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# Riemannian fast-marching on cartesian grids using Voronoi's first reduction of quadratic forms 

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#### Abstract

We address the numerical computation of distance maps with respect to Riemannian metrics of strong anisotropy. For that purpose we solve generalized eikonal equations, discretized using adaptive upwind finite differences on a Cartesian grid, in a single pass over the domain using a variant of the fast marching algorithm. The key ingredient of our PDE numerical scheme is Voronoi's first reduction, a tool from discrete geometry which characterizes the interaction of a quadratic form with an additive lattice. This technique, never used in this context, which is simple and cheap to implement, allows us to efficiently handle Riemannian metrics of eigenvalue ratio $10^{2}$ and more.

Two variants of the introduced scheme are also presented, adapted to sub-Riemannian and to Rander metrics, which can be regarded as degenerate Riemannian metrics and as Riemannian metrics perturbed with a drift term respectively. We establish the convergence of the proposed scheme and of its variants, with convergence rates. Numerical experiments illustrate the effectiveness of our approach in various contexts, in dimension up to five, including an original sub-Riemannian model related to the penalization of path torsion.


Keywords: Riemannian metric, Sub-Riemannian metric, Rander metric, Eikonal equation, Viscosity solution, Fast Marching Method, Voronoi Reduction.

## 1 Introduction

In this paper, we develop a new and efficient numerical method for the computation of distance maps with respect to anisotropic Riemannian metrics, sub-Riemannian metrics and Rander metrics. For that purpose we discretize generalized eikonal equations, also called static first order Hamilton Jacobi Bellman (HJB) Partial Differential Equations (PDEs), on a cartesian grid. The novelty of approach lies on a special representation of the Hamiltonian, via upwind finite finite differences on an adaptive stencil, which is designed using Voronoi's first reduction of quadratic forms [52] - a tool from discrete geometry mostly known for its applications in the study sphere packings and in number theory. For this reason, the method is referred to as Fast-Marching using Voronoi's First Reduction (FM-VR1).

Before entering the details of the addressed PDEs and of their discretisations, let us mention some of the potential applications. The standard eikonal equation reads $\|\mathrm{d} u\|=f$, with suitable boundary conditions, where $\mathrm{d} u$ denotes the differential of a function $u$ defined on a domain of $\mathbb{R}^{d}$. This PDE characterizes distance maps with respect an isotropic metric, defined locally as $f$-times

[^0]the euclidean metric, see [54] for a study and an overview of its numerous applications. Once the distance map is computed, globally optimal minimal paths (geodesics) w.r.t. the metric can be extracted by gradient descent, with numerous applications in e.g. image processing [44] or motion planning. This paper is devoted to the numerical solution of generalized eikonal equations, characterizing distance maps with respect to anisotropic metrics, of the three following types.

- A Riemannian metric on a domain of $\mathbb{R}^{d}$ is described by a field $\mathcal{M}$ of positive definite tensors, and gives rise to the generalized eikonal equation $\|\mathrm{d} u\|_{\mathcal{M}^{-1}}=1$. Numerical methods for Riemannian distance computation have applications in geometry processing [49], optics [28], statistics with the Fisher-Rao distance, ... In image processing and segmentation, anisotropic Riemannian metrics are often used to favor paths aligned with tubular structures of interest [29, 7, 15].
- A sub-Riemannian metric can be regarded as a degenerate Riemannian metric, which tensors have some infinite eigenvalues [38]. As a result, motion is only possible along a subspace of the tangent space, depending on the current position. This property is referred to as non-holonomy, and models for instance a robotic system with fewer controls than degrees of freedom. A fundamental instance is the Reeds-Shepp car model, posed on the configuration space $\mathbb{R}^{2} \times \mathbb{S}^{1}$, which can move forward and backward, rotate, but not translate sideways, see [51, 23] for a numerical study with applications to image segmentation and motion planning. A variant, presented in this paper 3.2 , is related to the penalization of path torsion.
- A Rander metric is defined locally as the sum of a Riemannian metric $\mathcal{M}$ and of a sufficiently small co-vector field $\hat{\eta}$, see [46] and $\$ 1.3$. These metrics are non-symmetric, thus define asymmetric distances, and give rise to the inhomogeneous generalization of eikonal equation $\|\mathrm{d} u-\hat{\eta}\|_{\mathcal{M}^{-1}}=1$. The travel-time of a boat subject to a drift due to water currents can be measured by integrating a Rander metric, see $\$ 3$, and its optimization is called Zermelo's problem [2, 12]. In image segmentation, the Chan-Vese energy of a region can be reformulated as the length of its contour measured w.r.t a Rander metric, see [14].

Our numerical approach has its limitations: it cannot address more general anisotropic metrics, such as those arising in seismic imaging [48], and it cannot handle domains discretized using triangulations or unstructured point sets as in [30]. Indeed, the algorithmic tools that we leverage [52] limit the scope of our method to eikonal equations whose hamiltonian has a quadratic structure, and to domains discretized using a cartesian grid. As often, efficiency is at the cost of specialization. In order to better describe the advantages and the specificities of our approach, let us formally state the addressed problem and review the existing methods.

This paper is devoted to the construction and analysis of a numerical scheme for computing the arrival times $u: \Omega \rightarrow \mathbb{R}$ of a front starting from the boundary of a domain $\Omega$, and propagating at unit speed w.r.t. a given metric $\mathcal{F}: T \Omega \cong \Omega \times \mathbb{R}^{d} \rightarrow[0, \infty]$ of one of the above three classes. Several classes of methods can be distinguished in the literature for such purposes.

- Eulerian schemes, such as the one presented in this paper, rely on a characterization of the arrival times as the unique viscosity solution 19 to the eikonal PDE , which reads

$$
\begin{equation*}
\forall \mathbf{p} \in \Omega, \mathcal{H}_{\mathbf{p}}(\mathrm{d} u(\mathbf{p}))=1 / 2, \quad \forall \mathbf{p} \in \partial \Omega, u(\mathbf{p})=0 \tag{1}
\end{equation*}
$$

where $\mathcal{H}$ denotes the Hamiltonian associated with the metric, i.e. the Legendre-Fenchel transform of the Lagrangian $\frac{1}{2} \mathcal{F}^{2}$. Finite differences are typically used for discretization 63 , 50, 5], although discontinuous Galerkin methods have recently been considered [32].

- Semi-Lagrangian schemes, rely on a self-consistency property of the arrival times referred to as Bellman's optimality principle: for any point and neighborhood $\mathbf{p} \in V(\mathbf{p}) \subseteq \Omega$, one has

$$
\begin{equation*}
u(\mathbf{p})=\min _{\mathbf{q} \in \partial V(\mathbf{p})}\left(u(\mathbf{q})+d_{\mathcal{F}}(\mathbf{q}, \mathbf{p})\right) \tag{2}
\end{equation*}
$$

where $d_{\mathcal{F}}$ denotes the path-length distance associated with the metric $\mathcal{F}$. In the semiLagrangian paradigm, a discrete counterpart of (2) is implemented numerically, by constructing polygonal stencils $V(\mathbf{p})$ with their vertices among the discretization points, interpolating the unknown $u$ on the facets of $\partial V(\mathbf{p})$, and locally approximating the distance by the metric $d_{\mathcal{F}}(\mathbf{q}, \mathbf{p}) \approx \mathcal{F}_{\mathbf{p}}(\mathbf{p}-\mathbf{q})$. A natural choice for $V(\mathbf{p})$ is the union of the triangles containing the vertex $\mathbf{p}$ in a given mesh of the domain $\Omega[9$. Following the discovery [30, 64 that a generalized acuteness property obeyed by the facets of $V(\mathbf{p})$ enables solving the discretized system in an efficient, single pass manner, see below, a number of more complex designs have been proposed [10, 57, 29, 1]. Constructions based on algorithmic geometry, introduced by the author in [34, 35, 36, allow satisfy the acuteness property for strongly anisotropic metrics while limiting the size and the cardinality of the stencils.

- Heat related methods solve a diffusion equation on a short time interval [21], or an elliptic equation with a small parameter [47], and exploit the relationship between the geodesic distance and the short time asymptotics of the heat kernel [62]. This approach is limited to Riemannian metrics, either isotropic or anisotropic [67]. Its efficiency is tied to the numerical cost of solving sparse linear systems discretizing a laplacian, which is often favorable over alternative methods, especially in dimension $d=2$, thanks to the existence of highly optimized linear algebra libraries. The method requires some parameter tuning, since it involves two small scales (in time and space), and looses accuracy or degenerates to a graph distance if their relative magnitude is incorrectly set [21].
- Path based techniques compute minimal geodesics directly, rather than front arrival times. Ray-tracing techniques solve Hamilton's Ordinary Differential Equations (ODE) of motion from a point $\mathbf{p}$ of interest, adjusting the initial velocity direction until the desired target is reached [13]. Path bending methods progressively deform a path joining two endpoints of interest, so as to obey Hamilton's ODEs [48]. Path based methods can be very accurate, by using high order ODE integration schemes, and do not suffer from the curse of dimensionality, since computational domain needs not be discretized. However, they lack robustness, have difficulty handling obstacles, and one usually cannot guarantee that the path found is globally the shortest one.

Among the first two classes of methods, Eulerian and semi-Lagrangian discretization schemes, a further distinction must be made depending on the numerical solver of the resulting coupled, non-linear system of equations resulting from the discretization of (1) or (2).

- Causal, single pass methods, such as the one proposed in this paper $\frac{1}{}$, rely on the fast marching method, a variant of Dijkstra's algorithm. This approach is computationally efficient, with complexity $\mathcal{O}(\lambda N \ln N)$ where $N$ is the number of discretization points, and $\lambda$ the average number of neighbors of each discretization point in the numerical scheme. However, it is only applicable if the discretization obeys a property referred to as causality, see Definition 2.1. Among Eulerian schemes, this property holds for the natural discretization of the isotropic eikonal equation [56], but could not be extended to anisotropic metrics

[^1]until the FM-VR1 presented here. In the case of semi-Lagrangian schemes, causality is related to a geometrical property of the stencil [30], see above.
Alternatively, Dijkstra's method can also be used directly to compute shortest paths on a network with sufficiently fine connectivity within the domain [39, 60, with edge-lengths defined according to the metric. The resulting an approximation of the front arrival times typically lacks accuracy, but adequate heuristics may improve it [11].

- Iterative methods apply Gauss-Siedel updates to the numerical solution, of the system of equations discretizing the problem of interest, until it meets a convergence criterion. The ordering of the discretization points in the Gauss-Siedel iterations may vary: fast sweeping methods alternate sweeps along the $2 d$ directions of the grid [68, whereas a priority queue is considered in [9]. The reported complexity of these methods is $\mathcal{O}\left(N^{1+1 / d}\right)$, where $N$ is the discrete domain cardinality and $d$ is the domain dimension. However the hidden constant depends on the problem instance, and may increase substantially if the metric is strongly anisotropic or if the minimal paths change direction often (due to obstacles or to non-uniformities in the metric) [7, 34]. Yet another approach is to introduce a time variable and solve $\partial_{t} u+\mathcal{H} u-1 / 2=0$, with suitable boundary conditions, which asymptotic steady state obeys the static equation (1), see [42, 32] and references therein.

Recapitulating, the FM-VR1 scheme introduced in this paper is the first numerical solver of eikonal equations that is simultaneously (i) Eulerian, (ii) solvable in a single pass ${ }^{2}$, and (iii) compatible with several classes of anisotropic metrics. For these reasons it is simple to implement, fast to solve numerically independently of the problem instance, and has a wide application scope and generalization potential, see e.g. [37] for a variant devoted to the global optimization of path energies involving curvature.

Contributions. We describe numerical schemes devoted to the computation of Riemannian, sub-Riemannian, and Rander distances, by solving the corresponding generalized eikonal equations. We prove convergence rates, based on the doubling of variables technique, see chapter 10 of [24] or [58], which applies rather directly in the Riemannian case but requires non-trivial adaptations in the sub-Riemannian and Rander cases. Numerical experiments illustrate the efficiency of our numerical schemes, in dimension $2 \leq d \leq 5$, and their potential applications in image segmentation and motion planning.

Outline. The rest of this introduction is devoted to general notations, to the description of Voronoi's first reduction which is a key ingredient of our discretization, and to elements of optimal control. The impatient reader may however jump to $\$ 1.1$, $\$ 1.2$ and $\$ 1.3$, where the numerical schemes are described and the convergence results stated, in the Riemannian, sub-Riemannian and Rander cases respectively. Convergence proofs are provided in $\mathbb{K}_{2}$, and in Appendices $A$ and $B$ respectively. Numerical experiments are presented in $\$ 3$.

General notations. The ambient space dimension is fixed and denoted by $d$. The Euclidean space and the Cartesian grid are respectively denoted

$$
\mathbb{E}:=\mathbb{R}^{d}, \quad \mathbb{L}:=\mathbb{Z}^{d}
$$

Let $\Omega \subseteq \mathbb{E}$ be a domain, assumed throughout the paper to be bounded; additional geometrical assumptions are required in some results. For any grid scale $h>0$ we let

$$
\begin{equation*}
\Omega_{h}:=\Omega \cap h \mathbb{L}, \quad \quad \partial \Omega_{h}:=(\mathbb{E} \backslash \Omega) \cap h \mathbb{L} \tag{3}
\end{equation*}
$$

[^2]Geometric points are denoted $\mathbf{p} \in \Omega$, vectors $\dot{\mathbf{p}} \in \mathbb{E}$, and co-vectors $\hat{\mathbf{p}} \in \mathbb{E}^{*}$. The symbol $\gamma$ is reserved for paths within $\bar{\Omega}$ and has the special convention that $\dot{\gamma}(t):=\frac{d}{d t} \gamma(t)$ denotes time derivation. We denote by $\|\dot{\mathbf{p}}\|$ the euclidean norm, by $(\dot{\mathbf{p}} \cdot \dot{\mathbf{q}})$ the scalar product, and by $\langle\hat{\mathbf{p}}, \dot{\mathbf{q}}\rangle$ the duality bracket, where $\dot{\mathbf{p}}, \dot{\mathbf{q}} \in \mathbb{E}$ are vectors and $\hat{\mathbf{p}} \in \mathbb{E}^{*}$ is a co-vector. Denote by $\mathrm{GL}(\mathbb{E}) \subseteq \mathrm{L}(\mathbb{E}, \mathbb{E})$ the group of invertible linear transformations, and by $\mathrm{GL}(\mathbb{L}) \subseteq \mathrm{GL}(\mathbb{E})$ the subgroup of those which leave the cartesian grid $\mathbb{L}$ invariant - equivalently their matrix has integer coefficients and determinant $\pm 1$. Denote by $\mathrm{S}(\mathbb{E}) \subseteq \mathrm{L}\left(\mathbb{E}, \mathbb{E}^{*}\right)$ the space of symmetric linear maps, by $\mathrm{S}^{+}(\mathbb{E})$ the subset of semi-definite ones, and by $\mathrm{S}^{++}(\mathbb{E})$ the positive definite ones. We adopt the notations

$$
\|\dot{\mathbf{p}}\|_{M}:=\sqrt{\langle M \dot{\mathbf{p}}, \dot{\mathbf{p}}\rangle}, \quad \hat{\mathbf{p}} \otimes \hat{\mathbf{p}} \in \mathrm{S}^{+}(\mathbb{E})
$$

for the norm of $\dot{\mathbf{p}} \in \mathbb{E}$ induced by $M \in \mathrm{~S}^{++}(\mathbb{E})$, and for the self outer product of $\hat{\mathbf{p}} \in \mathbb{E}^{*}$. The dual vector space $\mathbb{E}^{*}$ and the dual lattic $\epsilon^{3} \mathbb{L}^{*}$ can be naturally identified with their primal counterparts: $\mathbb{E}^{*} \cong \mathbb{E}$ and $\mathbb{L}^{*} \cong \mathbb{L}$ using the Euclidean structure, but the distinction is kept for clarity.

Voronoi's first reduction of quadratic forms. This tool originates from the field of lattice geometry, and is one of the key ingredients of our numerical scheme. It was introduced by Voronoi [66] with the purpose of classifying the equivalence classes of positive quadratic forms, i.e. elements of $\mathrm{S}^{++}(\mathbb{E})$, under the action of the group $G L(\mathbb{L})$, following a line of research dating back to Lagrange 31. The modern presentation of Voronoi's theory [52] involves Ryskov's convex polyhedron $\mathcal{P} \subseteq \mathrm{S}^{++}(\mathbb{E})$, and for each $D \in \mathrm{~S}^{++}\left(\mathbb{E}^{*}\right)$ a linear program $\mathcal{L}(D)$, defined as follows:

$$
\begin{equation*}
\mathcal{P}:=\left\{M \in \mathrm{~S}^{++}(\mathbb{E}) ; \forall \dot{\mathbf{e}} \in \mathbb{L} \backslash\{0\},\|\dot{\mathbf{e}}\|_{M} \geq 1\right\}, \quad \mathcal{L}(D):=\inf _{M \in \mathcal{P}} \operatorname{Tr}(M D) \tag{4}
\end{equation*}
$$

Introducing the duality bracket $\langle M, D\rangle:=\operatorname{Tr}(M D)$ between $\mathrm{S}(\mathbb{E})$ and $\mathrm{S}\left(\mathbb{E}^{*}\right)$, and observing that $\|\dot{\mathbf{e}}\|_{M}^{2}=\langle\langle M, \dot{\mathbf{e}} \otimes \dot{\mathbf{e}}\rangle\rangle$, one can rephrase Voronoi's optimization problem $\mathcal{L}(D)$ as follows

$$
\begin{equation*}
\text { minimize }\langle\langle M, D\rangle \text { subject to }\langle\langle M, \dot{\mathbf{e}} \otimes \dot{\mathbf{e}}\rangle \geq 1 \text { for all } \dot{\mathbf{e}} \in \mathbb{L} \backslash\{0\} . \tag{5}
\end{equation*}
$$

The vertices (resp. edges) of a polyhedron are its 0-dimensional (resp. 1-dimensional) facets.
Theorem (Voronoi, see [52]). The linear problem $\mathcal{L}(D)$ is feasible, for any $D \in \mathrm{~S}^{++}\left(\mathbb{E}^{*}\right)$, in the sense that the set of minimizers is non-empty and compact. In addition, the convex polytope $\mathcal{P} \subseteq \mathrm{S}(E)$ has a finite number of equivalence classes of vertices under the action of GL( $\mathbb{L})$.

The numerical method introduced in this paper, the FM-VR1, requires to solve one instance of (5) for each point of the discretization grid, for reasons explained in the next paragraph. Using a generic linear program solver for that purpose would be too slow to be practical, and we must rely on ad-hoc techniques leveraging the invariances of the problem. Selling's algorithm, see [53, 18 ] and Appendix D, serves that purpose in dimension $d \leq 3$, which is enough for our numerical experiments. In dimension $d>3$ one may solve (5) using a simplex-like method and relying on the classification of the vertices of the convex polyhedron $\mathcal{P}$, known in dimension $d \leq 8$ [17, 59], which classically are referred to as perfect quadratic forms and are studied for their relation with the densest periodic sphere packings [52].

Rather than its value or minimizer, we use the Karush-Kuhn-Tucker relations associated with the solution to the optimization problem (5) as the foundation our numerical scheme.

[^3]

Figure 1: Ellipsoid $\left\{\|\dot{\mathbf{p}}\|_{M} \leq 1\right\}$, and offsets appearing in the decomposition (6) of $D:=M^{-1}$, for some $M \in \mathrm{~S}^{++}\left(\mathbb{R}^{d}\right)$, in dimension $d=2$ (left) or $d=3$ (right).

These relations determine a decomposition of the input tensor $D$, see the next proposition, involving directions with integer entries corresponding to the active constraints in (5). Figure 1 illustrates, in dimension two and three, the close relationship between the anisotropy of the tensor $D$ and the directions of its decomposition, which locally define the stencil points of our adaptive discretization (17) of the eikonal PDE.

Proposition 1.1. The Kuhn-Tucker optimality conditions for the linear optimization problem $\mathcal{L}(D), D \in \mathrm{~S}^{++}\left(\mathbb{E}^{*}\right)$, imply that there exists $\left(\rho_{i}, \dot{\mathbf{e}}_{i}\right)_{i=1}^{d^{d^{\prime}}} \in\left(\mathbb{R}_{+} \times \mathbb{L}\right)^{d^{\prime}}$, where $d^{\prime}:=d(d+1) / 2$, such that

$$
\begin{equation*}
D=\sum_{1 \leq i \leq d^{\prime}} \rho_{i} \dot{\mathbf{e}}_{i} \otimes \dot{\mathbf{e}}_{i} . \tag{6}
\end{equation*}
$$

Furthermore there exists $C=C(d)$ such that for any $1 \leq i \leq d^{\prime}$

$$
\begin{equation*}
\left\|\dot{\mathbf{e}}_{i}\right\| \leq C \operatorname{Cond}(D)^{d-1}, \quad \quad \text { where } \operatorname{Cond}(D):=\sqrt{\|D\|\left\|D^{-1}\right\|} . \tag{7}
\end{equation*}
$$

Proof. The first point (6) follows from the feasibility of the linear program $\mathcal{L}(D)$, and its formulation (5). The number $d(d+1) / 2$ of contributions in this decomposition is the number of independent entries in a $d \times d$ symmetric matrix.

For proving the second point, the Euclidean space $\mathbb{E}$ is identified with its dual, which gives meaning to the trace $\operatorname{Tr}(M)$ and determinant $\operatorname{det}(M)$ of any $M \in \mathrm{~S}^{++}(\mathbb{E})$. Denote by $\left(M_{k}\right)_{k=1}^{K}$ a representative of each equivalence class of vertices of $\mathcal{P}$ under the action of of $\mathrm{GL}(\mathbb{L})$, see Theorem 1.1. Let $D \in \mathrm{~S}^{++}\left(\mathbb{E}^{*}\right)$ be arbitrary, and let $M$ be the minimizer of $\mathcal{L}(D)$. Then $M=A^{\mathrm{T}} M_{k} A$ for some $1 \leq k \leq K$ and some $A \in \operatorname{GL}(\mathbb{L})$. Thus $\operatorname{det}(M) \geq \min _{k=1}^{K} \operatorname{det}\left(M_{k}\right)=: c_{\Delta}>0$. On the other hand, $\operatorname{Tr}(M)\left\|D^{-1}\right\|^{-1} \leq \operatorname{Tr}(M D) \leq \operatorname{Tr}(\operatorname{Id} D) \leq d\|D\|$, by sub-optimality of $\operatorname{Id} \in \mathcal{P}$, hence $\operatorname{Tr}(M) \leq d \operatorname{Cond}(D)^{2}$. For any $1 \leq i \leq d^{\prime}$ one has $\left\|\dot{\mathbf{e}}_{i}\right\|_{M}=1$, since the corresponding constraint of the linear problem $\mathcal{L}(D)$ is active, hence as announced

$$
\left\|\dot{\mathbf{e}}_{i}\right\|^{2} \leq\left\|M^{-1}\right\|\left\|\dot{\mathbf{e}}_{i}\right\|_{M}^{2}=\left\|M^{-1}\right\| \leq \frac{\|M\|^{d-1}}{\operatorname{det}(M)} \leq \frac{\operatorname{Tr}(M)^{d-1}}{c_{\Delta}} \leq \frac{\left(d \operatorname{Cond}(D)^{2}\right)^{d-1}}{c_{\Delta}}
$$

The formula (6) is reminiscent of the decomposition $D=\sum_{i=1}^{d} \lambda_{i} \dot{\mathbf{v}}_{i} \otimes \dot{\mathbf{v}}_{i}$ of a symmetric tensor in terms of its normalized eigenvectors $\dot{\mathbf{v}}_{i}$ and of the associated eigenvalues $\lambda_{i}, 1 \leq i \leq d$. However the number of terms $d^{\prime} \neq d$ differs in (6), and most importantly the vectors $\dot{\mathbf{e}}_{i} \in \mathbb{L}$ have integer coefficients and can thus be used as offsets in a finite difference scheme on a grid, in contrast with the eigenvectors $\dot{\mathbf{v}}_{i} \in \mathbb{S}^{d-1}$ which do not belong to the grid unless $D$ is a diagonal matrix. In addition to Eikonal equations, which are the object of the present paper, the tensor decomposition (6) is used in [25] to design numerical schemes for two and three
dimensional anisotropic diffusion. An equivalent but two dimensional only concept is applied in [6] to Monge-Ampere equations, and in [8] to HJB PDEs of stochastic control. See 36] for estimates related to (7) in the average case upon random rotations of the tensor $D$, in dimension two.

Elements of optimal control. We refer to [3] for an overview of optimal control theory and its PDE formulations, and only introduce here the notations and definitions required for our purposes. Let $\mathfrak{C}(\mathbb{E})$ be the collection of compact and convex subsets of $\mathbb{E}$ containing the origin, equipped with the Hausdorff distance. Denote $\operatorname{Lip}(X, Y)$ the class of Lipschitz maps, with arbitrary Lipschitz constant, from a metric space $X$ to a metric space $Y$.

Definition 1.2. A family of controls is an element $\mathcal{B}$ of $\mathfrak{B}:=C^{0}(\bar{\Omega}, \mathfrak{C}(\mathbb{E}))$, which continuously associates to each point $\mathbf{p} \in \bar{\Omega}$ a control set $\mathcal{B}(\mathbf{p})$. A path $\gamma \in \operatorname{Lip}([0, T], \bar{\Omega})$, where $T \geq 0$, is said $\mathcal{B}$-controllable iff for almost every $t \in[0, T]$

$$
\begin{equation*}
\dot{\gamma}(t) \in \mathcal{B}(\gamma(t)), \quad \text { where } \dot{\gamma}(t):=\frac{d}{d t} \gamma(t) \tag{8}
\end{equation*}
$$

The minimal control time from $\mathbf{p} \in \bar{\Omega}$ to $\mathbf{q} \in \bar{\Omega}$, is defined as

$$
\begin{equation*}
T_{\mathcal{B}}(\mathbf{p}, \mathbf{q}):=\inf \{T \geq 0 ; \exists \gamma \in \operatorname{Lip}([0, T], \bar{\Omega}), \mathcal{B} \text {-controllable, } \gamma(0)=\mathbf{p}, \gamma(T)=\mathbf{q}\} \tag{9}
\end{equation*}
$$

The control sets corresponding to Riemannian, sub-Riemannian and Rander geometry are respectively ellipsoids, degenerate ellipsoids (with empty interior), and ellipsoids centered off the origin, see the illustrating figure. One easily shows that a minimal path from $\mathbf{p}$ to $\mathbf{q}$ exists as soon as $T_{\mathcal{B}}(\mathbf{p}, \mathbf{q})<\infty$, using Arzela-Ascoli's compactness theorem and the fact that $\Omega$ is bounded. See the appendices of [16, 23] for details, as well as related results such as the convergence of the control times and of the minimal paths associated with a converging family of controls under suitable assumptions. The above concepts can be rephrased in the framework of a local metric defined on the tangent space $\mathcal{F}: \bar{\Omega} \times \mathbb{E} \rightarrow[0, \infty]$ : given controls $\mathcal{B} \in \mathfrak{B}$, define for all $\mathbf{p} \in \bar{\Omega}$, $\dot{\mathbf{p}} \in \mathbb{E}$, and any path $\gamma \in \operatorname{Lip}([0,1], \bar{\Omega})$

$$
\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}):=\inf \{\lambda>0 ; \dot{\mathbf{p}} / \lambda \in \mathcal{B}(\mathbf{p})\}, \quad \quad \operatorname{Length}_{\mathcal{F}}(\gamma):=\int_{0}^{1} \mathcal{F}_{\gamma(t)}(\dot{\gamma}(t)) d t
$$

Note that these quantities can be infinite if the control sets have empty interior, such as in the sub-Riemannian case, and can be asymmetric $\left(\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}) \neq \mathcal{F}_{\mathbf{p}}(-\dot{\mathbf{p}})\right)$ if the control sets are not centered on the origin, as in the Rander case, see the illustrating figure and $\$ 1.2, \$ 1.3$, Conversely, the metric $\mathcal{F}$ uniquely determines the control sets $\mathcal{B}(\mathbf{p})=\left\{\dot{\mathbf{p}} \in \mathbb{E} ; \mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}) \leq 1\right\}$, and by time reparametrization the control time $T_{\mathcal{B}}(\mathbf{p}, \mathbf{q})$ from $\mathbf{p}$ to $\mathbf{q} \in \Omega$ is shown equal to the (quasi-)distance

$$
\begin{equation*}
d_{\mathcal{F}}(\mathbf{p}, \mathbf{q}):=\inf \left\{\operatorname{Length}_{\mathcal{F}}(\gamma) ; \gamma \in \operatorname{Lip}([0,1], \bar{\Omega}), \gamma(0)=\mathbf{p}, \gamma(1)=\mathbf{q}\right\} \tag{10}
\end{equation*}
$$

This paper is concerned with the exit time optimal control problem, which value function is defined for all $\mathbf{p} \in \bar{\Omega}$ by

$$
\begin{equation*}
u(\mathbf{q}):=\inf _{\mathbf{p} \in \partial \Omega} T_{\mathcal{B}}(\mathbf{p}, \mathbf{q}) \quad\left(=\inf _{\mathbf{p} \in \partial \Omega} d_{\mathcal{F}}(\mathbf{p}, \mathbf{q})\right) \tag{11}
\end{equation*}
$$



Figure: Examples of control sets for a (i) Riemannian, (ii) sub-Riemannian, and (iii) Rander metric. (iv) An admissible path, with tangents shown in red, w.r.t to some controls. (Illustration absent from journal version.)

The numerical computation of the function $u$ is the main topic of this paper. Under suitable assumptions [3], the function $u$ is the unique viscosity solution to the following HJB PDE involving the dual metric $\mathcal{F}^{*}: \bar{\Omega} \times \mathbb{E}^{*} \rightarrow \mathbb{R}_{+}$: for all $\mathbf{p} \in \Omega$

$$
\begin{equation*}
\mathcal{F}_{\mathbf{p}}^{*}(\mathrm{~d} u(\mathbf{p}))=1, \quad \quad \text { where } \mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}}):=\sup _{\mathbf{p} \neq 0} \frac{\langle\hat{\mathbf{p}}, \dot{\mathbf{p}}\rangle}{\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}})}, \tag{12}
\end{equation*}
$$

and $u(\mathbf{p})=0$ for all $\mathbf{p} \in \partial \Omega$. The formulations (12, left) and (1) left) of the eikonal PDE are equivalent, in view of the relation $\mathcal{H}=\frac{1}{2}\left(\mathcal{F}^{*}\right)^{2}$ between the Hamiltonian and the dual metric. Once $u$ is known, the shortest path from $\partial \Omega$ to $\mathbf{p} \in \Omega$ can be extracted by solving backwards in time the Ordinary Differential Equation (ODE)

$$
\begin{equation*}
\dot{\gamma}(t):=V(\gamma(t)), \quad \text { where } V(\mathbf{p}):=\mathrm{d} \mathcal{F}_{\mathbf{p}}^{*}(\mathrm{~d} u(\mathbf{p})), \tag{13}
\end{equation*}
$$

with final condition $\gamma(T)=\mathbf{p}$ where $T=u(\mathbf{p})$, see e.g. appendix C in [23]. In (13, right) the dual metric $\mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}})$ is differentiated w.r.t. the variable $\hat{\mathbf{p}}$. Note that $\mathrm{d} \mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}}) \in\left(\mathbb{E}^{*}\right)^{*} \cong \mathbb{E}$. For robust numerical geodesic backtracking it is essential to use an upwind estimation of the vector field $V(\mathbf{p})$, see Appendix E.

### 1.1 Riemannian metrics

A Riemannian metric on the bounded domain $\Omega \subseteq \mathbb{E}$, is described via a field of symmetric positive definite tensors $\mathcal{M} \in C^{0}\left(\bar{\Omega}, \mathrm{~S}^{++}(\mathbb{E})\right)$. The metric function $\mathcal{F}: \bar{\Omega} \times \mathbb{E} \rightarrow \mathbb{R}_{+}$has the expression

$$
\begin{equation*}
\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}):=\|\dot{\mathbf{p}}\|_{\mathcal{M}(\mathbf{p})} . \tag{14}
\end{equation*}
$$

Our objective is to compute the Riemannian distance $u: \bar{\Omega} \rightarrow \mathbb{R}$ to the boundary of $\Omega$, see (11), which is known to be the unique viscosity solution [20] to the Riemannian eikonal equation: for all $\mathbf{p} \in \Omega$

$$
\|\mathrm{d} u(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}=1 \quad \text { where } \mathcal{D}(\mathbf{p}):=\mathcal{M}(\mathbf{p})^{-1}
$$

and $u(\mathbf{p})=0$ for all $\mathbf{p} \in \partial \Omega$. Indeed, the dual to the Riemannian metric (14) reads $\mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}})=$ $\|\hat{\mathbf{p}}\|_{\mathcal{D}(\mathbf{p})}$. For each $\mathbf{p} \in \bar{\Omega}$, let $\left(\rho_{i}(\mathbf{p}), \dot{\mathbf{e}}_{i}(\mathbf{p})\right)_{i=1}^{d^{\prime}} \in\left(\mathbb{R}_{+} \times \mathbb{L}\right)^{d^{\prime}}$ be weights and offsets such that

$$
\begin{equation*}
\mathcal{D}(\mathbf{p})=\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p}) \dot{\mathbf{e}}_{i}(\mathbf{p}) \otimes \dot{\mathbf{e}}_{i}(\mathbf{p}) . \tag{15}
\end{equation*}
$$

In this paper, we advocate the use of Voronoi's first reduction of quadratic forms for obtaining the decomposition (15), see Proposition 1.2 . Our convergence results however only require to control the maximal stencil radius

$$
\begin{equation*}
r_{*}:=\max \left\{\left\|\dot{\mathbf{e}}_{i}(\mathbf{p})\right\| ; \mathbf{p} \in \bar{\Omega}, 1 \leq i \leq d^{\prime}\right\} \tag{16}
\end{equation*}
$$

If Proposition 1.2 is used for the stencil construction, then $r_{*}$ is by (7) bounded in terms of the maximal condition number of the metric, and the number of terms in (15) is $d^{\prime}=d(d+1) / 2$. For the sake of readability, we omit in the rest of the paper to write the dependence of the offset $\dot{\mathbf{e}}_{i}=\dot{\mathbf{e}}_{i}(\mathbf{p})$ on the point $\mathbf{p} \in \bar{\Omega}$. In the following, by $\max \{0, a, b\}^{2}$ we mean $(\max \{0, a, b\})^{2}$.

Theorem 1.3. Let $\mathcal{M} \in C^{0}\left(\bar{\Omega}, S^{++}(\mathbb{E})\right)$ be a Riemannian metric, and for all $\mathbf{p} \in \bar{\Omega}$ let $\mathcal{D}(\mathbf{p}):=$ $\mathcal{M}(\mathbf{p})^{-1}$ and $\left(\rho_{i}(\mathbf{p}), \dot{\mathbf{e}}_{i}(\mathbf{p})\right)_{i=1}^{d^{\prime}}$ be as in 15). Then for any $h>0$ there exists a unique solution $U_{h}: h \mathbb{L} \rightarrow \mathbb{R}$ to the following discrete problem: for all $\mathbf{p} \in \Omega_{h}$

$$
\begin{equation*}
h^{-2} \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p}) \max \left\{0, U_{h}(\mathbf{p})-U_{h}\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right), U_{h}(\mathbf{p})-U_{h}\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)\right\}^{2}=1 \tag{17}
\end{equation*}
$$

and $U_{h}(\mathbf{p})=0$ for all $\mathbf{p} \in \partial \Omega_{h}$. The solution $U_{h}$ can be computed via the fast-marching algorithm with complexity $\mathcal{O}\left(d^{\prime} N_{h} \ln N_{h}\right)$, where $N_{h}=\#\left(\Omega_{h}\right)$. If in addition the domain $\Omega$ satisfies an exterior cone condition, and if $\mathcal{M} \in \operatorname{Lip}\left(\bar{\Omega}, \mathrm{S}^{++}(\mathbb{E})\right)$, then for some constant $C=C(\mathcal{M}, \Omega)$ one has for all $h>0$

$$
\begin{equation*}
\max _{\mathbf{p} \in \Omega_{h}}\left|U_{h}(\mathbf{p})-u(\mathbf{p})\right| \leq C \sqrt{r_{*} h} . \tag{18}
\end{equation*}
$$

The estimate (18) outlines the importance of the stencil radius $r_{*}$, since it determines the effective scale $r_{*} h$ of the discretization and thus the accuracy of the numerical method. The construction of Proposition 1.2 is shown in [36] to minimize $r_{*}$, in dimension $d=2$. A convergence rate similar to (18) is obtained in [58] for the Ordered Upwind Method [57], a semi-Lagrangian solver of anisotropic eikonal equations. Note that the dependency of the constant $C=C(\Omega, \mathcal{M})$ in (18) with respect to the metric $\mathcal{M}$ is not explicited in Theorem 1.4. This point is analyzed in detail in the next sub-section, where we consider a family of increasingly anisotropic Riemannian metrics converging to a degenerate sub-Riemannian model.

Remark 1.4. The numerical scheme (17) relies on upwind finite differences, which are first order consistent with the absolute value of a directional derivative: for any $U \in C^{2}(\Omega), \mathbf{p} \in \Omega$, $\dot{\mathbf{e}} \in \mathbb{E}$, and any sufficiently small $h>0$

$$
\begin{aligned}
\max \{0, U(\mathbf{p})-U(\mathbf{p}-h \dot{\mathbf{e}}), U(\mathbf{p})-U(\mathbf{p}+h \dot{\mathbf{e}})\} / h & =\max \{0,\langle\mathrm{~d} U(\mathbf{p}), \dot{\mathbf{e}}\rangle,-\langle\mathrm{d} U(\mathbf{p}), \dot{\mathbf{e}}\rangle\}+\mathcal{O}(h), \\
& =|\langle\mathrm{d} U(\mathbf{p}), \dot{\mathbf{e}}\rangle|+\mathcal{O}(h) .
\end{aligned}
$$

The presence of " 0 " in the max, which may seem superfluous in view of the consistency analysis, is required for the monotony and causality of the numerical scheme, see Definition 2.1. In the related literature [50, 56], the above left upwind finite difference is often written in the following equivalent form: $\max \left\{\delta_{e}^{-} U(\mathbf{p}), \delta_{e}^{+} U(\mathbf{p})\right\}$ where denoting $a_{ \pm}:=\max \{0, \pm a\}$ one has

$$
\delta_{e}^{-} U(\mathbf{p}):=\left(\frac{U(\mathbf{p})-U(\mathbf{p}-h \dot{\mathbf{e}})}{h}\right)_{+}, \quad \delta_{e}^{+} U(\mathbf{p}):=\left(\frac{U(\mathbf{p}+h \dot{\mathbf{e}})-U(\mathbf{p})}{h}\right)_{-} .
$$

### 1.2 Sub-Riemannian metrics

We introduce a numerical approach to the computation of sub-Riemannian distances and geodesics, based on solving as described in $\$ 1.1$ the eikonal equations associated to a sequence of increasingly anisotropic approximate Riemannian metrics. This approach is related to [51], which however uses a different scheme for the Riemannian problems, and does not establish a convergence rate. More precisely our results apply to the slightly more general class of pre-Riemannian models.

Definition 1.5. A pre-Riemannian model on $\Omega$ is a finite family of vector fields $\dot{\omega}_{1}, \cdots, \dot{\omega}_{n} \in$ $\operatorname{Lip}(\bar{\Omega}, \mathbb{E})$. The control sets $\mathcal{B} \in \operatorname{Lip}(\bar{\Omega}, \mathfrak{C}(\mathbb{E}))$, and the semi-definite tensor field $\mathcal{D} \in \operatorname{Lip}\left(\bar{\Omega}, \mathrm{S}^{+}\left(\mathbb{E}^{*}\right)\right)$, for this model are defined for all $\mathbf{p} \in \bar{\Omega}$ by

$$
\mathcal{B}(\mathbf{p}):=\left\{\sum_{1 \leq i \leq n} \alpha_{i} \dot{\omega}_{i}(\mathbf{p}) ; \alpha \in \mathbb{R}^{n}, \sum_{1 \leq i \leq n} \alpha_{i}^{2} \leq 1\right\}, \quad \mathcal{D}(\mathbf{p}):=\sum_{1 \leq i \leq n} \dot{\omega}_{i}(\mathbf{p}) \otimes \dot{\omega}_{i}(\mathbf{p}) .
$$

A sub-Riemannian model [38] of step $k \geq 1$ is a pre-Riemannian model with the additional properties that the vector fields $\left(\dot{\omega}_{i}\right)_{1 \leq i \leq n}$ are smooth and that, together with their iterated commutators up to depth $k$, they span the tangent space $\mathbb{E}$ at each point $\mathbf{p} \in \bar{\Omega}$. The minimal control time $T_{\mathcal{B}}(\mathbf{p}, \mathbf{q})$ for a sub-Riemannian model is called the Carnot-Theodory distance, and by Chow's theorem it obeys $T_{\mathcal{B}}(\mathbf{p}, \mathbf{q}) \leq C\|\mathbf{p}-\mathbf{q}\|^{\frac{1}{k}}$, as $\mathbf{q} \rightarrow \mathbf{p} \in \Omega$. The distance $u$ to $\partial \Omega$ is the unique viscosity solution to the sub-Riemannian eikonal equation: $\|\mathrm{d} u(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}=1$ for all $\mathbf{p} \in \Omega$, and $u=0$ on $\partial \Omega$.

For better or worse, we do not use any techniques or results from sub-Riemannian geometry in this paper, but stick instead to the simpler pre-Riemannian concept. We do however make a further assumption.

Assumption 1.6. We fix a pre-Riemannian model $\left(\dot{\omega}_{i}\right)_{i=1}^{n}$, and assume that the exit time value function $u$, defined in (11), is bounded on $\Omega$. We further assume that the domain admits outward normals $\dot{\mathbf{n}}(\mathbf{p})$ with Lipschitz regularity on $\partial \Omega$, and that for each $\mathbf{p} \in \partial \Omega$ there exists $1 \leq i \leq n$ such that $\dot{\mathbf{n}}(\mathbf{p}) \cdot \dot{\omega}_{i}(\mathbf{p}) \neq 0$.

The finiteness of $u$ on $\Omega$ is a global controllability assumption, and it is obviously required if one intends to prove convergence rates of discrete approximations of $u$. The second assumption is related to short time local controllability at the boundary [3]. Together, these assumptions imply the Lipschitz regularity of $u$, see $\$$ A.1.
Definition 1.7. A completion of a pre-Riemannian model $\left(\dot{\omega}_{i}\right)_{i=1}^{n}$ is a second finite family of vector fields $\dot{\omega}_{1}^{*}, \ldots, \dot{\omega}_{n^{*}}^{*} \in \operatorname{Lip}(\bar{\Omega}, \mathbb{E})$, such that $\dot{\omega}_{1}(\mathbf{p}), \cdots, \omega_{n}(\mathbf{p}), \dot{\omega}_{1}^{*}(\mathbf{p}), \cdots, \dot{\omega}_{n^{*}}^{*}(\mathbf{p})$ spans $\mathbb{E}$ for each $\mathbf{p} \in \bar{\Omega}$.

For each $0<\varepsilon \leq 1$ the augmented pre-Riemannian model $\left(\omega_{1}, \cdots, \omega_{n}, \varepsilon \omega_{1}^{*}, \cdots, \varepsilon \omega_{n^{*}}^{*}\right)$ is equivalent (i.e. has the same control sets) to the Riemannian model of metric $\mathcal{M}_{\varepsilon}:=\mathcal{D}_{\varepsilon}^{-1}$, where pointwise on $\bar{\Omega}$

$$
\begin{equation*}
\mathcal{D}_{\varepsilon}=\mathcal{D}+\varepsilon^{2} \mathcal{D}^{*}, \quad \text { with } \mathcal{D}^{*}:=\sum_{1 \leq i \leq n^{*}} \dot{\omega}_{i}^{*} \otimes \dot{\omega}_{i}^{*} . \tag{19}
\end{equation*}
$$

In order to solve numerically the pre-Riemannian exit time problem, our strategy is to apply the scheme of Theorem 1.4 to the positive definite (but strongly anisotropic) Riemannian metric $\mathcal{M}_{\varepsilon}$, for small $\varepsilon>0$. Convergence towards the pre-Riemannian exit times $u: \Omega \rightarrow \mathbb{R}$ is established in the next theorem, when the relaxation parameter $\varepsilon$ and grid scale $h$ tend to 0 suitably.

Theorem 1.8. Consider a pre-Riemannian model $\dot{\omega}_{1}, \cdots, \dot{\omega}_{n} \in \operatorname{Lip}(\bar{\Omega}, \mathbb{E})$ obeying Assumption 1.7, and a completion $\dot{\omega}_{1}^{*}, \cdots, \dot{\omega}_{n}^{*}$. For each $0<\varepsilon \leq 1$ let $u_{\varepsilon}$ denote the distance to $\partial \Omega$ for the Riemannian metric $\mathcal{M}_{\varepsilon}$, and let $U_{h, \varepsilon}$ be the discrete solution of (17) with scale $h>0$. Then

$$
\begin{equation*}
\max _{\bar{\Omega}}\left|u-u_{\varepsilon}\right| \leq C \varepsilon, \quad \max _{\Omega_{h}}\left|u_{\varepsilon}-U_{h, \varepsilon}\right| \leq C^{\prime} \sqrt{r_{\varepsilon} h} \tag{20}
\end{equation*}
$$

where $r_{\varepsilon}$ denotes the maximal stencil radius for $\mathcal{M}_{\varepsilon}$, see (16), and where $C, C^{\prime}$ only depend on $\Omega,\left(\dot{\omega}_{i}\right)_{i=1}^{n}$, and $\left(\dot{\omega}_{i}^{*}\right)_{i=1}^{n^{*}}$. In particular $U_{h, \varepsilon} \rightarrow u$ uniformly as $\varepsilon \rightarrow 0$ and $h r_{\varepsilon} \rightarrow 0$.

By construction the condition number of the tensors $\mathcal{M}_{\varepsilon}$ is $\mathcal{O}\left(\varepsilon^{-1}\right)$, hence $r_{\varepsilon} \leq C \varepsilon^{-(d-1)}$ if Proposition 1.2 is used for the stencil construction. The convergence rate $\max _{\Omega_{h}}\left|U_{h, \varepsilon}-u\right| \leq$ $C h^{\frac{1}{d+1}}$ is thus ensured by choosing $\varepsilon=h^{\frac{1}{d+1}}$.

### 1.3 Rander geometry

Rander metrics are asymmetric metric $\{$, defined as the sum of a symmetric Riemannian part and of an anti-symmetric linear part 46. A Rander metric is thus described by a tensor field $\mathcal{M} \in C^{0}\left(\bar{\Omega}, \mathrm{~S}^{++}(\mathbb{E})\right)$, and a co-vector field $\hat{\eta} \in C^{0}\left(\bar{\Omega}, \mathbb{E}^{*}\right)$, subject to a compatibility condition:

$$
\begin{equation*}
\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}):=\|\dot{\mathbf{p}}\|_{\mathcal{M}(\mathbf{p})}+\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{p}}\rangle, \quad \text { where }\|\hat{\eta}(\mathbf{p})\|_{\mathcal{M}(\mathbf{p})^{-1}}<1 \tag{21}
\end{equation*}
$$

The smallness constraint $\left(21\right.$, right) ensures the positivity of the asymmetric norm $\mathcal{F}_{\mathbf{p}}(\cdot)$. The distance induced by a Rander metric is oriented: $d_{\mathcal{F}}(\mathbf{p}, \mathbf{q}) \neq d_{\mathcal{F}}(\mathbf{q}, \mathbf{p})$ in general.
Proposition 1.9. The distance $u$ to $\partial \Omega$, see (11), is the unique the viscosity solution to the inhomogeneous static first order HJB PDE

$$
\begin{equation*}
\|\mathrm{d} u(\mathbf{p})-\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}=1, \quad \text { where } \mathcal{D}(\mathbf{p})=\mathcal{M}(\mathbf{p})^{-1} \tag{22}
\end{equation*}
$$

for all $\mathbf{p} \in \Omega$, and $u=0$ on $\partial \Omega$.
Proof. It is known that $u$ obeys the eikonal $\operatorname{PDE} \mathcal{F}_{\mathbf{p}}^{*}(\mathrm{~d} u(\mathbf{p}))=1$, where $\mathcal{F}^{*}$ is the dual metric, see (12). Now for any $\hat{\mathbf{p}} \in \mathbb{E}^{*}$ observe the sequence of equivalences:

$$
\begin{aligned}
\mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}})=1 & \Leftrightarrow \exists \dot{\mathbf{p}} \in \mathbb{E} \backslash\{0\}, \hat{\mathbf{p}}=\mathrm{d} \mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}) \\
& \Leftrightarrow \exists \dot{\mathbf{p}} \in \mathbb{E} \backslash\{0\}, \hat{\mathbf{p}}=\mathcal{M}(\mathbf{p}) \dot{\mathbf{p}} /\|\dot{\mathbf{p}}\|_{\mathcal{M}(\mathbf{p})}+\hat{\eta}(\mathbf{p}) \\
& \Leftrightarrow\|\hat{\mathbf{p}}-\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}=1 .
\end{aligned}
$$

The first equivalence follows from convex duality $\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}})=\sup \left\{\langle\hat{\mathbf{p}}, \dot{\mathbf{p}}\rangle ; \mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}})=1\right\}$ and the enveloppe theorem, and the second one from the explicit expression (21) of $\mathcal{F}$.
Theorem 1.10. Let $(\mathcal{M}, \hat{\eta})$ be a Rander metric, and for all $\mathbf{p} \in \bar{\Omega}$ let $\mathcal{D}(\mathbf{p}):=\mathcal{M}(\mathbf{p})^{-1}$ and $\left(\rho_{i}(\mathbf{p}), \dot{\mathbf{e}}_{i}(\mathbf{p})\right)_{i=1}^{d^{\prime}}$ be as in 15$)$. Then for any $h>0$ there exists a unique solution to $U_{h}: h \mathbb{L} \rightarrow \mathbb{R}$ to the following discrete problem: for all $\mathbf{p} \in \Omega_{h}$
$\sum \lambda_{i}(\mathbf{p}) \max \left\{0, U_{h}(\mathbf{p})-U_{h}\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right)+h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle, U_{h}(\mathbf{p})-U_{h}\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)-h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle\right\}^{2}=h^{2}$,
and $U_{h}(\mathbf{p})=0$ for all $\mathbf{p} \in \partial \Omega_{h}$. If in addition $\Omega$ obeys an exterior cone condition, and $\mathcal{M}$ and $\eta$ have Lipschitz regularity, then for some $C=C(\Omega, \mathcal{M}, \hat{\eta})$ one has for all $h>0$

$$
\max _{\mathbf{p} \in \Omega_{h}}\left|U_{h}(\mathbf{p})-u(\mathbf{p})\right| \leq C \sqrt{r_{*} h} .
$$

[^4]The discretized PDE (23) cannot be solved using the Fast-Marching algorithm, contrary to the Riemannian case (17) and sub-Riemannian case, because the expression (23) may depend on non-causal, negative finite differences $U(\mathbf{p})-U(\mathbf{p}+\dot{\mathbf{e}})<0$ when $\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}\rangle>0$, in contradiction with Definition 2.1. For moderate anisotropies, good results are nevertheless obtained using Adaptive Gauss Siedel Iteration (AGSI), see [9] and $\S 3$. Alternatively, in dimension $d=2$, Rander distances can be computed via the single pass semi-Lagrangian method [35].

## 2 Convergence in the Riemannian case

This section is devoted to the proof of Theorem 1.4, which contains two parts: a claim of wellposedness for the system of equations discretizing the Riemannian eikonal PDE, and an error analysis as the grid scale is refined. For that purpose, two general and classical results are stated in 82.1 , and later specialized in $\$ 2.2$ to the model of interest.

### 2.1 Two general results

We formally introduce the concepts of monotone and causal finite difference schemes, and present (reformulations of) two classical results. Theorem 2.2 states that monotone schemes possess a unique solution, in the spirit of [61, 56, 41] and under adequate assumptions, which can be efficiently computed under the additional assumption of causality. Theorem 2.3 introduces a strategy for the numerical analysis, referred to as the doubling of variables argument and adapted from [24].

Definition 2.1. A (finite differences) scheme on a finite set $X$ is a continuous map $\mathfrak{F}: X \times$ $\mathbb{R} \times \mathbb{R}^{X} \rightarrow \mathbb{R}$. The scheme is said:

- Monotone, iff $\mathfrak{F}$ is non-decreasing w.r.t. the second and (each of the) third variables.
- Causal, iff $\mathfrak{F}$ only depends on the positive part of the third variable.

To the scheme is associated a function $\mathbb{R}^{X} \rightarrow \mathbb{R}^{X}$ still (abusively) denoted by $\mathfrak{F}$, and defined by

$$
(\mathfrak{F} U)(\mathbf{x}):=\mathfrak{F}\left(\mathbf{x}, U(\mathbf{x}),(U(\mathbf{x})-U(\mathbf{y}))_{\mathbf{y} \in X}\right),
$$

for all $\mathbf{x} \in X, U \in \mathbb{R}^{X}$. A discrete map $U \in \mathbb{R}^{X}$ is called a sub- (resp. strict sub-, resp. super-, resp. strict super-) solution of the scheme $\mathfrak{F}$ iff $\mathfrak{F} U \leq 0$ (resp. $\mathfrak{F} U<0$, resp. $\mathfrak{F} U \geq 0$, resp. $\mathfrak{F} U>0$ ) pointwise on $X$.

When the scheme $\mathfrak{F}$ is obvious from context, we simply speak of a sub- and super-solution.
Theorem 2.2 (Solving monotone schemes). Let $\mathfrak{F}$ be a monotone scheme on a finite set $X$ s.t.
(i) There exists a sub-solution $U^{-}$and a super-solution $U^{+}$to the scheme $\mathfrak{F}$.
(ii) Any super-solution to $\mathfrak{F}$ is the limit of a sequence of strict super-solutions.

Then there exists a unique solution $U \in \mathbb{R}^{X}$ to $\mathfrak{F} U=0$, and it satisfies $U^{-} \leq U \leq U^{+}$. If in addition the scheme is causal, then this solution can be obtained via the Dynamic-Programming algorithm, also called Dijkstra or Fast-Marching, with complexity $\mathcal{O}(M \ln N)$ where

$$
\begin{equation*}
N=\#(X), \quad M=\#(\{(x, y) \in X \times X ; \mathfrak{F} U(\mathbf{x}) \text { depends on } U(\mathbf{y})\}) . \tag{24}
\end{equation*}
$$

Proof. We provide for completeness the proof of existence and uniqueness, see [41] and [4] for closely related arguments in the discrete and continuous settings respectively. In contrast we refer to [61, 56] for the description of the fast marching algorithm.

Proof of uniqueness, via the comparison principle. Let $U^{+}$be a strict super-solution, and $U^{-}$ a sub-solution. Let $\mathbf{p} \in X$ be such that $U^{-}(\mathbf{p})-U^{+}(\mathbf{p})$ is maximal, so that $U^{-}(\mathbf{p})-U^{-}(\mathbf{q}) \geq$ $U^{+}(\mathbf{p})-U^{+}(\mathbf{q})$ for any $\mathbf{q} \in X$. Assuming for contradiction that $U^{-}(\mathbf{p}) \geq U^{+}(\mathbf{p})$ we obtain $0 \geq \mathfrak{F} U^{-}(\mathbf{p}) \geq \mathfrak{F} U^{+}(\mathbf{p})>0$ by monotony of the scheme and definition of a sub- and strict super-solution. This is a contradiction, hence $U^{-} \leq U^{+}$. Next using assumption (ii) we obtain that $U^{-} \leq U^{+}$still holds for any sub-solution $U^{-}$and any (possibly non-strict) super solution $U^{+}$. The uniqueness of the solution to $\mathfrak{F} U=0$ follows.

Proof of existence, by Perron's method. We prove that $U: X \rightarrow \mathbb{R}$, defined by $U(\mathbf{p}):=$ $\sup \{\tilde{U}(\mathbf{p}) ; \tilde{U}$ sub-solution $\}$ for all $\mathbf{p} \in X$, is a solution to the scheme $\mathfrak{F}$. By the previous argument one has $\tilde{U} \leq U^{+}$for any sub-solution $\tilde{U}$, and up to considering the sub-solution $\max \left\{\tilde{U}, U^{-}\right\}$we may as well assume $\tilde{U} \geq U^{-}$. Thus $U^{-} \leq \tilde{U} \leq U^{+}$and therefore $U^{-} \leq U \leq U^{+}$ by taking the pointwise supremum. Consider an arbitrary $\mathbf{p} \in X$, and let $\tilde{U}$ be a sub-solution such that $U(\mathbf{p})=\tilde{U}(\mathbf{p})$, which exists by continuity of $\mathfrak{F}$ and a compactness argument. By construction $U \geq \tilde{U}$, hence $\mathfrak{F} U(\mathbf{p}) \leq \mathfrak{F} \tilde{U}(\mathbf{p}) \leq 0$ by monotony of the scheme, hence $U$ is a sub-solution by arbitraryness of $\mathbf{p} \in X$. Furthermore, assume for contradiction that there exists $\mathbf{p}_{0} \in X$ such that $\mathfrak{F} U\left(\mathbf{p}_{0}\right)<0$, and define $U_{\varepsilon}\left(\mathbf{p}_{0}\right):=U\left(\mathbf{p}_{0}\right)+\varepsilon$ and $U_{\varepsilon}(\mathbf{p}):=U(\mathbf{p})$ for all $\mathbf{p} \in X \backslash\left\{\mathbf{p}_{0}\right\}$. Then $U_{\varepsilon}$ is a sub-solution for any sufficiently small $\varepsilon>0$, by monotony and continuity of the scheme $\mathfrak{F}$, thus $U\left(\mathbf{p}_{0}\right) \geq U_{\varepsilon}\left(\mathbf{p}_{0}\right)$ by construction which is a contradiction. Finally we obtain $\mathfrak{F} U=0$ identically on $X$, as announced.

The following result is a general strategy for proving convergence rates for discretizations of first order HJB PDEs, adapted from [24]. For completeness, the proof is presented in §C. The cartesian grid $h \mathbb{L}$ could be replaced with an arbitrary $h$-net of $E$, in other words a discrete set such that union of all balls of radius $h$ centered at the points of this set covers $\mathbb{E}$.

Theorem 2.3 (Doubling of variables argument). Let $u: \mathbb{E} \rightarrow \mathbb{R}$ be supported on a bounded domain $\bar{\Omega}$ and $C_{\text {Lip }}^{u}$-Lipschitz, and let $U_{h}: h \mathbb{L} \rightarrow \mathbb{R}$ be supported on $\Omega_{h}:=\Omega \cap h \mathbb{L}$. Given $\lambda \in[1 / 2,1[$ and $\delta>0$, define
$\bar{M}_{\lambda, \delta}:=\sup _{(\mathbf{p}, \mathbf{q}) \in h \mathbb{L} \times \mathbb{E}} \lambda U_{h}(\mathbf{p})-u(\mathbf{q})-\frac{1}{2 \delta}\|\mathbf{p}-\mathbf{q}\|^{2}, \quad \widetilde{M}_{\lambda, \delta}:=\sup _{(\mathbf{p}, \mathbf{q}) \in h \mathbb{L} \times \mathbb{E}} \lambda u(\mathbf{q})-U_{h}(\mathbf{p})-\frac{1}{2 \delta}\|\mathbf{p}-\mathbf{q}\|^{2}$,
and denote by $(\overline{\mathbf{p}}, \overline{\mathbf{q}}),(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \in(h \mathbb{L}) \times \mathbb{E}$ the point pairs where the maxima are respectively attained. Then $\|\mathbf{p}-\mathbf{q}\| \leq 4 C_{\mathrm{Lip}}^{u} \delta$. Assume furthermore that for some $C_{\mathrm{bd}}^{U}$, $C_{\mathrm{bd}}^{U}$ and $c_{\mathrm{bd}}^{U} \geq 4 C_{\mathrm{Lip}} \delta$ the following holds:
(i) None of the two maximal pairs $(\overline{\mathbf{p}}, \overline{\mathbf{q}})$ and $(\tilde{\mathbf{q}}, \tilde{\mathbf{q}})$ belongs to $\Omega_{h} \times \Omega$.
(ii) $\left|U_{h}(\mathbf{p})\right| \leq C_{\mathrm{bd}}^{U} d_{\partial \Omega}(\mathbf{p})+C_{\mathrm{bd}}^{\prime U} h$, for all $\mathbf{p} \in \Omega_{h}$ such that $d_{\partial \Omega}(\mathbf{p}) \leq c_{\mathrm{bd}}$.

Then one has with $C_{0}=4 C_{\mathrm{Lip}}^{u} \max \left\{C_{\mathrm{Lip}}^{u}, C_{\mathrm{bd}}^{U}\right\}$

$$
\begin{equation*}
\max _{\mathbf{p} \in \Omega_{h}}\left|u(\mathbf{p})-U_{h}(\mathbf{p})\right| \leq 2\left(C_{0} \delta+C_{\mathrm{bd}}^{\prime U} h+(1-\lambda) \max _{\bar{\Omega}}|u|\right) \tag{25}
\end{equation*}
$$

When applying Theorem 2.3 , to one of our specific models, Property (i) follows directly from the consistency of the discretization, while Property (ii) requires to establish a discrete counterpart of short-time local controllability at the boundary. In the Riemannian case, Property
(i) is established in Lemma 2.7, and Property (ii) in Proposition 2.10. Some adaptations of these arguments are required to establish property (ii) in the sub-Riemannian case, see Proposition A.7, and property (i) in the Rander case, see Lemma B. 3 .

Explicit expressions of the constants $C_{\mathrm{Lip}}^{u}, C_{\mathrm{bd}}^{U}, c_{\mathrm{bd}}^{U}$, are provided in terms of the model parameters $\mathcal{M}, \Omega$. The constant $C_{\mathrm{bd}}^{\prime U}$, also given explicitly, depends linearly on the stencil maximal radius: $C_{\mathrm{bd}}^{\prime U}=C_{\mathrm{bd}}^{\prime \prime U} r_{*}$, and property (i) is shown to hold provided

$$
\lambda \leq 1-C_{1} \delta-C_{2} r_{*} \frac{h}{\delta},
$$

where $C_{1}, C_{2}$ are again explicit constants depending only on $\mathcal{M}, \Omega$. Choosing $\lambda$ equal to this upper bound, and defining $\delta=\sqrt{r_{*} h}$, one gets the error estimate

$$
\begin{equation*}
\max _{\mathbf{p} \in \Omega_{h}}\left|u(\mathbf{p})-U_{h}(\mathbf{p})\right| \leq 2 C_{\mathrm{bd}}^{\prime \prime U} r_{*} h+2\left(C_{0}+\left(C_{1}+C_{2}\right)\|u\|_{\infty}\right) \sqrt{r_{*} h} . \tag{26}
\end{equation*}
$$

as announced in Theorems 1.4, 1.9 and 1.11 for Riemannian, sub-Riemannian, and Rander metrics respectively.

### 2.2 Application to the Riemannian case

We establish Theorem 1.4, on the discretization of Riemannian exit time problems, by specializing the general results of 2.1 . For that purpose we consider a discretization scheme $\mathfrak{F}_{h}$, on the finite domain $\Omega_{h}$, see (3), of the following form: for any $U: \Omega_{h} \rightarrow \mathbb{R}$ and any $\mathbf{p} \in \Omega_{h}$

$$
\begin{equation*}
\left(\mathfrak{F}_{h} U(\mathbf{p})\right)^{2}:=h^{-2} \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p}) \max \left\{0, U(\mathbf{p})-U\left(\mathbf{p}+\dot{\mathbf{e}}_{i}\right), U(\mathbf{p})-U\left(\mathbf{p}-\dot{\mathbf{e}}_{i}\right)\right\}^{2}, \tag{27}
\end{equation*}
$$

where $U$ is extended by zero on $h \mathbb{L} \backslash \Omega_{h}$. The next proposition implies, by Theorem 2.2 , the existence and uniqueness of a solution to the equation $\mathfrak{F}_{h} U-1 \equiv 0$, and the applicability of the Fast-Marching algorithm to compute it, as announced in Theorem 1.4 .

Recall that $r_{*}$ is the maximal stencil radius, as defined in (16). The square root of the largest eigenvalue among all tensors of a tensor field $\mathcal{M} \in C^{0}\left(\bar{\Omega}, \mathrm{~S}^{++}(\mathbb{E})\right)$ is denoted by

$$
\begin{equation*}
\lambda^{*}(\mathcal{M}):=\max _{\mathbf{p} \in \bar{\Omega},\|\dot{\mathbf{p}}\|=1}\|\dot{\mathbf{p}}\|_{\mathcal{M}(\mathbf{p})} . \tag{28}
\end{equation*}
$$

Proposition 2.4. Let $\mathcal{M} \in C^{0}\left(\bar{\Omega}, \mathrm{~S}^{++}(\mathbb{E})\right)$ be a Riemannian metric, and for all $\mathbf{p} \in \bar{\Omega}$ let $\mathcal{D}(\mathbf{p}):=\mathcal{M}(\mathbf{p})^{-1}$ and $\left(\rho_{i}(\mathbf{p}), \dot{\mathbf{e}}_{i}(\mathbf{p})\right)_{i=1}^{d^{\prime}}$ be as in 15). Then the scheme $\mathfrak{F}_{h}$ defined by (27) is monotone and causal. In addition:
(i) The null map $U=0$ satisfies $\mathfrak{F}_{h} U \equiv 0$, hence is a sub-solution to the scheme $\mathfrak{F}_{h}-1$.
(ii) Let $R>0$ be such that $\Omega$ is contained in the ball of radius $R-h r_{*}$ and centered at the origin, and let $U(\mathbf{p}):=R-\|\mathbf{p}\|$, for all $\mathbf{p} \in \Omega_{h}$. Then for all $\lambda \geq 0$, and all $\mathbf{p} \in \Omega_{h}$

$$
\mathfrak{F}_{h} U(\mathbf{p}) \geq\|\mathbf{p} /\| \mathbf{p}\| \|_{\mathcal{D}(\mathbf{p})},
$$

where $\mathbf{p} /\|\mathbf{p}\|$ can be replaced with an arbitrary unit vector in the case $\mathbf{p}=0$. As a result, $\lambda U$ is a super-solution to the scheme $\mathfrak{F}_{h}-1$ for any $\lambda \geq \lambda^{*}(\mathcal{M})$.
(iii) If $U$ is a super-solution to $\mathfrak{F}_{h}-1$, then $(1+\varepsilon) U$ is a strict super-solution for any $\varepsilon>0$.

Proof. The monotony and causality of the scheme $\mathfrak{F}_{h}$ immediately follow from its expression (27). Point (i) is trivial, and point (iii) follows from the homogeneity property $\mathfrak{F}_{h}(\lambda U)=\lambda \mathfrak{F}_{h} U$. In the rest of this proof, the point $\mathbf{p} \in \Omega$ is regarded both as a vector in $\mathbb{E}$ and as a co-vector in $\mathbb{E}^{*}$, thanks to the euclidean structure of $\mathbb{E}$. For point (ii), we obtain by convexity of the euclidean norm, for any $\mathbf{p}, \dot{\mathbf{e}} \in \mathbb{E}$.

$$
U(\mathbf{p})-U(\mathbf{p}+\dot{\mathbf{e}})=\|\mathbf{p}+\dot{\mathbf{e}}\|-\|\mathbf{p}\| \geq\langle\mathbf{p} /\|\mathbf{p}\|, \dot{\mathbf{e}}\rangle,
$$

where $\mathbf{p} /\|\mathbf{p}\|$ can be replaced with any unit vector if $\mathbf{p}=0$. Hence for all $\mathbf{p} \in \Omega_{h}$, as announced,

$$
\left(\mathfrak{F}_{h} U(\mathbf{p})\right)^{2} \geq \frac{1}{\|\mathbf{p}\|^{2}} \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p})\left\langle\mathbf{p}, \dot{\mathbf{e}}_{i}\right\rangle^{2}=\frac{\|\mathbf{p}\|_{\mathcal{D}(\mathbf{p})}^{2}}{\|\mathbf{p}\|^{2}}
$$

Finally $\mathfrak{F}_{h}(\lambda U)(\mathbf{p}) \geq 1$ if $\lambda \geq \lambda^{*}(\mathcal{M})$ by the 1 -homogeneity of $\mathfrak{F}$, and the observation that the least eigenvalue of $\mathcal{D}(\mathbf{p})$ is inverse of the largest eigenvalue of $\mathcal{M}(\mathbf{p})$.

In the rest of this section, we establish the properties required to apply the doubling of variables argument, Theorem 2.3, to prove the second part of Theorem 1.4. The following proposition immediately implies that the exit time value function, denoted hereafter by $u$, is $C_{\text {Lip }}^{u}$-Lipschitz, with $C_{\text {Lip }}^{u}:=\lambda^{*}(\mathcal{M})$.
Proposition 2.5. Let $\Omega \subseteq \mathbb{E}$ be an arbitrary bounded domain, equipped with a metric $\mathcal{F}$ : $\bar{\Omega} \times \mathbb{E} \rightarrow \mathbb{R}$ such that $\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}) \leq C_{0}\|\dot{\mathbf{p}}\|$ for any $\mathbf{p} \in \bar{\Omega}, \dot{\mathbf{p}} \in \mathbb{E}$. Then the distance $u$ from $\partial \Omega$ is $C_{0}$-Lipschitz.

Proof. Let $\mathbf{p}, \mathbf{q} \in \bar{\Omega}$, and let us prove that $u(\mathbf{q}) \leq u(\mathbf{p})+C_{0}\|\mathbf{p}-\mathbf{q}\|$. Let $\gamma:[0,1] \rightarrow \mathbb{E}$ be the parametrization of the line segment $[\mathbf{p}, \mathbf{q}]$ at constant speed. If this segment intersects $\partial \Omega$, then denoting $T \in[0,1]$ the largest time such that $\gamma(T) \in \partial \Omega$ one has $u(\mathbf{q}) \leq \operatorname{Length}_{\mathcal{F}}\left(\gamma_{[T, 1]}\right) \leq$ $C_{0}\|\mathbf{p}-\mathbf{q}\|$ as announced. Otherwise, denoting by $\tilde{\gamma}$ a minimal path from $\partial \Omega$ to $\mathbf{p}$ one has by path concatenation $u(\mathbf{q}) \leq$ length $_{\mathcal{F}}(\tilde{\gamma})+\operatorname{length}_{\mathcal{F}}(\gamma) \leq u(\mathbf{p})+C_{0}\|\mathbf{p}-\mathbf{q}\|$ as announced.

The rest of this section is split into two parts, devoted to proving assumptions (i) and (ii) of Theorem 2.3 , the doubling of variables argument; in the context of the Riemannian discretization scheme, thus by (26) concluding the proof of Theorem $\sqrt{1.4}$. For that purpose, we adopt the notations and other assumptions of Theorems 1.4 and 2.3. and we denote by $U_{h}$ the solution to the scheme $\mathfrak{F}_{h}-1$. In particular $\lambda \in[1 / 2,1[$ and $\delta>0$ are parameters from Theorem 2.3, and $(\overline{\mathbf{p}}, \overline{\mathbf{q}}),(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \in \mathbb{E} \times h \mathbb{L}$ are points pairs where the maxima $\bar{M}_{\lambda, \delta}, \widetilde{M}_{\lambda, \delta}$ are attained.

Establishing assumption (i) of Theorem 2.3. Our first lemma is a direct application of the definition of sub- and super-solutions of HJB PDEs and monotone discretization schemes.

Lemma 2.6. Let $\overline{\mathbf{w}}:=(\overline{\mathbf{p}}-\overline{\mathbf{q}}) / \delta$, and let $\bar{U}(\mathbf{p}):=\langle\overline{\mathbf{w}}, \mathbf{p}\rangle+\frac{1}{2 \delta}\|\mathbf{p}-\overline{\mathbf{p}}\|^{2}$ for all $\mathbf{p} \in h \mathbb{L}$. Then

$$
\begin{equation*}
\mathfrak{F}_{h} \bar{U}(\overline{\mathbf{p}}) \leq \lambda, \quad\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{q}})} \geq 1 \tag{29}
\end{equation*}
$$

Let $\tilde{\mathbf{w}}:=(\tilde{\mathbf{p}}-\tilde{\mathbf{q}}) / \delta$, and $\operatorname{let} \tilde{U}(\mathbf{p}):=\langle\tilde{\mathbf{w}}, \mathbf{p}\rangle-\frac{1}{2 \delta}\|\mathbf{p}-\tilde{\mathbf{p}}\|^{2}$ for all $\mathbf{p} \in h \mathbb{L}$. Then

$$
\begin{equation*}
\tilde{F}_{h} \tilde{U}(\tilde{\mathbf{p}}) \geq 1, \quad\|\tilde{\mathbf{w}}\|_{\mathcal{D}(\tilde{\mathbf{q}})} \leq \lambda . \tag{30}
\end{equation*}
$$

Here and below we regard $\overline{\mathbf{w}}$ and $\tilde{\mathbf{w}}$ as co-vectors, using the euclidean structure of $\mathbb{E}$.

Proof. Note that the scheme $\mathfrak{F}_{h}$ is here (slightly abusively) applied to the functions $\bar{U}$, $\tilde{U}$, which are non-zero over $h \mathbb{L} \backslash \Omega_{h}$. We focus on the proof of (29), the case of (30) being similar. By definition of $\bar{M}_{\lambda, \delta}$, the function

$$
\mathbf{p} \in h \mathbb{L} \mapsto \lambda U_{h}(\mathbf{p})-u(\overline{\mathbf{q}})-\frac{1}{2 \delta}\|\mathbf{p}-\overline{\mathbf{q}}\|^{2}=\lambda U_{h}(\mathbf{p})-\bar{U}(\mathbf{p})-K
$$

attains its maximum at $\overline{\mathbf{p}}$, where $K=\bar{u}(\overline{\mathbf{q}})-\frac{1}{2}\langle\overline{\mathbf{w}}, \overline{\mathbf{p}}+\overline{\mathbf{q}}\rangle$ is independent of the variable $\mathbf{p}$. Hence for all $\mathbf{p} \in h \mathbb{L}$
$\lambda U_{h}(\overline{\mathbf{p}})-\bar{U}(\overline{\mathbf{p}})-K \geq \lambda U_{h}(\mathbf{p})-\bar{U}(\mathbf{p})-K, \quad$ equivalently $U_{h}(\overline{\mathbf{p}})-U_{h}(\mathbf{p}) \geq U(\overline{\mathbf{p}}) / \lambda-\bar{U}(\mathbf{p}) / \lambda$.
By monotony of the scheme $\mathfrak{F}_{h}$, see Definition 2.1, we obtain $\mathfrak{F}_{h}(\bar{U} / \lambda)(\overline{\mathbf{p}}) \leq \mathfrak{F}_{h} U_{h}(\overline{\mathbf{q}})=1$, hence 29. left) by the homogeneity of $\mathfrak{F}_{h}$. Likewise, defining $\bar{u}(\mathbf{q}):=\langle\overline{\mathbf{w}}, \mathbf{q}\rangle-\frac{1}{2 \delta}\|\mathbf{q}-\overline{\mathbf{q}}\|^{2}$ for all $\mathbf{q} \in \mathbb{E}$, the function

$$
\mathbf{q} \in \mathbb{E} \mapsto \lambda U_{h}(\overline{\mathbf{p}})-u(\mathbf{q})-\frac{1}{2 \delta}\|\mathbf{p}-\overline{\mathbf{q}}\|^{2}=\bar{u}(\mathbf{q})-u(\mathbf{q})-K^{\prime}
$$

attains its minimum at $\overline{\mathbf{q}}$, where $K^{\prime}$ is the adequate constant. Since $u$ is a (super-)solution to the PDE (11), this implies $1 \leq\|\mathrm{d} \bar{u}(\overline{\mathbf{q}})\|_{\mathcal{D}(\overline{\mathbf{q}})}=\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{q}})}$, which concludes the proof.

The following lemma assumes for contradiction that $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \in \Omega_{h} \times \Omega$ and obtains estimates contradicting Lemma 2.6 established above, provided $\lambda$ is above a certain bound, which is assumed in the following. Therefore, arguing by contradiction, one must have $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \notin \Omega_{h} \times \Omega$, and likewise $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \notin \Omega_{h} \times \Omega$ by a similar argument, which establishes assumption (i) of Theorem 2.3. Let $C_{\text {Lip }}^{\mathcal{D}}$ be a constant such that for all $\mathbf{p}, \mathbf{q} \in \bar{\Omega}$ and all $\hat{\mathbf{p}} \in \mathbb{E}^{*}$

$$
\begin{equation*}
\left|\|\hat{\mathbf{p}}\|_{\mathcal{D}(\mathbf{p})}-\|\hat{\mathbf{p}}\|_{\mathcal{D}(\mathbf{q})}\right| \leq C_{\mathrm{Lip}}^{\mathcal{D}}\|\mathbf{p}-\mathbf{q}\|\|\hat{\mathbf{p}}\| . \tag{31}
\end{equation*}
$$

Such a constant exists by the Lipschitz regularity of the metric $\mathcal{M}$, assumed in Theorem 1.4
Lemma 2.7. Assume that $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \in \Omega_{h} \times \Omega$ and define $\overline{\mathbf{w}}$ and $\bar{U}$ as in Lemma 2.6. Then

$$
\begin{equation*}
\left|\mathfrak{F}_{h} \bar{U}(\overline{\mathbf{p}})-\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})}\right| \leq C_{1} r_{*} \frac{h}{\delta}, \quad \quad\left|\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{q}})}-\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})}\right| \leq C_{2} \delta \tag{32}
\end{equation*}
$$

with $C_{1}:=\lambda^{*}(\mathcal{D}) \sqrt{d}$ and $C_{2}:=\left(4 C_{\text {Lip }}^{u}\right)^{2} C_{\text {Lip }}^{\mathcal{D}}$. This contradicts 29) unless $\lambda \geq 1-C_{1} r_{*} h / \delta-C_{2} \delta$. The same estimates and conclusion hold for $\tilde{\mathbf{w}}$ and $\tilde{U}$ if $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \in \Omega_{h} \times \Omega$.

Proof. We focus on the case of $(\overline{\mathbf{p}}, \overline{\mathbf{q}})$, the second case of ( $\tilde{\mathbf{p}}, \tilde{\mathbf{q}})$ being similar, and begin with the proof of (32, left) which contains the key technical points. By definition of the quadratic function $\bar{U}$, one has

$$
\begin{equation*}
\max \left\{0, \bar{U}(\overline{\mathbf{p}})-\bar{U}\left(\overline{\mathbf{p}}+h \dot{\mathbf{e}}_{i}\right), \bar{U}(\overline{\mathbf{p}})-\bar{U}\left(\overline{\mathbf{p}}-h \dot{\mathbf{e}}_{i}\right)\right\}=h\left|\left\langle\overline{\mathbf{w}}, \dot{\mathbf{e}}_{i}\right\rangle\right|+\frac{h^{2}}{\delta}\left\|\dot{\mathbf{e}}_{i}\right\|^{2}, \tag{33}
\end{equation*}
$$

for any $1 \leq i \leq d^{\prime}$, where $\left(\rho_{i}(\overline{\mathbf{p}}), \dot{\mathbf{e}}_{i}\right)_{i=1}^{d^{\prime}}$ are the weights and offsets of the discretization scheme at $\overline{\mathbf{p}}$, see (15). Denote by $w, e \in \mathbb{R}^{d^{\prime}}$ the vectors of components, respectively, $w_{i}:=\left|\left\langle\overline{\mathbf{w}}, \dot{\mathbf{e}}_{i}\right\rangle\right|$, and $e_{i}:=\left\|\dot{\mathbf{e}}_{i}\right\|^{2}$, for all $1 \leq i \leq d^{\prime}$. Introduce also the semi-norm $\|z\|_{\overline{\mathbf{p}}}:=\sqrt{\sum_{i=1}^{d^{\prime}} \rho_{i}(\overline{\mathbf{p}}) z_{i}^{2}}$ on $\mathbb{R}^{d^{d^{\prime}}}$. Then by (33) and the consistency relation (6) one has

$$
\mathfrak{F}_{h} \bar{U}(\overline{\mathbf{p}})=\left\|w+\frac{h}{\delta} e\right\|_{\overline{\mathbf{p}}}, \quad\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})}=\|w\|_{\overline{\mathbf{p}}}
$$

and therefore $\left|\mathfrak{F}_{h} \bar{U}(\overline{\mathbf{p}})-\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})}\right| \leq \frac{h}{\delta}\|e\|_{\overline{\mathbf{p}}}$ by the triangular inequality. Finally observe that

$$
\|e\|_{\mathbf{p}}^{2}=\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\overline{\mathbf{p}})\left(\left\|\dot{\mathbf{e}}_{i}\right\|^{2}\right)^{2} \leq \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\overline{\mathbf{p}})\left\|\dot{\mathbf{e}}_{i}\right\|^{2} \max _{1 \leq i \leq d^{\prime}}\left\|\dot{\mathbf{e}}_{i}\right\|^{2}=\operatorname{Tr}(\mathcal{D}(\overline{\mathbf{p}})) r_{*}^{2},
$$

where $\operatorname{Tr}$ denotes the trace of a matrix. The announced result (32, left) then follows from $\operatorname{Tr}(\mathcal{D}(\mathbf{p})) \leq d\left(\lambda^{*}(\mathbf{p})\right)^{2}$.

The second estimate (32, right) follows from the Lipschitz regularity of the metric (31), together with the upper bound $\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\| \leq 4 C_{\text {Lip }}^{u} \delta$ established in Theorem 2.3, which implies $\|\overline{\mathbf{w}}\| \leq 4 C_{\text {Lip }}^{u}$. Combining these estimates with Lemma 2.6 yields

$$
\lambda+C_{1} r_{*} h / \delta \geq\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})} \geq 1-C_{2} \delta,
$$

which implies the announced lower bound for $\lambda$. The same estimates can be derived in the second case, and with Lemma 2.6 they imply $1-C_{1} r_{*} h / \delta \leq\|\tilde{\mathbf{w}}\|_{\mathcal{D}(\tilde{\mathbf{p}})} \leq \lambda+C_{2} \delta$ which yields the same lower bound for $\lambda$.

Establishing assumption (ii) of Theorem 2.3. The reste of this section is devoted to proving, in Proposition 2.10, an estimate on the growth of the discrete solution $U_{h}$ close to $\partial \Omega$, thus implying assumption (ii) of Theorem 2.3. A natural strategy would be to prove a global Lipschitz type estimate for the discrete solution $U_{h}$, as in e.g. [9, 34, but unfortunately the assumptions of Theorem 1.4 are too weak for that purpose, and actually we cannot exclude a staggered grid effect (never observed in practice) far from $\partial \Omega$. Instead, the idea underlying our proof is to construct from any point in $\mathbf{p}_{0} \in \Omega_{h}$ sufficiently close to $\partial \Omega$, a short chain of neighbors $\mathbf{p}_{1}, \cdots, \mathbf{p}_{n}$ ending in $\partial \Omega_{h}$ and connected by offsets of the numerical scheme $\mathbf{p}_{i+1}=$ $\mathbf{p}_{i}+h \dot{\mathbf{e}}_{i}\left(\mathbf{p}_{i}\right)$ which associated weights $\rho_{i}\left(\mathbf{p}_{i}\right)$ are positively bounded below. This chain is the discrete counterpart of a short time local control to the to boundary [3].

Our first step is to provide a precise definition to the exterior cone condition assumed in the statement of Theorem 1.4.

Definition 2.8 (Exterior cone condition). The domain $\Omega \subseteq \mathbb{E}$ obeys an exterior cone condition iff there exists constants $C_{\Omega}$ and $c_{\Omega}>0$ such that for all $h \leq c_{\Omega}$,

$$
\forall \mathbf{p} \in \partial \Omega, \exists \mathbf{q} \in B\left(\mathbf{p}, C_{\Omega} h\right), \text { such that } B(\mathbf{q}, h) \subseteq \mathbb{E} \backslash \Omega,
$$

where $B(\mathbf{q}, h)$ denotes the open ball of center $\mathbf{q}$ and radius $h$.
The next technical lemma compares the euclidean norm with its first order Taylor expansion.
Lemma 2.9. For any $\mathbf{p}, \dot{\mathbf{e}} \in \mathbb{E}$ with $\mathbf{p} \neq 0$, one has $\|\mathbf{p}+\dot{\mathbf{e}}\| \leq\|\mathbf{p}\|+(\mathbf{p} \cdot \dot{\mathbf{e}}) /\|\mathbf{p}\|+\|\dot{\mathbf{e}}\|^{2} /(2\|\mathbf{p}\|)$.
Proof. Multiplying both sides by $\|\mathbf{p}\|$ and rearranging terms the statement is found equivalent to $\|\mathbf{p}\|\|\mathbf{p}+\dot{\mathbf{e}}\| \leq \frac{1}{2}\left(\|\mathbf{p}\|^{2}+\|\mathbf{p}+\dot{\mathbf{e}}\|^{2}\right)$, equivalently to $0 \leq(\|\mathbf{p}\|-\|\mathbf{p}+\dot{\mathbf{e}}\|)^{2}$ which holds true.

In the following proposition, we let $\operatorname{Cond}(\mathcal{D}):=\max \{\operatorname{Cond}(\mathcal{D}(\mathbf{p})) ; \mathbf{p} \in \bar{\Omega}\}$, where the condition number of a symmetric matrix is defined in (7).

Proposition 2.10. Let $\mathbf{p} \in \Omega_{h}$, and let $\mathbf{q} \in \mathbb{E}$ be such that $\|\mathbf{p}-\mathbf{q}\| \geq C_{0} r_{*} h$, with $C_{0}:=$ $\operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}$. Then there exists $1 \leq i \leq d^{\prime}$ and a sign $s \in\{-1,1\}$ such that

$$
\begin{equation*}
U_{h}(\mathbf{p}) \leq U_{h}\left(\mathbf{p}+h s \dot{\mathbf{e}}_{i}\right)+h C_{1}\left\|\dot{\mathbf{e}}_{i}\right\|, \quad \text { and }\left\|\mathbf{p}+h s \dot{\mathbf{e}}_{i}-\mathbf{q}\right\| \leq\|\mathbf{p}-\mathbf{q}\|-h c_{2}\left\|\dot{\mathbf{e}}_{i}\right\|, \tag{34}
\end{equation*}
$$

with $C_{1}:=\lambda^{*}(\mathcal{M}) \sqrt{d^{\prime}}, c_{2}:=1 /\left(2 \operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}\right)$. This implies assumption (ii) of Theorem 2.3. with the constants $C_{\mathrm{bd}}^{U}=C_{1} / c_{2}, C_{\mathrm{bd}}^{\prime U}:=C_{\mathrm{bd}}^{U} C_{\Omega} C_{0} r_{*}$, and $c_{\mathrm{bd}}=+\infty$.

Proof. Denote by $\left(\lambda_{*}\right)^{2}$ and $\left(\lambda^{*}\right)^{2}$ the smallest and largest eigenvalue of $\mathcal{D}(\mathbf{p})$ respectively. Let $\mathbf{n}:=(\mathbf{q}-\mathbf{p}) /\|\mathbf{q}-\mathbf{p}\|$, regarded as a co-vector thanks to the euclidean structure of $\mathbb{E}$. Then

$$
\begin{equation*}
\lambda_{*}^{2} \leq\|\mathbf{n}\|_{\mathcal{D}(\mathbf{p})}^{2}=\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p})\left\langle\mathbf{n}, \dot{\mathbf{e}}_{i}\right\rangle^{2} . \tag{35}
\end{equation*}
$$

Fix $1 \leq i \leq d^{\prime}$ such that $\rho_{i}(\mathbf{p})\left\langle\mathbf{n}, \dot{\mathbf{e}}_{i}\right\rangle^{2} \geq \lambda_{*}^{2} / d^{\prime}$. Denote $\rho^{2}:=\rho_{i}(\mathbf{p})$, and $\dot{\mathbf{e}}:=s \dot{\mathbf{e}}_{i}$ where $s$ is the sign of $\left\langle\mathbf{n}, \dot{\mathbf{e}}_{i}\right\rangle$. One has, using that $\rho^{2} \dot{\mathbf{e}} \otimes \dot{\mathbf{e}} \preceq \mathcal{D}(\mathbf{p}) \preceq\left(\lambda^{*}\right)^{2}$ Id for the second inquality

$$
\begin{equation*}
\rho\langle\mathbf{n}, \dot{\mathbf{e}}\rangle \geq \lambda_{*} / \sqrt{d^{\prime}} \quad \rho\|\dot{\mathbf{e}}\| \leq \lambda^{*} \tag{36}
\end{equation*}
$$

By definition of the discretization scheme (27) one has $h^{-2} \rho^{2} \max \left\{0, U_{h}(\mathbf{p})-U_{h}(\mathbf{p}+\dot{\mathbf{e}})\right\}^{2} \leq$ $\left(\widetilde{F}_{h} U_{h}(\mathbf{p})\right)^{2}=1$, hence using (36, left) we obtain (34, left):

$$
\begin{equation*}
U(\mathbf{p})-U(\mathbf{p}+\dot{\mathbf{e}}) \leq \frac{h}{\rho} \leq h \frac{\|\dot{\mathbf{e}}\| \sqrt{d^{\prime}}}{\lambda_{*}} . \tag{37}
\end{equation*}
$$

By (36) and $\operatorname{Cond}(\mathcal{D}) \geq \lambda^{*} / \lambda_{*}$ one has $\langle\mathbf{n}, \dot{\mathbf{e}}\rangle \geq\|\dot{\mathbf{e}}\| /\left(\operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}\right)$. Denote $\mathbf{r}:=(\mathbf{p}-\mathbf{q}) / h$ and observe that $\|\mathbf{r}\| \geq C_{0} r_{*} \geq \operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}\|\dot{\mathbf{e}}\|$, by assumption and by definition of the max stencil radius $r_{*}$, see (16). Using Lemma 2.9 we obtain (34, right):

$$
\|\mathbf{r}+\dot{\mathbf{e}}\| \leq\|\mathbf{r}\|-\langle\mathbf{n}, \dot{\mathbf{e}}\rangle+\frac{\|\dot{\mathbf{e}}\|^{2}}{2\|\mathbf{r}\|} \leq\|\mathbf{r}\|-\frac{\|\dot{\mathbf{e}}\|}{\operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}}+\frac{\|\dot{\mathbf{e}}\|}{2 C_{0}} \frac{\|\dot{\mathbf{e}}\|}{r_{*}} \leq\|\mathbf{r}\|-\frac{\|\dot{\mathbf{e}}\|}{2 \operatorname{Cond}(\mathcal{D}) \sqrt{d^{\prime}}} .
$$

Finally, we conclude the proof of assumption (ii). Let $\mathbf{p}_{0} \in \Omega_{h}$. Let $\mathbf{q}_{*} \in \partial \Omega$ be the closest point to $\mathbf{p}_{0}$, and let $\mathbf{q} \in B\left(\mathbf{q}_{*}, C_{\Omega} C_{0} r_{*} h\right)$ be such that $B\left(\mathbf{q}, C_{0} r_{*} h\right) \subseteq \mathbb{E} \backslash \Omega$. By the above argument, there exists a finite sequence of points $\mathbf{p}_{1}, \cdots, \mathbf{p}_{k-1} \in \Omega_{h}, \mathbf{p}_{k} \in \partial \Omega_{h}$, such that $U\left(\mathbf{p}_{i}\right) \leq U\left(\mathbf{p}_{i+1}\right)+C_{1} \delta_{i}$ and $\left\|\mathbf{p}_{i+1}-\mathbf{q}\right\| \leq\left\|\mathbf{p}_{i}-\mathbf{q}\right\|-c_{2} \delta_{i}$, denoting $\delta_{i}:=\left\|\mathbf{p}_{i+1}-\mathbf{p}_{i}\right\|$, for all $0 \leq i<k$. Since $U\left(\mathbf{p}_{k}\right)=0$ we obtain $U\left(\mathbf{p}_{0}\right) \leq C_{1}\left(\delta_{0}+\cdots+\delta_{k-1}\right)$, and since $\left\|\mathbf{p}_{k}-\mathbf{q}\right\| \geq 0$ we obtain $c_{2}\left(\delta_{0}+\cdots+\delta_{k-1}\right) \leq\left\|\mathbf{p}_{0}-\mathbf{q}\right\| \leq\left\|\mathbf{p}_{0}-\mathbf{q}_{*}\right\|+C_{\Omega} C_{0} r_{*} h$. Hence finally, as announced

$$
U\left(\mathbf{p}_{0}\right) \leq\left(C_{1} / c_{2}\right)\left(d_{\partial \Omega}\left(\mathbf{p}_{0}\right)+C_{\Omega} C_{0} r_{*} h\right) .
$$

## 3 Numerical results

We illustrate the numerical methods introduced in this paper with a series of numerical experiments, involving Riemannian, sub-Riemannian and Rander metrics, in $83.1,43.2$ and 83.3 respectively. Open source numerical codes for the Riemannian and sub-Riemannian models 5 are available on the author's webpag $\AA^{6}$.

### 3.1 Riemannian examples

We validate our algorithm on several two and three dimensional Riemannian test cases, which are split into two groups. The problems of the first group - related to differential geometry and seismic imaging - feature smooth Riemannian metrics with pronounced yet bounded anisotropy, and accuracy is the main concern. The problems of the second group - related to tubular structure segmentation in medical image data - feature discontinuous Riemannian metrics and extreme anisotropies, so that robustness is the main concern.

[^5]Smooth Riemannian metrics. The first test, two dimensional and introduced in 65], is the computation of the distance from the origin on a parametric surface w.r.t. the Riemannian metric induced by the Euclidean metric on $\mathbb{R}^{3}$. The surface is described by the height map

$$
z(x, y):=(3 / 4) \sin (3 \pi x) \sin (3 \pi y),
$$

hence the Riemannian metric is $\mathcal{M}(x, y)=\operatorname{Id}+\nabla z(x, y) \nabla z(x, y)^{\mathrm{T}}$, which maximum condition number (7) is $\approx 5.1$. The parametrization domain is the unit square $[-0.5,0.5]^{2}$, rotated ${ }^{7}$ by the angle $\pi / 6$.

The second test, two dimensional and introduced in [57], is inspired by seismic imaging applications. Note that, admittedly, some authors have claimed that this application requires more complex types of Finslerian anisotropies [48. The Riemannian metric tensor $\mathcal{M}(x, y)$ has eigenvector $(1,(\pi / 2) \cos (4 \pi x))$ with eigenvalue $0.8^{-2}$. The second eigenvalue is $0.2^{-2}$, hence the condition number is 4 . The parametrization domain is $[-0.5,0.5]^{2}$, and the distance is computed from the origin.

The third test, introduced here for the first time, extends the seismic imaging inspired second test to three dimensions. The Riemannian metric tensor $\mathcal{M}(x, y, z)$ has eigenvector $(\cos (3 \pi(x+y)), \sin (3 \pi(2 x-y)), 0.5)$, with eigenvalue $0.2^{-2}$. The two other eigenvalues are equal to $0.8^{-2}$, hence the condition number is 4 . The domain is $[-0.5,0.5]^{3}$ and the distance is computed from the origin $(0,0,0)$.

The level sets of the distance maps associated to these three tests, computed numerically using the FM-VR1, and a number of the corresponding minimal geodesics, are presented in Figure 2. The accuracy and computation time in the two dimensional test cases are compared in 3.4 with several alternative numerical methods.

Discontinuous Riemannian metrics with extreme anisotropy. Anisotropic fast marching methods have shown their relevance for image segmentation methods based on minimal paths [7, 15]. In these applications, the metric often varies quickly, if not discontinuously, both in orientation and aspect ratio. For instance, the Riemannian metric is often designed to favor paths which remain close and tangent to a collection of thin tubular structures in the image.

We present two numerical experiments inspired by these applications, in two and three dimensions, which first appeared in [7] and [34] respectively. The Riemannian metric is Euclidean (identity matrix) except in the neighborhood of a curve $\Gamma$ embedded in the domain, where the metric is extremely anisotropic, with eigenvalues $\left(1,1 / 100^{2}\right)$ or $\left(1,1,1 / 50^{2}\right)$ in the two and three dimensional experiments respectively, and the tangent vector to the curve $\Gamma$ is an eigenvector for the small eigenvalue. See 34 for a complete description. The level sets of the distance maps associated to these two test cases, computed numerically using the FM-VR1, and some of the corresponding minimal geodesics, are illustrated on Figure 3 .

In these extreme test cases, the FM-VR1 behaves particularly well in terms of CPU time and accuracy, comparably to the FM-LBR which similarly uses an adaptive discretization strategy. In contrast, iterative numerical methods such as the AGSI 9, and fast marching methods based on less sophisticated stencil constructions such as [1] have be shown to fail on these types of benchmarks [7, 34].

[^6]

Figure 2: Numerics, illustrating $\$ 3.1$, for the three test cases involving smooth anisotropic Riemannian metrics. Top: level lines of the numerically computed distance map from the domain center. Bottom: backtracked minimal geodesics, from the domain center to various endpoints. Left: geodesic distance on a parametric surface, $283^{2}$ grid, 0.1 s CPU time. Center: test inspired by seismic data analysis, $193^{2}$ grid, 0.04 s CPU time. Right: likewise in 3D, on a $101^{3}$ grid, CPU time 5.02s. All CPU times measured on a 2.7Ghz laptop using a single thread.


Figure 3: Numerics, illustrating $\$ 3.1$, for test cases involving discontinuous and extremely anisotropic Riemannian metrics, in dimension $d \in\{2,3\}$, inspired by applications to tubular structure segmentation in medical data [7]. (i,iii) Level lines of the numerically computed distance from the image center. (ii,iv) Backtracked minimal geodesics, from the image center to various endpoints. Because of the chosen metric, these paths are concatenations of straight lines (euclidean geodesics), and of portions adjacent to a given spiraling curve $\Gamma$ (along which paths are favored by the metric anisotropy). Left: $201 \times 201$ grid, 0.03s CPU time. Right $201 \times 201 \times 272$ grid, 25 s CPU time.

### 3.2 Sub-Riemannian models

We consider several sub-Riemannian models, posed on the configuration space $\mathbb{M}:=\mathbb{R}^{d} \times \mathbb{S}^{d-1}$ of positions and orientations. Such configurations are denoted $\mathbf{p}=(\mathbf{x}, \mathbf{n}) \in \mathbb{M}$, and their tangent vectors $\dot{\mathbf{p}}=(\dot{\mathbf{x}}, \dot{\mathbf{n}}) \in T_{\mathbf{p}} \mathbb{M}$. For the simplicity of the exposition we regard $\mathbf{n}$ as a genuine unit vector in $\mathbb{R}^{d}$, so that $\langle\mathbf{n}, \dot{\mathbf{n}}\rangle=0$, although our numerical implementation relies on an angular parametrisation of the sphere $\mathbb{S}^{d-1}$.

We choose to describe the sub-Riemannian models of interest via an approximating family of Riemannian metrics $\left(\mathcal{M}_{\varepsilon}\right)_{\varepsilon>0}$, where $\varepsilon$ is a relaxation parameter. The orthogonal projection of a vector $\dot{\mathbf{x}} \in \mathbb{R}^{d}$, onto the hyperplane orthogonal to a given unit vector $\mathbf{n} \in \mathbb{S}^{d-1}$, is denoted

$$
P_{\mathbf{n}}(\dot{\mathrm{x}}):=\dot{\mathrm{x}}-\langle\mathbf{n}, \dot{\mathbf{x}}\rangle \mathbf{n} .
$$

The Reeds-Shepp model. This model, defined on $\mathbb{R}^{2} \times \mathbb{S}^{1}$, describes a car ${ }^{8}$, which state is described by a position $\mathbf{x} \in \mathbb{R}^{2}$, and an orientation $\mathbf{n}=(\cos \theta, \sin \theta) \in \mathbb{S}^{1}$, see Figure 4. The car can move forward and backward, rotate in either direction, but not move sideways. This model also plays a central role in the study of the visual cortex organisation and function, in which case it is referred to as the Petitot-Citti-Sarti model [43]. Recenty, data-driven variants of the Reeds-Shepp model and of its higher dimensional counterparts have been considered for tubular structure segmentation in medical image data [5, 23]. The Riemannian relaxations of this model's metric read, for each $\varepsilon>0$,

$$
\begin{equation*}
\|\dot{\mathbf{p}}\|_{\mathcal{M}_{\varepsilon}(\mathbf{p})}^{2}:=\mathcal{S}(\mathbf{p})^{-2}\left(\langle\mathbf{n}, \dot{\mathbf{x}}\rangle^{2}+\varepsilon^{-2}\left\|P_{\mathbf{n}}(\dot{\mathbf{x}})\right\|^{2}+\xi^{2}\|\dot{\mathbf{n}}\|^{2}\right) \tag{38}
\end{equation*}
$$

where $\mathcal{S}: \mathbb{M} \rightarrow] 0, \infty[$ is a point dependent speed function, with physical units [length]/[time], and $\xi$ is a parameter which has the dimension [length] of a radius of curvature. Parameters $\mathcal{S}$ and $\xi$ may be constant or variable over the domain, possibly dictated by the considered application in a data-driven manner [5, 23].

The Reeds-Shepp model is related to curvature penalization for the following reason: consider a smooth path $\mathbf{x}:[0, T] \rightarrow \mathbb{R}^{d}$, with non-vanishing speed. Then there exists a unique $\mathbf{n}:[0, T] \rightarrow$ $\mathbb{S}^{d-1}$, up to a global change of sign, such that the lifted path $t \in[0, T] \mapsto \gamma(t)=(\mathbf{x}(t), \mathbf{n}(t))$ has finite length with respect to the sub-Riemannian metric $\mathcal{M}_{0}$. Indeed, one must set $\mathbf{n}(t):=$ $\pm \dot{\mathbf{x}}(t) /\|\dot{\mathbf{x}}(t)\|$, so that $P_{\mathbf{n}(t)}(\dot{\mathbf{x}}(t))=0$ for all $t \in[0, T]$. Then, denoting by $\kappa(t):=\|\dot{\mathbf{n}}(t)\| /\|\dot{\mathbf{x}}(t)\|$ the curvature of the path $\mathbf{x}$, one obtains

$$
\int_{0}^{T}\|\dot{\gamma}(t)\|_{\mathcal{M}_{0}(t)} \mathrm{d} t=\int_{0}^{T} \sqrt{1+\xi^{2} \kappa(t)^{2}} \frac{\|\dot{\mathbf{x}}(t)\| \mathrm{d} t}{\mathcal{S}(\gamma(t))}
$$

Note that (contrary to what this discussion may suggest) the physical projections of geodesic paths for the sub-Riemannian metric $\mathcal{M}_{0}$ are only piecewise smooth typically, because they feature cusps, see Figures 4. 6, and the discussion in [23].

Some experiments involving two and three dimensional physical paths are presented in Figures 4, 5 and 6. Let us emphasize that we are here solving strongly anisotropic PDEs on three and five dimensional domains respectively. The control sets for the Reeds-Shepp model posed on $\mathbb{R}^{2} \times \mathbb{S}^{1}$ are illustrated on page ??. For the model posed on $\mathbb{R}^{3} \times \mathbb{S}^{2}$, the sphere $\mathbb{S}^{2}$ is parametrized using the Euler angles $\theta, \varphi) \mapsto(\cos \theta, \sin \theta \cos \varphi, \sin \theta \sin \varphi)$, from the flat domain $[0, \pi] \times[0,2 \pi]$ equipped with the adequate Riemannian metric and boundary conditions.

[^7]

Figure 4: Numerics for the Reeds-Shepp sub-Riemannian model, illustrating 3.2 (i) Explanation of the model: the parameter space is three dimensional, and the sub-Riemannian structure forces the car ground speed $\dot{\mathbf{x}}(t)$ to remain aligned with the direction $\mathbf{n}(t)=(\cos \theta(t), \sin \theta(t))$ defined by the third coordinate. In the next sub-figures, only the planar projection $\mathbf{x}(t)$ of the minimal paths $\gamma(t)=(\mathbf{x}(t), \mathbf{n}(t)) \in \mathbb{R}^{2} \times \mathbb{S}^{1}$ is shown. (ii) Projections of minimal paths in $[-0.5,0.5]^{2} \times \mathbb{S}^{1}$, from $(0,0,0)$ to various endpoints with imposed orientation $\theta=\pi / 4$, with model parameters $\varepsilon=0.1, \xi=0.3$. (iii) Comparison of the numerically backtracked paths (solid) with those obtained using an high order ODE shooting method based on the Hamilton equations of geodesics (40) (dashed blue). (iv) Some minimal paths (black) for the Reeds-Shepp model in the presence of obstacles (grayed). Domain $[0,1]^{2} \times \mathbb{S}^{1}$ discretized on a $90^{2} \times 60$ grid, CPU time 0.36 s. Model parameters $\varepsilon=0.1$, and $\xi=0.4$. Some orientations, arbitrary, are imposed at the geodesics endpoints.

The numerical results are similar to those obtained in [23] using the semi-lagrangian FMLBR, but computation times are substantially smaller for the 5D test case, see the discussion in \$3.1, by a factor 5 typically for the five dimensional test cases. Figure 4 illustrates the spatial projections in $\mathbb{R}^{2}$ of the minimal geodesics associated with the classical Reeds-Shepp model posed on $\mathbb{R}^{2} \times \mathbb{S}^{1}$, with and without obstacles, and in the latter case a comparison with high accuracy solutions obtained with an ODE shooting method.

A variant related to torsion penalization. We introduce a new sub-Riemannian model, which relaxed metric is defined for all $\mathbf{p}=(\mathbf{x}, \mathbf{n}) \in \mathbb{M}$ and $\dot{\mathbf{p}}=(\dot{\mathbf{x}}, \dot{\mathbf{n}}) \in T_{\mathbf{p}} \mathbb{M}$ by

$$
\begin{equation*}
\|\dot{\mathbf{p}}\|_{\widetilde{\mathcal{M}}_{\varepsilon}(\mathbf{p})}^{2}:=\mathcal{S}(\mathbf{p})^{-2}\left(\langle\mathbf{n}, \dot{\mathbf{x}}\rangle^{2}+\varepsilon^{-2}\left\|P_{\mathbf{n}}(\dot{\mathbf{x}})\right\|^{2}+\xi^{2}\|\dot{\mathbf{n}}\|^{2}\right), \tag{39}
\end{equation*}
$$

where again $\mathcal{S}: \mathbb{M} \rightarrow] 0, \infty[$ is the speed function, and $\xi$ has the dimension of a length. The model (39) favors paths which are possibly non-smooth but are embedded in smooth surfaces, a property that is relevant for certain tasks in medical data segmentation 60. Indeed the physical velocity $\dot{\mathbf{x}}$ is constrained by the cost of $\varepsilon^{-2}\langle\mathbf{n}, \dot{\mathbf{x}}\rangle^{2}$ to remain (approximatedly if $\varepsilon>0$ ) in the plane orthogonal to the vector $\mathbf{n}$, which variation is itself controlled by the cost of $\|\dot{\mathbf{n}}\|^{2}$. Note also that the most natural way to lift a physical curve $\mathbf{x}:[0, T] \rightarrow \mathbb{R}^{3}$ into $\gamma=(\mathbf{x}, \mathbf{n}):[0, T] \rightarrow$ $\mathbb{R}^{3} \times \mathbb{S}^{2}$ obeying the orthogonality constraint $\langle\mathbf{x}(t), \mathbf{n}(t)\rangle=0$ for all $t \in[0, T]$, is to define $\mathbf{n}(t):=(\dot{\mathbf{x}}(t) \times \ddot{\mathbf{x}}(t)) /\|\dot{\mathbf{x}}(t) \times \ddot{\mathbf{x}}(t)\|$. Then denoting by $\tau(t):=\|\dot{\mathbf{n}}(t)\| /\|\dot{\mathbf{x}}(t)\|$ the torsion of the path $\mathbf{x}$ one obtains

$$
\int_{0}^{T}\|\dot{\gamma}(t)\|_{\widetilde{\mathcal{M}}_{0}(t)} \mathrm{d} t=\int_{0}^{T} \sqrt{1+\xi^{2} \tau(t)^{2}} \frac{\|\dot{\mathbf{x}}(t)\| \mathrm{d} t}{\mathcal{S}(\gamma(t))}
$$

Nevertheless our model is only related to torsion penalization, and not equivalent to it, because there exists other lifts $\gamma=(\mathbf{x}, \mathbf{n}):[0, T] \rightarrow \mathbb{R}^{d} \times \mathbb{S}^{d-1}$ of the curve $\mathbf{x}:[0, T] \rightarrow \mathbb{R}^{d}$ obeying


Figure 5: Numerics, illustrating $\$ 3.2$, for sub-riemannian models posed on $\mathbb{R}^{3} \times \mathbb{S}^{2}$. Test case inspired by 3D tubular structure segmentation. (i) Contour plot of the speed function, which is high in the neighborhood of two curves of small curvature and small torsion respectively. (ii, iii) Position $\mathbf{x}(t)$ and orientation $\mathbf{n}(t)$, respectively, of the sub-Riemannian geodesics $\gamma(t)=$ ( $\mathbf{x}(t), \mathbf{n}(t))$ extracted with the Reeds-Shepp model (purple), and its torsion related variant (blue). Distinct paths are selected, along the appropriate vessel centerline. Domain $([0, \pi] \times[-1,1] \times$ $[-1,0.5]) \times \mathbb{P}^{2}$ discretized using a $(40 \times 20 \times 16) \times(5 \times 20)$ grid, $\varepsilon=0.2$. CPU time 6.6 s .
the required orthogonality constraint, and which energy could be smaller than the torsion based one.

On the experiment presented Figure 5 , the speed function $\left.\mathcal{S}: \mathbb{R}^{3} \times \mathbb{S}^{2} \rightarrow\right] 0, \infty[$ only depends on the physical position $\mathbf{x}$, and is small away from two curves $\Gamma_{1}, \Gamma_{2}$ of interest

$$
\mathcal{S}(\mathbf{x}, \mathbf{n}):=\max \left\{s, \exp \left(-\frac{\operatorname{dist}\left(\mathbf{x}, \Gamma_{1} \cup \Gamma_{2}\right)^{2}}{2 \sigma^{2}}\right)\right\}
$$

where $s=1 / 6$ and $\sigma=0.15$. The curves $\Gamma_{1}, \Gamma_{2}$ are parametrized by $t \in[0, \pi]$ as follows

$$
\gamma_{1}(t):=\left(t, \sin (t)^{2} \cos (4 t), 0\right), \quad \gamma_{2}(t):=\left(t, \sin (t)^{3} \cos (2 t), \sin (t)^{3} \sin (2 t)\right) .
$$

Hence $\Gamma_{1}$ has large curvature but no torsion, whereas $\Gamma_{2}$ has small curvature but some torsion. Using our anisotropic fast marching method, we compute the shortest path between the common endpoints $\mathbf{x}_{0}, \mathbf{x}_{1} \in \mathbb{R}^{3}$ of these curves, among all possible tangent directions $\mathbf{n}_{0}, \mathbf{n}_{1} \in \mathbb{S}^{2}$ at these endpoints. Figure 5 shows the level lines of the cost function $\mathcal{S}$, and the minimal geodesics corresponding to the two models (38) and (39), numerically computed using the FM-VR1. As could be expected, the torsion related model selects a path along $\Gamma_{1}$, whereas the Reeds-Shepp model selects a path along $\Gamma_{2}$.

Validation of the approach. We present on Figure 6 two empirical validations of our numerical approach to computing globally optimal geodesics for the five dimensional Reeds-Shepp model and its torsion related variant. We first show that the sub-Riemannian constraint, of collinearity $P_{\mathbf{n}}(\dot{\mathbf{x}})$ and orthogonality $\langle\mathbf{n}, \dot{\mathbf{x}}\rangle$ are approximately satisfied, despite their relaxation in (38) and (39), with $\varepsilon=0.1$. We then compare the obtained minimal paths with solutions of the Hamilton equations of geodesics

$$
\begin{equation*}
\frac{d \mathbf{p}}{d t}=-\frac{\partial \mathcal{H}}{\partial \hat{\mathbf{p}}}, \quad \frac{d \hat{\mathbf{p}}}{d t}=\frac{\partial \mathcal{H}}{\partial \mathbf{p}}, \quad \text { where } \mathcal{H}(\mathbf{p}, \hat{\mathbf{p}}):=\frac{1}{2}\left\langle\hat{\mathbf{p}}, \mathcal{D}_{0}(\mathbf{p}) \hat{\mathbf{p}}\right\rangle \tag{40}
\end{equation*}
$$

Here $\mathcal{D}_{0}$ denotes the inverse tensor to the sub-Riemannian metric (38) or (39), which is well defined when $\varepsilon=0$, in contrast to $\mathcal{M}_{0}$ itself. This ODE is solved using a fourth order RungeKutta method, and the initial conditions are adjusted using a Newton method to meet the desired endpoints.


Figure 6: Validation of the FM-VR1 numerical method applied to the sub-Riemannian ReedsShepp models (i, ii), and its torsion related variant (iii, iv), posed on $\mathbb{R}^{3} \times \mathbb{S}^{2}$. Parameters $\varepsilon=0.1$, $\xi=0.5$, constant cost function $\mathcal{S}$. (i, iii) The angular component $\mathbf{n}(t)$ of the minimal paths $\gamma(t)=(\mathbf{x}(t), \mathbf{n}(t)) \in \mathbb{R}^{3} \times \mathbb{S}^{2}$, illustrated with arrows, satisfies (approximately, as expected) the sub-Riemannian constraint (i) $\mathbf{n}(t) \times \dot{\mathbf{x}}(t)=0$ or (iii) $\langle\mathbf{n}(t), \dot{\mathbf{x}}(t)\rangle=0$. (ii, iv) Comparison of the backtracked geodesics (black) with the results of an ODE shooting method (blue) based on Hamilton's equations of geodesics.

### 3.3 Rander models

We consider some instances of Zermelo's navigation problem, which models a boat navigating on a body of water [12, 2]. The (motor) boat is capable of a certain maximum speed, in any direction, and inertia is not taken into account. However the boat is also subject to a drift due to current or wind, which in our experiments is variable over the domain (and constant in time). The goal is to move from one given point to another in minimal time.

Formally, let us denote by $\Omega \subseteq \mathbb{R}^{d}$ the domain, denote by $\dot{\eta}: \bar{\Omega} \rightarrow \mathbb{R}^{d}$ the drift, and assume that the maximum speed is 1 (in the Euclidean norm). The boat starts from anywhere on $\partial \Omega$, and all points of $\Omega$ are regarded as potential target points. The objective is thus to find for each $\mathbf{p} \in \Omega$ the minimal time $T=u(\mathbf{p}) \geq 0$ for which there exists a path $\gamma:[0, T] \rightarrow \bar{\Omega}$ such that $\gamma(0) \in \partial \Omega, \gamma(T)=\mathbf{p}$, and

$$
\|\dot{\gamma}(t)-\dot{\eta}(\gamma(t))\| \leq 1 \quad \text { a.e. } t \in[0, T] .
$$

We assume that $\|\dot{\eta}(\mathbf{x})\|<1$ for all $\mathbf{x} \in \mathbb{R}^{d}$, otherwise the system would not be locally controllable (the drift speed being larger than the maximum boat speed). Following [2] we reformulate this problem as a shortest path problem with respect to a Rander metric, of parameters ( $\left.\mathcal{D}^{-1}, \hat{\eta}\right)$ specified in the next proposition, see also $(22)$.

Proposition 3.1. The value function for this problem obeys $\|\mathrm{d} u(\mathbf{p})-\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}=1$ on $\Omega$ in the sense of viscosity solutions, and $u=0$ on $\partial \Omega$, where

$$
\mathcal{D}(\mathbf{p}):=\left(1-\|\dot{\eta}(\mathbf{p})\|^{2}\right)(1-\dot{\eta}(\mathbf{p}) \otimes \dot{\eta}(\mathbf{p})), \quad \quad \hat{\eta}(\mathbf{p})=\frac{-\dot{\eta}(\mathbf{p})}{1-\|\dot{\eta}(\mathbf{p})\|^{2}}
$$

Proof. The value function differential $\hat{\mathbf{p}}=\mathrm{d} u(\mathbf{p})$, where defined, obeys the equivalent constraints

$$
\sup _{\dot{\mathbf{v}} \in B}\langle\hat{\mathbf{p}}, \dot{\mathbf{v}}+\dot{\eta}\rangle \leq 1 \Leftrightarrow\|\hat{\mathbf{p}}\|+\langle\hat{\mathbf{p}}, \dot{\eta}\rangle \leq 1 \Leftrightarrow\|\hat{\mathbf{p}}\|^{2} \leq(1-\langle\hat{\mathbf{p}}, \dot{\eta}\rangle)^{2} \Leftrightarrow\|\hat{\mathbf{p}}-\hat{\eta}\|_{\mathcal{D}}^{2} \leq 1,
$$

where $B$ denotes the euclidean unit ball. The dependency of $\dot{\eta}, \hat{\eta}, \mathcal{D}$ to the base point $\mathbf{p}$ was omitted for readability. The leftmost identity follows from Bellman's optimality principle. The first equivalence is trivial, the second equivalence follows from the impossibility of $\|\hat{\mathbf{p}}\|=\langle\hat{\mathbf{p}}, \dot{\eta}\rangle-1$ (since $\|\dot{\eta}\|<1$ ), and the third equivalence results from direct computations using e.g. that $\mathcal{D} \hat{\eta}=-\left(1-\|\dot{\eta}\|^{2}\right) \dot{\eta}$.


Figure 7: Numerics, Illustrating $\$ 3.3$, of two instances of Zermelo's navigation problem in dimension two and three. Level lines of the distance map from the domain center (minimal travel time), and minimal geodesics to various endpoints, computed with the variant of the FM-VR1 adapted to Rander metrics. (i, ii) Grid size $201 \times 201$, CPU time 0.14s. (iii, iv) Grid size $101^{3}$, CPU time 7.8s.

|  | FM-VR1 |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1st order | 2nd order | FM-LBR | FM-8 | FE | MAOUM |  |
|  | Embedded surface distance test, $293 \times 293$ grid |  |  |  |  |  |  |
| CPU time | $0.10^{*}$ | $0.12^{*}$ | 0.20 | 0.21 | 1.44 | 1.31 |  |
| $L^{\infty}$ error | 5.8 | $0.22^{*}$ | 5.52 | 12.5 | 9.45 | 8.56 |  |
| $L^{1}$ error | 1.6 | 0.066 | 1.46 | 3.42 | 2.51 | 2.52 |  |
|  | Seismic inspired test, |  |  |  |  |  |  |
|  | $193 \times 193$ grid |  |  |  |  |  |  |
| CPU time | $0.042^{*}$ | $0.048^{*}$ | 0.076 | 0.079 | 0.77 | 0.36 |  |
| $L^{\infty}$ error | 4.5 | $0.15^{*}$ | 2.90 | 3.03 | 3.67 | 7.66 |  |
| $L^{1}$ error | 1.5 | 0.056 | 1.03 | 1.30 | 1.40 | 2.3 |  |

Figure 8: Comparison of the CPU time and accuracy of the proposed FM-VR1 with several alternatives in two Riemannian test cases, see $\$ 3.1$ and $\$ 3.4$. All errors multiplied by 100 for readability. Asterix *, see Remark 3.2.

We present a two dimensional experiment on $\Omega=] 0,1\left[^{2}\right.$, first introduced in [57], and a three dimensional generalization in $\Omega=] 0,1[3$. The drift has the explicit expression

$$
\dot{\eta}(\mathbf{x}):=\alpha \sin \left(4 \pi \mathbf{x}_{1}\right) \sin \left(4 \pi \mathbf{x}_{2}\right) \frac{\mathbf{x}}{\|\mathbf{x}\|} . \quad\left(\text { resp. } \dot{\eta}(\mathbf{x}):=\alpha \sin \left(4 \pi \mathbf{x}_{1}\right) \sin \left(4 \pi \mathbf{x}_{2}\right) \sin \left(4 \pi \mathbf{x}_{3}\right) \frac{\mathbf{x}}{\|\mathbf{x}\|} .\right)
$$

Recall that our numerical scheme (23) for eikonal equations of Rander type lacks the causality property, see Definition 2.1. The fast marching method is therefore not applicable, and we use instead Adaptive Gauss Siedel Iterations, in the spirit of [9]. Figure 7 illustrates the level lines of the distance map $u$, and some of the corresponding minimal geodesics, in the two and three dimensional test cases. The computation time and the $L^{\infty}$ and $L^{1}$ errors obtained with the two dimensional problem are presented in $\S \sqrt[3.4]{ }$, and compared with several alternative semilagrangian methods [9, 1, 35].

### 3.4 Comparison with alternative methods

We compare the accuracy and computation time of the FM-VR1 with several alternative methods proposed in the literature for solving anisotropic eikonal equations [61, 55, 1, 9, 34, 35]. As discussed in the introduction, these numerical methods can be divided into two groups: causal

|  | RD-VR1 |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1st order | 2nd order | FM-ASR | FM-8 | FE | MAOUM |
|  | Zermelo navigation problem, $285 \times 285$ grid |  |  |  |  |  |
| CPU time | $1.03\left(0.48^{*}\right)$ | $1.39\left(0.81^{*}\right)$ | 0.29 | 0.16 | 1.08 | 0.69 |
| $L^{\infty}$ error | 0.84 | 0.23 | 0.64 | 0.64 | 1.05 | 2.8 |
| $L^{1}$ error | 0.44 | 0.0095 | 0.13 | 0.11 | 0.41 | 0.17 |

Figure 9: Comparison of the CPU times and accuracy of the proposed RD-VR1 method with several alternatives in a test case involving a Rander metric, see $\$ 3.3$ and 3.4 All errors multiplied by 100 for readability. When testing second order accuracy, the seed point $(0,0)$ was replaced with a precomputed solution on the centered disk of radius 5 pixels. Asterix * first time using the AGSI as decribed in [9, second time with a variant which limits the front width to 10 pixels. See also remark 3.2 .


Figure 10: Numerical error as a function of gridsize for the two dimensional test cases. Second order convergence is achieved in the $L^{1}$ norm, and in the $L^{\infty}$ norm except for the Zermelo's navigation problem. See remark 3.2 for the experiment setup.

| First order |  |  |  |  |  |  | Second order |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Surface |  | Seismic |  | Zermelo |  |  | Surface |  | Seismic |  | Zermelo |  |
| n | $L^{1}$ | $L^{\infty}$ | $L^{1}$ | $L^{\infty}$ | $L^{1}$ | $L^{\infty}$ | n | $L^{1}$ | $L^{\infty}$ | $L^{1}$ | $L^{\infty}$ | $L^{1}$ | $L^{\infty}$ |
| 51 | 7.73 | 24.9 | 5.54 | 15.4 | 2.21 | 4.63 | 51 | 6.60 | 17.9 | 2.28 | 14.3 | 0.474 | 0.953 |
| 101 | 4.65 | 15.6 | 2.87 | 8.42 | 1.20 | 2.42 | 101 | 0.946 | 2.98 | 0.373 | 1.22 | 0.0887 | 0.657 |
| 201 | 2.40 | 8.28 | 1.46 | 4.35 | 0.62 | 1.21 | 201 | 0.149 | 0.527 | 0.050 | 0.142 | 0.0214 | 0.319 |
| 401 | 1.15 | 4.23 | 0.692 | 2.11 | 0.303 | 0.575 | 401 | 0.034 | 0.124 | 0.012 | 0.034 | 0.0044 | 0.162 |
| $\sigma$ | 1.06 | 0.97 | 1.08 | 1.04 | 1.05 | 1.08 | $\sigma$ | 2.12 | 2.09 | 2.10 | 2.07 | 2.28 | 0.98 |

Figure 11: Numerical error observed with the proposed schemes, for the three test cases, for several grid sizes $n \times n$. Last line: exponent such that $\operatorname{err}(n) \approx n^{-\sigma}$, obtained using $n \in$ $\{201,401\}$. All errors multiplied by 100 for readability.


Figure 12: (I,II,III) Numerical error using the second order scheme for the three tests, with resolution $n=273,273$ and 375 respectively. The numerical error is $\mathcal{O}\left(h^{2}\right)$, except on the pixels displayed white where it is typically $\mathcal{O}(h)$. These are located close to the boundary in the Riemannian cases, as well as along the cut locus and in other regions in the Rander case, which uses a non-causal scheme. (IV) Computational domain $[-0.5,0.5]^{2}$. For the tests involving second order finite differences, the exact solution is provided within the black disk of radius 0.05 , so as to mitigate the numerical error related to the point source singularity, and the gray region of width 0.05 is excluded from the computation of the $L^{\infty}$ error.
discretizations, solved via the single pass fast marching algorithm, and non-causal discretizations, solved via iterative methods. These groups are analogous, in the context of distance computation on graphs, to Dijkstra's algorithm and Bellman-Ford's algorithm respectively, in which case the counterpart of causality is the positivity of the edge weights, see $\$ 1$.

The test cases, taken from [55, are two dimensional and involve smooth metrics of Riemannian and Rander type with pronounced anisotropy. The good performance of the introduced discretization is confirmed, in terms of accuracy and CPU time, although it does not outperform e.g. the previous adaptive semi-Lagrangian schemes [34, 35] of the author.

Before proceeding, we would like to attract the attention of the reader to other qualities of the FM-VR1 scheme introduced in this paper, which are not put into light by this specific benchmark. Indeed, our discretization reliably handles much more extreme anisotropies than those considered here, see $\$ 3.1$, to the point that we can approximate sub-Riemannian metrics \$3.2. This is made possible by its full adaptivity, local and anisotropic. In addition, the Eulerian nature of the FM-VR1 makes its particularly simple to implement and cheap numerically, especially as dimension increases, in comparison with semi-Lagrangian methods which require to handle the complex combinatorics of the polyhedral neighborhood of each point, see (2). The numerical cost reduction is particularly evident in the five dimensional sub-Riemannian experiments, see $\$ 3.2$, which run approximately five times faster ${ }^{9}$ using the FM-VR1 than with the semi-Lagrangian implementation described in [23].

This benchmarks presented here extend previous works of the author published in [35, 34].

## Causal discretizations.

- Fast Marching using Voronoi's first reduction (FM-VR1), introduced in this paper.
- Fast-Marching using Lattice Basis Reduction (FM-LBR), and Fast-Marching using Adaptive Stencil Refinement (FM-ASR), introduced by the author in [34] and [35]. Like the FM-VR1, these are a single pass methods, which require a cartesian grid, and achieve their efficiency by the use of adaptive stencils built using techniques from lattice geometry. In contrast with the FM-VR1, these are semi-Lagrangian discretizations. The FM-LBR

[^8]applies to two and three dimensional Riemannian metrics, while the FM-ASR applies to two dimensional Finslerian metrics.

- The Monotone Acceptance Ordered Upwind method (MAOUM) [1] is a single pass semiLagrangian method using adaptive stencils. It differs from the FM-LBR and FM-ASR by its less sophisticated stencil construction, which produces large isotropic stencils, often at the expense of accuracy and complexity.
- Fast-Marching using 8-point stencils (FM-8) is the original semi-Lagrangian scheme 61] instantiated with non-adaptive stencils consisting of the 8 closest grid neighbors, see [29] for a three dimensional extension. This method is non-consistent for Riemannian metrics which condition number exceeds $1+\sqrt{2}$, because its stencils lack the acuteness property [55]. Hence convergence towards the continuous problem solution fails as the grid is refined in problem instances considered here. Nevertheless, the FM-8 is fast and its accuracy is surprisingly competitive at low grid sizes.

Non causal discretizations. We use Adaptive Gauss Siedel Iterations (AGSI ${ }^{10}$ ) to solve the following discretizations of the eikonal equation, which lack the causality property. We also report some computation times obtained by limiting the front width to 5 pixels.

- Rander Distances using Voronoi's First Reduction (RD-VR1), introduced in this paper.
- The Finite Element discretization (FE) of [9], a semi-Lagrangian discretization using nonadaptive stencils extracted from a triangulation of the domain, here by half-squares.

Remark 3.2 (Experiment setup, possible sources of bias). All CPU times obtained using a single thread. CPU times are empirical data, only indicative of general performance. CPU times for the FM-VR1 and RD-VR1 obtained on a 2.7 GHz core i7 laptop, whereas CPU times for the other methods were copied from previous works [34, 35] and obtained using a 2.4 GHz core 2 duo laptop. Numerical errors are with respect to a solution computed with the proposed algorithms on a fine grid of resolution $2001 \times 2001$, and then bilinearly interpolated. In the experiments involving the second order scheme, the boundary data is provided on a small disk of radius 0.05 rather than a single point source, and a boundary layer of width 0.05 is excluded from the $L^{\infty}$ error computation, as illustrated on Figure 12 (iv).

Discussion of accuracy. The numerical scheme introduced in this paper, the FM-VR1, belongs to the category of monotone discretizations of PDEs, also referred to as degenerate elliptic [41. This property ensures excellent stability and robustness properties, by the discrete comparison principle see Theorem 2.2, but limits the accuracy that can be achieved: monotone finite difference schemes are at most first order consistent for first order equations (such as the eikonal equation here considered), and at most second order for second order equations, see Theorem 4 in [41. In this paragraph, we discuss some techniques aiming at improving the accuracy of eikonal equation solvers, that have been proposed in the literature.

First, let us point out that the observed numerical error is $\mathcal{O}(h)$ in typical instances, see Figures 10 and 11 , despite the established rate of convergence being only $\mathcal{O}(\sqrt{h})$, see Theorem 1.4. In order to further reduce the numerical error, one possibility that we consider here is to

[^9]introduce second order accurate approximations of the directional derivatives of the solution, in the spirit of the Higher Accuracy Fast Marching Method (HAFMM) introduced in [54]:
\[

$$
\begin{equation*}
\langle\mathrm{d} u(\mathbf{p}), \dot{\mathbf{e}}\rangle=\frac{1}{h}(u(\mathbf{p}+h \dot{\mathbf{e}})-u(\mathbf{p}))-\frac{1}{2 h}(u(\mathbf{p})-2 u(\mathbf{p}+h \dot{\mathbf{e}})+u(\mathbf{p}+2 h \dot{\mathbf{e}}))+\mathcal{O}\left(h^{2}\right) . \tag{41}
\end{equation*}
$$

\]

This expression replaces the first order finite differences used in the FM-VR1 17), in some occurrences only and with enough caution and failsafe policies. Indeed, it breaks the monotony of the scheme. In order to preserve the convergence guarantees, the numerical method should fall back to the original first order upwind difference $(u(\mathbf{p}+h \dot{\mathbf{e}})-u(\mathbf{p})) / h$ whenever the second order corrective term exceeds the expected $\mathcal{O}(h)$ magnitud ${ }^{11}$. This ensures that the solution $U_{h}^{*}$ to the modified scheme obeys the original monotone scheme $\mathcal{F}_{h}$ defined in 17) up to an $\mathcal{O}(h)$ error: $\mathcal{F}_{h} U_{h}^{*}=1+\mathcal{O}(h)$. From this point, recalling the 2-homogeneity of $\mathcal{F}_{h}$ and using the discrete comparison principle, we obtain the global estimate $U_{h}^{*}=U_{h}+\mathcal{O}(h)$ which implies the convergence of $U_{h}^{*}$. The improved convergence rate of $U_{h}^{*}$, in comparison with the original scheme solution $U_{h}$, is only (often) observed numerically, and is not established by the analysis. We refer to [26] for a detailed analysis of related techniques in the (more complex) context of second order PDEs, in particular of the Monge-Ampere equation, using the concepts of filtered-scheme and of quasi-monotone discretization.

The solutions to eikonal equations typically have non-smooth singularities, which require a special treatment if second order convergence is to be achieved.

- The cut locus is the collection of points $\mathbf{p} \in \Omega$ reached by more than one minimal geodesic. The solution $u$ is non-differentiable on this set, denoted by $\mathcal{C}$ and which is typically a union of $(d-1)$-dimensional manifolds. Fortunately, because the minimal paths for the addressed problem do not cross $\mathcal{C}$ (geodesics loose optimality when they do so), the numerical error associated with this singularity does not excessively pollute the rest of the numerical solution.
An $\mathcal{O}(h)$ numerical error is in principle be expected at discretization points which stencil goes accross the non-differentiability set $\mathcal{C}$ of the solution. This is observed in the Rander test case, but surprisingly not in the Riemannian test cases where an $\mathcal{O}\left(h^{2}\right)$ only error is observed, see Figures 10 and 12. This empirical good surprise can be attributed to to the limited diffusivity of the chosen discretization, see the discussion in [56].
- Point source singularities. In the considered test cases, as in many applications, the null boundary condition $u\left(\mathbf{p}_{*}\right)=0$ is imposed at an isolated seed point, rather than on a full domain's boundary as in our theoretical results $\$ 1$. This produces a non-differentiable singularity with dominant term $\mathbf{q} \mapsto\left\|\mathbf{p}_{*}-\mathbf{q}\right\|_{\mathcal{M}\left(\mathbf{p}_{*}\right)}$ at the source point $\mathbf{p}_{*}$. Similar effects are encountered at the corners of obstacles [45]. Factoring techniques, additive or multiplicative, incorporate corrective terms in the numerical scheme based on the analytic expression of the singularity [33]. Another approach, used in our experiments on second order accuracy, is to pre-compute the solution $u$ on a finer grid in a small neighborhood of the singularity.
- Outflow boundary solutions. In the considered test cases, as in many applications, outflow boundary conditions are imposed on part of the domain's boundary, rather than null Dirichlet conditions as in our theoretical results \$1. As a result, some of the problem's

[^10]minimal paths are the (non-smooth) concatenation of usual geodesics and of parts of the domain's boundary. A similar issue occurs with non-constant Dirichlet boundary conditions unless they are $K$-Lipschitz with respect to the metric [9], with $K<1$. This phenomenon limits the smoothness of the problem solution, and raises numerical difficulties.

In our experiments using the second order scheme, we exclude a small band along $\partial \Omega$ from the $L^{\infty}$ error computation, where this phenomenon causes the error to be $\mathcal{O}(h)$ instead of $\mathcal{O}\left(h^{2}\right)$. This is sufficient in the Riemannian test cases, but not in the Rander case where the non-causality of the scheme causes the error to diffuse more inside the domain, see Figure 12 .

Finally, let us mention multi-stencil schemes [27] and other variants [22] of the fast marching method, which improve accuracy by the use of stencils wider than strictly necessary, featuring in particular the diagonal directions of the discretization grid.

## 4 Conclusion

In this paper, we introduced a new discretization of anisotropic eikonal equations on Cartesian grids. The discretization makes use of adaptive stencils built using a tool from discrete geometry: Voronoi's first reduction of quadratic forms.

A convergence proof is provided, with convergence rates, in the setting of Riemannian metrics, but also of sub-Riemannian and of asymmetric Rander metrics. Numerical experiments show that the method is particularly suitable for problems involving strong anisotropy, such as Riemannian tensor condition numbers of $\approx 10$ and more, and scales well up to dimension $d=5$.

Future directions of research include designing causal discretizations for non-symmetric Hamiltonians, addressing point sets more general than cartesian grids, either unstructured or obtained by gluing several grid patches of different scales, and developing applications to motion planning and image segmentation.

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## A Convergence in the sub-Riemannian case

This section is devoted to the proof of Theorem 1.9, namely the numerical analysis of the FM-VR1 scheme introduced in this paper in the sub-Riemannian (or pre-Riemannian) setting. The estimates of $u-u_{\varepsilon}$ and $u_{\varepsilon}-U_{h, \varepsilon}$, respectively related to the model relaxation and to its discretization, are presented separately in $\$$ A. 1 and $\$$ A.2. The arguments used to prove the Lipschitz regularity of the solution $u$ to the continuous problem, and to control the growth close to $\partial \Omega$ of the solution $U_{h, \varepsilon}$ to the relaxed and discretized problem, differ substantially from those used in the Riemannian case.

Before turning to these proofs, we recall two basic results on the regularity of the orthogonal projection onto a set, assumed respectively to be convex or to have a smooth boundary. Let $d_{B}(\mathbf{p}):=\min _{\mathbf{q} \in B}\|\mathbf{p}-\mathbf{q}\|$ denote the distance to a non-empty closed set $B \subseteq \mathbb{E}$, and let $\mathrm{P}_{B}(\mathbf{p})$ denote the minimizer $\mathbf{q} \in B$ for $d_{B}(\mathbf{p})$, when it is unique. The Haussdorff distance between two closed subsets of $\mathbb{E}$ is denoted $\mathcal{H}(\cdot, \cdot)$.

Proposition A.1. Let $\mathbf{p}, \mathbf{p}^{\prime} \in \mathbb{E}$, and let $B, B^{\prime} \subseteq \mathbb{E}$ be non-empty closed and convex. Then

$$
\begin{equation*}
\left\|\mathrm{P}_{B}(\mathbf{p})-\mathrm{P}_{B}\left(\mathbf{p}^{\prime}\right)\right\| \leq\left\|\mathbf{p}-\mathbf{p}^{\prime}\right\|, \quad\left\|\mathrm{P}_{B}(\mathbf{p})-\mathrm{P}_{B^{\prime}}(\mathbf{p})\right\| \leq \sqrt{\mathcal{H}\left(B, B^{\prime}\right)} \sqrt{d_{B}(\mathbf{p})+d_{B^{\prime}}(\mathbf{p})} \tag{42}
\end{equation*}
$$

Proof. The uniqueness and Lipschitz regularity of the projection onto a convex set (42, left), are extremely classical hence their proof is omitted. Proof of (42, right). Let $\mathbf{q}:=\mathrm{P}_{B}(\mathbf{p})$ and $\mathbf{q}^{\prime}:=\mathrm{P}_{B^{\prime}}(\mathbf{p})$. We first assume that $\mathbf{p} \neq \mathbf{q}$, and regard $\mathbf{q}-\mathbf{p}$ as a co-vector by Riez duality. Observe that $B$ is contained in the half space $H:=\{\mathbf{r} \in \mathbb{E} ;\langle\mathbf{q}-\mathbf{p}, \mathbf{r}-\mathbf{q}\rangle \geq 0\}$, hence

$$
\mathcal{H}\left(B, B^{\prime}\right) \geq d_{B}\left(\mathbf{q}^{\prime}\right) \geq d_{H}\left(\mathbf{q}^{\prime}\right)=\max \left\{0,\left\langle\frac{\mathbf{q}-\mathbf{p}}{\|\mathbf{q}-\mathbf{p}\|}, \mathbf{q}-\mathbf{q}^{\prime}\right\rangle\right\} \geq\left\langle\mathbf{q}-\mathbf{p}, \mathbf{q}-\mathbf{q}^{\prime}\right\rangle / d_{B}(\mathbf{p})
$$

Thus $d_{B}(\mathbf{p}) \mathcal{H}\left(B, B^{\prime}\right) \geq\left\langle\mathbf{q}-\mathbf{p}, \mathbf{q}-\mathbf{q}^{\prime}\right\rangle$, and this inequality also holds without the assumption $\mathbf{q} \neq \mathbf{p}$. Summing this identity with the similar one obtained exchanging the roles of $\left(B, B^{\prime}\right)$ and $\left(\mathbf{q}, \mathbf{q}^{\prime}\right)$ we obtain $\left(d_{B}(\mathbf{p})+d_{B^{\prime}}(\mathbf{p})\right) \mathcal{H}\left(B, B^{\prime}\right) \geq\left\langle(\mathbf{q}-\mathbf{p})+\left(\mathbf{p}-\mathbf{q}^{\prime}\right), \mathbf{q}-\mathbf{q}^{\prime}\right\rangle=\left\|\mathbf{q}-\mathbf{q}^{\prime}\right\|^{2}$ which is the announced result.

Here and below, slightly abusively, we regard normal vectors to $\partial \Omega$ as co-vectors.
Proposition A.2. Assume that the domain boundary $\partial \Omega$ admits outward normals $\mathbf{n}(\mathbf{q}), \mathbf{q} \in \partial \Omega$, which have $1 / R_{\Omega}$-Lipschitz regularity w.r.t. $\mathbf{q}$. Then $P_{\partial \Omega}(\mathbf{p})$ is uniquely defined for all $\mathbf{p} \in \mathbb{E}$ such that $d_{\partial \Omega}(\mathbf{p})<R_{\Omega}$. Furthermore $d_{\partial \Omega}(\mathbf{p}+\dot{\mathbf{e}})$, for $\mathbf{p}, \dot{\mathbf{e}} \in \mathbb{E}$, is either zero or obeys

$$
d_{\partial \Omega}(\mathbf{p}+\dot{\mathbf{e}}) \leq d_{\partial \Omega}(\mathbf{p})+\langle\mathbf{n}(\mathbf{p}), \dot{\mathbf{e}}\rangle+\frac{\|\dot{\mathbf{e}}\|^{2}}{2 R_{\Omega}}, \quad \text { where } \mathbf{n}(\mathbf{p}):=\mathbf{n}\left(\mathrm{P}_{B}(\mathbf{p})\right)
$$

Proof. The Lipschitz assumption on the normals implies, for any $\mathbf{q} \in \partial \Omega$, the inclusions $B(\mathbf{q}-$ $\left.R_{\Omega} \mathbf{n}(\mathbf{q}), R_{\Omega}\right) \subseteq \bar{\Omega}$ and $B\left(\mathbf{q}+R_{\Omega} \mathbf{n}(\mathbf{q}), R_{\Omega}\right) \subseteq \mathbb{E} \backslash \Omega$. Fix $\mathbf{p} \in \mathbb{E}$, and let $\mathbf{q} \in \partial \Omega$ be an arbitrary closest point to $\mathbf{p}$. The first inclusion implies the announced uniqueness when $d_{\partial \Omega}(\mathbf{p})<R_{\Omega}$, and the second inclusion, together with Lemma 2.9 applied to $\left\|\left(\mathbf{q}+R_{\Omega} \mathbf{n}(\mathbf{q})-\mathbf{p}\right)-\dot{\mathbf{e}}\right\|$, implies the distance estimate.

## A. 1 Estimating $u_{\varepsilon}-u$

In this subsection we bound, in the uniform norm, the difference between the value function $u$ of the pre-Riemannian problem, and the one $u_{\varepsilon}$ associated to the Riemannian approximation (19),
for any $0<\varepsilon \leq 1$. Assumption 1.7, on global controllability and short time local controllability at the boundary, is central in the proof. Related arguments can be found in 3], but the proof is provided for completeness and because in the process we establish estimates used in \$A.2. We use the notations and assumptions of Theorem 1.9. Let us introduce the control sets of the Riemannian relaxation

$$
\mathcal{B}_{\varepsilon}(\mathbf{p}):=\left\{\dot{\mathbf{p}} \in \mathbb{E} ;\|\dot{\mathbf{p}}\|_{\mathcal{M}_{\varepsilon}(\mathbf{p})} \leq 1\right\}
$$

for each $0<\varepsilon \leq 1, \mathbf{p} \in \bar{\Omega}$, with the convention that $\mathcal{B}_{0}$ denotes the pre-Riemannian control sets of Definition 1.6. Note the inclusion $\mathcal{B}_{\varepsilon}(\mathbf{p}) \subseteq \mathcal{B}_{\varepsilon^{\prime}}(\mathbf{p})$ for any $0 \leq \varepsilon \leq \varepsilon^{\prime} \leq 1$, which implies the pointwise inequalities $T_{\mathcal{B}_{\varepsilon}}(\mathbf{p}, \mathbf{q}) \geq T_{\mathcal{B}_{\varepsilon^{\prime}}}(\mathbf{p}, \mathbf{q})$ for the control times, and $u_{\varepsilon}(\mathbf{p}) \geq u_{\varepsilon^{\prime}}(\mathbf{p})$ for the exit times, for any $\mathbf{p}, \mathbf{q} \in \bar{\Omega}$.

Our first lemma establishes the Lipschitz regularity of the control sets $\mathcal{B}_{\varepsilon}(\mathbf{p})$ and of the tensors $\mathcal{D}_{\varepsilon}(\mathbf{p})$ with respect to the position $\mathbf{p} \in \bar{\Omega}$ and the relaxation parameter $\varepsilon \in[0,1]$. The control sets regularity allows to apply Gronwall's Lemma in the proof of Proposition A.4, whereas the tensors regularity is used in $\S$ A. 2 for establishing Assumption (i) of the doubling of variables argument Theorem 2.3. Denote by $A_{\varepsilon}(\mathbf{p})$ the matrix of columns $\dot{\omega}_{1}, \cdots, \dot{\omega}_{n}, \varepsilon \dot{\omega}_{1}^{*}, \cdots, \varepsilon \dot{\omega}_{n^{*}}^{*}$, for all $0 \leq \varepsilon \leq 1, \mathbf{p} \in \bar{\Omega}$. Since the these vector fields are Lipschitz and bounded, the matrix field $A_{\varepsilon}$ obeys for some constant $C_{\text {Lip }}^{\mathcal{D}}$ the following Lipschitz regularity property: for all $\mathbf{p}, \mathbf{q} \in \bar{\Omega}$ and all $\varepsilon, \varepsilon^{\prime} \in[0,1]$

$$
\begin{equation*}
\max \left\{\left\|A_{\varepsilon}(\mathbf{p})-A_{\varepsilon^{\prime}}(\mathbf{q})\right\|,\left\|A_{\varepsilon}^{\mathrm{T}}(\mathbf{p})-A_{\varepsilon^{\prime}}^{\mathrm{T}}(\mathbf{q})\right\|\right\} \leq C_{\mathrm{Lip}}^{\mathrm{D}}\left(\|\mathbf{p}-\mathbf{q}\|+\left|\varepsilon-\varepsilon^{\prime}\right|\right) . \tag{43}
\end{equation*}
$$

Recall that the operator norm of an $m \times n$ matrix $A$ is defined by $\|A\|:=\sup \{\|A \dot{\mathbf{r}}\| ; \dot{\mathbf{r}} \in$ $\left.\mathbb{R}^{n},\|\dot{\mathbf{r}}\| \leq 1\right\}$. One easily checks that $\|A\|=\left\|A^{\mathrm{T}}\right\|$ for any matrix $A$, hence the l.h.s. of (43) could be slightly simplified, but we prefer to emphasize the fact that both the regularity of the matrix field $A_{\varepsilon}$ and of its transpose are used.

Lemma A.3. One has the Lipschitz regularity: for all $\mathbf{p}, \mathbf{q} \in \bar{\Omega}, \varepsilon, \varepsilon^{\prime} \in[0,1]$, and all $\hat{\mathbf{r}} \in \mathbb{E}^{*}$

$$
\begin{align*}
\mathcal{H}\left(\mathcal{B}_{\varepsilon}(\mathbf{p}), \mathcal{B}_{\varepsilon^{\prime}}(\mathbf{q})\right) & \leq C_{\operatorname{Lip}}^{\mathcal{D}}\left(\|\mathbf{q}-\mathbf{p}\|+\left|\varepsilon-\varepsilon^{\prime}\right|\right)  \tag{44}\\
\left|\|\hat{\mathbf{r}}\|_{\mathcal{D}_{\varepsilon}(\mathbf{p})}-\|\hat{\mathbf{r}}\|_{\mathcal{D}^{\prime}}(\mathbf{q})\right| & \leq C_{\operatorname{Lip}}^{\mathcal{D}}\left(\|\mathbf{q}-\mathbf{p}\|+\left|\varepsilon-\varepsilon^{\prime}\right|\right)\|\hat{\mathbf{r}}\| . \tag{45}
\end{align*}
$$

Proof. For any maps $\varphi, \psi$ from an arbitrary space $X$ to $\mathbb{E}$, one has $\mathcal{H}(\varphi(X), \psi(X)) \leq \sup _{x \in X} \| \varphi(x)-$ $\psi(x) \|$. Observing that $\mathcal{B}_{\varepsilon}(\mathbf{p})=\left(A_{\varepsilon}(\mathbf{p})\right)(B)$, where $B$ is the unit ball of $\mathbb{R}^{n+n^{*}}$, we obtain

$$
\mathcal{H}\left(\mathcal{B}_{\varepsilon}(\mathbf{p}), \mathcal{B}_{\varepsilon^{\prime}}(\mathbf{q})\right) \leq \sup _{x \in B}\left\|A_{\varepsilon}(\mathbf{p}) x-A_{\varepsilon^{\prime}}(\mathbf{q}) x\right\|=\left\|A_{\varepsilon}(\mathbf{p})-A_{\varepsilon^{\prime}}(\mathbf{q})\right\| \leq C_{\text {Lip }}^{\mathcal{D}}\left(\|\mathbf{q}-\mathbf{p}\|+\left|\varepsilon-\varepsilon^{\prime}\right|\right)
$$

which is (44). Observing that $\mathcal{D}_{\varepsilon}(\mathbf{p})=A_{\varepsilon}(\mathbf{p}) A_{\varepsilon}(\mathbf{p})^{\mathrm{T}}$, see Definition 1.6, we obtain 45), since

$$
\left|\|\hat{\mathbf{r}}\|_{\mathcal{D}_{\varepsilon}(\mathbf{p})}-\|\hat{\mathbf{r}}\|_{\mathcal{D}_{\varepsilon^{\prime}}(\mathbf{q})}\right|=\left|\left\|A_{\varepsilon}^{\mathrm{T}}(\mathbf{p}) \hat{\mathbf{r}}\right\|-\left\|A_{\varepsilon^{\prime}}^{\mathrm{T}}(\mathbf{q}) \hat{\mathbf{r}}\right\|\right| \leq\left\|A_{\varepsilon}^{\mathrm{T}}(\mathbf{p})-A_{\varepsilon^{\prime}}^{\mathrm{T}}(\mathbf{q})\right\|\|\hat{\mathbf{r}}\| .
$$

Proposition A.4. Let $0 \leq \varepsilon_{0} \leq 1$, let $\gamma_{0}:\left[0, T_{0}\right] \rightarrow \bar{\Omega}$ be a $\mathcal{B}_{\varepsilon_{0}}$-admissible path, see Definition 1.3. and let $\mathbf{p}_{0}:=\gamma_{0}(0)$. Let $0 \leq \varepsilon_{1} \leq 1$, and let $\gamma_{1}:\left[0, T_{1}\right] \rightarrow \bar{\Omega}$ be a solution to the ODE

$$
\dot{\gamma}_{1}(t):=\mathrm{P}_{\mathcal{B}_{\varepsilon_{1}}\left(\gamma_{1}(t)\right)}\left(\dot{\gamma}_{0}(t)\right),
$$

with initial condition $\gamma_{1}(0)=\mathbf{p}_{1} \in \Omega$, where the final time $T_{1}$ is either $T_{0}$ or the time where $\gamma_{1}$ reaches $\partial \Omega$. Then $\gamma_{1}$ is $\mathcal{B}_{1}$-admissible, and for any $0 \leq t \leq T_{1}$ one has

$$
\left|\gamma_{0}(t)-\gamma_{1}(t)\right| \leq\left(\left|\mathbf{p}_{0}-\mathbf{p}_{1}\right|+\left|\varepsilon_{0}-\varepsilon_{1}\right|\right) \exp \left(C_{\operatorname{Lip}}^{\mathcal{B}} t\right)-\left|\varepsilon_{0}-\varepsilon_{1}\right| .
$$

Proof. The orthgonal projection $\mathrm{P}_{B}(\mathbf{p})$, of a given $\mathbf{p} \in \mathbb{E}$, depends continuously (in fact with $(1 / 2)$-Holder regularity) on the closed and convex set $B$, see Proposition A.1. The right hand side of A.4 therefore depends continuously on $\gamma_{1}(t)$, hence this ODE admits solutions by Peano's existence theorem. Note that the Picard-Lindelof/Cauchy-Lipschitz uniqueness theorem does not apply since it requires Lipschitz regularity of the r.h.s., but the lack of uniqueness is fortunately not an issue in this proof.

The $\mathcal{B}_{\varepsilon_{1}}$-admissibility of $\gamma_{1}$ holds by construction, and since $\dot{\gamma}_{0}(t) \in \mathcal{B}_{\varepsilon_{0}}\left(\gamma_{0}(t)\right)$ one has

$$
\left|\dot{\gamma}_{0}(t)-\dot{\gamma}_{1}(t)\right| \leq \mathcal{H}\left(\mathcal{B}_{\varepsilon_{0}}\left(\gamma_{0}(t)\right), \mathcal{B}_{\varepsilon_{1}}\left(\gamma_{1}(t)\right)\right) \leq C_{\operatorname{Lip}}^{\mathcal{B}}\left(\left|\varepsilon_{0}-\varepsilon_{1}\right|+\left|\gamma_{0}(t)-\gamma_{1}(t)\right|\right)
$$

for any $0 \leq t \leq T_{1}$. The announced estimate then follows from Gronwall's lemma.
The following lemma makes use of the transversality property in Assumption 1.7 to upper bound the exit time function $u_{\varepsilon}$ close to the domain boundary $\partial \Omega$. This property is equivalent to: $\|\mathbf{n}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2}=\sum_{i=1}^{n}\left\langle\mathbf{n}(\mathbf{p}), \dot{\omega}_{i}(\mathbf{p})\right\rangle^{2} \neq 0$ for all $\mathbf{p} \in \partial \Omega$, where $\mathbf{n}(\mathbf{p})$ denotes the outward normal to $\partial \Omega$. Hence denoting $\mathbf{n}(\mathbf{p}):=\mathbf{n}\left(\mathrm{P}_{\partial \Omega}(\mathbf{p})\right)$ for $\mathbf{p}$ close enough to $\partial \Omega$, as in Proposition A.2, there exists by continuity positive constants $c_{\mathcal{D}}, c_{\Omega}$ such that

$$
\begin{equation*}
\|\mathbf{n}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})} \geq c_{\mathcal{D}}, \text { for all } \mathbf{p} \in \bar{\Omega} \text { such that } d_{\partial \Omega}(\mathbf{p}) \leq c_{\Omega} \tag{46}
\end{equation*}
$$

In the next lemma and proposition we construct paths from an arbitrary point $\mathbf{p} \in \bar{\Omega}$ to $\partial \Omega$, whereas the original problem $(9)$ is to find a path from $\partial \Omega$ to $\mathbf{p}$. This change of orientation is only used for notational simplicity, and is valid since the paths can be reverse parametrized, and since the control sets are symmetric: $\mathcal{B}_{\varepsilon}(\mathbf{p})=\left\{-\dot{\mathbf{p}} ; \dot{\mathbf{p}} \in \mathcal{B}_{\varepsilon}(\mathbf{p})\right\}$, for all $\mathbf{p} \in \bar{\Omega}, \varepsilon \in[0,1]$.

Lemma A.5. For all $\mathbf{p} \in \bar{\Omega}$ such that $d_{\partial \Omega}(\mathbf{p}) \leq c_{\Omega}$, one has $u(\mathbf{p}) \leq d_{\partial \Omega}(\mathbf{p}) / c_{\mathcal{D}}$.
Proof. Let $D \Omega:=\left\{\mathbf{q} \in \bar{\Omega} ; d_{\partial \Omega}(\mathbf{q}) \leq c_{\Omega}\right\}$. Define a vector field $\dot{\mathbf{v}}: D \Omega \rightarrow \mathbb{E}$ by

$$
\dot{\mathbf{v}}(\mathbf{q}):=\frac{\mathcal{D}(\mathbf{q}) \mathbf{n}(\mathbf{q})}{\|\mathbf{n}(\mathbf{q})\|_{\mathcal{D}(\mathbf{q})}}=A_{0}(\mathbf{q}) \frac{A_{0}^{\mathrm{T}}(\mathbf{q}) \mathbf{n}(\mathbf{q})}{\left\|A_{0}^{\mathrm{T}}(\mathbf{q}) \mathbf{n}(\mathbf{q})\right\|}
$$

for all $\mathbf{q} \in D \Omega$, and note that $\dot{\mathbf{v}}(\mathbf{q}) \in \mathcal{B}(\mathbf{q})$. Consider the solution to the $\operatorname{ODE} \dot{\gamma}(t):=\dot{\mathbf{v}}(\gamma(t))$, with initial condition $\gamma(0)=\mathbf{p} \in D \Omega$, stopping at the time $T \in[0, \infty]$ when $\gamma$ leaves $D \Omega$. By construction, $\gamma$ is a $\mathcal{B}_{\varepsilon}$-admissible path, and for all $0 \leq t \leq T$ one has by Proposition A. 2

$$
\frac{d}{d t} d_{\partial \Omega}(\gamma(t))=-\langle\mathbf{n}(\gamma(t)), \dot{\mathbf{v}}(\gamma(t))\rangle=-\|\mathbf{n}(\gamma(t))\|_{\mathcal{D}(\gamma(t))} \leq-c_{\mathcal{D}}
$$

Therefore $T \leq d_{\partial \Omega} / c_{\mathcal{D}}$ and $\gamma(T) \in \partial \Omega$, hence $u(\mathbf{p}) \leq T$ and the announced result follows.
The following proposition establishes the Lipschitz regularity of $u_{\varepsilon}(\mathbf{p})$ with respect to both $\varepsilon \in[0,1]$ and $\mathbf{p} \in \bar{\Omega}$. Regularity w.r.t. $\varepsilon$ proves 20 , left) in Theorem 1.9 , which was the aim of this section. Regularity w.r.t. $\mathbf{p}$ is used in the next subsection to apply the doubling of variables techniques.

Proposition A.6. One has the Lipschitz regularity property

$$
\begin{equation*}
\left|u_{\varepsilon_{0}}\left(\mathbf{p}_{0}\right)-u_{\varepsilon_{1}}\left(\mathbf{p}_{1}\right)\right| \leq C_{\operatorname{Lip}}^{u}\left(\left\|\mathbf{p}_{0}-\mathbf{p}_{1}\right\|+\left|\varepsilon_{0}-\varepsilon_{1}\right|\right) \tag{47}
\end{equation*}
$$

for all $\mathbf{p}_{0}, \mathbf{p}_{1} \in \bar{\Omega}$ and all $\varepsilon_{0}, \varepsilon_{1} \in[0,1]$, where $C_{\mathrm{Lip}}^{u}:=\exp \left(C_{\mathrm{Lip}}^{\mathcal{B}}\|u\|_{\infty}\right) / c_{\mathcal{D}}$.

Proof. Assume w.l.o.g. that $u_{\varepsilon_{0}}\left(\mathbf{p}_{0}\right) \leq u_{\varepsilon_{1}}\left(\mathbf{p}_{1}\right)$. Let $\gamma_{0}:\left[0, T_{0}\right] \rightarrow \bar{\Omega}$ be an optimal $\mathcal{B}_{\varepsilon_{0}}$-admissible path from $\mathbf{p}_{0}$ to $\partial \Omega$, where $T_{0}:=u_{\varepsilon_{0}}\left(\mathbf{p}_{0}\right)$. Let $\gamma_{1}:\left[0, T_{1}\right] \rightarrow \bar{\Omega}$ be as in Proposition A.4. If $\gamma_{1}\left(T_{1}\right) \in \partial \Omega$, then $u_{1}\left(\mathbf{p}_{1}\right) \leq T_{1} \leq T_{0}$ and the result follows.

Otherwise by Proposition A.4 we get $d_{\partial \Omega}\left(\gamma_{1}\left(T_{0}\right)\right) \leq\left\|\gamma_{1}\left(T_{0}\right)-\gamma_{0}\left(T_{0}\right)\right\| \leq C\left(\left|\varepsilon_{0}-\varepsilon_{1}\right|+\| \mathbf{p}_{0}-\right.$ $\left.\mathbf{p}_{1} \|\right)$ with $C:=\exp \left(C_{\text {Lip }} T_{0}\right)$. Therefore

$$
u_{\varepsilon_{1}}\left(\mathbf{p}_{1}\right) \leq T_{0}+d_{\partial \Omega}\left(\gamma_{1}\left(T_{0}\right)\right) / c_{\mathcal{D}} \leq u_{\varepsilon_{0}}\left(\mathbf{p}_{0}\right)+C_{\operatorname{Lip}}^{u}\left(\left\|\mathbf{p}_{0}-\mathbf{p}_{1}\right\|+\left|\varepsilon_{0}-\varepsilon_{1}\right|\right)
$$

as announced, using Lemma A. 5 and recalling that $u_{\varepsilon_{1}} \leq u_{0}=u$ on $\bar{\Omega}$, and assuming that $d_{\partial \Omega}\left(\gamma_{1}\left(T_{0}\right)\right) \leq c_{\Omega}$. Thus 47) holds for all $\left(\mathbf{p}_{0}, \varepsilon_{0}\right),\left(\mathbf{p}_{1}, \varepsilon_{1}\right) \in \bar{\Omega} \times[0,1]$ obeying $\left\|\mathbf{p}_{0}-\mathbf{p}_{1}\right\|+$ $\left|\varepsilon_{0}-\varepsilon_{1}\right| \leq c_{\Omega} / C$. By defining $u_{\varepsilon}(\mathbf{p})=0$ for all $\mathbf{p} \in \mathbb{E} \backslash \Omega, \varepsilon \in[0,1]$, the result extends to all $\left(\mathbf{p}_{0}, \varepsilon_{0}\right),\left(\mathbf{p}_{1}, \varepsilon_{1}\right) \in \mathbb{E} \times[0,1]$ subject to the same closeness constraint, which finally can be removed since Lipschitz regularity on a convex set is a local property.

## A. 2 Estimating $U_{h, \varepsilon}-u_{\varepsilon}$

In this subsection we complete the proof of Theorem 1.9 by estimating the difference $\left|U_{h, \varepsilon}-u_{\varepsilon}\right|$ on $\Omega_{h}$, for any $\left.\left.h>0, \varepsilon \in\right] 0,1\right]$, using the doubling of variables technique Theorem 2.3 . We use the notations and assumptions of Theorems 1.9 and 2.3 , except of course assumptions (i) and (ii) of the latter which we intend to prove.

The first assumption of Theorem 2.3 is the Lipschitz regularity of the value function $u_{\varepsilon}$, which is established in Proposition A. 6 above, with a constant $C_{\text {Lip }}^{u}$ independent of $\varepsilon$. Note that, in contrast, naively adapting the Riemannian argument of Proposition 2.5 yields the Lipschitz constant $\lambda\left(\mathcal{M}_{\varepsilon}\right) \approx \varepsilon^{-1}$ exploding as $\varepsilon \rightarrow 0$, and thus unsuitable for proving Theorem 1.9.

The next step is to establish assumption (i) of Theorem [2.3. Lemma 2.6 from the Riemannian case applies without modification to $u_{\varepsilon}$ and $U_{h, \varepsilon}$ since it does not involve quantitative properties of the Riemannian metric field $\mathcal{M}_{\varepsilon}$. Lemma 2.7 from the Riemannian case also applies, with constants independent of $\varepsilon \in] 0,1]$. Indeed the dual tensors have the expression $\mathcal{D}_{\varepsilon}=\mathcal{D}+\varepsilon^{2} \mathcal{D}^{*}$, see (19), and therefore their max norm (28) is bounded $\lambda^{*}\left(\mathcal{D}_{\varepsilon}\right) \leq \lambda^{*}\left(\mathcal{D}_{1}\right)$ independently of $\varepsilon \in[0,1]$. The last ingredient used to prove assumption (i) in the Riemannian case is the Lipschitz regularity of the dual norms (31), which is established in (45) for the pre-Riemannian model with a constant $C_{\text {Lip }}^{\mathcal{D}}$ independent of $\varepsilon \in[0,1]$.

The following proposition establishes assumption (ii) of Theorem 2.3, a discrete counterpart of short time local controllability at the boundary, by adapting the arguments developed in the Riemannian case, see Proposition 2.10. Note the use Assumption 1.7, which is required due to the lack of uniform definiteness of the the tensors $\mathcal{D}_{\varepsilon}, 0<\varepsilon \leq 1$. The weights and offsets used in the decomposition (15) of $\mathcal{D}_{\varepsilon}(\mathbf{p})$ are denoted $\left(\rho_{i, \varepsilon}(\mathbf{p}), \dot{\mathbf{e}}_{i, \varepsilon}(\mathbf{p})\right)_{i=1}^{d^{\prime}}, \mathbf{p} \in \bar{\Omega}, 0<\varepsilon \leq 1$.
Proposition A.7. Let $\mathbf{p} \in \Omega_{h}$ be such that $d(\mathbf{p}, \partial \Omega) \leq c_{\Omega}$, and let $\left.\left.\varepsilon \in\right] 0,1\right]$. Then there exists $1 \leq i \leq d^{\prime}$ and a sign $s \in\{-1,1\}$ such that either $\mathbf{p}+h s \dot{\mathbf{e}}_{i, \varepsilon} \notin \Omega_{h}$ or

$$
U_{h, \varepsilon}(\mathbf{p}) \leq U_{h, \varepsilon}\left(\mathbf{p}+h s \dot{\mathbf{e}}_{i, \varepsilon}\right)+h C_{1}\left\|\dot{\mathbf{e}}_{i, \varepsilon}\right\|, \quad d\left(\mathbf{p}+h s \dot{\mathbf{e}}_{i, \varepsilon}, \partial \Omega\right) \leq d(\mathbf{p}, \partial \Omega)-h c_{2}\left\|\dot{\mathbf{e}}_{i, \varepsilon}\right\|,
$$

with $C_{1}:=\sqrt{d^{\prime}} / c_{\mathcal{D}}$ and $c_{2}:=c_{\mathcal{D}} /\left(2 \lambda^{*}\left(\mathcal{D}_{1}\right) \sqrt{d^{\prime}}\right)$. This implies assumption (ii) of Theorem 2.3. with the constants $C_{\mathrm{bd}}=C_{1} / c_{2}, C_{\mathrm{bd}}^{\prime}:=C_{\mathrm{bd}} r_{\varepsilon}$, and $c_{\mathrm{bd}}=c_{\mathcal{D}}$.
Proof. Using Assumption 1.7 and (46), one obtains a counterpart for (35)

$$
c_{\mathcal{D}}^{2} \leq\|\mathbf{n}(\mathbf{p})\|_{\mathcal{D}_{\varepsilon}(\mathbf{p})}^{2}=\sum_{1 \leq i \leq d^{\prime}} \rho_{i, \varepsilon}(\mathbf{p})\left\langle\mathbf{n}(\mathbf{p}), \dot{\mathbf{e}}_{i, \varepsilon}\right\rangle^{2},
$$

where $\mathbf{n}(\mathbf{p}):=\mathbf{n}\left(\mathrm{P}_{\partial \Omega}(\mathbf{p})\right)$. The proof is then similar to the one of Proposition 2.10, up to the replacement of Lemma 2.9 with Proposition A.2.

## B Convergence in the Rander case

This section is devoted to proof of Theorem 1.11, namely the numerical analysis of the Rander metric variant of our PDE discretization scheme, using its notations and assumptions. Consider the scheme $\mathfrak{F}_{h}$ on the discrete domain $\Omega_{h}$ defined for any $U: \Omega_{h} \rightarrow \mathbb{R}$ and $\mathbf{p} \in \Omega_{h}$ by $\left(\mathfrak{F}_{h} U(\mathbf{p})\right)^{2}:=$

$$
\begin{equation*}
h^{-2} \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p}) \max \left\{0, U(\mathbf{p})-U\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right)+h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle, U(\mathbf{p})-U\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)-h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle\right\}^{2} . \tag{48}
\end{equation*}
$$

By convention, $U$ is extended by 0 outside $\Omega_{h}$. Note that this scheme is non-causal as soon as some of the terms $h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle$ are non-zero, see Definition 2.1, in contrast with 27).

Proposition B.1. The scheme (48) is monotone. In addition:
(i) The null map $U=0$ satisfies $\mathfrak{F}_{h} U(\mathbf{p})=\|\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2}<1$ for all $\mathbf{p} \in \Omega_{h}$, hence it is a sub-solution for $\mathfrak{F}_{h}-1$.
(ii) Let $R>0$ be such that $\Omega$ is contained in the ball of radius $R-h r_{*}$, and let $U(\mathbf{p}):=R-\|\mathbf{p}\|$, for all $\mathbf{p} \in \Omega_{h}$. Then for all $\lambda \geq 0$, and all $\mathbf{p} \in \Omega_{h}$,

$$
\mathfrak{F}_{h}(\lambda U)(\mathbf{p}) \geq\|\lambda \mathbf{p} /\| \mathbf{p}\|-\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2},
$$

where $\mathbf{p} /\|\mathbf{p}\|$ can be replaced with an arbitrary unit vector in the case $\mathbf{p}=0$. Thus $\lambda U$ is a super-solution for all sufficiently large $\lambda$.
(iii) Let $U$ is a super-solution for $\mathfrak{F}_{h}-1$, and let $\mathbf{p}_{1}, \cdots, \mathbf{p}_{N}$ be the points of $\Omega_{h}$ ordered in such way that $U\left(\mathbf{p}_{1}\right) \leq \cdots \leq U\left(\mathbf{p}_{N}\right)$. For each $\varepsilon>0$ let $V_{\varepsilon}: \Omega_{h} \rightarrow \mathbb{R}$ be defined by $V_{\varepsilon}\left(\mathbf{p}_{i}\right)=U\left(\mathbf{p}_{i}\right)+\varepsilon-\varepsilon^{1+i}$. Then $V_{\varepsilon}$ is a strict super-solution to $\mathfrak{F}_{h}-1$ for all sufficiently small $\varepsilon$.

Proof. Point (i) follows from the identity $\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p})\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle^{2}=\|\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2}$, and the smallness assumption 21, right) on the co-vector field $\hat{\eta}$. Point (ii) is proved as in Proposition 2.4 Point (iii) is in contrast non-trivial. Let $U$ be a super-solution for $\mathfrak{F}_{h}-1$ and let $1 \leq k \leq N$. Denote by $m_{i}(\mathbf{p})$ the $i$-th maximum of three terms appearing in (48), so that $\mathfrak{F}_{h} U(\mathbf{p})=\sum_{i=1}^{d^{\prime}} \rho_{i}(\mathbf{p}) m_{i}(\mathbf{p})^{2}$ for each $\mathbf{p} \in \Omega_{h}$. Then one has the Taylor expansion

$$
\begin{equation*}
\mathfrak{F}_{h} V_{\varepsilon}\left(\mathbf{p}_{k}\right)-\mathfrak{F}_{h} U\left(\mathbf{p}_{k}\right)=2 \sum_{1 \leq i \leq d^{\prime}} \rho_{i}\left(\mathbf{p}_{k}\right) m_{i}\left(\mathbf{p}_{k}\right)\left(\varepsilon^{1+k_{i}}-\varepsilon^{1+k}\right)+\mathcal{O}\left(\varepsilon^{2+2 k_{*}}\right), \tag{49}
\end{equation*}
$$

where $k_{i}$ is an integer depending on $i \in\left\{1, \cdots, d^{\prime}\right\}$ and $k$, and chosen so that $\mathbf{p}_{k_{i}}=\mathbf{p}_{k}$, (resp. $\mathbf{p}_{k_{i}}=\mathbf{p}_{k}+h \dot{\mathbf{e}}_{i}$, resp. $\left.\mathbf{p}_{k_{i}}=\mathbf{p}_{k}-h \dot{\mathbf{e}}_{i}\right)$ if the maximum defining $m_{i}\left(\mathbf{p}_{k}\right)$ is achieved for the first (resp. second, resp. third) term. (If this point is outside $\Omega_{h}$, we let $k_{i}=0$. In case of a tie, the point with smallest index is chosen.) We also denoted $k_{*}:=\min \left\{k_{i} ; 1 \leq i \leq d^{\prime}\right\}$.

We prove below that $k_{i}<k$ for some $1 \leq i \leq d^{\prime}$, which by (49) implies that $\mathfrak{F}_{h} V_{\varepsilon}\left(\mathbf{p}_{k}\right)>$ $\mathfrak{F}_{h} U\left(\mathbf{p}_{k}\right) \geq 1$ for all sufficiently small $\varepsilon>0$ as announced. Assume for contradiction that $k_{i} \geq k$ for all $1 \leq i \leq d^{\prime}$, hence that $U\left(\mathbf{p}_{k_{i}}\right) \geq U\left(\mathbf{p}_{k}\right)$ and therefore that $m_{i}\left(\mathbf{p}_{k}\right) \leq\left|\left\langle\hat{\eta}\left(\mathbf{p}_{k}\right), \dot{\mathbf{e}}_{i}\right\rangle\right|$ for all $1 \leq i \leq d^{\prime}$. Then denoting $\mathbf{p}:=\mathbf{p}_{k}$ one obtains

$$
\mathfrak{F}_{h} U(\mathbf{p})=\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p}) m_{i}(\mathbf{p})^{2} \leq \sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p})\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle^{2}=\|\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2}<1,
$$

in contradiction with our assumption that $U$ is a super-solution for $\mathfrak{F}_{h}-1$. The result follows.

In the rest of this section, we establish the properties required to apply the doubling of variables argument Theorem 2.3, using its notations. The first ingredient is the Lipschitz regularity of the exit time value function $u$. By construction the Rander metric (21) satisfies for all $\mathbf{p} \in \bar{\Omega}$ and all $\dot{\mathbf{p}} \in \mathbb{E}$.

$$
\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}})=\|\dot{\mathbf{p}}\|_{\mathcal{M}(\mathbf{p})}+\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{p}}\rangle \leq\left(\lambda^{*}(\mathcal{M})+\|\hat{\eta}\|_{\infty}\right)\|\dot{\mathbf{p}}\| .
$$

Hence $u$ is $C_{\text {Lip }}^{u}$-Lipschitz by Proposition 2.5, as desired, with constant $C_{\text {Lip }}^{u}:=\lambda^{*}(\mathcal{M})+\|\hat{\eta}\|_{\infty}$ where the max-norm $\lambda^{*}(\mathcal{M})$ of a tensor field is defined in (28).

Establishing assumption (i) of Theorem 2.3. We proceed similarly to the Riemannian case 2.2 , starting with an extension of Lemma 2.6 to Rander metrics. The proof, left to the reader, is similar up to fact that the scheme $\mathfrak{F}_{h} U$ and the PDE operator $\|\mathrm{d} u-\hat{\eta}\|_{\mathcal{D}}$ are not homogeneous w.r.t. their respective argument $U$ and $u$. The variables $\lambda \in[1 / 2,1[, \delta>0$, $(\overline{\mathbf{p}}, \overline{\mathbf{q}}),(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \in \Omega_{h} \times \Omega$ are from Theorem 2.3 .
Lemma B.2. Let $\overline{\mathbf{w}}:=(\overline{\mathbf{p}}-\overline{\mathbf{q}}) / \delta$ and let $\bar{U}(\mathbf{p}):=\langle\overline{\mathbf{w}}, \mathbf{p}\rangle+\frac{1}{2 \delta}\|\mathbf{p}-\overline{\mathbf{p}}\|^{2}$, for all $\mathbf{p} \in \mathbb{L}_{h}$. Then

$$
\mathfrak{F}_{h}(\bar{U} / \lambda)(\overline{\mathbf{p}}) \leq 1 \quad\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{q}})\|_{\mathcal{D}(\overline{\mathbf{q}})} \geq 1
$$

Let $\tilde{\mathbf{w}}:=(\tilde{\mathbf{p}}-\tilde{\mathbf{q}}) / \delta$ and let $\tilde{U}(\mathbf{p}):=\langle\tilde{\mathbf{w}}, \mathbf{p}\rangle-\frac{1}{2 \delta}\|\mathbf{p}-\tilde{\mathbf{p}}\|^{2}$, for all $\mathbf{p} \in \mathbb{L}_{h}$. Then

$$
\mathfrak{F}_{h} \tilde{U}(\tilde{\mathbf{p}}) \geq 1 \quad\|\tilde{\mathbf{w}} / \lambda-\hat{\eta}(\tilde{\mathbf{q}})\|_{\mathcal{D}(\tilde{\mathbf{q}})} \leq 1
$$

Let $C_{\text {Lip }}^{\mathcal{D}}$ be a Lipschitz regularity constants for the tensors $\mathcal{D}$, in the sense of (31). Let also $C_{\text {Lip }}^{\hat{\eta}}$ and $c_{\mathcal{D}, \hat{\eta}}>0$ be such that for all $\mathbf{p}, \mathbf{q} \in \bar{\Omega}$

$$
\|\hat{\eta}(\mathbf{p})-\hat{\eta}(\mathbf{q})\| \leq C_{\operatorname{Lip}}^{\hat{\eta}}\|\mathbf{p}-\mathbf{q}\|, \quad \quad\|\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})} \leq 1-c_{\mathcal{D}, \hat{\eta}} .
$$

Similarly to the Riemannian case, we argue by contradiction to establish assumption (i). Indeed, a contraposition of the following lemma shows that, if parameter $\lambda$ is below a given bound, then $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \notin \Omega_{h} \times \Omega$ and $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \notin \Omega_{h} \times \Omega$, which is assumption (i) of Theorem 2.3.
Lemma B.3. Assume that $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \in \Omega_{h} \times \Omega$ and define $\overline{\mathbf{w}}$ and $\bar{U}$ as in Lemma 2.6. Then

$$
\begin{equation*}
\left.\mid \mathfrak{F}_{h} \overline{( } U / \lambda\right) \left.(\overline{\mathbf{p}})-\|\overline{\mathbf{w}} / \lambda-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})}\left|\leq C_{0} r_{*} \frac{h}{\delta}, \quad\right|\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{q}})\|_{\mathcal{D}(\overline{\mathbf{q}})}-\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})} \right\rvert\, \leq C_{1} \delta \tag{50}
\end{equation*}
$$

with $C_{0}:=\lambda^{*}(\mathcal{D}) \sqrt{d}$ and $C_{1}:=C_{\text {Lip }}^{\mathcal{D}}\left(4 C_{\text {Lip }}^{u}\right)\left(4 C_{\text {Lip }}^{u}+C_{\text {Lip }}^{\hat{\eta}}\right)+\lambda^{*}(\mathcal{D}) C_{\text {Lip }}^{\hat{\eta}}$. Assuming $\delta \leq c_{\mathcal{D}, \hat{\eta}} /\left(2 C_{1}\right)$ this implies $\lambda \geq 1-\frac{2}{c_{\mathcal{D}, \hat{\eta}}}\left(C_{0} r_{*} \frac{h}{\delta}+C_{1} \delta\right)$. The same estimates and conclusion hold for $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}})$.
Proof. We focus on the case of $(\overline{\mathbf{p}}, \overline{\mathbf{q}})$, the second case of $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}})$ being similar, and we begin with the proof of 50 , left). By definition of the quadratic function $\bar{U}$, one has for any $1 \leq i \leq d^{\prime}$

$$
\begin{align*}
& \max \left\{0, \bar{U}(\overline{\mathbf{p}}) / \lambda-\bar{U}\left(\overline{\mathbf{p}}+\dot{\mathbf{e}}_{i}\right) / \lambda+h\left\langle\hat{\eta}(\overline{\mathbf{p}}), \dot{\mathbf{e}}_{i}\right\rangle, \bar{U}(\overline{\mathbf{p}}) / \lambda-\bar{U}\left(\overline{\mathbf{p}}-\dot{\mathbf{e}}_{i}\right) / \lambda-h\left\langle\hat{\eta}(\overline{\mathbf{p}}), \dot{\mathbf{e}}_{i}\right\rangle\right\}  \tag{51}\\
& =h\left|\left\langle\overline{\mathbf{w}} / \lambda-\hat{\eta}(\overline{\mathbf{p}}), \dot{\mathbf{e}}_{i}\right\rangle\right|+\frac{h^{2}}{\delta}\left\|\dot{\mathbf{e}}_{i}\right\|^{2} .
\end{align*}
$$

From this point, the arguments developed in the Riemannian case apply without modification. The second estimate (50, right) follows from

$$
\left|\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{q}})\|_{\mathcal{D}(\overline{\mathbf{q}})}-\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})}\right| \leq C_{\operatorname{Lip}}^{\mathcal{D}}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|(\|\overline{\mathbf{w}}\|+\|\hat{\eta}(\overline{\mathbf{q}})\|)+\lambda^{*}(\mathcal{D}) C_{\operatorname{Lip}}^{\hat{\eta}}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|,
$$

combined with $\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\| \leq 4 C_{\mathrm{Lip}}^{u} \delta$, see Theorem 2.3 , hence $\|\overline{\mathbf{w}}\| \leq 4 C_{\mathrm{Lip}}^{u}$. Combining the two estimates $(32$ with Lemma B.2 and using the convexity of the norm, we obtain

$$
1+C_{0} r_{*} \frac{h}{\delta} \geq\|\overline{\mathbf{w}} / \lambda-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})} \geq\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})}+\frac{1 / \lambda-1}{\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\overline{\mathbf{p}})}}\langle\overline{\mathbf{w}}, \mathcal{D}(\mathbf{p})(\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}}))\rangle
$$

The scalar product in the above r.h.s. is bounded below as follows

$$
\begin{aligned}
2\langle\overline{\mathbf{w}}, \mathcal{D}(\mathbf{p})(\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}}))\rangle & =\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\mathbf{p})}^{2}+\|\overline{\mathbf{w}}\|_{\mathcal{D}(\overline{\mathbf{p}})}^{2}-\|\hat{\eta}(\mathbf{p})\|_{\mathcal{D}(\mathbf{p})}^{2} \\
& \geq\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\mathbf{p})}^{2}-\left(1-c_{\mathcal{D}, \eta}\right)^{2}
\end{aligned}
$$

Using $\|\overline{\mathbf{w}}-\hat{\eta}(\overline{\mathbf{p}})\|_{\mathcal{D}(\mathbf{p})} \geq 1-C_{1} \delta$, and assuming $2 C_{1} \delta \leq c_{\mathcal{D}, \eta}$ for the second inequality, we obtain

$$
1+C_{0} r_{*} \frac{h}{\delta} \geq\left(1-C_{1} \delta\right)+(1 / \lambda-1)\left(\left(1-C_{1} \delta\right)-\frac{\left(1-c_{\mathcal{D}, \eta}\right)^{2}}{1-C_{1} \delta}\right) \geq\left(1-C_{1} \delta\right)+(1-\lambda) \frac{c_{\mathcal{D}, \eta}}{2}
$$

We used the elementary inequalities $1 / \lambda-1 \geq 1-\lambda$ and $(1-c / 2)-(1-c)^{2} /(1-c / 2) \geq c / 2$. This implies the announced lower bound for $\lambda$.

Establishing assumption (ii) of Theorem 2.3. The case of Rander metrics only requires a minor adaptation of the Riemannian argument, presented Proposition 2.10 .

Proposition B.4. Let $\mathbf{p} \in \Omega_{h}$, and let $\mathbf{q} \in \mathbb{E}$ be such that $\|\mathbf{p}-\mathbf{q}\| \geq C_{0} r_{*} h$, where $C_{0}:=$ $\mu(\mathcal{D}) \sqrt{d^{\prime}}$. Then there exists $1 \leq i \leq d^{\prime}$ and a sign $s \in\{-1,1\}$ such that

$$
U_{h}(\mathbf{p}) \leq U_{h}\left(\mathbf{p}+h s \dot{\mathbf{e}}_{i}\right)+h C_{1}\left\|\dot{\mathbf{e}}_{i}\right\|, \quad \quad\left\|\mathbf{p}+h s \dot{\mathbf{e}}_{i}-\mathbf{q}\right\| \leq\|\mathbf{p}-\mathbf{q}\|-h c_{2}\left\|\dot{\mathbf{e}}_{i}\right\|
$$

with $C_{1}:=\lambda^{*}(\mathcal{M}) \sqrt{d^{\prime}}+\|\hat{\eta}\|_{\infty}, c_{2}:=1 /\left(2 \mu(\mathcal{D}) \sqrt{d^{\prime}}\right)$. This implies assumption (ii) of Theorem 2.3. with the constants $C_{\mathrm{bd}}^{U}=C_{1} / c_{2}, C_{\mathrm{bd}}^{\prime U}:=C_{\mathrm{bd}} C_{\Omega} C_{0} r_{*}$, and $c_{\mathrm{bd}}^{U}=+\infty$.

Proof. The arguments developed in the Riemannian case apply with the following adaptation of (37): by definition of the discretization scheme (48) one has $h^{-2} \rho^{2} \max \left\{0, U_{h}(\mathbf{p})-U_{h}(\mathbf{p}+\dot{\mathbf{e}})-\right.$ $h\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}\rangle\}^{2} \leq\left(\mathfrak{F}_{h} U_{h}(\mathbf{p})\right)^{2}=1$, hence

$$
U(\mathbf{p})-U(\mathbf{p}+\dot{\mathbf{e}}) \leq \frac{h}{\rho}+h\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}\rangle \leq h\|\dot{\mathbf{e}}\|\left(\frac{\sqrt{d^{\prime}}}{\lambda_{*}}+\|\hat{\eta}\|_{\infty}\right)
$$

## C Doubling of variables

We establish in this section the doubling of variables argument, presented in Theorem 2.3 and adapted from [24]. Since the domain $\Omega$ is by assumption bounded, its closure $\bar{\Omega}$ is compact, and its sampling $\Omega_{h}:=\Omega \cap h \mathbb{L}$ is finite. Since the functions $u$ and $U_{h}$ are supported on these sets, and since $u$ is $C_{\text {Lip }}$-Lipschitz hence continuous, the maxima $\bar{M}_{\lambda, \delta}$ and $\widetilde{M}_{\lambda, \delta}$ are well defined and attained, at some point pairs $(\overline{\mathbf{p}}, \overline{\mathbf{q}}),(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \in(h \mathbb{L}) \times \mathbb{E}$. Our first step is to establish the closeness of $\overline{\mathbf{p}}$ with $\overline{\mathbf{q}}$, and of $\tilde{\mathbf{p}}$ with $\tilde{\mathbf{q}}$, as announced in Theorem 2.3 .

Lemma C.1. Under the assumptions of Theorem 2.3, one has $\max \{\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|,\|\tilde{\mathbf{p}}-\tilde{\mathbf{q}}\|\} \leq 4 C_{\mathrm{Lip}} \delta$.

Proof. Using the optimality properties defining $(\overline{\mathbf{p}}, \overline{\mathbf{q}})$ and $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}})$, see Theorem 2.3, and comparing with the alternative point pairs $(\overline{\mathbf{p}}, \overline{\mathbf{p}})$ and ( $\tilde{\mathbf{p}}, \tilde{\mathbf{p}})$ respectively, one obtains

$$
\begin{aligned}
& \lambda U_{h}(\overline{\mathbf{p}})-u(\overline{\mathbf{q}})-\frac{1}{2 \delta}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|^{2}=\bar{M}_{\lambda, \delta} \geq \lambda U_{h}(\overline{\mathbf{p}})-u(\overline{\mathbf{p}}), \\
& \lambda u(\tilde{\mathbf{q}})-U_{h}(\tilde{\mathbf{p}})-\frac{1}{2 \delta}\|\tilde{\mathbf{p}}-\tilde{\mathbf{q}}\|^{2}=\widetilde{M}_{\lambda, \delta} \geq \lambda u(\tilde{\mathbf{p}})-U_{h}(\tilde{\mathbf{p}}) .
\end{aligned}
$$

The next line is obtained by first rearranging the terms of these inequalities, eliminating in particular the instances of $U_{h}$, and then using the Lipschitz regularity of $u$

$$
\frac{1}{2 \delta}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|^{2} \leq u(\overline{\mathbf{p}})-u(\overline{\mathbf{q}}) \leq C_{\operatorname{Lip}}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|, \quad \frac{1}{2 \delta}\|\tilde{\mathbf{p}}-\tilde{\mathbf{q}}\|^{2} \leq \lambda(u(\tilde{\mathbf{q}})-u(\tilde{\mathbf{p}})) \leq \lambda C_{\operatorname{Lip}}\|\tilde{\mathbf{q}}-\tilde{\mathbf{p}}\| .
$$

The announced result follows from these estimates and the upper bound $1 / \lambda \leq 2$.
The rest of this section is devoted to the proof of (25), and begins with an estimate of $u-U_{h}$ in terms of the suprema $\bar{M}_{\lambda, \delta}, \widetilde{M}_{\lambda, \delta}$ and of the max norm $\|u\|_{\infty}:=\sup _{\mathbb{E}}|u|$, which is well defined by continuity of $u$ and compactness of its support.

Lemma C.2. Under the assumptions of Theorem 2.3, one has

$$
\sup _{\mathbf{p} \in h \mathbb{L}}\left|u(\mathbf{p})-U_{h}(\mathbf{p})\right| \leq 2\left(\max \left\{\bar{M}_{\lambda, \delta}, \widetilde{M}_{\lambda, \delta}\right\}+(1-\lambda)\|u\|_{\infty}\right) .
$$

Proof. By the optimality properties of $(\overline{\mathbf{p}}, \overline{\mathbf{q}})$ and ( $\tilde{\mathbf{p}}, \tilde{\mathbf{q}})$ one obtains for any $\mathbf{p} \in \Omega_{h}$, respectively,

$$
\begin{aligned}
\bar{M}_{\lambda, \delta} & \geq \lambda U_{h}(\mathbf{p})-u(\mathbf{p}) \geq \lambda\left(U_{h}(\mathbf{p})-u(\mathbf{p})\right)-(1-\lambda)\|u\|_{\infty} . \\
\widetilde{M}_{\lambda, \delta} & \geq \lambda u(\mathbf{p})-U_{h}(\mathbf{p}) \geq\left(u(\mathbf{p})-U_{h}(\mathbf{p})\right)-(1-\lambda)\|u\|_{\infty} .
\end{aligned}
$$

The announced result follows from these one-sided estimates on $u-U_{h}$, and from $1 / \lambda \leq 2$.
The next paragraph establishes some conditional estimates on $\bar{M}_{\lambda, \delta}$ and $\widetilde{M}_{\lambda, \delta}$, depending on the location of the points $\overline{\mathbf{p}}, \overline{\mathbf{q}}, \tilde{\mathbf{p}}, \tilde{\mathbf{q}}$. If $\tilde{\mathbf{p}} \in h \mathbb{L} \backslash \Omega_{h}$, then $U_{h}(\tilde{\mathbf{p}})=0$ and we obtain

$$
\begin{equation*}
\widetilde{M}_{\lambda, \delta}=\lambda u(\tilde{\mathbf{q}})-\frac{1}{2 \delta}\|\tilde{\mathbf{p}}-\tilde{\mathbf{q}}\|^{2} \leq \lambda u(\tilde{\mathbf{q}}) \leq \lambda C_{\mathrm{Lip}} d_{\partial \Omega}(\tilde{\mathbf{q}}) \leq C_{\mathrm{Lip}}\|\tilde{\mathbf{p}}-\tilde{\mathbf{q}}\| \leq 4 C_{\mathrm{Lip}}^{2} \delta . \tag{52}
\end{equation*}
$$

We used successively the negativity of the quadratic term, the Lipschitz regularity of $u$ and the fact that it vanishes outside $\Omega$, the fact that $\tilde{\mathbf{p}} \in h \mathbb{L} \backslash \Omega_{h} \subseteq \mathbb{E} \backslash \Omega$, and the previously obtained estimate on $\tilde{\mathbf{p}}-\tilde{\mathbf{q}}$. Likewise, if $\overline{\mathbf{p}} \in h \mathbb{L} \backslash \Omega_{h}$ then $\bar{M}_{\lambda, \delta} \leq-u(\overline{\mathbf{p}}) \leq 4 C_{\text {Lip }}^{2} \delta$.

Next if $\overline{\mathbf{q}} \in \mathbb{E} \backslash \Omega$, then $u(\overline{\mathbf{q}})=0$ and we obtain similarly to (52)

$$
\bar{M}_{\lambda, \delta}=\lambda U_{h}(\overline{\mathbf{p}})-\frac{1}{2 \delta}\|\overline{\mathbf{p}}-\overline{\mathbf{q}}\|^{2} \leq U_{h}(\overline{\mathbf{p}}) \leq C_{\mathrm{bd}} d_{\partial \Omega}(\overline{\mathbf{p}})+C_{\mathrm{bd}}^{\prime} h \leq 4 C_{\mathrm{Lip}} C_{\mathrm{bd}} \delta+C_{\mathrm{bd}}^{\prime} h .
$$

We used the same arguments as in (52), except for the second inequality which is based on assumption (ii) of Theorem 2.3. Likewise, if $\tilde{\mathbf{q}} \in \mathbb{E} \backslash \Omega$ then $\widetilde{M}_{\lambda, \delta} \leq-U_{h}(\tilde{\mathbf{p}}) \leq C_{\mathrm{bd}}\left(h+4 C_{\mathrm{Lip}} \delta\right)$.

Following assumption (i) of Theorem 2.3, we assume that $(\overline{\mathbf{p}}, \overline{\mathbf{q}}) \notin \Omega_{h} \times \Omega$, thus either $\overline{\mathbf{p}} \notin \Omega_{h}$ or $\overline{\mathbf{q}} \notin \Omega$, which yields by the above arguments

$$
\bar{M}_{\lambda, \delta} \leq \max \left\{4 C_{\mathrm{Lip}}^{2} \delta, 4 C_{\mathrm{Lip}} C_{\mathrm{bd}} \delta+C_{\mathrm{bd}}^{\prime} h\right\}
$$

Likewise for $\widetilde{M}_{\lambda, \delta}$ using the assumption that $(\tilde{\mathbf{p}}, \tilde{\mathbf{q}}) \notin \Omega_{h} \times \Omega$. The announced result 25 follows from Lemma C. 2 and the these bounds on $\bar{M}_{\lambda, \delta}$ and $\widetilde{M}_{\lambda, \delta}$.

## D Selling's algorithm

This appendix is devoted to Selling's algorithm [18], which allows in dimension $d \leq 3$ to compute Voronoi's first reduction (5), and the related tensor decomposition (6) which is at the foundation of our numerical scheme. This algorithm is used in all the numerical experiments presented in this paper. For that purpose, we need to introduce a few concepts from lattice geometry. The vertices (resp. edges) of a polyhedron are its 0 -dimensional (resp. 1-dimensional) facets.

Definition D.1. A superbase of $\mathbb{Z}^{d}$ is a $(d+1)$-tuple $b=\left(\mathbf{v}_{0}, \cdots, \mathbf{v}_{d}\right) \in\left(\mathbb{Z}^{d}\right)^{d+1}$ such that $\left|\operatorname{det}\left(\mathbf{v}_{1}, \cdots, \mathbf{v}_{d}\right)\right|=1$ and $\mathbf{v}_{0}+\cdots+\mathbf{v}_{d}=0$. One also defines $M_{b}:=\frac{1}{2} \sum_{i=0}^{d} \mathbf{v}_{i} \mathbf{v}_{i}^{T} \in \mathrm{~S}^{++}(\mathbb{E})$.

In dimension $d \leq 3$, the vertices of Ryskov's polyhedron $\mathcal{P} \subseteq \mathrm{S}^{++}\left(\mathbb{R}^{d}\right)$, defined in (4), are the matrices $M_{b}$, where $b=\left(\mathbf{v}_{0}, \cdots, \mathbf{v}_{d}\right)$ is an arbitrary superbase of $\mathbb{Z}^{d}$. See [17, 59] for a similar classification, in dimension $d \leq 8$, of these vertices under the name of perfect forms. The vertices joined to $M_{b}$ by an edge of the polyhedron $\mathcal{P}$ are in dimension $d=2$ (resp. $d=3$ ) associated with the superbases

$$
\begin{equation*}
b^{\prime}=\left(-\mathbf{v}_{i}, \mathbf{v}_{j}, \mathbf{v}_{i}-\mathbf{v}_{j}\right) \quad\left(\text { resp. } b^{\prime}=\left(-\mathbf{v}_{i}, \mathbf{v}_{j}, \mathbf{v}_{k}+\mathbf{v}_{i}, \mathbf{v}_{l}+\mathbf{v}_{i}\right)\right) \tag{53}
\end{equation*}
$$

where $\{i, j, k\}$ (resp. $\{i, j, k, l\}$ ) is an arbitrary permutation of $\{0, \cdots, d\}$. In addition one has

$$
\begin{equation*}
\operatorname{Tr}\left(M_{b^{\prime}} D\right)=\operatorname{Tr}\left(M_{b} D\right)-C_{d}\left\langle\mathbf{v}_{i}, D \mathbf{v}_{j}\right\rangle \tag{54}
\end{equation*}
$$

where $C_{2}=2$ and $C_{3}=1$. Selling's algorithm, given an input matrix $D$ and a superbase $b=\left(\mathbf{v}_{0}, \cdots, \mathbf{v}_{d}\right)$, looks for indices $0 \leq i<j \leq d$ such that $\left\langle\mathbf{v}_{i}, D \mathbf{v}_{j}\right\rangle>0$. If any exists, then the superbase is replaced with (53) and the process is repeated. Otherwise the algorithm stops, and yields what is called a $D$-obtuse superbase, i.e. $\left\langle\mathbf{v}_{i}, D \mathbf{v}_{j}\right\rangle \leq 0$ for all $i \neq j$.

Clearly, this algorithm is equivalent to a walk on the graph defined by the vertices and the edges of Ryskov's polyhedron $\mathcal{M}_{d}$, in which the next vertex is a neighbor selected so as to reduce the objective function $M \mapsto \operatorname{Tr}(M D)$. Since Voronoi's first reduction (5) is a well posed linear program, the process eventually ends, similarly to the classical simplex algorithm. The vertex $M_{b}$ associated to the eventual $D$-obtuse superbase $b=\left(\mathbf{v}_{0}, \cdots, \mathbf{v}_{d}\right)$ is optimal, and the corresponding Karush-Kuhn-Tucker relations read

$$
D=-\sum_{0 \leq i<j \leq d}\left\langle\mathbf{v}_{i}, D \mathbf{v}_{j}\right\rangle \mathbf{e}_{i j} \otimes \mathbf{e}_{i j},
$$

where $\mathbf{e}_{i j}:=\mathbf{v}_{k}^{\perp}$ in dimension $d=2$ and with $\{i, j, k\}=\{0,1,2\}$, and $\mathbf{e}_{i j}:=\mathbf{v}_{k} \times \mathbf{v}_{l}$ in dimension $d=3$ and with $\{i, j, k, l\}=\{0,1,2,3\}$. See e.g. Lemma 4.4 in [37] for this matrix identity.

The numerical cost of Selling's algorithm, in our applications which involve condition numbers $\lesssim 100$, is low enough to be neglected. For application involving (much) higher condition numbers, it is natural to perform a preliminary basis reduction [40], with cost $\mathcal{O}(\ln \operatorname{Cond}(D))$, after what Selling's algorithm terminates in a single step, see Proposition 1 in [25].

## E Reconstruction of an upwind gradient, for geodesic extraction

Solving eikonal equations is, in many applications, only a means to extract minimal geodesics for the corresponding optimal control problem. These paths are the integral lines of the intrinsic gradient (defined w.r.t. the metric), of the eikonal PDE solution:

$$
\begin{equation*}
\dot{\gamma}(t):=V(\gamma(t)), \quad \text { where } V(\mathbf{p}):=\mathrm{d} \mathcal{F}_{\mathbf{p}}^{*}(\mathrm{~d} u(\mathbf{p})) \tag{55}
\end{equation*}
$$

See the discussion \$1, or Appendix C of [23] We describe in this appendix a robust estimator of the vector field $V$, defined directly from the PDE scheme, and in terms of upwind quantities.

In the case of Riemannian metrics on has $V(\mathbf{p})=\mathcal{D}(\mathbf{p}) \mathrm{d} u(\mathbf{p})$, as follows from the dual metric expression $\mathcal{F}_{\mathbf{p}}^{*}(\hat{\mathbf{p}})=\|\hat{\mathbf{p}}\|_{\mathcal{D}(\mathbf{p})}$ and the eikonal equation. The following proposition, generalizing a similar result in the isotropic case in [50], extracts this direction from the FM-VR1 scheme.

Proposition E.1. For each $\mathbf{p} \in \Omega_{h}$ one has the first order approximation

$$
\begin{equation*}
\mathcal{D}(\mathbf{p}) \mathrm{d} u(\mathbf{p})=h^{-1} \sum_{1 \leq i \leq d^{\prime}} \varepsilon_{i} \rho_{i}(\mathbf{p}) \max \left\{0, u(\mathbf{p})-u\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right), u(\mathbf{p})-u\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)\right\} \dot{\mathbf{e}}_{i}+o(h) \tag{56}
\end{equation*}
$$

where for each $1 \leq i \leq d^{\prime}$ one defines $\varepsilon_{i}=0$ (resp. $\varepsilon_{i}=-1$, resp. $\varepsilon_{i}=1$ ) if the $i$-th maximum is $0\left(\operatorname{resp} . u(\mathbf{p})-u\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right)\right.$, resp. $\left.u(\mathbf{p})-u\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)\right)$.

Proof. We may assume that $u$ is linear, $u(\mathbf{p})=\left\langle\hat{\eta}, \mathbf{p}-\mathbf{p}_{0}\right\rangle$, in which case (56) becomes by 15 )

$$
h^{-1} \sum_{1 \leq i \leq d^{\prime}} \varepsilon_{i} \rho_{i}(\mathbf{p})\left|\left\langle\hat{\eta}, h \dot{\mathbf{e}}_{i}\right\rangle\right| \dot{\mathbf{e}}_{i}=\sum_{1 \leq i \leq d^{\prime}} \rho_{i}(\mathbf{p})\left\langle\hat{\eta}, \dot{\mathbf{e}}_{i}\right\rangle \dot{\mathbf{e}}_{i}=\mathcal{D}(\mathbf{p}) \hat{\eta}=\mathcal{D}(\mathbf{p}) \mathrm{d} u(\mathbf{p})
$$

Rander metrics take the general form $\mathcal{F}(\dot{\mathbf{p}})=\|\dot{\mathbf{p}}\|_{\mathcal{M}}+\langle\hat{\eta}, \dot{\mathbf{p}}\rangle$, omitting for readability the dependency of $\mathcal{F}, \mathcal{M}$ and $\hat{\eta}$ on the base point $\mathbf{p}$. Interestingly, algebraic structure of the dual metric is similar: $\mathcal{F}^{*}(\hat{\mathbf{p}})=\|\hat{\mathbf{p}}\|_{\mathcal{M}^{*}}+\left\langle\hat{\mathbf{p}}, \dot{\eta}^{*}\right\rangle$, where $\left(\mathcal{M}^{*}, \dot{\eta}^{*}\right)$ is related to $(\mathcal{M}, \hat{\eta})$ by the relation:

$$
\left(\begin{array}{cc}
\mathcal{M} & \hat{\eta}^{\mathrm{T}}  \tag{57}\\
\hat{\eta} & 1
\end{array}\right)\left(\begin{array}{cc}
\mathcal{M}^{*} & \dot{\eta}^{*} \\
\left(\dot{\eta}^{*}\right)^{\mathrm{T}} & 1
\end{array}\right)=\alpha^{2} \mathrm{Id}
$$

for some $\alpha=\alpha(\mathcal{M}, \hat{\eta})$. See also [35] for equivalent explicit formulas. As a result, the intrinsic gradient 55, right) reads

$$
V(\mathbf{p})=\frac{\mathcal{M}^{*}(\mathbf{p}) \mathrm{d} u(\mathbf{p})}{\|\mathrm{d} u(\mathbf{p})\|_{\mathcal{M}^{*}(\mathbf{p})}}+\dot{\eta}^{*}(\mathbf{p})
$$

Finally, an estimation of the differential $\mathrm{d} u(\mathbf{p})$ is deduced as in the Riemannian case from the numerical scheme, inferred from the analoguous first order approximation: $\mathcal{D}(\mathbf{p})(\mathrm{d} u(\mathbf{p})-\hat{\eta}(\mathbf{p}))=$
$h^{-1} \sum_{1 \leq i \leq d^{\prime}} \varepsilon_{i} \rho_{i}(\mathbf{p}) \max \left\{0, u(\mathbf{p})-u\left(\mathbf{p}+h \dot{\mathbf{e}}_{i}\right)+h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle, u(\mathbf{p})-u\left(\mathbf{p}-h \dot{\mathbf{e}}_{i}\right)-h\left\langle\hat{\eta}(\mathbf{p}), \dot{\mathbf{e}}_{i}\right\rangle\right\} \dot{\mathbf{e}}_{i}+o(h)$.


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[^1]:    ${ }^{1}$ Except for the numerical scheme devoted to Rander metrics

[^2]:    ${ }^{2}$ Except the variant devoted to Rander metrics.

[^3]:    ${ }^{3}$ Formally, $\mathbb{L}^{*}=\left\{\hat{\mathbf{p}} \in \mathbb{E}^{*} ; \forall \dot{\mathbf{p}} \in \mathbb{L},\langle\hat{\mathbf{p}}, \dot{\mathbf{p}}\rangle \in \mathbb{Z}\right\}$.

[^4]:    ${ }^{4} \mathrm{~A}$ metric lacking the symmetry property, i.e. $\mathcal{F}_{\mathbf{p}}(\dot{\mathbf{p}}) \neq \mathcal{F}_{\mathbf{p}}(-\dot{\mathbf{p}})$ for some point $\mathbf{p} \in \bar{\Omega}$ and vector $\dot{\mathbf{p}} \in \mathbb{E}$, is usually referred to as a quasi-metric. The quasi- prefix is however dropped in this paper, for consistency with the Riemannian and sub-Riemannian cases.

[^5]:    ${ }^{5}$ Numerical codes for Rander metrics, which are more experimental, are available on demand.
    ${ }_{6}$ https://github.com/Mirebeau

[^6]:    ${ }^{7}$ Without this rotation, the riemannian metric anisotropy is mostly aligned with the coordinate axes, which makes the problem significantly easier numerically, see the numerical experiments in 34 .

[^7]:    ${ }^{8}$ The Reeds-Shepp model more faithfully describes a a wheelchair actually, whereas an actual car is better described by the combination of a wheelchair and of a trailer.

[^8]:    ${ }^{9}$ Two dimensional instances do not show similar CPU time reduction, because numerical cost is in that case dominated by the handling of a priority queue, rather than by geometrical computations.

[^9]:    ${ }^{10}$ Note that in [34, 35] the acronym AGSI refers to discretization FE of 9. This denomination is not proper, since it ties the PDE discretization with the numerical solver used for the discrete system of equations.

[^10]:    ${ }^{11}$ The implementation details of the fast marching method suggest additional restrictions, discussed in [54], such as requiring upwinding $u(\mathbf{p}+2 \dot{\mathbf{e}}) \leq u(\mathbf{p})$, and limiting the use of second order finite differences to a recomputation stage executed when a point of the front propagation is accepted and its value is frozen.

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    ${ }^{13} \mathrm{PhD}$ student under the direction of R. Duits at TU/e University, Eindhoven.
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