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# Technical Appendix on Sparse Bayesian Regression 

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#### Abstract

We report the technical details for a sparse bayesian approach to regression. It can be seen as an extension of the Relevance Vector Machine of Tipping et al [1] to a more general setting that can handle vector-valued regression and generic quadratic priors.


## 1 Quadratic Energies \& Marginal Likelihood of the Data

We want to minimize in $\boldsymbol{w}$ the following cost:

$$
\begin{equation*}
E(\boldsymbol{w})=(\boldsymbol{t}-\Phi \boldsymbol{w})^{\top} \beta \mathbf{H}(\boldsymbol{t}-\Phi \boldsymbol{w})+\boldsymbol{w}^{\top}(\mathbf{A}+\lambda \mathbf{R}) \boldsymbol{w} \tag{1}
\end{equation*}
$$

and jointly optimize in $\mathbf{A}=\operatorname{diag}\left(\mathbf{A}_{i}\right), \beta, \lambda . \boldsymbol{w} \mid \boldsymbol{t}, \mathbf{A}, \lambda, \beta$ follows a Gaussian distribution $\mathcal{N}(\mu, \Sigma)$, where

$$
\begin{equation*}
\boldsymbol{\mu}=\Sigma \Phi^{\top} \beta \mathbf{H} \boldsymbol{t}, \quad \Sigma=\left(\mathbf{A}+\lambda \mathbf{R}+\Phi^{\top} \beta \mathbf{H} \Phi\right)^{-1} \tag{2}
\end{equation*}
$$

A key element is that the distribution of $\boldsymbol{t} \mid \mathbf{A}$ is also Gaussian, $\boldsymbol{t} \mid \mathbf{A} \sim \mathcal{N}(\mathbf{0}, \mathbf{C})$, with

$$
\begin{equation*}
\mathbf{C}=(\beta \mathbf{H})^{-1}+\Phi(\mathbf{A}+\lambda \mathbf{R})^{-1} \Phi^{\top} \tag{3}
\end{equation*}
$$

Indeed, we see that:

$$
\begin{aligned}
p(\boldsymbol{t} \mid \mathbf{A}) & =\int p(\boldsymbol{t} \mid \boldsymbol{w}) p(\boldsymbol{w} \mid \mathbf{A}) d \boldsymbol{w} \\
& \propto \int \exp -\frac{1}{2}(\boldsymbol{t}-\Phi \boldsymbol{w})^{\top} \beta \mathbf{H}(\boldsymbol{t}-\Phi \boldsymbol{w}) \cdot \exp -\frac{1}{2} \boldsymbol{w}^{\top}(\mathbf{A}+\lambda \mathbf{R}) \boldsymbol{w} \cdot d \boldsymbol{w} \\
& \propto \exp -\frac{1}{2} \boldsymbol{\mu}^{\top} \Sigma^{-1} \boldsymbol{\mu} \cdot \exp -\frac{1}{2} \boldsymbol{t}^{\top} \beta \mathbf{H} \boldsymbol{t} \cdot \int \exp -\frac{1}{2}(\boldsymbol{w}-\boldsymbol{\mu})^{\top} \Sigma^{-1}(\boldsymbol{w}-\boldsymbol{\mu}) d \boldsymbol{w}
\end{aligned}
$$

where the integral sums to 1 . Thus using (2),

$$
\left.p(\boldsymbol{t} \mid \mathbf{A}) \propto \exp -\frac{1}{2} \boldsymbol{t}^{\boldsymbol{\top}}(\beta \mathbf{H}) \boldsymbol{t}-\frac{1}{2} \boldsymbol{t}^{\boldsymbol{\top}}(\beta \mathbf{H}) \Phi \Sigma \Phi^{\boldsymbol{\top}}(\beta \mathbf{H})\right) \boldsymbol{t}
$$

$\boldsymbol{t} \mid \mathbf{A}$ is Gaussian since the distribution is proportional to a Gaussian, and by identification it ensues that $\mathbf{C}^{-1}=\beta \mathbf{H}-(\beta \mathbf{H}) \Phi \Sigma \Phi^{\top}(\beta \mathbf{H})$. We then get the desired result using a matrix inversion identity. The fast RVM algorithm proceeds by iteratively implementing a single change to one of the $\mathbf{A}_{i}$ 's on the block-diagonal matrix $\mathbf{A}$; specifically the one that maximizes the increase of a quantity known as the evidence, $\log p(\boldsymbol{t} \mid \mathbf{A})$, then re-estimating the parameters of the conditional posterior $\boldsymbol{w} \mid \boldsymbol{t}, \mathbf{A}, \lambda, \beta$ using (2). The algorithm starts with all $\mathbf{A}_{i}$ set to $\infty$. The computations involve rank-one "block" updates; it also turns out that the optimal $\mathbf{A}_{i}$ 's are rank one matrices (so we actually have rank-one updates).

## 2 Computation of the gain in evidence for a given action

We want to evaluate the change in $\log p\left(t \mid \mathbf{A}_{-i}, \mathbf{A}_{i}\right)$ when a single prior weight $\mathbf{A}_{i}$ is changed. Recall that $\boldsymbol{t} \mid \mathbf{A}_{-i}, \mathbf{A}_{i} \sim \mathcal{N}(\mathbf{0}, \mathbf{C})$, thus

$$
\begin{equation*}
\log p\left(\boldsymbol{t} \mid \mathbf{A}_{-i}, \mathbf{A}_{i}\right)=-1 / 2\left\{\log |\mathbf{C}|+\boldsymbol{t}^{\top} \mathbf{C}^{-1} \boldsymbol{t}\right\} \tag{4}
\end{equation*}
$$

Now let's single out the contribution of $\mathbf{A}_{i}$, for each of the two terms above. Noting the form of eq. (3), let $\mathbf{L}=(\mathbf{A}+\lambda \mathbf{R})^{-1}$. We can single out the contribution of the $i$ th basis to $\mathbf{L}$ first, as a rank one update:

$$
\mathbf{L}=\left(\begin{array}{cc}
\mathbf{L}_{-i} & 0  \tag{5}\\
0 & 0
\end{array}\right)+\mathrm{U}_{i} l_{i} \mathrm{U}_{i}^{\top}
$$

with $\mathrm{U}_{i}^{\top}=\left(\left(\lambda \mathbf{L}_{-i} \mathbf{R}_{i}\right)^{\top} \mathrm{I}^{\top}\right)$ and $l_{i}=\left\{\mathbf{A}_{i}+\lambda \mathbf{R}_{i i}-\left(\lambda \mathbf{R}_{i}\right)^{\top} \mathbf{L}_{-i}\left(\lambda \mathbf{R}_{i}\right)\right\}^{-1} \triangleq\left\{\mathbf{A}_{i}+\right.$ $\left.\kappa_{i}\right\}^{-1}$. Note also that any basis for which $\mathbf{A}_{j}=\infty$ can be disregarded, since its corresponding line and column in $\mathbf{L}$ and $\mathbf{L}_{-i}$ will be null. We can interpret $\mathbf{L}$ as a square matrix of dimension the number of active bases (including the basis under consideration), and the algorithmic complexity of matrix operations involving $\mathbf{L}$ or $\mathbf{L}_{-i}$ will indeed be related to the reduced sized of these matrices. This is also true for $\mathrm{U}_{i}$, as a direct consequence, and for all of the other quantities involved.

Injecting (5) into (3) gives a decomposition of $\mathbf{C}$ into the sum of a term that does not depend on the $i$ th basis and of a rank-one term:

$$
\begin{equation*}
\mathbf{C}=\mathbf{C}_{-i}+\Phi^{i} \mathrm{U}_{i} l_{i} \mathrm{U}_{i}^{\top} \Phi^{i^{\top}} \tag{6}
\end{equation*}
$$

$\Phi$ is superscripted with $i$ to recall that along with all the other active bases, the $i$ th basis $\phi_{i}$ is present in this matrix. Using rank-one updates for, respectively, the determinant and the inverse, and letting $\mathbf{C}_{-i}^{-1} \triangleq\left(\mathbf{C}_{-i}\right)^{-1}$, we get the two following expressions:

$$
\begin{gather*}
|\mathbf{C}|=\left|\mathbf{C}_{-i}\right|\left|l_{i}\right|\left|l_{i}^{-1}+\mathrm{U}_{i}^{\top} \Phi^{i^{\top}} \mathbf{C}_{-i}^{-1} \Phi^{i} \mathrm{U}_{i}\right|  \tag{7}\\
\boldsymbol{t}^{\top} \mathbf{C}^{-1} \boldsymbol{t}=\boldsymbol{t}^{\top}\left(\mathbf{C}_{-i}^{-1}-\mathbf{C}_{-i}^{-1} \Phi^{i} \mathrm{U}_{i}\left\{l_{i}^{-1}+\mathrm{U}_{i}^{\top} \Phi^{i \top} \mathbf{C}_{-i}^{-1} \Phi^{i} \mathrm{U}_{i}\right\}^{-1} \mathrm{U}_{i}^{\top} \Phi^{i^{\top}} \mathbf{C}_{-i}^{-1}\right) \boldsymbol{t} \tag{8}
\end{gather*}
$$

These quantities rewrite more compactly if we introduce appropriate notations. Namely, let $q_{j}(i) \triangleq \phi_{j}^{\top} \mathbf{C}_{-i}^{-1} \boldsymbol{t} \in \mathbb{R}^{d}$ and $s_{j k}(i) \triangleq \phi_{j}^{\top} \mathbf{C}_{-i}^{-1} \phi_{k} \in \mathcal{M}_{d, d}$. The concatenation of these $q_{j}$ 's for all active bases plus the basis under scrutiny (total of m bases), a.k.a $\boldsymbol{q}_{i} \in \mathbb{R}^{d \times m}$, will come in helpful. Similarly $s_{i} \in \mathcal{M}_{d \times m, d \times m}$ will denote the matrix with $s_{k l}(i)$ as $(k, l)$ th coefficient, where indices span the set of active bases plus the $i$ th basis. Now, let $q^{i} \triangleq \mathrm{U}_{i}^{\top} \Phi^{i \top} \mathbf{C}_{-i}^{-1} \boldsymbol{t}=\mathrm{U}_{i}^{\top} \boldsymbol{q}_{i} \in \mathbb{R}^{d}$, and $s^{i} \triangleq \mathrm{U}_{i}^{\top} \Phi^{i^{\top}} \mathbf{C}_{-i}^{-1} \Phi^{i} \mathrm{U}_{i}=\mathrm{U}_{i}^{\top} s_{i} \mathrm{U}_{i} \in \mathcal{M}_{d, d}$. With these notations in hand and recalling that $l_{i}^{-1}=\mathbf{A}_{i}+\kappa_{i}$, we can rewrite eq. (7) and eq. (8) as:

$$
\begin{gather*}
\log |\mathbf{C}|=\log \left|\mathbf{C}_{-i}\right|-\log \left|\mathbf{A}_{i}+\kappa_{i}\right|+\log \left|\mathbf{A}_{i}+\kappa_{i}+s^{i}\right|  \tag{9}\\
\boldsymbol{t}^{\top} \mathbf{C}^{-1} \boldsymbol{t}=\boldsymbol{t}^{\boldsymbol{\top}} \mathbf{C}_{-i}^{-1} \boldsymbol{t}-q^{i \top}\left\{\mathbf{A}_{i}+\kappa_{i}+s^{i}\right\}^{-1} q^{i} \tag{10}
\end{gather*}
$$

Ignoring the terms that do not depend on the $i$ th basis, we see that the contribution to the evidence of the model for a given value of $\mathbf{A}_{i}$ is directly related to:

$$
\begin{equation*}
l\left(\mathbf{A}_{i}\right)=\log \left|\mathbf{A}_{i}+\kappa_{i}\right|-\log \left|\mathbf{A}_{i}+\kappa_{i}+s^{i}\right|+q^{i^{\top}}\left\{\mathbf{A}_{i}+\kappa_{i}+s^{i}\right\}^{-1} q^{i} \tag{11}
\end{equation*}
$$

Naturally if $\lambda=0$ (no additional regularization) this comes down to the regular RVM, with $q^{i}=q_{i}, s^{i}=s_{i i}$ and $\kappa_{i}=0$.

## 3 Maximization of the gain in evidence

If $q^{i} q^{i \top}-s^{i}$ has no positive eigenvalue, the maximum $\mathbf{A}_{i}$ lies at infinity and the basis should remain inactive, or be removed. Otherwise the gradient of Eq. (11) provides ground to look for rank-one maximizers $\mathbf{A}_{i}=\alpha_{i} \eta_{i} \eta_{i}^{\top}$. To that aim we compute $n_{i}=$ $s^{i^{-1}} q^{i}$ and

$$
\begin{equation*}
a_{i}=\frac{\left(n_{i}^{\top} s^{i} n_{i}\right)^{2}}{\left(n_{i}^{\top} q^{i}\right)^{2}-n_{i}^{\top} s^{i} n_{i}}-n_{i}^{\top} \kappa_{i} n_{i} \tag{12}
\end{equation*}
$$

If $a_{i} \geq 0$ the maximizer is given by $\alpha_{i}=a_{i}$ and $\eta_{i}=n_{i}$. Otherwise ( $a_{i}<0$ ) we set $\alpha_{i}=0$ and numerically solve over the optimal orientation $\eta_{i}$. This latter case arises when the regularization level alone is sufficient to make additional "shrinkage" unnecessary.

## 4 Update of $\boldsymbol{\lambda}$

We derive an update rule via an expectation-maximization procedure. Knowing $\boldsymbol{w}$, it would be straightforward to derive an estimate of $\lambda$ by maximizing the log-likelihood of $\boldsymbol{w}$ or the posterior of $\lambda$ given $\boldsymbol{w}$. However $\boldsymbol{w}$ is hidden in our model. Instead, we try to maximize the log-likelihood on average (i.e. to minimize the average loss): $\max _{\lambda} \mathbb{E}_{\boldsymbol{w} \sim \mathcal{N}(\boldsymbol{\mu}, \Sigma)}\left[\log p(\boldsymbol{t}, \boldsymbol{w} \mid \mathbf{A}, \lambda, \beta) \mid \mathbf{A}^{*}, \beta^{*}\right]$, where $\mathbf{A}^{*}$ and $\beta^{*}$ are our current estimates of the respective quantities. Discarding terms which are constants of $\lambda$, we obtain the following:

$$
\begin{equation*}
\lambda^{*}=\underset{\lambda}{\arg \max }-\frac{\lambda}{2} \operatorname{tr}(\Sigma \mathbf{R})+\frac{1}{2} \log |\mathbf{A}+\lambda \mathbf{R}|-\frac{\lambda}{2} \boldsymbol{\mu}^{\top} \mathbf{R} \boldsymbol{\mu} \tag{13}
\end{equation*}
$$

Deriving leads to:

$$
\begin{equation*}
\frac{\partial}{\partial \lambda} f(\lambda) \propto \operatorname{tr}\left(\{\mathbf{A}+\lambda \mathbf{R}\}^{-1} \mathbf{R}\right)-\operatorname{tr}(\Sigma \mathbf{R})-\boldsymbol{\mu}^{\top} \mathbf{R} \boldsymbol{\mu} \tag{14}
\end{equation*}
$$

This is a decreasing function of $\lambda$ and thus has at most one zero. If $\frac{\partial}{\partial \lambda} f(\lambda)$ is negative at the origin, $\lambda^{*}=0$. Otherwise, we optimize by using the Newton method on a log scale. This is motivated by the fact that the function of interest is not only decreasing, but also smooth and convex. Lastly note that $\{\mathbf{A}+\lambda \mathbf{R}\}^{-1} \mathbf{R}=\left\{\mathbf{R}^{-1} \mathbf{A}+\lambda \mathrm{I}\right\}^{-1}$, so we can compute the eigenvalues $\delta_{k}$ of $\mathbf{A}^{1 / 2} \mathbf{R}^{-1} \mathbf{A}^{1 / 2}$ once and rely on the fact that $\operatorname{tr}\left(\{\mathbf{A}+\lambda \mathbf{R}\}^{-1} \mathbf{R}\right)=\sum_{k} 1 /\left(\delta_{k}+\lambda\right)$ to avoid repeated matrix inversions.

## References

1. Tipping, M.E., Faul, A.C., et al.: Fast marginal likelihood maximisation for sparse bayesian models. In: Workshop on artificial intelligence and statistics. Volume 1., Jan (2003)
