|$, and the ray $\left[x_{\text {new }} x\right)$
contains the bisector of the angle $\angle z_{1} x z_{2}$, the triangles $z_{1} x x_{\text {new }}$ and $z_{2} x x_{\text {new }}$ are equal. Therefore $\angle z_{1} x_{\text {new }} x_{\varepsilon}+\angle z_{2} x_{\text {new }} x_{\varepsilon}=\pi$ and $\left|z_{1} x_{\text {new }}\right|=\left|z_{2} x_{\text {new }}\right|$; let us denote the last value by $l_{\text {new }}$. Let us denote by $\gamma$ the angle $z_{1} x_{\text {new }} x_{\varepsilon}$ and write the cosine theorems for the triangles $z_{1} x_{\varepsilon} x_{\text {new }}$ and $z_{2} x_{\varepsilon} x_{\text {new }}$ and use the fact that $\left|x_{\varepsilon} x_{\text {new }}\right| \leq\left|x x_{\varepsilon}\right|=O(\varepsilon)$ :

$$
\begin{aligned}
& \left|z_{1} x_{\varepsilon}\right|=\sqrt{l_{\text {new }}^{2}+\left|x_{\varepsilon} x_{\text {new }}\right|^{2}-2 l_{\text {new }}\left|x_{\varepsilon} x_{\text {new }}\right| \cos \gamma}=l_{\text {new }}-\left|x_{\varepsilon} x_{\text {new }}\right| \cos \gamma+o(\varepsilon), \\
& \left|z_{2} x_{\varepsilon}\right|=\sqrt{l_{\text {new }}^{2}+\left|x_{\varepsilon} x_{\text {new }}\right|^{2}+2 l_{\text {new }}\left|x_{\varepsilon} x_{\text {new }}\right| \cos \gamma}=l_{\text {new }}+\left|x_{\varepsilon} x_{\text {new }}\right| \cos \gamma+o(\varepsilon) .
\end{aligned}
$$

Thus,

$$
\begin{equation*}
\left|z_{1} x_{\varepsilon}\right|+\left|z_{2} x_{\varepsilon}\right|=2 l_{\text {new }}+o(\varepsilon) . \tag{4.2}
\end{equation*}
$$

Since $\left|x y_{1}\right|=\left|x_{\varepsilon} y_{1}^{\varepsilon}\right|=r$, and $\left|x x_{\varepsilon}\right|,\left|y_{1} y_{1}^{\varepsilon}\right|=O(\varepsilon)$, the angle between the lines $\left(x y_{1}\right)$ and $\left(x_{\varepsilon} y_{1}^{\varepsilon}\right)$ is $O(\varepsilon)$, its cosine is equal to $1-o(\varepsilon)$, therefore projection length the segment $\left[x_{\varepsilon} y_{1}^{\varepsilon}\right]$ to the straight line $\left(x y_{1}\right)$ is equal to $r-o(\varepsilon)$. Let $y^{\prime}$ be the projection of the point $y_{1}^{\varepsilon}$ onto the straight line $\left(x y_{1}\right)$. Due to the smoothness of $M$, the length of the segment $y_{1} y^{\prime}$ is equal to $\varepsilon \cos \alpha+o(\varepsilon)$, therefore

$$
\left|x_{n e w} x\right|=\varepsilon \cos \alpha+o(\varepsilon)
$$

It remains to write the cosine theorem for the triangle $z_{1} x x_{n e w}$ :

$$
\begin{equation*}
l_{\text {new }}=\sqrt{l^{2}+(\varepsilon \cos \alpha+o(\varepsilon))^{2}+2 l(\varepsilon \cos \alpha+o(\varepsilon)) \cos \beta}=l+\varepsilon \cos \alpha \cos \beta+o(\varepsilon) . \tag{4.3}
\end{equation*}
$$

Combining (4.2) and 4.3), we get

$$
\Gamma(\varepsilon)-\Gamma(0)=2 l_{\text {new }}+o(\varepsilon)-2 l=2 \varepsilon \cos \alpha \cos \beta+o(\varepsilon)
$$

that is, the desired derivative is equal to

$$
\Gamma^{\prime}(0)=2 \cos \alpha \cos \beta .
$$

Cases 3 and 4. In these cases $k=2$. Point $x$ lies at the intersection of $\partial B_{r}\left(y_{1}\right) \cap \partial B_{r}\left(y_{2}\right)$. Note that the circles $\partial B_{r}\left(y_{1}\right)$ and $\partial B_{r}\left(y_{2}\right)$ intersect at two points. Indeed, suppose that this is not the case, that is, the circles touch at point $x$. If $d=2$, then $x$ lies on the segment $\left[z_{1} z_{2}\right]$; in this case, the length of $\Sigma$ can be reduced by removing the connected component $S \backslash\{x\}$ containing one of the $z_{i}$ that has no entering points. So $d=1$. In addition, the segment $x z_{1}$ lies on the common tangent to the circles. But this contradicts the fact that $\Sigma$ is a minimizer: consider a point $x^{\prime}$ on the interval $\left[x z_{1}\right]$ such that $\left|x x^{\prime}\right|=\chi$ for small $\chi>0$. Since $x^{\prime}$ lies on the common tangent to the circles, $\left|x^{\prime} y_{1}\right|=r+o(\chi)$, $\left|x^{\prime} y_{2}\right|=r+o(\chi)$. Then, for a sufficiently small $\chi$, the length of $\Sigma$ can be reduced by replacing the segment $\left[x x^{\prime}\right]$ with segments connecting $x^{\prime}$ with the circles $B_{r}\left(y_{1}\right)$ and $B_{r}\left(y_{2}\right)$.

Since $\Sigma$ is connected, lies in $N$ and does not intersect $B_{r}\left(y_{1}\right) \cup B_{r}\left(y_{2}\right)$, then $x$ is that point from $\partial B_{r}\left(y_{1}\right) \cap \partial B_{r}\left(y_{2}\right)$, which lies on the same side relative to the line $\left(y_{1} y_{2}\right)$ as the points $z_{i}, 1 \leq i \leq d$. Let $\varepsilon$ be so small that the circles $\partial B_{r}\left(y_{1}^{\varepsilon}\right)$ and $\partial B_{r}\left(y_{2}^{\varepsilon}\right)$ also intersect at two points. Then $x_{\varepsilon}$ is that point in $\partial B_{r}\left(y_{1}^{\varepsilon}\right) \cap \partial B_{r}\left(y_{2}^{\varepsilon}\right)$ that lies in the same half-plane relative to the straight line $\left(y_{1}^{\varepsilon} y_{2}^{\varepsilon}\right)$, as the points $z_{i}, 1 \leq i \leq d$. Let us find $x_{\varepsilon}$ explicitly (see the left side of Fig. 4.17).

Triangles $x y_{1} y_{2}$ and $x_{\varepsilon} y_{1}^{\varepsilon} y_{2}^{\varepsilon}$ are isosceles with side $r$; let $\angle x y_{1} y_{2}=\angle x y_{2} y_{1}=: \alpha$, and also $\angle x_{\varepsilon} y_{1}^{\varepsilon} y_{2}^{\varepsilon}=$ $\angle x_{\varepsilon} y_{2}^{\varepsilon} y_{1}^{\varepsilon}=: \alpha_{\varepsilon}$.

Let us introduce the following coordinates: midpoint $o$ of the segment $\left[y_{1} y_{2}\right.$ ] be the origin of coordinates; $X$ axis is aligned with the beam $\left[y_{2} y_{1}\right)$; the $Y$ axis is codirected with the ray $[o x)$. Then

$$
o=(0,0), \quad x=(0, r \sin \alpha), \quad y_{1}=(r \cos \alpha, 0), \quad y_{2}=(-r \cos \alpha, 0)
$$




Figure 4.17: Finding coordinates of $x_{\varepsilon}$ in the cases 3 and 4

Let the angle between straight line $\left(y_{1} y_{2}\right)$ and $M$ at point $y_{1}$ be equal to $\delta$. Then

$$
y_{1}^{\varepsilon}=(r \cos \alpha+\varepsilon \cos \delta+o(\varepsilon), \varepsilon \sin \delta+o(\varepsilon)) .
$$

Therefore, by the Pythagorean theorem

$$
\left|y_{2} y_{1}^{\varepsilon}\right|=\sqrt{(2 r \cos \alpha+\varepsilon \cos \delta+o(\varepsilon))^{2}+(\varepsilon \sin \delta+o(\varepsilon))^{2}}=2 r \cos \alpha+\varepsilon \cos \delta+o(\varepsilon) .
$$

Let $o_{\varepsilon}$ be the midpoint of the segment $\left[y_{1}^{\varepsilon} y_{2}^{\varepsilon}\right]$. Then

$$
o_{\varepsilon}=\left(\frac{\varepsilon \cos \delta}{2}+o(\varepsilon), \frac{\varepsilon \sin \delta}{2}+o(\varepsilon)\right) .
$$

By definition of cosine

$$
\alpha_{\varepsilon}=\arccos \left(\frac{\left|y_{2} o_{\varepsilon}\right|}{r}\right)=\arccos \left(\cos \alpha+\frac{\varepsilon \cos \delta}{2 r}+o(\varepsilon)\right)=\alpha-\frac{\cos \delta}{2 r \sin \alpha} \varepsilon+o(\varepsilon) .
$$

Let us denote by $\Delta$ the oriented angle $\angle y_{1} y_{2} y_{1}^{\varepsilon}$ (that is, for negative $\varepsilon$ we have $\Delta<0$ ). By the sine theorem for the triangle $y_{1} y_{2} y_{1}^{\varepsilon}$,

$$
\frac{\varepsilon-o(\varepsilon)}{\sin \Delta}=\frac{\left|y_{1} y_{2}\right|}{\sin (\delta-\Delta+o(\varepsilon))} \geq\left|y_{1} y_{2}\right|, \quad \text { that is } \quad \Delta=o(\varepsilon)
$$

(here we use the fact that $\angle y_{2} y_{1}^{\varepsilon} y_{1}=\delta-\Delta+o(\varepsilon)$, since $M$ is a smooth curve). Hence,

$$
\Delta=\sin \Delta+o(\varepsilon)=\frac{\varepsilon \sin (\delta+O(\varepsilon))}{\left|y_{1} y_{2}\right|}=\frac{\varepsilon \sin \delta}{2 r \cos \alpha}+o(\varepsilon) .
$$

Counting the angles in an isosceles triangle $x y_{2} x_{\varepsilon}$ gives

$$
\angle x y_{2} x_{\varepsilon}=\alpha-\alpha_{\varepsilon}-\Delta=\left(\frac{\cos \delta}{2 r \sin \alpha}-\frac{\sin \delta}{2 r \cos \alpha}\right) \varepsilon+o(\varepsilon)=\frac{\cos (\alpha+\delta)}{r \sin (2 \alpha)} \varepsilon+o(\varepsilon) .
$$

Now it's clear that

$$
\left|x x_{\varepsilon}\right|=2 r \sin \frac{\angle x y_{2} x_{\varepsilon}}{2}=\frac{\cos (\alpha+\delta)}{\sin (2 \alpha)} \varepsilon+o(\varepsilon)
$$

and the angle between the segment $\left[x x_{\varepsilon}\right]$ and the axis $X$ (the straight line passing through the point $x$ and parallel to $\left(y_{1} y_{2}\right)$ ) (see the right side of Fig. 4.17) is equal

$$
\pi-\alpha-\frac{\pi-\angle x y_{2} x_{\varepsilon}}{2}=\frac{\pi}{2}-\alpha+\frac{\cos (\alpha+\delta)}{2 r \sin (2 \alpha)} \varepsilon+o(\varepsilon)=\frac{\pi}{2}-\alpha+o(1)
$$

Case 3. Let $d=1, k=2$. Let us denote by $\beta$ the angle between $\left[z_{1} x\right]$ and the $X$ axis (see the right side of Fig. 4.17). Then

$$
\angle z_{1} x x_{\varepsilon}=\frac{3 \pi}{2}-\alpha-\beta+o(1)
$$

By the cosine theorem for the triangle $z_{1} x x_{\varepsilon}$

$$
\begin{gathered}
\left|z_{1} x_{\varepsilon}\right|=\sqrt{\left|z_{1} x\right|^{2}-2\left|x x_{\varepsilon}\right|\left|z_{1} x\right| \cos \angle z_{1} x x_{\varepsilon}+\left|x x_{\varepsilon}\right|^{2}}=\left|z_{1} x\right|-\left|x x_{\varepsilon}\right| \cos \angle z_{1} x x_{\varepsilon}+o(\varepsilon)= \\
\left|z_{1} x\right|+\frac{\cos (\alpha+\delta) \sin (\alpha+\beta)}{\sin (2 \alpha)} \varepsilon+o(\varepsilon) .
\end{gathered}
$$

Then

$$
\Gamma(\varepsilon)-\Gamma(0)=\left|z_{1} x_{\varepsilon}\right|-\left|z_{1} x\right|=\frac{\cos (\alpha+\delta) \sin (\alpha+\beta)}{\sin (2 \alpha)} \varepsilon+o(\varepsilon)
$$

that is, the derivative is equal

$$
\Gamma^{\prime}(0)=\frac{\cos (\alpha+\delta) \sin (\alpha+\beta)}{\sin (2 \alpha)}
$$

Case 4. Let $d=2, k=2$. As in case 3 , let $\beta$ denote the angle between $\left[z_{1} x\right]$ and the $X$ axis; similarly, we denote by $\gamma$ the angle between $\left[z_{2} x\right]$ and the $X$ axis. Repeating the reasoning from the previous case separately for the segment $\left[z_{1} x\right]$ and for the segment $\left[z_{2} x\right]$, we obtain the value of the derivative

$$
\Gamma^{\prime}(0)=\frac{\cos (\alpha+\delta)}{\sin (2 \alpha)}(\sin (\alpha+\beta)+\sin (\alpha+\gamma))
$$

Transitions between the cases. Let in the second case the angle between the segments $\left[z_{1} x\right]$ and $\left[z_{2} x\right]$ be equal to $\frac{2 \pi}{3}$. Then for $\varepsilon>0$ the angle between the segments $\left[z_{1} x_{\varepsilon}\right]$ and $\left[z_{2} x_{\varepsilon}\right]$ is less than $\frac{2 \pi}{3}$. If we replace $\left[z_{1} x_{\varepsilon}\right] \cup\left[x_{\varepsilon} z_{2}\right]$ with the Steiner tree for the triangle $z_{1} x_{\varepsilon} z_{2}$, the length changes to $o(\varepsilon)$. Thus, we can consider this case as a degenerate first case. In this case, the value of the derivative will not change, since

$$
2 \cos \beta \cos \alpha=\cos \alpha \text { with } \beta=\pi / 3
$$

Under a similar assumption, the fourth case can be considered as a degenerate third case, and the value of the derivative will again remain unchanged:

$$
\begin{gathered}
(\sin (\alpha+\beta)+\sin (\alpha+\gamma)) \frac{\cos (\alpha+\delta)}{\sin (2 \alpha)}=2 \sin \left(\frac{2 \alpha+\beta+\gamma}{2}\right) \cos \left(\frac{\beta-\gamma}{2}\right) \frac{\cos (\alpha+\delta)}{\sin (2 \alpha)}= \\
\sin \left(\alpha+\beta+\frac{\pi}{3}\right) \frac{\cos (\alpha+\delta)}{\sin (2 \alpha)}
\end{gathered}
$$

for $\gamma-\beta=2 \pi / 3$.

## Structural statements.

Proposition 4.6.2. Let $x \in \Sigma \backslash N_{r}$ be an energetic point, $y(x) \in M$ be any of the points corresponding to it. Then the derivative of the length $\Sigma$ in a neighborhood of the point $x$ as $y$ moves along $M$ is non-negative.

Proof. Let us assume the opposite. Let $d, z_{i}$ be taken from the definition of derivative. Let us consider

$$
\Sigma_{\varepsilon}=\Sigma \backslash\left(\bigcup_{i=1}^{d}\left[z_{i} x\right]\right) \cup\left(\bigcup_{i=1}^{d}\left[z_{i} x_{\varepsilon}\right]\right) .
$$

For a sufficiently small $\varepsilon>0$, the length of $\Sigma_{\varepsilon}$ is less than the length of $\Sigma$, the set $\Sigma_{\varepsilon}$ is connected, and $F_{M}\left(\Sigma_{\varepsilon}\right) \leq r$. This contradicts the fact that $\Sigma$ is a minimizer.

Proposition 4.6.3. Let $y \in M$ correspond to two energetic points $x_{1}, x_{2} \in \Sigma \backslash N_{r}$. According to the Proposition 4.6.1, points $x_{1}$ and $x_{2}$ lie on opposite sides of the straight line $\left(y n_{y}\right)$. Then the derivatives of the length $\Sigma$ in the vicinity of the points $x_{1}$ and $x_{2}$ as $y$ moves along $M$ are equal.

Proof. Let us assume the contrary, without loss of generality, that the derivative of the length $\Sigma$ in the neighborhood of the point $x_{1}$ is greater than the derivative of the length $\Sigma$ in the neighborhood of the point $x_{2}$. Let $d_{1}, z_{i}^{1}$ and $d_{2}, z_{i}^{2}$ be taken from from the definition of the derivatives in the neighborhood of $x_{1}$ and $x_{2}$, respectively. Let us consider

$$
\Sigma_{\varepsilon}=\Sigma \backslash\left(\bigcup_{j=1}^{2} \bigcup_{i=1}^{d_{j}}\left[z_{i}^{j} x_{j}\right]\right) \cup\left(\bigcup_{j=1}^{2} \bigcup_{i=1}^{d_{j}}\left[z_{i}^{j}\left(x_{j}\right)_{\varepsilon}\right]\right)
$$

For a sufficiently small $\varepsilon>0$, the length of $\Sigma_{\varepsilon}$ is less than the length of $\Sigma, \Sigma_{\varepsilon}$ is connected, and $F_{M}\left(\Sigma_{\varepsilon}\right) \leq r$. This contradicts the fact that $\Sigma$ is a minimizer.

### 4.7 Horseshoe theorem

### 4.7.1 Sketch of the proof

The proof of Theorem 4.3 .2 consists of two main steps. The first one is to show that a minimizer is the union of chords of $M_{r}$, arcs of $M_{r}$ and closures of connected components of $N \backslash M_{r}$, which are local Steiner trees with at most 4 terminals. Moreover the graph, whose vertices are these Steiner trees and edges connect the trees which are connected by arc of chord, is a finite path.

On the second step we enclose this path into a cycle and consider the inner region $T$ (see Fig. 4.18). Then we compare the turn of $S$ and $\overline{B_{r}(S)} \cap M$ where $S$ is an arc of $\Sigma \cap M_{r}$ or a connected component $\partial T \backslash M_{r}$. It turns out that for every $S$ one has

$$
\begin{equation*}
\operatorname{turn}(S \cap \partial T) \geq \operatorname{turn}\left(\overline{B_{r}(S)} \cap M\right) \tag{4.4}
\end{equation*}
$$

Also one can take into account the boundary terms to get the following key inequality

$$
2 \pi=\operatorname{turn}(\partial T)=\sum_{S} \operatorname{turn}(S \cap \partial T)+\text { boundary terms } \geq \sum_{S} \operatorname{turn}\left(\overline{B_{r}(S)} \cap M\right) \geq \operatorname{turn}(M)=2 \pi
$$

Thus all the inequalities in (4.4) are equalities. The proof of (4.4) immediately imply that this is possible only for a horseshoe.


Figure 4.18: Figure to the construction of $T$.

### 4.7.2 Lemmas for the first step

In the sequel the union of the closures of all connected components of $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$ is denoted by $\Sigma_{r}$. Recall that by Lemma 4.6.1 and Theorem 4.2.1 the set $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$ is a finite union of line segments. Note that the number of line segments in $\Sigma$ sometimes it can be infinite, see Theorem 4.2.3.

Lemma 4.7.1. Let $M$ be a convex closed curve with minimal radius of curvature $R$ and $\Sigma$ be a minimizer with the energy $r<R$. Then the length of each line segment in $\Sigma_{r}$ does not exceed $a_{M}(r)$ for some $a_{M}(r) \leq 2 r$. For the circumference $\partial B_{R}(o)$ one can take $a_{\partial B_{R}(o)}(r)=2 r \sqrt{1-\frac{r^{2}}{4 R^{2}}}$.

Proof of Lemma 4.7.1. Proof of (i): No change in the set $\operatorname{Int}\left(\Sigma \cap N_{r}\right)$ influences the value of $F_{M}(\Sigma)$, so if we take the closure $S$ of any connected component of $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$ and substitute it by a Steiner tree connecting $S \cap M_{r}$ (which must be nonempty if $\Sigma \cap \operatorname{Int}\left(N_{r}\right) \neq \emptyset$ because of connectedness of $\Sigma$ and the requirement $F_{M}(\Sigma) \leq r$ which gives $\Sigma \backslash \operatorname{Int}\left(N_{r}\right) \neq \emptyset$ ), then the length of the resulting set should remain the same by optimality of $\Sigma$, and thus $S$ is itself a Steiner tree connecting $S \cap M_{r}$ as claimed.

Proof of (it): Recall that $\Sigma=E_{\Sigma} \sqcup X_{\Sigma} \sqcup S_{\Sigma}$, where $X_{\Sigma}$ is a discrete set of points, $S_{\Sigma}$ consists of Steiner trees (hence of line segments) and $E_{\Sigma} \subset M_{r}$ by Lemma 4.6.6(iii).

Proof of (iit): Remove an arbitrary open line segment $\Delta$ from the set $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$. The value of $F_{M}$ does not change, i.e. $F_{M}(\Sigma \backslash \Delta)=F_{M}(\Sigma)$, and by the absence of cycles $\Sigma \backslash \Delta$ splits into two connected components $\Sigma_{1}$ and $\Sigma_{2}$, so that $\Sigma \backslash \Delta=\Sigma_{1} \sqcup \Sigma_{2}$ ( $\Sigma$ is closed, so $\Sigma_{1}, \Sigma_{2}$ are closed too). Obviously $M \subset \overline{B_{r}\left(\Sigma_{1}\right)} \cup \overline{B_{r}\left(\Sigma_{2}\right)}$. Then by connectedness of $M$ there is such a point $a \in M$ that $a \in \overline{B_{r}\left(\Sigma_{1}\right)} \cap \overline{B_{r}\left(\Sigma_{2}\right)}$, but then there are points $b \in \overline{\Sigma_{1}}$ and $c \in \overline{\Sigma_{2}}$ such that $|a b| \leq r,|a b| \leq r$. Hence
the distance between $\Sigma_{1}$ and $\Sigma_{2}$ does not exceed $|b c| \leq 2 r$ but the length of the deleted segment $\Delta$ does not exceed the distance between the $\Sigma_{1}$ and $\Sigma_{2}$ in view of optimality of $\Sigma$ (otherwise one could connect $\Sigma_{1}$ with $\Sigma_{2}$ with a shorter segment). we let then $a_{M}(r)$ be the supremum of $|b c|$ over all the possible choices of $\Delta$, so that we have proven $a_{M}(r) \leq 2 r$.

In the case $M=\partial B_{R}(o)$ the length of the segment [bc] reaches its maximal value when [bc] is a chord and $|a b|=|a c|=r$. Then we can calculate the maximal value of length of $[b c]$ in this case:

$$
\sin \frac{\angle a o c}{2}=\frac{|a c|}{2|o c|}=\frac{r}{2 R},
$$

so that

$$
|b c|=2|o c| \sin \angle a o c=4|o c| \sin \frac{\angle a o c}{2} \cos \frac{\angle a o c}{2}=2 r \sqrt{1-\frac{r^{2}}{4 R^{2}}} .
$$

Lemma 4.7.2. Let $\Sigma$ be an r-minimizer for a closed convex curve $M$ with minimal radius of curvature $R>2 a_{M}(r)+r$, where $a_{M}$ is defined in Lemma 4.7.1). Then $\Sigma$ has no Steiner point in $\operatorname{Int}\left(N_{r}\right) \cup$ $\left(S_{\Sigma} \cap N_{r}\right)$ and moreover $\Sigma \cap N_{r}$ consists of chords of $M_{r}$ with disjoint interiors.
Proof. Assume the contrary i.e. that $\Sigma$ has a Steiner point $x \in \operatorname{Int}\left(N_{r}\right) \cup\left(S_{\Sigma} \cap N_{r}\right)$. By the condition on the raduis of curvature there is a point $o \in N$ such that $X \in B_{R}(o)$ and $B_{R}(o) \subset \operatorname{Int}(N)$ (hence $B_{R-r}(o) \subset \operatorname{Int}\left(N_{r}\right)$, and in particular, $\left.o \in \operatorname{Int}\left(N_{r}\right)\right)$. Recall that as defined in Lemma 4.7.1 $\Sigma_{r}$ is the union of the closures of all connected components of $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$. Now denote by $x_{0}$ one of the Steiner points of $\Sigma_{r} \cup\left(S_{\Sigma} \cap M_{r}\right)$ nearest to $o$, and let $t:=\left|o x_{0}\right|$. we claim that $x_{0} \in \operatorname{Int}\left(N_{r}\right)$. In fact, otherwise $x_{0} \in M_{r}$ and hence

$$
t=\operatorname{dist}\left(o, M_{r}\right)=\operatorname{dist}(o, M)-r \geq R-r>3.98 r,
$$

but $x_{0}$ is a Steiner point, hence, in view of the smoothness and convexity of $M_{r}$ there are two line segments $\left[x_{0} i\right] \subset \Sigma, i=1,2$ at angle $2 \pi / 3$ with respect to each other, intersecting $\operatorname{Int}\left(N_{r}\right)$. Suppose without loss of generality that $\angle o x_{0} z_{1} \leq \pi / 3$. Then $z_{1} \in B_{t}(o) \subset \operatorname{Int}\left(N_{r}\right)$, since otherwise there is an $y \in\left[x_{0} z_{1}\right] \cap \partial B_{t}(o) \subset \Sigma \cap \partial B_{t}(o)$ such that the line segment $\left[x_{0} y\right] \subset \Sigma$ is a chord of $\partial B_{t}(o)$, which provides the estimate

$$
\left|x_{0} y\right|=2 t \cos \angle o x_{0} z_{1} \geq t>3.98 r
$$

contrary to Lemma 2.7 (iii), this contradiction proving the claim.
Let $\Sigma^{\prime}$ stand for the closure of the connected component of $\Sigma \cap \operatorname{Int}\left(N_{r}\right)$ containing $x_{0}$. By the structure of a Steiner tree since $x_{0}$ belongs to $\operatorname{Int}\left(N_{r}\right)$ then there are three maximal line segments of $\Sigma^{\prime}$ starting from $x_{0}$. Consider such a pair of them $\left[x_{0} x_{-1}\right],\left[x_{0} x_{1}\right]$ that the point $o$ belongs to the angle $\angle x_{-1} x_{0} x_{1}$ (not excluding the case it belongs to one of the sides of this angle). Recall that $\angle x_{-1} x_{0} x_{1}=2 \pi / 3$. Also note that points $x_{-1}, x_{1}$ lie outside of $B_{t}(o)$. Hence either $\left[x_{0} x_{1}\right]$ or $\left[x_{0} x_{-1}\right]$ intersects $B_{t}(o)$. We assume without loss of generality that it is $\left[x_{0} x_{1}\right]$. Denote the intersection of the segment $\left[x_{0} x_{1}\right]$ and the circumference $\partial B_{t}(o)$ by $c$.
we claim that $t \leq a_{M}(r)$. Supposing the contrary, since $\left|x_{0} c\right| \leq a_{M}(r)$ and $\left|o x_{0}\right|=|o c|=t>$ $a_{M}(r) \geq\left|x_{0} c\right|$, we have $\angle o x_{0} c>\pi / 3$, hence the segment $\left[x_{0} x_{-1}\right]$ also intersects $B_{t}(o)$. Denote the intersection of the segment $\left[x_{0} x_{-1}\right]$ with $\partial B_{t}(o)$ by $d$ and note that also $\angle o x_{0} d>\pi / 3$, and hence $\angle c x_{0} d>2 \pi / 3$ which contradicts the local optimality of $\Sigma$, showing the claim.

Note that $x_{1}, x_{-1}$ belong to $\operatorname{Int}\left(N_{r}\right)$ because $R-r>2 a_{M}(r) \geq t+a_{M}(r)$, and hence $x_{1}, x_{-1}$ are Steiner points. Also by Lemma 4.7.1 the lengths $\left[x_{0} x_{-1}\right]$ and $\left[x_{0} x_{1}\right]$ do not exceed $a_{M}(r)$. Consider a regular hexagon $P$ with sidelength $a_{M}(r)$ such that $x_{0}$ is a vertex of $P$ and the segments $\left[x_{0} x_{1}\right]$, $\left[x_{0} x_{-1}\right.$ ] belong to two sides of $P$. The following assertions hold.

- $\operatorname{diam} P=2 a_{M}(r)$.
- The line segment $\left[o x_{0}\right]$ splits the angle $\angle x_{-1} x_{0} x_{1}=2 \pi / 3$ in two angles, at least one of them is acute. Denote the latter angle by $\angle o x_{0} b$, where $b$ is the corresponding vertex of $P$ (so that $\left.\left|x_{0} b\right|=a_{M}(r)\right)$. Then the angle $\angle o b x_{0}$ is also acute because $\left|o x_{0}\right|=t \leq a_{M}(r)=\left|x_{0} b\right|$. Therefore the perpendicular from $o$ to the line $\left(x_{0} b\right)$ intersects the latter inside $\left[x_{0} b\right]$, so that $o$ is inside the square built on $\left[x_{0} b\right]$. But this square is a subset of $P$ hence $o \in P$.
- The above assertions imply that $P \subset \overline{B_{2 a_{M}(r)}(o)}$, and hence $P \subset \operatorname{Int}\left(N_{r}\right)$.

Now let us pick such vertices $x_{-2}$ and $x_{2}$ that $\left[x_{1} x_{2}\right],\left[x_{-1} x_{-2}\right] \subset \Sigma_{r}$ and $o$ belongs to both angles $\angle x_{0} x_{1} x_{2}$ and $\angle x_{0} x_{-1} x_{-2}$. Clearly $x_{2}, x_{-2} \in P \subset \operatorname{Int}\left(N_{r}\right)$ so they again are Steiner points. Let us define the points $x_{3}, x_{-3}$ in the same way: $\left[x_{2} x_{3}\right],\left[x_{-2} x_{-3}\right] \in \Sigma_{r}$ and $o$ belongs to the angles $\angle x_{1} x_{2} x_{3}$ and $\angle x_{-1} x_{-2} x_{-3}$. Points $x_{3}, x_{-3}$ also belong to $P$, hence to $\operatorname{Int}\left(N_{r}\right)$, hence they also are Steiner points. The six constructed line segments belong to $\operatorname{Int}\left(N_{r}\right)$, so there is no endpoint there. continuing inductively this construction, we arrive at two paths in $P \subset \operatorname{Int}\left(N_{r}\right)$ : one path (starting from $x_{0}, x_{1}, x_{2}, x_{3} \ldots$ ) turns left every time and the other one (starting from $x_{0}, x_{-1}, x_{-2}, x_{-3} \ldots$ ) turns right every time. Thus $\Sigma \cap P \subset \Sigma \cap \operatorname{Int}\left(N_{r}\right)$ contains a cycle or an endpoint of $\Sigma \operatorname{in} \operatorname{Int}\left(N_{r}\right)$, but both cases are impossible for a Steiner tree by the absence of cycles and Lemma 4.7.1.

Let us recall several definitions we used in the second part of Section 4.6.3. Consider the set of the closures of connected components of $\Sigma \backslash N_{r}$ and denote it by $V_{C}(G)$ (further it will be associated with a subset of the vertex set of a graph). Note that $\Sigma$ is connected (and does not reduce to a single point), so every $S \in V_{C}(G)$ has positive length. In our setting $M$ is compact, thus every $\Sigma$ has finite length, hence the set $V_{C}(G)$ is at most countable.

Consider an arbitrary $S \in V_{C}(G)$. Note that by connectedness of $S$ the set $\overline{B_{r}(S)} \cap M$ is always a closed arc. We denote it by $Q_{S}$.

Consider the set of all maximal arcs of $M_{r}$ in the set $\Sigma$, which are not contained in the closure of a connected component of $\Sigma \backslash N_{r}$. Let us denote by $V_{A}(G)$ the subset of such arcs having an energetic point in their interior. Note that if $M$ is not strictly convex, then an $\operatorname{arc}[\breve{b c}]$ of $M_{r}$ can be a chord of $M_{r}$. In this situation if $] \breve{b c}\left[\right.$ has no energetic point then we will consider it as a chord of $M_{r}$ : note that if $\Sigma \backslash] \breve{b c}\left[\right.$ does not cover $q_{x} \in M$ for some $\left.x \in\right] b c[$, then $x$ is energetic; thus if $\left.] \breve{b c}\right]$ has no energetic point then $[b c]=[\breve{b c}]$ has all the properties of a standard chord of $M_{r}$.

Obviously, an $\operatorname{arc}[\breve{b c}] \in V_{A}(G)$ of $M_{r}$ covers an $\operatorname{arc} q_{[\breve{b c}]}:=\left[q_{b} q_{c}\right]$ of $M$, where $q_{b}, q_{c} \in M$ are the unique points such that $\operatorname{dist}\left(b, q_{b}\right)=\operatorname{dist}\left(c, q_{c}\right)=r$.

Definition 4.7.1. Denote by $n(S)$ and $m(S)$ the numbers of energetic and entering points in $S$, respectively.

Theorem 4.6.1 says $n(S) \leq 2, m(S) \leq 3$ and $S$ is a locally minimal tree for its energetic and entering points. We need to strengthen this statement using the bound curvature.

Lemma 4.7.3. Let $M$ be a closed convex curve with minimal radius of curvature $R>2 a_{M}(r)+r$, $\Sigma$. Let $S$ be the closure of a connected component of $\Sigma \backslash N_{r}$. Then $m(S) \leq 2$.

By the previous Lemma, $S$ is a locally minimal tree for at most $n(S)+m(S) \leq 4$ points. All the possible combinatorial types of such networks are listed in Figures 4.33 and 4.34 .

Proof of Lemma 4.7.3. Let $S \in V_{C}(G)$ be the closure of a connected component of $\Sigma \backslash N_{r}$.
Assume the contrary i.e. the existence of at least three different entering points in $S$. Let us denote them $i_{1}, i_{2}$ and $i_{3}$ such that $q_{i_{2}} \in\left[q_{i_{1}} q_{i_{3}}\right] \subset Q_{S}$. Note that $i_{2}$ cannot be energetic, because $q_{i_{2}}$ is not an end of $Q_{S}$. So $i_{2}$ has such a neighbourhood $U$ that $U \cap \Sigma$ is a segment or a regular tripod; by Lemma 4.7 .2 it is a segment.

We claim that $\Sigma$ contains a chord $\left[i_{2} j\right]$ of $M_{r}$. It is true if $\Sigma$ is not tangent to $M_{r}$ at $i_{2}$. Now, let $\Sigma$ be tangent to $M_{r}$ at $i_{2}$, so $i_{2}$ belongs to two closures of different connected components of $\Sigma \backslash N_{r}$; one of them is $S$; denote the second one by $S^{\prime}$. Let $P_{1}$ be the region bounded by the arc $\left[\breve{\left.i_{1} i_{2}\right]}\right.$ of $M_{r}$ (choosing in such a way that $P_{1}$ does not contain $N_{r}$ ) and the unique path between $i_{1}$ and $i_{2}$ in $S$. Define $P_{3}$ analogously (with $i_{3}$ in place of $i_{1}$ ). Obviously, $S^{\prime} \subset P_{1}$ or $S^{\prime} \subset P_{3}$. Hence $q\left(S^{\prime}\right) \subset q(S)$ and replacing $S^{\prime}$ in $\Sigma$ by a Steiner tree for $S^{\prime} \cap M_{r}$ we get a connected competitor to $\Sigma$ still covering $M$. Also, any Steiner tree for $S^{\prime} \cap M_{r}$ belongs to $N_{r}$ by the convexity of $M_{r}$, so this replacement decreases the length, which is impossible. Hence, we get the claim, i.e. there is a chord $\left[i_{2} j\right] \subset \Sigma$ of $M_{r}$.

Then $\left|i_{2} j\right| \leq\left|i_{1} j\right|$ (otherwise we can replace $\left[i_{2} j\right]$ by $\left[i_{1} j\right]$ in $\Sigma$ producing the competitor of strictly lower length), and analogously $\left|i_{2} j\right| \leq\left|i_{3} j\right|$. Note that $j \notin S$ because $\Sigma$ has no loops. One can see that points $i_{1}, i_{2}, i_{3}, j$ belong to $M_{r}$ in the natural (clockwise) order otherwise the arc $Q_{S_{j}}$ is a subset of $Q_{S}$, where $S_{j}$ is the closure of the connected component of $\Sigma \backslash N_{r}$ containing $j$, which is impossible.

Hence $\left|j i_{2}\right|$ is at least the diameter $d$ of the maximal ball inscribed in $N_{r}$ and touching $M_{r}$ at point $i_{2}$, i.e. the double inradius of $M_{r}$. Since $d \geq 2(R-r)$, we have $\left|j i_{2}\right| \geq 2(R-r)>2 r$ contradicting Lemma 4.7.1, showing the claim $m(S) \leq 2$.

Definition 4.7.2. Under conditions of Theorem 4.3.2 consider the following abstract graph $G=$ $(V(G), E(G))$ (recall that the set of vertices $V(G)=V_{C}(G) \sqcup V_{A}(G)$; by Lemma 4.6.10 it is finite), where the set of edges $E(G)$ is defined as follows:

- in the case $S_{1}, S_{2} \in V_{C}(G)$ there is an edge between them if they are connected in $\Sigma$ by a chord of $M_{r}$ or if $S_{1} \cap S_{2} \neq \emptyset$;
- in the case $S_{1} \in V_{C}(G),[\breve{b c}] \in V_{A}(G)$ there is an edge between $S_{1}$ and $[\breve{b c}]$, if $S_{1} \cap[\breve{b c}] \neq \emptyset$;
- and finally in the case $\left[\breve{b_{1} c_{1}}\right],\left[\breve{b_{2} c_{2}}\right] \in V_{A}(G)$ there is no edge between them.

Corollary 4.7.1. Under conditions of Theorem 4.3.2 graph $G$ has no cycles; it has exactly two vertices of degree 1 and all the other vertices have degree 2 . In other words $G$ is a path with at least one edge.

Proof. First, by Lemma 4.6.10 the graph is finite. By Lemma 4.6 .11 every chord of $M_{r}$ in $\Sigma$ connects exactly two vertices in $V(G)$. Thus, the inequality $m(S) \leq 2$ (Lemma 4.7.3) implies $\operatorname{deg}(v) \leq 2$ for $v \in V_{C}(G)$; for $v \in V_{A}(G)$ the inequality $\operatorname{deg}(v) \leq 2$ holds by Lemma 4.6.12.

Note that if $\left(S_{1}, S_{2}\right) \in E(G)$ then there is a path between $S_{1}$ and $S_{2}$ in $\Sigma$ not intersecting other sets $S \in V(G), S \notin\left\{S_{1}, S_{2}\right\}$. It means that if $G$ has a cycle $c$ then so has $\Sigma$, contradicting to the absence of cycles. Moreover, the path between two points in $\Sigma$ belonging to two different vertices of $V(G)$ naturally induces a path in $G$ (in fact, if a path in $\Sigma$ connects two different vertices $S_{1}, S_{2} \in V(G)$ without touching other vertices, then $\left(S_{1}, S_{2}\right) \in E(G)$; therefore for a generic path in $\Sigma$ connecting two different vertices of $G$ it is enough to split it in a finite number of paths connecting different vertices in $G$ and not passing throw other vertices). Therefore, connectedness of $\Sigma$ gives us that $G$ is connected. we conclude that $G$ is a path.


Figure 4.19: Picture to Lemma 4.7.4

Now we have to show that $\# V(G)>1$. Suppose the contrary, i.e. $V(G)=\{v\}$. If $v \in V_{C}(G)$, then $m(v)=0$, so $v$ is a segment that is impossible. Otherwise $v$ is an arc, but $q_{v}=M$, so $v=M_{r}$ contains a loop. We got again a contradiction with the absence of cycles.

Thus under conditions of Theorem 4.3 .2 there are two connected components of $\Sigma \backslash N_{r}$ with one entering point; these components correspond to the leaves of our graph. we call them ending components and denote by $S_{l}$ and $S_{r}$ (calling them left and right respectively); the other components will be called middle components.

By Lemma 4.6.9 every point of $M$ is covered by at most two sets from $V(G)$. By Corollary 4.7.1 graph $G$ is a path, so if $S_{1}, S_{2}$ are connected by an edge in $G$, then $Q_{S_{1}} \cap Q_{S_{2}} \neq \emptyset$. Moreover, the same reasoning gives $Q_{S_{l}} \cap Q_{S_{r}} \neq \emptyset$, because otherwise there would be some part of $M$ not covered by $\Sigma$.

Lemma 4.7.4. The arcs $Q_{S_{l}}$ and $Q_{S_{r}}$ have disjoint interiors.
Denote by $a$ an arbitrary point of the intersection of $Q_{S_{l}}$ and $Q_{S_{r}}$ (see Fig. 4.18); by Lemma 4.7.4 there are at most 2 such points. Consider the set $\hat{\Sigma}:=\Sigma \cup\left[a s_{l}\right] \cup\left[a s_{r}\right]$, where $\left[a s_{l}\right]$ and $\left[a s_{r}\right]$ are segments of length $r$ connecting $a$ with $s_{l}$ and $s_{r}$ respectively. From the absence of cycles and the fact that $B_{r}(a) \cap \Sigma=\emptyset$, the set $\hat{\Sigma}$ bounds the unique region which we further denote by $T$ (see Fig 4.18).

Previous lemmas give us the following corollary.
Corollary 4.7.2. The boundary of $T$ is a closed curve consisting of a finite number of arcs of $M_{r}$ and a finite number of line segments.

Consider the behavior of the tangent line to the boundary of $T$. Corollary 4.7.2 and Lemma 4.6.13 imply that all points where tangent direction is discontinuous (i.e. points where the tangent line to $\partial T$ does not exist) except $a$ belong to connected components of $\Sigma \backslash N_{r}$.

Proof of Lemma 4.7.4. Recall that $m\left(S_{l}\right)=m\left(S_{r}\right)=1$. Denote the ends of $Q_{S_{l}}$ and $Q_{S_{r}}$ in the following way: $Q_{S_{l}}=\left[q_{l}^{S_{l}} q_{r}^{S_{l}}\right], Q_{S_{r}}=\left[q_{l}^{S_{r}} q_{r}^{S_{r}}\right]$. Suppose the contrary, i.e. that $\left.q_{r}^{S_{l}} \in\right] q_{l}^{S_{r}} q_{r}^{S_{r}}\left[, q_{l}^{S_{r}} \in\right.$ $] q_{l}^{S_{l}} q_{r}^{S_{l}}\left[\right.$. Suppose that $n\left(S_{l}\right)=2$ or $n\left(S_{r}\right)=2$ (let $n\left(S_{l}\right)=2$, the case $n\left(S_{r}\right)=2$ is completely analogous). Then by Lemma 4.6.6(v) there is an energetic point of $S_{l}$ corresponding to the point $q_{r}^{S_{l}}$.

But $B_{r}\left(q_{r}^{S_{l}}\right) \cap \Sigma \neq \emptyset$, because $\left.q_{r}^{S_{l}} \in\right] q_{l}^{S_{r}} q_{r}^{S_{r}}\left[=Q_{S_{r}}\right.$. So we have a contradiction with the assumption $n\left(S_{l}\right)=2$, and hence $S_{l}$ coincides with the segment $\left[c_{l} v_{l}\right]$. Clearly, $v_{l}, c_{l}$ and $q_{l}^{S_{l}}$ lie on the same line (otherwise one can replace $\left[v_{l} V^{\prime}\right]$ by the part of the segment $\left[V^{\prime} q_{l}^{S_{l}}\right]$, where $V^{\prime}:=\partial B_{\varepsilon}\left(v_{l}\right) \cap\left[v_{l} c_{l}\right]$ producing a competitor of strictly lower length). Hence $\left[c_{l} v_{l}\right]$ is tangent to $B_{r}\left(Q_{r}^{S_{l}}\right)$ (see Fig. 4.19).

Let $w_{l}$ be such a point of $\left[c_{l} v_{l}\right]$ that $\operatorname{dist}\left(w_{l}, q_{r}^{S_{l}}\right)=r$ and $w_{r}$ be such a point of $\left[c_{r} v_{r}\right]$ that $\operatorname{dist}\left(w_{r}, q_{l}^{S_{r}}\right)=r$. Note that the points $c_{l}, v_{l}, q_{l}^{S_{l}}$ lie on the same line, so $\operatorname{dist}\left(w_{l} q_{l}^{S_{l}}\right) \geq r=$ $\operatorname{dist}\left(w_{l}, q_{r}^{S_{l}}\right)$, so $\angle q_{r}^{S_{l}} q_{l}^{S_{l}} w_{l} \leq \angle q_{l}^{S_{l}} q_{r}^{S_{l}} w_{l}$. The segment $\left[c_{l} v_{l}\right]$ is tangent to $B_{r}\left(Q_{r}^{S_{l}}\right)$, hence $\left(q_{r}^{S_{l}} w_{l}\right) \perp$ $\left(v_{l} c_{l}\right)$. Calculating angles in triangle $\Delta q_{r}^{S_{l}} q_{l}^{S_{l}} w_{l}$ we have $\angle q_{r}^{S_{l}} q_{l}^{S_{l}} w_{l} \leq \pi / 4$. Obviously, $\angle q_{r}^{S_{r}} q_{l}^{S_{l}} w_{l} \leq$ $\angle q_{r}^{S_{l}} q_{l}^{S_{l}} w_{l}$, so $\angle q_{r}^{S_{r}} q_{l}^{S_{l}} w_{l} \leq \pi / 4$. By symmetry we have inequality $\angle q_{l}^{S_{l}} q_{r}^{S_{r}} w_{r} \leq \pi / 4$. Denote by $o$ the intersection point of $\left(v_{l} c_{l}\right)$ and $\left(v_{r} c_{r}\right)$. From the triangle $\Delta q_{r}^{S_{r}} q_{l}^{S_{l}} o$ we have $\angle q_{r}^{S_{r}} o q_{l}^{S_{l}} \geq \pi / 2$.

Note that $2 r>\left|w_{l} w_{r}\right| \geq\left|c_{l} c_{r}\right|$ and $\angle q_{r}^{S_{r}} o q_{l}^{S_{l}}=\angle c_{l} o c_{r} \geq \pi / 2$. It means that $\left|c_{l} O\right|<2 r$ and $\left|c_{r} O\right|<2 r$. Hence the intersection point of the rays $\left[v_{l} c_{l}\right)$ and $\left[v_{r} c_{r}\right)$ belongs to $N_{r}$, that contradicts the optimality of $\Sigma$.

### 4.7.3 Central lemma

Now we are ready to state the central Lemma. Figure 4.18 should simplify the reading of its statement.
Lemma 4.7.5. Under conditions of Theorem 4.3.2 let $\Sigma$ be a minimizer, $S \in V(G)$ be the closure of a connected component of $\Sigma \backslash N_{r}$ or an arc of $M_{r}$. Then the following assertions hold.

- If $S$ is a middle component or an arc of $M_{r}$ then $\operatorname{turn}\left(Q_{S}\right) \leq \operatorname{turn}(S)$. The equality holds if and only if $S$ is an arc of $M_{r}$.
- If $S$ is an ending component then for the left and the right components we have

$$
\begin{aligned}
\operatorname{turn}\left(Q_{S_{l}}\right) & \leq \operatorname{turn}\left(S_{l}\right)+\angle\left(\left[c_{l} s_{l}\right),\left[s_{l} a\right)\right)+\angle\left(\left[s_{l} a\right), a\right), \\
\operatorname{turn}\left(Q_{S_{r}}\right) & \leq \angle\left(a,\left[a s_{r}\right)\right)+\angle\left(\left[a s_{r}\right),\left[s_{r} c_{r}\right)\right)+\operatorname{turn}\left(S_{r}\right),
\end{aligned}
$$

where $A$ stands for the tangent ray to $M$ at the point a directed from the left to the right (see Fig. 4.18, angles $\angle\left(\left[s_{l} a\right), A\right), \angle\left(A,\left[a s_{r}\right)\right)$ are marked red) and $c_{i}$ is the branching point of $S_{i}$ if $S_{i}$ is a tripod and the entering point of $S_{i}$ in other cases, where $i \in\{l, r\}$ (the definition is correct by Lemma 4.7.3). The equality holds if and only if $S$ is a segment of the tangent line to $M_{r}$.

Remark 4.7.1. If in Lemma 4.7.5 we assume that $\Sigma$ has no Steiner points in $N_{r}$ then it is enough to request the inequality $r<R / 2.9$ (see proof of Lemma 4.7.5, case 1a).

Proof of Lemma 4.7.5. Obviously, if $S$ is an arc, then the compared values are equal.
It suffices thus to consider the case when $S$ is the closure of a connected component of $\Sigma \backslash N_{r}$. Denote by $q_{l}$ and $q_{r}$ the ends of $Q_{S}$. Let $o$ be an intersection point of the normals to $M$ at points $q_{l}$ and $q_{r}$. It exists unless turn $\left(Q_{S}\right)=0$ in which case the claim is obvious. Note that turn $\left(Q_{S}\right)=\angle q_{l} o q_{r}$ and denote for brevity thus value by $\gamma$. Also one has $\left|q_{l} o\right| \geq R,\left|q_{r} o\right| \geq R$. Note that Lemmas 4.7.2 and 4.7.3 as well as Corollary 4.7.1 hold true when $R>2 a_{M}(r)+r$ which is guaranteed when $R>5 r$ (or $R>4.98 r$ in the case when $M$ is a circumference of radius $R$ ), i.e. under the conditions of the statement being proven.

By Lemma 4.7.3 $S$ is a locally minimal tree for at most $n(S)+m(S) \leq 4$ points. All the possible combinatorial types of such networks are listed in Figures 4.33 and 4.34 . Note that if $S$ is a middle component then $m(S)=2$, otherwise $m(S)=1$. Let us analyze all the possible types one by one, first when $S$ is a middle component, then for $S$ an ending component.

1. Let $S$ be a middle component. By Lemma 4.6 .6 (iv) $S$ is a locally minimal tree, moreover it has two entering points (if one, then it is an ending component) and one or two energetic points.
(a) The case $n=2, m=2$, the combinatorial type (a) on Fig. 4.34 (see Fig. 4.20). Denote


Figure 4.20: Picture to the case 1a: middle component, $n=2, m=2$.
the Steiner points of $S$ by $v_{l}$ and $v_{r}$. In this case $\operatorname{turn}(S)=\pi / 3+\pi / 3=2 \pi / 3$. Assuming the contrary (it means that $\gamma \geq 2 \pi / 3$ ) and connecting $o$ with $q_{l}$ and $q_{r}$, we get a (non convex) pentagon $q_{l} v_{l} v_{r} q_{r} o$ with two angles equal to $4 \pi / 3$ and one angle at least $2 \pi / 3$, which is impossible.
(b) The case $n=2, m=2$, the combinatorial type (b) on Fig. 4.34 (see Fig. 4.21). Note that


Figure 4.21: A general picture to the case 1 b Figure 4.22: A marginal picture to the case 1 b ; middle component, $n=2, m=2$.
in this case there exists a Steiner point adjacent to both entering points, and also there exists a Steiner point (we call it $b$ ) adjacent to both energetic points. Clearly $\operatorname{turn}(S)=\pi / 3$.


Figure 4.23: Picture to the case 1c middle component, $n=2, m=2$.

Let us prove that $\operatorname{turn}\left(Q_{S}\right)<\pi / 3$. we evaluate the arc of $M$ bounded by continuations of segments starting from $b$. Clearly this arc is maximal when $b$ belongs to $M_{r}$ (it is the marginal case). Hence it is enough to look at the angle in $N \backslash N_{r}$ of size $2 \pi / 3$ with vertex $b$ on $M_{r}$. It is well-known that the arc is maximal when $S$ is tangent to $M_{r}$ and when $M$ is a circumference. In this case the normal to $M_{r}$ at $b$ splits the angle $\angle q_{l} b q_{r}=2 \pi / 3$ in two angles: one of size $\pi / 2$ and another of size $\pi / 6$ (see Fig. 4.22), so that the size of the arc is

$$
\arccos \left(1-\frac{1}{\delta}\right)+\frac{\pi}{6}-\arcsin \left(\frac{1}{2}\left(1-\frac{1}{\delta}\right)\right)
$$

where $\delta:=R / r$, hence it is strictly less than $\pi / 3$ for $\delta \geq 2.9$.
(c) The case $n=2, m=2$, the combinatorial type (c) on Fig. 4.34.

There are two possibilities for $S$ in this case, see Fig. 4.23 and Fig. 4.24.
The case on Fig. 4.24 can be reduced to the previous case 1b. Obviously, $\operatorname{turn}(S)=\pi / 3$. Let us fix the entering points $y_{l}, y_{r}$ and the left energetic point $w_{l}$ and move the right energetic point $w_{r}$ to the right (in the direction of the ray $\left[w_{l} w_{r}\right)$ ). Then at some time the combinatorial type changes to (b) on Fig. 4.34, during this process $\operatorname{turn}(S)=\pi / 3$, and $\operatorname{turn}\left(Q_{S}\right)$ grows, but turn $\left(Q_{S}\right) \leq \pi / 3$. By case 1 b .
The case on Fig. 4.23 denote the energetic points of $S$ by $w_{l}$ and $w_{r}$, and the entering points by $y_{l}, y_{r}$ respectively, and the branching point by $v_{l}$ (without loss of generality it is connected with $w_{l}$ and $\left.y_{l}\right)$. Let $2 \beta:=\angle v_{l} w_{r} y_{r}$, and note that $\angle y_{l} v_{l} w_{r}=2 \pi / 3$. Then $\operatorname{turn}(S)=(\pi-2 \pi / 3)+(\pi-2 \beta)=4 \pi / 3-2 \beta$. Assume the contrary (i.e. in


Figure 4.24: Picture to the case 1 C] middle component, $n=2, m=2$.
this case $\gamma \geq 4 \pi / 3-2 \beta$ ) and call $l$ the point of intersection of $\left(q_{l} w_{l}\right)$ and $\left(q_{r} w_{r}\right)$. By Lemma 4.6.6(v)(b) $\angle l w_{r} v_{l}=\angle y_{r} w_{r} v_{l} / 2=\beta$. Then

$$
\pi-\pi / 3-\beta=\angle q_{l} l q_{r}>\angle q_{l} o q_{r}=\gamma
$$

(the first equality coming from $\Delta v_{l} w_{r} l$ ) which implies

$$
\gamma \geq 4 \pi / 3-2 \beta>2 \pi / 3-\beta=\angle q_{l} l q_{r}>\gamma
$$

a contradiction.
(d) The case $n=2, m=2$, the combinatorial type (d) on Fig. 4.34 (see Fig. 4.25). Denote the energetic points of $S$ by $w_{l}$ and $w_{r}$, and the entering points by $y_{l}, y_{r}$ respectively. Let $2 \alpha:=\angle y_{l} w_{l} w_{r}, 2 \beta:=\angle w_{l} w_{r} y_{r}$. Then $\operatorname{turn}(S)=(\pi-2 \alpha)+(\pi-2 \beta)$. Assume the contrary (it means that $\gamma \geq 2 \pi-2 \alpha-2 \beta$ ) and denote by $l$ the point of intersection of $\left(q_{l} w_{l}\right)$ and $\left(q_{r} w_{r}\right)$. By Lemma 4.6.6(v)(b) $\angle l w_{l} w_{r}=\angle y_{l} w_{l} w_{r} / 2=\alpha, \angle l w_{r} w_{l}=\angle y_{r} w_{r} w_{l} / 2=\beta$. Then

$$
\pi-\alpha-\beta=\angle q_{l} l q_{r}>\angle q_{l} o q_{r}=\gamma
$$

(the first equality coming from $\Delta w_{l} w_{r} l$ ) which implies

$$
\gamma \geq 2 \pi-2 \alpha-2 \beta>\pi-\alpha-\beta=\angle q_{l} l q_{r}>\gamma
$$

a contradiction.
(e) The case $n=1, m=2$, the combinatorial type (b) on Fig. 4.33 (see Fig. 4.26, Fig. 4.27, Fig. 4.28.
Clearly, $\operatorname{turn}(S)=\pi / 3$. To prove the statement, assume the contrary (i.e. $\gamma \geq \pi / 3$ ) and as in the previous case connect $o$ with $q_{l}$ and $q_{r}$. Denote the energetic point of $S$ by $w$. Let us consider three subcases:

- the point $w$ covers both $q_{r}$ and $q_{l}$ (see Fig. 4.26);
- the point $w$ covers $q_{l}$ and $q_{r}$ is covered by an entering point (see Fig. 4.27);


Figure 4.25: Picture to the case 1d: middle component, $n=2, m=2$.


Figure 4.26: Picture to the case 1e: middle component, $m=2, n=1$.


Figure 4.27: Picture to the case 1e: middle component, $m=2, n=1$.


Figure 4.28: Picture to the case 1e: middle component, $m=2, n=1$.

- $w$ covers $q_{l}$ and $q_{r}$ is covered by $h \in S \backslash\left(M_{r} \cup w\right)$ (see Fig. 4.28).

In the SUBCASE (i) $\left|w q_{r}\right|=\left|w q_{l}\right|=r$. Let us connect $o$ with $w$, and note that the angle $\angle q_{l} O q_{r}=\gamma$ splits into two parts; let us pick the largest one (without loss of generality it is $\angle w o q_{r}$ ). Consider the triangle $\Delta o q_{r} w$ with side $\left|o q_{r}\right| \geq R$ and acute angle ( $\alpha$ on the Fig. 4.26) at least $\pi / 6$ against the side $\left|w q_{r}\right|=r$. Recalling that $R>2 r$ and denoting by $\beta:=\angle o w q_{r}$, by the law of sines for triangle $\Delta o q_{r} w$ we get

$$
\sin \beta=\frac{\left|o q_{r}\right|}{r} \sin \alpha \geq \frac{R}{2 r}>1
$$

a contradiction.
In the SUbCASE (ii) $q_{r}$ is covered by the entering point $i$. Then $(c i)$ is perpendicular to $\left(i q_{r}\right)$, where $c$ is the branching point of $S$, so points $q_{r}, o, i$ lie on the same line. Consider the sum of the angles in the non convex quadrilateral $q_{l} c i o$ : it is $\angle q_{l}+\angle c+\angle i+\angle o \geq$ $\angle q_{l}+4 \pi / 3+\pi / 2+\pi / 3>2 \pi$, a contradiction.
In the subcase (iii) $q_{r}$ is covered by $\left.h \in\right] c i[$, where $c$ is the branching point of $S, i$ is an entering point of $S$. Note that (ci) is perpendicular to $\left(h q_{r}\right)$; points $q_{l}, w, c$ lie on the same line. Consider the sum of the angles in the non convex pentagon $q_{l} c h q_{r} o$ : it is $\angle q_{l}+\angle c+\angle h+\angle q_{r}+\angle o \geq \angle q_{l}+4 \pi / 3+3 \pi / 2+\angle q_{r}+\pi / 3>3 \pi$, a contradiction.
(f) The last case $n=1, m=2$, the combinatorial type (c) on Fig. 4.33 (see Fig. 4.29). Then $S$ consists of two segments, i.e. $S=[b w] \cup[w d]$, where $b, d \in M_{r}$ are entering points, $w$ is energetic and $\angle b w d \geq 2 \pi / 3$. In this case $\operatorname{turn}(S)=\pi-\angle b w d$.
First, connect $o$ with $q_{l}$ and $q_{r}$ then denote $k_{l}=\left[o q_{l}\right] \cap M_{r}$ and $k_{r}=\left[o q_{r}\right] \cap M_{r}$. Now consider the convex quadrilateral $P=k_{l} o k_{r} w$. The sum of the angles $\angle k_{l}+\angle k_{r}+\angle w$ of $P$ is at least $\pi / 2+\angle b w d+\pi / 2$, so that the remaining angle (which is equal to $\gamma$ ) is at most $\pi-\angle b w d=\operatorname{turn}(S)$ as claimed.
If one has the equality then both $[b w]$ and $[w d]$ are tangent to $M_{r}$, but $w$ is not energetic point in this case, because $q_{l}$ is covered by $b=k_{l}, q_{r}$ is covered by $d=k_{r}$, so we got a contradiction.


Figure 4.29: Picture to the case 1f: middle component, $m=2, n=1$.

2. Let $S$ be an ending component (without loss of generality let it be the left one, so $q_{r}=A$ ). Recall that c denotes the branching point if $S$ is a tripod and the entering point if $S$ is a seqment. Then there are two options:
(a) The case $n=1$, $m=1$, the combinatorial type (a) on Fig. 4.33 (see Fig. 4.30). In this case $S=\left[c s_{l}\right]$, where $c \in M_{r},\left|s_{l} q_{r}\right|=r$, and $\operatorname{turn}(S)=0$. Denote by $k$ such a point that $k \in\left[o q_{l}\right)$ and $\angle o q_{r} k=\pi / 2$. Define the points $l:=\left[s_{l} c\right) \cap\left(o q_{l}\right)$ and $p:=\left[c s_{l}\right) \cap\left(q_{r} k\right)$, and introduce the angles $\alpha:=\angle p s_{l} q_{r}$ and $\beta:=\angle s_{l} q_{r} k$.
The following two situations have to be considered. Note that $\left|s_{l} q_{l}\right|=r$, otherwise one can replace $\left[c s_{l}\right] \cap B_{\varepsilon}\left(s_{l}\right)$ in $\Sigma$ by the part $[d f]$ of the segment $\left[d q_{r}\right]$ where $d=\left[c s_{l}\right] \cap \partial B_{\varepsilon}\left(s_{l}\right)$, $f$ is the point satisfying $\operatorname{dist}\left(f, q_{r}\right)=r$, producing the competitor of strictly lower length.

- case $\angle c s_{l} q_{r} \leq \pi$ (see the top picture on Fig. 4.30).

Then $\angle\left(\left[a s_{l}\right), A\right)=\beta$ and $\angle\left(\left[c s_{l}\right),\left[s_{l} a\right)\right)=\alpha$, so that

$$
\operatorname{turn}(S)+\angle\left(\left[c s_{l}\right),\left[s_{l} a\right)\right)+\angle\left(\left[s_{l} a\right), A\right)=\alpha+\beta
$$

Note that $\angle s_{l} p k=\alpha+\beta$ and $\angle o k q_{r}=\pi / 2-\gamma$. If $\alpha+\beta \leq \gamma$ (contrary to the claim being proven), then $\angle o k p+\angle k p s_{l}<\pi / 2$ so $\angle k l p>\pi / 2$, which is impossible because then $\left|c q_{l}\right|<\left|s_{l} q_{l}\right|$ which contradicts $\left|s_{l} q_{l}\right|=r,\left|c q_{l}\right| \geq r$.

- case $\angle c s_{l} q_{r}>\pi$ (see the bottom picture on Fig. 4.30). In this case $\angle\left(\left[s_{l} a\right), A\right)=\beta$


Figure 4.30: Picture to the case 2a: ending component, $n=1, m=1$.


Figure 4.31: Picture to the case 2b; ending component, $n=2, m=1$.
and $\angle\left(\left[c s_{l}\right),\left[s_{l} a\right)\right)=-\alpha$, so that

$$
\operatorname{turn}(S)+\angle\left(\left[c s_{l}\right),\left[s_{l} a\right)\right)+\angle\left(\left[s_{l} a\right), A\right)=\beta-\alpha
$$

and we know that $\angle k p c=\beta-\alpha$. If $\beta-\alpha \leq \gamma$ (the contrary to the claim being proven), then $\angle o k p+\angle k p c<\pi / 2$, which is impossible because then $\left|c q_{l}\right|<\left|s_{l} q_{l}\right|$ which contradicts $\left|s_{l} q_{l}\right|=r,\left|c q_{l}\right| \geq r$.
(b) The case $n=2, m=1$, the combinatorial type (b) on Fig. 4.33 (see Fig. 4.31).

Note that $S$ is a tripod: $S=[b c] \cup[c w] \cup\left[c s_{l}\right] \subset \overline{\left(N \backslash N_{r}\right)}$, where $b \in M_{r}$. Let us prove that $q_{r}=\left[c s_{l}\right) \cap M$ and $q_{l}=[c w) \cap M$. Suppose the contrary i.e. without loss of generality $c, s_{l}$, and $q_{r}$ do not lie on the same line. Let us pick a sufficiently small $\varepsilon>0$ and denote by $j$ the intersection point of $\partial B_{\varepsilon}\left(s_{l}\right)$ with $\left[c s_{l}\right]$. Then one may replace $\left[j s_{l}\right]$ by $[j i]$ in $\Sigma$, where $i$ stands for the intersection point of $\partial B_{r}\left(q_{r}\right)$ with $\left[j q_{r}\right]$. Clearly the resulting set covers $Q_{s_{l}}$, so it has the same energy $F_{M}$; by the triangle inequality it has strictly lower length, so we got a contradiction.
Note that $\left|s_{l} q_{r}\right|=r=\left|w q_{l}\right| ; B_{r}\left(q_{r}\right) \cap \Sigma=B_{r}\left(q_{l}\right) \cap \Sigma=\emptyset$. Let $k \in\left[o q_{l}\right)$ be the point satisfying $\left(q_{r} k\right) \perp\left(o q_{r}\right)$. Then $\alpha:=\operatorname{turn}(S)=\angle\left([b c),\left[c q_{r}\right)\right)=\pi / 3, \angle\left(\left[c s_{l}\right),\left[s_{l} a\right)\right)=0$


Figure 4.32: Picture to the case 2c ending component, $n=2, m=1$.
and $\beta:=\left(\left[c q_{r}\right),\left[q_{r} k\right)\right)=\angle\left(\left[s_{l} a\right), A\right)$. we have to show $\alpha+\beta>\gamma$. Let $p$ be the point of intersection of $\left(k q_{r}\right]$ and $[b c)$. Then $\angle o k p=\pi / 2-\gamma$ and $\angle k p c=\alpha+\beta$. Assume the contrary, i.e. $\alpha+\beta \leq \gamma$. Then $\angle o k p+\angle k p c \leq \pi / 2$ hence $\angle k l p \geq \pi / 2$, where $l$ is the point of intersection of $(b c)$ and $(o k)$, but since $\angle q_{l} c l=2 \pi / 3$, then the sum of the angles of the triangle $\Delta c l q_{l}$ exceeds $\pi$, which is impossible.
(c) The case $n=2, m=1$, the combinatorial type (c) on Fig. 4.33 (see Fig. 4.32). In this case $a=q_{r}, s_{l}=w_{r}$. Denote $\angle\left(\left[c w_{r}\right),\left[w_{r} q_{r}\right)\right)$ by $\alpha, \angle\left(\left[s_{l} a\right), A\right)=\angle\left(\left[w_{r} q_{r}\right), A\right)$ by $\beta$, clearly $\operatorname{turn}(S)=\alpha+\beta, \operatorname{turn}\left(Q_{S}\right)=\gamma$. Let $l$ be the point of intersection of $\left(w_{r} c\right)$ and $\left(q_{l} o\right)$. Suppose the contrary, i.e. $\gamma \geq \alpha+\beta$. Then

$$
\begin{aligned}
& \angle w_{r} l q_{l}=\pi-\angle w_{r} l o=\pi-\left(2 \pi-\angle l w_{r} q_{r}-\angle w_{r} q_{r} o-\angle q_{r} o l\right)= \\
& \pi-(2 \pi-(\pi-\alpha)-(\pi / 2-\beta)-\gamma)=\pi / 2-\beta-\alpha+\gamma \geq \pi / 2
\end{aligned}
$$

which is impossible because then $\left|c q_{l}\right|<\left|s_{l} q_{l}\right|$, which contradicts $\left|s_{l} q_{l}\right|=r,\left|c q_{l}\right| \geq r$.

### 4.7.4 Finishing the proof

Now the proof of Theorem 4.3.2 is just few lines.


Figure 4.33: Locally miminal trees for sets of 2 and 3 points.


Figure 4.34: Locally miminal trees for sets of 4 points.

Proof of Theorem4.3.2. By Lemma 4.7.1 $2 a_{M}(r)+r<5 r$ for general $M$, and $2 a_{M}(r)+r<4.98 r$ when $M$ is the circumference. Note that

$$
2 \pi=\operatorname{turn}(\partial T)=\sum_{S \in V(G)} \operatorname{turn}(S)+\angle\left(\left[c_{l} s_{l}\right),\left[s_{l} a\right)\right)+\angle\left(\left[s_{l} a\right), a\right)+\angle\left(\left[a s_{r}\right),\left[s_{r} c_{r}\right)\right)+\angle\left(a,\left[a s_{r}\right)\right)
$$

by Lemma 4.6.12 and Lemma 4.6.13, and also $\operatorname{turn}(M)=2 \pi$. Hence by Lemma 4.7.5

$$
\begin{aligned}
2 \pi & =\sum_{S \in V(G)} \operatorname{turn}(S)+\angle\left(\left[c_{l} s_{l}\right),\left[s_{l} a\right)\right)+\angle\left(\left[s_{l} a\right), a\right)+\angle\left(\left[a s_{r}\right),\left[s_{r} c_{r}\right)\right)+\angle\left(a,\left[a s_{r}\right)\right) \\
& \geq \sum_{S \in V(G)} \operatorname{turn}\left(Q_{S}\right) \geq \operatorname{turn}(M)=2 \pi
\end{aligned}
$$

Thus all the inequalities in Lemma 4.7 .5 are equalities. Summing up, every global minimizer $\Sigma$ consists of arcs of $M_{r}$ and segments of tangent lines to $M_{r}$, i.e. components of the combinatorial type (a) on Fig. 4.33, tangent to $M_{r}$. Every vertex, corresponding to a component of the combinatorial type (a) on Fig. 4.33 has degree 1 in $G$. Thus $\Sigma$ has the unique arc of $M_{r}$, and because of the absence of loops it cannot coincide with $M_{r}$. By Lemma 4.6.12 every maximal arc $[b c] \in V_{A}(G)$ is connected in the graph $G$ with two vertices, corresponding to connected components of $\Sigma \backslash N_{r}$. Hence any minimizer is a horseshoe.

Proof of Corollary 4.3.1. Let $\hat{\Sigma}$ be a local minimizer in the sense of Definition 4.1.2. Suppose the claim is false, i.e.

$$
\begin{equation*}
\mathcal{H}^{1}(\hat{\Sigma})-\mathcal{H}^{1}(\Sigma)<(R-5 r) / 2 \tag{4.5}
\end{equation*}
$$

and $\hat{\Sigma}$ is not a horseshoe. Suppose first that $\hat{\Sigma}_{r}$ contains no line segment of length exceeding

$$
a_{M}^{\prime}(r):=2 r+\mathcal{H}^{1}(\hat{\Sigma})-\mathcal{H}^{1}(\Sigma)<2 r+(R-5 r) / 2 .
$$

Then Lemma 4.7.2 remains true for this situation with $a_{M}^{\prime}$ instead of $a_{M}$, because $2 a_{M}^{\prime}(r)+r<R$. Lemma 4.7.3 also remains true with $a_{M}^{\prime}(r)$ instead of $a_{M}$ by the same reason. we may repeat now line by line the proof of Theorem 4.3.2 without any change because all the arguments used in this proof as well as in Lemma 4.7.5 are local, except the Lemma 4.7 .2 and Lemma 4.7 .3 which hold true with $a_{M}^{\prime}$ instead of $a_{M}$. This proves that $\hat{\Sigma}$ is a horseshoe in the considered case.

On the other hand it is impossible to $\hat{\Sigma}_{r}$ to have a segment of length at least $a_{M}^{\prime}(r)$, otherwise using the replacement from Lemma 4.7.1 and get a contradiction with (4.5).

### 4.8 Open questions

### 4.8.1 Regularity

The first question, especially if the answer in negative, might be difficult.
Question 4.8.1. Does there exist a nonplanar maximal distance minimizer with infinite number of branching points?

An easier question should be to construct an example of a minimizer with a branching point, whose neighbourhood does not coincide with a regular tripod:

Question 4.8.2. To construct a (nonplanar) maximal distance minimizer $\Sigma$ containing a locally nonplanar branching point $x$, i.e. for every $\varepsilon>0$ the set $B_{\varepsilon}(x) \cap \Sigma$ does not belong to a plane.

Thus the question if there exists a nonplanar maximal distance minimizer with an infinite number of points with three tangent rays also makes sense.

The following question asks if one-sided tangents should have continuity from the corresponding side.

Question 4.8.3. Does Lemma 4.2.1 holds for a d-dimensional maximal distance minimizer?
All the questions in this subsection can be also asked for a local minimizers.

### 4.8.2 Explicit solutions

Recall that the horseshoe conjecture is still open.
Question 4.8.4. Find maximal distance minimizers for a circumference of radius $4.98 r>R>r$.
At the same time, the statement of Theorem 4.3 .2 does not hold for a general $M$ if the assumption on the minimal radius of curvature is omitted as we show below.

Define a stadium to be the boundary of the $R$-neighborhood of a segment. By the definition, a stadium has the minimal radius of curvature $R$. Let us show that if $R<1.75 r$ and a stadium is long enough, then there is the connected set $\Sigma^{\prime}$ that has the length smaller than an arbitrary horseshoe and covers $M$.


Figure 4.35: Horseshoe is not a minimizer for long enough stadium with $R<1.75 r$.

Define $\Sigma_{0}$ to be the locally minimal tree depicted in Fig. 4.35. Let $\Sigma^{\prime}$ consist of copies of $\Sigma_{0}$, glued at points $a$ and $b$ along the stadium. Note that $F_{M}\left(\Sigma^{\prime}\right) \leq r$ by the construction. In the case $R<1.75 r$ the length of $\Sigma_{0}$ is strictly smaller than $2|a b|$. Thus for a long enough stadium $\Sigma^{\prime}$ has length $\alpha L+O(1)$, where $L$ is the length of the stadium and $\alpha<2$ is a constant depending on $\Sigma_{0}$ and $R$. On the other hand, any horseshoe has length $2 L+O(1)$.

This example leads to the following problems.
Question 4.8.5. Find the minimal $\alpha$ such that Theorem 4.3.2 holds with the replacement of 5 r with $\alpha r$.

Question 4.8.6. Describe the set of r-minimizers for a given stadium.
Analogously to the stadium case one can easily show that for some sufficiently small $\frac{\left|a_{1} a_{2}\right|}{\left|a_{2} a_{3}\right|}<1$ and some $r>0$ a minimizer should have another topology than depicted at Fig. 4.7.

Also one may consider the following relaxation of Problem 4.8.6.

Question 4.8.7. Fix a real $a>2 r$. Let $M(l)$ be the union of two sides of length $l$ of a rectangle $a \times l$ and $\Sigma(l)$ be a minimizer for $M(l)$. Find

$$
s(a):=\lim _{l \rightarrow \infty} \frac{\mathcal{H}^{1}(\Sigma(l))}{l} .
$$

If $a>10 r$ one may add up $M(l)$ to a stadium and use Theorem 4.3.2 to get $s(a)=2$.

### 4.8.3 Uniqueness

Recall that if $\Sigma$ be an $r$-minimizer for some $M$, then it is a minimizer for $\overline{B_{r}(\Sigma)}$. This motivates the following question.

Question 4.8.8. Let $\Sigma$ be an r-minimizer for some $M$. Is $\Sigma$ the unique $r$-minimizer for $\overline{B_{r}(\Sigma)}$ ?
A weaker form of this question is if we replace $r$ with some positive $r_{0}<r$ in the hypothesis. Again we are interested whether Proposition 4.5.4 holds in larger dimensions.

Question 4.8.9. Fix $d \geq 3$ and $n \geq 4$. Find the Hausdorff dimension of $d$-dimensional $n$-point ambiguous configurations $M$ (as a subset of $\mathbb{R}^{d n}$ ).

A weaker question is to determine whether the set of $d$-dimensional $n$-point ambiguous configurations has measure zero.

## Chapter 5

## Johnson-type graphs

This chapter is based on papers [15] and [23]. We consider a family of distance graphs in $\mathbb{R}^{d}$ and find its independence numbers in some cases.

Define the graph $J_{ \pm}(d, k, t)$ in the following way: the vertex set consists of all vectors from $\{-1,0,1\}^{d}$ with exactly $k$ nonzero coordinates; edges connect the pairs of vertices with scalar product $t$. We find the independence number of $J_{ \pm}(d, k, t)$ for an odd negative $t$ and $d>d_{0}(k, t)$.

### 5.1 Basics

We start with common definitions. Let $G=(V, E)$ be a graph. A subset $I$ of vertices of $G$ is independent if no edge connects vertices of $I$. The independence number of a graph $G$ is the maximal size of an independent set in $G$; we denote it by $\alpha(G)$.

Generalized Johnson graphs are the graphs $J(d, k, t)$ defined as follows: the vertex set consists of vectors from the hypercube $\{0,1\}^{d}$ with exactly $k$ nonzero coordinates, edges connect vertices with scalar product $t$ (so $J(d, k, t)$ is nonempty if $k<d$ and $2 k-d \leq t<k$ ). Generalized Kneser graphs $K(d, k, t)$ have the same vertex set but the edges connect vertices with scalar product at most $t$.

Now we introduce the main hero of the chapter. Define graphs $J_{ \pm}(d, k, t)$ as follows: the vertex set consists of vectors from $\{-1,0,1\}^{d}$ with exactly $k$ nonzero coordinates, edges connect vertices with scalar product $t$. The graph $J_{ \pm}(d, k, t)$ is nonempty if $k<d$ and $-k \leq t<k$, and also if $k=d$ and $d-t$ is even. If $t=-k$, then the graph $J_{ \pm}(d, k, t)$ is a matching. Note that the edges connect vertices of the Euclidean distance $\sqrt{2(k-t)}$, which means that $J_{ \pm}(d, k, t)$ is a distance graph.

Finally, define $K_{ \pm}(d, k, t)$ as the graph which splits the vertex set with $J_{ \pm}(d, k, t)$ but the edges connect vertices with scalar product at most $t$.

The support of a vertex is the set of its non-zero coordinates. For $k=2$ we use the notation $a^{i} b^{j}$ for a vertex with support $\{a, b\}$ and signs $i, j \in\{+,-\}$ on coordinates $a$ and $b$, respectively; We use similar notation for $k=3$.

An automorphism of a graph is a bijection from a set of vertices onto itself that preserves adjacency. A graph is called vertex-transitive if for any vertices $u$ and $v$ there is a graph automorphism that takes $u$ to $v$.

Finally, let $m(a, b)$ be the number of the most significant unequal digit in the binary notation of the numbers $a$ and $b$ (bits are numbered starting from one).

### 5.1.1 Independence and chromatic numbers of $J(d, k, t)$ and $K(d, k, t)$

Independent sets in these families of graphs are classical combinatorial objects. Indeed, we have a natural bijection between the set of $k$-subsets of $[d]$ and $V[J(d, k, t)]=V[K(d, k, t)]$. The celebrated Erdős-Ko-Rado theorem [37] determines all maximal independent sets in $J(d, k, 0)=K(d, k, 0)$. A natural generalization was done by Erdôs and Sós, who introduce "forbidden intersection problem", which involves finding the independence numbers of graphs $J(d, k, t)$. Then the Frankl-Wilson theorem [52], the Frankl-Füredi theorem [44] and the Ahlswede-Khachatryan Complete Intersection Theorem [1] answered a lot of questions about the size and the structure of maximal independent sets in the graphs $J(d, k, t)$ and $K(d, k, t)$.

On the other hand a lot of questions in combinatorial geometry are related to embeddings of these graphs into $\mathbb{R}^{d}$. Frankl and Wilson [52] used the graphs $J(d, k, t)$ to get an exponential lower bound on the chromatic number of the Euclidean space (Nelson-Hadwiger problem); Kahn and Kalai [68] used them to disprove Borsuk's conjecture.

Let us describe the picture for some small $k$ and $t$. Erdős, Ko and Rado [37] proved that $d \geq 2 k$ implies

$$
\alpha[J(d, k, 0)]=\binom{d-1}{k-1} .
$$

Then Lovász [87] proved Kneser's conjecture, namely that $\chi[J(d, k, 0)]=d-2 k+2$ for $d \geq 2 k$. The following result was introduced to get a constructive bound on the Ramsey number.

Proposition 5.1.1 (Nagy, [94]). Let $d=4 s+t$, where $0 \leq t \leq 3$. Then

$$
\alpha[J(d, 3,1)]= \begin{cases}d & \text { if } t=0, \\ d-1 & \text { if } t=1, \\ d-2 & \text { if } t=2 \text { or } 3 .\end{cases}
$$

Then Larman and Rogers [82] used the bound $\chi[J(d, 3,1)] \geq \frac{|V[J(d, 3,1)]|}{\alpha[J(d, 3,1)]}$ to show that the chromatic number of the Euclidean space is at least quadratic in the dimension (initially it was proposed by Erdős and Sós). It turns out that the chromatic number of $J(d, 3,1)$ is very close to $\frac{|V[J(d, 3,1)]|}{\alpha[J(d, 3,1)]}$ (and sometimes is equal to this ratio).

Theorem 5.1.1 (Balogh-Kostochka-Raigorodskii [4]). Consider $l \geq 2$. If $d=2^{l}$, then

$$
\chi[J(d, 3,1)] \leq \frac{(d-1)(d-2)}{6}
$$

If $d=2^{l}-1$, then

$$
\chi[J(d, 3,1)] \leq \frac{d(d-1)}{6}
$$

Finally, for an arbitrary d

$$
\chi[J(d, 3,1)] \leq \frac{(d-1)(d-2)}{6}+\frac{11}{2} n .
$$

Tort [120] proved that for $d \geq 6$,

$$
\chi[K(d, 3,1)]=\left[\frac{(d-1)^{2}}{4}\right] .
$$

Zakharov [127] showed that the existence of Steiner systems (see Subsection 5.2.6) implies that

$$
\chi[J(d, k, t)] \leq(1+o(1)) \frac{(k-t-1)!}{(2 k-2 t-1)!} d^{k-t}
$$

for fixed $k>t$. In general $\chi[J(d, k, t)]=\Theta\left(d^{t+1}\right)$ for $k>2 t+1$ and $\chi[J(d, k, t)]=\Theta\left(d^{k-t}\right)$ for $k \leq 2 t+1$.

### 5.1.2 Known facts about the graphs $J_{ \pm}(d, k, t)$ and $K_{ \pm}(d, k, t)$

From a geometrical point of view $J_{ \pm}(d, k, t)$ is a natural generalization of $J(d, k, t)$. Raigorodskii [107, 108 used the graphs $J_{ \pm}(d, k, t)$ to significantly refine the asymptotic lower bounds in the Borsuk's problem and the Nelson-Hadwiger problem.

Unfortunately, there is no general method to find the independence number of $J_{ \pm}(d, k, t)$ even asymptotically. One of the reasons is that the known answers have varied and sometimes rather complicated structures. For instance the proof of the following result analogous to Proposition 5.1.1 is relatively long and the answer is quite surprising.

Theorem 5.1.2 (Cherkashin-Kulikov-Raigorodskii, [16]). For $d \geq 1$ define $c(d)$ as follows:

$$
c(d)=\left\{\begin{array}{ll}
0 & \text { if } n \equiv 0 \\
1 & \text { if } n \equiv 1 \\
2 & \text { if } n \equiv 2 \text { or } 3
\end{array} \quad(\bmod 4)\right.
$$

Then

$$
\alpha\left[J_{ \pm}(d, 3,1)\right]=\max \{6 d-28,4 d-4 c(d)\}
$$

In recent papers [46, 47, 49] Frankl and Kupavskii generalized the Erdős-Ko-Rado theorem for some subgraphs of $J_{ \pm}(d, k, t)$. We need additional definitions.

$$
\begin{gathered}
V_{k, l}:=\left\{v \in\{-1,0,1\}^{d} \mid v \text { has exactly } k^{\prime} 1^{\prime} \text { and exactly } l^{\prime}-1^{\prime}\right\} . \\
J(d, k, l, t):=\left(V_{k, l},\left\{\left(v_{1}, v_{2}\right) \mid\left\langle v_{1}, v_{2}\right\rangle=t\right\}\right)
\end{gathered}
$$

Theorem 5.1.3 (Frankl-Kupavskii, [46]). For $2 k \leq n \leq k^{2}$ the equality

$$
\alpha[J(d, k, 1,-2)]=k\binom{d-1}{k}
$$

holds. In the case $d>k^{2}$ the following equality holds

$$
\alpha[J(d, k, 1,-2)]=k\binom{k^{2}-1}{k}+\sum_{i=k^{2}}^{d-1}\binom{i}{k} .
$$

Paper [47] deals with a more generic problem.

Theorem 5.1.4 (Frankl-Kupavskii, 47]). For $2 k \leq d$ the following bounds hold

$$
\binom{d}{k+l}\binom{k+l-1}{l-1} \leq \alpha[J(d, k, l,-2 l)] \leq\binom{ d}{k+l}\binom{k+l-1}{l-1}+\binom{d}{2 l}\binom{2 l}{l}\binom{d-2 l-1}{k-l-1}
$$

In the case $2 k \leq n \leq 3 k-l$ the following equality holds

$$
\alpha[J(d, k, l,-2 l)]=\frac{k}{d}\left|V_{k, l}\right| .
$$

To introduce the next result, we will need the following definition.

## Definition 5.1.1.

$$
S(d, D):= \begin{cases}\left.\sum_{\substack{m \\ j=0 \\ d-1 \\ d \\ j \\ j \\ m}}\right)+\sum_{j=0}^{m}\binom{d}{j} & \text { if } D=2 m \\ \text { if } D=2 m+1\end{cases}
$$

In 45 (see 48 for a version with a fixed mistake) Frankl and Kupavskii determined the independence number of $K_{ \pm}(d, k, t)$ for $d>d_{0}(k, t)$ and found the asymptotics of the independence number of $J_{ \pm}(d, k, t)$ if $t<0$ and $d>d_{0}(k, t)$.

Theorem 5.1.5 (Frankl-Kupavskii, 48). For any $k \in \mathbb{N}$ and $d \geq n\left(k_{0}\right)$ we have:

1. $\alpha\left[K_{ \pm}(d, k, t)\right]=\binom{d-t-1}{k-t-1} \quad$ for $-1 \leq t \leq k-1$,
2. $\alpha\left[K_{ \pm}(d, k, t)\right]=S(k,|t|-1)\binom{d}{k} \quad$ for odd $t$ such that $-k-1 \leq t<0$,
3. $\alpha\left[K_{ \pm}(d, k, t)\right]=\alpha\left[J\left(d, k-\frac{|t|}{2}, \frac{|t|}{2}, t\right)\right]+S(k,|t|-2)\binom{d}{k} \quad$ for even $t$ such that $-k-1 \leq t<0$.

Theorem 5.1.6 (Frankl-Kupavskii, 45]). For any $k \in \mathbb{N}, t<0$ and $d>d_{0}(k, t)$ we have

$$
\alpha\left[J_{ \pm}(d, k, t)\right] \leq S(k,|t|-1)\binom{d}{k}+O\left(d^{k-1}\right)
$$

The main technique in the Frankl-Kupavskii theorems is shifting. It turns out that shifting can not increase a scalar product, so it preserves the independence property of a set in a Kneser-type graph. Unfortunately, the latter does not hold for Johnson-type graphs. Using additional arguments one can derive weaker results which are tight only in asymptotics. But it looks impossible to find the independence number of $J_{ \pm}(d, k, t)$ for $t>-k$ using shifting.

### 5.1.3 Results

Let $J(d, k$, even $)$ be a graph with the vertex set $\{0,1\}^{d}$, where edges connect vertices with even scalar product (dote that each vertex has a loop if $k$ is even). Define $J(d, k, o d d)$ in a similar way. Let $J_{ \pm}(d, k$, even $)$ and $J_{ \pm}(d, k$, odd $)$ be defined analogously to $J(d, k$, even $)$ and $J(d, k, o d d)$.

Observation 5.1.1. If $d>d_{0}(k)$, then

$$
\begin{aligned}
\alpha\left[J_{ \pm}(d, k, \text { even })\right] & =2^{k} \alpha[J(d, k, \text { even })] \\
\alpha\left[J_{ \pm}(d, k, \text { od })\right] & =2^{k} \alpha[J(d, k, \text { od } d)] .
\end{aligned}
$$

For $d>d_{0}(k)$ the exact values of $\alpha\left[J_{ \pm}(d, k\right.$, even $\left.)\right]$ and $\alpha\left[J_{ \pm}(d, k, o d d)\right]$ are determined in Theorem 5.2.6.

Proof of Observation 5.1.1. Let prt stand for odd or even.
To prove the lower bounds consider an arbitrary maximal independent set $I$ in $J(d, k, p r t)$. Then all the vertices on the supports from $I$ form an independent set $I_{ \pm}$in $J_{ \pm}(d, k, p r t)$. So

$$
\alpha\left[J_{ \pm}(d, k, p r t)\right]=2^{k} \alpha[J(d, k, p r t)] .
$$

The upper bounds simply follow from Lemma 5.2.1, since $J(d, k, p r t)$ is a subgraph of $J_{ \pm}(d, k, p r t)$.

Observation 5.1.2. For every $d \geq k$ we have

$$
\alpha\left[J_{ \pm}(d, k, k-1)\right]=2^{k} \alpha[J(d, k, k-1)] .
$$

Note that $\alpha[J(d, k, k-1)]$ is the size of a largest partial Steiner $(d, k, k-1)$-system. In particular, if the divisibility conditions hold, then $\alpha[J(d, k, k-1)]=\binom{d}{k-1} / k$ (see Subsection 5.2.6. .

Proof of Observation 5.1.2. Since $J(d, k, k-1)$ is a subset of $J_{ \pm}(d, k, k-1)$, by Lemma 5.2.1 we have

$$
\alpha\left[J_{ \pm}(d, k, k-1)\right] \leq 2^{k} \alpha[J(d, k, k-1)] .
$$

To prove the lower bound consider an arbitrary maximal independent set $I$ in the graph $J(d, k, k-$ $1)$. Then all the vertices on the supports from $I$ form an independent set $I_{ \pm}$in $J_{ \pm}(d, k, k-1)$.

We use the Katona averaging method and Reed-Solomon codes to prove the following theorem.
Theorem 5.1.7 (Cherkashin-Kiselev [15]). Suppose that $d>k 2^{k+1}$. Then

$$
\alpha\left[J_{ \pm}(d, k,-1)\right]=\binom{d}{k} .
$$

Theorem 5.1.7 can be generalized as follows.
Theorem 5.1.8 (Cherkashin-Kiselev [15]). Suppose that tis a negative odd number, $d>d_{0}(k)$. Then

$$
\alpha\left[J_{ \pm}(d, k, t)\right]=S(k,|t|-1)\binom{d}{k}
$$

where $S$ is defined in Definition 5.1.1.
The next theorem is a consequence of Theorems 5.2.1 and 5.1.7.
Theorem 5.1.9 (Cherkashin-Kiselev [15]). Let $d>\frac{9}{2} k^{3} 2^{k}$. Then

$$
\alpha\left[J_{ \pm}(d, k, 0)\right]=2\binom{d-1}{k-1}
$$

One can extract a stability version of the previous theorem from its proof.
The support of a vertex $v$ is the set of nonzero coordinates of $v$; we denote it by $\operatorname{supp} v$. Let $\mathcal{H}_{k}=\left(V_{k}, E_{k}\right)$ be a $k$-graph such that

$$
V_{k}:=\bigcup_{u \in[d]}\left\{u^{+}, u^{-}\right\}, \quad E_{k}:=\left\{\left.A \in\binom{V(\mathcal{H})}{k} \right\rvert\,\left\{u^{+}, u^{-}\right\} \not \subset A \text { for every } u\right\}
$$

There is a natural bijection between $E_{k}$ and $V\left(J_{ \pm}(d, k, t)\right)$. Introduce notion signplace for a vertex of $\mathcal{H}_{k}$ and place for a pair of vertices $\left\{u^{+}, u^{-}\right\}, u \in[d]$; note that the latter definition does not depend on $k$.

Corollary 5.1.1. Suppose that $I$ is an independent set in $J_{ \pm}(d, k, 0)$ and no place intersects all the vertices of $I$. Then

$$
|I| \leq C(k)\binom{d}{k-2}
$$

Let us proceed with the chromatic numbers in some corner cases. To warm up, let us correctly color $J_{ \pm}(d, 2,-1)$ in $2\left\lceil\log _{2} d\right\rceil+2$ colors. Let the first and second colors get vertices with non-negative and non-positive values, respectively. Only vertices of the form $a^{+} b^{-}$remain. Let us color the vertex $a^{+} b^{-}$with the color $m(a, b)$ if in the bit $m(a, b)$ the number $a$ has 1 , and the number $b$, respectively, has 0 ; let us paint the vertex $a^{-} b^{+}$in the color $-m(a, b)$. It is easy to see that all the vertices are colored, and each edge connects vertices of different colors. The total is just $2\left\lceil\log _{2} d\right\rceil+2$ colors.

The following theorem shows that the asymptotic behavior of the chromatic number is approximately two times less than in the example given.

Theorem 5.1.10 (Cherkashin [23]). For all $d \geq 2$ the inequalities are satisfied

$$
\log _{2} n \leq \chi\left(J_{ \pm}[d, 2,-1]\right) \leq \log _{2} d+\left(\frac{1}{2}+o(1)\right) \log _{2} \log _{2} d
$$

In the case $k=3, t=-1$ the picture is asymptotically the same.
Theorem 5.1.11 (Cherkashin [23]). For some positive constants $c, C$ and arbitrary $d>3$ the following inequalities hold:

$$
c \log _{2} d \leq \chi\left(J_{ \pm}[d, 3,-1]\right) \leq C \log _{2} d
$$

And for $k=3, t=-2$ we have an interesting picture.
Theorem 5.1.12 (Cherkashin [23]). For all $d \geq 3$ the inequalities are satisfied

$$
\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil \leq \chi\left(J_{ \pm}[d, 3,-2]\right) \leq 4\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil+6
$$

### 5.2 Tools

### 5.2.1 Trivial bounds on the chromatic numbers

Let $k$ be fixed and $t$ be negative. Then all vectors with non-negative coordinates form an independent set $I$ with a fraction of vertices $\frac{1}{2^{k}}$. Thus, for any Johnson type graph $G=J_{ \pm}(d, k, t)$ the classical inequality

$$
\begin{equation*}
\chi(G) \geq \frac{|V(G)|}{\alpha(G)} \tag{5.1}
\end{equation*}
$$

gives a lower bound for the chromatic number not exceeding $2^{k}$. We will see later that this inequality is often very inaccurate.

On the other hand, all graphs considered in the section are vertex-transitive. Consequently, for $G=J_{ \pm}(d, k, t)$ the inequality holds, which is true for all vertex-transitive graphs [86]

$$
\begin{equation*}
\chi(G) \leq(1+\ln \alpha(G)) \frac{|V(G)|}{\alpha(G)} \leq C(k) \log _{2} d \tag{5.2}
\end{equation*}
$$

where $C(k)$ is a constant depending only on $k$. We will see that sometimes this estimate gives the exact order of growth in $n$, in particular for $J_{ \pm}(d, 3,-1)$.

### 5.2.2 Katona averaging method

Properties of a graph with a rich group of automorphisms sometimes can be established via consideration of a proper subgraph. We say that a graph $G$ is vertex-transitive if for every vertices $v_{1}, v_{2}$, $G$ has an automorphism $f$ such that $f\left(v_{1}\right)=v_{2}$. The following lemma is a special case of Lemma 1 from [70].

Lemma 5.2 .1 (Katona, [70]). Let $G=(V, E)$ be a vertex-transitive graph. Let $H$ be a subgraph of G. Then

$$
\frac{\alpha(G)}{|V(G)|} \leq \frac{\alpha(H)}{|V(H)|}
$$

For example Lemma 5.2.1 immediately implies that for every fixed $k, t$ the following decreasing sequences converge

$$
a_{n}:=\frac{\alpha\left[J_{ \pm}(d, k, t)\right]}{\left|V\left[J_{ \pm}(d, k, t)\right]\right|} \quad \text { and } \quad b_{n}:=\frac{\alpha\left[K_{ \pm}(d, k, t)\right]}{\left|V\left[K_{ \pm}(d, k, t)\right]\right|}
$$

as $J_{ \pm}(d-1, k, t)$ and $K_{ \pm}(d-1, k, t)$ are isomorphic to subgraphs of $J_{ \pm}(d, k, t)$ and $K_{ \pm}(d, k, t)$, respectively, and both $J_{ \pm}(d, k, t)$ and $K_{ \pm}(d, k, t)$ graphs are clearly vertex-transitive.

Also since $J(d, k, t)$ is a subgraph of $J_{ \pm}(d, k, t)$, Lemma 5.2.1 implies

$$
\frac{\alpha\left[J_{ \pm}(d, k, t)\right]}{\left|V\left[J_{ \pm}(d, k, t)\right]\right|} \leq \frac{\alpha[J(d, k, t)]}{|V[J(d, k, t)]|},
$$

which gives by $\left|V\left[J_{ \pm}(d, k, t)\right]\right|=2^{k}\binom{d}{k}=2^{k}|V[J(d, k, t)]|$ the following bound

$$
\begin{equation*}
\alpha\left[J_{ \pm}(d, k, t)\right] \leq 2^{k} \alpha[J(d, k, t)] \tag{5.3}
\end{equation*}
$$

It turns out that bound (5.3) is rarely close to the optimal. On the other hand sometimes it is tight, for instance in Propositions 5.1.1 and 5.1.2.

### 5.2.3 Nontrivial intersecting families

A family of sets $\mathcal{A}$ is intersecting if every $a, b \in A$ have nonempty intersection. A transversal is a set that intersects each member of $\mathcal{A}$.

Theorem 5.2.1 (Erdős-Lovász, [38]). Let $\mathcal{A}$ be an intersecting family consisting of $k$-element sets. Then at least one of the following statements is true:
(i) $\mathcal{A}$ has a transversal of size at most $k-1$;
(ii) $|\mathcal{A}| \leq k^{k}$.

One can find better bounds in the case (ii) [3, 22, 43, 128]. In particular, for $k=3$ it is known that $3^{3}=27$ in (ii) can be replaced with 10 and this result is sharp [50].

Theorem 5.2.2 (Deza, [32]). Let $\mathcal{A}$ be a family of $k$-element sets such that $\left|A \cap A^{\prime}\right|$ is the same for all distinct $A, A^{\prime} \in \mathcal{A}$. Then at least one of the following statements is true:
(i) $A \cap A^{\prime}$ is the same for all distinct $A, A^{\prime} \in \mathcal{A}$;
(ii) $|\mathcal{A}| \leq k^{2}-k+1$.

### 5.2.4 An isodiametric inequality

Define the Hamming distance between two subsets of $[d]$ as the size of their symmetric difference. The Hamming distance between two vectors $v_{1}, v_{2} \in\{-1,0,1\}^{d}$ is the number of coordinates that differ between $v_{1}$ and $v_{2}$. The diameter of a family $\mathcal{A} \subset 2^{[d]}$ or $\mathcal{A} \subset\{-1,0,1\}^{d}$ is the maximal distance between its members.

Theorem 5.2.3 (Kleitman, [75]). Let $\mathcal{A} \subset 2^{[d]}$ be a family with diameter at most $D$ for $d>D$. Then

$$
|\mathcal{A}| \leq S(d, D)
$$

where $S$ is defined in Definition 5.1.1.
Theorem 5.2 .3 is sharp: in the case of even $D$ the equality holds for the family $\mathcal{K}(d, D):=$ $\left\{A \subset[d]:|A| \leq \frac{D}{2}\right\}$ and in the case of odd $D$ the equality holds for the family $\mathcal{K}_{x}(d, D):=\{A \subset$ $\left.[d]:|A \backslash\{x\}| \leq \frac{D}{2}\right\}$ for some fixed $x \in[d]$.

Moreover, in [42] Frankl proved the following stability result. Let $A \Delta B$ stand for the symmetric difference of the sets $A$ and $B$. We say that a family $\mathcal{A}^{\prime}$ is a translate of a family $\mathcal{A}$ if $\mathcal{A}^{\prime}=\{A \Delta T: A \in$ $\mathcal{K}(d, D)\}$ for some $T \subset[d]$.

Theorem 5.2.4 (Frankl, 42]). Let $\mathcal{A} \subset 2^{[d]}$ be a family with the diameter at most $D$ and $|\mathcal{A}|=$ $S(d, D)$ for $d \geq D+2$. Then in the case of even $D$ family $\mathcal{A}$ is a translate of $\mathcal{K}(d, D)$ and in the case of odd $D$ the family $\mathcal{A}$ is a translate of $\mathcal{K}_{y}(d, D)$.

### 5.2.5 Simple hypergraphs and Reed-Solomon codes

A hypergraph $H=(V, E)$ is a collection of (hyper)edges $E$ on a finite set of vertices $V$. A hypergraph is called $k$-uniform if every edge has size $k$. A hypergraph is simple if every two edges share at most one vertex. The following construction is a special case of Reed-Solomon codes ([89], Chapter 10); it is also known as Kuzjurin's construction [80].

Fix a prime $p>k$ and let the vertex set $V$ be the union of $k$ disjoint copies of a field with $p$ elements $\mathbb{F}=G F(p)$; call them $\mathbb{F}_{1}, \ldots, \mathbb{F}_{k}$. Consider the following system of linear equations

$$
\sum_{i=1}^{k} i^{j} x_{i}=0, \quad j=0,1, \ldots, k-3
$$

over $\mathbb{F}_{p}$. The solutions $\left\{x_{1}, \ldots x_{k}\right\} \in \mathbb{F}_{1} \sqcup \cdots \sqcup \mathbb{F}_{k}$, where $x_{i} \in \mathbb{F}_{i}$, form the edge set $E$. Fixing two arbitrary variables there is a unique solution over $\mathbb{F}_{p}$, because the corresponding square matrix is a Vandermonde matrix with nonzero determinant. It means that there are $p^{2}$ different solutions and $\left|e_{1} \cap e_{2}\right| \leq 1$ for every distinct $e_{1}, e_{2} \in E$. Summing up, $H_{p}(k):=(V, E)$ is a $p$-regular $k$-uniform simple hypergraph with $|V|=p k$ and $|E|=p^{2}$.

A $k$-uniform hypergraph is $b$-simple if every two edges share at most $b$ vertices. The same construction with $k-b-1$ equations gives an example of a $k$-uniform $b$-simple hypergraph $H(p, k, b)$.

Further we use regularity of $H=H(p, k, b)$ in the following sense. Consider an arbitrary vertex subset $A$ of size $b$. If $A$ contains at most 1 vertex from every copy of $\mathbb{F}_{p}$, then $H$ has exactly $p$ hyperedges containing $A$; otherwise $H$ contains no such edges. Slightly abusing the notation we say that $b$-degree of $H$ is $p$.

### 5.2.6 Steiner systems

A Steiner system with parameters $d, k$ and $l$ is a collection of $k$-subsets of $[d]$ such that every $l$-subset of $[d]$ is contained in exactly one set of the collection. There are some obvious necessary 'divisibility conditions' for the existence of Steiner ( $d, k, l$ )-system:

$$
\binom{k-i}{l-i} \text { divides }\binom{d-i}{k-i} \text { for every } 0 \leq i \leq k-1 .
$$

In a breakthrough paper [71] Keevash proved the existence of Steiner ( $d, k, l$ )-systems for fixed $k$ and $l$ under the divisibility conditions and for $d>d_{0}(k, l)$ (different proofs can be found in [56, 72]).

Partial Steiner system. When the divisibility conditions do not hold we are still able to construct a large partial Steiner system, that is, a collection of $k$-subsets of $[d]$ such that every $l$-subset of $[d]$ is contained in at most one set of the collection. Rödl confirmed a conjecture of Erdős and Hanani and proved the following theorem.

Theorem 5.2.5 (Rödl, [110]). For every fixed $k$ and $l<k$, and for every $d$ there exists a partial ( $d, k, l$ )-system with

$$
(1-o(1))\binom{d}{l} /\binom{k}{l}
$$

$k$-subsets.
Later the result was refined in [59, 73, 77]. Also it follows from the mentioned results on Steiner systems.

### 5.2.7 Families with even or odd intersections

Recall that $J(d, k$, even $)$ and $J(d, k, o d d)$ were defined in Subsection 5.1.3. Frankl and Tokushige determined the independence numbers of these graphs.

Theorem 5.2.6 (Frankl-Tokushige, [51]). Let $d \geq d_{0}(k)$. Then

$$
\begin{array}{rlr}
\alpha[J(d, k, \text { odd })] & =\binom{\lfloor d / 2\rfloor}{ k / 2} & \text { for even } k, \\
\alpha[J(d, k, \text { even })] & =\binom{\lfloor(d-1) / 2\rfloor}{(k-1) / 2} & \text { for odd } k .
\end{array}
$$

In the case when $k$ is even, the equality is achieved for the following family: we split [d] into pairs and take all sets consisting of $k / 2$ pairs. In the case when $k$ is odd we also add a fixed point $x \in[d]$ to each constructed set.

### 5.3 Examples

Let us start with a simple example which is rarely close to the independence number.
Example 5.3.1. Let $t<0, k>|t|$. Then $\alpha\left[J_{ \pm}(d, k, t)\right] \geq 2^{|t|-1}\binom{d}{k}$.
Proof. Fix an ordering of the coordinates. Take all vertices of $J_{ \pm}(d, k, t)$ with the first $k-|t|+1$ nonzero coordinates equal to 1 . Any two such vertices can have different signs on at most $|t|-1$ positions, therefore their scalar product is at least $-|t|+1=t+1$.

The following example is a part of Theorem 5.1.5.
Example 5.3.2. For any $t<0$ and $k>|t|$ we have

$$
\alpha\left[J_{ \pm}(d, k, t)\right] \geq S(k,|t|-1)\binom{d}{k}
$$

and for even $t$ we also have

$$
\alpha\left[J_{ \pm}(d, k, t)\right] \geq S(k,|t|-1)\binom{d}{k}+\binom{k-1}{|t| / 2}
$$

Proof. We start with the first bound for the case of odd $t$. Let $I_{o d d}$ be the set of all vertices of $J_{ \pm}(d, k, t)$ with at most $(|t|-1) / 2$ negative entries. Each $k$-set is the support of exactly

$$
\sum_{j=0}^{(|t|-1) / 2}\binom{k}{j}=S(k,|t|-1)
$$

vertices in $I_{o d d}$. Any two vectors in $I_{o d d}$ may differ in at most $2(|t|-1) / 2=|t|-1$ coordinates, so their scalar product is at least $t+1$, and $I_{\text {odd }}$ is an independent set of the desired size.

Now we deal with the case of even $t$. Fix an ordering of the coordinates. For every $k$-set $f$ add to $I_{\text {even }}$ all the vertices with support $f$ and with at most $|t| / 2-1$ negative entries on $f$ and all the vertices with -1 on the last coordinate of $f$ and exactly $|t| / 2-1$ other negative coordinates. Then each $k$-set is the support of exactly

$$
\sum_{j=0}^{|t| / 2-1}\binom{k}{j}+\binom{k-1}{|t| / 2-1}=S(k,|t|-1)
$$

vertices in $I_{\text {even }}$. Assume that $I_{\text {even }}$ is not independent, i.e. the scalar product of some $v_{1}, v_{2} \in I_{\text {even }}$ is equal to $t$. Then $v_{1}$ and $v_{2}$ together have at least $|t|$ negative entries. Hence both $v_{1}$ and $v_{2}$ have exactly $|t| / 2$ negative entries, so both $v_{1}$ and $v_{2}$ have -1 at the last coordinates $x_{1}$ and $x_{2}$ of $\operatorname{supp} v_{1}$ and $\operatorname{supp} v_{2}$, respectively. But then both $v_{1}$ and $v_{2}$ can not have +1 at coordinates $x_{2}$ and $x_{1}$ respectively, so the scalar product is at least $t+1$. This contradiction shows that $I_{\text {even }}$ is an independent set of the desired size.

Now we proceed to the second bound. Let us add to $I_{\text {even }}$ all the vertices on the lexicographically first support $\{1, \ldots, k\}$ with exactly $|t| / 2$ negative entries and having +1 at the $k$-th coordinate. Obviously the resulting set $I$ has the claimed size. By definition, no edge connects two vertices from $I$ on the support $\{1, \ldots, k\}$.

Consider a vertex $v$ from $I_{\text {even }}$ and a vertex $u \in I \backslash I_{\text {even }}$. Note that $u$ and $v$ together have at most $|t|$ negative entries. Since the largest coordinate of $\operatorname{supp} v$ is greater than $k$ and $v$ has -1 in this coordinate, the scalar product of $u$ and $v$ is at least $t+1$. Thus $I$ is independent.

Example 5.3.3. For $t \geq 0$ we have

$$
\alpha\left[J_{ \pm}(d, k, t)\right] \geq 2 \alpha[J(d, k, t)]
$$

Proof. Let $I \subset V[J(d, k, t)]$ be an independent set of size $\alpha[J(d, k, t)]$. Define $I_{ \pm}$as a subset of $V\left[J_{ \pm}(d, k, t)\right]$ consisting of vertices with all positive or all negative entries on every support $f=\operatorname{supp} v$, $v \in I$. It is easy to see that the subset $I_{ \pm}$is independent in $J_{ \pm}(d, k, t)$.

### 5.4 Proofs

### 5.4.1 Proof of Theorem 5.1.7

We start with the lower bound. One can take the vertices only with non-negative coordinates (so exactly one vertex on each support is taken); obviously the scalar product of such vertices is always non-negative, so

$$
\alpha\left[J_{ \pm}(d, k,-1)\right] \geq\binom{ d}{k}
$$

Now we will show the upper bound. Denote $G:=J_{ \pm}(d, k,-1)$. Fix a prime $p, d /(2 k) \leq p \leq n / k$ (so by the statement of the theorem $p>2^{k}$ ), and let $H:=H_{p}(k)$ (see Subsection 5.2.5) be a $p$-regular $k$-uniform simple hypergraph with $V(H) \subset[d]$. Define graph $G[H]$ as a subgraph of $G$, consisting of the vertices with support on edges of $H$. So we have

$$
|V(G[H])|=2^{k}|E(H)| .
$$

Fix an independent set $I$ in $G[H]$; consider the set $X \subset[d]$ of coordinates on which the vertices from $I$ have both signs. Denote by supp $I$ the set of all supports of vertices from $I(\operatorname{supp} I \subset E(H))$ and for a given $e \in E(H)$ put $e_{X}:=e \cap X$.

Note that $I$ has at most $2^{\left|e_{X}\right|}$ vertices on the support $e\left(\left|e_{X}\right|\right.$ might be zero). Hence

$$
\begin{equation*}
|I| \leq \sum_{e \in \operatorname{supp} I} 2^{\left|e_{X}\right|} \leq \sum_{e \in E(H):\left|e_{X}\right|=0} 2^{\left|e_{X}\right|}+\sum_{e \in E(H):\left|e_{X}\right|>0} 2^{\left|e_{X}\right|} \leq\left|\left\{e \in E(H):\left|e_{X}\right|=0\right\}\right|+\sum_{e \in E(H):\left|e_{X}\right|>0} 2^{\left|e_{X}\right|} \tag{5.4}
\end{equation*}
$$

Let us show that $e_{X}$ form a disjoint cover of $X$. Suppose the contrary, i.e. there are $e, f \in \operatorname{supp} I$ such that $e_{X} \cap f_{X} \neq \emptyset$. Since the hypergraph $H$ is simple, and $e, f$ correspond to its hyperedges, we have $\left|e_{X} \cap f_{X}\right|=|e \cap f|=1$. Put $\{u\}:=e \cap f$. By the definition of $X$ there are vertices $v_{1}$, $v_{2} \in I$ having different signs on $u$. Since $I$ is independent and any two different supports intersect in
at most 1 coordinate, $v_{1}$ and $v_{2}$ have the same support (say, not $f$ ). So every vertex of $G[H]$ with support $f$ forms an edge in $G[H]$ with one of $v_{1}$ or $v_{2}$, thus $I$ is not independent; contradiction.

So $\sum\left|e_{X}\right|=|X|$. Since the sequence $2^{k} / k, k \geq 1$, is non-decreasing and

$$
\frac{a_{1}+a_{2}+\ldots+a_{t}}{b_{1}+b_{2}+\ldots+b_{t}} \leq \max \left(\frac{a_{1}}{b_{1}}, \frac{a_{2}}{b_{2}}, \ldots, \frac{a_{t}}{b_{t}}\right),
$$

we have

$$
\begin{equation*}
\sum_{e \in E(H):\left|e_{X}\right|>0} 2^{\left|e_{X}\right|} \leq \frac{|X|}{k} 2^{k} . \tag{5.5}
\end{equation*}
$$

By definition, every $e \in\left\{e \in E(H):\left|e_{X}\right|=0\right\}$ has empty intersection with $X$. Since $H$ is $p$ regular, $X$ intersects at least $\frac{p|X|}{k}$ edges of $H$ (since every $k$-edge is counted at most $k$ times), so

$$
\begin{equation*}
\left|\left\{e \in E(H):\left|e_{X}\right|=0\right\}\right| \leq|E(H)|-\frac{p|X|}{k} \tag{5.6}
\end{equation*}
$$

Summing up, by (5.4, (5.5), 5.6) and the choice of $p$, we have

$$
|I| \leq|E(H)|-\frac{|X|}{k} p+\frac{|X|}{k} 2^{k} \leq|E(H)|
$$

which implies $\alpha(G[H]) \leq|E(H)|$, hence

$$
\frac{V(G[H])}{\alpha(G[H])} \geq 2^{k}
$$

By the definition $G[H]$ is a subgraph of the graph $G$, so Lemma 5.2.1 finishes the proof.
For some $k$ one can choose a smaller $H$ and require a weaker inequality for $d$, for instance in the case $k=3$ (see Subsection 5.5.2).

### 5.4.2 Proof of Theorem 5.1.8

This is a generalization of the proof of Theorem 5.1.7. The lower bound is provided in the first part of Example 5.3.2

Denote $G=J_{ \pm}(d, k, t)$ during the proof. The case $t=-k$ is obvious, because $J_{ \pm}(d, k,-k)$ is a matching. From now $|t| \leq k-1$. Fix $d$ and a prime $p \leq n / k$ to be large enough. Let $H=H(p, k,|t|)$ (see Subsection 5.2.5) be a $k$-uniform $|t|$-simple hypergraph with $|t|$-codegree $p$. Fix an embedding of $V(H)$ into $[d]$.

Define $G[H]$ as a subgraph of $G$, consisting of all the vertices with support on edges of $H$. Fix an independent set $I$ in $G[H]$.

Let an object $O$ be a pair of opposite vectors $\{o,-o\}$ with support of size $|t|$ with $\{0, \pm 1\}$ entries. Let $\mathcal{X}$ be the set of objects $O=\{o,-o\}$ such that $\left(v_{1}, o\right)=\left(v_{2},-o\right)=|t|$ for some vertices $v_{1}, v_{2} \in I$ (this means that $v_{1}$ and $v_{2}$ coincide on supp $o$ with $o$ and $-o$, respectively).

Let $E_{\text {tight }}$ be the set consisting of such edges $e \in E(H)$ that diam $I[e]<|t|$, where $I[e]$ stands for the set of vertices of $I$ with the support $e$. Put $E_{\text {wide }}:=E(H) \backslash E_{\text {tight }}$. Let $I_{\text {tight }}$ and $I_{\text {wide }}$ stand for the sets of vertices of $I$ with the support from $E_{\text {tight }}$ and $E_{\text {wide }}$ respectively. Then $|I|=\left|I_{\text {tight }}\right|+\left|I_{\text {wide }}\right|$.

Consider an arbitrary support $e \in E_{\text {wide }}$; by the definition of $E_{\text {wide }}$ there are an object $X=$ $\{x,-x\} \in \mathcal{X}$ and vertices $v_{1}, v_{2} \in I[e]$, such that $\left(v_{1}, x\right)=\left(v_{2},-x\right)=|t|$. Since $I$ is independent and $H$ is $|t|$-simple, distinct $e_{1}$ and $e_{2} \in E_{\text {wide }}$ cannot lead to the same $X \in \mathcal{X}$, so

$$
\begin{equation*}
|\mathcal{X}| \geq\left|E_{\text {wide }}\right| . \tag{5.7}
\end{equation*}
$$

For every $e \in E_{\text {tight }}$ we have $\operatorname{diam} I[e]<|t|$, thus Theorem 5.2 .3 implies

$$
\begin{equation*}
|I[e]| \leq S(k,|t|-1) \tag{5.8}
\end{equation*}
$$

Let us study the case of equality in (5.8). Fix a support $f \subset[d],|f|=k$, and consider a family $\mathcal{A} \subset\{-1,1\}^{f}$ with diameter at most $|t|-1$ and size $S(k,|t|-1)$. Also consider an object $O=\{o,-o\}$ such that $\operatorname{supp} O \subset f$ (recall that $|\operatorname{supp} O|=|t|)$. By the pigeon-hole principle and the oddity of $t$, one of $o,-o$ has at most $(|t|-1) / 2$ negative entries. Thus there is a vector $v$ from $\mathcal{K}(k,|t|-1)$ such that $(v, o)=|t|$ or $(v,-o)=|t|$. By Theorem 5.2.4 $\mathcal{A}$ is a translate of $\mathcal{K}(k,|t|-1)$, so the previous conclusion also holds for $\mathcal{A}$.

Fix an object $X \in \mathcal{X}$ and consider an arbitrary support $e \in E_{\text {tight }}$ containing supp $X$. Assume that $(v, x)= \pm t$ for some $v \in I[e]$. Consider a support $g \in E_{\text {wide }}$ such that there are $u_{1}, u_{2} \in I[g]$, satisfying $\left(u_{1}, x\right)=\left(u_{2},-x\right)=t(g$ exists because $X \in \mathcal{X})$. Since $H$ is $|t|$-simple, $\left(u_{1}, v\right)=t$ or $\left(u_{2}, v\right)=t$; a contradiction. By Theorem 5.2.4 we can refine the bound (5.8) in this case:

$$
\begin{equation*}
|I[e]| \leq S(k,|t|-1)-1 . \tag{5.9}
\end{equation*}
$$

By the construction of $H$ for every $X \in \mathcal{X}, \operatorname{supp} X$ is contained in exactly $p$ edges of $H$ (because it is contained in at least one edge). Every edge of $H$ is the support of $\binom{k}{|t|} 2^{|t|-1}$ objects, so is counted above at most $\binom{k}{|t|} 2^{|t|-1}$ times. By 5.7 at most $|\mathcal{X}|$ of the edges are wide. So the refined bound (5.9) is applicable to at least

$$
\frac{p|\mathcal{X}|}{\binom{k}{|t|}^{|t|-1}}-|\mathcal{X}|
$$

tight edges. Then

$$
\left|I_{\text {tight }}\right| \leq S(k,|t|-1)\left|E_{\text {tight }}\right|-\frac{p|\mathcal{X}|}{\binom{k}{|t|} 2^{|t|-1}}+|\mathcal{X}| .
$$

On the other hand there is a straightforward bound

$$
\left|I_{\text {wide }}\right| \leq 2^{k}\left|E_{\text {wide }}\right| \leq 2^{k}|\mathcal{X}|
$$

Putting it all together

$$
\begin{equation*}
|I|=\left|I_{\text {tight }}\right|+\left|I_{\text {wide }}\right| \leq S(k,|t|-1)\left|E_{\text {tight }}\right|-\frac{p|\mathcal{X}|}{\binom{k}{|t|} 2^{|t|-1}}+\left(2^{k}+1\right)|\mathcal{X}| . \tag{5.10}
\end{equation*}
$$

For a large $d$ (then $p$ is also large enough) the inequality (5.10) implies

$$
\alpha(G[H]) \leq S(k,|t|-1)|E(H)| .
$$

By the definition $G[H]$ is a subgraph $G$ and

$$
\frac{\alpha(G[H])}{V(G[H])} \leq \frac{S(k,|t|-1)}{2^{k}}
$$

so Lemma 5.2.1 finishes the proof.

### 5.4.3 Proof of Theorem 5.1.9

Consider an arbitrary independent set $I$ in the graph $J_{ \pm}(d, k, 0)$. Note that supports of the vertices of $I$ form an intersecting family; denote it by $F$. Let $U$ be a minimal (by inclusion) transversal of $F$. As $U$ is minimal, for every coordinate $a \in U$ there is a vertex $x_{a} \in I$, such that supp $x_{a} \cap U=\{a\}$.

In the case $|U|>1$ we can consider the set

$$
C:=U \cup \operatorname{supp} x_{a} \cup \operatorname{supp} x_{b}
$$

for two different $a, b \in U$. Note that $|C| \leq 3 k$ and every $f \in F$ intersects $C$ in at least two places (suppose that $|f \cap U|=1$, then it should intersect either $\left(\operatorname{supp} x_{a}\right) \backslash U$ or $\left.\left(\operatorname{supp} x_{b}\right) \backslash U\right)$. Hence

$$
|I| \leq 2^{k}\binom{|C|}{2}\binom{d}{k-2}<2^{k} \frac{9 k^{2}}{2}\binom{d}{k-2}
$$

Recall that $d>\frac{9}{2} k^{3} 2^{k}$, so

$$
2^{k} \frac{9 k^{2}}{2}\binom{d}{k-2}<2^{k} \frac{9 k^{2}}{2} \frac{d^{k-2}}{(k-2)!}<\frac{d}{k-1} \frac{d^{k-2}}{(k-2)!}<2\binom{d-1}{k-1}
$$

The remaining case is $|U|=1$, say $U=\{u\}$. Consider only vertices containing $u^{+}$, by Theorem 5.1.7 we have at most $\binom{d-1}{k-1}$ such vertices. The same bound for $u^{-}$gives the desired bound.

Example 3 and Erdős-Ko-Rado theorem give a lower bound.

### 5.4.4 Proof of Corollary 5.1.1

Let us repeat the proof of Theorem 5.1.9. Let $I$ be an arbitrary independent set in $J_{ \pm}(d, k, 0)$. Then

$$
|I|<2^{k} \frac{9 k^{2}}{2}\binom{d}{k-2}
$$

or the family of all supports of vertices from $I$ has a transversal of size 1. The first possibility implies

$$
|I| \leq C(k)\binom{d}{k-2}
$$

the latter one contradicts the condition of the corollary.

### 5.4.5 Proof of Theorem 5.1.10

Proof. Let us start with the lower estimate. Consider a subset of vertices containing coordinates of different signs; let us call it $V_{ \pm}$. Suppose that we covered $V_{ \pm}$with independent sets $I_{1}, \ldots, I_{q}$. Consider the set $I_{j}$. If at some coordinate $a$ vertices from $I_{j}$ take values of both signs, then all vertices from $I_{j}$ whose support contains $a$ are $a^{+} b^{-}$and $a^{-} b^{+}$for some coordinate $b$. Thus, on coordinate $b$, the vertices from $I_{j}$ also take values of different signs; let us call such coordinates diverse, and the remaining coordinates positive and negative, respectively.

Let us define an auxiliary graph $G_{j}$ whose vertices are coordinates, and an edge connects a pair of coordinates if this pair is the support of a vertex from $I_{j}$. Note that $G_{j}$ is a bipartite graph: indeed, the support either contains positive and negative coordinates (since we consider only vertices from $V_{ \pm}$) or two different ones; We showed above that $G_{j}$ is a matching on various coordinates.

Since $I_{1}, \ldots, I_{q}$ cover $V_{ \pm}$, the graphs $G_{j}$ cover the complete graph $K_{[d]}$ on the set of coordinates. It is well known that such coverage requires at least $\log _{2} d$ bipartite subgraphs. Indeed, if $F_{i}:=$ $G_{1} \cup \cdots \cup G_{i}$, then $\alpha\left(F_{i+1}\right) \geq \alpha\left(F_{i}\right) / 2$; on the other hand, $\alpha\left(K_{[d]}\right)=1$.

Let us move on to an example. We will color $J_{ \pm}(d, 2,-1)$ with $2 m+2$ colors for $n \leq\binom{ 2 m+1}{m}$. Let us associate with each coordinate an $m$-element subset $[2 m+1]$; let us denote the matching by $f$. For $1 \leq i \leq 2 m+1$ the color $I_{i}$ consists of vertices $a^{+} b^{+}$for which $i \in f(a), f(b)$, of vertices $a^{+} b^{-}$for which $i \in f(a), i \notin f(b)$ and vertices $a^{-} b^{-}$for which $i \notin f(a), f(b)$. Note that all vertices of the form $a^{+} b^{-}$are covered by these colors. Indeed, for any m-element subsets $f(a), f(b)$ there is an element $i$ that belongs to $f(a)$ but not to $f(b)$. Similarly, all vertices of the form $a^{-} b^{-}$are covered with colors, since for any two $m$-element subsets of a $(2 m+1)$-element set there is an element from their common complement. The last color with number $2 m+2$ contains all vertices of the form $a^{+} b^{+}$.

Since $J_{ \pm}\left(d_{1}, 2,-1\right)$ is a subgraph of $J_{ \pm}\left(d_{2}, 2,-1\right)$ for $d_{1}<d_{2}$, it remains to check that the inequality $d \geq\binom{ 2 m-1}{m-1}$ implies the inequality $2 m+2 \leq \log _{2} d+(1+o(1)) \frac{1}{2} \log _{2} \log _{2} d$. This is true because

$$
\binom{2 m-1}{m-1}=\frac{1}{2}\binom{2 m}{m}=\Omega\left(\frac{4^{m}}{\sqrt{m}}\right)
$$

### 5.4.6 Proof of Theorem 5.1.11

Let $I$ be an independent set of the graph $J_{ \pm}(d, 3,-1)$. We call a coordinate $i \in[d]$ diverse if the vertices from $I$ take both non-zero values on $i$. We call a vertex $v \in I$ special if the support of $v$ contains a diverse coordinate.

Lemma 5.4.1. Let I be an independent set of the graph $J_{ \pm}(d, 3,-1)$ for which the $t$ coordinates are diverse. Then

$$
|I| \leq 8 t(d-2)+\binom{d-t}{3}
$$

Proof. Consider a varied coordinate $i$. Let $I_{i}$ be a subset of vertices $I$ whose support contains $i$. Then $I_{i}$ contains vertex $v_{+}$with support $\{i, a, b\}$ and sign + on $i$, as well as vertex $v_{-}$with support $\{i, c, e\}$ and the sign - on $i$ (the sets $\{a, b\}$ and $\{c, e\}$ must intersect). Then any vertex $v \in I_{i}$ with sign + on $i$ intersects $\{c, e\}$ (otherwise there is an edge between $v_{+}$and $v_{-}$, which contradicts the independence of $I$ ); similarly, any $v$ with a $-\operatorname{sign}$ on $i$ intersects $\{a, b\}$. Thus

$$
\left|I_{i}\right| \leq 8(d-2) \quad \text { and } \quad \bigcup_{1 \leq i \leq t} I_{i} \leq 8 t(d-2)
$$

We counted all vertices for which at least one coordinate is diverse. There are vertices left whose support lies on $d-t$ uniform coordinates. There are no more than $\binom{d-t}{3}$.
Corollary 5.4.1. Let I be an independent set of the graph $J_{ \pm}(d, 3,-1)$. Then the number of special vertices does not exceed $\mathrm{cn}^{2}$.

The upper bound follows from the fact that the graph $J_{ \pm}(d, 3,-1)$ is vertex-transitive, as was shown in the introduction.

Let us move on to the lower estimate. Consider the partition of the set of vertices $V$ of the graph $J_{ \pm}(d, 3,-1)$ into independent sets $I_{1}, \ldots, I_{k}$. Consider the set $V_{ \pm} \subset V$, consisting of vertices that have coordinates of different signs; the cardinality of $V_{ \pm}$is $6\binom{d}{3}$. Note that for each pair of
coordinates $(i, j)$ there are $4(d-2)$ vertices from $V_{ \pm}$in which $i$ and $j$ have different signs. We say that a set of vertices splits a pair of coordinates $(i, j)$ if the set contains all these $4(d-2)$ vertices. Then, by Corollary 5.4.1, the special vertices of all independent sets in the union split $O(k d)$ pairs of coordinates.

Completing the proof is similar to Theorem 5.1.10. Let independent sets $I_{1}, \ldots, I_{q}$ cover the graph $G$. Consider a bipartite graph $F_{i}$ on a set of coordinates whose parts are uniform coordinates $I_{i}$ of positive and negative signs. As was shown in the proof of Theorem5.1.10, the graph $F:=F_{1} \cup \cdots \cup F_{q}$ has an independent set of size at least $d / 2^{q}$. For $q<\frac{1}{3} \log _{2} d$ we have

$$
\alpha(F) \geq d^{2 / 3}
$$

Then all pairs of coordinates on an independent set of size $d^{2 / 3}$ must be separated by different coordinates; there are asymptotically more of them than $k d$, which is a contradiction.

### 5.4.7 Proof of Theorem 5.1.12

Let us start with the lower estimate. Let us assume that we have covered all the vertices of the graph with independent sets $I_{1}, \ldots, I_{q}$. Let us fix the order of coordinates $[d]$ and consider only the subset of vertices $V_{\text {alt }}$ in which the signs alternate, that is, vertices with support $\{a, b, c\}$, where $a<b<c$, have the form $a^{+} b^{-} c^{+}$and $a^{-} b^{+} c^{-}$.

Consider an auxiliary graph $H$, the vertices of which are unordered pairs of coordinates, and the edges are drawn between pairs of the form $\{a, b\}$ and $\{b, c\}$, where $a<b<c$. Then, for each edge $H$, the union of vertices, as a support, contains two vertices of the graph $G$ from $V_{a l t}$; and vice versa, the support of any vertex from $V_{a l t}$ is uniquely obtained as the union of vertices corresponding to the edge $H$.
Lemma 5.4.2. For each $1 \leq j \leq q$, the edges corresponding to the supports of $I_{j}$ form a bipartite subgraph in $H$.

Let us call this subgraph $H_{j}$.
Proof. Let us put two labels for each vertex from $I_{j} \cap V_{a l t}$ : if the vertex has the form $a^{+} b^{-} c^{+}$(in accordance with the order on $[d])$, then pairs of coordinates $\{a, b\}$ and $\{b, c\}$ receive labels $L$ and $R$, respectively, and if the form is $a^{-} b^{+} c^{-}$, then vice versa. Let some pair of coordinates $\{e, f\}$ receive both label $L$ and label $R$ from vertices $v_{1}$ and $v_{2}$. Then the supports of $v_{1}$ and $v_{2}$ coincide, otherwise their scalar product is equal to -2 , which contradicts the independence of $I_{j}$. But then the support of any other vertex from $I_{j} \cap V_{\text {alt }}$ does not contain a pair of coordinates $\{e, f\}$, otherwise an edge appears with either $v_{1}$ or $v_{2}$.

Now we define vertices $H$ only with the label $L$ in one part, and only with the label $R$ in the other. The subgraph on these vertices is bipartite. Vertices with both labels, as we showed above, are pendant, so adding them leaves the graph bipartite.

Let complete bipartite subgraphs on parts of subgraphs $H_{i}$ partition $V(H)$ into independent (in the graph $H$ ) sets $J_{1}, \ldots, J_{w}$.

Lemma 5.4.3. Let $J$ be an independent set of the graph $H$. Then there is a partition $[d]=B \sqcup E$ such that for any vertex $\{b, e\} \in J$ we have $b \in B, e \in E$ and $b<e$.
Proof. No coordinate can be the first at vertex $j_{1} \in J$ and the second at $j_{2} \in J$, since $J$ is an independent set, and such $j_{1}$ and $j_{2}$ would form an edge. This allows us to define $B$ as the set of first coordinates of the vertices $j \in J$, and $E$ as $[d] \backslash B$.

For each graph $J_{i}$, consider the partition $B_{i} \sqcup E_{i}$ from Lemma 5.4.2. Let some pair of coordinates not lie in different parts of any partition $B_{i} \sqcup E_{i}$. Then the corresponding vertex of the graph $H$ does not lie in the union of independent sets $J_{1}, \ldots, J_{w}$, a contradiction. Consider the union $F$ of bipartite graphs on the parts $B_{i}, E_{i}$ for $1 \leq i \leq w$. All parts of $F$ have size 1 , hence $w \geq\left\lceil\log _{2} d\right\rceil$, similar to the corresponding part of the proof of Theorem 5.1.10.

It turns out that subgraphs $H_{1}, \ldots, H_{q}$ partition $V(H)$ into at least $\left\lceil\log _{2} d\right\rceil$ sets, that is, subgraphs at least $\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil$, which completes the proof of the lower bound.

Let us have a deal with the upper bound and demonstrate the coloring of the graph in

$$
4\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil+6
$$

colors.
First, let us color all vertices of the form $a^{+} b^{-} c^{+}$, where $a<b<c$, in $2\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil$ colors. Let us assign the color $\operatorname{sign}(m(a, b)-m(b, c)) \cdot m(m(a, b), m(b, c))$ to such a vertex. Note that for natural numbers $x<y<z$ the inequality $m(x, y) \neq m(y, z)$ holds; indeed: from $x<y$ it follows that in the bit $m(x, y)$ the number $x$ takes the value 0 , and the number $y$ takes 1 ; similar reasoning for $y$ and $z$ entails a contradiction.

Let there be a pair of vertices of color $j$ on supports $\{a, b, c\}$, where $a<b<c$, and $\{b, c, e\}$ with scalar product -2 . Then, since the first vertex has the form $a^{+} b^{-} c^{+}$, the second vertex has the signs $b^{+} c^{-}$, hence the coordinates are of order $a<b<c<e$. Since the vertices are of the same color, the expressions $m(a, b)-m(b, c)$ and $m(b, c)-m(c, e)$ have the same sign. Let $m(a, b)<m(b, c)<m(c, e)$, then $m(m(a, b), m(b, c))=m(m(b, c), m(c, e))$ means that $m(b, c)$ has a one on one side and a zero on the other side; contradiction. The case $m(a, b)>m(b, c)>m(c, e)$ is treated in the same way.

Similarly, one can color vertices of the form $a^{-} b^{+} c^{-}(a<b<c)$ in $2\left\lceil\log _{2}\left\lceil\log _{2} d\right\rceil\right\rceil$ colors.
Finally, the last six colors consist of vertices of the form $a^{+} b^{+} c^{+}, a^{+} b^{+} c^{-}, a^{+} b^{-} c^{-}, a^{-} b^{+} c^{+}, a^{-} b^{-} c^{+}$ and $a^{-} b^{-} c^{-}$, respectively, where $a<b<c$. A direct check shows that with such a coloring there are no edges of the same color.

### 5.5 Independent numbers in the case $k \leq 3$

We have implemented Östergård algorithm [97] to find independence numbers of several small graphs. All the calculations were done on a standard laptop in a few hours. The source can be found in [74].

### 5.5.1 The case $k=2$

The case $t=-1$. By simple calculations we have

$$
\alpha\left[J_{ \pm}(2,2,-1)\right]=\alpha\left[J_{ \pm}(3,2,-1)\right]=4, \quad \alpha\left[J_{ \pm}(4,2,-1)\right]=8, \quad \alpha\left[J_{ \pm}(5,2,-1)\right]=10 .
$$

In Section 5.2 .2 we show that the sequence

$$
\frac{\alpha\left[J_{ \pm}(d, 2,-1)\right]}{\left|V\left[J_{ \pm}(d, 2,-1)\right]\right|}
$$

is non-increasing, so

$$
\alpha\left[J_{ \pm}(d, 2,-1)\right]=\binom{d}{2}
$$

for $d \geq 5$.

The case $t=0$. It is straightforward to check that

$$
\alpha\left[J_{ \pm}(2,2,0)\right]=2, \quad \alpha\left[J_{ \pm}(3,2,0)\right]=\alpha\left[J_{ \pm}(4,2,0)\right]=6
$$

For the case $d>4$ we can repeat the proof of the Theorem 5.1.9 and show that $\alpha\left[J_{ \pm}(d, 2,0)\right]=2(d-1)$.
The case $t=1$. From Proposition 5.1 .2 we have

$$
\begin{array}{lr}
\alpha\left[J_{ \pm}(d, 2,1)\right]=2 n & \text { for even } d \\
\alpha\left[J_{ \pm}(d, 2,1)\right]=2(d-1) & \text { for odd } d .
\end{array}
$$

### 5.5.2 The case $k=3, t=-1$

Proposition 5.5.1. Let $d \geq 7$. Then

$$
\alpha\left[J_{ \pm}(d, 3,-1)\right]=\binom{d}{3} .
$$

Proof. Fano plane is the projective plane over $G F(2)$ i.e. the following simple 3 -graph on 7 vertices

$$
\{1,2,3\},\{1,4,7\},\{1,5,6\},\{2,4,6\},\{2,5,7\},\{3,4,5\},\{3,6,7\}
$$

Consider an arbitrary embedding $F$ of the Fano plane into $V\left[J_{ \pm}(d, 3,-1)\right]$. As usual consider the subgraph $G[F]$; it has $7 \cdot 2^{3}=56$ vertices. One may check by hands or via computer that $\alpha(G[F])=7$. By Lemma 5.2.1

$$
\alpha\left[J_{ \pm}(d, 3,-1)\right] \leq\binom{ d}{3}
$$

On the other hand, Example 5.3.1 implies $\alpha\left[J_{ \pm}(d, 3,-1)\right]=\binom{d}{3}$.
By the computer calculations we have

$$
\alpha\left[J_{ \pm}(6,3,-1)\right]=21>\binom{6}{3}=20
$$

so Proposition 5.5.1 is sharp. Also

$$
\alpha\left[J_{ \pm}(5,3,-1)\right]=14, \quad \alpha\left[J_{ \pm}(4,3,-1)\right]=8, \quad \alpha\left[J_{ \pm}(3,3,-1)\right]=2 .
$$

### 5.5.3 The case $k=3, t=0$

By the computer calculations we have

$$
\begin{gathered}
\alpha\left[J_{ \pm}(3,3,0)\right]=\alpha\left[J_{ \pm}(4,3,0)\right]=8, \quad \alpha\left[J_{ \pm}(5,3,0)\right]=20, \quad \alpha\left[J_{ \pm}(6,3,0)\right]=32 \\
\alpha\left[J_{ \pm}(7,3,0)\right]=\alpha\left[J_{ \pm}(8,3,0)\right]=\alpha\left[J_{ \pm}(9,3,0)\right]=56 .
\end{gathered}
$$

Proposition 5.5.2. Let $d \geq 9$. Then

$$
\alpha\left[J_{ \pm}(d, 3,0)\right]=2\binom{d-1}{2}
$$

Proof. The example is inherited from Theorem 5.1.9.
Let us proceed with the upper bound. For the case $d=9$ the computer calculations give us the desired result. Let us repeat the proof of Theorem 5.1.9, updating it for small values of $d$. Let $I$ be a maximal independent set in $G:=J_{ \pm}(d, 3,0)$.

Clearly supports of vertices of $I$ form a 3 -uniform intersecting family. Theorem 5.2 .1 states that an intersecting family either contains at most 27 sets or has a 2 -transversal. It is known [50] that the constant 27 can be refined to 10 .

In the first case the family of supports has no 2-transversal. Then $|I| \leq 8 \cdot 10$, which is enough for $d>10$. Assume the contrary to the statement in the case $d=10$, id est $|I|>72$. It implies that vertices in $I$ have exactly 10 different supports. Suppose that every pair of supports splits exactly one vertex. Then by Theorem 5.2.2 all the supports have one common vertex, so at least $1+2 \cdot 10>10$ coordinates are required. Thus there are supports $f_{1}, f_{2}$ such that $\left|f_{1} \cap f_{2}\right|=2$. The initial graph $G$ has 16 vertices with supports $f_{1}$ and $f_{2}$; by the equality $\alpha\left[J_{ \pm}(4,3,0)\right]=8, I$ has at least 8 missing vertices on these supports. This refines the bound $|I| \leq 80$ to the desired $|I| \leq 72=2\binom{9}{2}$.

In the second case we have a one-point transversal set, say $U=\{u\}$. Let $I_{\text {sign }}$ be a set of vertices from $I$ containing $u^{\text {sign }}$, where sign $\in\{+,-\}$. Clearly $|I|=\left|I_{+}\right|+\left|I_{-}\right|$. After removing coordinate $u$ from every vertex, $I_{+}$becomes an independent set in $J_{ \pm}(d-1,2,-1)$. By Subsection 5.5.1 $\left|I_{+}\right| \leq\binom{ d-1}{2}$. The same bound for $I_{-}$finishes the proof in this case.

In the last case we have a transversal set of size 2 , say $\{a, b\}$. Let $I_{a}$ be the set of vertices of $I$ containing $a$ and not containing $b, I_{b}$ is defined analogously. Both $I_{a}$ and $I_{b}$ are nonempty, otherwise there is a one-point transversal set which is the previous case. Define $I_{a b}=I \backslash I_{a} \backslash I_{b}$. Computer calculations show that for $d=10$ we have at most 48 vertices in an independent set with such conditions.

Let $d$ be greater than 10 ; for every set $A \subset[d]$, such that $|A|=10$ and $a, b \in A$, we have $\alpha(G[A]) \leq$ 48 (here $G[A]$ stands for the subgraph of $G$ containing all the vertices $v$ such that supp $v \subset A$ ). Define $I[A]$ as the set of vertices $i$ from $I$ such that $\operatorname{supp} i \subset A$; note that $I[A]$ is an independent set. Every vertex from $I_{a b}$ belongs to $\binom{d-3}{7}$ different $A$, every vertex from $I_{a} \cup I_{b}$ belongs to $\binom{d-4}{6}$ different $A$. Summing up inequalities $|I[A]| \leq \alpha(G[A]) \leq 48$ over all choices of $A$ we got

$$
\binom{d-3}{7}\left|I_{a b}\right|+\binom{d-4}{6}\left(\left|I_{a}\right|+\left|I_{b}\right|\right) \leq 48\binom{d-2}{8}
$$

which is equivalent to

$$
\frac{d-3}{7}\left|I_{a b}\right|+\left(\left|I_{a}\right|+\left|I_{b}\right|\right) \leq \frac{48}{56}(d-2)(d-3) .
$$

Finally,

$$
|I|=\left|I_{a b}\right|+\left|I_{a}\right|+\left|I_{b}\right| \leq \frac{d-3}{7}\left|I_{a b}\right|+\left(\left|I_{a}\right|+\left|I_{b}\right|\right) \leq \frac{48}{56}(d-2)(d-3)<2\binom{d-1}{2} .
$$

### 5.5.4 The case $k=3, t=-2$

Example 5.3.2 gives us a lower bound $\alpha\left[J_{ \pm}(d, 3,-2)\right] \geq 2\binom{d}{2}+2$. Note that the Katona averaging method does not give an exact result because of the additional term of a smaller order of growth.

First, note that Theorem 5.1.6 in this case gives the bound

$$
\alpha\left[J_{ \pm}(d, 3,-2)\right] \leq 2\binom{d}{3}+8\binom{d}{2}
$$

Indeed, let $I$ be an independent set in $J_{ \pm}(d, 3,-2)$. We call a vertex $v \in I$ bad if there is another vertex with the same support which differs in exactly two places. Otherwise we call a vertex good. From Theorem 5.2.3 there are at most $2\binom{d}{3}$ good vertices.

Let us show that the number of bad vertices is at most $8\binom{d}{2}$. Indeed, each bad vertex has a pair of signplaces such that antipodal pair of signplaces contained in another vertex. But then all vertices containing one of these two pairs of signplaces must have the same third place therefore there are at most $8\binom{d}{2}$ bad vertices.

Using more accurate double counting we can prove the following upper bound.
Proposition 5.5.3. For $d \geq 6$ we have

$$
\alpha\left[J_{ \pm}(d, 3,-2)\right] \leq 2\binom{d}{3}+\frac{8}{3}\binom{d}{2} .
$$

Proof. A pair of vertices $v, w \in I$ is called tangled if these vertices have the same support and differ exactly at two places. Define the weight $c_{I}(v, i, j)$, where $v \in I$ and $i, j \in v$, in the following way:

$$
c_{I}(v, i, j)=\left\{\begin{array}{l}
1, \text { if } v \text { does not have tangled vertices in } G \\
2, \text { if } v \text { has a tangled vertex in } G \text { which differs at places } i, j, \\
0.5, \text { otherwise. }
\end{array}\right.
$$

Note that for a vertex $v$ sum of corresponding weights is at least 3 . Let $d_{i, j}$ be the sum of weight of vertices containing places $i$ and $j$ and let us estimate an upper bound for $d_{i, j}$. Then there are three cases which depend on whether there are tangled vertices containing places $i, j$ and whether these vertices have antipodal signs on places $i, j$.

In the first case there are no tangled vertices in $I$ which differ in places $i, j$. Then for any place $l$ the total weight of vertices with support $\{i, j, l\}$ is at most 2 . Then $d_{i, j} \leq 2(d-2)$. In the second case there are tangled vertices in $I$ which contain all four pairs of signplaces on places $i, j$. Then there are at most 8 vertices containing these places and $d_{i, j} \leq 16$.

In the last case there are two vertices in $I$ which are antipodal on places $i, j$ and there are no vertices in $I$ which contain one of the pairs of signplaces on places $i, j$. Then there are at most 4 vertices which differ in places $i, j$ and their total weight is at most 8 . The rest of vertices containing places $i, j$ have the same signs on these places therefore their total weight is at most $2(d-2)$.

Therefore, $d_{i, j} \leq 2 n+4$ and

$$
3|I| \leq \sum_{1 \leq i<j \leq n} d_{i, j} \leq\binom{ d}{2}(2 n+4)=6\binom{d}{3}+8\binom{d}{2} .
$$

### 5.6 Open questions

It seems very challenging to find a general method providing the independence number of $J_{ \pm}(d, k, t)$. Here we discuss questions that seem for us both interesting and relatively easy.

Small values of the parameters. The smallest interesting case is $J_{ \pm}(d, 3,-2)$. We hope that for $d>d_{0}$ Example 5.3.2 is the best possible, i.e.

$$
\alpha\left[J_{ \pm}(d, 3,-2)\right]=\alpha\left[K_{ \pm}(d, 3,-2)\right]=2\binom{d}{3}+2
$$

Recall that the last equality is established by Theorem 5.1.5.
Another small case leads to the following conjecture.
Conjecture 5.6.1. Let $d>d_{0}$ be an even number. Then

$$
\alpha\left[J_{ \pm}(d, 4,1)\right]=2 n(d-2)
$$

Obviously $\alpha\left[J_{ \pm}(d, 4,1)\right] \geq \alpha\left[J_{ \pm}(d, 4, o d d)\right]=2 n(d-2)$ (see Proposition 5.1.2).
Chromatic numbers. Usually finding or evaluating the chromatic number is a more complicated problem than finding or evaluating the independence number. In particular Lovász [87] proved Kneser's conjecture on the chromatic number of $K(d, k, 0) 17$ year after Erdős, Ko and Rado determined the independence number of this graph.

In the setting of this chapter we have

$$
c(k, t) n \leq \frac{\left|V\left[J_{ \pm}(d, k, t)\right]\right|}{\alpha\left[J_{ \pm}(d, k, t)\right]} \leq \chi\left[J_{ \pm}(d, k, t)\right] \leq \frac{\left|V\left[J_{ \pm}(d, k, t)\right]\right|}{\alpha\left[J_{ \pm}(d, k, t)\right]} \log \left|V\left[J_{ \pm}(d, k, t)\right]\right| \leq C(k, t) n \log n
$$

for some positive constants $c(k, t), C(k, t)$. The second inequality holds since $J_{ \pm}(d, k, t)$ is a vertextransitive graph (see [86]).

Recall that Theorem 5.1.12 shows that for $k=3, t=-2$ the chromatic number of a Johnson-type graph may not coincide with simple general bounds.

Difference between $J_{ \pm}(d, k, t)$ and $K_{ \pm}(d, k, t)$. It turns out that for a negative odd $t$ Theorems 5.1.5 and 5.1.8 give

$$
\alpha\left[J_{ \pm}(d, k, t)\right]=\alpha\left[K_{ \pm}(d, k, t)\right] .
$$

Does it hold for all negative $t$ ? Do we have

$$
\chi\left[J_{ \pm}(d, k, t)\right]=\chi\left[K_{ \pm}(d, k, t)\right]
$$

in this case?
The general comparison of the behavior of independence numbers and chromatic numbers of these graphs is also of interest.

## Chapter 6

## Chromatic numbers of 2-dimensional spheres

In 1976 Simmons conjectured that every coloring of a 2-dimensional sphere of radius strictly greater than $1 / 2$ in three colors has a pair of monochromatic points at distance 1 apart. Paper [21] proves this conjecture and we repeat the proof here.

### 6.1 Introduction

A coloring of a given set $M$ is a map from $M$ to the set of colors. A coloring of a subset $M$ of a metric space is proper if no pair of monochromatic points lie at distance 1 apart. The minimum number of colors that admits a proper coloring of $M$ is called the chromatic number of $M$; we denote it by $\chi(M)$. In the case of $M \subset \mathbb{R}^{d}$, the distance typically comes from the induced Euclidean metric on $M$.

A slightly different point of view is to consider a unit distance graph $G(M)$ : the points of $M$ are the vertices of $G(M)$ and edges connect points at unit distance apart. By definition, $\chi(M)=\chi(G(M))$. The de Bruijn-Erdős theorem states that if $\chi(M)$ is finite then there is a finite subgraph $H$ of $G(M)$ such that $\chi(H)=\chi(G(M))$.

Denote by $S^{2}(r)$ the two-dimensional sphere of radius $r$ in $\mathbb{R}^{3}$ centered at the origin. Let $\chi\left(S^{2}(r)\right)$ be the chromatic number of $S^{2}(r)$ with respect to the Euclidean metric. Obviously, if $r<1 / 2$ and $r=1 / 2$ then the chromatic number is equal to 1 and 2 , respectively. Note that for any $r>\frac{1}{2}$ there is $r_{1}<r$ such that $S^{1}\left(r_{1}\right)$ contains an odd cycle. Since $S^{1}\left(r_{1}\right) \subset S^{2}(r)$, we obtain that $\chi\left(S^{2}(r)\right) \geq 3$. G. Simmons [116] proved that

$$
\chi\left(S^{2}(r)\right) \geq 4 \quad \text { for } \quad r \geq \frac{\sqrt{3}}{3}
$$

In the proof, Simmons constructs certain subgraphs of $G\left(S^{2}(r)\right)$ that contain triangles. Obviously, for smaller values of the radius $G\left(S^{2}(r)\right)$ is triangle-free, and so other ideas are needed.

Then L. Lovász 88 generalized the odd cycle construction to an arbitrary dimension, showing that for every $d \geq 3$ there exists a family of strongly self-dual polytopes inscribed in $S^{d-1}(r)$ whose graphs of diameters have chromatic number $d+1$ and that $r$ can be arbitrarily close to $\frac{1}{2}$. In our notation this result can be formulated as follows:

Theorem 6.1.1 (Lovász, 88]). For every $d \geq 2$ there exists a monotonically decreasing sequence $r_{k}^{(d)}, k=1,2, \ldots$, such that

$$
\lim _{k \rightarrow \infty} r_{k}^{(d)}=\frac{1}{2} \quad \text { and } \quad \chi\left(S^{d-1}\left(r_{k}^{(d)}\right)\right) \geq n+1
$$

Since $S^{d-1}\left(r_{1}\right) \subset S^{n}(r)$ for $r_{1} \leq r$, we get the following inequality.

## Corollary 6.1.1.

$$
\chi\left(S^{d-1}(r)\right) \geq n \quad \text { for } \quad r>\frac{1}{2}
$$

Some sources state that the chromatic number of a two-dimensional sphere $S^{2}(r)$ is known only for $r \leq \frac{1}{2}$ and for $r=\frac{\sqrt{2}}{2}$ [67, 90]. But it should be clarified that the equality $\chi\left(S^{2}(r)\right)=n+1=4$ is true for $r \in\left\{r_{k}^{(3)}\right\} \cap\left(\frac{1}{2}, \frac{\sqrt{3-\sqrt{3}}}{2}\right]$. Explicit formulas for algebraic numbers $r_{k}^{(3)}$, if such exist, seem to be too complicated, but it is not difficult to compute $r_{k}^{(3)}$ for a given $k$ with an arbitrary precision by approximately solving a certain optimization problem. For example, the first non-trivial construction in the case of a two-dimensional sphere corresponds to a unit distance embedding of the Grötzsch graph at $r=0.54003829 \ldots$

It is worth noting that chromatic numbers in high dimensions were studied using algebraic, topological and combinatorial methods. A.M. Raigorodskii [109] showed that for every fixed $r>1 / 2$ the chromatic number of an $d$-dimensional sphere grows exponentially with $d$. O. Kostina [76] refined asymptotic lower bounds. R. Prosanov [105] gave a new asymptotic upper bound. The paper of A. Kupavskii [79] contains several results on the number of different colors on a sphere of given radius in every proper coloring of $\mathbb{R}^{d}$.

A lot of results on colorings of 2-dimensional spheres were obtained by Simmons [116]. Recent discovery of a 5-chromatic unit distance subgraph of the Euclidean plane 31] spurred interest in the topic and in particular to the chromatic number of a 2-dimensional sphere.

Among the other results, Voronov, Neopryatnaya, and Dergachev [122] constructed several 5chromatic subgraphs of 2-dimensional spheres, which lead to the bounds

$$
\begin{aligned}
& \chi\left(S^{2}\left(r_{1}\right)\right) \geq 5 \quad \text { where } \quad r_{1}=\cos \frac{3 \pi}{10}=\frac{\sqrt{5-\sqrt{5}}}{2 \sqrt{2}}=0.58778 \ldots \\
& \chi\left(S^{2}\left(r_{2}\right)\right) \geq 5 \quad \text { where } \quad r_{2}=\cos \frac{\pi}{10}=\frac{\sqrt{5+\sqrt{5}}}{2 \sqrt{2}}=0.95105 \ldots
\end{aligned}
$$

The paper [117] contains a family of proper colorings of $S^{2}(r)$ spheres in 7 colors, provided $r$ is large enough.

The following statement was formulated by Simmons as a conjecture [116]. The proof of Simmons' conjecture is the main result of the chapter.

Theorem 6.1.2 (Cherkashin-Voronov [21]). For every $r>\frac{1}{2}$ we have

$$
\chi\left(S^{2}(r)\right) \geq 4
$$

We note that for $\frac{1}{2}<r \leq \frac{\sqrt{3-\sqrt{3}}}{2}=0.563 \ldots$ a proper 4-coloring of $S^{2}(r)$ can be obtained from a partition of the sphere into four equal spherical triangles [116]. It implies the following corollary.

Corollary 6.1.2. $\chi\left(S^{2}(r)\right)=4$ for $\frac{1}{2}<r \leq \frac{\sqrt{3-\sqrt{3}}}{2}=0.563 \ldots$.

### 6.2 Proof of Theorem 6.1.2

Recall that for $r \geq \frac{\sqrt{3}}{3}$ the statement was proved in [116].
Here is the sketch of the proof. Fix $r \in\left(\frac{1}{2}, \frac{\sqrt{3}}{3}\right)$. Suppose that there is a proper 3-coloring of the sphere $S^{2}(r)$. Further arguments consist of two steps. In the first step we use the BorsukUlam theorem to show that every color is dense in the sphere. Consider a graph $G_{k}$ with vertices $x_{1}, \ldots x_{2 k+1}, y_{1}, \ldots, y_{2 k+1}$ and edges $\left\{\left(y_{i}, y_{i+1}\right),\left(x_{i}, y_{i}\right): 1 \leq i \leq 2 k+1\right\}$ (where indices are modulo $2 k+1$ ), i.e. an odd cycle with attached pendant vertices. We provide an explicit representation of $G_{k}$ as a unit distance subgraph of the sphere. The second step is to show that this embedding is stable under small perturbations of $x_{i}$. Then one can move every $x_{i}$ at a red point, which forces the odd cycle on vertices $y_{i}$ to be colored in the remaining two colors. The contradiction proves the theorem.

Note that the idea of attaching an odd cycle to a finite set $A$ in order to exclude the possibility of $A$ to be monochromatic was used in a series of papers devoted to the existence of planar unit distance graphs with chromatic number 4 and arbitrarily large girth [61, [118, 126]. The key twist in step 2 is to find the required embedding of $G_{k}$ implicitly, i.e. the corresponding $A$ is not a constructive set. Similar ideas were used in 69.

### 6.2.1 Step 1. Each color is a dense set

All the distances are considered in the metrics induced from Euclidean space $\mathbb{R}^{3}$, the distance between $x$ and $y$ is denoted by $\|x-y\|$.

Fix $r \in\left(\frac{1}{2}, \frac{\sqrt{3}}{3}\right)$ and consider $S^{2}(r)$. Suppose that there is a proper coloring of $S^{2}(r)$ in three colors. Consider the unit distance graph $G=G\left(S^{2}(r)\right)$. Then the neighborhood of a vertex $x$ in graph $G$ forms a circle of diameter $d=\frac{\sqrt{4 r^{2}-1}}{r}$ in the sphere. It is worth noting that every circle in the sphere has two centers at a pair of antipodal points and hence it has two radii; a circle of diameter $d=\frac{\sqrt{4 r^{2}-1}}{r}$ has radii 1 and $\rho=\sqrt{4 r^{2}-1}$ in the induced metric. Since $r<\frac{\sqrt{3}}{3}$, the smaller radius is $\rho$, so we refer to $\rho$ as the radius and $-x$ as the center of the circle. Vice versa, any circle of radius $\rho$ is a graph-neighborhood of some vertex of $G$, and hence contains points of at most two colors.

We need the following technical statement.
Lemma 6.2.1. Let $D \subset S^{2}(r) \times S^{2}(r)$ be a set of pairs $(x, y)$ such that $0<\|x-y\|<d$. Then

- for every $(x, y) \in D$ there are two circles of radius $\rho$ containing $x$ and $y$. One may denote their centers by $c_{r}$ and $c_{l}$ in such a way that the triple of radius-vectors $\left(x, y, c_{r}\right)$ is right-handed and the triple $\left(x, y, c_{l}\right)$ is left-handed.
- The functions $c_{r}(x, y)$ and $c_{l}(x, y)$ from $D$ to $S^{2}(r)$ are continuous.

In what follows, we will call a circle passing through the points $x, y$ with center $c$ right-handed if the triple $(x, y, c)$ is right-handed, and left-handed otherwise.

Let $C_{\text {red }}, C_{\text {blue }}, C_{\text {green }}$ be the sets of red, blue and green points, respectively. A chromaticity of a point $x$ is the number of sets $\overline{C_{r e d}}, \overline{C_{\text {blue }}}, \overline{C_{\text {green }}}$ containing $x$ (as usual, $\bar{T}$ stands for the closure of a set $T$ ). A set $T \subset S^{2}(r)$ is called dense if $\bar{T}=S^{2}(r)$. Let $B_{\rho}(x)$ denote the set of points $y \in S^{2}(r)$ such that $\|x-y\|<\rho$, i.e. an open ball of radius $\rho$ and diameter $d$.

Lemma 6.2.2. If some open ball of diameter d contains points of all three colors then each of $C_{r e d}$, $C_{\text {blue }}, C_{\text {green }}$ is dense in the sphere.

Proof. Consider points $x \in C_{r e d}, y \in C_{\text {blue }}$ and $z \in C_{\text {green }}$ inside a ball $K_{0}$ of diameter $d$. Then one can continuously move $K_{0}$ to a ball $K$ containing two points (say, $x$ and $y$ ) on the boundary; at the first such moment the point $z$ lies inside $K$. The circle $\partial K$ contains blue and red points, and so it is colored in blue and red only. Hence, it contains a point $u$ lying in the closures of $C_{\text {red }}$ and $C_{\text {blue }}$; without loss of generality, assume that point $u$ is red. A red-green circle (right-handed, guaranteed to exist by Lemma 6.2.1) of diameter $d$ containing $z$ and $u$ and a blue-green circle (left-handed) with the diameter $d$ containing $z$ and blue point $u^{\prime}$ in a small neighborhood of $u$ intersect in a green point $v$. Note that if $u=u^{\prime}$ then $v=u=u^{\prime}$. Hence, due to the continuity of circles in Lemma 6.2.1, $v$ may be arbitrarily close to $u$ with a proper choice of $u^{\prime}$ (see Fig. 6.1). It implies that the chromaticity of $u$ is three.


Figure 6.1: Finding a point with chromaticity 3 in Lemma 6.2.2
Since $u$ has chromaticity 3 , a small neighborhood of $u$ contains a point $a \neq u$ with the chromaticity at least 2. Suppose that $a$ has chromaticity 2 (say, $a$ does not lie in $\overline{C_{\text {green }}}$ ) and $\|a-u\|<d$. Consider a green point $b$ in a small neighborhood of $u$. Consider a red point $e$ and a blue point $f$ in a small neighborhood of $a$. Then the right-handed circle containing $b$ and $e$ is red-green and the left-handed circle containing $b$ and $f$ is blue-green, so they intersect in a green point $g$. Since the neighborhoods can be chosen arbitrarily small, $g$ can be arbitrarily close to $a$. Hence, $a$ has chromaticity 3 , a contradiction.

Thus, we have shown that if a point with the chromaticity 3 and a point with the chromaticity at least 2 lie at a distance smaller than $d$, then they both have chromaticity 3 .


Figure 6.2: Propagation of 3-chromaticity along a circle in Lemma 6.2.2

Now let $x_{1}$ and $x_{2}$ be points of chromaticity 3 such that $\left\|x_{1}-x_{2}\right\|<d$. We claim that any point on a circle $L$ of diameter $d$ containing $x_{1}$ and $x_{2}$ has chromaticity three. By the previous argument,
it is enough to show that the chromaticity is at least 2 . Without loss of generality, a triple $\left(x_{1}, x_{2}, c\right)$ is left-handed, where $c$ is the center of $L$ on the sphere. Arguing indirectly, assume that a point $y_{1} \in L$ has a small red neighborhood $U_{y_{1}}$. Choose a blue point $u_{1}$ in a small neighborhood of $x_{1}$ and a green point $v_{1}$ in a small neighborhood of $x_{2}$ (see Fig. 6.2. By Lemma 6.2.1 the left-handed circle of diameter $d$ passing through blue point $u_{1}$, green point $v_{1}$ is close to $L$ so it intersects red set $U_{y_{1}}$; this contradiction shows that every point on $L$ has chromaticity 3 .

Let $q$ be an arbitrary point of $S^{2}(r)$. Consider a path $q_{0}, q_{1} \ldots q_{t}=q$ such that $q_{0} \in L$ and $\left\|q_{i+1}-q_{i}\right\|<\rho$ for $0 \leq i \leq t-1$. A circle $L_{1}$ of diameter $d$ that passes through $q_{1}$ and $q_{0}$ intersects $L$ in two points, so by the previous argument every point (in particular, $q_{1}$ ) of $L_{1}$ has chromaticity 3 . By induction, a circle $L_{i+1}$ of diameter $d$ that passes through $q_{i+1}$ and $q_{i}$ intersects $L_{i}$ in two points, so every point in $L_{i+1}$ (in particular $q_{i+1}$ ) has chromaticity 3 . So $q=q_{t}$ also has chromaticity 3 . Since $q \in S^{2}(r)$ was arbitrary, every point of $S^{2}(r)$ has chromaticity 3.

Suppose that the condition of Lemma 6.2 .2 does not hold, i.e.

$$
\text { every open ball of diameter } d \text { contains points of at most two colors. }
$$

Consider a continuous function

$$
f: S^{2}(r) \rightarrow \mathbb{R}^{2}, \quad f(x)=\left(\operatorname{dist}\left(x, \overline{C_{\text {red }}}\right), \operatorname{dist}\left(x, \overline{C_{\text {blue }}}\right)\right)
$$

where dist $(\cdot)$ stands for the distance between a point and a set in $\mathbb{R}^{3}$. By the Borsuk-Ulam theorem there exists $x^{*} \in S^{2}(r)$ such that $f\left(x^{*}\right)=f\left(-x^{*}\right)$. We have to deal with three cases.


Figure 6.3: Case 1

Case 1: $f\left(x^{*}\right)=(0,0)$. Without loss of generality, the point $x^{*}$ is blue. One may pick a red point $z$, which is arbitrarily close to $x^{*}$. If $\left\|x^{*}-z\right\|<\rho$, then the intersection of circles of unit Euclidean radius with centers $x^{*}$ and $z$ consists of two green points $y_{1}, y_{2}$ belonging to the circle of radius $\rho$ centered at $-x^{*}$. Hence, one can cover a small neighborhood of $-x^{*}$ and $y_{1}$ by a ball of diameter $d$. Every neighborhood of $-x^{*}$ contains red and blue points; point $y_{1}$ is green (see Fig. 6.3). We have a contradiction with assumption ( $\star$ ).

Case 2: $f\left(x^{*}\right)=(a, b), a, b>0$. Then both points $x^{*},-x^{*}$ are green. We may swap blue and green colors to reduce the situation to the next case with the same $x^{*}$.

Case 3: $f\left(x^{*}\right)=(a, 0), a>0$. We claim that $a>\rho$. Assume the contrary, i.e. $x^{*} \in \overline{C_{b l u e}}$ and for every $\eta>0$ there is a red point $z=z_{\eta}$ such that $\left\|x^{*}-z\right\| \leq \rho+\eta$. Note that if $x^{*}$ is green, then it contradicts $(\star)$, so $x^{*}$ is blue. There are distinct points $y_{1}, y_{2} \in \overline{B_{\rho}\left(-x^{*}\right)}$ such that $\left\|x^{*}-y_{1}\right\|=\left\|x^{*}-y_{2}\right\|=\left\|z-y_{1}\right\|=\left\|z-y_{2}\right\|=1$. Since $x^{*}$ is blue and $z$ is red $y_{1}, y_{2} \in \overline{C_{\text {green }}}$. Recall that $f\left(-x^{*}\right)=f\left(x^{*}\right)$, so there is a point $z^{\prime} \in \overline{C_{r e d}} \cap \overline{B_{\rho}\left(-x^{*}\right)}$. Let $y^{\prime} \in\left\{y_{1}, y_{2}\right\}$ be such that $z^{\prime},-x^{*}$ and $y^{\prime}$ do not lie on a great circle of $S^{2}(r)$. Then for a small enough $\eta$ the neighborhoods of $-x^{*}, y^{\prime}$ and $z^{\prime}$ can be covered by a ball of diameter $d$. This is a contradiction with ( $\star$ ).

So the set $\overline{B_{\rho}\left(x^{*}\right)} \cup \overline{B_{\rho}\left(-x^{*}\right)}$ is colored with blue and green.
Lemma 6.2.3. The bipartite subgraph of $S^{2}(r)$ with parts $\overline{B_{\rho}\left(x^{*}\right)}$ and $\overline{B_{\rho}\left(-x^{*}\right)}$ is connected.
Proof. Any point $x \in \overline{B_{\rho}\left(x^{*}\right)}$ has a common neighbor with $x^{*}$ since the corresponding unit circles intersect. So $\overline{B_{\rho}\left(x^{*}\right)}$ belong to the same connected component; the same holds for $\overline{B_{\rho}\left(-x^{*}\right)}$. There is an edge between $\overline{B_{\rho}\left(x^{*}\right)}$ and $\overline{B_{\rho}\left(-x^{*}\right)}$, and so the subgraph is connected.

By Lemma 6.2.3, one can color $\overline{B_{\rho}\left(x^{*}\right)} \cup \overline{B_{\rho}\left(-x^{*}\right)}$ in two colors in the unique way (up to symmetry): the first part is blue and the second one is green. Then the distance from $x^{*}$ and $-x^{*}$ to $\overline{C_{b l u e}}$ is zero and nonzero simultaneously.

This contradiction implies that each color is dense in the sphere.

### 6.2.2 Step 2. Stability of embedding

In this section we will need the implicit function theorem [78] in the following weakened formulation.
Theorem 6.2.1. Let $F: \mathbb{R}^{2 s} \rightarrow \mathbb{R}^{s}$ be a continuously differentiable function,

$$
F=F(X, Y)=F\left(x_{1}, \ldots, x_{s} ; y_{1}, \ldots, y_{s}\right),
$$

and at some point $X=a, Y=b$ the following conditions are satisfied

$$
F(a, b)=0, \quad \operatorname{det}\left(\frac{\partial F(X, Y)}{\partial Y}\right)_{X=a, Y=b} \neq 0
$$

Then there exists $\eta>0$ such that the system of equations $F(X, Y)=0$ is solvable in $Y$ for any $X$ satisfying the condition $\|X-a\|<\eta$.

Recall that $G_{k}$ is an odd cycle of length $m=2 k+1$ with an extra pendant (leaf) vertex attached to each vertex of the cycle. In particular, $G_{k}$ has $2 m$ vertices and $2 m$ edges.

Denote by $y_{1}, \ldots, y_{m}$ the points of $S^{2}(r)$ that correspond to the cycle vertices and by $x_{1}, \ldots, x_{m}$ the points of $S^{2}(r)$ that correspond to the pendant vertices. For convenience, let us put $X=\left(x_{1}, \ldots, x_{m}\right)$ and $Y=\left(y_{1}, \ldots, y_{m}\right)$ the vectors of dimension $s=3 m$ containing all coordinates. Then the embedding of $G_{k}$ can be given by the pair $(X, Y)$.

Lemma 6.2.4. Fix the radius $r \in\left(\frac{1}{2}, \frac{\sqrt{3}}{3}\right)$. Then if $k$ is large enough, there exists a unit distance embedding $(X, Y)$ of $G_{k}$ into $S^{2}(r)$ and a constant $\eta>0$ such that for any $\tilde{X}$ satisfying $\|\tilde{X}-X\|<\eta$ there exists $Y$ such that $(\tilde{X}, \tilde{Y})$ is a "perturbed" unit distance embedding of $G_{k}$.

In other words, for any sufficiently small perturbation of pendant vertices, it is possible to find the embedding of the cycle vertices.

Proof. We provide the desired unit distance embedding explicitly. In what follows we slightly abuse the notation and write $x_{i}$ and $y_{i}$ for a vertex of the graph, the corresponding point on $S^{2}(r)$, and its 3 -dimensional vector representation. Consider the system of equations defining the embedding $G_{k}$ in $S^{2}(r)$ :

$$
\left\{\begin{array}{l}
f_{i}=\left\|y_{i}\right\|^{2}-r^{2}=0, \quad 1 \leq i \leq m ;  \tag{6.1}\\
f_{i+m}=\left\|y_{i}-y_{i+1}\right\|^{2}-1=0, \quad 1 \leq i \leq m-1 ; \\
f_{2 m}=\left\|y_{m}-y_{1}\right\|^{2}-1=0 ; \\
f_{i+2 m}=\left\|x_{i}-y_{i}\right\|^{2}-1=0, \quad 1 \leq i \leq m
\end{array}\right.
$$

Next, we will be interested in the family of embeddings, the $k=2$ case of which is depicted on Fig. 6.4.


Figure 6.4: Unit distance embedding of $G_{k}$, the $k=2$ case
Note that (6.1) allows $x_{i}$ to lie in $\mathbb{R}^{3}$, not only $S^{2}(r)$, but the cycle $y_{1}, \ldots, y_{m}$ must lie on the sphere.

One can consider the function corresponding to the left-hand side of the system 6.1).

$$
F=\left(f_{1}, \ldots, f_{3 m}\right)=F\left(x_{11}, x_{12}, x_{13}, \ldots, x_{m 3} ; y_{11}, \ldots, y_{m 3}\right)
$$

Suppose that the Jacobian matrix $J=\left(\frac{\partial F}{\partial Y}\right)$ is nondegenerate,

$$
\operatorname{det} J=\operatorname{det}\left(\frac{\partial F}{\partial Y}\right) \neq 0
$$

then the statement of the lemma follows from Theorem 6.2.1. The rest of the proof is devoted to the calculation of this determinant.

The matrix $J$ has the following form (recall that $x_{i}$ and $y_{i}$ are $1 \times 3$ vectors):

$$
J(X, Y)=2\left(\begin{array}{cccccc}
y_{1} & 0 & 0 & 0 & \ldots & 0 \\
0 & y_{2} & 0 & 0 & \ldots & 0 \\
0 & 0 & y_{3} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & y_{m} \\
y_{1}-y_{2} & y_{2}-y_{1} & \ldots & 0 & \ldots & 0 \\
0 & y_{2}-y_{3} & y_{3}-y_{2} & 0 & \ldots & 0 \\
\vdots & & & & & \\
y_{1}-y_{m} & 0 & \ldots & \ldots & 0 & y_{m}-y_{1} \\
y_{1}-x_{1} & 0 & \ldots & \ldots & 0 & 0 \\
0 & y_{2}-x_{2} & 0 & \ldots & 0 & 0 \\
\vdots & & \vdots & & \vdots & \\
0 & \ldots & 0 & \ldots & 0 & y_{m}-x_{m}
\end{array}\right) .
$$

Subtracting some rows from each other, we get

$$
\operatorname{det} J=2^{3 m} \operatorname{det}\left(\begin{array}{cccccc}
y_{1} & 0 & 0 & 0 & \ldots & 0 \\
0 & y_{2} & 0 & 0 & \ldots & 0 \\
0 & 0 & y_{3} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & y_{m} \\
y_{2} & y_{1} & \ldots & 0 & \ldots & 0 \\
0 & y_{3} & y_{2} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
y_{m} & \ldots & 0 & \ldots & 0 & y_{1} \\
x_{1} & 0 & 0 & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & x_{m}
\end{array}\right)=(-1)^{s} 2^{3 m} \operatorname{det}\left(\begin{array}{cccccc}
y_{1} & 0 & 0 & 0 & \ldots & 0 \\
x_{1} & 0 & 0 & 0 & \ldots & 0 \\
y_{2} & y_{1} & 0 & \ldots & \ldots & 0 \\
0 & y_{2} & 0 & 0 & \ldots & 0 \\
0 & x_{2} & 0 & 0 & \ldots & 0 \\
0 & y_{3} & y_{2} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & y_{m-1} \\
0 & 0 & 0 & 0 & \ldots & y_{m} \\
0 & 0 & 0 & 0 & \ldots & x_{m} \\
y_{m} & 0 & 0 & \ldots & 0 & y_{1}
\end{array}\right) \text {, }
$$

where the term $(-1)^{s}, s \in\{0,1\}$ is responsible for the parity of the permutation of the rows. Since we are not interested in the sign of the determinant, there is no point in evaluating the parity.

Then, expanding the determinant by the last row and rearranging the rows, we obtain $\operatorname{det} J=$

$$
=(-1)^{s} 2^{3 m} \operatorname{det}\left(\begin{array}{cccccc}
y_{1} & 0 & 0 & 0 & \ldots & 0 \\
x_{1} & 0 & 0 & 0 & \ldots & 0 \\
y_{2} & y_{1} & 0 & \ldots & \ldots & 0 \\
0 & y_{2} & 0 & 0 & \ldots & 0 \\
0 & x_{2} & 0 & 0 & \ldots & 0 \\
0 & y_{3} & y_{2} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & y_{m-1} \\
0 & 0 & 0 & 0 & \ldots & y_{m} \\
0 & 0 & 0 & 0 & \ldots & x_{m} \\
0 & 0 & 0 & \ldots & 0 & y_{1}
\end{array}\right)+(-1)^{s} 2^{3 m} \operatorname{det}\left(\begin{array}{cccccc}
y_{m} & 0 & 0 & \ldots & 0 & 0 \\
y_{1} & 0 & 0 & 0 & \ldots & 0 \\
x_{1} & 0 & 0 & 0 & \ldots & 0 \\
y_{2} & y_{1} & 0 & \ldots & \ldots & 0 \\
0 & y_{2} & 0 & 0 & \ldots & 0 \\
0 & x_{2} & 0 & 0 & \ldots & 0 \\
0 & y_{3} & y_{2} & 0 & \ldots & 0 \\
\vdots & & & \ddots & & \\
0 & 0 & 0 & 0 & \ldots & y_{m-1} \\
0 & 0 & 0 & 0 & \ldots & y_{m} \\
0 & 0 & 0 & 0 & \ldots & x_{m}
\end{array}\right)=
$$

$$
=(-1)^{s} 2^{3 m}\left(V_{1} \ldots V_{m}+V_{1}^{\prime} \ldots V_{m}^{\prime}\right)
$$

where

$$
V_{i}=\operatorname{det}\left(\begin{array}{c}
y_{i} \\
x_{i} \\
y_{i+1}
\end{array}\right)=-\operatorname{det}\left(\begin{array}{c}
y_{i} \\
y_{i+1} \\
x_{i}
\end{array}\right), \quad V_{i}^{\prime}=\operatorname{det}\left(\begin{array}{c}
y_{i} \\
y_{i+1} \\
x_{i+1}
\end{array}\right)
$$

Here in the determinant calculations we used the fact that after decomposing the last row into a sum, each of the summands becomes block-triangular. In addition, note that the cyclic permutation of matrix rows does not change the sign of the determinant, since $m$ is odd.

Now we fix the following embedding (Fig. 6.4). Let vertices $y_{i}$ lie in the plane $z=h$ (and form a regular $m$-gon), and vertices $x_{i}$ lie in the plane $z=-h$ (and also form a regular $m$-gon). Note that the radius of the circumcircle of the $m$-gon is greater than $\frac{1}{2}$, and $r^{2}<\frac{1}{3}$, hence

$$
\begin{equation*}
h<\left(\frac{1}{3}-\frac{1}{4}\right)^{1 / 2}=\frac{1}{2 \sqrt{3}}<\frac{1}{2} \tag{6.2}
\end{equation*}
$$

Denote by $U_{m}$ the rotation matrix by an angle $2 \pi / m$ counterclockwise around $z$-axis. Then $y_{i+1}$ $=U_{m} y_{i}, x_{i+1}=U_{m} x_{i}$. Hence, all $V_{i}$ coincide and all $V_{i}^{\prime}$ also coincide; put $V=V_{i}$ and $V^{\prime}=V_{i}^{\prime}$. Hence

$$
\operatorname{det} J=(-1)^{s} 2^{3 m}\left(V^{m}+\left(V^{\prime}\right)^{m}\right)
$$

We claim that

$$
V+V^{\prime}=\operatorname{det}\left(\begin{array}{c}
y_{1} \\
y_{2} \\
x_{2}-x_{1}
\end{array}\right) \neq 0
$$

Indeed, since $y_{13}=y_{23}=h, x_{13}=x_{23}=-h$, the equality

$$
\alpha y_{1}+\beta y_{2}+\gamma\left(x_{2}-x_{1}\right)=0
$$

implies $\alpha=-\beta$, i.e.

$$
\begin{equation*}
\alpha\left(y_{1}-y_{2}\right)=\gamma\left(x_{1}-x_{2}\right) \tag{6.3}
\end{equation*}
$$

Recall that $\left\|y_{1}-y_{2}\right\|=\left\|x_{1}-x_{2}\right\|=1$, so $\alpha= \pm \gamma$.
Since both sets of points $\mathcal{X}=\left\{x_{1}, \ldots, x_{m}\right\}, \mathcal{Y}=\left\{y_{1}, \ldots, y_{m}\right\}$ form vertices of congruent regular $m$-gons, in the case $\alpha=\gamma$, we have $x_{1}-x_{2}=y_{1}-y_{2}$ and the projections of $x_{i}$ and $y_{i}$ on the plane $z=0$ coincide, $i=1,2, \ldots, m$, and taking into account (6.2), we have

$$
\left\|x_{1}-y_{1}\right\|=2 h<1
$$

In the case $\alpha=-\gamma$, we have $x_{1}-x_{2}=y_{2}-y_{1}$ and the sets $\mathcal{X}$ and $\mathcal{Y}$ are symmetric about the origin. Then $x_{1} x_{2} y_{1} y_{2}$ is a rectangle, and

$$
\left\|x_{1}-y_{1}\right\|^{2}>\left\|x_{1}-x_{2}\right\|^{2}+4 h^{2}>1
$$

In both cases we got a contradiction.
Then the equation (6.3) does not hold and so $V+V^{\prime} \neq 0$. Hence

$$
\operatorname{det} J=(-1)^{s} 2^{3 m}\left(V^{m}+\left(V^{\prime}\right)^{m}\right) \neq 0
$$

as required.

### 6.3 Open questions

Is the chromatic number of $S^{2}(r)$ "almost monotonically" increasing with $r$ ? Id est, is the chromatic number monotonic except for an at most countable set of values $r$ ? Recall that the known results (see Table 1) allow for such possibility.

| $r$ | Estimate for $\chi(r)=\chi\left(S^{2}(r)\right)$ | Source |
| :---: | :---: | :---: |
| $r<1 / 2$ | $\chi(r)=1$ |  |
| $r=1 / 2$ | $\chi(r)=2$ |  |
| $\frac{1}{2}<r \leq \frac{\sqrt{3-\sqrt{3}}}{2}$ | $\chi(r)=4$ | Corollary 1 |
| $r>\frac{\sqrt{3-\sqrt{3}}}{2}$ | $\chi(r) \geq 4$ | Theorem 2 |
| $r=\frac{\sqrt{5-\sqrt{5}}}{2 \sqrt{2}}$ | $\chi(r) \geq 5$ | $[122]$ |
| $r=\frac{1}{\sqrt{2}}$ | $\chi(r)=4$ | $[116, ~ 57]$ |
| $r=\frac{\sqrt{5+\sqrt{5}}}{2 \sqrt{2}}$ | $\chi(r) \geq 5$ | $[122]$ |
| $r \leq \frac{1}{\sqrt{3}}$ | $\chi(r) \leq 5$ | $[116,90]$ |
| $r \leq \sqrt{3} / 2$ | $\chi(r) \leq 6$ | $[90$ |
| $r \geq 12.44$ | $\chi(r) \leq 7$ | $[117]$ |
| $r>1 / 2$ | $\chi(r) \leq 15$ | $[28,[106]$ |

Table 6.1: Lower and upper estimates for $\chi\left(S^{2}(r)\right)$.

Is there a proper coloring of $S^{2}(r)$ in $\chi\left(S^{2}(r)\right)$ colors such that every color is dense? It is interesting that all known upper bounds are given by explicit colorings in which every color is a finite union of regions bounded by piecewise-continuous curves.

What is the minimum number of vertices in a subgraph $G$ of a sphere $S^{2}(r)$ with $\chi(G)=$ $\chi\left(S^{2}(r)\right)$ ? By the de Bruijn-Erdős theorem, this number is finite. Note that the proof of Theorem 2 does not give any finite 4 -chromatic unit distance graph.


Figure 6.5: 4-coloring of the sphere. Here $s_{0} \rightarrow 0$ as $r \rightarrow 1 / 2$
Let us focus on the case $r=1 / 2+\varepsilon, \varepsilon \rightarrow 0$. Then the sphere can be colored in 4 colors in the way shown in Figure 6.5. Let us denote by $s_{0}$ the area of the spherical cap of color 0. Observe that
$s_{0}=4 \pi \varepsilon+o(\varepsilon)$, and thus, via averaging, we have the lower bound $n_{4}(r) \geq c \varepsilon^{-1}$ for some $c>0$, where $n_{4}(r)$ is the minimum number of vertices in a 4 -chromatic unit distance graph. Can this obvious bound be refined?

## Chapter 7

## Chromatic numbers of 3-dimensional slices

We follow the paper [123] and prove that for an arbitrary $\varepsilon>0$ holds

$$
\chi\left(\mathbb{R}^{3} \times[0, \varepsilon]^{6}\right) \geq 10
$$

where $\chi(M)$ stands for the chromatic number of an (infinite) graph with the vertex set $M$ and the edge set consists of pairs of points at the distance 1 apart.

### 7.1 Introduction

We study colorings of a set Slice $(d, k, \varepsilon)=\mathbb{R}^{d} \times[0, \varepsilon]^{k}$ in a finite number of colors with the forbidden distance 1 between monochromatic points; further such sets are called slices. Slightly abusing the notation we say that $d$ is the dimension of a slice.

Define graph $G(d, k, \varepsilon)$, which vertices are the point of Slice $(d, k, \varepsilon)$ and edges connect points at the Euclidean distance 1 apart. Put

$$
\chi[\operatorname{Slice}(d, k, \varepsilon)]:=\chi[G(d, k, \varepsilon)],
$$

where $\chi(H)$ is the chromatic number of a graph $H$. Obviously for every positive $\varepsilon$ one has

$$
\chi\left(\mathbb{R}^{d}\right) \leq \chi[\operatorname{Slice}(d, k, \varepsilon)] \leq \chi\left(\mathbb{R}^{d+k}\right)
$$

Since $\chi\left(\mathbb{R}^{d}\right)=(3+o(1))^{d}$ (see [82]), the chromatic number of a slice is finite. So by the de BruijnErdős theorem it is achieved on a finite subgraph.

### 7.1.1 Nelson-Hadwiger problem and its planar generalizations

In this notation the classical Nelson-Hadwiger problem is to determine $\chi[\operatorname{Slice}(2,0,0)]$, but as usual we write $\chi\left(\mathbb{R}^{2}\right)$ for this quantity. The best known bounds up to the date are

$$
5 \leq \chi\left(\mathbb{R}^{2}\right) \leq 7
$$

The upper bound is a classical coloring of a regular hexagon tiling due to Isbell. The lower bound were obtained by de Grey [31] in 2018, breaking a 70 year-old record (another constructions are contained in [60, 40, 122, 102, 103]).

The study of slice colorings started at [69] with the following main result.

Theorem 7.1.1. For every positive $\varepsilon$ holds

$$
6 \leq \chi[\operatorname{Slice}(2,2, \varepsilon)]
$$

Theorem 7.1.1 is a strengthening of the result $\chi_{\varepsilon}\left(\mathbb{R}^{2}\right) \geq 6$ (Currie-Eggleton [30]), where $\chi_{\varepsilon}$ stands for the minimal number of colors, for which there is a coloring of plane without a pair of monochromatic points at the distance in the range $[1,1+\varepsilon]$. Exoo [39] conjectured that for every $\varepsilon>0$ holds $\chi_{\varepsilon}\left(\mathbb{R}^{2}\right) \geq 7$. Recently Voronov [124] proved this conjecture.

On the other hand Isbell' coloring implies inequality

$$
\chi_{\varepsilon}\left(\mathbb{R}^{2}\right) \leq 7
$$

for $\varepsilon<0.13 \ldots$. As a corollary, for every $k$ there is $\varepsilon_{k}>0$ such that for every positive $\varepsilon<\varepsilon_{k}$ holds

$$
\chi[\operatorname{Slice}(2, k, \varepsilon)] \leq 7
$$

### 7.1.2 The chromatic numbers of real 3-dimensional slices

First recall the best bounds on $\chi\left(\mathbb{R}^{3}\right)$. The best known lower bound $\chi\left(\mathbb{R}^{3}\right) \geq 6$ is due to Nechush$\tan$ [95]. The upper bound $\chi\left(\mathbb{R}^{3}\right) \leq 15$ is obtained independently by Coulson [28] and by Radoičić and Tóth [106].

The main result of this section is the following theorem.
Theorem 7.1.2 (Cherkashin-Kanel-Belov-Strukov-Voronov [123]). There is $\varepsilon_{0}>0$, such that for every positive $\varepsilon<\varepsilon_{0}$ holds

$$
10 \leq \chi[\operatorname{Slice}(3,6, \varepsilon)] \leq 15
$$

The upper bound immediately follows from the coloring of $\chi\left(\mathbb{R}^{3}\right)$ from [28, 106], similarly to 2 -dimensional case. The lower bound requires somewhat more complicated arguments than in two dimensions.

The following theorem is a quantitative strengthening of Theorem 10 from 69 and is of an independent interest.
Theorem 7.1.3 (Cherkashin-Kanel-Belov-Strukov-Voronov [23]). Let $T \subset \mathbb{R}^{d}$ be a regular simplex with the edge length $a=\sqrt{2 d(d+1)}$. Then every proper coloring of $\mathbb{R}^{d}$ in a finite number of colors contains a point from $T$ belonging to the closures of at least $d+1$ colors.
Corollary 7.1.1. For every positive $\varepsilon^{\prime}$ holds

$$
\chi_{\varepsilon^{\prime}}\left(\mathbb{R}^{3}\right) \geq 10
$$

Indeed, the orthogonal projection of a unit distance graph from the proof of Theorem 7.1.2 for a fixed $\varepsilon$ has distances between adjacent vertices in the range $\left[\sqrt{1-6 \varepsilon^{2}}, 1\right]$.

### 7.1.3 The chromatic numbers of 2-dimensional rational slices

Denote by $[0, \varepsilon]_{\mathbb{Q}}$ the set of rational numbers from $[0, \varepsilon]$. In paper 69 it is shown that

$$
\chi\left(\mathbb{Q} \times[0, \varepsilon]_{\mathbb{Q}}^{3}\right)=3 .
$$

Benda and Perles [7] show that $\chi\left(\mathbb{Q}^{4}\right)=4$. Thus the chromatic number of $\mathbb{Q}^{2} \times[0, \varepsilon]_{\mathbb{Q}}^{2}$ is at most 4 .
Proposition 7.1.1. For every $\varepsilon>0$ holds

$$
\chi\left(\mathbb{Q}^{2} \times[0, \varepsilon]_{\mathbb{Q}}^{2}\right)=4
$$

### 7.2 Notation and auxiliary lemmas

Here and after we focus on the following Slice $(3,6, \varepsilon) \subset \mathbb{R}^{9}$. By $S_{r}^{d}(x)$ we denote a $d$-dimensional sphere of the radius $r$ and centered at $x$.
Definition 7.2.1. An attached sphere of a simplex with vertices $\left\{v_{i}\right\}_{1 \leq i \leq k}, 3 \leq k \leq 4$ is a set of points at the distance 1 from each $v_{i}$ :

$$
S\left(v_{1}, \ldots, v_{k} ; 1\right):=\bigcap_{i} S\left(v_{i} ; 1\right) \subset \mathbb{R}^{9}
$$

Note that if the radius $r$ of a circumscribed $(k-2)$-dimensional sphere $v_{1}, \ldots, v_{k}$ is smaller than 1 , then $S\left(v_{1}, \ldots, v_{k} ; 1\right)$ is a $(9-k)$-dimensional sphere with the radius $\sqrt{1-r^{2}}$.
Definition 7.2.2. A t-equator of a sphere $S$ is a subsphere of the dimension $t$ which radius is equal to the radius of $S$.

As usual, $\bar{T}$ stands for the closure of a set $T$.
Definition 7.2.3. Let a metric $X$ be colored in a finite number of colors; denote these colors by $C_{1}, \ldots, C_{m}$. A chromaticity of a point $x \in X$ is the number of sets $\overline{C_{i}}, 1 \leq i \leq m$ containing $x$.
Lemma 7.2.1 (Knaster-Kuratowski-Mazurkiewicz). Suppose that $(d-1)$-dimensional simplex is covered by closed sets $X_{1}, \ldots, X_{d}$ in such a way that every face $I \subset[d]$ is contained in the union of $X_{i}$ over $i \in I$. Then all sets $X_{i}$ have a common point.

The following lemma is a spherical analogue of the planar lemma from [69]. The proof is also analogous; we provide it in the interest of completeness.
Lemma 7.2.2. Let $S_{r}^{2}$ be a sphere of radius $r>\sqrt{\frac{1}{2}}$, $\varepsilon$ be a positive number, and $Q \subset S_{r}^{2}$ be a $\varepsilon$-neighbourhood of a curve $\xi \subset S_{r}^{2}$, such that

$$
\operatorname{diam} \xi>\frac{\sqrt{4 r^{2}-1}}{r}
$$

Then $\chi(Q) \geq 3$.
Proof. Without loss of generality $\varepsilon<1$. Denote by $G(Q)$ the corresponding graph; we are going to find an odd cycle in $G(Q)$.

Consider a point $u \in \xi$. Since $\operatorname{diam} \xi>\frac{\sqrt{4 r^{2}-1}}{r}=\operatorname{diam} S(u ; 1)$, the intersection of $S(u ; 1)$ and $\xi$ is non-empty. Let $v \in S(u ; 1) \cap \xi$; consider such points $v_{1}, v_{2}, v_{3}, v_{4} \in S_{r}^{2}$ that $\left\|u-v_{1}\right\|=1$; $\left\|v_{i}-v_{i+1}\right\|=1 ; i=1,2,3$. If the angles at the vertices of polygonal chain $\operatorname{vuv}_{1} v_{2} v_{3} v_{4}$ are at most $\frac{\varepsilon}{2}$, then $\left\|v-v_{1}\right\|<\frac{\varepsilon}{2},\left\|u-v_{2}\right\|<\frac{\varepsilon}{2},\left\|v-v_{3}\right\|<\varepsilon,\left\|u-v_{4}\right\|<\varepsilon$, and hence $v_{i} \in Q, i=1,2,3,4$.
Note that

$$
\begin{aligned}
& l_{1}=\left\|u-v_{2}\right\| \in\left[0 ; 2 \sin \frac{\varepsilon}{4}\right], \\
& l_{2}=\left\|v_{2}-v_{4}\right\| \in\left[0 ; 2 \sin \frac{\varepsilon}{4}\right]
\end{aligned}
$$

can be chosen arbitrarily, and the oriented angle between vectors $\overrightarrow{v_{2} u}$ and $\overrightarrow{v_{2} v_{4}}$ can be independently chosen from $\left[-\frac{\varepsilon}{4} ; \frac{\varepsilon}{4}\right]$. Fix the line containing vector $\overrightarrow{v_{2}} \vec{u}$; one may choose it orthogonal to $u v$. Then a set of all possible $v_{4}$ contains a rhombus centered at $u$ with the side length $2 \sin \frac{\varepsilon}{4}$ and the angle $\frac{\varepsilon}{2}$. Then $G(Q)$ contains a path of length 4 between $u$ and an arbitrary point from a $\gamma$-neighbourhood of $u$, where $\gamma=\sin \frac{\varepsilon}{2} \sin \frac{\varepsilon}{4}$.

Thus one may move from $u$ to $v$ along $\xi$ by steps of size at most $\gamma$. Every such step corresponds to a path of length 4 in $G(Q)$; since $v$ is adjacent to $u$ we find a desired odd cycle in $G(Q)$.


Figure 7.1: A path of length four between $u$ and $v_{4}$.
Lemma 7.2.3. Suppose that a sphere $S_{r}^{2} \subset \mathbb{R}^{3}, r>\sqrt{\frac{1}{2}}$ has a proper coloring a finite number of colors. Then it has a point with the chromaticity at least 3.
Proof. Note that for $r>\sqrt{\frac{1}{2}}$

$$
\frac{\sqrt{4 r^{2}-1}}{r}<2 r .
$$

By compactness of the sphere it is sufficient to show that there is a spherical ball with an arbitrarily small radius containing points of at least 3 colors. Suppose the contrary: there is a proper coloring of the sphere and $\varepsilon>0$ such that every spherical ball with the radius $\varepsilon$ is colored in at most 2 colors. Consider a partition of the sphere onto cells such that every cell contains a ball with the radius $\delta=\varepsilon_{0} / 100$ and is contained in a ball of the radius $\varepsilon_{0} / 10$. Then every cell contains points of at most two colors, moreover all the adjacent cells are colored in the same two colors.

Consider an arbitrary cell with two colors (say, colors 1 and 2). Let $A_{0}$ be the region which is maximal by inclusion, that contains cells with one- or two-colored cells of colors 1 and 2 . The diameter of $A_{0}$ is smaller than $\frac{\sqrt{4 r^{2}-1}}{r}<2 r$ otherwise, by looking at any path between diametrally opposed points, we have a contradiction with Lemma 7.2 .2 . Let us consider the outer boundary $p$ of the reigon $A_{0}$. Every cell adjacent to $p$ is adjacent to some cell not in $A_{0}$, hence it is monochromatic; moreover, colors of all cells from $A_{0}$ along $p$ are the same, otherwise there would be a ball that contains cells of two different colors and cell not from $A_{0}$, which contradicts our assumption. So we may assume that all cells from $A_{0}$ along $p$ are of color 1 . Same argument shows that cells not from $A_{0}$ along $p$ cannot contain two different colors that are not 1 or 2 , and cannot contain the color 2 . Therefore all cells adjacent to $p$ are colored in colors 1 or 3 (maybe both). Consider the region $A_{1}$, that contains cells along $p$ of colors 1 and 3 , and is maximal by inclusion. We can apply to $A_{1}$ the smae argument, and by induction we obtain the sequence of 2 -colored regions $A_{i}$. Note that (spherical) diameter of $A_{i}$ is increasing by at least $\delta$ each step, so eventually we obtain the contradiction to Lemma 7.2.2.

The proof of the main result require technical statement on stability of circumscribed circle of a triangle with vertices of a form $\left(0,0,0, b_{1}, \ldots, b_{6}\right)$ with respect to a shifts by vectors from the main subspace $\mathbb{R}^{3}$, i.e. vectors of type $(p, q, r, 0, \ldots, 0)$. Such shifts will be called orthogonal. The next lemma will be applied for the case of $S^{5}$, but we prove it in the general case.

Lemma 7.2.4. Suppose that several points are chosen on $S^{k}$ so that minimal distance between two chosen points is $\Omega\left(\mathrm{m}^{-2}\right)$. Then there is the triangle $\mathcal{T}$ which vertices are amongst the chosen points satisfying the following condition. Every orthogonal shift of it vertices by $O\left(m^{-3}\right)$ causes change of radius $R$ of circumscribed circle of $\mathcal{T}$ by $O\left(\frac{R}{m^{2}}\right)$ and shift of its center by $O\left(\frac{R}{m}\right)$.

Proof. Let us find a triangle $\mathcal{T}_{0}$ from selected points with heights $\Omega\left(m^{-2}\right)$. Assume the contrary, id est that there is no such triangle. Let us consider the maximum distance between these points; say, it is achieved between points $A$ and $B$. Then all other points should lie in the $o\left(m^{-2}\right)$-neighborhood of the great circle $A B$ (any great circle $A B$, if the points $A$ and $B$ were diametrically opposite): indeed, otherwise the height from point $C$ to $A B$ is equal to $\Omega\left(m^{-2}\right)$ and it is the smallest of the heights of triangle $A B C$, since points $A$ and $B$ were chosen at the maximum distance and $A B C$ is suitable for the role of $\mathcal{T}_{0}$.

Let us consider the projections of the selected points onto the great circle $A B$ (they are uniquely determined). Since the pairwise distances between the selected points are equal to $\Omega\left(m^{-2}\right)$, and the points lie in the $o\left(m^{-2}\right)$-neighborhood of the great circle $A B$, the projections are separated from each other by at least by $\Omega\left(m^{-2}\right)$. One of the arcs $A B$ contains the projection of at least $m_{1} \geq m / 2$ points. Let us number the points according to the projection on this arc $A B$; let $C$ be the point with the number $\left[m_{1} / 2\right]$. Then $A C$ and $B C$ are equal to $\Omega(1 / m)$. Let $O$ be the circumcenter of triangle $A B C$. Let us denote the lengths of the sides $A B, B C, A C$ by $c, a, b$, respectively; let the lengths of the heights be equal to $h_{a}, h_{b}, h_{c}$.

It is clear that $\angle A C B$ is the largest of the angles of triangle $A B C$ and

$$
\angle A C B \leq \angle O C A+\angle O C B=\arccos \frac{b}{2 R}+\arccos \frac{a}{2 R} \leq 2 \arccos \frac{\Omega(1 / m)}{2 R}
$$

since the triangles $A C O$ and $B C O$ are isosceles. Then

$$
2 \arccos \frac{\Omega(1 / m)}{2 R}=\pi-\frac{\Omega(1 / m)}{R} .
$$

Therefore, $\sin \angle A C B=\sin (\pi-\angle A C B)=\Omega(1 / m)$. Since $A B$ is the longest side, the height from point $C$ is the smallest. Then, since the sine of at least one of the angles $A$ and $B$ is also equal to $\Omega(1 / m)$, the height from point $C$ is equal to $\Omega\left(1 / m^{2}\right)$.

The triangle $\mathcal{T}_{0}=A B C$ has been found; let us show that it is suitable as $\mathcal{T}$. Let us keep the notation for the parameters of the triangle $\mathcal{T}_{0}$ introduced above. Let the shifted points be $A^{\prime}, B^{\prime}$ and $C^{\prime}$. Let us denote by $\Delta q$ the change in the value of $q$ during the transition from $A B C$ to $A^{\prime} B^{\prime} C^{\prime}$.

Let us show that an orthogonal shift of the ends of the segment $y_{1} y_{2}$ by $O\left(m^{-3}\right)$ changes (increases) the length of the segment $l$ by $O\left(m^{-6} l^{-1}\right)$. Let us denote the shifted points $z_{1}, z_{2}$, respectively. Due to orthogonality, $\left(y_{i}-y_{j}, z_{j}-y_{j}\right)=0$. The square of the new length is

$$
\begin{gathered}
\left(z_{1}-z_{2}, z_{1}-z_{2}\right)=\left\|\left(z_{1}-y_{1}\right)+\left(y_{1}-y_{2}\right)+\left(y_{2}-z_{2}\right)\right\|^{2}= \\
=\left(z_{1}-y_{1}, z_{1}-y_{1}\right)+\left(y_{1}-y_{2}, y_{1}-y_{2}\right)+\left(z_{2}-y_{2}, z_{2}-y_{2}\right)-2\left(z_{1}-y_{1}, z_{2}-y_{2}\right) ;
\end{gathered}
$$

That is, the difference in the squares of the lengths is estimated as

$$
\left(z_{1}-z_{2}, z_{1}-z_{2}\right)-\left(y_{1}-y_{2}, y_{1}-y_{2}\right)=O\left(m^{-6}\right)
$$

It remains to apply the difference of squares formula.
It turns out that $\Delta a=O\left(m^{-6} a^{-1}\right)$, similarly for other sides. Let $H$ be the base of the height $C H$, since $A B$ is the greatest, $H$ belongs to the segment $A B$. Note that the length of the height $h_{c}$ cannot decrease, and on the other hand, the distance from the shifted vertex $C$ to the point that is projected into $H$ changes by no more than $O\left(m^{-6} h_{c}^{-1}\right)$, and the length of the new height $h_{c}^{\prime}$ does not exceed this distance. Let $S$ be the area of triangle $A B C$, then

$$
\Delta S=O\left(c \Delta h_{C}+\Delta c \cdot h_{C}\right)=O\left(h_{c} m^{-6} c^{-1}\right)+O\left(c m^{-6} h_{c}^{-1}\right),
$$

hence,

$$
\frac{\Delta S}{S}=O\left(m^{-6} c^{-2}\right)+O\left(m^{-6} h_{c}^{-2}\right)=O\left(m^{-2}\right)
$$

Using the well-known formula

$$
R=\frac{a b c}{4 S}
$$

we get

$$
\begin{gathered}
\Delta R=O\left(\max \left(\frac{\Delta a \cdot b c}{S}, \frac{\Delta b \cdot a c}{S}, \frac{\Delta c \cdot a b}{S}, \frac{a b c \Delta S}{S^{2}}\right)\right)= \\
O\left(\max \left(\frac{\Delta a}{a} R, \frac{\Delta b}{b} R, \frac{\Delta c}{c} R, \frac{\Delta S}{S} R\right)\right)=O\left(\frac{R}{m^{2}}\right)
\end{gathered}
$$

which is what was required.
Let us limit the shift of the center of the circumscribed circle when changing along one coordinate. We showed above that the heights and sides of a triangle change slightly when the vertices are orthogonally shifted by $O\left(\mathrm{~m}^{-3}\right)$, which means it will be possible to repeat the following estimate several times.

Let us consider three-dimensional Cartesian coordinates in which $C$ is the center, the plane $A B C$ is generated by the first two coordinates, and the latter corresponds to the infinitesimal shift. Then the normal to the plane $A B C$ has the form

$$
\bar{v}_{1}=\overline{A C} \times \overline{B C}=(0,0,2 S) .
$$

Then the normal to the plane $A^{\prime} B^{\prime} C^{\prime}$ is equal to

$$
\bar{v}_{2}=A^{\prime} C^{\prime} \times B^{\prime} C^{\prime}=\left(A_{y} B_{z}^{\prime}-A_{z}^{\prime} B_{y}, A_{x} B_{z}^{\prime}-A_{z}^{\prime} B_{x}, 2 S\right) .
$$

Without loss of generality, $a \geq b$ and then $\left|A_{x}\right|,\left|B_{x}\right|,\left|A_{y}\right|,\left|B_{y}\right| \leq a$. Therefore $\left|A_{y} B_{z}^{\prime}-A_{z}^{\prime} B_{y}\right|, \mid A_{x} B_{z}^{\prime}-$ $A_{z}^{\prime} B_{x} \mid=O\left(a m^{-3}\right)$. Recall that $S=0.5 a h_{a}=\Omega\left(a m^{-2}\right)$, which implies $\left|A_{y} B_{z}^{\prime}-A_{z}^{\prime} B_{y}\right|, \mid A_{x} B_{z}^{\prime}-$ $A_{z}^{\prime} B_{x} \mid=O\left(S m^{-1}\right)$. Let us estimate the angle $\phi$ between the planes $A C B$ and $A^{\prime} B^{\prime} C^{\prime}$ :

$$
\cos \phi=\frac{\left(v_{1}, v_{2}\right)}{\sqrt{\left(v_{1}, v_{1}\right) \cdot\left(v_{2}, v_{2}\right)}}=\frac{4 S^{2}}{2 S \cdot \sqrt{4 S^{2}+O\left(S^{2} m^{-2}\right)}}=1-O\left(m^{-2}\right), \quad \phi=O\left(m^{-1}\right)
$$

Consequently, the change in the center does not exceed $O(R \sin \phi)=O(R / m)$.

### 7.3 Proof of Theorem 7.1.3

Suppose the contrary. Let $C_{i}$ - be a set of points $\mathbb{R}^{d}$ of color $i, 1 \leq i \leq m$. Define

$$
C_{i}^{*}:=\overline{\overline{\operatorname{Int}} \overline{C_{i}}} \quad \text { (the closure of the interior of the closure). }
$$

Split every $C_{i}^{*}$ into connected components (with respect to the standard topology):

$$
C_{i}^{*}=\bigcup_{\alpha \in A_{i}} D_{\alpha}
$$

Put also $\left\{D_{\alpha}\right\}=\bigcup_{i=1}^{m} \bigcup_{\alpha \in A_{i}} D_{\alpha}$.
(i) Sets $C_{i}^{*}$ cover $\mathbb{R}^{d}$. Suppose the contrary, i.e. $\exists v: \forall i \quad v \notin C_{i}^{*}$. Then there is an open ball $B(v ; \eta)$ such that

$$
B(v ; \eta) \cap C_{i}^{*}=\emptyset ; \quad B(v ; \eta) \subset \bigcup C_{i} .
$$

Consider an arbitrary ball

$$
B^{1} \subset B(v ; \eta) \backslash \bar{C}_{1} .
$$

Then $B^{1}$ cannot be a subset of $\bar{C}_{i}$, otherwise the intersection of the interior of $\bar{C}_{i}$ and $B(v ; \eta)$ is nonempty. Define a sequence of balls

$$
B^{k+1} \subset B^{k} \backslash \bar{C}_{k}
$$

Note that points of $B^{m+1}$ do not belong to any $\bar{C}_{i}$, which is a contradiction.
(ii) Suppose that a point $v \in T$ belongs to exactly $k$ sets $C_{i}^{*}$. Assume that $k \leq n$ (otherwise the chromaticity of $v$ is at least $d+1$ ). Then for every $\mu_{0}>0$ there is $\mu<\mu_{0}$ such that the sphere $S(v ; 1-\mu)$ does not intersect at least one of those $k$ sets.


Figure 7.2: Illustration to item (ii).
We can assume without loss of generality that $v \in C_{i}^{*}, 1 \leq i \leq k$. Suppose the contrary, i.e. there is such a $\mu_{0}>0$ that for every $\mu \in\left(0, \mu_{0}\right)$ holds

$$
S(v ; 1-\mu) \cap C_{i}^{*} \neq \emptyset, \quad 1 \leq i \leq k .
$$

By the definition of $C_{i}^{*}$ any neighbourhood of an arbitrary point $x \in C_{i}^{*}$ contains a point from $\operatorname{Int} \bar{C}_{i}$. Hence, the set

$$
\mathcal{M}_{0}:=\left\{\mu \in\left(0, \mu_{0}\right) \mid \exists S(v ; 1-\mu) \cap \operatorname{Int} \bar{C}_{i}=\emptyset\right\}
$$

is closed and nowhere dense.
Fix some $\mu \in\left(0, \mu_{0}\right) \backslash \mathcal{M}_{0}$. One may choose points $x_{1}, \ldots, x_{k}$ in such a way that

$$
x_{i} \in S(v ; 1-\mu) \cap \operatorname{Int} \bar{C}_{i}, \quad 1 \leq i \leq k
$$

and $\left\{v, x_{1}, \ldots, x_{k}\right\}$ are in a general position (i.e. all the simplices are non-degenerate). Consider any $\eta>0$ such that $B\left(x_{i} ; \eta\right) \subset C_{i}^{*}, 1 \leq i \leq k$. Put

$$
z \in B(0 ; \eta) ; \quad y_{i}=x_{i}+z
$$

Define

$$
\begin{gathered}
w_{0}=\phi\left(x_{1}, \ldots, x_{k}\right):=\operatorname{Argmin}_{u \in U}\|u-v\|, \quad U=\bigcap_{1 \leq i \leq k} S\left(y_{i} ; 1\right), \\
w_{1}=\phi\left(y_{1}, \ldots, y_{k}\right) .
\end{gathered}
$$

By the construction the color of $w_{1}$ differs from the colors of $y_{1}, \ldots, y_{k}$.
In a small neighbourhood of $\left\{y_{i}\right\}$ function $w(\cdot)$ is properly defined and continuous. Choose points

$$
y_{i}^{\prime} \in C_{i}, \quad 1 \leq i \leq k,
$$

for which exists $w_{2}=\phi\left(y_{1}^{\prime}, \ldots, y_{k}^{\prime}\right)$. Then

$$
w_{2} \in \bigcup_{j=k+1}^{m} C_{j} .
$$

At the same time the quantity

$$
\delta\left(y_{1}^{\prime}, \ldots, y_{k}^{\prime}\right)=\max _{1 \leq i \leq k}\left\|y_{i}^{\prime}-y_{i}\right\|
$$

can be chosen arbitrarily small and hence

$$
w_{1} \in \bigcup_{j=k+1}^{m} \overline{C_{j}}
$$

Since $z \in B(0 ; \eta)$ was chosen arbitrarily

$$
B(w ; \eta) \subset \bigcup_{j=k+1}^{m} \overline{C_{j}}
$$

Hence an arbitrary neighbourhood of $w_{0}$ has an inner point of at least one set $\overline{C_{j}}, k+1 \leq j \leq m$, so $w_{0} \in C_{j}^{*}$ for some $j$. Note that $w_{0}=\phi\left(y_{1}, \ldots, y_{k}\right) \rightarrow v$ with $\mu \rightarrow 0$, thus $v$ belongs to at least one of sets $C_{j}^{*}, k+1 \leq j \leq m$, which contradicts to the initial assumption.
(iii) There is a cover of $T$ by sets from $\left\{D_{\alpha}\right\}$, such that every set from the cover is contained in a closed unit ball. By (ii) every point is covered by at least one set that satisfies the condition. Axiom of choice finishes the proof of the item. For every color $i$ denote by $\left\{D_{\beta}^{(i)}\right\}$ the chosen sets.
(iv) There is a finite cover of $T$ by closed sets $D_{i k}^{\prime}, 1 \leq i \leq m, 1 \leq k \leq K_{i}$, such that every set from the cover is the union of some sets from $\left\{D_{\alpha}\right\}$ and also is contained in a closed unit ball.

For every $i, 1 \leq i \leq m$ consider a sequence $v_{1}^{i}, v_{2}^{i}, \cdots \in \bigcup D_{\beta}^{(i)}$ such that

$$
\begin{gathered}
\gamma\left(v_{j}^{i}\right)=i \\
v_{s+1}^{i} \in \bigcup D_{\beta}^{(i)} \backslash \bigcup_{1 \leq j \leq s} B\left(v_{j}^{i} ; 1\right)
\end{gathered}
$$

Let the sequence be maximal (with respect to inclusion). The pairwise distances $v_{j}^{i}, j=1,2, \ldots$ are at least 1 , so the sequence is finite because $T$ is bounded. Now let us define

$$
D_{i k}^{\prime}=B\left(v_{k}^{i} ; 1\right) \cap\left(\bigcup D_{\beta}^{(i)}\right) \backslash \bigcup_{1 \leq j \leq k-1} D_{i j}^{\prime} .
$$



Figure 7.3: Illustration to item (iv). The construction of sets $D_{i k}^{\prime}$

Every set $D_{i k}^{\prime}$ is separated from other connected components of $C_{i}^{*}$ by a neighbourhood of a sphere, without points from $C_{i}^{*}$ (see Fig. 7.3). Thus these sets are closed.

Come back to the main construction and note that every set $D_{i k}^{\prime}$ cannot intersect every face of simplex $T$, because it is contained in an open unit ball while the inner radius of $T$ is equal to 1 . Split the cover $\mathcal{D}^{\prime}=\bigcup_{i}\left\{D_{i k}^{\prime}\right\}$ into $d+1$ subfamilies, in the way that every set subfamily consists of sets that do not intersect a face of $T$. Clearly there is a bijection between subfamilies and the vertices of $T$. Let $X_{i}, i=1, \ldots, d+1$ be the unions of sets over corresponding subfamilies. By Lemma 7.2.1 sets $X_{i}$ have a common point $x_{*}$, and thus an arbitrary neighbourhood of $x_{*}$ intersects with at least $d+1$ sets from $\left\{D_{\alpha}\right\}$. They belong to at least $d+1$ different $\left\{C_{i}^{*}\right\}$, because $\left\{D_{\alpha}\right\}$ are connected components. Hence, $x^{*}$ has the chromaticity at least $d+1$.

### 7.4 Proof of Theorem 7.1.2

Outline of the proof. Suppose the contrary to the statement. First, we find points $v_{1}, v_{2}, v_{3}, v_{4}$ of different colors, such that the intersection $I$ of attached sphere $S\left(v_{1}, v_{2}, v_{3}, v_{4} ; 1\right)$ and the slice contains 2-equator $S^{2}$ of the sphere and we also require the radius of the sphere to be close to 1 .

Then $I$ is close (in the sense of Hausdorff distance) to $S_{1-\eta}^{2} \times[0, \varepsilon]^{3}$, where $\eta$ is small enough. Then one can follow the arguments from [69], that were applied in the case of 2-dimensional slices. Note that the sets of colors of $I$ and $\left\{v_{1}, v_{2}, v_{3}, v_{4}\right\}$ are disjoint.

Let us find points $v_{5}, v_{6}, v_{7} \in I$, such that an equator of the corresponding attached sphere belongs to the slice. The attached sphere of $v_{1}, \ldots, v_{7}$ has an equator belonging to the slice, so the intersection of $v_{1}, \ldots, v_{7}$ contains a spherical neighbourhood of a circle. It requires 3 additional colors in addition to the colors of points $v_{1}, \ldots, v_{7}$.

Step 1. Finding of points $v_{1}, v_{2}, v_{3}, v_{4}$, which attached sphere has the radius closed to 1 and the great circle belonging to the slice.

This requires the 3 -dimensional subspace $U$ spanned by $v_{1}, v_{2}, v_{3}, v_{4}$, to be "almost orthogonal" to the main subspace $\mathbb{R}^{3}$, and the circumradius of the simplex $v_{1} v_{2} v_{3} v_{4}$ in $U$ to be small enough.


Figure 7.4: Construction of a rainbow 10-point set.
Consider the standard Cartesian coordinate system in slice $\mathbb{R}^{3} \times[0, \varepsilon]^{6}$ :

$$
v=\left(x_{1}, x_{2}, x_{3}, y_{1}, \ldots, y_{6}\right), \quad x_{i} \in \mathbb{R}, \quad y_{i} \in[0, \varepsilon] .
$$

For a given point $v=\left(x_{1}, x_{2}, x_{3}, y_{1}, \ldots, y_{6}\right)$ define projections

$$
p_{R}(v)=\left(x_{1}, x_{2}, x_{3}, 0, \ldots, 0\right) \quad \text { and } \quad p_{\varepsilon}(v)=\left(0,0,0, y_{1}, \ldots, y_{6}\right)
$$

Consider sphere $S:=S_{\varepsilon_{1}}^{5}$ of the radius $\varepsilon_{1}<\varepsilon / 2$ centered at $(0,0,0, \varepsilon / 2, \varepsilon / 2, \ldots, \varepsilon / 2)$; note that $S \subset(0,0,0) \times[0, \varepsilon]^{6}$. Let $T \subset \mathbb{R}^{3}$ be an arbitrary regular tetrahedron with the edge length $2 \sqrt{6}$ and the center at the origin and $u$ be an arbitrary point of sphere $S$. By Lemma 7.1.3 every set $T \times\{u\} \subset T \times S$ has a point with chromaticity at least 4 .

Fix the parameters $\delta, h>0$, which values will be chosen later.
Consider a set of points $U=\left\{u_{1}, \ldots, u_{m}\right\} \subset S$ such that $\left\|u_{i}-u_{j}\right\| \geq \delta, i \neq j$ and $m$ is maximal. Obviously $m=\Omega\left(\delta^{-5}\right)$. Match every point $u_{i} \in U$ with an arbitrary point $u_{i}^{*} \in T \times\left\{u_{i}\right\}$ with chromaticity at least 4 .

Consider how $T$ is cut by a cubic mesh with edge length $h$ :

$$
T_{i, j, k}:=\bigsqcup_{i, j, k} T \cap Z_{i, j, k} ; \quad Z_{i, j, k}:=[i h,(i+1) h) \times[j h,(j+1) h) \times[k h,(k+1) h),
$$

where $i, j, k$ are integers. Since $T \subset \mathbb{R}^{3}$ is bounded, one has

$$
\#\left\{(i, j, k): T \cap Z_{i, j, k} \neq \emptyset\right\}=O\left(h^{-3}\right)
$$

Consider points $w_{i}=p_{R}\left(u_{i}^{*}\right) \in T$. Put $h=\delta^{3 / 2}$. There is a cell $T_{a, b, c}$ such that it contains at least

$$
m=\frac{\Omega\left(\delta^{-5}\right)}{O\left(\delta^{-9 / 2}\right)}=\Omega\left(\delta^{-\frac{1}{2}}\right)
$$

points from $\left\{w_{i}\right\}$. Note that $h=O\left(m^{-3}\right), \delta=O\left(m^{-2}\right)$ and

$$
\operatorname{diam} T_{a, b, c} \leq \sqrt{3} h=\sqrt{3} \delta^{3 / 2}=O\left(m^{-3}\right)
$$

Now apply Lemma 7.2 .4 for these $m$ points. It gives a triangle $\mathcal{T}=w_{1} w_{2} w_{3}$ such that its arbitrary small orthogonal shift, in particular the triangle $u_{1}^{*} u_{2}^{*} u_{3}^{*}$ has circumradius at most $\left(1 / 4+O\left(m^{-2}\right)\right) \varepsilon$ and its circumcircle $\omega$ belongs to the slice. Let us construct a (five-dimensional) sphere $S^{*}$ on $\omega$ as the diameter and choose $v_{4}$ as the most distant point from the plane $u_{1}^{*} u_{2}^{*} u_{3}^{*}$ on the sphere $S^{*}$. Then the simplex $u_{1}^{*} u_{2}^{*} u_{3}^{*} v_{4}$ is a non-degenerate simplex whose circumscribed sphere belongs to the interior of the slice.

It remains to choose in arbitrarily small neighborhoods of the points $u_{1}^{*}, u_{2}^{*}$ and $u_{3}^{*}$ the points $v_{1}, v_{2}$ and $v_{3}$, respectively, in such a way that the points $v_{1}, v_{2}, v_{3}$ and $v_{4}$ have pairwise different colors.

Step 2. Finding in sphere $S\left(v_{1}, v_{2}, v_{3}, v_{4} ; 1\right)$ points $v_{5}, v_{6}, v_{7}$ of different colors such that attached (2dimensional) sphere $S\left(v_{1}, \ldots, v_{7} ; 1\right)$ has a 2 -equator belonging to the slice. Note that $S\left(v_{1}, \ldots, v_{7}, 1\right)$ is the intersection of $S\left(v_{1}, v_{2}, v_{3}, v_{4} ; 1\right)$ and $S\left(v_{5}, v_{6}, v_{7} ; 1\right)$. A proper choice of $\varepsilon_{1}, h$ makes radii of the spheres and the distance between its centers close to 1 . Then the radius of $S\left(v_{1}, \ldots, v_{7}, 1\right)$ tends to $\frac{\sqrt{3}}{2}>\frac{1}{2}$.

Suppose the intersection of an attached sphere with the slice

$$
M:=S\left(v_{1}, v_{2}, v_{3}, v_{4} ; 1\right) \bigcap \mathbb{R}^{3} \times[0, \varepsilon]^{6}=S_{1-\eta}^{5} \bigcap \mathbb{R}^{3} \times[0, \varepsilon]^{6}
$$

is properly colored and the equator $M_{E}=S_{1-\eta}^{2}$ belongs to the slice.
Let $H \subset \mathbb{R}^{9}$ be the 6 -dimensional subspace containing $S_{1-\eta}^{5}$. Consider a coordinates in $H$ such that $M_{E}$ belongs to the subspace spanned by the first 3 coordinates. For every point $u \in M_{E}$, $u=\left(u_{1}, u_{2}, u_{3}, 0,0,0\right)$ consider a sphere

$$
S^{2}(u ; \nu)=\left\{\sqrt{1-\nu^{2}} u+\xi \mid \quad \xi=\left(0,0,0, \xi_{4}, \xi_{5}, \xi_{6}\right) ;\|\xi\|=\nu\right\}
$$

Note that $S^{2}(u ; \nu)$ is a subset of $M$ when $\nu$ is small enough.
For every $u$ consider the following regular pentagon belonging to $S^{2}(u ; \nu)$ :

$$
w_{u, k}=\left(u_{1}, u_{2}, u_{3}, \cos \frac{2 \pi k}{5} \nu, \sin \frac{2 \pi k}{5} \nu, 0\right), \quad k=1, \ldots, 5
$$

Let $u$ be a point. If one can find among $w_{u, 1}, \ldots, w_{u, 5}$ points of three different colors, then they can be taken as $v_{5}, v_{6}, v_{7}$. Otherwise vertices of every pentagon are colored in at most 2 different colors, i.e. there is a color with at least three representatives. Call this color dominating at $u$.

Consider an auxiliary coloring $\pi$ in which every point of $M_{E}$ has its dominating color. Let us show that $\pi$ is proper. Indeed if the distance between $p, q \in M_{E}$ is equal to 1 , then $\left\|w_{p, k}-w_{q, k}\right\|=1$ for every $k$, so by the pigeonhole principle dominating colors at $p$ and $q$ are different.

By Lemma 7.2 .3 sphere $M_{E}$ has a point $u^{*}$ with chromaticity at least 3 , i.e. an arbitrary neighbourhood of $u^{*}$ has three points of different dominating colors. Then one may choose from corresponding pentagons 3 points of different colors in a way that chosen points lie in three small neighbourhoods of points $w_{u^{*}, 1}, \ldots, w_{u^{*}, 5}$. Every triangle with vertices in these points is non-degenerate, and has sides of length at least $\nu$.

Step 3. Recall that every point from $v_{1}, v_{2}, v_{3}, v_{4}$ and every point from $v_{5}, v_{6}, v_{7}$ lie at the distance 1 apart. Moreover, $v_{1}, v_{2}, v_{3}, v_{4}$ have pairwise different colors; the same holds for $v_{5}, v_{6}$, $v_{7}$. Moreover, by Lemma 7.2 .2 (applied to equator that lies in the slice) the intersection of attached sphere $S\left(v_{1}, \ldots, v_{7} ; 1\right)$ and the slice has the chromatic number at least 3 . Hence we show that a proper coloring of the slice requires at least $4+3+3=10$ colors, as desired.

### 7.5 Proof of Proposition 7.1.1

Consider the following 4 points in $\mathbb{Q}^{2} \times[0, \varepsilon]_{\mathbb{Q}}^{2}$ :

$$
\begin{gather*}
A=(0,0,0,0),  \tag{7.1}\\
B=\left(q, \frac{1}{2}, \alpha, \beta\right), \quad C=\left(q,-\frac{1}{2}, \alpha, \beta\right),  \tag{7.2}\\
D=(2 q, 0,0,0) . \tag{7.3}
\end{gather*}
$$

So we have

$$
\begin{equation*}
|A B|^{2}=|A C|^{2}=|B D|^{2}=|C D|^{2}=q^{2}+\frac{1}{4}+\alpha^{2}+\beta^{2} \tag{7.4}
\end{equation*}
$$

Our goal is to choose numbers $q \in \mathbb{Q}$ and $\alpha, \beta \in[0, \varepsilon]_{\mathbb{Q}}$ such that expression (7.4) is equal to 1 . Let $q=a / 2 b$, where $a$ and $b$ are some integers. Then we need

$$
\begin{equation*}
\alpha^{2}+\beta^{2}=\frac{3}{4}-\frac{a^{2}}{(2 b)^{2}}=\frac{3 b^{2}-a^{2}}{4 b^{2}} . \tag{7.5}
\end{equation*}
$$

It is enough for $(a, b)$ to satisfy

$$
\begin{equation*}
3 b^{2}-a^{2}=2 \tag{7.6}
\end{equation*}
$$

so if $b$ is large enough, we can put $\alpha=\beta=\frac{1}{2 b}$.
Let us construct a series of solutions to (7.6) as follows. Given the solution $\left(a_{n}, b_{n}\right)$, we build next pair as

$$
\begin{equation*}
\left(a_{n+1}, b_{n+1}\right)=\left(7 a_{n}+12 b_{n}, 4 a_{n}+7 b_{n}\right) \tag{7.7}
\end{equation*}
$$

One can check that $\left(a_{n+1}, b_{n+1}\right)$ is a solution to (7.6) by straightforward computation and use of assumption that so is $\left(a_{n}, b_{n}\right)$. Now by taking $\left(a_{0}, b_{0}\right)=(1,1)$ we get an infinite sequence of solutions with $b_{n}$ strictly increasing without limit. So for any given $\varepsilon$ there is some $n_{\varepsilon}$ such that for $n>n_{\varepsilon}$

$$
\begin{equation*}
\frac{3 b_{n}^{2}-a_{n}^{2}}{4 b_{n}^{2}}=\frac{1}{2 b_{n}^{2}}<2 \varepsilon^{2} \tag{7.8}
\end{equation*}
$$

which implies $1 / 2 b<\varepsilon$
Now we are going to find such integers $x$ and $y$ that

$$
\begin{equation*}
x \cdot \frac{a_{n}}{b_{n}}+y \cdot \frac{a_{n+1}}{b_{n+1}}=1 \tag{7.9}
\end{equation*}
$$

or

$$
\begin{equation*}
x \cdot a_{n} b_{n+1}+y \cdot a_{n+1} b_{n}=b_{n} b_{n+1} . \tag{7.10}
\end{equation*}
$$

So existence of such $x$ and $y$ is equivalent to

$$
\begin{equation*}
\operatorname{gcd}\left(a_{n} b_{n+1}, a_{n+1} b_{n}\right) \mid b_{n} b_{n+1} . \tag{7.11}
\end{equation*}
$$

It is sufficient to show that $\operatorname{gcd}(\ldots)=1$ :

$$
\begin{gather*}
\operatorname{gcd}\left(a_{n} b_{n+1}, a_{n+1} b_{n}\right) \mid\left(a_{n} b_{n+1}-a_{n+1} b_{n}\right)  \tag{7.12}\\
a_{n} b_{n+1}-a_{n+1} b_{n}=a_{n}\left(4 a_{n}+7 b_{n}\right)-b_{n}\left(7 a_{n}+12 b_{n}\right)=4 a_{n}^{2}-12 b_{n}^{2}=-4\left(3 b_{n}^{2}-a_{n}^{2}\right)=-8 . \tag{7.13}
\end{gather*}
$$

And from (7.7) it is clear that

$$
\begin{align*}
a_{n+1} & \equiv a_{n} \equiv \ldots \equiv a_{0}=1 \quad(\bmod 2)  \tag{7.14}\\
b_{n+1} & \equiv b_{n} \equiv \ldots \equiv b_{0}=1 \quad(\bmod 2) \tag{7.15}
\end{align*}
$$

So $\operatorname{gcd}(\ldots)=1$ as required.
Finally, let $\chi\left(\mathbb{Q}^{2} \times[0, \varepsilon]_{\mathbb{Q}}^{2}\right)=3$. Then points $A$ and $D$ have the same color. Hence, each point at the distance $k \cdot a_{n} / b_{n}+l \cdot a_{n+1} / b_{n+1}$ (where $1 / 2 b_{n}^{2}<2 \varepsilon^{2}$ and $k, l$ are integers) from 0 has the same color. Taking $k=x$ and $l=y$, one can obtain that $(1,0,0,0)$ has the same color. A contradiction.

Remark 7.5.1. Recursion formula (7.7) was obtained the following way. Consider a ring $\mathbb{Z}[\sqrt{3}]$. It has the norm

$$
N(a+b \sqrt{3})=(a+b \sqrt{3})(a-b \sqrt{3})=a^{2}-3 b^{2} .
$$

Then (7.6) transforms to an equation $N(\alpha)=-2$. Norm is multiplicative: $N(\alpha \beta)=N(\alpha) N(\beta)$ for any $\alpha, \beta \in \mathbb{Z}[\sqrt{3}]$. So if $N(\alpha)=-2$ and $N(\zeta)=1$, then $N(\alpha \zeta)=-2$. For (7.7) one can take $\zeta=7+4 \sqrt{3}$.

### 7.6 Further questions

Question 7.6.1. Let $\mathcal{M}_{d}$ be a family of compact convex set $\mathbb{R}^{d}$ such that a proper coloring of any $\mathbb{R}^{d}$ have a point of chromaticity at least $d+1$ in every $M \in \mathcal{M}_{d}$. Evaluate $V_{d}^{*}=\inf _{M \in \mathcal{M}_{d}} \operatorname{Vol} M$ from above.

Theorem 7.1.3 gives the bound $V_{d}^{*} \leq \frac{\sqrt{d+1}}{d!\sqrt{2^{d}}} \cdot(\sqrt{2 d(d+1)})^{d}=\frac{\sqrt{d^{d}(d+1)^{d+1}}}{d!}$.

## Bibliography

[1] Rudolf Ahlswede and Levon H. Khachatrian. The complete intersection theorem for systems of finite sets. European Journal of Combinatorics, 18(2):125-136, 1997.
[2] Enrique G. Alvarado, Bala Krishnamoorthy, and Kevin R. Vixie. The maximum distance problem and minimal spanning trees. International Journal of Analysis and Applications, 19(5):633659, 2021.
[3] Andrii Arman and Troy Retter. An upper bound for the size of a $k$-uniform intersecting family with covering number $k$. Journal of Combinatorial Theory, Series A, 147:18-26, 2017.
[4] József Balogh, Alexandr Kostochka, and Andrei Raigorodskii. Coloring some finite sets in $\mathbb{R}^{n}$. Discussiones Mathematicae Graph Theory, 33(1):25-31, 2013.
[5] Mikhail Basok, Danila Cherkashin, Nikita Rastegaev, and Yana Teplitskaya. On uniqueness in Steiner problem. International Mathematics Research Notices, page rnae025, 2024.
[6] Mikhail Basok, Danila Cherkashin, and Yana Teplitskaya. Inverse maximal and average distance minimizer problems. arxiv preprint arXiv:2212.01903, 2022.
[7] Miro Benda and Micha Perles. Colorings of metric spaces. Geombinatorics, 9(3):113-126, 2000.
[8] Marc Bernot, Vicent Caselles, and Jean-Michel Morel. Optimal transportation networks: models and theory. Springer, 2008.
[9] Edward Bierstone and Pierre D. Milman. Semianalytic and subanalytic sets. Publications Mathématiques de l'Institut des Hautes Études Scientifiques, 67(1):5-42, 1988.
[10] Marcus Brazil, Ronald L. Graham, Doreen A. Thomas, and Martin Zachariasen. On the history of the Euclidean Steiner tree problem. Archive for History of Exact Sciences, 68:327-354, 2014.
[11] Rainer E. Burkard, Tibor Dudás, and Thomas Maier. Cut and patch Steiner trees for ladders. Discrete Mathematics, 161(1-3):53-61, 1996.
[12] G. Buttazzo and E. Stepanov. Optimal transportation networks as free Dirichlet regions for the Monge-Kantorovich problem. Ann. Sc. Norm. Super. Pisa Cl. Sci. (5), 2(4):631-678, 2003.
[13] Giuseppe Buttazzo, Edouard Oudet, and Eugene Stepanov. Optimal transportation problems with free Dirichlet regions. In Variational methods for discontinuous structures, pages 41-65. Springer, 2002.
[14] D. D. Cherkashin, A. S. Gordeev, G. A. Strukov, and Y. I. Teplitskaya. Maximal distance minimizers for a rectangle. arXiv preprint arXiv:2106.00809, 2021.
[15] Danila Cherkashin and Sergei Kiselev. Independence numbers of johnson-type graphs. Bulletin of the Brazilian Mathematical Society, New Series, 54(3):30, 2023.
[16] Danila Cherkashin, Anatoly Kulikov, and Andrei Raigorodskii. On the chromatic numbers of small-dimensional Euclidean spaces. Discrete Applied Mathematics, 243:125-131, 2018.
[17] Danila Cherkashin and Fedor Petrov. Branching points in the planar Gilbert-Steiner problem have degree 3. To appear in Pure and Applied Functional Analysis, 2024.
[18] Danila Cherkashin and Yana Teplitskaya. On the horseshoe conjecture for maximal distance minimizers. ESAIM: Control, Optimisation and Calculus of Variations, 24(3):1015-1041, 2018.
[19] Danila Cherkashin and Yana Teplitskaya. An overview of maximal distance minimizers problem. arXiv preprint arXiv:2212.05607, 2022.
[20] Danila Cherkashin and Yana Teplitskaya. A self-similar infinite binary tree is a solution to the Steiner problem. Fractal and Fractional, 7(5):414, 2023.
[21] Danila Cherkashin and Vsevolod Voronov. On the chromatic number of 2-dimensional spheres. Discrete ${ }^{6}$ Computational Geometry, 71:467-479, 2024.
[22] Danila D. Cherkashin. About maximal number of edges in hypergraph-clique with chromatic number 3. Moscow Journal of Combinatorics and Number Theory, 1(3):3-11, 2011.
[23] D.D. Cherkashin. On the chromatic numbers of Johnson-type graphs. Zapiski Nauchnuh Seminarov POMI, 518:192-200, 2022. [English translation will appear in Journal of Mathematical Sciences].
[24] D.D. Cherkashin, A.S. Gordeev, G.A. Strukov, and Y.I. Teplitskaya. On minimizers of the maximal distance functional for a planar convex closed smooth curve. arXiv preprint arXiv, 2020.
[25] F. R. K. Chung and R. L. Graham. Steiner trees for ladders. Annals of Discrete Mathematics, 2:173-200, 1978.
[26] Earl A. Coddington and Norman Levinson. Theory of ordinary differential equations. Tata McGraw-Hill Education, 1955.
[27] Maria Colombo, Antonio De Rosa, and Andrea Marchese. On the well-posedness of branched transportation. Communications on Pure and Applied Mathematics, 74(4):833-864, 2021.
[28] David Coulson. A 15-colouring of 3-space omitting distance one. Discrete mathematics, 256(1):83-90, 2002.
[29] R. Courant and H. Robbins. What is mathematics? London: Oxford University Press, 1941.
[30] James D. Currie and Roger B. Eggleton. Chromatic properties of the Euclidean plane. arXiv preprint arXiv:1509.03667, 2015.
[31] Aubrey D. N. J. de Grey. The chromatic number of the plane is at least 5. Geombinatorics, 25(1):18-31, 2018.
[32] M. Deza. Solution d'un problème de Erdős-Lovász. Journal of Combinatorial Theory, Series B, 16(2):166-167, 1974.
[33] Ding-Zhu Du, Frank K. Hwang, and J. F. Weng. Steiner minimal trees for regular polygons. Discrete © Computational Geometry, 2:65-84, 1987.
[34] Herbert Edelsbrunner, Alexandr O. Ivanov, and Roman N. Karasev. Current open problems in discrete and computational geometry. Modelirovanie i Analiz Informatsionnyh Sistem, 19(5):517, 2012.
[35] Herbert Edelsbrunner and Nataliya Strelkova. On the configuration space of Steiner minimal trees. arXiv preprint arXiv:1906.06577, 2019.
[36] Herbert Edelsbrunner and Nataliya P. Strelkova. On the configuration space of Steiner minimal trees. Russian Mathematical Surveys, 67(6):1167-1168, 2012.
[37] P. Erdôs, C. Ko, and R. Rado. Intersection theorems for systems of finite sets. Quart. J. Math. Oxford Ser.(2), 12:313-320, 1961.
[38] Paul Erdős and László Lovász. Problems and results on 3-chromatic hypergraphs and some related questions. Infinite and finite sets, 10(2):609-627, 1975.
[39] Geoffrey Exoo. ع-unit distance graphs. Discrete \& Computational Geometry, 33(1):117-123, 2005.
[40] Geoffrey Exoo and Dan Ismailescu. The chromatic number of the plane is at least 5: A new proof. Discrete $\mathcal{G}$ Computational Geometry, pages 1-11, 2018.
[41] José F. Fernando. On the set of local extrema of a subanalytic function. Collectanea mathematica, 71(1):1-24, 2020.
[42] Peter Frankl. A stability result for families with fixed diameter. Combinatorics, Probability and Computing, 26(4):506-516, 2017.
[43] Peter Frankl. A near-exponential improvement of a bound of Erdős and Lovász on maximal intersecting families. Combinatorics, Probability and Computing, pages 1-7, 2019.
[44] Peter Frankl and Zoltán Füredi. Forbidding just one intersection. Journal of Combinatorial Theory, Series A, 39(2):160-176, 1985.
[45] Peter Frankl and Andrey Kupavskii. Intersection theorems for $\{0, \pm 1\}$-vectors and $s$-crossintersecting families. Moscow Journal of Combinatorics and Number Theory, 2(7):91-109, 2017.
[46] Peter Frankl and Andrey Kupavskii. Erdős-Ko-Rado theorem for $\{0, \pm 1\}$-vectors. Journal of Combinatorial Theory, Series A, 155:157-179, 2018.
[47] Peter Frankl and Andrey Kupavskii. Families of vectors without antipodal pairs. Studia Scientiarum Mathematicarum Hungarica, 55(2):231-237, 2018.
[48] Peter Frankl and Andrey Kupavskii. Correction to the article Intersection theorems for $(0, \pm 1)$ vectors and $s$-cross-intersecting families. Moscow Journal of Combinatorics and Number Theory, 8(4):389-391, 2019.
[49] Peter Frankl and Andrey Kupavskii. Intersection theorems for (- 1, 0, 1)-vectors. European Journal of Combinatorics, page 103830, 2024.
[50] Peter Frankl, Katsuhiro Ota, and Norihide Tokushige. Covers in uniform intersecting families and a counterexample to a conjecture of Lovász. Journal of Combinatorial Theory, Series A, 74(1):33-42, 1996.
[51] Peter Frankl and Norihide Tokushige. Uniform eventown problems. European Journal of Combinatorics, 51:280-286, 2016.
[52] Peter Frankl and Richard M. Wilson. Intersection theorems with geometric consequences. Combinatorica, 1(4):357-368, 1981.
[53] Michael R. Garey, Ronald L. Graham, and David S. Johnson. The complexity of computing Steiner minimal trees. SIAM Journal on Applied Mathematics, 32(4):835-859, 1977.
[54] Edgar N. Gilbert. Minimum cost communication networks. Bell System Technical Journal, 46(9):2209-2227, 1967.
[55] Edgar N. Gilbert and Henry O. Pollak. Steiner minimal trees. SIAM Journal on Applied Mathematics, 16(1):1-29, 1968.
[56] S. Glock, D. Kühn, A. Lo, and D. Osthus. The existence of designs via iterative absorption. Memoirs of American Mathematical Society, 284(1406), 2023.
[57] Christopher David Godsil and Joseph Zaks. Colouring the sphere. University of Waterloo research report, CORR 88-12, 1988.
[58] Alexey Gordeev and Yana Teplitskaya. On regularity of maximal distance minimizers in $\mathbb{R}^{n}$. To appear in The Annali della Scuola Normale Superiore di Pisa, Classe di Scienze, 2024.
[59] David A. Grable. More-than-nearly-perfect packings and partial designs. Combinatorica, 19(2):221-239, 1999.
[60] Marijn J. H. Heule. Computing small unit-distance graphs with chromatic number 5. Geombinatorics, 28(1):32-50, 2018.
[61] Robert Hochberg and Paul O'Donnell. Some 4-chromatic unit-distance graphs without small cycles. Geombinatorics, 5(4):137-141, 1996.
[62] Frank K. Hwang. A linear time algorithm for full Steiner trees. Operations Research Letters, 4(5):235-237, 1986.
[63] Frank K. Hwang, Dana S. Richards, and Pawel Winter. The Steiner tree problem, volume 53. Elsevier, 1992.
[64] Frank K. Hwang and J. F. Weng. The shortest network under a given topology. Journal of Algorithms, 13(3):468-488, 1992.
[65] Alexandr O. Ivanov and Alexey A. Tuzhilin. Uniqueness of Steiner minimal trees on boundaries in general position. Sbornik: Mathematics, 197(9):1309-1340, 2006.
[66] Vojtěch Jarník and Miloš Kössler. On minimal graphs containing $n$ given points. Časopis pro Pěstováni Matematiky a Fysiky, 63(8):223-235, 1934.
[67] Tommy R. Jensen and Bjarne Toft. Graph coloring problems. John Wiley \& Sons, 2011.
[68] Jeff Kahn and Gil Kalai. A counterexample to Borsuk's conjecture. Bulletin of the American Mathematical Society, 29(1):60-62, 1993.
[69] A. Kanel-Belov, V. Voronov, and D. Cherkashin. On the chromatic number of an infinitesimal plane layer. St. Petersburg Mathematical Journal, 29(5):761-775, 2018.
[70] Gyula O.H. Katona. Extremal problems for hypergraphs. In Combinatorics, pages 215-244. Springer, 1975.
[71] Peter Keevash. The existence of designs. arXiv preprint arXiv:1401.3665, 2014.
[72] Peter Keevash. The existence of designs II. arXiv preprint arXiv:1802.05900, 2018.
[73] Jeong Han Kim. Nearly optimal partial Steiner systems. Electronic Notes in Discrete Mathematics, 7:74-77, 2001.
[74] Sergey Kiselev. Supplementary files to the paper "Independence numbers of Johnson-type graphs", 2020. github.com/shuternay/Johnson-independence-numbers.
[75] Daniel J. Kleitman. On a combinatorial conjecture of Erdős. Journal of Combinatorial Theory, 1(2):209-214, 1966.
[76] Olga A. Kostina. On lower bounds for the chromatic number of spheres. Mathematical Notes, 105(1-2):16-27, January 2019.
[77] Alexandr V. Kostochka and Vojtech Rödl. Partial Steiner systems and matchings in hypergraphs. Random Structures \& Algorithms, 13(3-4):335-347, 1998.
[78] Steven G. Krantz and Harold R. Parks. The implicit function theorem: history, theory, and applications. Springer Science \& Business Media, 2002.
[79] Andrei Kupavskii. On the colouring of spheres embedded in $\mathbb{R}^{n}$. Sbornik: Mathematics, 202(6):859-886, 2011.
[80] Nikolai N. Kuzjurin. On the difference between asymptotically good packings and coverings. European Journal of Combinatorics, 16(1):35-40, 1995.
[81] Sergei K. Lando and Alexander K. Zvonkin. Graphs on Surfaces and Their Applications. Encyclopaedia of Mathematical Sciences. Springer Berlin Heidelberg, 2010.
[82] David G. Larman and C. Ambrose Rogers. The realization of distances within sets in Euclidean space. Mathematika, 19(1):1-24, 1972.
[83] Antoine Lemenant. About the regularity of average distance minimizers in $\mathbb{R}^{2}$. J. Convex Anal, 18(4):949-981, 2011.
[84] Antoine Lemenant. A presentation of the average distance minimizing problem. Journal of Mathematical Sciences, 181(6):820-836, 2012.
[85] Peter Lippmann, Enrique Fita Sanmartín, and Fred A. Hamprecht. Theory and approximate solvers for branched optimal transport with multiple sources. Advances in Neural Information Processing Systems, 35:267-279, 2022.
[86] László Lovász. On the ratio of optimal integral and fractional covers. Discrete mathematics, 13(4):383-390, 1975.
[87] László Lovász. Kneser's conjecture, chromatic number, and homotopy. Journal of Combinatorial Theory, Series A, 25(3):319-324, 1978.
[88] László Lovász. Self-dual polytopes and the chromatic number of distance graphs on the sphere. Acta Sci. Math.(Szeged), 45(1-4):317-323, 1983.
[89] F. J. MacWilliams and N. J. A. Sloane. The theory of error-correcting codes, volume 16. Elsevier, 1977.
[90] Greg Malen. Measurable colorings of $S_{r}^{2}$. Geombinatorics, 24(4):172-180, 2015.
[91] Zdzislaw A. Melzak. On the problem of Steiner. Canad. Math. Bull, 4(2):143-148, 1961.
[92] M. Miranda, Jr., E. Paolini, and E. Stepanov. On one-dimensional continua uniformly approximating planar sets. Calc. Var. Partial Differential Equations, 27(3):287-309, 2006.
[93] Sunra J. N. Mosconi and Paolo Tilli. 「-convergence for the irrigation problem. J. Convex Anal, 12(1):145-158, 2005.
[94] Zsigmond Nagy. A certain constructive estimate of the Ramsey number. Matematikai Lapok, 23(301-302):26, 1972.
[95] Oren Nechushtan. On the space chromatic number. Discrete Mathematics, 256(1-2):499-507, 2002.
[96] Konstantin I. Oblakov. Non-existence of distinct codirected locally minimal trees on a plane. Moscow University Mathematics Bulletin, 64(2):62-66, 2009.
[97] Patric R. J. Östergård. A fast algorithm for the maximum clique problem. Discrete Applied Mathematics, 120(1-3):197-207, 2002.
[98] Emanuele Paolini and Eugene Stepanov. Qualitative properties of maximum distance minimizers and average distance minimizers in $\mathbb{R}^{n}$. Journal of Mathematical Sciences, 122(3):3290-3309, 2004.
[99] Emanuele Paolini and Eugene Stepanov. Existence and regularity results for the Steiner problem. Calculus of Variations and Partial Differential Equations, 46(3-4):837-860, 2013.
[100] Emanuele Paolini and Eugene Stepanov. On the Steiner tree connecting a fractal set. arXiv preprint arXiv:2304.01932, 2023.
[101] Emanuele Paolini, Eugene Stepanov, and Yana Teplitskaya. An example of an infinite Steiner tree connecting an uncountable set. Advances in Calculus of Variations, 8(3):267-290, 2015.
[102] Jaan Parts. The chromatic number of the plane is at least 5 - a human-verifiable proof. Geombinatorics, 29(4):137-166, 2020.
[103] Jaan Parts. Graph minimization, focusing on the example of 5-chromatic unit-distance graphs in the plane. Geombinatorics, 30(2):77-102, 2020.
[104] Christian Pommerenke. Boundary behaviour of conformal maps, volume 299 of Grundlehren der mathematischen Wissenschaften. Springer Berlin Heidelberg, 1992.
[105] Roman Prosanov. Chromatic numbers of spheres. Discrete Mathematics, 341(11):3123-3133, November 2018.
[106] Radoš Radoičić and Géza Tóth. Note on the chromatic number of the space. In Discrete and computational geometry, Algorithms and Combinatorics book series, volume 25, pages 695-698. Springer, 2003.
[107] Andrei M. Raigorodskii. On the chromatic number of a space. Russian Mathematical Surveys, 55(2):351-352, 2000.
[108] Andrei M. Raigorodskii. Borsuk's problem and the chromatic numbers of some metric spaces. Russian Mathematical Surveys, 56(1):103, 2001.
[109] Andrei M. Raigorodskii. On the chromatic numbers of spheres in $\mathbb{R}^{n}$. Combinatorica, 32(1):111123, 2012.
[110] Vojtěch Rödl. On a packing and covering problem. European Journal of Combinatorics, 6(1):6978, 1985.
[111] Joachim H. Rubinstein and Doreen A. Thomas. Graham's problem on shortest networks for points on a circle. Algorithmica, 7:193-218, 1992.
[112] Joachim H. Rubinstein, Doreen A. Thomas, and Nicholas C. Wormald. Steiner trees for terminals constrained to curves. SIAM Journal on Discrete Mathematics, 10(1):1-17, 1997.
[113] Filippo Santambrogio and Paolo Tilli. Blow-up of optimal sets in the irrigation problem. The Journal of Geometric Analysis, 15(2):343-362, 2005.
[114] Rolf Schneider. Convex bodies: the Brunn-Minkowski theory. Number 151 in Encyclopedia of Mathematics and its Applications. Cambridge university press, 2014.
[115] Isaak J. Schoenberg. Metric spaces and positive definite functions. Transactions of the American Mathematical Society, 44(3):522-536, 1938.
[116] Gustavus J. Simmons. The chromatic number of the sphere. Journal of the Australian Mathematical Society, 21(4):473-480, 1976.
[117] T. Sirgedas. The surface of a sufficiently large sphere has chromatic number at most 7. Geombinatorics, 30:138-151, 2021.
[118] Alexander Soifer. The mathematical coloring book: Mathematics of coloring and the colorful life of its creators. Springer Science \& Business Media, 2008.
[119] Paolo Tilli. Some explicit examples of minimizers for the irrigation problem. J. Convex Anal, 17(2):583-595, 2010.
[120] J. R. Tort. Un problème de partition de l'ensemble des parties à trois éléments d'un ensemble fini. Discrete Mathematics, 44(2):181-185, 1983.
[121] Marcus G. Volz, Marcus Brazil, Charl J. Ras, Konrad J. Swanepoel, and Doreen A. Thomas. The Gilbert arborescence problem. Networks, 61(3):238-247, 2013.
[122] V. A. Voronov, A. M. Neopryatnaya, and E. A. Dergachev. Constructing 5-chromatic unit distance graphs embedded in the Euclidean plane and two-dimensional spheres. Discrete Mathematics, 345(12):113106, 2022.
[123] V.A. Voronov, A.Ya. Kanel-Belov, G.A. Strukov, and D.D. Cherkashin. On the chromatic numbers of 3-dimensional slices. Zapiski Nauchnuh Seminarov POMI, 518:94-113, 2022. [English translation will appear in Journal of Mathematical Sciences].
[124] Vsevolod Voronov. On the chromatic number of the plane with an arbitrarily short interval of forbidden distances. arXiv preprint arXiv:2304.10163, 2023.
[125] Endre Weiszfeld and Frank Plastria. On the point for which the sum of the distances to $n$ given points is minimum. Annals of Operations Research, 167:7-41, 2009.
[126] Nicholas Wormald. A 4-chromatic graph with a special plane drawing. Journal of the Australian Mathematical Society, 28(1):1-8, 1979.
[127] D. Zakharov. Chromatic numbers of Kneser-type graphs. Journal of Combinatorial Theory, Series A, 172:105188, 2020.
[128] Dmitrii Zakharov. On the size of maximal intersecting families. Combinatorics, Probability and Computing, 33(1):32-49, 2024.

