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# Low-rank tensor reconstruction of concentrated densities with application to Bayesian inversion

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## Low-rank tensor reconstruction of concentrated densities with application to Bayesian inversion

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ABSTRACT. This paper presents a novel method for the accurate functional approximation of possibly highly concentrated probability densities. It is based on the combination of several modern techniques such as transport maps and low-rank approximations via a nonintrusive tensor train reconstruction. The central idea is to carry out computations for statistical quantities of interest such as moments based on a convenient reference density for which accurate numerical methods can be employed. Since the exact transport from target to reference can usually not be determined exactly, one has to cope with a perturbed reference density due to a numerically approximated transport map. By the introduction of a layered approximation and appropriate coordinate transformations, the problem is split into a set of independent approximations in seperately chosen orthonormal basis functions, combining the notions h- and p-refinement (i.e. "mesh size" and polynomial degree). An efficient low-rank representation of the transport is achieved via the Variational Monte Carlo (VMC) method. This nonintrusive regression technique reconstructs the map in the tensor train format. An a priori convergence analysis with respect to the error terms introduced by the different (deterministic and statistical) approximations in the Hellinger distance and the Kullback-Leibler divergence is derived. Important applications are presented and in particular the context of Bayesian inverse problems is illuminated which is a main motivation for the developed approach. Several numerical examples illustrate the efficacy with densities of different complexity and degrees of perturbation of the transport to the reference density. The (superior) convergence is demonstrated in comparison to Monte Carlo and Markov-Chain Monte Carlo methods.

#### 1. OVERVIEW

We derive a novel numerical method for the functional representation of complicated (in particular highly concentrated) probability densities. This difficult task usually is attacked with Markov-Chain Monte Carlo (MCMC) methods which yield samples of the posterior. Despite their popularity, the convergence rate of these methods is ultimately limited by the employed Monte Carlo sampling technique, see e.g. [10] for recent multilevel techniques in this context. Moreover, practical issues e.g. regarding the initial number of samples (burn-in) arise.

In this work, we propose a new approach based on *function space representations with efficient surro- gate models* in several instances. This is motivated by our previous work on adaptive low-rank approximations of solutions of parametric random PDEs with Adaptive Stochastic Galerkin FEM (ASGFEM, see e.g. [17, 14]) and in particular the sampling-free Bayesian inversion presented in [15] where the setting of uniform random variables was examined. A generalization to the important case of Gaussian random variables turns out to be non-trivial from a computational point of view due to the difficulties of representing highly concentrated densities in a compressing tensor format which in fact is required in order to cope with the high dimensionality of the problem. As a consequence, we develop a discretization approach which takes into account the potentially problematic structure of the probability density at hand by a combination of several transformations and approximations that can be chosen adaptively to counteract the interplay of the employed numerical approximations. With the computed functional representation of the density, the evaluation of moments or other statistical quantities of interest can be carried out efficiently and with high accuracy.

A central idea of the method is to obtain a map which transports the target density to some convenient reference density. For the latter, accurate numerical methods are assumed to be available, hence simplifying any subsequent computations once such a transport map is determined. Transport maps for probablity densities are a classical topic in mathematics. They are under active research in particular in the area of optimal transport [48, 41] and also have become popular in current machine learning research [47, 39, 9]. A main application we have in mind is Bayesian inversion where, given a prior

density and some observations of the forward model, the posterior density of the sought parameters should be determined. In this context, the rescaling approaches in [42, 43] based on the Laplace approximation can be considered as transport maps of a certain (affine) form. More general transport maps have been examined extensively in [19, 36, 37] and other works of the group. In the optimal case that the exact transport is known and a chosen multivariate Gaussian reference density, a subsequent polar transformation results in 1D approximation problems of analytically smooth functions. By projection onto a trigonometric or polynomial basis, a very accurate functional representation of the density could then be obtained. However, it is obvious that by assuming an accurate transport, the difficulty is shifted from coping with a complicated density to the computation of a complicated map which of course usually is also infeasible. We hence suppose that only an inexact transport is available, leading to perturbed densities. The degree of perturbation has then to be coped with in the subsequent approximation steps. Similar to ideas in adaptive FEM, in addition to the selection of (local) approximation spaces of a certain degree ("p-refinement"), we introduce a spatial decomposition of the density representation into layers ("h-refinement") around some center of mass of the considered density. This enables to exploit the decay behaviour of the approximated density, in the best case in a one dimensional low-rank setting. Overall, this "hp-refinement" allows to balance incaccuracies and hence perturbations of the reference density by putting more effort into the discretization part. One hence has the freedom to decide whether more effort should be invested into computing an exact transport map or into a more elaborate discretization (with more layers and larger bases) of the transported and transformed density.

For eventual computations with the devised (possibly high-dimensional) functional density representation, an efficient representation format is required. In our context, hierarchical tensors and in particular tensor trains (TT) prove to be advantageous, cf. [2, 35]. These compressing formats enable to alleviate the curse of dimensionality under suitable conditions and allow for efficient evluations of very high-dimensional objects. We aim to obtain a low-rank tensor representation of the respective transport map, for each layer of the discretization. In certain cases it can be shown that these maps are even of rank one. In more general cases, a low-rank representability may be observed numerically. To generate a tensor train representation of the maps (in the extended form, i.e. coupled with a function basis), the Variational Monte Carlo (VMC) method [18] is employed. It basically is a tensor regression approach based on function samples for which a convergence analysis is available. Notably, depending on the chosen loss functional, it leads to the best approximation in the respective model space. It has previously been examined in the context of random PDEs in [18] as an alternative nonintrusive numerical approach to Adaptive Stochastic Galerkin FEM in the TT format [17, 14]. The approximation of [14] is in fact used in one of the presented examples for Bayesian inversion with the random Darcy equation with lognormal coefficient. We note that surrogate models of the forward model have been used in the context of MCMC e.g. in [30] and tensor representations (obtained by cross approximation) were used in [11] to improve the efficiency of MCMC sampling.

The derivation of our method is supported by an a priori convergence analysis with respect to the Hellinger distance and the Kullback-Leibler divergence. In the analysis, different error sources have to be considered, in particular a layer truncation error depending on decay properties of the density, low-rank representability, rank truncation errors and the perturbations introduced by inexact transport maps, leading to perturbed reference densities. Moreover, the VMC error analysis [18] comprising statistical estimation and numerical approximation errors is adjusted to be applicable to the devised approach. While not usable for an a posterior error control in its current initial form, the derived analysis leads the way to more elaborate results for this promising method in future research.

With the constructed functional density surrogate, sampling-free computations of statistical quantities of interest such as moments or marginals become feasible by fast tensor contractions, even for highly concentrated or (depending on the available transport map) nonlinearly transformed high-dimensional

densities. For non-smooth quantities simple sample generation based on inverse transformation sampling approaches from the uniform distribution are easily possible due to the compactness of the discretisation layers.

While several assumptions have to be satisfied for this to work most efficiently, the approach is rather general and can be further adapted to the considered problem. Moreover, it should be emphasized that by obtaining a functional representation, structural properties of the density at hand (in particular smoothness, sparsity, low-rank approximability and decay behaviour in different parameters) can be exploited in a much more extensive way than what is possible with sampling based methods such as MCMC, leading to more accurate and faster converging statistical computations. We note that the perturbed posterior surrogate can be used to efficiently generate samples for rejection sampling or within a MCMC scheme. Since the perturbed transport can be seen as a preconditioner, the sample generation can be based on the perturbed prior. These samples can then be pushed forward to the posterior. As a prospective extension, the constructed posterior density can directly be used in a Stochastic Galerkin FEM based on the integral structure, closing the loop of forward and inverse problem, resulting in the inferred forward problem with model data determined by Bayesian inversion from the observed data.

The structure of the paper is as follows. Section 2 is concerned with the representation of probability densities and introduces a relation between a target and a reference density. Such a transport map can be determined numerically by approximation in a chosen class of functions and with an assumed structure, leading to an the concept of perturbed reference densities. To counter the perturbation, a layered truncated discretization is introduced. An efficient low-rank representation of the mappings is described in Section 3 where the tensor train format is discussed. In order to obtain this nonintrusively, the Variational Monte Carlo (VMC) tensor reconstruction is reviewed. A priori convergence results with respect to the Hellinger distance and Kullback-Leibler divergence are derived in Section 4. For practical purposes, the proposed method is described in terms of an algorithm in Section 5. Possible applications we have in mind are examined in Section 6. In particular, the setting of Bayesian inverse problems is recalled. Moreover, the computation of moments and marginals is scrutinized. Section 7 illustrates the performance of the proposed method. In addition to an examination of the numerical sensitivity of the accuracy with respect to the perturbation of the transport maps, a typical model problem from Uncertainty Quantification (UQ) is depicted, namely the identification of a parametrization for the random Darcy equation with lognormal coefficient given as solution of a stochastic Galerkin FEM.

#### 2. Density representation

The aim of this section is to introduce the central ideas of the proposed approximation of densities. For this task, two established concepts are reviewed, namely *transport maps* [19, 5], which are closely related to the concept of optimal transport [48, 41], and *hierarchical low-rank tensor representations* [35, 26, 2]. By the combination of these techniques, assuming the access to a suitable transformation, the developed approach yields a functional representation of the density in a format which is suited to computations with high-dimensional functions. In particular, we are able to handle highly concentrated posterior densities, e.g. appearing in the context of Bayesian inverse problems. While transport maps on their own in principle enable the generation of samples of some target distribution, the combination with a functional low-rank representation allows for integral quantities such as (centered) moments to become computable. Given an approximate transport map, the low-rank representation can be seen as a further approximation step (improving the inaccuracy of the used transport) to gain direct access to the target density.

Consider a target measure  $\pi$  with Radon-Nikodym derivative with respect to to the Lebesgue measure  $\lambda$  denoted as f with support in  $\mathbb{R}^d$ ,  $d < \infty$ , i.e.

(1) 
$$f(y) := \frac{\mathrm{d}\pi}{\mathrm{d}\lambda}(y), \quad y \in Y := \mathbb{R}^d.$$

In the following we assume that point evaluations of f are available up to a multiplicative constant, motivated by the framework of Bayesian posterior density representation with unknown normalization constant. Furthermore, let  $\pi_0$  be some reference measure exhibiting a Radon-Nikodym derivative with respect to to the Lebesgue measure denoted as  $f_0$ . This is motivated by the prior measure and density in the context of Bayesian inference.

2.1. **Transport Maps.** The notion of density transport is classical and with optimal transport has become a popular field recently, see e.g. [48, 41]. It has been employed to improve numerical approaches for Bayesian inverse problems e.g. in [19, 5, 11]. Similar approaches are discussed in terms of sample transport e.g. for Stein's method [31, 8] or multi-layer maps [5]. We review the properties required for our approach in what follows. Note that since our target application is Bayesian inversion, instead of the more general reference and target densities, we usually use the terms prior and posterior.

Let  $X:=\mathbb{R}^d$  and assume that there exists an exact transport map

$$(2) T: X \to Y,$$

which is a diffeomorphism<sup>1</sup> that relates  $\pi$  and  $\pi_0$  via pullback, i.e.

(3) 
$$f_0(x) = f(T(x))|\det \mathcal{J}_T(x)|, \quad x \in X.$$

Then, computations might be carried out in terms of the measure  $\pi_0$ , which is commonly assumed to be of a simpler structure. For instance the moment computation with respect to some multiindex  $\alpha$  reads as follows,

(4) 
$$\int_{Y} y^{\alpha} d\pi(y) = \int_{Y} T(x)^{\alpha} d\pi_{0}(x) = \int_{Y} T(x)^{\alpha} f_{0}(x) d\lambda(x).$$

Note that the computation of the right-hand side in (4) may still be a challenging task depending on the actual structure of T. In [19] T is expanded in chaos polynomials with respect to  $\pi_0$ . From a practical point of view, this provides access to lower-order moments using orthogonality of the underlying polynomial system.

Here we follow an alternative strategy with the aim to efficiently compute moments of some target density based on a functional representation. Notably we assume a convenient (simple) structures of T with the potential drawback of reduced accuracy, i.e. an inexact (pull-back) transport from the target to an auxiliary (instead of the exact reference) density. Motivated by the Bayesian context, we call such a pull-back of some posterior density the *perturbed prior density*, see Section 2.2. Given a simple transport structure, the possibly demanding computational task is shifted to the accurate approximation of the perturbed prior. For this, there is justified hope of feasibility in some appropriate (alternative) coordinate system. In order to tackle moment computations, other posterior statistics or to generate posterior samples, we hence devise a numerical approach that enables a *workload balancing between the reconstruction of some problem-dependent transport structure and the accurate evaluation of the perturbed prior*. In the following we list some examples of transport maps.

 $<sup>^{</sup>m 1}$ note that the requirements on T can be weakened, e.g. to local Lipschitz

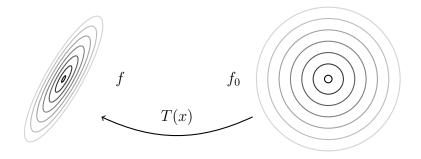


FIGURE 1. Illustration of affine transport: translation, rotation and rescaling.

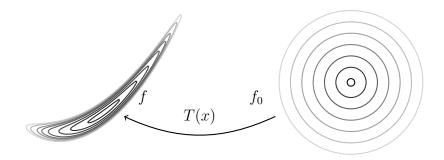


FIGURE 2. Illustration of quadratic transport: affine properties and bending.

2.1.1. Affine transport. In [42, 43] the authors employ an affine linear preconditioning for acceleration of MCMC or sparse-grid integration in the context of highly informative and concentrated Bayesian posterior densities, using a s.p.d. matrix  $H \in \mathbb{R}^{d,d}$  and  $M \in \mathbb{R}^d$ . In the mentioned articles, up to a multiplicative constant, H corresponds to the inverse square root of the Hessian at the MAP M, i.e. the location of the local optimum of an Laplace approximation of the posterior density. This rather simple construction, under the assumption of an unimodal density, leads to stable numerical algorithms for the computation of quantities of interest as the posterior mass concentrates. In the framework of push-forward of a reference density  $f_0$  to a target density  $f_0$  this concept coincides with an affine transport

$$(5) y = T(x) = Hx + M, \quad x \in X.$$

In the transport setting H and M may be computed for instance via some minimisation of the Kullback-Leibler divergence as in [19]. Note that H and M do not necessarily have to be the inverse square root of the Hessian or the MAP. Figure 1 illustrates the concept of an affine transport.

2.1.2. Quadratic transport. A more general class of polynomial transport exhibits the form

(6) 
$$T(x) = \frac{1}{2}x : A : x + Hx + M, \quad x \in X,$$

with  $A \in \mathbb{R}^{d,d,d}$ ,  $H \in \mathbb{R}^{d,d}$ ,  $M \in \mathbb{R}^d$ . Such a quadratic transport may be used for simple curve structures as depicted in Figure 2.

2.1.3. More general transport maps. The parametrisation of transport maps can be chosen quite liberally as long as certain criteria are satisfied, which are either directly imposed in the ansatz space  $\mathcal T$  of the maps or added as constraints during optimization. In particular, the approximate transport map has to be invertible, which can be ensured by requiring a positive Jacobian. A commonly used

measure for optimization is the Kullback-Leibler divergence<sup>2</sup> leading to the optimization problem

(7) 
$$\min_{T \in \mathcal{T}} d_{\mathrm{KL}}(Y; T\pi_0, \pi)$$
 such that  $\det \nabla T > 0$   $\pi$ -almost everywhere.

Several suggestions regarding simplifications and special choices of function spaces  $\mathcal{T}$  such as smooth triangular maps based on higher-order polynomials or radial basis functions can for instance be found in the review article [19]. An interesting idea is to subdivide the task into the iterative computation of simple correction maps which are then composed as proposed in [5]. We again emphasize that while an accurate transport map is desirable, any approximation of such a map can in principle be used with the proposed method. In fact one can decide whether it is beneficial to spend more effort on the approximation of the perturbed density or on a better representation of the transport.

2.2. Inexact transport and the perturbed prior. In general, the transport map T is unknown or difficult to determine and hence has to be approximated by some  $\tilde{T}\colon X\to Y$ , e.g. using a polynomial chaos representation with respect to  $\pi_0$  [19] or with a more advanced composition of simple maps in a reduced space such as in [5]. As a consequence, it holds

(8) 
$$\int_{Y} y^{\alpha} d\pi(y) \approx \int_{X} \tilde{T}(x)^{\alpha} d\pi_{0}(x)$$

subject to the accuracy of the involved approximation of T. One can also view  $\tilde{T}$  as the push-forward of some measure  $\tilde{\pi}_0$  with density  $\tilde{f}_0$  to  $\pi$  given by

(9) 
$$\tilde{f}_0(x) = f(\tilde{T}(x))|\det \mathcal{J}_{\tilde{T}}(x)|.$$

We henceforth refer to (9) as the auxiliary reference density or, motivated by Bayes framework, as the *perturbed prior density*. Using this construction, the moment computation reads

(10) 
$$\int_{Y} y^{\alpha} d\pi(y) = \int_{X} \tilde{T}(x)^{\alpha} d\tilde{\pi}_{0} = \int_{X} \tilde{T}(x)^{\alpha} \tilde{f}_{0}(x) d\lambda(x).$$

If one would know  $\tilde{f}_0$ , by (9) and (10) one would also have access to the exact posterior.

Equation (10) is the starting point of a follow-up strategy by approximating  $\tilde{f}_0$  in another coordinate system which is better adapted to the structure of the approximate (perturbed) prior. Consider a (fixed) diffeomorphism

(11) 
$$\Phi \colon \hat{X} \subset \mathbb{R}^d \to X, \quad \hat{x} \mapsto x = \Phi(\hat{x})$$

with Jacobian  $\hat{x} \mapsto |\det \mathcal{J}_{\Phi}(\hat{x})|$  and define the *perturbed transformed prior* 

(12) 
$$\hat{f}_0 \colon \hat{X} \mapsto \mathbb{R}_+, \quad \hat{x} \mapsto \hat{f}_0(\hat{x}) := \tilde{f}_0(\Phi(\hat{x})).$$

In case (12) can be approximated accurately by some function  $\tilde{f}_0^h$  then

(13) 
$$\int_{Y} y^{\alpha} d\pi(y) \approx \int_{\hat{Y}} \tilde{T}(\Phi(\hat{x}))^{\alpha} \tilde{f}_{0}^{h}(\hat{x}) |\det \mathcal{J}_{\Phi}(\hat{x})| d\lambda(\hat{x})$$

with accuracy determined only by the approximation quality of  $\tilde{f}_0^h$ . Thus, (12) and (13) enable a balancing between the construction of the transport map approximation  $\tilde{T}$  of T to shift its complexity given the underlying diffeomorphism  $\Phi$  to the approximation of (12) in a new coordinate system intrinsic to  $\hat{X}$ .

<sup>&</sup>lt;sup>2</sup>although in machine learning Wasserstein or Sinkhorn distances have become very popular when so-called normalising flows are computed

The construction of  $\tilde{T}$  and a suiteable map in (11) may be used to obtain a convenient transformed auxillary reference density giving in (12). An approximation thereof can be significantly simpler compared to a possibly complicated and concentrated target density f or the computation of the exact transport T. This e.g. is satisfied if

- $f_0$  is a Gaussian density and  $\tilde{T}$  maps f to  $\tilde{f}_0$  which is in some sense near to a Gaussian density. In this case,  $\Phi$  from (11) may be chosen as the d-dimensional spherical transformation and extended low-rank tensor formats are employed as approximation class, see Section 3. In this setting, the introduction of an adapted coordinate system allows to shift the exponential decay to the one dimensional radial parameter. The accuracy of an approximation can then be improved easily by additional h-refinements as described in Section 2.3.
- The reference density  $f_0$  has a complicated form and might be replaced by  $\tilde{f}_0$  to become computationally accessible.

In the following we state an important property that needs to be fulfilled by the perturbed prior  $f_0$  in order to lead to a convergent method with the employed approximations.

**Definition 2.1.** (outer polynomial exponential decay) A function  $\tilde{f}_0\colon X\to\mathbb{R}^+$  has outer polynomial exponential decay if there exists a simply connected compact  $K\subset X$  with, a polynomial  $\pi^+$  being positive on  $X\setminus K$  and a C>0, s.t.

(14) 
$$\tilde{f}_0(x) \le C \exp\left(-\pi^+(x)\right), \quad x \in X \setminus K.$$

2.3. Layer based representation. To further refine and motivate the notion of an adapted coordinate system, let  $L \in \mathbb{N}$  and  $(X^{\ell})_{\ell=1}^{L}$  be pairwise disjoint domains in X s.t.

$$K := \bigcup_{\ell=1}^{L} \overline{X^{\ell}}$$

is simply connected and compact and define  $X^{L+1}:=X\setminus K$ . Then, for given  $L\in\mathbb{N}$  we may decompose the perturbed prior  $\tilde{f}_0$  as

(16) 
$$\tilde{f}_0(x)=\sum_{\ell=1}^{L+1}\tilde{f}_0^\ell(x)\quad \text{with}\quad \tilde{f}_0^\ell:=\chi_\ell\tilde{f}_0,$$

where  $\chi_\ell$  denotes the indicator function on  $X^\ell$ . Moreover, for any tensor set  $\hat{X}^\ell := \mathop{\textstyle \bigvee}_{i=1}^d \hat{X}_i^\ell$  and diffeomorphism  $\Phi^\ell \colon \hat{X}^\ell \mapsto X^\ell$ ,  $1 \le \ell \le L+1$ , we may represent the *localised perturbed prior*  $\tilde{f_0}^\ell$  as a pull-back function

(17) 
$$\tilde{f}_0^\ell = \hat{f}_0^\ell \circ \varPhi^{\ell-1},$$

where  $\hat{f}_0^\ell$  is a map defined on  $\hat{X}^\ell$  as in (12). We consider the following illustrative example.

#### Example 2.2. (multivariate polar transformation)

The d-dimensional spherical coordinate system allows for simple layer layouts in terms of hyperspherical shells. In particular, for  $\ell=1,\ldots,L+1<\infty$ , with  $0=\rho_1<\rho_2<\ldots<\rho_{L+1}<\rho_{L+2}=\infty$ , let

(18) 
$$\hat{X}^{\ell} := [\rho_{\ell}, \rho_{\ell+1}] \times [0, 2\pi] \times \sum_{i=2}^{d-2} [0, \pi], \quad X^{\ell} := B_{\rho_{\ell+1}}(0) \setminus B_{\rho_{\ell}}(0) \subset X,$$

i.e.  $\hat{X}^{\ell}$  and  $X^{\ell}$  denote the corresponding adopted (transformed) and the original parameter space, respectively. Then, for  $\hat{x}=(\rho,\theta_0,\pmb{\theta})\in\hat{X}$ ,  $\pmb{\theta}=(\theta_1,\ldots,\theta_{d-2})$ , the polar transformation  $\Phi^{\ell}:\hat{X}^{\ell}\to \mathbb{R}$ 

 $X^{\ell}$  reads

(19) 
$$\Phi^{\ell}(\hat{x}) = \Phi^{\ell}(\rho, \theta_0, \boldsymbol{\theta}) = \rho \begin{bmatrix} \cos \theta_0 \sin \theta_1 \sin \theta_2 \cdots \sin \theta_{d-3} \sin \theta_{d-2} \\ \sin \theta_0 \sin \theta_1 \sin \theta_2 \cdots \sin \theta_{d-3} \sin \theta_{d-2} \\ \cos \theta_1 \sin \theta_2 \cdots \sin \theta_{d-3} \sin \theta_{d-2} \\ \cos \theta_2 \cdots \sin \theta_{d-3} \sin \theta_{d-2} \\ \vdots \\ \cos \theta_{d-3} \sin \theta_{d-2} \\ \cos \theta_{d-2} \end{bmatrix}$$

Moreover, the Jacobian is given by

(20) 
$$\det \mathcal{J}_{\Phi^{\ell}}(\rho, \theta_0, \boldsymbol{\theta}) = \rho^{d-1} \prod_{i=1}^{d-2} \sin^i \theta_i.$$

The advantage of employing this layer based coordinate change is twofold.

- The definition on a finite set of layers enables a representation of the density on bounded domains. Even though the effective last remainder layer is unbounded, we assume that K is sufficiently large to cover a high percentage of the probability mass of  $\tilde{f}_0$ .
- The chosen coordinate system can favour specific approximation schemes. For instance, we present a class of low-rank formats which are able to represent with linear complexity with respect to dimension *d* while circumventing typical disadvantages such as the requirement of large ranks.

Up to this point, the choice of transformation  $\Phi^{\ell}$ ,  $\ell=1,\ldots,L+1$ , is fairly general. However, for the further development of the method we assume the following property.

#### Definition 2.3. (rank 1 stability)

Let  $\mathcal{X}, \hat{\mathcal{X}} = \underset{i=1}{\overset{\hat{d}}{\times}} \hat{\mathcal{X}}_i \subset \mathbb{R}^d$  be open and bounded sets. A diffeomorphism  $\Phi \colon \hat{\mathcal{X}} \mapsto \mathcal{X}$  is called rank 1 stable if  $\Phi$  and the absolute value of its Jacobian  $\det \mathcal{J}_{\Phi}$  have rank 1, i.e. there exists univariate functions  $\Phi_i \colon \hat{\mathcal{X}}_i \to \mathcal{X}$ ,  $h_i \colon \hat{\mathcal{X}} \to \mathbb{R}$ ,  $i=1,\ldots,d$ , such that for  $\hat{x} \in \hat{\mathcal{X}}$ 

(21) 
$$\Phi(\hat{x}) = \prod_{i=1}^d \Phi_i(\hat{x}_i), \quad |\det \mathcal{J}_{\Phi}(\hat{x})| = \prod_{i=1}^d h_i(\hat{x}_i).$$

Proposition 2.4. The multivariate polar coordinate transformation from Example 2.2 is rank 1 stable.

Due to the notion of rank 1 stable transformations, the map  $\hat{x} \mapsto T(\Phi(\hat{x}))$  in (13) inherits the rank structure of T, see section 3. Furthermore, since the Jacobian  $\hat{x} \mapsto |\det \mathcal{J}_{\Phi}(\hat{x})|$  is rank 1, we can construct tensorized orthonormal basis functions which may be used to approximate the perturbed transformed prior in (12).

**Remark 2.5.** The described concept can be extended to any rank  $r \in \mathbb{N}$  Jacobians of  $\Phi$ , i.e.

(22) 
$$|\det \mathcal{J}_{\varPhi}(\hat{x})| = \sum_{k=1}^{r} \prod_{i=1}^{d} h_{i,k}(\hat{x}_{i}).$$

Motivated by the right hand side in (13), one might use several approximations of the perturbed transformed prior  $\tilde{f}_0 \circ \Phi$  in r distinct tensorized spaces, each associated to the rank 1 weight  $\prod_{i=1}^d h_{i,k}$ .

2.4. **Layer truncation.** This paragraph is devoted to the treatment of the last (remainder or truncation) layer introduced in (16) with the aim to suggest some approximation choices.

If  $\tilde{f}_0$  is represented in the layer format (16), it is convenient to simply extend the function to zero after layer  $L \in \mathbb{N}$ . By this, the remaining (possibly small) probability mass is neglected. Such a procedure is typically employed in numerical applications and does not impose any computational issues since events on the outer truncated domain are usually exponentially unlikely for truncation value chosen sufficiently large. Nevertheless, in order to present a rigorous treatment, we require properties like absolute continuity, which would be lost by using a cut-off function. Inspired by [43] regarding the information limit of unimodal posterior densities<sup>3</sup>, we suggest a Gaussian approximation for the last layer L+1 on the unbounded domain  $X^{L+1}$ , i.e. for some s.p.d.  $\Sigma \in \mathbb{R}^{d,d}$  and  $\mu \in \mathbb{R}^d$  we define the *hybrid representation of the perturbed prior* by

(23) 
$$\tilde{f}_0^{\text{Trun}}(x) := \frac{1}{C_L^{<} + C_L^{>}} \left\{ \begin{array}{l} \tilde{f}_0^{\ell}(x), & x \in X^{\ell}, \ell = 1, \dots, L, \\ f_{\Sigma, \mu}(x), & x \in X^{L+1}. \end{array} \right.$$

Here.

(24) 
$$C_L^{<} := \int\limits_{X \setminus K} f_{\Sigma,\mu}(x) \,\mathrm{d}\lambda(x),$$

(25) 
$$C_L^> := \sum_{\ell=1}^L \int\limits_{X^\ell} \tilde{f}_0^\ell(x) \,\mathrm{d}\lambda(x),$$

and  $f_{\Sigma,\mu}$  denotes the Gaussian probability density function with mean  $\mu$  and covariance matrix  $\Sigma$ .

Remark 2.6. A good choice for  $\mu$  and  $\Sigma$  would be the mean and covariance of the exact perturbed prior  $\tilde{f}_0$ , which however is not accessible a priori. Thus, in numerical simulations one may choose  $\mu$  and  $\Sigma$  as (centralised) moments of the normalised truncated perturbed prior density  $f_0^{\tilde{\text{Trun}}}|_K$  or as the MAP point and the corresponding square root of the numerically computed Hessian as a result of an optimization algorithm on  $\tilde{f}_0$ .

**Lemma 2.7.** (truncation error) For  $\mu \in \mathbb{R}^d$  and  $\Sigma \in \mathbb{R}^{d,d}$  let  $\tilde{f}_0$  have outer polynomial exponential decay with positive polynomial  $\tilde{\pi}^+$  and  $\tilde{C}>0$  with  $K=\overline{B_R(\mu)}$  for some R>0. Then, for  $C_\Sigma=1/2\lambda_{\min}(\Sigma^{-1})$  there exists  $C=C(\tilde{C},\Sigma,d,C_\Sigma)>0$  s.t. it holds

(26) 
$$\|\tilde{f}_0 - \tilde{f}_0^{\text{Trun}}\|_{L^1(X \setminus K)} \le C \left( \|\exp(-\tilde{\pi}^+)\|_{L^1(X \setminus K)} + \Gamma\left(d/2, C_{\Sigma}R^2\right) \right)$$

and

(27) 
$$\left| \int_{X \setminus K} \log \left( \frac{\tilde{f}_0}{f_{\Sigma,\mu}} \right) \tilde{f}_0 \, \mathrm{d}x \right| \leq \int_{X \setminus K} \left( \frac{1}{2} \|x\|_{\Sigma^{-1}}^2 + \tilde{\pi}^+(x) \right) e^{-\tilde{\pi}^+(x)} \, \mathrm{d}\lambda(x)$$

with the incomplete Gamma function  $\Gamma$ .

*Proof.* The proof follows immediately from the definition of  $ilde{f}_0^{\mathrm{Trun}}$ .

In the case that the perturbed prior is close to a Gaussian standard normal distribution, it holds  $c \approx 1$ .

Note that the normalisation constant in (24) may exhibit an analytical form whereas computing (25) suffers from the curse of dimensionality and is in general not available. To circumvent this issue and render further use of the representation (23) feasible, we introduce a suitable low-rank approximation in the next section.

<sup>&</sup>lt;sup>3</sup>A results of [43] is that under suitable conditions the posterior distribution converges to a Gaussian in the limit of zero noise and infinite measurements.

#### 3. LOW-RANK TENSOR TRAIN FORMAT

The computation of high-dimensional integrals and efficient construction of surrogates is a challenging task with a multitude of approaches. Some of these techniques are sparse grid methods [6, 23], collocation [20, 33, 22] or modern sampling techniques [24, 40, 32]. As motivated by equation (25), we aim for a model to adequately approximate the localised perturbed prior maps  $\tilde{f}_0^\ell$ . The idea to introduce an adopted coordinate system is to enable the exploitation of additional structure for low-rank representations such as the tensor train (TT) format [35, 27, 26]. Hence, this section is concerned with the concept of low-rank tensor trains in possibly infinite dimensional spaces. We highlight a black-box ("non-intrusive") sample-based approach to obtain such a representation of arbitrary maps via the so-called Variational Monte Carlo method [18].

Let  $\hat{X} = \bigotimes_{i=1}^d \hat{X}_i$  be a tensor space of separable Banach spaces  $\hat{X}_i, i \in [d] := \{1, \dots, d\}$ , and consider a map  $g \colon \hat{X} \to \mathbb{R}$ . The function g can be represented in the *tensor train* (TT) format if there exists a *rank vector*  $\mathbf{r} = (r_1, \dots, r_{d-1}) \in \mathbb{N}^{d-1}$  and univariate functions  $g^i[k_{i-1}, k_i] \colon \hat{X}_i \to \mathbb{R}$  for  $k_i \in [r_i], i \in [d]$ , such that for all  $\hat{x} \in \hat{X}$ 

(28) 
$$g(\hat{x}) = \sum_{k=1}^{r} \prod_{i=1}^{d} g^{i}[k_{i-1}, k_{i}](\hat{x}_{i}) \quad \text{with} \quad \mathbf{k} := (k_{1}, \dots, k_{d-1}).$$

For ease of notation it is convenient to set  $k_0=k_d=1$ . In the forthcoming sections we consider weighted tensorized Lebesgue spaces. In particular, for a non-negative weight function  $w\colon \hat{X}\to \mathbb{R}$  with  $w=\bigotimes_{i=1}^d w_i, w\in L^1(\hat{X})$ , define the tensorization of  $L^2(\hat{X},w)=\bigotimes_{i=1}^d L^2(\hat{X}_i,w_i)$  by

(29) 
$$V(\hat{X}) := L^2(\hat{X}, w) = \left\{ v \colon \hat{X} \to \mathbb{R} \mid \|v\|_{\mathcal{V}}^2 := \int_{\hat{X}} v(\hat{x})^2 w(\hat{x}) \, \mathrm{d}\lambda(\hat{x}) < \infty \right\}.$$

We assume that there exists an accessible orthonormal complete basis  $\{P_k^i:k\in\mathbb{N}\}$  in  $L^2(\hat{X}_i,w_i)$  for every  $i\in[d]$  which is known a priori. For discretisation purposes, we introduce the finite dimensional subspaces

(30) 
$$\mathcal{V}_{i,n_i} := \overline{\operatorname{span}\left\{P_1^i, \dots, P_{n_i}^i\right\}} \subseteq L^2(\hat{X}_i, w_i), \quad i = 1, \dots, d, n_i \in \mathbb{N}.$$

On these we formulate the extended tensor train format in terms of the coefficient tensors

(31) 
$$G^i: [r_{i-1}] \times [n_i] \times [r_i] \to \mathbb{R}, (k_{i-1}, j, k_i) \mapsto G^i[k_{i-1}, j, k_i], i \in [d],$$

such that every univariate function  $g^i \in \mathcal{V}_{i,n_i}$  can be written as

(32) 
$$g^{i}[k_{i-1}, k_{i}](\hat{x}_{i}) = \sum_{i=1}^{n_{i}} G^{i}[k_{i-1}, j, k_{i}] P_{j}^{i}(\hat{x}_{i}) \quad \text{for } \hat{x} \in \hat{X}_{i}.$$

In contrast to the full tensor format for which the function

(33) 
$$g \in \mathcal{V}_{\Lambda} := \bigotimes_{i=1}^{d} \mathcal{V}_{i,n_i} \subseteq \mathcal{V}(\hat{X})$$

can be expressed by a high dimensional algebraic tensor  $G\colon \Lambda:= \mathop{\textstyle \times}_{i=1}^d [n_i] \to \mathbb{R}$  and tensorized spanning functions  $P_\alpha:= \mathop{\textstyle \bigotimes}_{i=1}^d P_{\alpha_i}$  for  $\alpha=(\alpha_1,\ldots,\alpha_d)\in \Lambda$  such that

(34) 
$$g(\hat{x}) = \sum_{\alpha \in A} G[\alpha_1, \dots, \alpha_d] \prod_{i=1}^d P_{\alpha_i}(\hat{x}_i),$$

the format given by (28) and (32) admits a linear structure in the dimension. More precisely, the memory complexity of  $\mathcal{O}(\max\{n_1,\ldots,n_d\}^d)$  in (34) reduces to

(35) 
$$\mathcal{O}(\max\{r_1,\ldots,r_{d-1}\}^2 \cdot d \cdot \max\{n_1,\ldots,n_d\}).$$

This observation raises the question of expressibility for certain classes of functions and the existence of a low-rank vector r where  $\max\{r_1,\ldots,r_{d-1}\}$  is sufficiently small for practical computations. This issue is e.g. addressed in [44, 1, 25] under certain assumptions on the regularity and in [21, 35, 3, 18] explicit (algorithmic) constructions of the format are discussed even in case that g has no analytical representation.

For later reference we define the finite dimensional low-rank manifold of rank r tensor trains by

$$\mathcal{M}_r := \{ g \in \mathcal{V}(\hat{X}) \mid g \text{ as in } (28) \text{ with } g^i \text{ as in } (32) \}.$$

This is an embedded manifold in the finite full tensor space  $V_A$  from (33) admitting the cone property. We also require the concept of the algebraic (full) tensor space

$$\mathbb{T} := \left\{ G \colon \mathbb{N}^d \to \mathbb{R} \right\}$$

and the corresponding low-rank form for given  $oldsymbol{r} \in \mathbb{N}^{d-1}$  defined by

(38) 
$$\mathbb{TT}_{r} := \left\{ G \colon \Lambda \to \mathbb{R} \mid G[\alpha] = \sum_{k=1}^{r} \prod_{i=1}^{d} G[k_{i-1}, \alpha_{i}, k_{i}] \right\}.$$

Without going into detail, we mention the higher order singular value decomposition (HOSVD), which is used to decompose a full algebraic tensor into a low-rank tensor train. The algorithm is based on successive unfoldings of the full tensor into matrices, which are orthogonalized and possibly truncated by a singular value decomposition, see [34] for details. This algorithm enables us to state the following Lemma.

**Lemma 3.1** ([34, Theorem 2.2]). For any  $g \in \mathcal{V}_{\Lambda}$  and  $r \in \mathbb{R}^{d-1}$  there exists an extended low-rank tensor train  $g_r \in \mathcal{M}_r$  with

(39) 
$$||g - g_r||_{\mathcal{V}(\hat{X})}^2 \le \sum_{i=1}^{d-1} \sigma_i^2,$$

where  $\sigma_i$  is the distance of the i-th unfolding matrix of the coefficient tensor of g in the HOSVD to its best rank  $r_i$  approximation in the Frobenius norm.

*Proof.* The proof follows from the best approximation result of the usual matrix SVD with respect to the Frobenius norm and the orthonormality of the chosen basis.  $\Box$ 

3.1. Tensor train regression by Variational Monte Carlo. We review the sampling-based Variational Monte Carlo (VMC) method presented in [18] which is employed to obtain TT representations of the local maps  $\Phi^\ell$  as in (17). The approach generalizes the concept of randomised tensor completion [16] and its analysis relies on the theory of statistical learning, leading to a priori convergence results. It can also be seen as a generalized tensor least squares technique. An alternative cross-interpolation method for probability densities is presented in [11]. The technique relies on active sampling along adaptively chosen index fibers to reduce the amount of evaluations of g.

For the VMC framework, consider the *model class*  $\mathcal{M}_r(\underline{c}, \overline{c}) \subset \mathcal{M}_r$  of truncated rank  $r \in \mathbb{R}^{d-1}$  tensor trains which is given for  $0 \leq \underline{c} < \overline{c} \leq \infty$  by

$$\mathcal{M}_{\boldsymbol{r}}(\underline{c}, \overline{c}) := \left\{ g \in \mathcal{M}_{\boldsymbol{r}} \mid \underline{c} \leq g(\hat{x}) \leq \overline{c} \quad \text{a.e. in} \quad \hat{X} \right\}.$$

The model class  $\mathcal{M}_r(\underline{c}, \overline{c})$  is a finite subset of the truncated nonlinear space

$$\mathcal{V}(\hat{X}, \underline{c}, \overline{c}) := \{ v \in L^2(\hat{X}, w) \mid \underline{c} \le v(\hat{x}) \le \overline{c} \quad \text{a.e. in} \quad \hat{X} \} \subseteq \mathcal{V}(\hat{X}),$$

which we equip with the metric  $d_{\mathcal{V}(\hat{X},c,\bar{c})}(v,w) := \|v-w\|_{\mathcal{V}}$ .

Alternative, we may characterize  $\mathcal{M}_r(\underline{c}, \overline{c})$  and  $\mathcal{V}(\hat{X}, \underline{c}, \overline{c})$ , by shifting the bounds as constraints of the coefficients. Here, we write for  $\ell^2(\mathbb{T}) := \{G \in \mathbb{T} \mid \sum_{\alpha \in \mathbb{N}^d} G[\alpha]^2 < \infty\}$ 

$$(42) \qquad \mathcal{V}(\hat{X},\underline{c},\overline{c}) = \left\{ v(\hat{x}) = \sum_{\alpha \in \mathbb{N}^d} G[\alpha] \cdot P_\alpha(\hat{x}) \mid G \in \ell^2(\mathbb{T}), \ \underline{F}(G) \geq 0, \ \overline{F}(G) \leq 0 \right\},$$

$$\mathcal{M}(\underline{c}, \overline{c}) = \left\{ v(\hat{x}) = \sum_{\alpha \in A} G[\alpha] \cdot P_{\alpha}(\hat{x}) \mid G \in \mathbb{TT}_{r}, \ \underline{F}^{r}(G) \geq 0, \ \overline{F}^{r}(G) \leq 0 \right\},$$

for constrain functions  $\underline{F}, \overline{F} \colon \ell^2(\mathbb{T}) \to \mathbb{R}$  and  $\underline{F}^r, \overline{F}^r \colon \ell^2(TT_r) \to \mathbb{R}$  implicitly bounding the coefficient tensors. Note that due to the orthonormality of  $\{P_\alpha\}_{\alpha \in \mathbb{N}^d}$  in  $\mathcal{V}(\hat{X})$  it holds for every  $v \in \mathcal{V}(\hat{X})$ 

$$\|v\|_{\mathcal{V}} = \|G\|_{\ell^2(\mathbb{T})} \quad \text{with} \quad v = \sum_{\alpha} G[\alpha] P_{\alpha} \in \mathcal{V}.$$

Additionally, we define a loss function  $\iota\colon \mathcal{V}(\hat{X},\underline{c},\overline{c})\times\hat{X}\to\mathbb{R}$  such that  $\iota(\cdot,\hat{x})$  is continuous for almost all  $\hat{x}\in\hat{X}$  and  $\iota(v,\cdot)$  is integrable with respect to the weight function w of  $\mathcal{V}(\hat{X})$  for every  $v\in\mathcal{V}(\hat{X},\underline{c},\overline{c})$ . Then, we consider the cost functional  $\mathscr{J}:\mathcal{V}(\hat{X},\underline{c},\overline{c})\to\mathbb{R}$  given by

(45) 
$$\mathscr{J}(v) := \int_{\hat{X}} \iota(v, \hat{x}) w(\hat{x}) d\lambda(\hat{x}).$$

To further analyse the approximability in the given tensor train format using sampling techniques we present the following properties of two commonly accepted error measures for probability density functions.

**Lemma 3.2.** (KL loss compatibility) Let  $h^* \in \mathcal{V}(\hat{X},0,c^*)$  for  $c^* < \infty$  and  $0 < \underline{c} < \overline{c} < \infty$ . Then

$$\mathcal{V}(\hat{X},\underline{c},\overline{c}) \ni g \mapsto \iota(g,\hat{x}) = \iota(g,\hat{x},h^*) := -\log(g(x))h^*(x)$$

is uniformly bounded and Lipschitz continuous on  $\mathcal{M}_r(\underline{c},\overline{c})$  if  $P_\alpha\in L^\infty(\hat{X})$  for every  $\alpha\in\Lambda$ . Furthermore,  $\mathscr{J}$  is globally Lipschitz continuous on the metric space  $(\mathcal{V}(\hat{X},\underline{c},\overline{c}),d_{\mathcal{V}(\hat{X},\underline{c},\overline{c})})$ .

*Proof.* The loss  $\iota$  is bounded on  $\mathcal{M}_r(\underline{c}, \overline{c})$  since  $0 < \underline{c} < \overline{c} < \infty$ . Let  $g_1, g_2 \in \mathcal{V}_r(\hat{X}, \underline{c}, \overline{c})$  with coefficient tensors  $G_1$  and  $G_2 \in \mathbb{TT}_r$ , then

(47) 
$$|\iota(g_1, \hat{x}) - \iota(g_2, \hat{x})| \leq \sup_{\xi \in [\underline{c}, \overline{c}]} \left\{ \frac{1}{\xi} \right\} \sup_{\hat{x} \in \hat{X}} \{h^*(\hat{x})\} |g_1(\hat{x}) - g_2(\hat{x})|.$$

The global Lipschitz continuity of  $\mathscr{J}$  follows by using (47) and

$$|\mathscr{J}(g_1) - \mathscr{J}(g_2)| \le C_L ||g_1 - g_2||_{L^1(\hat{X}, w)} \le CC_L d_{\mathcal{V}(\hat{X}, c, \bar{c})}(g_1, g_2),$$

with a constant C related to the embedding of  $L^2(\hat{X},w)$  into  $L^1(\hat{X},w)$ . In case that  $g_1,g_2$  additionally be in  $\mathcal{M}_{\boldsymbol{r}}(\underline{c},\overline{c})$  due to Parseval's identity and the finite dimensionality of  $\mathcal{M}_{\boldsymbol{r}}(\underline{c},\overline{c})$  there exists  $c=c\left(\sup_{\alpha\in A}\|P_\alpha\|_{L^\infty(\hat{X})}\right)>0$  such that

$$(49) |g_1(x) - g_2(x)| \le c ||G_1 - G_2||_{\ell^2(\mathbb{T})} = c ||g_1 - g_2||_{\mathcal{V}} = c \, d_{\mathcal{V}(\hat{X},\underline{c},\overline{c})}(g_1,g_2),$$

which yields the Lipschitz continuity on  $\mathcal{M}_r(\underline{c}, \overline{c})$ . Now let  $g_1, g_2 \in \mathcal{V}(\hat{X}, \underline{c}, \overline{c})$ . The global Lipschitz continuity of  $\mathscr{J}$  follows by using (47) and

$$|\mathcal{J}(g_1) - \mathcal{J}(g_2)| \le C_L ||g_1 - g_2||_{L^1(\hat{X}, w)} \le CC_L d_{\mathcal{V}(\hat{X}, c, \overline{c})}(g_1, g_2),$$

with a constant C related to the embedding of  $L^2(\hat{X}, w)$  into  $L^1(\hat{X}, w)$ .

**Lemma 3.3.** ( $L^2$ -loss compatibility) Let  $h^* \in \mathcal{V}(\hat{X}, 0, \overline{c})$  for  $\overline{c} < \infty$  Then

(51) 
$$\mathcal{V}(\hat{X}, 0, \bar{c}) \ni g \mapsto \iota(g, \hat{x}) = \iota(g, \hat{x}, h^*) := |g(\hat{x}) - h^*(\hat{x})|^2$$

is uniformly bounded and Lipschitz continuous on  $\mathcal{M}_r(0, \overline{c})$  provided  $P_\alpha \in L^\infty(\hat{X})$  for every  $\alpha \in \Lambda$ .

*Proof.* Let  $g_1, g_2 \in \mathcal{V}(\hat{X}, 0, \overline{c})$ . Then

$$|\iota(g_1, \hat{x}) - \iota(g_2, \hat{x})| \le |g_1(\hat{x}) - g_2(\hat{x})| \cdot |g_2(\hat{x}) + g_2(\hat{x})| + 2|g_1(\hat{x}) - g_2(\hat{x})|h^*(\hat{x}).$$

Due to  $\overline{c} < \infty$  the Lipschitz property follow as in the proof of Lemma 3.2 if  $g_1, g_2$  in  $\mathcal{M}_r(\underline{c}, \overline{c})$ .

To examine the VMC convergence in our setting, we recall the analysis of [18] in a slightly more general manner. The target objective of the method is to find a minimizer

(53) 
$$v^* \in \operatorname{argmin}_{v \in \mathcal{V}(\hat{X}, c, \bar{c})} \mathscr{J}(v).$$

Due to the infinite dimensional setting we confine the minimization problem in (53) to our model class  $\mathcal{M}=\mathcal{M}_r(\underline{c},\overline{c})$ . This yields the minimization problem

(54) 
$$\operatorname{find} v_{\mathcal{M}}^* \in \operatorname{argmin}_{v \in \mathcal{M}} \mathscr{J}(v).$$

A crucial step is then to consider the empirical functional instead of the integral in  $\mathscr{J}$ , namely

(55) 
$$\mathscr{J}_N(v) := \frac{1}{N} \sum_{k=1}^N \iota(v; \hat{x}^k),$$

with independent samples  $\{\hat{x}^k\}_{k\leq N}$  distributed according to the measure  $w\lambda$  with a (possibly rescaled) weight function w with respect to the Lebesgue measure  $\lambda$ . The corresponding empirical optimization problem then takes the form

(56) 
$$\operatorname{find} v_{\mathcal{M},N}^* \in \operatorname{argmin}_{v \in \mathcal{M}} \mathcal{J}_N(v).$$

The analysis examines different errors with respect to  $h^* \in \mathcal{V}(\hat{X},0,\overline{c})$  defined by

(57) 
$$\mathcal{E} := \left| \mathscr{J}(h^*) - \mathscr{J}\left(v_{\mathcal{M},N}^*\right) \right| \qquad \qquad (VMC \text{ error})$$

(58) 
$$\mathcal{E}_{\mathrm{app}} := |\mathscr{J}(h^*) - \mathscr{J}(v_{\mathcal{M}}^*)| \qquad \qquad (\mathsf{approximation} \ \mathsf{error})$$

(59) 
$$\mathcal{E}_{gen} := \left| \mathcal{J} \left( v_{\mathcal{M}}^* \right) - \mathcal{J} \left( v_{\mathcal{M},N}^* \right) \right| \qquad \text{(generalization error)}.$$

By a simple splitting, the VMC error can be bounded by the approximation and the generalization error, namely

(60) 
$$\mathcal{E} \leq \mathcal{E}_{app} + \mathcal{E}_{gen}$$
.

For our application,  $\mathcal E$  corresponds to a (truncated) Kullback-Leibler divergence and to  $\|.\|^2_{\mathcal V}$ . Due to the global Lipschitz property on  $\mathcal V(\hat X,\underline c,\overline c)$  with  $\underline c>0$  in the setting of (46) or  $\underline c\geq 0$  as in (51), the approximation error can be bounded by the best approximation in  $\mathcal M$ . In particular it exists C>0, such that

(61) 
$$\mathcal{E}_{\text{app}} \le C \inf_{v \in \mathcal{M}} \|h^* - v\|_{\mathcal{V}(\hat{X})}^2.$$

We note that such an estimation by the best approximation in  $\mathcal M$  with respect to the  $\mathcal V(\hat X)$ -norm may not be required in the context of Kullback-Leibner divergence, if one is interested directly in the

best approximation within the latter. Then the  $\underline{c}>0$  assumption can be relaxed in the construction of  $\mathcal{V}(\hat{X},\underline{c},\overline{c})$ , as no global Lipschitz continuity of  $\mathscr{J}$  in Lemma 3.2 is required. Thus, more naturally the subspace of  $\mathcal{V}(\hat{X},0,\overline{c})$  of absolutly continuous functions with respect to  $h^*$  may be considered instead.

It remains to bound the statistical generalization error  $\mathcal{E}_{\mathrm{gen}}$ . For this the notion of covering numbers is required. Let  $(\Omega, \mathcal{F}, \mathbb{P})$  be an abstract probability space.

**Definition 3.4.** (covering number) Let  $\epsilon > 0$ . The covering number  $\nu(\mathcal{M}, \epsilon)$  denotes the minimal number of open balls of radius  $\epsilon$  with respect to the metric  $d_{\nu(\hat{X}, \epsilon; \bar{\epsilon})}$  needed to cover  $\mathcal{M}$ .

**Lemma 3.5.** Let  $\iota$  be defined as in (46) or (51). Then there exist  $C_1, C_2 > 0$  only depending on the constants of uniform bound and Lipschitz continuity on  $\mathcal M$  given in Lemma 3.2 resp. 3.3 such that for  $\epsilon > 0$  and  $N \in \mathbb N$  denoting the number of samples in the empirical cost functional in (55) it holds

$$(62) \qquad \mathbb{P}[\mathcal{E}_{gen} > \epsilon] \le 2\nu(\mathcal{M}, C_2^{-1}\epsilon)\delta(1/4\epsilon, N) \quad \text{with} \quad \delta(\epsilon, N) \le 2\exp(-2\epsilon^2 N/C_1^2).$$

*Proof.* The claim follows immediately from Lemma 3.2 resp. 3.3 and [18, Thm. 4.12, Cor. 4.19].

**Remark 3.6** (choice of  $\underline{c}, \overline{c}$  and  $\hat{X}$ ). Due to the layer based representation in (16) and (23) on each layer  $\hat{X}^{\ell} = \Phi^{-1}(X^{\ell})$  we have the freedom to choose  $\underline{c}$  separately. In particular, assuming that the perturbed prior  $\tilde{f}_0$  decays per layer, we can choose  $\underline{c}$  according to the decay and with this control the constant in (47).

#### 4. ERROR ESTIMATES

This section is devoted to the derivation of a priori error estimates for the previously introduced construction in terms of the Hellinger distance and Kullback Leibler divergence. Concepts arising in Variational Monte Carlo from Section 3.1 are linked to the employed layer structure and the results for two loss functions are collected in the concluding theorems.

Recall that our goal is to approximate the perturbed prior  $\tilde{f}_0$  given some transport  $\tilde{T}$  represented by a function  $\tilde{f}_0^{\rm Trun,TT}$  defined by

(63) 
$$\tilde{f}_0^{\text{Trun,TT}}(x) := \frac{1}{C_L^{\leq} + C_L^{\geq,\text{TT}}} \left\{ \begin{array}{l} \tilde{f}_0^{\ell,\text{TT}}(x), & x \in X^{\ell}, \ell = 1, \dots, L, \\ f_{\Sigma,\mu}(x), & x \in X^{L+1}. \end{array} \right.$$

with

$$C_L^< = \int\limits_{X\backslash K} f_{\Sigma,\mu}(x)\,\mathrm{d}\lambda(x) \quad \text{and} \quad C_L^{>,\mathrm{TT}} := \sum_{\ell=1}^L \int\limits_{X^\ell} \tilde{f}_0^{\ell,\mathrm{TT}}(x)\,\mathrm{d}\lambda(x).$$

Here  $\tilde{f}_0^{\ell,\mathrm{TT}} = \hat{f}_0^{\ell,\mathrm{TT},N_\ell} \circ \varPhi^{\ell^{-1}}$  is the pullback of a function  $\hat{f}_0^{\ell,\mathrm{TT},N_\ell}$  in  $\mathcal{M}^\ell = \mathcal{M}(\underline{c}_\ell,\overline{c}_\ell)$  over  $\hat{X}^\ell$ . Analog to the empirical minimization problem (55) with  $w_\ell = |\det \mathcal{J}_{\varPhi^\ell}|$  we choose  $\hat{f}_0^{\ell,\mathrm{TT},N_\ell}$  as

(64) 
$$\hat{f}_0^{\ell, \text{TT}, N_\ell} \in \operatorname{argmin}_{v \in \mathcal{M}^\ell} \frac{1}{N_\ell} \sum_{k=1}^{N_\ell} \iota(v, \hat{x}^k, \hat{f}_0),$$

with samples  $\{\hat{x}^k\}_{k=1}^{N_\ell}$  drawn from the (possibly rescaled) finite measure  $w_\ell \lambda$ .

The connection to the actual approximation of the target density f given by

(65) 
$$\tilde{f}^{TT} := \tilde{f}_0^{\text{Trun,TT}} \circ \tilde{T}^{-1} \otimes |\mathcal{J}_{\tilde{T}^{-1}}|$$

is reviewed in the following. We refer to Figure 3 for a visual presentation of the involved objects, approximations and transformations.

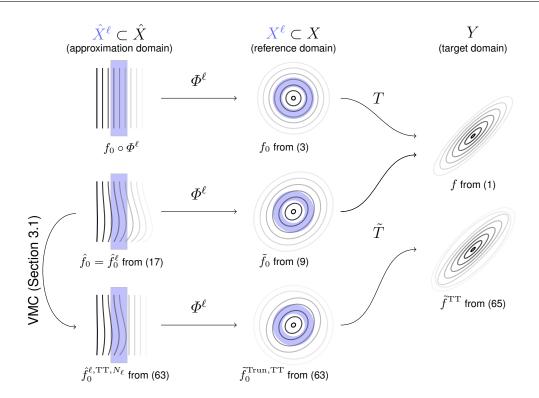


FIGURE 3. Overview of the presented method sketching the different involved transformations and approximations with references to the respective equations.

We first consider the relation of a target density f and its perturbed prior  $\tilde{f}_0$ . Since the transport  $\tilde{T}$  maps X to Y, an error functional  $d(Y;\cdot,\cdot)$  has to satisfy

(66) 
$$d\left(Y; f, \tilde{f}^{TT}\right) = d\left(X; \tilde{f}_0, \tilde{f}_0^{\text{Trun,TT}}\right).$$

This stability property ensures that control of the error of the approximation in terms of the perturbed prior with respect to  $d(X;\cdot,\cdot)$  transfers directly to f. Note that this criterion is canonical as passing to the image space of some measurable function is fundamental in probability theory.

Prominent measures of distinctness for two absolutely continuous Lebesgue probability density functions  $h_1$  and  $h_2$  on some measurable space Z are the Hellinger distance

(67) 
$$d_{\text{Hell}}(Z, h_1, h_2) = \int_{Z} \left( \sqrt{h_1}(z) - \sqrt{h_2}(z) \right)^2 d\lambda(z),$$

and the Kullback-Leibler divergence

(68) 
$$d_{\mathrm{KL}}(Z, h_1, h_2) = \int\limits_{Z} \log \left(\frac{h_1(z)}{h_2(z)}\right) h_1(z) \,\mathrm{d}\lambda(z).$$

For the Hellinger distance, the absolute continuity assumption can be dropped from an analytical point of view. Observe that both  $d_{\rm Hell}$  and  $d_{\rm KL}$  both satisfy (66).

#### Lemma 4.1. It holds

(69) 
$$d_{\text{Hell}}(Y; f, \tilde{f}^{TT}) = d_{\text{Hell}}(X; \tilde{f}_0, \tilde{f}_0^{\text{Trun,TT}}),$$

(70) 
$$d_{\mathrm{KL}}(Y; f, \tilde{f}^{\mathrm{TT}}) = d_{\mathrm{KL}}(X; \tilde{f}_0, \tilde{f}_0^{\mathrm{Trun, TT}}).$$

Proof. We only show (70) since (69) follows by similar arguments. By definition

(71) 
$$d_{\mathrm{KL}}(Y; f, \tilde{f}^{\mathrm{TT}}) = \int_{Y} \log \left( \frac{f(y)}{\tilde{f}^{\mathrm{TT}}(y)} \right) f(y) \, \mathrm{d}\lambda(y)$$

and the introduction of the transport map  $\tilde{T}$  yields the claim

(72) 
$$\int_{X} \log \left( \frac{f \circ \tilde{T}(x)}{\tilde{f}^{\mathrm{TT}} \circ \tilde{T}(x)} \cdot \frac{|\det \mathcal{J}_{\tilde{T}}(x)|}{|\det \mathcal{J}_{\tilde{T}}(x)|} \right) \tilde{f}_{0}(x) \, \mathrm{d}\lambda(x) = d_{\mathrm{KL}}(X; \tilde{f}_{0}, \tilde{f}_{0}^{\mathrm{Trun, TT}}).$$

Collecting the previous results and notations the following assumption turns out to be required.

**Assumption 4.2.** For a target density  $f\colon Y\to\mathbb{R}_+$  and a transport map  $\tilde{T}\colon X\to Y$ , there exists a simply connected compact set K such that  $\tilde{f}_0=(f\circ T)\otimes|\det\mathcal{J}_T|\in L^2(K)$  has outer polynomial exponential decay with polynomial  $\pi^+$  on  $X\setminus K$ . Consider the symmetric positive definite matrix  $\Sigma\in\mathbb{R}^{d,d}$  and  $\mu\in\mathbb{R}^d$  as the covariance and mean for the outer approximation  $f_{\Sigma,\mu}$ . Furthermore, let  $K=\bigcup_{\ell=1}^L\overline{X^\ell}$ , with  $X^\ell$  being the image of a rank 1 stable diffeomorphism  $\Phi^\ell\colon\hat{X}^\ell\to X^\ell$  for every  $\ell=1,\ldots,L$ .

We can now formulate the main theorems of this section regarding the convergence of the approximation.

**Theorem 4.3.** (A priori estimates in Hellinger distance) Let Assumption 4.2 hold and let a sequence of sample sizes  $(N^\ell)_{\ell=1}^L \subset \mathbb{N}$  be given. For every  $\ell=1,\ldots,L$  consider bounds  $0<\underline{c}^\ell<\overline{c}^\ell<\infty$  and let  $\tilde{f}^{\mathrm{TT}}$  be defined as in (65). Then there exist constants  $C,C_{\Sigma},C^\ell,C^\ell_\iota>0,\ell=1,\ldots,L$  such that

(73) 
$$d_{\text{Hell}}(Y, f, \tilde{f}^{\text{TT}}) \leq C \left( \sum_{\ell=1}^{L} \left( \mathcal{E}_{\text{best}}^{\ell} + \mathcal{E}_{\text{sing}}^{\ell} + \mathcal{E}_{\text{gen}}^{\ell} \right) + \mathcal{E}_{\text{trun}} \right).$$

Here,  $\mathcal{E}_{\mathrm{best}}^{\ell}$  denotes the error of the best approximation  $v_{A}^{\ell}$  to  $\hat{f}_{0}^{\ell}$  in the full truncated polynomial space  $\mathcal{V}_{A}^{\ell}(\underline{c}^{\ell},\overline{c}^{\ell})=\mathcal{V}_{A}^{\ell}\cap\mathcal{V}(\hat{X}^{\ell},\underline{c}^{\ell},\overline{c}^{\ell})$  given by

$$\mathcal{E}_{\text{best}}^{\ell} := \|\hat{f}_{0}^{\ell} - v_{\Lambda}^{\ell}\|_{\mathcal{V}(\hat{X}^{\ell})} = \inf_{v^{\ell} \in \mathcal{V}_{\Lambda}^{\ell}(\underline{c}^{\ell}, \overline{c}^{\ell})} \|\hat{f}_{0}^{\ell} - v^{\ell}\|_{\mathcal{V}(\hat{X}^{\ell})},$$

 $\mathcal{E}_{ ext{sing}}^\ell$  is the low-rank approximation error of the algebraic tensor associated to  $v_A^\ell$  and the truncation error  $\mathcal{E}_{ ext{trun}}$  is given by

$$\mathcal{E}_{\mathrm{trun}}^2 := \|\exp\left(-\pi^+\right)\|_{L^1(X\setminus K)} + \Gamma\left(d/2, C_{\Sigma}R^2\right).$$

Furthermore, for any  $(\epsilon^\ell)_{\ell=1}^L \subset \mathbb{R}_+$  the generalization errors  $\mathcal{E}^\ell_{\mathrm{gen}}$  can be bounded in probability

$$\mathbb{P}(\mathcal{E}_{\text{gen}}^{\ell} > \epsilon^{\ell}) \leq 2\nu(\mathcal{M}^{\ell}, C^{\ell}\epsilon^{\ell})\delta^{\ell}(1/4\epsilon^{\ell}, N^{\ell})$$

with  $\nu$  denoting the covering number from Definition 3.4 and  $\delta^\ell(\epsilon,N) \leq 2\exp(-2\epsilon^2N/C_\iota^\ell)$ .

*Proof.* We note that the Hellinger distance can be bounded by the  $L^2$  norm. Note that  $|\sqrt{a}-\sqrt{b}| \le \sqrt{|a-b|}$  for  $a,b\ge 0$  and with Lemma 4.1 it holds

$$d_{\text{Hell}}(Y; f, \tilde{f}^{\text{TT}}) = d_{\text{Hell}}(X; \tilde{f}_0, \tilde{f}_0^{\text{Trun,TT}})$$

$$\leq \|\tilde{f}_0 - \tilde{f}_0^{\text{Trun,TT}}\|_{L^1(K)} + \|\tilde{f}_0 - \tilde{f}_0^{\text{Trun,TT}}\|_{L^1(X \setminus K)}.$$

Since  $K=\cup_{\ell=1}^L X^\ell$  and  $X^\ell$  bounded, there exist constants  $C(X^\ell)>0$ ,  $\ell=1,\dots,L$  such that

$$\|\tilde{f}_{0} - \tilde{f}_{0}^{\text{Trun,TT}}\|_{L^{1}(K)} = \sum_{\ell=1}^{L} \|\tilde{f}_{0} - \tilde{f}_{0}^{\text{Trun,TT}}\|_{L^{1}(X^{\ell})}$$

$$\leq \sum_{\ell=1}^{L} C(X_{\ell}) \|\tilde{f}_{0} - \tilde{f}_{0}^{\text{Trun,TT}}\|_{L^{2}(X^{\ell})}.$$

Moreover, by construction

(74) 
$$\|\tilde{f}_0 - \tilde{f}_0^{\text{Trun,TT}}\|_{L^2(X^{\ell})} = \|\hat{f}_0^{\ell} - \hat{f}_0^{\ell,\text{TT},N_{\ell}}\|_{\mathcal{V}(\hat{X}^{\ell})}.$$

Then the claim follows by application of Lemmas 2.7, 3.1 and 3.5 together with equation (60).  $\Box$ 

Theorem 4.4. (A priori estimates in Kullback Leibler divergence) Let Assumption 4.2 hold and let a sequence of sample sizes  $(N^\ell)_{\ell=1}^L \subset \mathbb{N}$  be given. For every  $\ell=1,\ldots,L$  consider bounds  $0<\underline{c}^\ell<\overline{c}^\ell<\infty$  and let  $\widetilde{f}^{\mathrm{TT}}$  be defined as in (65). Then there exists constants  $C,C_\Sigma,C^\ell,C^\ell_\iota>0$ ,  $\ell=1,\ldots,L$  such that

(75) 
$$d_{\mathrm{KL}}(Y, f, \tilde{f}^{\mathrm{TT}}) \leq C \left( \sum_{\ell=1}^{L} \left( \mathcal{E}_{\mathrm{best}}^{\ell} + \mathcal{E}_{\mathrm{sing}}^{\ell} + \mathcal{E}_{\mathrm{gen}}^{\ell} \right) + \mathcal{E}_{\mathrm{trun}} \right).$$

Here,  $\mathcal{E}^\ell_{\mathrm{best}}$  denotes the error of the best approximation  $v^\ell_\Lambda$  to  $\hat{f}^\ell_0$  in the full truncated polynomial space  $\mathcal{V}^\ell_\Lambda(c^\ell,\overline{c}^\ell)=\mathcal{V}^\ell_\Lambda\cap\mathcal{V}(\hat{X}^\ell,c^\ell,\overline{c}^\ell)$  given by

$$\mathcal{E}_{\text{best}}^{\ell} := \|\hat{f}_0^{\ell} - v_{\Lambda}^{\ell}\|_{\mathcal{V}(\hat{X}^{\ell})} = \inf_{v^{\ell} \in \mathcal{V}_{\Lambda}^{\ell}(\underline{c}^{\ell}, \overline{c}^{\ell})} \|\hat{f}_0^{\ell} - v^{\ell}\|_{\mathcal{V}(\hat{X}^{\ell})},$$

 $\mathcal{E}_{ ext{sing}}^\ell$  is the low-rank approximation error of the algebraic tensor associated to  $v_A^\ell$  and the truncation error  $\mathcal{E}_{ ext{trun}}$  is given by

$$\mathcal{E}_{\text{trun}} := \int_{X \setminus K} \left( \frac{1}{2} \|x\|_{\Sigma^{-1}}^2 + \tilde{\pi}^+(x) \right) e^{-\tilde{\pi}^+(x)} \, \mathrm{d}\lambda(x).$$

Furthermore, for any  $(\epsilon^\ell)_{\ell=1}^L \subset \mathbb{R}_+$  the generalization errors  $\mathcal{E}_{\mathrm{gen}}^\ell$  can be bounded in probability

$$\mathbb{P}(\mathcal{E}_{\text{gen}}^{\ell} > \epsilon^{\ell}) \le 2\nu(\mathcal{M}^{\ell}, C^{\ell} \epsilon^{\ell}) \delta^{\ell}(1/4\epsilon^{\ell}, N^{\ell})$$

with  $\nu$  denoting the covering number from Definition 3.4 and  $\delta^\ell(\epsilon,N) \leq 2\exp(-2\epsilon^2N/C_\iota^\ell)$ .

Proof. Using Lemma 4.1 and the construction (63) it holds

(76) 
$$d_{KL}(Y; f, \tilde{f}^{TT}) = \sum_{\ell=1}^{L} \int_{X^{\ell}} \log \frac{\tilde{f}_{0}}{\tilde{f}_{0}^{\ell, TT}} \tilde{f}_{0} d\lambda(x) + \int_{X \setminus K} \log \frac{\tilde{f}_{0}}{f_{\Sigma, \mu}} \tilde{f}_{0} d\lambda(x).$$

By Lemma 2.7 we can bound the integral over  $X \setminus K$  by the truncation error  $\mathcal{E}_{trun}$ . Employing the loss function and cost functional of Lemma 3.2 yields

(77) 
$$\int_{X^{\ell}} \log \frac{\tilde{f}_0}{\tilde{f}_0^{\ell, \text{TT}}} \tilde{f}_0 d\lambda(x) \leq \mathcal{E}_{\text{app}}^{\ell} + \mathcal{E}_{\text{gen}}^{\ell}.$$

The claim eventually follows by application of Lemmas 3.1 and 3.5 together with (60).

#### 5. ALGORITHM

Since a variety of techniques are employed in the density discretisation, this section provides an exemplary algorithmic workflow to illustrate the required steps in practical applications (see also Figure 1 for a sketch of the components of the method). The general method to obtain a representation of the density (1) by its auxiliary reference (9) is summarized in Algorithm 1. Based on this, the computation of possible quantities of interest such as moments (10) or marginals are considered in Sections 6.3 and 6.4, respectively. In the following we briefly describe the involved algorithmic procedures.

Computing the transformation. Obtaining a suitable transport map is a current research topic and discussed in various publications. We hence do not cover it in detail here but only provide some references. In Section 2.1, two naive options are introduced. In the numerical applications, we employ affine transport and also illustrate the capabilities of quadratic transport in a two-dimensional example. For the affine linear transport we utilize a semi-Newton optimizer to obtain the maximum value of f and an approximation of the Hessian at the optimal value, see Section 2.1.1. For the construction of a quadratic transport we rely on the library  ${\tt TransportMaps}$  [4]. We summarize the task to provide the (possibly inexact) transport map in the function

(78) 
$$\tilde{T} \leftarrow \text{ComputeTransport}[f].$$

In the following paragraphs we assume  $\Phi^\ell$  to be the multivariate polar transformation as in Example 2.2, defined on the corresponding hyperspherical shells  $\hat{X}^\ell$ . We refer to  $\hat{X}_1^\ell$  as the *radial dimension* and  $\hat{X}_i^\ell$  as the *angular dimensions* for  $1 < i \le d$ . Note that the loop over  $\ell = 1, \ldots, L$  can be parallelized in practice.

**Generating an orthonormal basis**. To obtain suitable finite dimensional subspaces one has to introduce spanning sets that allow for an efficient computation of e.g. moments (4) and the optimization of the functional (45). Given a fixed dimension vector  $\boldsymbol{n}^\ell \in \mathbb{N}^d$  for the current  $\hat{X}^\ell$ ,  $\ell=1,\ldots,L$ , and by the chosen parametrization via  $\Phi^\ell$  introducing the weight  $w^\ell$ , the function

(79) 
$$\mathcal{P}^{\ell} = \{\mathcal{P}_{i}^{\ell}\}_{i=1}^{d} \leftarrow \text{GenerateONB}[\hat{X}^{\ell}, \boldsymbol{n}^{\ell}, w^{\ell}, \tau_{\text{GS}}]$$

can be split into three distinct algorithmic parts as follows.

■ 1st coordinate  $\hat{x}_1$ : The computation of an orthonormal polynomial basis  $\{P_{1,\alpha}^\ell\}_\alpha$  with respect to the weight  $w_1^\ell(\hat{x}_1) = \hat{x}_1^{d-1}$  in the radial dimension by a stabilized Gram-Schmidt method. This is numerically unstable since the involved summations cause cancellation. As a remedy, we define arbitrary precision polynomials with a significant digit length  $\tau_{\rm mant}$  to represent polynomial coefficients. By this, point evaluations of the orthonormal and computations of integrals of the form

(80) 
$$\int_{\hat{X}_{1}^{\ell}} \hat{x}_{1}^{m} P_{1,\alpha}^{\ell}(\hat{x}_{1}) \hat{x}_{1}^{d-1} d\lambda(\hat{x}_{1}), \quad m \in \mathbb{N},$$

e.g. required for computing moments with polynomial transport, can be realized exactly. The length  $au_{\rm mant}$  is set to 100 in the numerical examples and the additional run-time is negligible as the respective calculations can be precomputed.

**2**nd coordinate  $\hat{x}_2$ : Since  $\hat{X}_2^{\ell} = [0, 2\pi]$  and to preserve periodicity, we employ trigonometric polynomials given by

(81) 
$$P_{2,j}^{\ell}(\hat{x}_2) = \begin{cases} \frac{1}{\sqrt{2\pi}}, & j = 1\\ \frac{\sin(\frac{j}{2}\hat{x}_2)}{\sqrt{\pi}}, & j \text{ even} \\ \frac{\cos(\frac{j-1}{2}\hat{x}_2)}{\sqrt{\pi}}, & j > 1 \text{ odd.} \end{cases}$$

Note that here the weight function is constant, i.e.  $w_2^{\ell}(\hat{x}_2) \equiv 1$ , and the defined trigonometric polynomials are orthonormal in  $L^2(\hat{X}_2^{\ell})$ .

coordinate  $\hat{x}_3,\ldots,\hat{x}_d$ : On the remaining angular dimensions  $i=3,\ldots,d$ , we employ the usual Gram-Schmidt orthogonalization algorithm on  $[0,\pi]$  with weight function  $w_i^\ell(\hat{x}_i)=\sin^i(\hat{x}_i)$ , based on polynomials.

Fortunately, the basis for dimensions  $1 < i \le d$  coincides on every layer  $\ell = 1, \ldots, L$ . It hence can be computed just once and passed to the individual process handling the current layer. Only the basis in the radial dimension needs to be adjusted to  $\hat{X}^{\ell}$ . The parameter  $\tau_{\rm GS}$  collects all tolerance parameters for the applied numerical quadrature and the significant digit length  $\tau_{\rm mant}$ .

**Generation of Samples**. To generate samples on  $\hat{X}^\ell$  with respect to the weight function  $w^\ell$ , we employ inverse transform sampling. For this the weight function is rescaled to have unit norm in  $L^1(\hat{X}^\ell)$ . Then, the involved inverse cumulative distribution functions can be computed analytically. We denote the generation process of  $N \in \mathbb{N}$  samples as the function

(82) 
$$\mathcal{S}^{\ell} := \left\{ \left( \hat{x}^s, \hat{f}_0^{\ell}(\hat{x}^s) \right) \right\}_{s=1}^N \leftarrow \text{GenerateSamples}[\hat{f}_0^{\ell}, \hat{X}^{\ell}, w^{\ell}, N].$$

**Reconstruction of a Tensor Train surrogate**. The VMC reconstruction approach of Section 3 is summarized in the function

(83) 
$$\left\{\hat{F}_{0,i}^{\ell,\mathrm{TT}}\right\}_{i=1}^{d} \leftarrow \mathrm{ReconstructTT}[\mathcal{S}^{\ell}, \mathcal{P}^{\ell}, \boldsymbol{r}^{\ell}, \tau_{\mathrm{Recon}}].$$

The tensor components  $\hat{F}_{0,i}^{\ell,\mathrm{TT}}$  are associated with the corresponding basis  $\mathcal{P}_i^\ell$  to form a rank  $r^\ell$  extended tensor train as defined in (28) and (32). The additional parameter  $\tau_{\mathrm{Recon}}$  is a "wild-card" for all parameters that determine the VMC algorithm.

The method basically involves the optimization of a loss functional over the set of tensor trains with rank (at most)  $r^\ell$ . In the presented numerical computations we consider a mean-square loss and the respective empirical approximation based on a current sample set  $\mathcal{S}^\ell$ . The tensor optimization, based on a rank adaptive, alternating direction fitting (ADF) algorithm, is implemented in the xerus library [28] and wrapped in the ALEA framework [13]. Additionally, the machine learning framework PyTorch [38] can be utilized in ALEA to minimize the empirical cost functional from (55) by a wide class of state-of-the-art stochastic optimizers. The latter enables stochastic gradient methods to optimize the tensor coefficients as known from machine learning applications. Having this setting in mind, the actual meaning of the parameter  $\tau_{\rm Recon}$  depends on the chosen optimizer. In this article we focus on the ADF implementation and initialize e.g. the starting rank, the number of iteration of the ADF and a target residual norm.

#### 6. APPLICATIONS

In the preceding sections the creation of surrogate models of rather arbitrary probability density functions were developed. Using this, in the following we focus on actual applications where such a representation is beneficial. We start with the framework of Bayesian inverse problems with target density (1) corresponding to the Lebesgue posterior density. Subsequently, we cover the computation of moments and marginals as important applications.

6.1. **Bayesian inversion.** This section is devoted to a brief review of the Bayesian paradigm. We recall the general formalism and highlight the notation with setup of Section 2 in mind. We closely follow the presentation in [15] and refer to [46, 7, 29] for a comprehensive overview.

#### Algorithm 1 Tensor train surrogate creation of perturbed prior

Input: Lebesgue target density 
$$f\colon \mathbb{R}^d \to \mathbb{R}_+$$
 (1)

tensor spaces 
$$\left\{\hat{X}^{\ell}\right\}_{\ell=1}^{L}$$
, with  $\hat{X}^{\ell}=\mathop{\textstyle imes_{i=1}}^{d}\hat{X}_{i}^{\ell}$  (17)

coordinate transformations 
$$\Phi^{\ell} \colon \hat{X}^{\ell} \to X^{\ell} \subset \mathbb{R}^{d}$$
 (21)

with rank-1 Jacobians 
$$w^\ell := |\det \left[ \mathcal{J}_{\varPhi^\ell} \right] | \colon \hat{X}^\ell o \mathbb{R}$$

with rank-1 Jacobians 
$$w^\ell := |\det \left[ \mathcal{J}_{\varPhi^\ell} \right]| \colon \hat{X}^\ell \to \mathbb{R}$$
 basis dimensions  $(\boldsymbol{n}^1, \dots, \boldsymbol{n}^L)$ ,  $\boldsymbol{n}^\ell \in \mathbb{N}^d$  for  $\ell = 1, \dots, L$  (32)

sample size  $N_{\ell} \in \mathbb{N}$ ,  $\ell = 1, \dots, L$  for level-wise reconstruction

tensor train ranks 
$$(\boldsymbol{r}^1,\ldots,\boldsymbol{r}^L)$$
,  $\boldsymbol{r}^\ell\in\mathbb{N}^{d-1}$ , for  $\ell=1,\ldots,L$  (28)

Gram-Schmidt tolerance parameter  $au_{\mathrm{GS}}$ 

tensor reconstruction parameter  $au_{
m Recon}$ 

Output: Level-wise low-rank approximation of perturbed prior

Diffeomorphism  $\tilde{T} \leftarrow \text{ComputeTransport}[f]$ 

for  $\ell = 1, \dots, L$ , (in parallel) do

- ullet Set transformed *perturbed prior*  $\hat{f}_0^\ell(\hat{x}) := \left(f \circ \tilde{T} \otimes |\det \mathcal{J}_{\tilde{T}}|\right) \circ \Phi^\ell(\hat{x}), \ \hat{x} \in \hat{X}^\ell$
- ullet Build one-dimensional ONB  $\mathcal{P}_i^\ell$  of  $\mathcal{V}_{i,n_i^\ell}\subseteq L^2(\hat{X}_i^\ell,w_i^\ell)$  for  $i=1,\dots,d$

$$\mathcal{P}^{\ell} = \{\mathcal{P}_i^{\ell}\}_{i=1}^d \leftarrow \text{GenerateONB}[\hat{X}^{\ell}, \boldsymbol{n}^{\ell}, w^{\ell}, \tau_{\text{GS}}]$$

ullet Generate samples with respect to the weight  $w^\ell$ 

$$\mathcal{S}^{\ell} := \left\{ \left( \hat{x}^s, \hat{f}_0^{\ell}(\hat{x}^s) \right) \right\}_{s=1}^{N} \leftarrow \text{GenerateSamples}[\hat{f}_0^{\ell}, \hat{X}^{\ell}, w^{\ell}, N]$$

ullet Reconstruct TT surrogate  $\hat{f}_0^{\ell,\mathrm{TT}}\colon \hat{X}^\ell o \mathbb{R}$ 

$$\left\{ \tilde{F}_{0,i}^{\ell,\mathrm{TT}} \right\}_{i=1}^{d} \leftarrow \mathrm{ReconstructTT}[\mathcal{S}^{\ell}, \mathcal{P}^{\ell}, \boldsymbol{r}^{\ell}, \tau_{\mathrm{Recon}}]$$

Equip tensor components with basis

$$\begin{array}{rcl} \hat{f}_{0}^{\ell,\mathrm{TT}}(\hat{x}) & := & \sum_{\pmb{k}}^{\pmb{r}^{\ell}} \prod_{i=1}^{d} \hat{f}_{0,i}^{\ell,\mathrm{TT}}[k_{i-1},k_{i}](\hat{x}_{i}) \\ \text{where } \hat{f}_{0,i}^{\ell,\mathrm{TT}}[k_{i-1},k_{i}](\hat{x}_{i}) & := & \sum_{j=1}^{n_{j}^{\ell}} \hat{F}_{0,i}^{\ell,\mathrm{TT}}[k_{i-1},\mu_{i},k_{i}] P_{i,j}^{\ell}(\hat{x}_{i}) \\ \text{end for} \\ \text{return } \left\{\tilde{f}_{\ell}\right\}_{l=1}^{L} \end{array}$$

Let Y, V and  $\mathcal{Y}$  denote separable Hilbert spaces equipped with norms  $\|\cdot\|_H$  and inner products  $\langle \cdot, \cdot \rangle_H$ ,  $H \in \{Y, V, \mathcal{Y}\}$ . The uncertain quantity  $y \in Y$  is tied to the model output  $q \in V$  by the forward map

(84) 
$$G: Y \to V, \quad \theta \mapsto q(y) := G(y).$$

The usual forward problem reads

(85) Given 
$$y \in Y$$
, find  $q \in V$ .

In contrast to this, the inverse problem is defined by

(86) Given observations of 
$$q$$
, find  $y \in Y$ .

The term *observations* is determined by a bounded linear operator  $\mathcal{O} \colon V \to \mathcal{Y}$  that describes the measurement process of the quantity q. In practical applications this could be direct observations at sensor points or averaged values from monitoring devices, i.e. one models  $\mathcal{Y} = \mathbb{R}^J$  for some  $J \in \mathbb{N}$ .

Classically, the (deterministic) quantification problem (86) is not well-posed. To overcome this, a problem regularization of some kind is required. The chosen probabilistic approaches introduces a random measurable additive noise  $\eta\colon (\Omega,\mathcal{U},\mathbb{P})\to (\mathcal{Y},\mathcal{B}(\mathcal{Y}))$ , with law  $\mathcal{N}(0,C_0)$  for some symmetric positive definite covariance operator  $C_0$  on  $\mathcal{Y}$ , and compose it with the observed system response

(87) 
$$\delta = (\mathcal{O} \circ G)(y) + \eta =: \mathcal{G}(y) + \eta \quad \text{where } \mathcal{G} : Y \to \mathcal{Y}.$$

As a consequence, the quantities y, q and  $\delta$  become random variables over a probability space  $(\Omega, \mathcal{F}, \mathbb{P})$  with values in Y, V and  $\mathcal{Y}$ , respectively. In [46] mild conditions on the forward operator are derived to show a continuous version of Bayes formula which yields the existence and uniqueness of the Radon-Nikodym derivative of the (posterior) measure  $\pi_{\delta}$  of the conditional random variable  $y|\delta$  with respect to a prior measure  $\pi_0$  of y. More precisely, by the Gaussian noise assumption on  $\eta$  and assuming independence with respect to y, both measures  $\pi_0$  and  $\pi_{\delta}$  on Y are related by the *Bayesian potential* 

(88) 
$$\Psi(y,\delta) := \frac{1}{2} \langle C_0^{-1}(\delta - \mathcal{G}(y)), \delta - \mathcal{G}(y) \rangle_{\mathcal{Y}}$$

in the sense that

(89) 
$$\frac{\mathrm{d}\pi_{\delta}}{\mathrm{d}\pi_{0}}(y) = Z^{-1} \exp\left(-\Psi(y,\delta)\right), \quad Z := \mathbb{E}_{\pi_{0}}\left[\exp\left(-\Psi(y,\delta)\right)\right].$$

Note that we interchangeably write y as an element of Y and the corresponding random variable with values in Y.

6.2. **Bayesian inversion for parametric PDEs.** Random partial differential equations (PDEs), i.e. PDEs with random (in our case parametric) data, play an important role in the popular research area of Uncertainty Quantification (UQ). As a prominent benchmark example, we consider the structural ground water flow model, also called the Darcy problem, as e.g. examined in [12, 17, 14]. In this linear second order PDE model, the forward operator G in (84) on some domain  $D \subset \mathbb{R}^d$ , d=1,2,3 is determined by a forcing term  $f \in L^2(D)$  and the random/parametric quantity  $a(y) \in L^\infty(D)$ , which for almost every  $y \in Y$  models a conductivity or permeability coefficient. The physical system is described by

(90) 
$$-\operatorname{div}\left(a(y)\nabla g(y)\right) = f \text{ in } D, \quad g(y)|_{\partial D} = 0$$

and the solution  $q(y) \in V = H^1_0(D)$  corresponds to the system response G(y) = q(y). Pointwise solvability of (90) for almost every  $y \in Y$  is guaranteed by a Lax-Milgram argument and existence and uniqueness of a uniform solution  $q \colon D \times Y \to \mathbb{R}$  for several choices of a as well as a comprehensive overview on parametric PDEs can be found in [45].

For the applications in this article we employ a truncated log-normal coefficient field

(91) 
$$a(y) = \exp\left(\sum_{k=1}^{d} a_k y_k\right)$$

for some fixed  $(a_k)_{k=1}^d$  with  $a_k \in L^2(D)$  and the image of some random variable with law  $\mathcal{N}(0,I)$  denoted by  $y=(y_k)_{k=1}^d \in Y$ . Together with synthetic point observations of q at nodes  $\delta=(\delta_1,\ldots,\delta_J)$ 

in D corresponding to some unknown  $q(y^*)$ ,  $y^* \in Y$  and related by (87), we consider the Bayesian posterior density (89). Following the initial notation in (1), we set

$$(92) f(y) = Z^{-1} \mathrm{d}\pi_{\delta}(y) \mathrm{d}\pi_{0}(y)$$

as the Lebesgue density of the target measure  $\pi$  on Y.

6.3. **Moment computation.** In this section we discuss the computation of moments for the presented layer-based format with low-rank tensor train approximations. In particular we are interested in an efficient generation of the map

(93) 
$$\alpha \mapsto \int_{V} y^{\alpha} f(y) d\lambda(y)$$
 for multiindex  $\alpha = (\alpha_k)_k \in \mathbb{N}_0^d$ .

Given some transport  $\tilde{T}\colon X\to Y$  with an associated perturbed prior  $\tilde{f}_0=(f\circ \tilde{T})\otimes |{\rm det}\,\mathcal{J}_{\tilde{T}}|$  by integral transformation it holds

(94) 
$$\int_{Y} y^{\alpha} f(y) d\lambda(y) = \int_{X} \tilde{T}(x)^{\alpha} \tilde{f}_{0}(x) d\lambda(x).$$

Motivated by the preceding sections fix  $1 \leq \ell \leq L$  and assume tensor spaces  $\hat{X}^\ell, X^\ell$  such that a layer based splitting can be employed to obtain integrals over  $X^\ell$  of the form

(95) 
$$\int_{Y} y^{\alpha} f(y) d\lambda(y) = \sum_{\ell=1}^{L} \int_{Y^{\ell}} \tilde{T}(x)^{\alpha} \tilde{f}_{0}(x) dx.$$

Note that we neglect the remaining unbounded layer  $X^{L+1}$  since for moderate  $|\alpha|$  and  $\operatorname{vol}(\bigcup_{\ell=1}^L X^\ell)$  sufficiently large, the contribution to the moment does not have a significant influence on the overall approximation. Additionally, a rank 1 stable diffeomorphism  $\Phi^\ell\colon \hat{X}^\ell\mapsto X^\ell$  is assumed for which there exist univariate functions  $\Phi^\ell_{,j}\colon \hat{X}^\ell_j\to X^\ell$  with  $\Phi^\ell_{,j}=(\Phi^\ell_{i,j})_{i=1}^d$  and  $h_j\colon \hat{X}^\ell_j\to \mathbb{R}$  for every  $j=1,\ldots,d$  such that

6.3.1. Moments under affine transport. Let  $H=[h_{ki}]_{k,i=1}^d=[h_1,h_2,\ldots,h_d]\in\mathbb{R}^{d,d}$  be a symmetric positive definite matrix and  $M=(M_i)_{i=1}^d\in\mathbb{R}^d$  such that the considered transport map takes the form

(97) 
$$\tilde{T}(\cdot) = H \cdot + M.$$

With the multinomial coefficient for  $j \in \mathbb{N}$ ,  $\beta \in \mathbb{N}_0^d$  with  $j = |\beta|$  given by

$$\begin{pmatrix} j \\ \boldsymbol{\beta} \end{pmatrix} := \frac{j!}{\beta_1! \cdot \ldots \cdot \beta_d!},$$

the computation of moments corresponds to the multinomial theorem as seen in the next lemma.

**Lemma 6.1.** Let  $k \in \mathbb{N}$  with  $1 \le k \le d$  and  $\alpha_k \in \mathbb{N}_0$ . It holds

(98) 
$$[H\Phi^{l}(\hat{x}) + M]_{k}^{\alpha_{k}} = \sum_{j_{k}=0}^{\alpha_{k}} \sum_{|\boldsymbol{\beta}_{k}|=j_{k}} C_{k}^{H}[j_{k}, \alpha_{k}, \boldsymbol{\beta}_{k}] \prod_{j=1}^{d} \boldsymbol{\Phi}_{j}^{\boldsymbol{\beta}_{k}}(\hat{x}_{j}),$$

where the high-dimensional coefficient  $C_k^H$  is given by

(99) 
$$C_k^H[j_k, \alpha_k, \boldsymbol{\beta}_k] := \begin{pmatrix} \alpha_k \\ j_k \end{pmatrix} c_k^{\alpha_k - j_k} \begin{pmatrix} j_k \\ \boldsymbol{\beta}_k \end{pmatrix} h_k^{\boldsymbol{\beta}_k},$$

with 
$$c_k := \sum\limits_{i=1}^d h_{ki} M_i$$
 and

(100) 
$$m{\Phi}_j^{m{eta}_k} := [\Phi_{1,j}^\ell(\hat{x}_j), \dots, \Phi_{d,j}^\ell(\hat{x}_j)]^{m{eta}_k}.$$

Proof. Note that

$$[H\Phi^{\ell}(\hat{x}) + M)]_{k}^{\alpha_{k}} = \sum_{j_{k}=0}^{\alpha_{k}} {\alpha_{k} \choose j_{k}} c_{k}^{\alpha_{k}-j_{k}} \left( \sum_{i=1}^{d} h_{ki} \prod_{j=1}^{d} \Phi^{\ell}_{ij}(\hat{x}_{j}) \right)^{j_{k}}.$$

The statement follows by the multinomial theorem since

$$\left(\sum_{i=1}^d h_{ki} \prod_{j=1}^d \Phi_{ij}^{\ell}(\hat{x}_j)\right)^{j_k} = \sum_{|\boldsymbol{\beta}_k| = j_k} \begin{pmatrix} j_k \\ \boldsymbol{\beta}_k \end{pmatrix} \left(\prod_{i=1}^d h_{ki}^{(\boldsymbol{\beta}_k)_i}\right) \left(\prod_{j=1}^d \prod_{i=1}^d \Phi_{ij}^{\ell}(\hat{x}_j)^{(\boldsymbol{\beta}_k)_i}\right). \quad \Box$$

Generalizing Lemma 6.1 to multi-indices  $oldsymbol{lpha} \in \mathbb{N}_0^d$  yields

(101) 
$$[H\Phi^{\ell}(\hat{x}) + M)]^{\alpha} = \sum_{j=0}^{\alpha} \sum_{(|\beta_k|)_k = j} \left( \prod_{k=1}^{d} C_k^H[j_k, \alpha_k, \beta_k] \right) \prod_{j=1}^{d} \Phi_j^{\sum_{k=1}^{d} \beta_k}(\hat{x}_j),$$

where the abbreviation  $\sum\limits_{(|m{eta}_k|)_k=m{j}}:=\sum\limits_{|m{eta}_1|=j_1}\dots\sum\limits_{|m{eta}_d|=j_d}$  is used.

By exploitation of the layerwise tensor train representation of  $f_\ell$  from (63) and using the rank-1 stable map (96), the high-dimensional integral over  $X^\ell$  reduces to

(102) 
$$\int\limits_{Y} \tilde{T}(x)^{\alpha} \tilde{f}_0(x) d\lambda(x) = \sum_{j=0}^{\alpha} \sum_{(|\beta_k|)_k = j} \left( \prod_{k=1}^d C_k^H[j_k, \alpha_k, \beta_k] \right) \sum_{k=0}^{r_\ell} \times$$

(103) 
$$\times \prod_{i=1}^{d} \int_{\hat{X}_{i}} \left[ \hat{f}_{\ell,i}[k_{i-1}, k_{i}] \otimes \boldsymbol{\varPhi}_{i}^{\sum\limits_{k=1}^{d} \boldsymbol{\beta}_{k}} \otimes h_{i} \right] (\hat{x}_{i}) \, \mathrm{d}\hat{x}_{i}.$$

Note that the right-hand side is composed via decoupled one dimensional integrals only. We point out that while the structure is simplified, the definition of  $\Phi_j$  in (100) a priori results in several integrals (indexed by  $\sum_{k=1}^d \beta_k$ ). These integrals, whose number depends on the cardinality of  $\alpha$ , have to be computed. This simplifies in several cases, e.g. when  $\Phi^\ell$  transforms the spherical coordinate system to Cartesian coordinates.

**Moment computation using spherical coordinates**. In the special case that  $\Phi^{\ell}$  is the multivariate polar transformation of Example 2.2, the number of distinct computation of integrals from (103) reduces significantly. Recall that  $\hat{x}_1 = \rho$ ,  $\hat{x}_{2:d} = \theta = (\theta_0, \dots, \theta_{d-2})$  and let  $\beta_i^k := (\beta_k)_i$  as the i-th entry of

 $\beta_k$ . We find that

(104) 
$$\mathbf{\Phi}_{1}^{\sum_{k=1}^{d} \beta_{k}}(\rho) = \rho^{|\mathbf{j}|},$$

$$\mathbf{\Phi}_{2}^{k=1}(\theta_{0}) = \cos^{\left(\sum_{k=1}^{d} \beta_{k}^{k}\right)}(\theta_{0}) \sin^{\left(\sum_{k=1}^{d} \beta_{2}^{k}\right)}(\theta_{0}),$$

$$\mathbf{\Phi}_{i+2}^{\stackrel{d}{\stackrel{}{\sim}} \boldsymbol{\beta}_k}(\boldsymbol{\theta}_i) = \sin^{\left(\sum\limits_{l=1}^i\sum\limits_{k=1}^d\beta_l^k\right)}(\boldsymbol{\theta}_i)\cos^{\left(\sum\limits_{k=1}^d\beta_{i+1}^k\right)}(\boldsymbol{\theta}_i), \qquad 1 \le i \le d-2.$$

It is notable that the exponential complexity due to the indexing of  $\sum_{k=1}^d \beta_k$  reduces to linear complexity in  $|\alpha|$ . More precisely, the amount of exponents in (104) - (106) is linear in the dimensions since the sums only depend on  $|\alpha|$ , leading to  $\mathcal{O}(|\alpha|d)$  different integrals that may be precomputed for each pair  $(k_{i-1},k_i)$ . This exponential complexity in the rank vanishes in the presence of an approximation basis associated with each coordinate dimension as defined in Section 3.

6.4. **Computation of marginals.** In probability theory and statistics, marginal distributions and especially marginal probability density functions provide insights into an underlying joint density by means of lower dimensional functions that can be visualized. The computation of marginal densities is a frequent problem encountered e.g. in parameter estimation and when using sampling techniques since histograms and corner plots provide easy access to (in general high-dimensional) integral quantities.

In contrast to the Markov chain Monte Carlo algorithm, the previously presented method of a layer based surrogate for the Lebesgue density function  $f\colon Y=\mathbb{R}^d\to\mathbb{R}$  allows for a functional representation and approximation of marginal densities without additional evaluations of f.

For simplicity, for  $y \in Y$  and  $i=1,\ldots,d$  define  $y_{-i}=(y_1,\ldots,y_{i-1},y_{i+1},\ldots y_d)$  as the marginalized variable where the i-th component is left out and  $f(y_{-i},y_i):=f(y)$ . Then, for given  $i=1,\ldots,d$ , the i-th marginal density reads

(107) 
$$df_i(y_i) := \int_{\mathbb{R}^{d-1}} f(y_{-i}, y_i) d\lambda(y_{-i}).$$

Computing this high-dimensional integral by quadrature or sampling is usually infeasible and the transport map approach as given by (4) fails since the map  $T\colon X\to Y$  can not be used directly in (107). Alternatively, the we can represent  $\mathrm{d} f_i\colon\mathbb{R}\to\mathbb{R}$  in a given orthonormal basis  $\{\varphi_j\}_{j=1}^{N_\varphi}$  and consider

(108) 
$$\mathrm{d}f_i(y_i) = \sum_{i=1}^{N_\varphi} \beta_j \varphi_j(y_i),$$

where  $\beta_j$ ,  $j=1,\ldots,N_{arphi}$  denotes the  $L^2(\mathbb{R})$  projection coefficient

(109) 
$$\beta_j := \int_{\mathbb{P}} \varphi_j(y_i) \mathrm{d}f_i(y_i) \mathrm{d}\lambda(y_i).$$

With this the transport map approach can in fact be employed and arguments as described in the computation of moments in Section 6.3 are valid.

A convenient basis is given by monomials since (109) then simplifies to

(110) 
$$\beta_j = \int_{\mathbb{R}^d} y_i^j f(y) \mathrm{d}\lambda(y),$$

which is the moment corresponding to the multi-index  $\alpha=(\alpha_k)_{k=1}^d\in\mathbb{N}^d$  with  $\alpha_k=\delta_{k,j}$ . Alternativly indicator functions may be considered in the spirit of histograms.

6.5. More general quantities of interest. The main question of applicability of a novel presented surrogate is probably if for a given quantity of interest (QoI)  $Q: Y \to \mathbb{R}$  the expectation

(111) 
$$\mathbb{E}\left[Q\right] = \int_{V} Q(y)f(y)d\lambda(y)$$

can be computed efficiently.

In the previous subsections we already discussed this issue for moments in Section 6.3 and basis representations of marginals in Section 6.4 since in those cases the structure of Q allows for direct computations of the integrals via tensor contractions. For more involved choices of the QoI we suggest a universal sampling approach by repeated evaluation of the low-rank surrogate. More precisely, by application of the transformation approach we can approximate

(112) 
$$\mathbb{E}[[Q] \approx \sum_{\ell=1}^{L} \int_{\hat{X}^{\ell}} Q \circ \tilde{T} \circ \Phi^{\ell}(\hat{x}) \tilde{f}_{0}^{\ell, \mathrm{TT}}(\hat{x}) | \det \left[ \mathcal{J}_{\Phi^{\ell}} \right](\hat{x}) | \mathrm{d}\lambda(\hat{x})$$

and replace the integrals over  $\hat{X}^\ell$  by Monte Carlo estimates with samples according to the (normalized) weight  $|\det{[\mathcal{J}_{\varPhi^\ell}]}|$ . Those samples can be obtained by uniform sampling on the tensor spaces  $\hat{X}^\ell$  and the inverse transform approach as mentioned in the paragraph **Generating Samples** of Section 5. Alternatively, efficient MCMC sampling by marginalisation can be employed [49].

#### 7. NUMERICAL VALIDATION AND APPLICATIONS

This section is devoted to a numerical validation of the proposed Algorithm 1 using various types of transformations T while employing it with practical applications. We focus on three example settings. The first consists of an artificial Gaussian posterior density, which could be translated to a linear forward model and Gaussian prior assumptions in the Bayesian setting. Second, we study the approximation behaviour under non-exact transport and conclude as a third setting with an actual Bayesian inversion application governed by the log-normal Darcy flow problem of Section 6.2.

7.1. **Validation experiment 1: Gaussian density.** In this experiment we confirm the theoretical results from Section 4 and verify the numerical algorithm. Even though the approximation of a Gaussian density is not a challenging task for the proposed algorithm, it can be seen as the most basic illustration and reveals interesting properties of the approximation.

We consider the target or posterior density determined by a Gaussian density with covariance matrix  $\Sigma \in \mathbb{R}^{d,d}$  and mean  $\mu \in \mathbb{R}^d$  as

(113) 
$$\frac{\mathrm{d}\pi}{\mathrm{d}\lambda}(x) = f(x) = C \exp\left(-\frac{1}{2}\|x - \mu\|_{\Sigma^{-1}}^2\right),$$

where  $C=(2\pi)^{-d/2}\det \varSigma^{-1/2}$  is the normalizing factor of the multivariate Gaussian density. We set the covariance operator such that the Gaussian density belongs to uncorrelated random variables, i.e.  $\varSigma$  exhibits a diagonal structure, and it holds for some  $0<\sigma\ll 1$  that  $\varSigma=\sigma^2I$ . This Gaussian setting has several benefits as a validation setting. On the one hand, we have explicit access to the quantities that are usually of interest in Bayesian inference like the mean, covariance, normalisation constant and marginals. On the other hand, the optimal transport to a standard normal density

(114) 
$$f_0(x) = (2\pi)^{-d/2} \exp\left(-\frac{1}{2}||x||^2\right)$$

is given by an affine linear function, defined via mean  $\mu$  and covariance  $\Sigma$ . We subsequently employ the multivariate polar transformation from Example 2.2 and expect a rank-1 structure in the reconstruction of the local approximations of the (perturbed) prior.

dimension	$\sigma^2 = 10^{-2}$	$\sigma^2 = 10^{-4}$	$\sigma^2 = 10^{-6}$	$\sigma^2 = 10^{-8}$
2	$5.24 \cdot 10^{-11}$	$1.09 \cdot 10^{-10}$	$2.8 \cdot 10^{-11}$	$9.3 \cdot 10^{-11}$
4	$2.21 \cdot 10^{-10}$	$4.57 \cdot 10^{-10}$	$5.48 \cdot 10^{-10}$	$3.4 \cdot 10^{-10}$
6	$5.01 \cdot 10^{-11}$	$9.5 \cdot 10^{-11}$	$7.49 \cdot 10^{-11}$	$6.19 \cdot 10^{-10}$
8	$1.48 \cdot 10^{-11}$	$8.21 \cdot 10^{-10}$	$2.99 \cdot 10^{-10}$	$2.1 \cdot 10^{-10}$
10	$2.91 \cdot 10^{-9}$	$9.61 \cdot 10^{-10}$	$4.43 \cdot 10^{-11}$	$2.46 \cdot 10^{-9}$

TABLE 1. Numerical approximation of Z in the Gaussian example. Absolute error of normalization constant computed via TT surrogate to Z=1.

The remainder of this section considers multiple choices of  $\sigma \in \mathbb{R}$  and  $d \in \mathbb{N}$  and highlights the stability of our method under decreasing variance (i.e. with higher density concentration) and increasing dimension. This is accomplished by comparing approximations with their exact counterparts. More specifically, the error of the normalization constant is observed, i.e.,

(115) 
$$\operatorname{err}_{Z} := |1 - Z_{h}|,$$

the relative  $\ell^2$  error of the mean and covariance

(116) 
$$\operatorname{err}_{\mu} := \|\mu - \mu_h\|_{\ell^2(\mathbb{R}^d)} \|\mu\|_{\ell^2(\mathbb{R}^d)}^{-1}, \quad \operatorname{err}_{\Sigma} := \|\Sigma - \Sigma_h\|_{\ell^2(\mathbb{R}^{d,d})} \|\Sigma\|_{\ell^2(\mathbb{R}^{d,d})}^{-1}$$

and the deviation in terms of the Kullback-Leibler divergence (68). Computing the Kullback-Leibler divergence is accomplished by Monte Carlo samples  $(x_i)_{i=1}^{N_{\rm KL}}$  of the posterior (i.e. in this case the multivariate Gaussian posterior) to compute the empirical approximation

(117) 
$$d_{\mathrm{KL}}(\pi, \pi_h) = \int_{\mathbb{R}^d} \log \left( \frac{f(x)}{f_h(x)} \right) f(x) d\lambda(x) \approx \frac{1}{N_{\mathrm{KL}}} \sum_{i=1}^{N_{\mathrm{KL}}} \log \left( \frac{f(x_i)}{f_h(x_i)} \right).$$

The index h generically denotes the employed approximation (63). In the numerical experiments the convergence of these error measures are depicted with respect to the amount of calls to the forward model (i.e. the Gaussian posterior density), the disrectization of the radial component  $\rho \in [0,\infty)$  in the polar coordinate system and the number of samples on each layer  $X^\ell$ ,  $\ell=1,\ldots,L$  for fixed  $L\in\mathbb{N}$ .

In Table 7.1 we show  $\operatorname{err}_Z$  for different choices of  $\sigma$  and d. The experiment comprises radial discretizations  $0=\rho_0<\rho_1<\ldots<\rho_L=10$  with L=19 equidistantly chosen layers and 1000 samples of  $f_0$  on each resulting subdomain  $X^\ell$ . The generated basis (79) contains polynomials of maximal degree 7 in  $\rho_\ell,\ell=0,\ldots,L$  and constant functions in every angular direction. The choice of constant functions relies on the assumption that the perturbed prior that has to be approximated corresponds to the polar transformation of (114), which is a function in  $\rho$  only. Additional numerical test show that even much fewer samples and a larger basis lead to the the assumed rank-1 structure.

Figure 4 depicts the error measures for increasing amount of layers, which increases the sample sizes linearly. Since drawing N=100 samples on each layer, going from L to L+1 samples implies additional N calls of the function f.

In Figure 5 we compare the number of calls of the posterior density f explicitly. Here, the presented low-rank surrogate is again constructed on an increasing amount of layers, whereas the Monte Carlo estimates are computed using a Markov-Chain Monte Carlo algorithm and subsequent empirical integration of the error quantity.

7.2. **Validation experiment 2: Perturbation of exact transport.** In the following experiment we consider the so-called "banana example" as posterior density. For this let  $f_0$  be the density of a standard

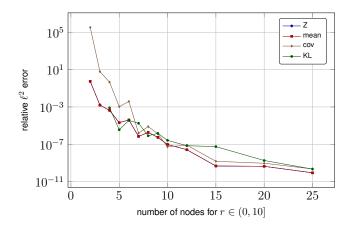


FIGURE 4. Approximation results for (1) normalisation constant  $(err_Z)$ , (2) mean  $(err_\mu)$ , (3) covariance  $(err_\Sigma)$  and (4) Kullback-Leibler divergence  $(d_{KL})$ . For the Gaussian setting we set d=6, mean= 1,  $\sigma^2=1e-6$  and take 100 samples on each layer for the surrogate construction. The Kullback-Leibler divergence is computed with  $N_{\rm KL}=10^4$  samples.

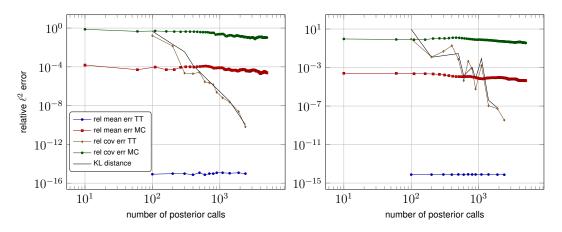


FIGURE 5. Example: Gaussian density with d=3 (left) and d=10 (right) for  $\mu=1$  and noise level  $\sigma=10^{-7}$ . We compare an MCMC approximation and our reconstruction setting in terms of function evaluations. The results for  $\mathrm{err}_{\mu},\,\mathrm{err}_{\Sigma}$  and  $d_{\mathrm{KL}}$  are shown. Note that the exact mean and covariance are known in this example.

normal Gaussian measure and let  $T_{\Sigma}$  be the affine transport of  $\mathcal{N}(0,I)$  to the Gaussian measure  $\mathcal{N}(0,\Sigma)$ . Furthermore, set

(118) 
$$T_2(x) = \begin{pmatrix} x_1 \\ x_2 - (x_1^2 + 1) \end{pmatrix}.$$

The exact transport T from  $\mathcal{N}(0,I)$  to the curved and concentrated banana distribution with density f is then given by

(119) 
$$T(x) = T_2 \circ T_{\Sigma}(x), \quad \Sigma = \begin{pmatrix} 1 & 0.9 \\ 0.9 & 1 \end{pmatrix}.$$

Note that the employed density can be transformed into a Gaussian density using a quadratic transport function. For this experiment, we employ transport maps  $\tilde{T}$  of varying accuracy for the pull-back of the posterior density to a standard Gaussian. In particular we use an affine transport  $\tilde{T}_1$  and a quadratic transport  $T_2$  to build an approximation of the optimal map  $\tilde{T}$  as convex combination

(120) 
$$\tilde{T}(x) = (1-t) \, \tilde{T}_1(x) + t \, T(x), \quad t \in [0,1].$$

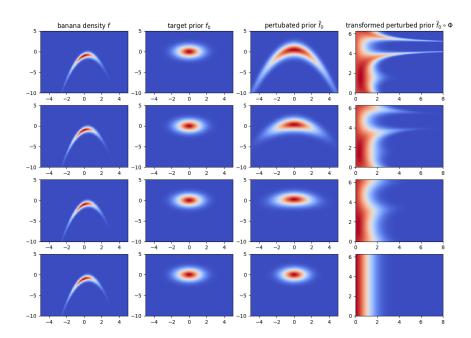


FIGURE 6. Illustration of the effect of the perturbed transport in (120) for t=0,0.25,0.5,1. (top to bottom).

For t=1, the transport map is optimal since it generates the desired reference density and for  $0 \le t < 1$  a perturbed prior density can be assumed with perturbation magnitude depending on the choice of t. The impact of the perturbed transport is visualized in Figure 6.

In Figure 7 we show the impact of an inexact transport on the approximation results in terms of  $\mathrm{err}_\mu$  and  $\mathrm{err}_\Sigma$ . For the considered target density, mean and covariance are known analytically and hence no reference sampling has to be carried out. We additionally employ an MCMC sampling to show the improvement due to the additional low-rank reconstruction. For the optimal transport map one observes that the surrogate reconstruction reduces to the approximation of a rank-1 Gaussian density, which can be done efficiently with few evaluations of f. If the transport becomes less accurate, additional samples are required to ensure at least comparable or better approximation results than MCMC.

7.3. Bayesian inversion with log-normal Darcy forward model. Revisiting the example of Section 6.2, we consider the elliptic diffusion problem with a log-normal random parametric permeability coefficient. The considered field in  $L^2(Y, L^\infty(D))$  takes the form

(121) 
$$a(x,y) = \exp\left(\sum_{i=1}^{d} a_k(x)y_k\right)$$

where the  $y_k$  correspond to random variables with law  $\mathbb{N}(0,1)$  and the  $L^2(D)$  basis functions are planar Fourier cosine modes. A detailed description and an adaptive Galerkin approach to solve the forward problem can be found in [14]. For the inverse problem, the observation operator is modelled by J=144 equidistantly distributed observations in  $D=[0,1]^2$  of the solution  $q(y^*)\in H^1_0(D)$  for some  $y^*\in Y=\mathbb{R}^d$ , which is drawn from a standard normal distribution. Additionally, the observations are perturbed by some centered Gaussian noise with covariance  $\sigma I$  with  $\sigma=10^{-7}$ .

To obtain the desired error measures, we employ reference computations that involve adaptive quadrature for the two dimensional example in Figure 8 and Markov-Chain Monte Carlo integration with  $10^6$  steps of the chain and a burn-in time of 1000 samples for the experiment in Figure 9. For the reconstruction algorithm on every layer we employ 100 samples each, which is depicted in Figure 9.

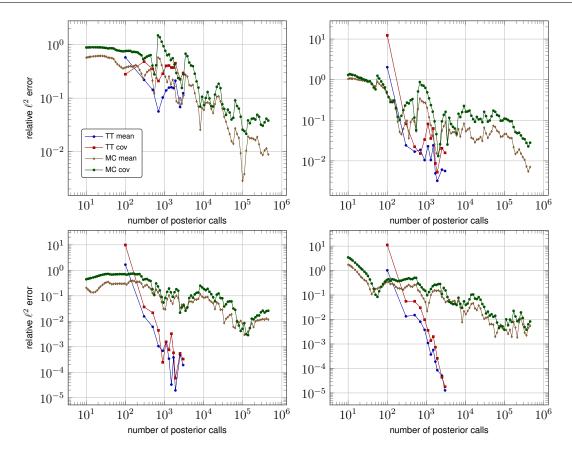


FIGURE 7. Convex combination of affine and quadratic transport for banana posterior. Affine linear transport (t=0 top left), (t=0.25) transport (top right), (t=0.5) transport (bottom left) and exact quadratic transport (t=1, bottom right). Error quantities  $\exp_{\mu}$  and  $\exp_{\Sigma}$  for the employed tensor train surrogate and a Markov-Chain Monte Carlo approximation in terms of the number of calls to the posterior function. The surrogate is reconstructed from 100 samples per layer yielding a tensor with radial basis up to polynomial degree 9 and Fourier modes up to degree 20.

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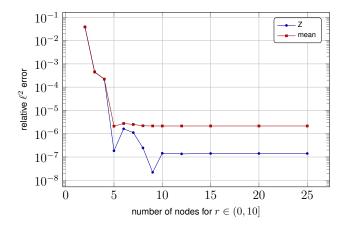


FIGURE 8. Comparison of the computed reference and the low-rank surrogate of (1) normalisation constant  $(\operatorname{err}_Z)$ , and (2) mean  $(\operatorname{err}_\mu)$ . For the Darcy setting with d=2 we observe 144 nodes in the physical domain. The measurements are perturbed by Gaussian noise with deviation  $\eta=1e-7$ . We employ an adaptive quadrature in the two dimensional space to obtain the reference quantities. The stagnation of the curves are due to non-optimal reference solutions or putting it more precicely, the TT approximation yields equivalent results to adaptive quadrature when taking 5 nodes of refinement.

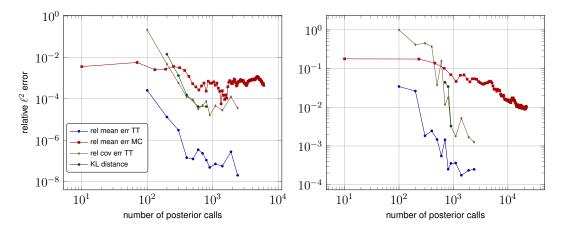


FIGURE 9. Darcy example with d=2 (left) and d=10 (right). Comparison of an MCMC method and the low-rank surrogate for computing the mean error  $(\mathrm{err}_{\mu})$  with respect to the number of calls to the solution of the forward problem. The reference mean is computed with  $10^6$  MCMC samples. Additionally the KL divergence is shown, which is computed using empirical integration.

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