
A Unifying Framework of Anytime Sparse Gaussian Process Regression Models with Stochastic Variational Inference for Big Data

Trong Nghia Hoang[†]
Quang Minh Hoang[†]
Kian Hsiang Low[†]

NGHIAHT@COMP.NUS.EDU.SG
HQMINH@COMP.NUS.EDU.SG
LOWKH@COMP.NUS.EDU.SG

[†]Department of Computer Science, National University of Singapore, Republic of Singapore

Abstract

This paper presents a novel unifying framework of anytime *sparse Gaussian process regression* (SGPR) models that can produce good predictive performance fast and improve their predictive performance over time. Our proposed unifying framework reverses the variational inference procedure to theoretically construct a non-trivial, concave functional that is maximized at the predictive distribution of any SGPR model of our choice. As a result, a stochastic natural gradient ascent method can be derived that involves iteratively following the stochastic natural gradient of the functional to improve its estimate of the predictive distribution of the chosen SGPR model and is guaranteed to achieve asymptotic convergence to it. Interestingly, we show that if the predictive distribution of the chosen SGPR model satisfies certain decomposability conditions, then the stochastic natural gradient is an unbiased estimator of the exact natural gradient and can be computed in constant time (i.e., independent of data size) at each iteration. We empirically evaluate the trade-off between the predictive performance vs. time efficiency of the anytime SGPR models on two real-world million-sized datasets.

1. Introduction

A *Gaussian process regression* (GPR) model is a Bayesian nonparametric model for performing nonlinear regression that provides a Gaussian predictive distribution with formal measures of predictive uncertainty. The expressivity of a *full-rank GPR* (FGPR) model, however, comes at a cost of cubic time in the size of the data, thus rendering it computationally impractical for training with massive datasets. To improve its scalability, a number of *sparse GPR* (SGPR) models (Lázaro-Gredilla et al., 2010; Quiñero-Candela

& Rasmussen, 2005; Snelson & Ghahramani, 2007; Titsias, 2009) exploiting low-rank approximate representations have been proposed, many of which share a similar structural assumption of conditional independence (albeit of varying degrees) based on the notion of inducing variables (Section 2) and consequently incur only linear time in the data size. The work of Quiñero-Candela & Rasmussen (2005) has in fact presented a unifying view of such SGPR models, which include the *subset of regressors* (SoR) (Smola & Bartlett, 2001), *deterministic training conditional* (DTC) (Seeger et al., 2003), *fully independent training conditional* (FITC) (Snelson & Ghahramani, 2005), *fully independent conditional* (FIC), *partially independent training conditional* (PITC) (Schwaighofer & Tresp, 2003), and *partially independent conditional* (PIC) (Snelson & Ghahramani, 2007) approximations. To scale up these SGPR models further for performing real-time predictions necessary in many time-critical applications and decision support systems (e.g., ocean sensing (Cao et al., 2013; Dolan et al., 2009; Low et al., 2008; 2009; 2011; 2012; Podnar et al., 2010), traffic monitoring (Chen et al., 2012; 2013b; 2015; Hoang et al., 2014a;b; Low et al., 2014a;b; Ouyang et al., 2014; Xu et al., 2014; Yu et al., 2012)), the work of Gal et al. (2014) has parallelized DTC while that of Chen et al. (2013a) has parallelized FITC, FIC, PITC, and PIC to be run on multiple machines. The recent work of Low et al. (2015) has produced a spectrum of SGPR models with PIC and FGPR at the two extremes that are also amenable to parallelization on multiple machines. Ideally, these parallel SGPR models can reduce the incurred time of their centralized counterparts by a factor close to the number of machines. In practice, since the number of machines is limited due to budget constraints, their incurred time will still grow with an increasing size of data. Like their centralized counterparts, they can be trained using all the data.

A more affordable alternative is to instead train a SGPR model in an anytime fashion with a small, randomly sampled subset of the data at each iteration, which requires only a single machine. To the best of our knowledge, the

only notable anytime SGPR model (Hensman et al., 2013) exploits a result of Titsias (2009) that DTC can alternatively be obtained using variational inference by minimizing the *Kullback-Leibler* (KL) distance between the variational approximation and the GP posterior distribution of some latent variables given the data, from which a *stochastic natural gradient ascent* (SNGA) method can be derived to achieve an asymptotic convergence of its predictive performance to that of DTC while incurring constant time per iteration. This anytime variant of DTC promises a huge speedup if the number of sampled subsets of data needed for convergence is much smaller than the total number of possible disjoint subsets that can be formed and sampled from all the data. But, it can be observed in our experiments (Section 5) that DTC often does not predict as well as the other SGPR models (except SoR) encompassed by the unifying view of Quiñonero-Candela & Rasmussen (2005) because it imposes the most restrictive structural assumption (Snelson, 2007). This motivates us to consider the possibility of constructing an anytime variant of any SGPR model of our choice whose derived SNGA method can achieve an asymptotic convergence of its predictive performance to that of the chosen SGPR model while preserving constant time per iteration. However, no alternative formulation based on variational inference exists for any SGPR model other than DTC in order to derive such a SNGA method.

To address the above challenge, this paper presents a novel unifying framework of anytime SGPR models that can produce good predictive performance fast and improve their predictive performance over time. Our proposed unifying framework, perhaps surprisingly, reverses the variational inference procedure to theoretically construct a non-trivial, concave functional (i.e., of distributions) that is maximized at the predictive distribution of any SGPR model of our choice (Section 3). Consequently, a SNGA method can be derived that involves iteratively following the stochastic natural gradient of the functional to improve its estimate of the predictive distribution of the chosen SGPR model and is guaranteed to achieve asymptotic convergence to it. Interestingly, we show that if the predictive distribution of the chosen SGPR model satisfies certain decomposability conditions (e.g., DTC, FITC, PIC), then the stochastic natural gradient is an unbiased estimator of the exact natural gradient and can be computed in constant time (i.e., independent of data size) at each iteration (Section 4). We empirically evaluate the trade-off between the predictive performance vs. time efficiency of the anytime SGPR models spanned by our unifying framework (i.e., including state-of-the-art anytime variant of DTC (Hensman et al., 2013)) on two real-world million-sized datasets (Section 5).

2. Background and Notations

Full-Rank Gaussian Process Regression (FGPR). Let \mathcal{X}

be a set representing the input domain such that each d -dimensional input feature vector $\mathbf{x} \in \mathcal{X}$ is associated with a latent output variable $f_{\mathbf{x}}$. Let $\{f_{\mathbf{x}}\}_{\mathbf{x} \in \mathcal{X}}$ denote a *Gaussian process* (GP), that is, every finite subset of $\{f_{\mathbf{x}}\}_{\mathbf{x} \in \mathcal{X}}$ follows a multivariate Gaussian distribution. Then, the GP is fully specified by its *prior* mean $\mathbb{E}[f_{\mathbf{x}}]$, which we assume to be zero for notational simplicity, and covariance $k_{\mathbf{x}\mathbf{x}'} \triangleq \text{cov}[f_{\mathbf{x}}, f_{\mathbf{x}'}]$ for all $\mathbf{x}, \mathbf{x}' \in \mathcal{X}$. Given a column vector $\mathbf{y}_{\mathcal{D}} \triangleq (y_{\mathbf{x}'})_{\mathbf{x}' \in \mathcal{D}}^{\top}$ of noisy observed outputs $y_{\mathbf{x}'} \triangleq f_{\mathbf{x}'} + \varepsilon$ for some set $\mathcal{D} \subset \mathcal{X}$ of training inputs where $\varepsilon \sim \mathcal{N}(0, \sigma_n^2)$ and σ_n^2 is the noise variance, a FGPR model can perform probabilistic regression by providing a GP predictive distribution $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}) = \mathcal{N}(K_{\mathbf{x}\mathcal{D}}(K_{\mathcal{D}\mathcal{D}} + \sigma_n^2 I)^{-1} \mathbf{y}_{\mathcal{D}}, k_{\mathbf{x}\mathbf{x}} - K_{\mathbf{x}\mathcal{D}}(K_{\mathcal{D}\mathcal{D}} + \sigma_n^2 I)^{-1} K_{\mathcal{D}\mathbf{x}})$ of the latent output $f_{\mathbf{x}}$ for any test input $\mathbf{x} \in \mathcal{X}$ where $K_{\mathbf{x}\mathcal{D}} \triangleq (k_{\mathbf{x}\mathbf{x}'})_{\mathbf{x}' \in \mathcal{D}}$, $K_{\mathcal{D}\mathcal{D}} \triangleq (k_{\mathbf{x}'\mathbf{x}''})_{\mathbf{x}', \mathbf{x}'' \in \mathcal{D}}$, and $K_{\mathcal{D}\mathbf{x}} \triangleq K_{\mathbf{x}\mathcal{D}}^{\top}$. Computing the GP predictive distribution incurs $\mathcal{O}(|\mathcal{D}|^3)$ time due to inversion of $K_{\mathcal{D}\mathcal{D}}$, hence causing the FGPR model to scale poorly in the size $|\mathcal{D}|$ of data.

Sparse Gaussian Process Regression (SGPR). To improve the scalability of the FGPR model, the SGPR models encompassed by the unifying view of Quiñonero-Candela & Rasmussen (2005) exploit a vector $\mathbf{f}_{\mathcal{U}} \triangleq (f_{\mathbf{x}'})_{\mathbf{x}' \in \mathcal{U}}^{\top}$ of $|\mathcal{U}|$ inducing output variables for some small set $\mathcal{U} \subset \mathcal{X}$ of inducing inputs (i.e., $|\mathcal{U}| \ll |\mathcal{D}|$) for approximating the GP predictive distribution $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}})$. Specifically, they share a similar structural assumption (Snelson & Ghahramani, 2007) that the joint distribution of $f_{\mathbf{x}}$ and $\mathbf{f}_{\mathcal{D}} \triangleq (f_{\mathbf{x}'})_{\mathbf{x}' \in \mathcal{D}}^{\top}$ conditioned on $\mathbf{f}_{\mathcal{U}}$ factorizes across a pre-defined partition of the input domain \mathcal{X} into P disjoint subsets $\mathcal{X}_1, \dots, \mathcal{X}_P$. That is, supposing $\mathbf{x} \in \mathcal{X}_P$,

$$p(f_{\mathbf{x}}, \mathbf{f}_{\mathcal{D}} | \mathbf{f}_{\mathcal{U}}) = p(f_{\mathbf{x}} | \mathbf{f}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \prod_{i=1}^P p(\mathbf{f}_{\mathcal{D}_i} | \mathbf{f}_{\mathcal{U}}) \quad (1)$$

where $\mathbf{f}_{\mathcal{D}_i} \triangleq (f_{\mathbf{x}'})_{\mathbf{x}' \in \mathcal{D}_i}^{\top}$ denotes a column vector of latent outputs for the disjoint subset $\mathcal{D}_i \triangleq \mathcal{X}_i \cap \mathcal{D} \subset \mathcal{D}$ of training inputs for $i = 1, \dots, P$. Using (1), the GP predictive distribution $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}})$ reduces to

$$p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}) = \int p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}}) d\mathbf{f}_{\mathcal{U}} \quad (2)$$

$$\simeq \int q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) q^*(\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{U}} \quad (3)$$

where $\mathbf{y}_{\mathcal{D}_P} \triangleq (y_{\mathbf{x}'})_{\mathbf{x}' \in \mathcal{D}_P}^{\top}$ is a vector of noisy observed outputs for the subset \mathcal{D}_P of training inputs, derivation of (2) is in Appendix C.1, and $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$ and $p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ are, respectively, approximated by $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$ and $q^*(\mathbf{f}_{\mathcal{U}})$ in (3), as discussed in Remarks 1 and 2 below.

Remark 1. PIC sets $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$ as the exact test conditional $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$. The other SGPR models have additionally assumed conditional independence of $f_{\mathbf{x}}$ and $\mathbf{f}_{\mathcal{D}_P}$ given $\mathbf{f}_{\mathcal{U}}$, thus resulting in $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) = p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ and $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \triangleq p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ (e.g., see eq. 5 in (Titsias,

2009) for the case of DTC). Then, $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_F}, \mathbf{f}_{\mathcal{U}})$ can be computed in $\mathcal{O}(|\mathcal{U}|^3)$ time and space (i.e., $|\mathcal{D}_F| = \mathcal{O}(|\mathcal{U}|)$ for the case of PIC), as shown in Appendix D.4.

Remark 2. The work of Titsias (2009) has approximated $p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ with a choice of $q^*(\mathbf{f}_{\mathcal{U}})$ whose resulting predictive distribution (3) coincides with that of DTC. Interestingly, other choices of $q^*(\mathbf{f}_{\mathcal{U}})$ can be derived to induce the predictive distributions (3) of the other SGPR models and computed in $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time, as shown in Appendix D.1.

Instead of computing $q^*(\mathbf{f}_{\mathcal{U}})$ directly that becomes prohibitively expensive (i.e., $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time) for massive datasets, we will derive a stochastic natural gradient ascent method that is guaranteed to achieve asymptotic convergence to $q^*(\mathbf{f}_{\mathcal{U}})$ of any SGPR model of our choice while incurring only $\mathcal{O}(|\mathcal{U}|^3)$ time per iteration (Section 4), thus producing an anytime variant of the chosen SGPR model.

Variational Inference for DTC. Variational inference can be used to derive $q^*(\mathbf{f}_{\mathcal{U}})$ of DTC as follows: A variational approximation to the posterior distribution of some latent output variables (e.g., $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$) can be derived analytically by minimizing their KL distance, provided that it factorizes in some way or has some parametric form that is inexpensive to evaluate (Bishop, 2006). The work of Titsias (2009) parameterizes the variational approximation $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$ to the GP posterior distribution $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ by

$$q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \triangleq p(\mathbf{f}_{\mathcal{D}} | \mathbf{f}_{\mathcal{U}}) q(\mathbf{f}_{\mathcal{U}}) \quad (4)$$

where $p(\mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}})$ is the exact training conditional (8) and $q(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu, \Sigma)$. The result below reveals how μ and Σ , which depend on training data $(\mathcal{D}, \mathbf{y}_{\mathcal{D}})$, can be selected to minimize the KL distance $D_{\text{KL}}(q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \| p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}}))$:

Lemma 1 For any PDF $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$ and $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$,

$$\log p(\mathbf{y}_{\mathcal{D}}) = L(q) + D_{\text{KL}}(q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \| p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}}))$$

where the functional $L(q)$ is defined as

$$L(q) \triangleq \int q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \log \left(\frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})}{q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})} \right) d\mathbf{f}_{\mathcal{D}} d\mathbf{f}_{\mathcal{U}}. \quad (5)$$

Its proof is in Appendix C.2. Lemma 1 implies that minimizing $D_{\text{KL}}(q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \| p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}}))$ is equivalent to maximizing $L(q)$ since $p(\mathbf{y}_{\mathcal{D}})$ is constant with respect to $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$. Using the parameterization in (4), $L(q)$ becomes a concave function in μ and Σ that is maximized when its gradient is zero. So, $q(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu, \Sigma)$ can be optimized by solving for μ and Σ such that $\partial L/\partial \mu = 0$ and $\partial L/\partial \Sigma = 0$.

Remark 1. From Lemma 1, since $D_{\text{KL}}(\cdot|\cdot)$ is non-negative, $\log p(\mathbf{y}_{\mathcal{D}}) \geq L(q)$, which recovers the variational lower bound of Titsias (2009) by setting $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ as the GP joint distribution.

Remark 2. The work of Titsias (2009) is originally intended to jointly optimize $q(\mathbf{f}_{\mathcal{U}})$, inducing inputs \mathcal{U} , and hyperpa-

rameters of $k(\cdot, \cdot)$. In the context of our work here, by assuming \mathcal{U} and the hyperparameters to be given, the optimal $q(\mathbf{f}_{\mathcal{U}}) \equiv q^*(\mathbf{f}_{\mathcal{U}})$ induces the predictive distribution (3) of DTC when $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_F}, \mathbf{f}_{\mathcal{U}}) \triangleq p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ (Titsias, 2009).

3. Reverse Variational Inference

This section introduces a novel, interesting use of variational inference, which we term *reverse variational inference*, to theoretically construct a concave, differential functional $L(q)$ that is maximized at $q(\mathbf{f}_{\mathcal{U}}) \equiv q^*(\mathbf{f}_{\mathcal{U}})$ of any SGPR model of our choice; we call this requirement **R1**. The functional $L(q)$ allows us to derive a *stochastic natural gradient ascent* (SNGA) method (Section 4) that takes small, iterative steps in the direction of the stochastic natural gradient of $L(q)$ to improve its estimate $q(\mathbf{f}_{\mathcal{U}})$ of $q^*(\mathbf{f}_{\mathcal{U}})$ and is guaranteed to achieve asymptotic convergence of $q(\mathbf{f}_{\mathcal{U}})$ to $q^*(\mathbf{f}_{\mathcal{U}})$ if the step sizes are scheduled appropriately (Robbins & Monro, 1951). If the stochastic natural gradient of $L(q)$ can be computed in constant time (i.e., independent of data size $|\mathcal{D}|$) and we call this requirement **R2**, then such a SNGA method is desirable in practice due to its anytime behavior of improving the estimation of $q^*(\mathbf{f}_{\mathcal{U}})$ over time. In Section 4, we will establish sufficient conditions for $q^*(\mathbf{f}_{\mathcal{U}})$ to satisfy **R2**, thus entailing a unifying framework of anytime SGPR models.

Constructing $L(q)$ to Satisfy R1. Let $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$ be factorized according to (4) and $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ where μ^* and Σ^* depend on training data $(\mathcal{D}, \mathbf{y}_{\mathcal{D}})$. Our key idea is to derive a joint distribution $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ such that $L(q)$ is maximized at $q(\mathbf{f}_{\mathcal{U}}) \equiv q^*(\mathbf{f}_{\mathcal{U}})$ of any SGPR model of our choice, which remains largely unexplored except for DTC: A result of Titsias (2009) has established that $q^*(\mathbf{f}_{\mathcal{U}})$ of DTC (Appendix D.1.3) maximizes $L(q)$ when $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ coincides with the GP joint distribution, hence satisfying **R1** for DTC only. Such a functional $L(q)$ is then shown by Hensman et al. (2013) to satisfy **R2** and can consequently be exploited for deriving a SNGA method to produce an anytime variant of DTC. However, this work neither extends nor discusses how $L(q)$ can be derived for other choices of $q^*(\mathbf{f}_{\mathcal{U}})$ and the conditions under which they will satisfy **R1** and **R2**. We will address both these issues in this section and the next, respectively. In addition, the predictive performance of their SNGA method is severely limited by the highly restrictive structural assumption of DTC. Finally, their anytime variant of DTC turns out to be a special case spanned by our unifying framework of anytime SGPR models (Section 4).

For the rest of this section, we will first evaluate $L(q)$ to a concave function in μ and Σ (i.e., Theorems 1, 2, and 3) subject to our factorization of $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$ in (4) and $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ in (10). Then, we will show how the parameters defining $p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ can be appropriately selected such that the induced $L(q)$ (5) is maximized at

$q^*(\mathbf{f}_U) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ of our choice (i.e., Theorem 4).

Theorem 1 Let $q(\mathbf{f}_D | \mathbf{f}_U) \triangleq q(\mathbf{f}_D, \mathbf{f}_U) / q(\mathbf{f}_U)$. Then,

$$L(q) = \int q(\mathbf{f}_U) L_U(q) d\mathbf{f}_U - D_{\text{KL}}(q(\mathbf{f}_U) \| p(\mathbf{f}_U)) \quad (6)$$

where the functional $L_U(q)$ is defined as

$$L_U(q) \triangleq \int q(\mathbf{f}_D | \mathbf{f}_U) \log \frac{p(\mathbf{f}_D, \mathbf{y}_D | \mathbf{f}_U)}{q(\mathbf{f}_D | \mathbf{f}_U)} d\mathbf{f}_D. \quad (7)$$

Its proof is in Appendix C.3. The parameterization of $q(\mathbf{f}_U) \triangleq \mathcal{N}(\mu, \Sigma)$ and factorization of $q(\mathbf{f}_D, \mathbf{f}_U)$ using (4) entail $q(\mathbf{f}_D | \mathbf{f}_U)$ being set as the exact training conditional:

$$q(\mathbf{f}_D | \mathbf{f}_U) \triangleq p(\mathbf{f}_D | \mathbf{f}_U) = \mathcal{N}(P_{DU} \mathbf{f}_U, K_{DD} - Q_{DD}) \quad (8)$$

where $P_{DU} \triangleq K_{DU} K_{UU}^{-1}$ and $Q_{DD} \triangleq K_{DU} K_{UU}^{-1} K_{UD}$. Using (8), $L_U(q)$ is reduced to a quadratic function of \mathbf{f}_U :

Theorem 2 By substituting (8) into $L_U(q)$ (7),

$$L_U(q) = -\frac{1}{2\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top P_{DU} \mathbf{f}_U + \frac{1}{\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top \mathbf{y}_D + C \quad (9)$$

where the constant C absorbs all terms independent of \mathbf{f}_U .

Its proof is in Appendix B.1. Then, we factorize

$$p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D) \triangleq p(\mathbf{y}_D | \mathbf{f}_D) p(\mathbf{f}_D | \mathbf{f}_U) p(\mathbf{f}_U) \quad (10)$$

where $p(\mathbf{y}_D | \mathbf{f}_D) = \mathcal{N}(\mathbf{f}_D, \sigma_n^2 I)$, $p(\mathbf{f}_D | \mathbf{f}_U)$ is the exact training conditional (8), and we let $p(\mathbf{f}_U) \triangleq \mathcal{N}(\nu, \Lambda^{-1})$ (instead of defining $p(\mathbf{f}_U)$ as a GP prior by Titsias (2009) that makes $p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D)$ a GP joint distribution) where Λ denotes a precision matrix. Then, by applying Theorems 1 and 2, $L(q)$ becomes a concave function in both μ and Σ :

Theorem 3 By substituting $q(\mathbf{f}_U) \triangleq \mathcal{N}(\mu, \Sigma)$, $p(\mathbf{f}_U) \triangleq \mathcal{N}(\nu, \Lambda^{-1})$, and $L_U(q)$ (9) into $L(q)$ (6),

$$L(q) = -\frac{1}{2} \mu^\top \Psi \mu - \frac{1}{2} \text{tr}(\Psi \Sigma) + \frac{1}{2} \log |\Sigma| + \mu^\top \left((1/\sigma_n^2) P_{DU}^\top \mathbf{y}_D + \Lambda \nu \right) + C' \quad (11)$$

where $\Psi \triangleq (1/\sigma_n^2) P_{DU}^\top P_{DU} + \Lambda$ and the constant C' absorbs all terms independent of μ and Σ .

Its proof is in Appendix B.2. Using Theorem 3, the conditions for the parameters ν and Λ defining $p(\mathbf{f}_U)$ (or, equivalently, $p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D)$) can be determined such that $L(q)$ is maximized at $q(\mathbf{f}_U) \equiv q^*(\mathbf{f}_U)$ by making its derivatives with respect to μ and Σ go to zero at $(\mu, \Sigma) = (\mu^*, \Sigma^*)$:

Theorem 4 If ν and Λ satisfy the following conditions:

$$\Lambda \nu + \frac{1}{\sigma_n^2} P_{DU}^\top \mathbf{y}_D = \left(\frac{1}{\sigma_n^2} P_{DU}^\top P_{DU} + \Lambda \right) \mu^* \quad (12)$$

$$\text{and } \Lambda = \Sigma^{*-1} - \frac{1}{\sigma_n^2} P_{DU}^\top P_{DU}, \quad (13)$$

then $L(q)$ is maximized at $q(\mathbf{f}_U) \equiv q^*(\mathbf{f}_U)$.

Its proof is in Appendix B.3.

Remark. (12) and (13) define the space of feasible pairs (ν, Λ) guaranteeing that $L(q)$ is maximized at $(\mu, \Sigma) = (\mu^*, \Sigma^*)$. Interestingly, it is not necessary to explicitly solve for (ν, Λ) in order to construct $L(q)$ that is maximized at $q(\mathbf{f}_U) \equiv q^*(\mathbf{f}_U)$, as shown in (34) in Appendix B.3.

4. Anytime Sparse GP Regression Models

Using Theorem 4, a gradient ascent method that is guaranteed to achieve asymptotic convergence of (μ, Σ) to (μ^*, Σ^*) can now be derived. Specifically, it starts with randomly initialized $(\mu, \Sigma) = (\mu^0, \Sigma^0)$ and iterates the following gradient ascent update until convergence:

$$\mu^{t+1} = \mu^t + \rho_t \frac{\partial L}{\partial \mu}(\mu^t, \Sigma^t), \quad \Sigma^{t+1} = \Sigma^t + \rho_t \frac{\partial L}{\partial \Sigma}(\mu^t, \Sigma^t) \quad (14)$$

where ρ_t is the step size and $\frac{\partial L}{\partial \mu}(\mu^t, \Sigma^t)$ and $\frac{\partial L}{\partial \Sigma}(\mu^t, \Sigma^t)$ denote, respectively, $\partial L / \partial \mu$ and $\partial L / \partial \Sigma$ (i.e., see (35) and (36) in Appendix B.3 for their expressions) being evaluated at $(\mu, \Sigma) = (\mu^t, \Sigma^t)$. This method is guaranteed to converge if (a) $\sum_t \rho_t = +\infty$ and (b) $\sum_t \rho_t^2 < +\infty$, which is a well-known result in optimization. For example, one possible schedule is $\rho_t = \rho_0 / (1 + \tau \rho_0 t)^\kappa$ where τ , κ , and ρ_0 are determined empirically. However, evaluating the exact gradient ($\partial L / \partial \mu, \partial L / \partial \Sigma$) requires computing $q^*(\mathbf{f}_U)$ directly that incurs $\mathcal{O}(|\mathcal{D}| |\mathcal{U}|^2)$ time (Appendix D.1), which is prohibitively expensive for massive datasets.

Stochastic Gradient Ascent (SGA). To sidestep the above scalability issue, we adopt the *stochastic gradient ascent* (SGA) method (Robbins & Monro, 1951) that replaces the exact gradient in (14) with its stochastic gradient ($\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma$). The key idea is to iteratively compute ($\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma$) in an efficient manner by randomly sampling a small block of data of size $|\mathcal{U}|$ whose incurred time per iteration is independent of the data size $|\mathcal{D}|$. We will prove in Theorem 5 later that such a stochastic gradient is an unbiased estimator of the exact gradient. As a result, (14) is also guaranteed to converge using the above schedule of $\{\rho_t\}_t$. To derive ($\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma$), the following decomposability conditions for (μ^*, Σ^*) are necessary:

Decomposability Conditions. Let $F'(\mathcal{U})$ and $G'(\mathcal{U})$ ($F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i})$ and $G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i})$) denote arbitrary functions depending on \mathcal{U} (\mathcal{U} and $(\mathcal{D}_i, \mathbf{y}_{\mathcal{D}_i})$) only. The decomposability conditions for (μ^*, Σ^*) are

$$\Sigma^{*-1} = F'(\mathcal{U}) + \sum_{i=1}^P F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}), \quad (15)$$

$$\Sigma^{*-1} \mu^* = G'(\mathcal{U}) + \sum_{i=1}^P G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}). \quad (16)$$

Remark 1. Though (15) and (16) may appear rather awkward when viewed using the moment parameterization

$q^*(\mathbf{f}_U) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$, they can alternatively be perceived as simple additive decomposability of the natural parameters $\theta_1 \triangleq \Sigma^{*-1}\mu^*$ and $\theta_2 \triangleq -(1/2)\Sigma^{*-1}$, which define the canonical parameterization of $q^*(\mathbf{f}_U)$ (Appendix E).

Remark 2. Interestingly, this canonical view reveals a systematic way to construct new SGPR models from existing ones that satisfy (15) and (16): Given a set $\{q_m^*(\mathbf{f}_U)\}_{m=1}^M$ of M SGPR models specified by their respective canonical parameterizations $\{(\theta_{1,m}, \theta_{2,m})\}_{m=1}^M$ satisfying (15) and (16), if they share the same test conditional $q^*(f_x | \mathcal{Y}_{\mathcal{D}_P}, \mathbf{f}_U)$ in (3), then any new SGPR model constructed with $\hat{\theta}_1 \triangleq \sum_{m=1}^M \alpha_m \theta_{1,m}$ and $\hat{\theta}_2 \triangleq \sum_{m=1}^M \alpha_m \theta_{2,m}$ (i.e., $\{\alpha_m\}_{m=1}^M$ is a set of linear coefficients) also satisfies (15) and (16).

In practice, the decomposability conditions (15) and (16) are satisfied by many SGPR models that share a similar structural assumption of conditional independence in (1) such as SoR, DTC, FITC, FIC, PITC, and PIC. For interested readers, $\{F(\mathcal{U}, \mathcal{Y}_{\mathcal{D}_i})\}_{i=1}^P$, $\{G(\mathcal{U}, \mathcal{Y}_{\mathcal{D}_i})\}_{i=1}^P$, $F'(\mathcal{U})$, and $G'(\mathcal{U})$ of these SGPR models are derived in Appendix D.2. If our choice of (μ^*, Σ^*) (i.e., $q^*(\mathbf{f}_U)$) satisfies (15) and (16), then the stochastic gradient $(\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma)$ is an unbiased estimator of the exact gradient $(\partial L / \partial \mu, \partial L / \partial \Sigma)$:

Theorem 5 *Let \mathcal{S} be a set of i.i.d. samples (i.e., $|\mathcal{S}| > 0$) drawn from a uniform distribution over $\{1, 2, \dots, P\}$ and*

$$\frac{\partial \hat{L}}{\partial \mu} \triangleq G'(\mathcal{U}) - F'(\mathcal{U})\mu + \frac{P}{|\mathcal{S}|} \sum_{s \in \mathcal{S}} G(\mathcal{U}, \mathcal{Y}_{\mathcal{D}_s}) - F(\mathcal{U}, \mathcal{Y}_{\mathcal{D}_s})\mu, \quad (17)$$

$$\frac{\partial \hat{L}}{\partial \Sigma} \triangleq \frac{1}{2}\Sigma^{-1} - \frac{1}{2}F'(\mathcal{U}) - \frac{P}{2|\mathcal{S}|} \sum_{s \in \mathcal{S}} F(\mathcal{U}, \mathcal{Y}_{\mathcal{D}_s}). \quad (18)$$

If (μ^*, Σ^*) satisfies (15) and (16), then $\mathbb{E}[\partial \hat{L} / \partial \mu] = \partial L / \partial \mu$ and $\mathbb{E}[\partial \hat{L} / \partial \Sigma] = \partial L / \partial \Sigma$.

Its proof is in Appendix B.4.

Remark. By assuming $|\mathcal{D}_i| = \mathcal{O}(|\mathcal{U}|)$ for $i = 1, \dots, P$, computing stochastic gradient $(\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma)$ (i.e., (17) and (18)) incurs time independent of data size $|\mathcal{D}|$, in particular, $\mathcal{O}(|\mathcal{S}||\mathcal{U}|^3)$ time for SoR, DTC, FITC, FIC, PITC, and PIC (Appendix D.3) that reduces to $\mathcal{O}(|\mathcal{U}|^3)$ time by setting $|\mathcal{S}| = 1$ in our experiments. Also, since (μ, Σ) (i.e., $q(\mathbf{f}_U)$) is readily available from the SGA update (14), the prediction time (i.e., time incurred to analytically integrate $q^*(f_x | \mathcal{Y}_{\mathcal{D}_P}, \mathbf{f}_U)$ with $q^*(\mathbf{f}_U) \equiv q(\mathbf{f}_U)$ in (3)) is independent of $|\mathcal{D}|$ (Appendix D.4). So, if the number t of iterations of SGA update is much smaller than $\min(|\mathcal{D}|/|\mathcal{U}|, P)$, then the anytime variants spanned by our unifying framework achieve a huge computational gain (i.e., $\mathcal{O}(t|\mathcal{U}|^3)$ time) over their corresponding SGPR models that incur $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time (Appendix D.1).

Stochastic Natural Gradient Ascent (SNGA). As the standard gradient of a function (e.g., $L(q)$) only points in

the direction of the steepest ascent when the space of its parameters (e.g., (μ, Σ)) is Euclidean (Amari, 1998), the SGA update (14) has implicitly defined the parameter space of $q(\mathbf{f}_U)$ using the Euclidean distance between two candidate parameters, which unfortunately appears to be a poor dissimilarity measure between their corresponding distributions (Hoffman et al., 2013). To capture a more meaningful notion of dissimilarity, the parameter space of $q(\mathbf{f}_U)$ is redefined using the symmetrized KL distance, which is a natural dissimilarity measure between two probability distributions (Hoffman et al., 2013). This motivates the use of the natural gradient of $L(q)$ in the Euclidean space that can be equivalently considered its standard gradient in the redefined parameter space implementing the symmetrized KL distance (Amari, 1998). Such a natural gradient of $L(q)$ will be used to derive the *stochastic natural gradient ascent* (SNGA) method. Intuitively, SNGA can be regarded as another version of SGA that operates in a different parameter space defined with a different distance metric. Therefore, both converge to the same optimal parameters, although SNGA is empirically demonstrated to converge faster than SGA (Amari, 1998) when an objective function $L(q)$ is optimized with respect to a parameterized distribution $q(\mathbf{f}_U)$. This is expected since the symmetrized KL distance is more accurate than the Euclidean distance in measuring the dissimilarity between parameterized distributions.

To derive the natural gradient of $L(q)$, the moment parameterization of $q(\mathbf{f}_U) \triangleq \mathcal{N}(\mu, \Sigma)$ is first replaced by its canonical counterpart $q(\mathbf{f}_U | \theta)$:

$$q(\mathbf{f}_U | \theta) \triangleq \mathcal{N}(\mu, \Sigma) = h(\mathbf{f}_U) \exp(\theta^\top \mathbf{T}(\mathbf{f}_U) - A(\theta))$$

where $\mathbf{T}(\mathbf{f}_U) \triangleq (\mathbf{f}_U; \text{vec}(\mathbf{f}_U \mathbf{f}_U^\top))$, $h(\mathbf{f}_U) \triangleq (2\pi)^{-|\mathcal{U}|/2}$, $A(\theta)$ is simply a normalizing function guaranteeing that $q(\mathbf{f}_U | \theta)$ integrates to unity, and the natural parameters $\theta \triangleq (\theta_1; \text{vec}(\theta_2))$ where $\theta_1 = \Sigma^{-1}\mu$ and $\theta_2 = -(1/2)\Sigma^{-1}$. In particular, the metric distance defining the parameter space is given by the Riemannian metric tensor $H(\theta)$ (Amari, 1998) that corresponds to the identity matrix when the Euclidean metric is used. Otherwise, when the parameter space implements the symmetrized KL distance, the work of Hoffman et al. (2013) has shown that $H(\theta)$ is defined by the Fisher information matrix (Amari, 1998):

$$H(\theta) \triangleq -\mathbb{E}_{\mathbf{f}_U | \theta} \left[\frac{\partial^2 \log q(\mathbf{f}_U | \theta)}{\partial \theta \partial \theta^\top} \right] = \frac{\partial^2 A(\theta)}{\partial \theta \partial \theta^\top}. \quad (19)$$

The last equality is formally verified in Appendix E.2. Let $\partial L / \partial \theta$ be the standard gradient of $L(q)$ with respect to θ . Then, its natural gradient is defined as $\partial \bar{L} / \partial \theta \triangleq H(\theta)^{-1} \partial L / \partial \theta$. To express $\partial \bar{L} / \partial \theta$ in terms of μ and Σ , let $\eta \triangleq [\eta_1; \text{vec}(\eta_2)]$ where $\eta_1 \triangleq \mu$ and $\eta_2 \triangleq \mu \mu^\top + \Sigma$. It can be verified that $\mathbb{E}[\mathbf{T}(\mathbf{f}_U)] = \eta$ (Appendix E.1), which implies $\partial \eta / \partial \theta = H(\theta)$ (Appendix E.3). Using this result,

$$\frac{\partial \bar{L}}{\partial \theta} \triangleq H(\theta)^{-1} \frac{\partial L}{\partial \theta} = H(\theta)^{-1} \frac{\partial \eta}{\partial \theta} \frac{\partial L}{\partial \eta} = \frac{\partial L}{\partial \eta}. \quad (20)$$

The last equality is due to $\partial \eta / \partial \theta = H(\theta)$. So, the natural gradient can be evaluated by taking the derivative of $L(q)$ with respect to η (20). To simplify the derivation, the partial derivatives of $L(q)$ are taken with respect to η_1 and η_2 instead of differentiating it with η directly. To achieve this, $L(q)$ is first represented as a function of η_1 and η_2 :

$$\begin{aligned} L(q) &= \frac{1}{2} \log |\eta_2 - \eta_1 \eta_1^\top| + \eta_1^\top \left(\frac{1}{\sigma_n^2} P_{\mathcal{D}\mathcal{U}}^\top \mathcal{Y}_{\mathcal{D}} + \Lambda \nu \right) \\ &\quad - \frac{1}{2} \eta_1^\top \Psi \eta_1 - \frac{1}{2} \text{tr}(\Psi \eta_2 - \Psi \eta_1 \eta_1^\top) + C' \end{aligned}$$

which can be straightforwardly verified using (11) and η 's definition. The natural gradient of $L(q)$ is then given by

$$\frac{\partial L}{\partial \eta_1} = -(\eta_2 - \eta_1 \eta_1^\top)^{-1} \eta_1 + \frac{1}{\sigma_n^2} P_{\mathcal{D}\mathcal{U}}^\top \mathcal{Y}_{\mathcal{D}} + \Lambda \nu, \quad (21)$$

$$\frac{\partial L}{\partial \eta_2} = \frac{1}{2} \left((\eta_2 - \eta_1 \eta_1^\top)^{-1} - \Psi \right). \quad (22)$$

Finally, note that if (ν, Λ) is chosen to satisfy (12) and (13) to guarantee that $L(q)$ is maximized at $q(\mathbf{f}_{\mathcal{U}}) \equiv q^*(\mathbf{f}_{\mathcal{U}})$ (Theorem 4), then $\Psi = \Sigma^{*-1}$ and $(1/\sigma_n^2) P_{\mathcal{D}\mathcal{U}}^\top \mathcal{Y}_{\mathcal{D}} + \Lambda \nu = \Sigma^{*-1} \mu^*$. In addition, by definition, $\theta_1 = (\eta_2 - \eta_1 \eta_1^\top)^{-1} \eta_1$ and $\theta_2 = -(1/2)(\eta_2 - \eta_1 \eta_1^\top)^{-1}$. Hence, (21) and (22) can be rewritten as

$$\frac{\partial L}{\partial \eta_1} = \Sigma^{*-1} \mu^* - \theta_1 \text{ and } \frac{\partial L}{\partial \eta_2} = -\theta_2 - \frac{1}{2} \Sigma^{*-1}. \quad (23)$$

So, if (μ^*, Σ^*) satisfies the decomposability conditions (15) and (16), then it is possible to derive a stochastic natural gradient that is an unbiased estimator of the exact natural gradient (23), as formalized in the result below:

Theorem 6 *Let \mathcal{S} be a set of i.i.d. samples (i.e., $|\mathcal{S}| > 0$) drawn from a uniform distribution over $\{1, 2, \dots, P\}$ and*

$$\frac{\partial \hat{L}}{\partial \eta_1} \triangleq G'(\mathcal{U}) - \theta_1 + \frac{P}{|\mathcal{S}|} \sum_{s \in \mathcal{S}} G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s}), \quad (24)$$

$$\frac{\partial \hat{L}}{\partial \eta_2} \triangleq -\theta_2 - \frac{1}{2} F'(\mathcal{U}) - \frac{P}{2|\mathcal{S}|} \sum_{s \in \mathcal{S}} F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s}). \quad (25)$$

If (μ^, Σ^*) satisfies (15) and (16), then $\mathbb{E}[\partial \hat{L} / \partial \eta_1] = \partial L / \partial \eta_1$ and $\mathbb{E}[\partial \hat{L} / \partial \eta_2] = \partial L / \partial \eta_2$.*

Its proof is in Appendix B.5. The gradient ascent update in (14) can now be revised to

$$\theta_1^{t+1} = \theta_1^t + \rho_t \frac{\partial \hat{L}}{\partial \eta_1}(\theta_1^t, \theta_2^t), \quad \theta_2^{t+1} = \theta_2^t + \rho_t \frac{\partial \hat{L}}{\partial \eta_2}(\theta_1^t, \theta_2^t) \quad (26)$$

such that the parameters (μ, Σ) of $q(\mathbf{f}_{\mathcal{U}})$ can be recovered from its natural parameters θ by setting $(\mu^t, \Sigma^t) = (-(1/2)(\theta_2^t)^{-1} \theta_1^t, -(1/2)(\theta_2^t)^{-1})$. Therefore, if $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ is selected as that of DTC, then (26) recovers the SNGA method of Hensman et al. (2013) to produce an anytime variant of DTC, which is a special case spanned by our unifying framework of anytime SGPR models.

5. Experiments and Discussion

This section empirically evaluates the predictive performance and time efficiency of anytime SGPR models¹ such as the anytime variants of PIC and FITC, which we, respectively, call PIC+ and FITC+, and the state-of-the-art anytime variant of DTC (Hensman et al., 2013), which we name DTC+², spanned by our unifying framework on two real-world datasets of a few million in size:

(a) The *EMULATE mean sea level pressure* (EMSLP) dataset (Ansell et al., 2006) of size 1278250 spans a 5° lat.-lon. grid bounded within lat. 25-70N and lon. 70W-50E from 1900 to 2003. Each input denotes a 6-dimensional feature vector of latitude, longitude, year, month, day, and incremental day count (starting from 0 on first day). The output is the mean sea level pressure (Pa).

(b) The AIRLINE dataset contains 2055733 records of information about every commercial flight in the USA from January to April 2008. The input denotes a 8-dimensional feature vector of the age of the aircraft (i.e., no. of years in service), travel distance (km), airtime, departure and arrival time (min.) as well as day of the week, day of the month, and month. The output is the delay time (min.) of the flight.

Both datasets are modeled using GPs whose prior covariance is defined by the squared exponential covariance function $k_{xx'} \triangleq \sigma_s^2 \exp(-0.5(\mathbf{x} - \mathbf{x}')^\top \Upsilon^{-2}(\mathbf{x} - \mathbf{x}'))$ with a diagonal matrix Υ of d length-scale components and signal variance σ_s^2 being its defining hyperparameters. These hyperparameters together with the noise variance σ_n^2 are learned by generalizing the distributed, variational DTC-like learning framework of Gal et al. (2014) to account for the more relaxed structural assumptions of PIC and FITC. Such a generalization can then handle massive datasets by distributing the computational load of learning the hyperparameters of PIC and FITC among parallel computing nodes; its details are deferred to a separate paper since the focus of our work here is on scaling up the existing SGPR models while assuming that the hyperparameters are learned in advance. On a separate note, learning these hyperparameters in an anytime fashion is highly non-trivial and beyond the scope of this paper, which we intend to pursue in the future as a continuation of our current work.

For each dataset, 5% is randomly selected and set aside as test data \mathcal{S} . The remaining data (i.e., training data) is partitioned into P blocks using k -means (i.e., $k = P$). All experiments are run on a Linux system with Intel® Xeon® E5-2670 at 2.6GHz with 96 GB memory.

¹In the case of performing multiple predictions for single test inputs, the predictive means of FITC and DTC coincide with that of FIC and SoR, respectively.

²DTC+ coincides with stochastic variational inference for GPs in (Hensman et al., 2013), as discussed in Section 3.

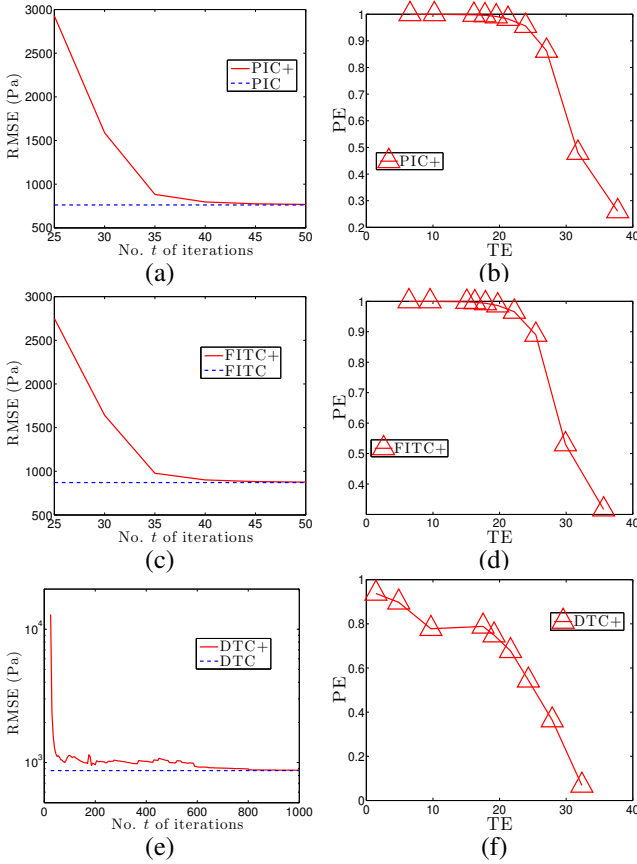


Figure 1. Graphs of RMSEs achieved by (a) PIC+, (c) FITC+, and (e) DTC+ vs. number t of iterations, and graphs of predictive efficiency (PE) vs. time efficiency (TE) showing the anytime efficiencies of (b) PIC+, (d) FITC+, and (f) DTC+ with $|\mathcal{U}| = 512$ inducing outputs and $P = 1000$ blocks for EMSLP dataset.

Four performance metrics are used to evaluate the anytime SGPR models: (a) *Root mean square error* (RMSE): $\sqrt{|\mathcal{S}|^{-1} \sum_{\mathbf{x} \in \mathcal{S}} (y_{\mathbf{x}} - \mu_{\mathbf{x}|\mathcal{D}})^2}$, (b) *mean negative log probability* (MNLPL): $0.5|\mathcal{S}|^{-1} \sum_{\mathbf{x} \in \mathcal{S}} ((y_{\mathbf{x}} - \mu_{\mathbf{x}|\mathcal{D}})^2 / \sigma_{\mathbf{xx}|\mathcal{D}} + \log(2\pi\sigma_{\mathbf{xx}|\mathcal{D}}))$, (c) incurred time, and (d) anytime efficiency demonstrating the trade-off between *time efficiency* (TE) vs. *predictive efficiency* (PE). Formally, TE (PE) is defined as the incurred time (RMSE) of the SGPR model divided by that of its anytime variant. Intuitively, increasing TE (i.e., by decreasing the number of iterations of SNGA update) reduces the incurred time of an anytime variant of the SGPR model at the cost of degrading its PE.

EMSLP Dataset. Figs. 1 and 2 show results of RMSEs, incurred times, and anytime efficiencies of PIC+, FITC+, and DTC+ averaged over 5 random instances with varying number t of iterations. It can be observed from Figs. 1a, 1c, and 1e that the RMSEs of PIC+, FITC+, and DTC+ consistently converge to within 0.75% of that of PIC (RMSE of 762.263 Pa), FITC (RMSE of 870.857 Pa), and DTC (RMSE of 870.878 Pa), respectively. The results of their MNLPLs show similar convergence behavior, as detailed in

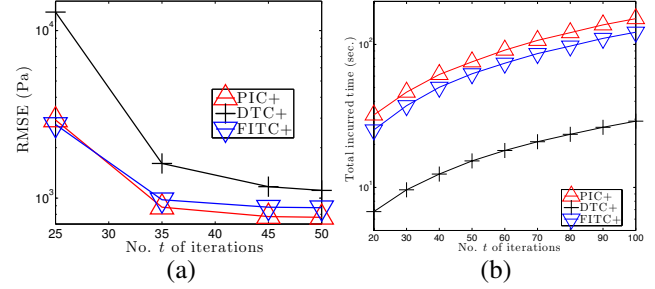


Figure 2. Graphs of (a) RMSEs and (b) total incurred times of PIC+, FITC+, and DTC+ vs. number t of iterations with $|\mathcal{U}| = 512$ inducing outputs and $P = 1000$ blocks for EMSLP dataset.

Appendix A. This corroborates our theoretical results in Section 4 that the anytime variants spanned by our unifying framework can achieve asymptotic convergence to the predictive distributions of their corresponding SGPR models. In particular, the RMSEs of PIC+ and FITC+ decrease quickly with an increasing number t of iterations of SNGA update and converge after 50 iterations, which demonstrate their scalability to massive datasets. In fact, during these first 50 iterations, the RMSEs achieved by PIC+ and FITC+ are significantly lower than that achieved by DTC+, as observed in Fig. 2a. On the other hand, the RMSE of DTC+ decreases more gradually and can only converge after 1000 iterations. This inferior predictive performance of DTC+ may be caused by its more restrictive structural assumption of deterministic relation between the training and inducing outputs (Appendix D.1), thus making it perform less robustly among heterogeneous datasets.

It can also be observed from Fig. 2a that the superior predictive performance (i.e., lower RMSE) of PIC+ over FITC+ becomes more pronounced with an increasing number t of iterations, which is expected: PIC+ imposes a more relaxed structural assumption of conditional independence than FITC+. For example, unlike FITC+, PIC+ does not assume conditional independence between the test and training outputs given the inducing outputs. Fig. 2b shows linear increases of total incurred time in the number t of iterations for PIC+, FITC+, and DTC+. Our experiments reveal that PIC+, FITC+, and DTC+ incur, respectively, an average of 1.53, 1.15, and 0.32 seconds per update iteration. So, PIC+ and FITC+ take ~ 76.5 and ~ 57.5 seconds to converge after 50 iterations while DTC+ takes ~ 320 seconds to converge after 1000 iterations.

Figs. 1b, 1d, and 1f reveal how the predictive efficiencies of the anytime SGPR models can be traded off to improve their time efficiencies to meet the real-time requirement in time-critical applications. It can be observed that both PIC+ and FITC+ can achieve a speedup of 22-24 (i.e., TE = 25) while preserving 96% of the predictive efficiencies of PIC and FITC (i.e., PE = 0.95); in other words, the RMSEs achieved by PIC+ and FITC+ are only $1/0.95 \approx 1.05$ times larger than that achieved by PIC and FITC. On the

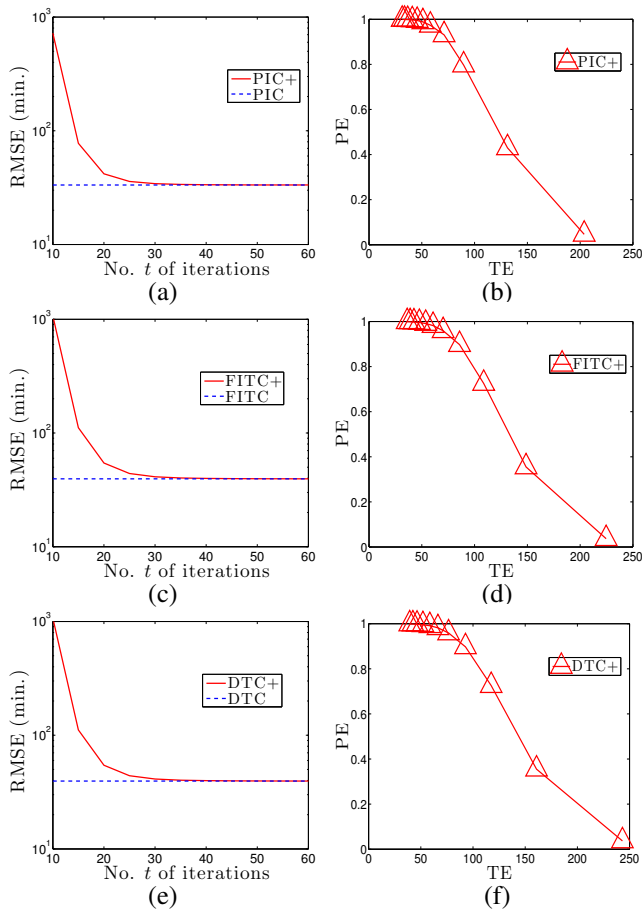


Figure 3. Graphs of RMSEs achieved by (a) PIC+, (c) FITC+, and (e) DTC+ vs. number t of iterations, and graphs of predictive efficiency (PE) vs. time efficiency (TE) showing the anytime efficiencies of (b) PIC+, (d) FITC+, and (f) DTC+ with $|\mathcal{U}| = 100$ inducing outputs and $P = 2000$ blocks for AIRLINE dataset.

other hand, with a speedup of 22, DTC+ can only reach 68% of the predictive efficiency of DTC.

AIRLINE Dataset. Figs. 3 and 4 show results of RMSEs, incurred times, and anytime efficiencies of PIC+, FITC+, and DTC+ averaged over 5 random instances with varying number t of iterations. The observations are mostly similar to that of the EMSLP dataset: From Figs. 3a, 3c, and 3e, the RMSEs of PIC+, FITC+, and DTC+ converge to within 0.04% of that of PIC (RMSE of 33.3515 min.), FITC (RMSE of 39.5302 min.), and DTC (RMSE of 39.5310 min.), respectively. The same observation can be made regarding the results of their MNLPs, as detailed in Appendix A. The RMSEs of PIC+, FITC+, and DTC+ decrease rapidly with an increasing number t of iterations and converge after 60 iterations. During these first 60 iterations, the RMSE achieved by PIC+ is much lower than that achieved by FITC+ and DTC+, as observed in Fig. 4a; this was previously explained in the discussion on the experimental results for EMSLP dataset. Fig. 4b shows linear increases of total incurred time in the number t of

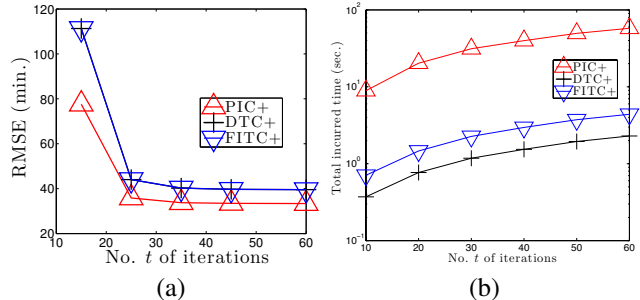


Figure 4. Graphs of (a) RMSEs and (b) total incurred times of PIC+, FITC+, and DTC+ vs. number t of iterations with $|\mathcal{U}| = 100$ and $P = 2000$ for the AIRLINE dataset.

iterations for PIC+, FITC+, and DTC+. Our experiments reveal that PIC+, FITC+, and DTC+ incur an average of 0.97, 0.08, and 0.04 seconds per iteration of SNGA update. So, it takes less than 1 minute for PIC+, FITC+, and DTC+ to converge after 60 iterations. Figs. 3b, 3d, and 3f reveal that PIC+, FITC+, and DTC+ can achieve a speedup of 50 (i.e., TE = 50) while preserving almost 100% of the predictive efficiencies of PIC, FITC, and DTC (i.e., PE = 1). But, as observed in Fig. 4a, PIC+ outperforms FITC+ and DTC+ by a huge margin; the same observation can be made for the EMSLP dataset (Fig. 2a). Hence, PIC+ offers the best predictive performance, anytime efficiency, and robustness in both EMSLP and AIRLINE datasets.

6. Conclusion and Future Work

This paper describes a novel unifying framework of anytime SGPR models (e.g., PIC+, FITC+, DTC+) that can produce good predictive performance fast and trade off between predictive performance vs. time efficiency. After applying our reverse variational inference procedure, a stochastic natural gradient ascent method can be derived that is guaranteed to achieve asymptotic convergence to the predictive distribution of any SGPR model of our choice. We prove that if the predictive distribution of the chosen SGPR model satisfies certain decomposability conditions, then the stochastic natural gradient is an unbiased estimator of the exact natural gradient and can be computed in constant time at each iteration. Empirical evaluation on two real-world million-sized datasets show that PIC+ outperforms FITC+ and state-of-the-art DTC+ (Hensman et al., 2013) in terms of predictive performance and anytime efficiency. A limitation of our unifying framework is that though it can produce the anytime variants of many existing SGPR models (Quiñero-Candela & Rasmussen, 2005), it does not cover some recent ones like (Lázaro-Gredilla et al., 2010; Low et al., 2015). So, in our future work, we will extend our framework to address this limitation as well as to learn the hyperparameters in an anytime fashion.

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A. MNLP Performances of Anytime SGPR Models for EMSLP and AIRLINE Datasets

Similar to Figs. 1 and 3 that show the convergence of the RMSEs of PIC+, FITC+, and DTC+ to that of PIC, FITC, and DTC, it can also be observed from Fig. 5 that the MNLPs of PIC+, FITC+, and DTC+ converge rapidly to that of PIC, FITC and DTC. Specifically, the MNLPs of PIC+, FITC+, and DTC+, respectively, converge to within 1.79% of that of PIC (MNLP of 8.3493), FITC (MNLP of 8.3973), and DTC (MNLP of 8.3973) for the EMSLP dataset. For the AIRLINE dataset, the MNLPs of PIC+, FITC+, and DTC+ converge to within 0.53% of that of PIC (MNLP of 5.6672), FITC (MNLP of 6.5211), and DTC (MNLP of 6.5212). This corroborates both our theoretical results in Section 4 and empirical analysis in Section 5.

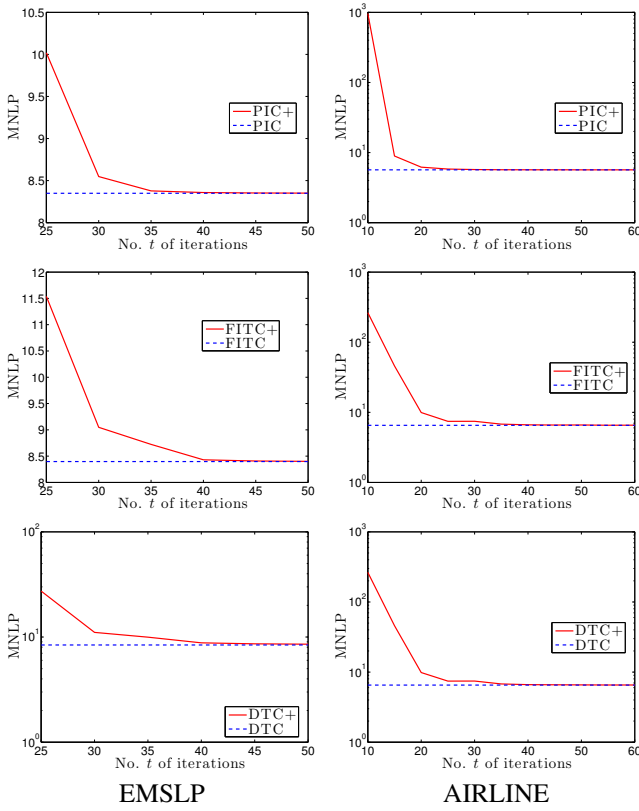


Figure 5. Graphs of MNLPs achieved by PIC+, FITC+, and DTC+ vs. number t of iterations with $|\mathcal{U}| = 512$ inducing outputs and $P = 1000$ blocks for EMSLP dataset (left) and with $|\mathcal{U}| = 100$ inducing outputs and $P = 2000$ blocks for AIRLINE dataset (right).

B. Proofs of Main Results

B.1. Proof of Theorem 2

Since $q(\mathbf{f}_D|\mathbf{f}_U) \triangleq p(\mathbf{f}_D|\mathbf{f}_U)$, (8) can be plugged into (7) to obtain

$$\begin{aligned} L_U(q) &= \int p(\mathbf{f}_D|\mathbf{f}_U) \log \frac{p(\mathbf{f}_D, \mathbf{y}_D|\mathbf{f}_U)}{p(\mathbf{f}_D|\mathbf{f}_U)} d\mathbf{f}_D \\ &= \int p(\mathbf{f}_D|\mathbf{f}_U) \log p(\mathbf{y}_D|\mathbf{f}_D) d\mathbf{f}_D. \end{aligned} \quad (27)$$

By applying $p(\mathbf{y}_D|\mathbf{f}_D) = \mathcal{N}(\mathbf{f}_D, \sigma_n^2 \mathbf{I})$ to (27),

$$\begin{aligned} L_U(q) &= \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U} \left[-\frac{1}{2\sigma_n^2} (\mathbf{y}_D - \mathbf{f}_D)^\top (\mathbf{y}_D - \mathbf{f}_D) \right] + C_1 \\ &= \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U} \left[-\frac{1}{2\sigma_n^2} \mathbf{f}_D^\top \mathbf{f}_D + \frac{1}{\sigma_n^2} \mathbf{f}_D^\top \mathbf{y}_D \right] + C_2 \end{aligned} \quad (28)$$

where C_1 absorbs all terms independent of \mathbf{f}_D and (28) is obtained by $C_2 \triangleq C_1 + (-1/2\sigma_n^2) \mathbf{y}_D^\top \mathbf{y}_D$. Since $p(\mathbf{f}_D|\mathbf{f}_U) = \mathcal{N}(P_{DU}\mathbf{f}_U, K_{DD} - Q_{DD})$ (8), the following Gaussian identities $\mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D] = P_{DU}\mathbf{f}_U$ and $\mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D^\top \mathbf{f}_D] = \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D]^\top \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D] + \text{tr}(K_{DD} - Q_{DD})$ can be applied to (28) to obtain

$$\begin{aligned} L_U(q) &= -\frac{1}{2\sigma_n^2} \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D^\top \mathbf{f}_D] + \frac{1}{\sigma_n^2} \mathbb{E}_{\mathbf{f}_D|\mathbf{f}_U}[\mathbf{f}_D]^\top \mathbf{y}_D + C_2 \\ &= -\frac{1}{2\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top P_{DU} \mathbf{f}_U + \frac{1}{\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top \mathbf{y}_D + C \end{aligned}$$

where $C \triangleq C_2 + (-1/2\sigma_n^2) \text{tr}(K_{DD} - Q_{DD})$.

B.2. Proof of Theorem 3

By plugging (9) into (6),

$$\begin{aligned} L(q) &= \int q(\mathbf{f}_U) L_U(q) d\mathbf{f}_U - D_{\text{KL}}(q(\mathbf{f}_U)||p(\mathbf{f}_U)) \\ &= \mathbb{E}_{\mathbf{f}_U} \left[-\frac{1}{2\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top P_{DU} \mathbf{f}_U + \frac{1}{\sigma_n^2} \mathbf{f}_U^\top P_{DU}^\top \mathbf{y}_D \right] - \\ &\quad D_{\text{KL}}(q(\mathbf{f}_U)||p(\mathbf{f}_U)) + C. \end{aligned} \quad (29)$$

Since $q(\mathbf{f}_U) \triangleq \mathcal{N}(\mu, \Sigma)$, the following Gaussian identities $\mathbb{E}_{\mathbf{f}_U}[\mathbf{f}_U] = \mu$ and $\mathbb{E}_{\mathbf{f}_U}[\mathbf{f}_U^\top P_{DU}^\top P_{DU} \mathbf{f}_U] = \mu^\top P_{DU}^\top P_{DU} \mu + \text{tr}(P_{DU}^\top P_{DU} \Sigma)$ can be applied to (29) to obtain

$$\begin{aligned} L(q) &= -\frac{1}{2\sigma_n^2} (\mu^\top P_{DU}^\top P_{DU} \mu + \text{tr}(P_{DU}^\top P_{DU} \Sigma)) + \\ &\quad \frac{1}{\sigma_n^2} \mu^\top P_{DU}^\top \mathbf{y}_D - D_{\text{KL}}(q(\mathbf{f}_U)||p(\mathbf{f}_U)) + C. \end{aligned} \quad (30)$$

We will now express $D_{\text{KL}}(q(\mathbf{f}_U)||p(\mathbf{f}_U))$ in (30) as a function of μ and Σ . Recall that $p(\mathbf{f}_U) \triangleq \mathcal{N}(\nu, \Lambda^{-1})$, which implies

$$-\log p(\mathbf{f}_U) = \frac{1}{2} (\mathbf{f}_U - \nu)^\top \Lambda (\mathbf{f}_U - \nu) + C_3 \quad (31)$$

where C_3 absorbs all terms independent of $\mathbf{f}_{\mathcal{U}}$. Then,

$$\begin{aligned}
 D_{\text{KL}}(q(\mathbf{f}_{\mathcal{U}})||p(\mathbf{f}_{\mathcal{U}})) &= \int q(\mathbf{f}_{\mathcal{U}}) \log \frac{q(\mathbf{f}_{\mathcal{U}})}{p(\mathbf{f}_{\mathcal{U}})} d\mathbf{f}_{\mathcal{U}} \\
 &= -\mathbb{H}[q(\mathbf{f}_{\mathcal{U}})] + \mathbb{E}_{\mathbf{f}_{\mathcal{U}}}[-\log p(\mathbf{f}_{\mathcal{U}})] \\
 &= \mathbb{E}_{\mathbf{f}_{\mathcal{U}}} \left[\frac{1}{2} (\mathbf{f}_{\mathcal{U}} - \nu)^\top \Lambda (\mathbf{f}_{\mathcal{U}} - \nu) \right] - \\
 &\quad \frac{1}{2} \log |\Sigma| + C_3 - C_4 \\
 &= \mathbb{E}_{\mathbf{f}_{\mathcal{U}}} \left[\frac{1}{2} \mathbf{f}_{\mathcal{U}}^\top \Lambda \mathbf{f}_{\mathcal{U}} - \nu^\top \Lambda \mathbf{f}_{\mathcal{U}} \right] - \\
 &\quad \frac{1}{2} \log |\Sigma| + C_5. \tag{32}
 \end{aligned}$$

The second equality is due to

$$\mathbb{H}[q(\mathbf{f}_{\mathcal{U}})] \triangleq - \int q(\mathbf{f}_{\mathcal{U}}) \log q(\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{U}}.$$

The third equality follows directly from (31) and $\mathbb{H}[q(\mathbf{f}_{\mathcal{U}})] = (1/2) \log |\Sigma| + C_4$ where C_4 absorbs all terms independent of μ and Σ . The last equality (32) is obtained by $C_5 \triangleq C_3 - C_4 + (1/2)\nu^\top \Lambda \nu$, which does not depend on μ and Σ . By applying the following Gaussian identities $\mathbb{E}_{\mathbf{f}_{\mathcal{U}}}[\mathbf{f}_{\mathcal{U}}] = \mu$ and $\mathbb{E}_{\mathbf{f}_{\mathcal{U}}}[\mathbf{f}_{\mathcal{U}}^\top \Lambda \mathbf{f}_{\mathcal{U}}] = \mu^\top \Lambda \mu + \text{tr}(\Lambda \Sigma)$ to (32),

$$\begin{aligned}
 D_{\text{KL}}(q(\mathbf{f}_{\mathcal{U}})||p(\mathbf{f}_{\mathcal{U}})) &= \frac{1}{2} \mu^\top \Lambda \mu + \frac{1}{2} \text{tr}(\Lambda \Sigma) - \mu^\top \Lambda \nu - \\
 &\quad \frac{1}{2} \log |\Sigma| + C_5. \tag{33}
 \end{aligned}$$

Plugging (33) into (30) gives (11) where $C' \triangleq C - C_5$ is independent of μ and Σ .

B.3. Proof of Theorem 4

By plugging (13) into the definition of $\Psi \triangleq (1/\sigma_n^2) P_{\mathcal{D}\mathcal{U}}^\top P_{\mathcal{D}\mathcal{U}} + \Lambda$ in Theorem 3, $\Psi = \Sigma^{*-1}$. By substituting $\Psi = \Sigma^{*-1}$ and (12) into (11) in Theorem 3,

$$\begin{aligned}
 L(q) &= -\frac{1}{2} \mu^\top \Sigma^{*-1} \mu - \frac{1}{2} \text{tr}(\Sigma^{*-1} \Sigma) + \\
 &\quad \frac{1}{2} \log |\Sigma| + \mu^\top \Sigma^{*-1} \mu^* + C'. \tag{34}
 \end{aligned}$$

Differentiating both sides of (34) with respect to μ and Σ gives

$$\frac{\partial L}{\partial \mu} = -\Sigma^{*-1} \mu + \Sigma^{*-1} \mu^*, \tag{35}$$

$$\frac{\partial L}{\partial \Sigma} = -\frac{1}{2} \Sigma^{*-1} + \frac{1}{2} \Sigma^{-1}. \tag{36}$$

By setting $\partial L / \partial \mu = 0$ and $\partial L / \partial \Sigma = 0$, $(\mu, \Sigma) = (\mu^*, \Sigma^*)$. In addition, since $L(q)$ (34) is concave in both μ and Σ and $\partial L / \partial \mu$ and $\partial L / \partial \Sigma$ are, respectively, independent of Σ and μ , it is clear that $L(q)$ is maximized when its derivatives go to zero.

B.4. Proof of Theorem 5

Since \mathcal{S} be a set of i.i.d. samples (i.e., $|\mathcal{S}| > 0$) drawn from a uniform distribution over $\{1, 2, \dots, P\}$,

$$\begin{aligned}
 \mathbb{E}[F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s})] &= \sum_{i=1}^P \Pr(s = i) F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \\
 &= \sum_{i=1}^P \frac{1}{P} F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \\
 &= \frac{1}{P} \sum_{i=1}^P F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}). \tag{37}
 \end{aligned}$$

Applying a similar argument to $G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s})$,

$$\mathbb{E}[G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s})] = \frac{1}{P} \sum_{i=1}^P G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}). \tag{38}$$

From (37) and (38),

$$\begin{aligned}
 &\mathbb{E}[G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s}) - F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s}) \mu] \\
 &= \mathbb{E}[G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s})] - \mathbb{E}[F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_s}) \mu] \\
 &= \frac{1}{P} \sum_{i=1}^P G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) - F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \mu. \tag{39}
 \end{aligned}$$

By taking the expectation on both sides of (17) and applying (39),

$$\begin{aligned}
 \mathbb{E} \left[\frac{\partial \widehat{L}}{\partial \mu} \right] &= G'(\mathcal{U}) + \sum_{i=1}^P G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) - \\
 &\quad \left(F'(\mathcal{U}) + \sum_{i=1}^P F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \right) \mu. \tag{40}
 \end{aligned}$$

By plugging (15) and (16) into the RHS of (40),

$$\mathbb{E} \left[\frac{\partial \widehat{L}}{\partial \mu} \right] = \Sigma^{*-1} \mu^* - \Sigma^{*-1} \mu = \frac{\partial L}{\partial \mu} \tag{41}$$

where the last equality follows from (35). Similarly, by taking the expectation on both sides of (18) and applying (37),

$$\begin{aligned}
 \mathbb{E} \left[\frac{\partial \widehat{L}}{\partial \Sigma} \right] &= \frac{1}{2} \left(\Sigma^{-1} - \left(F'(\mathcal{U}) + \sum_{i=1}^P F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \right) \right) \\
 &= \frac{1}{2} (\Sigma^{-1} - \Sigma^{*-1}) = \frac{\partial L}{\partial \Sigma} \tag{42}
 \end{aligned}$$

where the last two equalities follow directly from (15) and (36), respectively.

B.5. Proof of Theorem 6

This proof can be constructed by reiterating the exact arguments used in the proof of Theorem 5 (Appendix B.4). By taking the expectation on both sides of (24) and applying (38),

$$\mathbb{E} \left[\frac{\partial \widehat{L}}{\partial \eta_1} \right] = -(\eta_2 - \eta_1 \eta_1^\top)^{-1} \eta_1 + G'(\mathcal{U}) + \sum_{i=1}^P G(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}). \quad (43)$$

Then, plugging (16) into the RHS of (43) recovers (23), which implies $\mathbb{E}[\partial \widehat{L}/\partial \eta_1] = \partial L/\partial \eta_1$. Similarly, by taking the expectation on both sides of (25) and applying (37),

$$\mathbb{E} \left[\frac{\partial \widehat{L}}{\partial \eta_2} \right] = \frac{1}{2} (\eta_2 - \eta_1 \eta_1^\top)^{-1} - \frac{1}{2} \left(F'(\mathcal{U}) + \sum_{i=1}^P F(\mathcal{U}, \mathbf{y}_{\mathcal{D}_i}) \right). \quad (44)$$

Finally, plugging (15) into the RHS of (44) recovers (23), which implies $\mathbb{E}[\partial \widehat{L}/\partial \eta_2] = \partial L/\partial \eta_2$.

C. Proofs of Auxiliary Results

C.1. Derivation of (2)

Marginalizing out $f_{\mathbf{x}}$ from both sides of (1),

$$p(\mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) = \prod_{i=1}^P p(\mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}). \quad (45)$$

Using law of total probability,

$$p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}) = \int p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}}) d\mathbf{f}_{\mathcal{U}}. \quad (46)$$

So, to derive (2), it suffices to prove that $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) = p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$. To do this, note that

$$p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) = \frac{p(f_{\mathbf{x}}, \mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}})}{p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}})}. \quad (47)$$

To simplify (47), its denominator is first factorized:

$$\begin{aligned} p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) &= \int p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{D}}) p(\mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}} \\ &= \int \prod_{i=1}^P p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{D}_i}) p(\mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}} \\ &= \prod_{i=1}^P \int p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{D}_i}) p(\mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}_i} \\ &= \prod_{i=1}^P p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) \end{aligned} \quad (48)$$

such that the second equality follows from (45) and the factorization of $p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{D}}) = \prod_{i=1}^P p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{D}_i})$. Likewise, the numerator on the RHS of (47) can be factorized:

$$\begin{aligned} &p(f_{\mathbf{x}}, \mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) \\ &= \int p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{D}}) p(f_{\mathbf{x}}, \mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}} \\ &= \int p(\mathbf{y}_{\mathcal{D}_P}|\mathbf{f}_{\mathcal{D}_P}) p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) p(\mathbf{f}_{\mathcal{D}_P}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}_P} \times \\ &\quad \prod_{i=1}^{P-1} \int p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{D}_i}) p(\mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}_i} \\ &= \int p(f_{\mathbf{x}}, \mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{D}_P}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}_P} \times \\ &\quad \prod_{i=1}^{P-1} \int p(\mathbf{y}_{\mathcal{D}_i}, \mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) d\mathbf{f}_{\mathcal{D}_i} \\ &= p(f_{\mathbf{x}}, \mathbf{y}_{\mathcal{D}_P}|\mathbf{f}_{\mathcal{U}}) \prod_{i=1}^{P-1} p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) \\ &= p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \prod_{i=1}^P p(\mathbf{y}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) \end{aligned} \quad (49)$$

such that the second equality follows from the above factorization of $p(\mathbf{y}_{\mathcal{D}}|\mathbf{f}_{\mathcal{D}})$ and (1). Thus, plugging (48) and (49) into (47) yields

$$p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) = p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}). \quad (50)$$

Plugging (50) into (46) concludes our derivation of (2).

C.2. Proof of Lemma 1

For all $\mathbf{f}_{\mathcal{D}}$ and $\mathbf{f}_{\mathcal{U}}$, $p(\mathbf{y}_{\mathcal{D}}) = p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})/p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$, which directly implies

$$\log p(\mathbf{y}_{\mathcal{D}}) = \log \frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})}{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})}. \quad (51)$$

Let $q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})$ be an arbitrary *probability density function* (PDF). Integrating both sides of (51) with q gives

$$\log p(\mathbf{y}_{\mathcal{D}}) = \int q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \log \frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})}{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})} d\mathbf{f}_{\mathcal{D}} d\mathbf{f}_{\mathcal{U}}. \quad (52)$$

Using the identity $\log(ab) = \log a + \log b$,

$$\log \frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})}{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})} = \log \frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})}{q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})} + \log \frac{q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})}{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})}.$$

Plugging this into (52),

$$\begin{aligned} \log p(\mathbf{y}_{\mathcal{D}}) &= \int q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \log \frac{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}})}{q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})} d\mathbf{f}_{\mathcal{D}} d\mathbf{f}_{\mathcal{U}} + \\ &\quad \int q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}) \log \frac{q(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}})}{p(\mathbf{f}_{\mathcal{D}}, \mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})} d\mathbf{f}_{\mathcal{D}} d\mathbf{f}_{\mathcal{U}}. \end{aligned} \quad (53)$$

Using the definition of $L(q)$ and $D_{\text{KL}}(\cdot|\cdot)$, Lemma 1 follows directly from (53).

C.3. Proof of Theorem 1

Using $q(\mathbf{f}_D, \mathbf{f}_U) = q(\mathbf{f}_D|\mathbf{f}_U) q(\mathbf{f}_U)$, $L(q)$ (5) can be rewritten as

$$L(q) = \int_{\mathbf{f}_U} q(\mathbf{f}_U) \left[\int_{\mathbf{f}_D} q(\mathbf{f}_D|\mathbf{f}_U) \log \frac{p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D)}{q(\mathbf{f}_D, \mathbf{f}_U)} d\mathbf{f}_D \right] d\mathbf{f}_U. \quad (54)$$

Since $p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D) = p(\mathbf{f}_D, \mathbf{y}_D|\mathbf{f}_U) p(\mathbf{f}_U)$ and $q(\mathbf{f}_D, \mathbf{f}_U) = q(\mathbf{f}_D|\mathbf{f}_U) q(\mathbf{f}_U)$,

$$\log \frac{p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D)}{q(\mathbf{f}_D, \mathbf{f}_U)} = \log \frac{p(\mathbf{f}_D, \mathbf{y}_D|\mathbf{f}_U)}{q(\mathbf{f}_D|\mathbf{f}_U)} - \log \frac{q(\mathbf{f}_U)}{p(\mathbf{f}_U)}. \quad (55)$$

Thus, in order to simplify (54), we first use (55) to simplify

$$\begin{aligned} & \int_{\mathbf{f}_D} q(\mathbf{f}_D|\mathbf{f}_U) \log \frac{p(\mathbf{f}_D, \mathbf{f}_U, \mathbf{y}_D)}{q(\mathbf{f}_D, \mathbf{f}_U)} d\mathbf{f}_D \\ &= \int_{\mathbf{f}_D} q(\mathbf{f}_D|\mathbf{f}_U) \log \frac{p(\mathbf{f}_D, \mathbf{y}_D|\mathbf{f}_U)}{q(\mathbf{f}_D|\mathbf{f}_U)} d\mathbf{f}_D - \log \frac{q(\mathbf{f}_U)}{p(\mathbf{f}_U)} \\ &= L_U(q) - \log \frac{q(\mathbf{f}_U)}{p(\mathbf{f}_U)}. \end{aligned} \quad (56)$$

Plugging (56) into (54),

$$\begin{aligned} L(q) &= \int_{\mathbf{f}_U} q(\mathbf{f}_U) \left(L_U(q) - \log \frac{q(\mathbf{f}_U)}{p(\mathbf{f}_U)} \right) d\mathbf{f}_U \\ &= \int_{\mathbf{f}_U} q(\mathbf{f}_U) L_U(q) d\mathbf{f}_U - D_{\text{KL}}(q(\mathbf{f}_U) \| p(\mathbf{f}_U)). \end{aligned} \quad (57)$$

D. SGPR Models Satisfying Decomposability Conditions

This section demonstrates how the approximation $q^*(\mathbf{f}_U) \approx p(\mathbf{f}_U|\mathbf{y}_D)$ (i.e., (μ^*, Σ^*)) employed by SoR, DTC, FITC, FIC, PITC, and PIC can satisfy the decomposability conditions listed in (15) and (16). For clarity, we begin by deriving $q^*(\mathbf{f}_U)$ of these SGPR models in Appendix D.1.

D.1. Characterizing SGPR Models using (3)

D.1.1. PIC

The *partially independent conditional* (PIC) approximation (Snelson & Ghahramani, 2007) jointly specifies its test and training conditionals via the factorization in (1), which reduces the GP predictive distribution to (2), as shown in Appendix C.1. Since (1) is the only assumption made in PIC, its approximations $q^*(f_{\mathbf{x}}|\mathbf{y}_{D_P}, \mathbf{f}_U)$ and $q^*(\mathbf{f}_U)$ in (3) coincide with $p(f_{\mathbf{x}}|\mathbf{y}_{D_P}, \mathbf{f}_U)$ and $p(\mathbf{f}_U|\mathbf{y}_D)$, respectively. We will now derive $q^*(\mathbf{f}_U) \triangleq p(\mathbf{f}_U|\mathbf{y}_D)$ and prove that computing $q^*(\mathbf{f}_U)$ incurs $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time. Since $p(\mathbf{f}_{D_i}|\mathbf{f}_U)$ is the exact training conditional for the subset \mathcal{D}_i of training inputs,

$$p(\mathbf{f}_{D_i}|\mathbf{f}_U) = \mathcal{N}(K_{\mathcal{D}_i U} K_{U U}^{-1} \mathbf{f}_U, K_{\mathcal{D}_i \mathcal{D}_i} - Q_{\mathcal{D}_i \mathcal{D}_i}) \quad (58)$$

where

$$Q_{\mathcal{D}_i \mathcal{D}_i} \triangleq K_{\mathcal{D}_i U} K_{U U}^{-1} K_{U \mathcal{D}_i},$$

$K_{\mathcal{D}_i \mathcal{D}_i} \triangleq (k_{\mathbf{x}' \mathbf{x}''})_{\mathbf{x}', \mathbf{x}'' \in \mathcal{D}_i}$, $K_{U U} \triangleq (k_{\mathbf{x}' \mathbf{x}''})_{\mathbf{x}', \mathbf{x}'' \in \mathcal{U}}$, $K_{\mathcal{D}_i U} \triangleq (k_{\mathbf{x}' \mathbf{x}''})_{\mathbf{x}' \in \mathcal{D}_i, \mathbf{x}'' \in \mathcal{U}}$, and $K_{U \mathcal{D}_i} \triangleq K_{\mathcal{D}_i U}^\top$. Using (45) and (58), the training conditional $p(\mathbf{f}_D|\mathbf{f}_U)$ can be more compactly written as

$$p(\mathbf{f}_D|\mathbf{f}_U) = \mathcal{N}(K_{\mathcal{D} U} K_{U U}^{-1} \mathbf{f}_U, R) \quad (59)$$

where R is a block-diagonal matrix constructed from the P diagonal blocks of $K_{\mathcal{D} \mathcal{D}} - Q_{\mathcal{D} \mathcal{D}}$, each of which is $K_{\mathcal{D}_i \mathcal{D}_i} - Q_{\mathcal{D}_i \mathcal{D}_i}$ for $i = 1, \dots, P$. By combining with $p(\mathbf{f}_U) = \mathcal{N}(\mathbf{0}, K_{U U})$, the joint prior $p(\mathbf{f}_U, \mathbf{f}_D)$ can be expressed as

$$\mathcal{N}\left(\mathbf{0}, \begin{pmatrix} K_{U U} & K_{U \mathcal{D}} \\ K_{\mathcal{D} U} & Q_{\mathcal{D} \mathcal{D}} + R \end{pmatrix}\right).$$

Then, since $p(\mathbf{y}_D|\mathbf{f}_D) = \mathcal{N}(\mathbf{f}_D, \sigma_n^2 I)$, the joint prior $p(\mathbf{f}_U, \mathbf{y}_D)$ (Rasmussen & Williams, 2006) can be derived to be

$$\mathcal{N}\left(\mathbf{0}, \begin{pmatrix} K_{U U} & K_{U \mathcal{D}} \\ K_{\mathcal{D} U} & Q_{\mathcal{D} \mathcal{D}} + \Gamma \end{pmatrix}\right) \quad (60)$$

where $\Gamma \triangleq R + \sigma_n^2 I$. Using (60), $q^*(\mathbf{f}_U) = p(\mathbf{f}_U|\mathbf{y}_D) = \mathcal{N}(\mu^*, \Sigma^*)$ where

$$\mu^* \triangleq K_{U \mathcal{D}} (Q_{\mathcal{D} \mathcal{D}} + \Gamma)^{-1} \mathbf{y}_D, \quad (61)$$

$$\Sigma^* \triangleq K_{U U} - K_{U \mathcal{D}} (Q_{\mathcal{D} \mathcal{D}} + \Gamma)^{-1} K_{\mathcal{D} U}. \quad (62)$$

Note that

$$(Q_{\mathcal{D} \mathcal{D}} + \Gamma)^{-1} = \Gamma^{-1} - \Gamma^{-1} K_{\mathcal{D} U} \Phi K_{U \mathcal{D}} \Gamma^{-1} \quad (63)$$

follows directly from the matrix inversion lemma where $\Phi \triangleq (K_{U U} + K_{U \mathcal{D}} \Gamma^{-1} K_{\mathcal{D} U})^{-1}$. Using (63),

$$\begin{aligned} \mu^* &= K_{U \mathcal{D}} (Q_{\mathcal{D} \mathcal{D}} + \Gamma)^{-1} \mathbf{y}_D \\ &= K_{U \mathcal{D}} (\Gamma^{-1} - \Gamma^{-1} K_{\mathcal{D} U} \Phi K_{U \mathcal{D}} \Gamma^{-1}) \mathbf{y}_D \\ &= (\Phi^{-1} - K_{U \mathcal{D}} \Gamma^{-1} K_{\mathcal{D} U}) \Phi K_{U \mathcal{D}} \Gamma^{-1} \mathbf{y}_D \\ &= K_{U U} \Phi K_{U \mathcal{D}} \Gamma^{-1} \mathbf{y}_D \end{aligned} \quad (64)$$

where the last equality follows from the definition of Φ . Likewise,

$$\begin{aligned} \Sigma^* &= K_{U U} - K_{U \mathcal{D}} (Q_{\mathcal{D} \mathcal{D}} + \Gamma)^{-1} K_{\mathcal{D} U} \\ &= K_{U U} - K_{U U} \Phi K_{U \mathcal{D}} \Gamma^{-1} K_{\mathcal{D} U} \\ &= K_{U U} \Phi (\Phi^{-1} - K_{U \mathcal{D}} \Gamma^{-1} K_{\mathcal{D} U}) \\ &= K_{U U} \Phi K_{U U} \end{aligned} \quad (65)$$

where the last equality again follows from the definition of Φ . As such, μ^* (64) and Σ^* (65) can be computed in $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time instead of $\mathcal{O}(|\mathcal{D}|^3)$ time incurred by (61)

and (62). To see this, we can set $P = \mathcal{O}(|\mathcal{D}|/|\mathcal{U}|)$ and $|\mathcal{D}_i| = \mathcal{O}(|\mathcal{U}|)$ for $i = 1, \dots, P$ to guarantee that computing

$$K_{\mathcal{U}\mathcal{D}}\Gamma^{-1}\mathbf{y}_{\mathcal{D}} = \sum_{i=1}^P K_{\mathcal{U}\mathcal{D}_i}(K_{\mathcal{D}_i\mathcal{D}_i} - Q_{\mathcal{D}_i\mathcal{D}_i} + \sigma^2 I)^{-1}\mathbf{y}_{\mathcal{D}_i} \quad (66)$$

and

$$K_{\mathcal{U}\mathcal{D}}\Gamma^{-1}K_{\mathcal{D}\mathcal{U}} = \sum_{i=1}^P K_{\mathcal{U}\mathcal{D}_i}(K_{\mathcal{D}_i\mathcal{D}_i} - Q_{\mathcal{D}_i\mathcal{D}_i} + \sigma^2 I)^{-1}K_{\mathcal{D}_i\mathcal{U}} \quad (67)$$

incur $\mathcal{O}(|\mathcal{U}|^3 + \sum_{i=1}^P(|\mathcal{D}_i|^3 + |\mathcal{D}_i|^2|\mathcal{U}| + |\mathcal{D}_i||\mathcal{U}|^2)) = \mathcal{O}(|\mathcal{U}|^3 + (|\mathcal{D}|/|\mathcal{U}|)|\mathcal{U}|^3) = \mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time since $|\mathcal{U}| \ll |\mathcal{D}|$. It follows that Φ can be computed in $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2 + |\mathcal{U}|^3) = \mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time. Computing $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*) = p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ (i.e., (64) and (65)) therefore incurs $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time.

D.1.2. PITC, FITC, AND FIC

The only difference between *partially independent training conditional* (PITC) approximation (Quiñero-Candela & Rasmussen, 2005) and PIC (Snelson & Ghahramani, 2007) is that the former has additionally assumed conditional independence of $f_{\mathbf{x}}$ and $\mathbf{f}_{\mathcal{D}_P}$ given $\mathbf{f}_{\mathcal{U}}$, thus resulting in $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \triangleq p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) = p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ in (2). However, the training conditional $p(\mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}})$ remains the same as that of PIC (59), which implies that PIC and PITC share the same approximation $q^*(\mathbf{f}_{\mathcal{U}}) = p(\mathbf{f}_{\mathcal{U}}|\mathbf{y}_{\mathcal{D}})$ as it is derived independently from the test conditional $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$. The time needed to compute $q^*(\mathbf{f}_{\mathcal{U}})$ of PITC is therefore the same as that of PIC.

On the other hand, *fully independent training conditional* (FITC) approximation (Snelson & Ghahramani, 2005) is a special case of PITC (Quiñero-Candela & Rasmussen, 2005) whereby R in (59) is now a diagonal matrix constructed from the diagonal of $K_{\mathcal{D}\mathcal{D}} - Q_{\mathcal{D}\mathcal{D}}$. As a result, computing $K_{\mathcal{U}\mathcal{D}}\Gamma^{-1}\mathbf{y}_{\mathcal{D}}$ (66) and $K_{\mathcal{U}\mathcal{D}}\Gamma^{-1}K_{\mathcal{D}\mathcal{U}}$ (67) incur $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time. It follows that Φ can be computed in $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time. Computing $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ (i.e., (64) and (65)) of FITC therefore incurs $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time.

Finally, *fully independent conditional* (FIC) approximation differs from FITC by assuming an additional factorization across all the latent outputs for the test inputs, which only affects the test conditional. The training conditional of FIC remains the same as that of FITC, which implies that they share the same approximation $q^*(\mathbf{f}_{\mathcal{U}})$ as it is derived independently from the test conditional. The time needed to compute $q^*(\mathbf{f}_{\mathcal{U}})$ of FIC is therefore the same as that of FITC.

D.1.3. DTC AND SoR

The only difference between *deterministic training conditional* (DTC) approximation (Seeger et al., 2003) and PITC (Quiñero-Candela & Rasmussen, 2005) is that the former replaces $p(\mathbf{f}_{\mathcal{D}_i}|\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(K_{\mathcal{D}_i\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, K_{\mathcal{D}_i\mathcal{D}_i} - Q_{\mathcal{D}_i\mathcal{D}_i})$ (58) with $\mathcal{N}(K_{\mathcal{D}_i\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, \mathbf{0})$. This effectively replaces the training conditional $p(\mathbf{f}_{\mathcal{D}}|\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(K_{\mathcal{D}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, R)$ (59) with $\mathcal{N}(K_{\mathcal{D}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, \mathbf{0})$. So, R in the definition of Γ (60) is replaced by $\mathbf{0}$. The resulting Γ is plugged into the expression of Φ (63). Finally, the resulting Φ is plugged into (64) and (65) to derive $q^*(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu^*, \Sigma^*)$ of DTC where

$$\mu^* = \frac{1}{\sigma_n^2}K_{\mathcal{U}\mathcal{U}} \left(K_{\mathcal{U}\mathcal{U}} + \frac{1}{\sigma_n^2}K_{\mathcal{U}\mathcal{D}}K_{\mathcal{D}\mathcal{U}} \right)^{-1} K_{\mathcal{U}\mathcal{D}}\mathbf{y}_{\mathcal{D}} \quad (68)$$

$$\Sigma^* = K_{\mathcal{U}\mathcal{U}} \left(K_{\mathcal{U}\mathcal{U}} + \frac{1}{\sigma_n^2}K_{\mathcal{U}\mathcal{D}}K_{\mathcal{D}\mathcal{U}} \right)^{-1} K_{\mathcal{U}\mathcal{U}}. \quad (69)$$

Computing $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ (i.e., (68) and (69)) of DTC therefore incurs $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time.

Lastly, *subset of regressors* (SoR) (Smola & Bartlett, 2001) only differs from DTC (Seeger et al., 2003) by replacing its exact test conditional $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(K_{\mathbf{x}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, k_{\mathbf{x}\mathbf{x}} - Q_{\mathbf{x}\mathbf{x}})$ with $\mathcal{N}(K_{\mathbf{x}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}\mathbf{f}_{\mathcal{U}}, \mathbf{0})$ (Quiñero-Candela & Rasmussen, 2005). However, this change in test conditional does not affect the training conditional of DTC, which solely determines $q^*(\mathbf{f}_{\mathcal{U}})$, as explained previously in Appendix D.1.2. As a result, DTC and SoR share the same $q^*(\mathbf{f}_{\mathcal{U}})$. Computing $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ (i.e., (68) and (69)) of SoR therefore incurs $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time.

D.2. SGPR Models Satisfying Decomposability Conditions

This section demonstrates how $q^*(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu^*, \Sigma^*)$ induced by SoR, DTC, FITC, FIC, PITC, and PIC in Section D.1 can satisfy the decomposability conditions in (15) and (16). To simplify the analysis here, let us first recall that (a) PIC and PITC share the same $q^*(\mathbf{f}_{\mathcal{U}})$ (Appendix D.1.2), (b) SoR and DTC also induce the same $q^*(\mathbf{f}_{\mathcal{U}})$, and (c) FITC and FIC are special cases of PITC and PIC, respectively. Thus, it suffices to show that $q^*(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu^*, \Sigma^*)$ induced by DTC (Appendix D.2.1) and PIC (Appendix D.2.2) satisfy the decomposability conditions in (15) and (16) here.

D.2.1. DTC SATISFYING DECOMPOSABILITY CONDITIONS

According to Appendix D.1.3, DTC's induced approximation $q^*(\mathbf{f}_U) = \mathcal{N}(\mathbf{f}_U|\mu^*, \Sigma^*)$ of $p(\mathbf{f}_U|\mathbf{y}_D)$ is given by

$$\mu^* = \frac{1}{\sigma_n^2} K_{UU} \Phi K_{UD} \mathbf{y}_D, \quad (70)$$

$$\Sigma^* = K_{UU} \Phi K_{UU}, \quad (71)$$

where $\Phi \triangleq (K_{UU} + (1/\sigma_n^2) K_{UD} K_{DU})^{-1}$. Then,

$$\begin{aligned} \Sigma^{*-1} &= K_{UU}^{-1} \Phi^{-1} K_{UU}^{-1} \\ &= K_{UU}^{-1} \left(\frac{1}{\sigma_n^2} K_{UD} K_{DU} + K_{UU} \right) K_{UU}^{-1} \\ &= \frac{1}{\sigma_n^2} K_{UU}^{-1} K_{UD} K_{DU} K_{UU}^{-1} + K_{UU}^{-1} \end{aligned} \quad (72)$$

where the first equality follows directly from (71). Now,

$$K_{UD} K_{DU} = \sum_{i=1}^P K_{UD_i} K_{D_i U}. \quad (73)$$

By plugging (73) into (72),

$$\Sigma^{*-1} = \sum_{i=1}^P \left(\frac{1}{\sigma_n^2} K_{UU}^{-1} K_{UD_i} K_{D_i U} K_{UU}^{-1} \right) + K_{UU}^{-1}. \quad (74)$$

Also, by plugging (70) and (71) into $\Sigma^{*-1} \mu^*$ yields

$$\begin{aligned} \Sigma^{*-1} \mu^* &= \frac{1}{\sigma_n^2} K_{UU}^{-1} K_{UD} \mathbf{y}_D \\ &= \sum_{i=1}^P \frac{1}{\sigma_n^2} K_{UU}^{-1} K_{UD_i} \mathbf{y}_{D_i}. \end{aligned} \quad (75)$$

Finally, by letting $F'(\mathcal{U}) \triangleq K_{UU}^{-1}$ and $F(\mathcal{U}, \mathbf{y}_{D_i}) \triangleq (1/\sigma_n^2) K_{UU}^{-1} K_{UD_i} K_{D_i U} K_{UU}^{-1}$ in (74) and $G(\mathcal{U}, \mathbf{y}_{D_i}) \triangleq (1/\sigma_n^2) K_{UU}^{-1} K_{UD_i} \mathbf{y}_{D_i}$ and $G'(\mathcal{U}) = \mathbf{0}$ in (75), (15) and (16) result.

D.2.2. PIC SATISFYING DECOMPOSABILITY CONDITIONS

According to Appendix D.1.1, PIC (Snelson & Ghahramani, 2007) induces the following approximation $q^*(\mathbf{f}_U) = \mathcal{N}(\mathbf{f}_U|\mu^*, \Sigma^*)$ of $p(\mathbf{f}_U|\mathbf{y}_D)$:

$$\mu^* = K_{UU} \Phi K_{UD} \Gamma^{-1} \mathbf{y}_D, \quad (76)$$

$$\Sigma^* = K_{UU} \Phi K_{UU} \quad (77)$$

where $\Phi = (K_{UU} + K_{UD} \Gamma^{-1} K_{DU})^{-1}$ and $\Gamma = R + \sigma_n^2 I$. Using (67),

$$\begin{aligned} \Phi^{-1} &= K_{UU} + K_{UD} \Gamma^{-1} K_{DU} \\ &= K_{UU} + \sum_{i=1}^P K_{UD_i} \Gamma_i^{-1} K_{D_i U} \end{aligned} \quad (78)$$

where $\Gamma_i \triangleq K_{D_i D_i} - Q_{D_i D_i} + \sigma^2 I$. Then, plugging (78) into (77) yields

$$\begin{aligned} \Sigma^{*-1} &= K_{UU}^{-1} \Phi^{-1} K_{UU}^{-1} \\ &= K_{UU}^{-1} + \sum_{i=1}^P K_{UU}^{-1} K_{UD_i} \Gamma_i^{-1} K_{D_i U} K_{UU}^{-1}. \end{aligned} \quad (79)$$

Thus, setting $F'(\mathcal{U}) \triangleq K_{UU}^{-1}$ and $F(\mathcal{U}, \mathbf{y}_{D_i}) \triangleq K_{UU}^{-1} K_{UD_i} \Gamma_i^{-1} K_{D_i U} K_{UU}^{-1}$ in (79) gives (15). On the other hand,

$$\begin{aligned} \Sigma^{*-1} \mu^* &= K_{UU}^{-1} K_{UD} \Gamma^{-1} \mathbf{y}_D \\ &= K_{UU}^{-1} \sum_{i=1}^P K_{UD_i} \Gamma_i^{-1} \mathbf{y}_{D_i} \\ &= \sum_{i=1}^P K_{UU}^{-1} K_{UD_i} \Gamma_i^{-1} \mathbf{y}_{D_i} \end{aligned} \quad (80)$$

where the first equality follows directly from (76) and (77) and the second equality is due to (66). Finally, setting $G(\mathcal{U}, \mathbf{y}_{D_i}) \triangleq K_{UU}^{-1} K_{UD_i} \Gamma_i^{-1} \mathbf{y}_{D_i}$ and $G'(\mathcal{U}) = \mathbf{0}$ in (80) gives (16).

D.3. Time Complexity of Computing Stochastic Gradient ($\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma$) (i.e., (17) and (18))

By assuming $|\mathcal{D}_i| = \mathcal{O}(|\mathcal{U}|)$ for $i = 1, \dots, P$, computing $F'(\mathcal{U})$, $F(\mathcal{U}, \mathbf{y}_{D_s})$, and $G(\mathcal{U}, \mathbf{y}_{D_s})$ incurs $\mathcal{O}(|\mathcal{U}|^3)$ time for SoR, DTC, FITC, FIC, PITC, and PIC (Appendix D.2). So, computing stochastic gradient ($\partial \hat{L} / \partial \mu, \partial \hat{L} / \partial \Sigma$) (i.e., (17) and (18)) incurs $\mathcal{O}(|S||\mathcal{U}|^3)$ time for SoR, DTC, FITC, FIC, PITC, and PIC that reduces to $\mathcal{O}(|\mathcal{U}|^3)$ time by setting $|S| = 1$ in our experiments. If the number t of iterations of SGA update is much smaller than $\min(|\mathcal{D}|/|\mathcal{U}|, P)$, then the anytime variants spanned by our unifying framework achieve a huge computational gain (i.e., $\mathcal{O}(t|\mathcal{U}|^3)$ time) over their corresponding SGPR models that incur $\mathcal{O}(|\mathcal{D}||\mathcal{U}|^2)$ time. In practice, t can also be used to trade off between time efficiency and the accuracy of estimating $q^*(\mathbf{f}_U)$.

D.4. Time Complexity of Prediction

This section analyzes the time incurred to predict f_x using (3) by assuming that $q(\mathbf{f}_U) = \mathcal{N}(\mu, \Sigma)$ is readily available from the SGA update (14) (Section 4). In particular, this includes the time incurred to evaluate the test conditional $q^*(f_x|\mathbf{y}_{D_P}, \mathbf{f}_U)$ and the time incurred to analytically integrate it with $q^*(\mathbf{f}_U) \equiv q(\mathbf{f}_U)$ in (3) to obtain the approximate predictive distribution $q^*(f_x) \approx p(f_x|\mathbf{y}_D)$ (3). To achieve this, we analytically derive $q^*(f_x)$ in terms of μ and Σ , as detailed next.

D.4.1. PIC

Recall that $q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \triangleq p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$. To derive $p(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}})$, we use the fundamental definition of GP (Rasmussen & Williams, 2006) to give the following closed-form expression for the Gaussian joint distribution $p(f_{\mathbf{x}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P})$:

$$\mathcal{N}\left(\mathbf{0}, \begin{bmatrix} k_{\mathbf{x}\mathbf{x}} & K_{\mathbf{x}\mathcal{U}} & K_{\mathbf{x}\mathcal{D}_P} \\ K_{\mathcal{U}\mathbf{x}} & K_{\mathcal{U}\mathcal{U}} & K_{\mathcal{U}\mathcal{D}_P} \\ K_{\mathcal{D}_P\mathbf{x}} & K_{\mathcal{D}_P\mathcal{U}} & K_{\mathcal{D}_P\mathcal{D}_P} \end{bmatrix}\right). \quad (81)$$

Integrating (81) with $p(\mathbf{y}_{\mathcal{D}_P}|\mathbf{f}_{\mathcal{D}_P}) = \mathcal{N}(\mathbf{f}_{\mathcal{D}_P}, \sigma_n^2 I)$ gives the following closed-form expression for $p(f_{\mathbf{x}}, \mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P})$:

$$\mathcal{N}\left(\mathbf{0}, \begin{bmatrix} k_{\mathbf{x}\mathbf{x}} & K_{\mathbf{x}\mathcal{U}} & K_{\mathbf{x}\mathcal{D}_P} \\ K_{\mathcal{U}\mathbf{x}} & K_{\mathcal{U}\mathcal{U}} & K_{\mathcal{U}\mathcal{D}_P} \\ K_{\mathcal{D}_P\mathbf{x}} & K_{\mathcal{D}_P\mathcal{U}} & K_{\mathcal{D}_P\mathcal{D}_P} + \sigma_n^2 I \end{bmatrix}\right).$$

Now, $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}) = \mathcal{N}(\mathbb{E}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}], \mathbb{V}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}])$ where

$$\mathbb{E}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}] = [K_{\mathbf{x}\mathcal{U}} \ K_{\mathbf{x}\mathcal{D}_P}] \begin{bmatrix} K_{\mathcal{U}\mathcal{U}} & K_{\mathcal{U}\mathcal{D}_P} \\ K_{\mathcal{D}_P\mathcal{U}} & K_{\mathcal{D}_P\mathcal{D}_P} + \sigma_n^2 I \end{bmatrix}^{-1} \begin{bmatrix} \mathbf{f}_{\mathcal{U}} \\ \mathbf{y}_{\mathcal{D}_P} \end{bmatrix},$$

$$\mathbb{V}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}] = k_{\mathbf{x}\mathbf{x}} - [K_{\mathbf{x}\mathcal{U}} \ K_{\mathbf{x}\mathcal{D}_P}] \begin{bmatrix} K_{\mathcal{U}\mathcal{U}} & K_{\mathcal{U}\mathcal{D}_P} \\ K_{\mathcal{D}_P\mathcal{U}} & K_{\mathcal{D}_P\mathcal{D}_P} + \sigma_n^2 I \end{bmatrix}^{-1} \begin{bmatrix} K_{\mathcal{U}\mathbf{x}} \\ K_{\mathcal{D}_P\mathbf{x}} \end{bmatrix}.$$

To simplify the above expressions, let

$$J \triangleq \begin{bmatrix} K_{\mathcal{U}\mathcal{U}} & K_{\mathcal{U}\mathcal{D}_P} \\ K_{\mathcal{D}_P\mathcal{U}} & K_{\mathcal{D}_P\mathcal{D}_P} + \sigma_n^2 I \end{bmatrix}^{-1} = \begin{bmatrix} J_{\mathcal{U}\mathcal{U}} & J_{\mathcal{U}\mathcal{D}_P} \\ J_{\mathcal{D}_P\mathcal{U}} & J_{\mathcal{D}_P\mathcal{D}_P} \end{bmatrix} \quad (82)$$

where $J_{\mathcal{U}\mathcal{U}}$, $J_{\mathcal{U}\mathcal{D}_P}$, $J_{\mathcal{D}_P\mathcal{U}}$ and $J_{\mathcal{D}_P\mathcal{D}_P}$ can be derived by applying the matrix inversion lemma for partitioned matrices directly. Then, plugging (82) into the expression of $\mathbb{E}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}]$,

$$\mathbb{E}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}] = B\mathbf{f}_{\mathcal{U}} + \ell \quad (83)$$

where $B \triangleq K_{\mathbf{x}\mathcal{U}}J_{\mathcal{U}\mathcal{U}} + K_{\mathbf{x}\mathcal{D}_P}J_{\mathcal{D}_P\mathcal{U}}$ and $\ell \triangleq (K_{\mathbf{x}\mathcal{U}}J_{\mathcal{U}\mathcal{D}_P} + K_{\mathbf{x}\mathcal{D}_P}J_{\mathcal{D}_P\mathcal{D}_P})\mathbf{y}_{\mathcal{D}_P}$. Then, $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}) = \mathcal{N}(B\mathbf{f}_{\mathcal{U}} + \ell, \mathbb{V}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}])$. As such, $q^*(f_{\mathbf{x}})$ can be analytically derived by integrating $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P})$ with $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$:

$$q^*(f_{\mathbf{x}}) = \int \mathcal{N}(B\mathbf{f}_{\mathcal{U}} + \ell, \mathbb{V}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}]) \mathcal{N}(\mu, \Sigma) d\mathbf{f}_{\mathcal{U}} = \mathcal{N}(B\mu + \ell, \mathbb{V}[f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}] + B\Sigma B^\top). \quad (84)$$

Since $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$ is already available from the SGA update (14) (Section 4) and $|\mathcal{D}_P| = \mathcal{O}(|\mathcal{U}|)$, computing $q^*(f_{\mathbf{x}})$ (84) incurs $\mathcal{O}(|\mathcal{U}|^3)$ time and space.

D.4.2. PITC, FITC, FIC AND DTC

Since we only analyze the complexity of single-input prediction at \mathbf{x}_* , the approximated testing conditionals of PITC, FITC, FIC and DTC are set to the same exact conditional of GP (Quiñero-Candela & Rasmussen, 2005; Snelson & Ghahramani, 2007):

$$q^*(f_{\mathbf{x}}|\mathbf{y}_{\mathcal{D}_P}, \mathbf{f}_{\mathcal{U}}) \triangleq p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mathbf{f}_{\mathcal{U}}, k_{\mathbf{x}\mathbf{x}} - Q_{\mathbf{x}\mathbf{x}}) \quad (85)$$

where $P_{\mathbf{x}\mathcal{U}} \triangleq K_{\mathbf{x}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}$ and $Q_{\mathbf{x}\mathbf{x}} \triangleq K_{\mathbf{x}\mathcal{U}}K_{\mathcal{U}\mathcal{U}}^{-1}K_{\mathcal{U}\mathbf{x}}$. As such, $q^*(f_{\mathbf{x}})$ can be analytically derived by integrating $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ with $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$:

$$q^*(f_{\mathbf{x}}) = \int \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mathbf{f}_{\mathcal{U}}, k_{\mathbf{x}\mathbf{x}} - Q_{\mathbf{x}\mathbf{x}}) \mathcal{N}(\mu, \Sigma) d\mathbf{f}_{\mathcal{U}} = \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mu, k_{\mathbf{x}\mathbf{x}} - Q_{\mathbf{x}\mathbf{x}} + P_{\mathbf{x}\mathcal{U}}\Sigma P_{\mathbf{x}\mathcal{U}}^\top). \quad (86)$$

Since $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$ is already available from the SGA update (14) (Section 4), computing $q^*(f_{\mathbf{x}})$ (86) incurs $\mathcal{O}(|\mathcal{U}|^3)$ time and space.

D.4.3. SoR

SoR (Smola & Bartlett, 2001) further simplifies the test conditional in (85) by additionally imposing a deterministic relation between $f_{\mathbf{x}}$ and $\mathbf{f}_{\mathcal{U}}$:

$$q^*(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}}, \mathbf{y}_{\mathcal{D}_P}) \triangleq \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mathbf{f}_{\mathcal{U}}, \mathbf{0}). \quad (87)$$

As such, $q^*(f_{\mathbf{x}})$ can be analytically derived by integrating $p(f_{\mathbf{x}}|\mathbf{f}_{\mathcal{U}})$ with $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$:

$$q^*(f_{\mathbf{x}}) = \int \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mathbf{f}_{\mathcal{U}}, \mathbf{0}) \mathcal{N}(\mu, \Sigma) d\mathbf{f}_{\mathcal{U}} = \mathcal{N}(P_{\mathbf{x}\mathcal{U}}\mu, P_{\mathbf{x}\mathcal{U}}\Sigma P_{\mathbf{x}\mathcal{U}}^\top). \quad (88)$$

Since $q(\mathbf{f}_{\mathcal{U}}) = \mathcal{N}(\mu, \Sigma)$ is already available from the SGA update (14) (Section 4), computing $q^*(f_{\mathbf{x}})$ (88) incurs $\mathcal{O}(|\mathcal{U}|^3)$ time and space.

E. Canonical Parameterization of Gaussian Distribution

This section features the canonical parameterization of Gaussian distribution and highlights some of its properties that have been previously used in Section 4. For clarity and ease of reading, we demonstrate this in the context of $q(\mathbf{f}_{\mathcal{U}})$, which is originally specified using the moment parameterization $q(\mathbf{f}_{\mathcal{U}}) \triangleq \mathcal{N}(\mu, \Sigma)$. In particular, let $\theta \triangleq [\theta_1; \text{vec}(\theta_2)]$ where $\theta_1 \triangleq \Sigma^{-1}\mu$ and $\theta_2 \triangleq -(1/2)\Sigma^{-1}$. We begin with the following re-parameterization of $q(\mathbf{f}_{\mathcal{U}}|\theta)$ with respect to θ :

$$q(\mathbf{f}_{\mathcal{U}}|\theta) = \mathcal{N}(\mu, \Sigma) = h(\mathbf{f}_{\mathcal{U}}) \exp(\theta^\top \mathbf{T}(\mathbf{f}_{\mathcal{U}}) - A(\theta)) \quad (89)$$

where $\mathbf{T}(\mathbf{f}_U) \triangleq (\mathbf{f}_U; \text{vec}(\mathbf{f}_U \mathbf{f}_U^\top))$, $h(\mathbf{f}_U) \triangleq (2\pi)^{-|U|/2}$, and the normalizing function $A(\theta)$ is defined as

$$\begin{aligned} A(\theta) &\triangleq -\frac{1}{2} \text{tr}(\theta_2 \theta_1 \theta_1^\top) - \frac{1}{2} \log |-2\theta_2| \\ &= \frac{1}{2} \mu^\top \Sigma^{-1} \mu + \frac{1}{2} \log |\Sigma^{-1}|. \end{aligned} \quad (90)$$

Alternatively, we can define $Z(\theta)$ such that $A(\theta) = \log Z(\theta)$ and rewrite (89) as

$$q(\mathbf{f}_U|\theta) = \frac{1}{Z(\theta)} h(\mathbf{f}_U) \exp(\theta^\top \mathbf{T}(\mathbf{f}_U)). \quad (91)$$

Since $\int q(\mathbf{f}_U|\theta) d\mathbf{f}_U = 1$, (91) implies

$$Z(\theta) = \int h(\mathbf{f}_U) \exp(\theta^\top \mathbf{T}(\mathbf{f}_U)) d\mathbf{f}_U. \quad (92)$$

Using the above results, we are now able to verify the identities employed in Section 4, as detailed in the subsections below.

E.1. Evaluating $\eta \triangleq \mathbb{E}[\mathbf{T}(\mathbf{f}_U)]$

This section will show how $\mathbb{E}[\mathbf{T}(\mathbf{f}_U)]$ can be derived as a function of μ and Σ . To achieve this, we first use the definition of $\mathbf{T}(\mathbf{f}_U)$ to obtain

$$\begin{aligned} \mathbb{E}[\mathbf{T}(\mathbf{f}_U)] &= (\mathbb{E}[\mathbf{f}_U]; \mathbb{E}[\text{vec}(\mathbf{f}_U \mathbf{f}_U^\top)]) \\ &= (\mathbb{E}[\mathbf{f}_U]; \text{vec}(\mathbb{E}[\mathbf{f}_U \mathbf{f}_U^\top])) \\ &= [\mu; \text{vec}(\mu \mu^\top + \Sigma)] \end{aligned} \quad (93)$$

where the second equality follows from the definition of vec and expectation of vector while the last equality is derived using the Gaussian identities $\mathbb{E}[\mathbf{f}_U] = \mu$ and $\mathbb{E}[\mathbf{f}_U \mathbf{f}_U^\top] = \mu \mu^\top + \Sigma$ since $q(\mathbf{f}_U) = \mathcal{N}(\mu, \Sigma)$. In particular, if $\eta_1 \triangleq \mu$ and $\eta_2 \triangleq \mu \mu^\top + \Sigma$, then $\eta \triangleq [\eta_1; \text{vec}(\eta_2)] = \mathbb{E}[\mathbf{T}(\mathbf{f}_U)]$.

E.2. Evaluating $H(\theta)$

This section focuses on explicitly representing the Fisher information $H(\theta)$ in terms of $A(\theta)$. To achieve this, note that

$$\begin{aligned} \frac{\partial \log q(\mathbf{f}_U|\theta)}{\partial \theta} &= \frac{\partial}{\partial \theta} (\theta^\top \mathbf{T}(\mathbf{f}_U) - A(\theta)) \\ &= \mathbf{T}(\mathbf{f}_U) - \frac{\partial A(\theta)}{\partial \theta}. \end{aligned} \quad (94)$$

By applying the identity $\partial f(x)/\partial x^\top = (\partial f(x)/\partial x)^\top$ to (94),

$$\begin{aligned} \frac{\partial \log q(\mathbf{f}_U|\theta)}{\partial \theta^\top} &= \mathbf{T}(\mathbf{f}_U)^\top - \left(\frac{\partial A(\theta)}{\partial \theta} \right)^\top \\ &= \mathbf{T}(\mathbf{f}_U)^\top - \frac{\partial A(\theta)}{\partial \theta^\top}. \end{aligned} \quad (95)$$

Then, differentiating both sides of (95) with respect to θ yields

$$\frac{\partial^2 \log q(\mathbf{f}_U|\theta)}{\partial \theta \partial \theta^\top} = -\frac{\partial^2 A(\theta)}{\partial \theta \partial \theta^\top}. \quad (96)$$

Plugging (96) into the definition of $H(\theta)$ reveals (19).

E.3. Proof of $\partial \eta / \partial \theta = H(\theta)$

Let us differentiate $A(\theta)$ with respect to θ :

$$\begin{aligned} \frac{\partial A(\theta)}{\partial \theta} &= \frac{\partial \log Z(\theta)}{\partial \theta} = \frac{1}{Z(\theta)} \frac{\partial Z(\theta)}{\partial \theta} \\ &= \frac{1}{Z(\theta)} \int h(\mathbf{f}_U) \frac{\partial}{\partial \theta} \exp(\theta^\top \mathbf{T}(\mathbf{f}_U)) d\mathbf{f}_U \\ &= \frac{1}{Z(\theta)} \int h(\mathbf{f}_U) \exp(\theta^\top \mathbf{T}(\mathbf{f}_U)) \mathbf{T}(\mathbf{f}_U) d\mathbf{f}_U \\ &= \int q(\mathbf{f}_U|\theta) \mathbf{T}(\mathbf{f}_U) d\mathbf{f}_U = \mathbb{E}[\mathbf{T}(\mathbf{f}_U)] = \eta. \end{aligned} \quad (97)$$

Thus, (97) implies $\eta^\top = \partial A(\theta) / \partial \theta^\top$. Hence, differentiating both sides of this equality with respect to θ yields $\partial \eta^\top / \partial \theta = \partial^2 A(\theta) / \partial \theta \partial \theta^\top = H(\theta)$. Since $H(\theta)$ is symmetric, it follows that $\partial \eta / \partial \theta = H(\theta)$.