# An Econometric Model of International Long-run Growth Dynamics Appendix 

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## Appendix 1 Data

The data are annual values of real per-capita GDP for 113 countries spanning the 118-year period 1900-2017, taken from the Penn Word Table 9.1 (Feenstra, et al. (2015) and Inklar and Woltjer (2019)) and Bolt et al. (2018) Maddison Project Database. GDP is measured at constant 2011 national prices, expressed in U.S. dollars.

From the Penn World Table, we used the real GDP series rgdp ${ }^{n a}$, which measures real GDP at constant national prices, and is obtained from national accounts data for each country. We also extracted the population series population from the PWT. From these, we computed $r g d p n a p c=r g d p^{n a} /$ population for each year and country.

From the Maddison Project Database we extracted rgpdnapc, which is the same concept as $r g p d^{n a}$ in PWT, but expressed in per-capita terms. We also extracted the population series pop.

We then linked the series for each country as follows as follows: Letting $i$ denote country, $t$ denote year, and $t_{i}$ denote the first year that rgdpnapc is available in the PWT. Then

$$
\operatorname{rgdpnapc}_{i, t}^{\text {linked }}=\left\{\begin{array}{c}
\operatorname{rgdpnapc}_{i, t}^{P W T} \text { for } t \geq t_{i} \\
\operatorname{grpdnapc}_{i, t}^{\text {Maddison }} \times\left(\operatorname{rgdpnapc}_{i, t_{i}}^{P W T} / \operatorname{rgdpnapc}_{i, t_{i}}^{\text {Maddison }}\right) \text { for } t<t_{i}
\end{array}\right.
$$

Population was linked analogously. Results in the paper are based on these linked series.
The file TableA1_List_of_Countries.xlsx lists the 113 countries used in the analysis, each country's 2017-value of per-capita real GDP and population, and start and end dates. The file TableA2_Excluded_Countries. xlsx lists the 69 countries in the PWT that were excluded because of a short sample (less than 50 years), a small 2017 population (less than 3 million) or both.

The sample ends in 2017 (the last year available in PWT 9.1) and begins in 1900, where the (arbitrary) start date means that the sample contains at most 118 years of data. While data are available for some countries in earlier years from the Maddison Project Database, concerns about temporal stability led us to start the sample in 1900.

## Appendix 2 Computations

## Appendix 2.1 Unbalanced Panel

Müller and Watson $(2008$, 2018) extract low-frequency trends by projecting the original time series $y_{t}$ onto a finite set of trigonometric low-frequency regressors. In particular, for time series that have a deterministic trend, their suggestion for the regressors are the (asymptotic version) of the eigenvectors associated with the largest eigenvalues of a demeaned and detrended random walk. The variation in a random walk is dominated by low-frequency components (cf. Phillips (1998)). Thus, high variance linear combinations of a (demeaned and detrended) random walk extract low frequency variability, and the magnitude of the variance becomes an indicator for the associated frequency.

To be specific, let $V$ be the $T \times T$ covariance matrix of a random walk with unit variance and initialized at $t=0$, that is, $V_{i j}=\min (i, j)$; let $R$ denote the $T \times 2$ matrix with $t$ th row equal to $(1, t-(T+1) / 2)$. The full sample regressors that extract variation below the frequency cut-off $2 q / T$ are then given by $q+1$ vectors: Two corresponding to the columns of $R$, plus the $q-1$ eigenvectors of the $T \times T$ matrix

$$
\begin{equation*}
M V M \quad \text { with } M=I_{T}-R\left(R^{\prime} R\right)^{-1} R^{\prime} \tag{1}
\end{equation*}
$$

corresponding to the $q-1$ largest eigenvalues. As discussed in Müller and Watson (2008), the cut-off formula $2 q / T$ requires counting the the linear trend in the second column of $R$ as part of the $q$ non-constant vectors that extract variation. For future reference, denote the $q-1$ largest eigenvalue of $M V M$ by $\kappa_{q}$. Note that the regressors are orthogonal to each other by construction.

We now generalize this approach further and make it operational in an unbalanced panel. Let $J_{i}$ be a diagonal $T \times T$ matrix whose $t$ th diagonal element is equal to one if the data $y_{i, t}$ for country $i$ is observed at time $t$, and zero otherwise. Let $R_{i}=J_{i} R$, a $T \times 2$ matrix with zeros in all rows where $y_{i, t}$ is unobserved. The regressors for series $y_{i, t}$ are then given by $q_{i}+1$ weights: The two columns of $R_{i}$, and the $q_{i}-1$ eigenvectors of

$$
\begin{equation*}
M_{i} V M_{i} \text { with } M_{i}=J_{i}-R_{i}\left(R_{i}^{\prime} R_{i}\right)^{-1} R_{i}^{\prime} \tag{2}
\end{equation*}
$$

with eigenvalues larger or equal than $\kappa_{q}$. Note that the eigenvectors of $M_{i} V M_{i}$ are equal to zero by construction in all rows where $y_{i, t}$ is unobserved. Also note that if $y_{i, t}$ is observed
for a single spell of $T_{i}$ consecutive periods, then this construction yields the same set of regressors that would have been obtained from (1) for a fully observed time series of length $T_{i}$ with cut-off frequency $2 q_{i} / T_{i} \approx 2 q / T$.

## Appendix 2.2 Low-Frequency Trends and their Asymptotic Distribution

Let the $q_{i}+1$ vector $Y_{i}$ be the OLS coefficients of a regression of $\left\{y_{i, t}\right\}_{t=1}^{T}$ on the $q_{i}+1$ regressors obtained in this manner. The low-frequency panel analysis treats $Y_{i}$ as the only observation from country $i$, ignoring all higher frequency sample variation contained in the original series $\left\{y_{i, t}\right\}_{t=1}^{T}$.

The unbalanced panel introduces a complication in preserving the factor structure (1) for the low frequency trends for $y_{i, t}$ and the common factor $f_{t}$. We procede by introducing a common higher dimensional "baseline" low-frequency trend $X_{i}$, of which $Y_{i}$ is a partial observation (even if there is no missing data for country $i$ ). Specifically, consider the space spanned by the $\tilde{q}+1$ regressors extracted from (1), with $\tilde{q} \gg q$. For sufficiently large $\tilde{q}$, we can approximate the $q_{i}+1$ regressors obtained from (2) arbitrarily well, that is, the regressors obtained from (2) can be written as a linear combination of the $\tilde{q}+1$ baseline regressors. For computational efficiency, it is not attractive to make $\tilde{q}$ too large; we choose $\tilde{q}=q+15$ throughout. We then use ordinary least squares regressions (with $T$ observations) to find the best linear combination to approximate the $q_{i}+1$ regressors as linear combinations of the $\tilde{q}+1$ baseline regressors (and to make this approximation more accurate, we treat as missing any isolated observed value, that is a value $y_{i, t}$ where both $y_{i, t-1}$ and $y_{i, t+1}$ are missing). The smallest $R^{2}$ across all regressors and countries in this construction is larger than 0.973.

With this approximation in place, the relationship between the observed $q_{i}+1$ vector $Y_{i}$ and the latent augmented $\tilde{q}+1$ trend $X_{i}$ is given by

$$
\begin{equation*}
Y_{i}=B_{i}^{\prime} X_{i} \tag{3}
\end{equation*}
$$

where $B_{i}$ is the $\tilde{q}+1 \times q_{i}+1$ matrix obtained from the OLS regressions just described.
The linear structure of the model of equations (1)-(7) then yields a corresponding linear structure in the $\tilde{q}+1$ dimensional low-frequency trends

$$
\begin{equation*}
X_{i}=F+C_{i}, \quad i=1, \ldots, n \tag{4}
\end{equation*}
$$

$$
\begin{align*}
C_{i} & =\mu_{c} \iota_{1}+\lambda_{c, i} G_{J(i)}+U_{c, i}, \quad i=1, \ldots, n  \tag{5}\\
G_{j} & =\lambda_{g, j} H_{K(j)}+U_{g, j}, \quad j=1, \ldots, 25  \tag{6}\\
H_{k} & =U_{h, k}, \quad k=1, \ldots, 10  \tag{7}\\
F & =\iota_{1} f_{0}+\iota_{2} \mu_{m}+S_{m}+A  \tag{8}\\
Y^{0} & =\sum_{i=1}^{n} w_{i} X_{i}-F \tag{9}
\end{align*}
$$

where $\iota_{1}$ and $\iota_{2}$ in (5) and (8) are the first and second column of $I_{\tilde{q}+1}$, respectively, and $F$, $C_{i}, U_{c, i}, U_{g, j}, U_{h, k}, G_{j}, H_{k}, S_{m}$ and $A$ are the $\tilde{q}+1$ low-frequency trends obtained from $f_{t}$, $c_{i, t}, u_{c, i, t}, u_{g, j, t}, u_{h, k, t}, g_{j, t}, h_{k, t}, m_{t}$ and $\Delta a_{t}$, respectively. Equation (9) here represents the in-sample prior that $f_{t}$ is close to the population weighted average of OECD countries' $y_{i, t}$, implemented via a dummy observation prior where $Y^{0} \sim \mathcal{N}(0, \Delta)$ with $\Delta=0.01^{2} I_{\tilde{q}+1}$ in the prior, and the observed realization is $Y^{0}=0$. The population weights $w_{i}$ in (9) sum to one and are set to zero for non-OECD countries.

Now as discussed in Section 4.2, the focus on low-frequency variation offers the advantage that the large sample distribution of the low-frequency trends $U_{c, i}, U_{G, j}, U_{h, k}, S_{m}$ and $A$ of the driving innovations of the model are approximately Gaussian, with the large sample covariance matrix a function of the low-frequency scale and persistence parameters. For instance, recall from (8) that the $u_{t}$ terms in (2)-(4) are modelled as a weighted sum of two independent low-frequency (=local-to-unity) stationary $\operatorname{AR}(1)$ processes with unit longrun variance. Under such asymptotics, the low-frequency trend $U$ computed from $\left\{u_{t}\right\}_{t=1}^{T}$ satisfies $U \stackrel{a}{\sim} \mathcal{N}\left(0, \sigma_{u}^{2} \Sigma_{U}\left(\theta^{u}\right)\right)$ with $\theta^{u}=\left(\rho_{1}, \rho_{2}, \zeta\right)$, that is, the limiting Gaussian distribution only depends on low-frequency parameters and not, say, on the short-run dynamics of the disturbances driving the local-to-unity processes. In other words, under such asymptotics, $U$ is distributed as if the underlying $u_{t}$ was a sum of two independent Gaussian AR (1) processes with coefficient $\rho_{1}$ and $\rho_{2}$, respectively. Thus, an asymptotically justified approximation to the the $\tilde{q}+1 \times \tilde{q}+1$ covariance matrix $\Sigma_{U}\left(\theta^{u}\right)$ is simply given by

$$
\Sigma_{U}\left(\theta^{u}\right)=\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \tilde{R}^{\prime} \tilde{\Sigma}_{u}\left(\theta^{u}\right) \tilde{R}\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1}
$$

where $\tilde{R}$ is the $T \times(\tilde{q}+1)$ dimensional matrix of baseline regressors, and $\tilde{\Sigma}^{u}\left(\theta_{u}\right)$ is the $T \times T$ covariance matrix of a $\left(\zeta, \sqrt{1-\zeta^{2}}\right)$ weighted average of two independent stationary $\operatorname{AR}(1)$ processes with coefficients $\rho_{1}$ and $\rho_{2}$ and innovation variance equal to $\left(1-\rho_{1}^{2}\right)^{-1 / 2}$
and $\left(1-\rho_{2}^{2}\right)^{-1 / 2}$, respectively. The same arguments also yield

$$
\begin{align*}
& S_{m} \stackrel{a}{\sim} \mathcal{N}\left(0, \sigma_{m}^{2} \Sigma_{m}\left(\rho_{m}\right)\right)  \tag{10}\\
& A \stackrel{a}{\sim} \mathcal{N}\left(0, \sigma_{\Delta a}^{2} \Sigma_{a}\right) \tag{11}
\end{align*}
$$

where

$$
\begin{aligned}
\Sigma_{m}\left(\rho_{m}\right) & =\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \tilde{R}^{\prime} \tilde{\Sigma}_{m}\left(\rho_{m}\right) \tilde{R}\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \\
\Sigma_{a} & =\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \tilde{R}^{\prime} V \tilde{R}\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1}
\end{aligned}
$$

with $\tilde{\Sigma}_{m}\left(\rho_{m}\right)$ the $T \times T$ covariance matrix of the partial sum of a stationary Gaussian meanzero $\operatorname{AR}(1)$ model with coefficient $\rho_{m}$ and innovation variance equal to $\left(1-\rho_{m}^{2}\right)^{-1}$ (so that the unconditional variance of the $\operatorname{AR}(1)$ process is normalized to be equal to unity), and $V$ is the random-walk covariance matrix introduced earlier.

Further note that the joint large-sample distribution of the in-sample low-frequency trends and out-of-sample long-horizon forecasts is again normal, as discussed and exploited in Müller and Watson (2016). Taking the example of the forecast of $u_{T+\lfloor h T\rfloor}$ for some sample size independent $h>0$, the joint $(\tilde{q}+2) \times(\tilde{q}+2)$ approximate covariance matrix of $\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \tilde{R}^{\prime} u_{1: T}$ and $u_{T+\lfloor h T\rfloor}$ may be computed via

$$
\tilde{W}^{e \prime} \tilde{\Sigma}_{u}^{e}\left(\theta^{u}\right) \tilde{W}^{e} \text { with } \tilde{W}^{e \prime}=\left(\begin{array}{r}
\left(\tilde{R}^{\prime} \tilde{R}\right)^{-1} \tilde{R}^{\prime}  \tag{12}\\
0 \\
0
\end{array} \cdots \frac{\cdots}{0}\right)
$$

where $\tilde{W}^{e}$ is $(T+\lfloor h T\rfloor) \times(\tilde{q}+2)$ and $\tilde{\Sigma}_{u}^{e}\left(\theta^{u}\right)$ is the $(T+\lfloor h T\rfloor) \times(T+\lfloor h T\rfloor)$ covariance matrix of a $\left(\zeta, \sqrt{1-\zeta^{2}}\right)$ weighted average of two independent stationary $\operatorname{AR}(1)$ processes with coefficients $\rho_{1}$ and $\rho_{2}$ and innovation variance equal to $\left(1-\rho_{1}^{2}\right)^{-1 / 2}$ and $\left(1-\rho_{2}^{2}\right)^{-1 / 2}$, respectively. The same approximation readily extends to the forecast of $f_{T+\lfloor h T\rfloor}=\sum_{t=1}^{T+\lfloor h T\rfloor}\left(m_{s}+\Delta a_{s}\right)$.

## Appendix 2.3 Gibbs Sampler

The posterior results in the paper are computed from a (thinned) version of 150,000 Gibbs draws, after discarding 25,000 for burn-in. Given our choice of prior, each step of the Gibbs sampler is conjugate. Further, the dimension reduction to $\tilde{q}+1$ latent vectors, and the discrete parameter space allows precomputations of nearly all relevant $(\tilde{q}+1) \times(\tilde{q}+1)$ matrices. Finally, many of the computations within one Gibbs draw can be parallelized, as
they involve $n, 25$ or at least 10 independent components. In combination, this leads to very fast computations: In the baseline specification, the 175,000 draws take about 4 minutes on a dual 12 core workstation in a Fortran implementation.

The tight prior variance $\Delta$ of the dummy observation $Y^{0}$ in (9) prevents large moves of $F$; in order to accelerate convergence to reasonable starting values, we linearly decrease the variance of $Y^{0}$ from $1000 \Delta$ to $\Delta$ during the burn-in phase.

To complete the description of the model, recall that the innovations $\left\{U_{c, i}\right\}_{i=1}^{n},\left\{U_{g, j}\right\}_{j=1}^{25}$ and $\left\{U_{h, k}\right\}_{k=1}^{10}$ are independent with distribution

$$
\begin{align*}
U_{c, i} & \sim \mathcal{N}\left(0, \omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)\right)  \tag{13}\\
U_{g, j} & \sim \mathcal{N}\left(0, \omega^{2} \kappa_{g, j}^{2}\left(1-\lambda_{g, j}^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)\right)  \tag{14}\\
U_{h, k} & \sim \mathcal{N}\left(0, \omega^{2} \kappa_{h, k}^{2} \Sigma_{U}\left(\theta_{h, k}^{u}\right)\right) \tag{15}
\end{align*}
$$

where

- $\left\{\lambda_{c, i}\right\}_{i=1}^{n}$ and $\left\{\lambda_{g, j}\right\}_{j=1}^{25}$ are i.i.d. draws from the discrete distribution with 25 support points $\left\{\lambda^{l}\right\}_{l=1}^{25}$ and priors $p_{c}^{\lambda}=\left(p_{c, 1}^{\lambda}, \ldots, p_{c, 25}^{\lambda}\right)$ and $p_{g}^{\lambda}=\left(p_{g, 1}^{\lambda}, \ldots, p_{g, 25}^{\lambda}\right)$, respectively
- $\left\{\theta_{c, i}^{u}\right\}_{i=1}^{n},\left\{\theta_{g, j}^{u}\right\}_{j=1}^{25}$ and $\left\{\theta_{h, k}^{u}\right\}_{k=1}^{10}$ are i.i.d. draws from the discrete distribution with 100 support points $\left\{\theta^{l}\right\}_{l=1}^{100}$ and priors $p_{c}^{u}=\left(p_{c, 1}^{u}, \ldots, p_{c, 100}^{u}\right), p_{g}^{u}=\left(p_{g, 1}^{u}, \ldots, p_{g, 100}^{u}\right)$ and $p_{h}^{u}=\left(p_{h, 1}^{u}, \ldots, p_{h, 100}^{u}\right)$, respectively
- $\left\{\kappa_{c, i}\right\}_{i=1}^{n},\left\{\kappa_{g, j}\right\}_{j=1}^{25}$, and $\left\{\kappa_{h, k}\right\}_{k=1}^{10}$ are i.i.d. draws from the discrete distribution with 25 support points $\left\{\kappa^{l}\right\}_{l=1}^{25}$ and priors $p_{c}^{\kappa}=\left(p_{c, 1}^{\kappa}, \ldots, p_{c, 25}^{\kappa}\right), p_{g}^{\kappa}=\left(p_{g, 1}^{\kappa}, \ldots, p_{g, 25}^{\kappa}\right)$ and $p_{h}^{\kappa}=\left(p_{h, 1}^{\kappa}, \ldots, p_{h, 25}^{\kappa}\right)$, respectively.
and the various $p$ are realizations of the corresponding Dirichlet prior.
The state of the sampler thus consists of $\left\{X_{i}\right\}_{i=1}^{n} \in \mathbb{R}^{n(\tilde{q}+1)},\left(F, S_{m}\right) \in \mathbb{R}^{2(\tilde{q}+1)},\left\{G_{j}\right\}_{j=1}^{25} \in$ $\mathbb{R}^{25(\tilde{q}+1)},\left\{H_{k}\right\}_{k=1}^{10} \in \mathbb{R}^{10(\tilde{q}+1)}, \mu_{c} \in \mathbb{R},\left\{\lambda_{c, i}\right\}_{i=1}^{n} \in \mathbb{R}^{n},\left\{\lambda_{g, j}\right\}_{j=1}^{25} \in \mathbb{R}^{25},\left\{\theta_{c, i}^{u}\right\}_{i=1}^{n} \in \mathbb{R}^{3 n}$, $\left\{\theta_{g, j}^{u}\right\}_{j=1}^{25} \in \mathbb{R}^{75},\left\{\theta_{h, k}^{u}\right\}_{k=1}^{10} \in \mathbb{R}^{30}, \omega^{2} \in \mathbb{R},\left\{\kappa_{c, i}\right\}_{i=1}^{n} \in \mathbb{R}^{n},\left\{\kappa_{g, j}\right\}_{j=1}^{25} \in \mathbb{R}^{25},\left\{\kappa_{h, k}\right\}_{k=1}^{10} \in \mathbb{R}^{10}$, $\{J(i)\}_{i=1}^{n} \in \mathbb{N}^{n},\{K(j)\}_{i=1}^{25} \in \mathbb{N}^{25},\left(f_{0}, \mu_{m}\right) \in \mathbb{R}^{2}, \sigma_{\Delta a}^{2} \in \mathbb{R}, \sigma_{m} \in \mathbb{R}, \rho_{m} \in \mathbb{R}, p_{c}^{\lambda} \in \mathbb{S}^{25}$, $p_{g}^{\lambda} \in \mathbb{S}^{25}, p_{c}^{\theta} \in \mathbb{S}^{100}, p_{g}^{\theta} \in \mathbb{S}^{100}, p_{h}^{\theta} \in \mathbb{S}^{100}, p_{c}^{\kappa} \in \mathbb{S}^{25}, p_{g}^{\kappa} \in \mathbb{S}^{25}, p_{h}^{\kappa} \in \mathbb{S}^{25}$, where $\mathbb{S}^{m}$ is the simplex of dimension $m$. In each step of the Gibbs sampler described in detail below, one of these states is drawn conditional on the current value of all other states. We checked the
correctness of the sampler (and code) using the Geweke (2004) test as described in Müller and Watson (2019), who also provide additional details on the conjugate distributions exploited in the Gibbs steps below. The long-run forecasts are easily generated conditional on the full state of the sampler from the conditional normal distribution with parameters implied by (12) and the analogous expression for the forecast of $f_{t}$.

One draw of the Gibbs sampler consists of the following 28 steps.

1. $\left\{X_{i}\right\}_{i=1}^{n}$

From (4), (5) and (9), drawing $\left\{X_{i}\right\}_{i=1}^{n}$ amounts to drawing from the conditional distribution of $\left\{C_{i}\right\}_{i=1}^{n}$ given observations $\left\{B_{i}^{\prime} C_{i}\right\}_{i=1}^{n}$ and $Y^{0}=\sum_{i=1}^{n} w_{i} C_{i}$, where $C_{i}=$ $X_{i}-F$. Let $\mu_{C}=\left(\mu_{c} \iota_{1}^{\prime}+\lambda_{c, 1} G_{J(1)}^{\prime}, \ldots, \mu_{c} \iota_{1}^{\prime}+\lambda_{c, n} G_{J(n)}^{\prime}\right)^{\prime} \in \mathbb{R}^{n(\tilde{q}+1)}$ be the conditional mean of $C=\left(C_{1}^{\prime}, \ldots, C_{n}^{\prime}\right)^{\prime}, \Sigma$ the block diagonal conditional covariance matrix of $U_{c}$ with $i$ th block equal to $\omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right), B$ the $n(\tilde{q}+1) \times \sum_{i=1}^{n}\left(q_{i}+1\right)$ matrix such that $Y=\left(Y_{1}^{\prime}, \ldots, Y_{n}^{\prime}\right)=B^{\prime} X$ with $X=\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)^{\prime}$, and $w=\left(w_{1}, \ldots, w_{n}\right)^{\prime} \otimes I_{\tilde{q}+1} \in$ $\mathbb{R}^{n(\tilde{q}+1) \times(\tilde{q}+1)}$. The conditional joint distribution of $C, Y_{0}$ and $B^{\prime} C$ is then given by

$$
\left(\begin{array}{c}
C \\
Y^{0} \\
B^{\prime} C
\end{array}\right) \sim \mathcal{N}\left(\left(\begin{array}{c}
\mu_{C} \\
w^{\prime} \mu_{C} \\
B^{\prime} \mu_{C}
\end{array}\right),\left(\begin{array}{ccc}
\Sigma & \cdot & \cdot \\
w^{\prime} \Sigma & w^{\prime} \Sigma w+\Delta & \cdot \\
B^{\prime} \Sigma & B^{\prime} \Sigma w & B^{\prime} \Sigma B
\end{array}\right)\right)
$$

so that with $m=\mu_{C}+\Sigma B\left(B^{\prime} \Sigma B\right)^{-1} B^{\prime}\left(C-\mu_{C}\right)$ and $V=\Sigma-\Sigma B\left(B^{\prime} \Sigma B\right)^{-1} B^{\prime} \Sigma$,

$$
\binom{C}{Y^{0}} \left\lvert\, B^{\prime} C \sim \mathcal{N}\left(\binom{m}{w^{\prime} m},\left(\begin{array}{cc}
V & \cdot \\
w^{\prime} V & w^{\prime} V w+\Delta
\end{array}\right)\right) .\right.
$$

Thus,

$$
\begin{equation*}
C \mid\left(B^{\prime} C, Y^{0}\right) \sim \mathcal{N}\left(m+V w\left(w^{\prime} V w+\Delta\right)^{-1}\left(Y^{0}-w^{\prime} m\right), V-V^{\prime} w\left(w^{\prime} V w+\Delta\right)^{-1} w^{\prime} V\right) . \tag{16}
\end{equation*}
$$

In order to avoid manipulating the large matrices in this expression directly, we generate a draw from this distribution as follows: Draw $Z_{0} \sim \mathcal{N}\left(\mu_{C}, \Sigma\right)$, and note that

$$
Z_{1}=Z_{0}-\Sigma B\left(B^{\prime} \Sigma B\right)^{-1} B^{\prime}\left(Z_{0}-C\right) \sim \mathcal{N}(m, V)
$$

The draw of $Z_{0}$ and the computation of $Z_{1}$ can be implemented independently for each $\tilde{q}+1$ block, and each $\tilde{q}+1 \times \tilde{q}+1$ block of $B\left(B^{\prime} \Sigma B\right)^{-1} B^{\prime}$ can be precomputed (up to scale). Then draw $\varepsilon \sim \mathcal{N}\left(Y^{0}, \Delta\right)$ independent of $Z_{0}$, and note that

$$
Z_{1}-V w\left(w^{\prime} V w+\Delta\right)^{-1} w^{\prime}\left(Z_{1}-\varepsilon\right) \sim \mathcal{N}\left(m+V w\left(w^{\prime} V w+\Delta\right)^{-1}\left(Y^{0}-w^{\prime} m\right), V-V^{\prime} w\left(w^{\prime} V w+\Delta\right)^{-1} w^{\prime} V\right)
$$

which is a draw from the target (16), and the evaluation of $V w$ and $w^{\prime} V w$ simply consists of computing weighted sums of precomputed $\tilde{q}+1 \times \tilde{q}+1$ blocks.
2. $\left(F, S_{m}\right)$

Let $\mu_{F}=\iota_{1} f_{0}+\iota_{2} \mu_{m}, \Sigma_{F}=\sigma_{m}^{2} \Sigma_{m}\left(\rho_{m}\right)+\sigma_{\Delta a}^{2} \Sigma_{a}$ and $e=(1, \ldots, 1)^{\prime} \otimes I_{\tilde{q}+1} \in$ $\mathbb{R}^{n(\tilde{q}+1) \times(\tilde{q}+1)}$. Then from the independent Gaussian kernels (4) and (9), we obtain that the conditional posterior distribution of $F$ is

$$
F-\mu_{F} \sim \mathcal{N}\left(m_{F}, V_{F}\right)
$$

where $V_{F}=\left(\Sigma_{F}^{-1}+e^{\prime} \Sigma^{-1} e+\Delta^{-1}\right)^{-1}$ and $m_{F}=V_{F}\left(e^{\prime} \Sigma^{-1}\left(X-\mu_{C}-e \mu_{F}\right)+\Delta^{-1}\left(w^{\prime} X-\right.\right.$ $\left.Y^{0}-\mu_{F}\right)$ ). Furthermore, from (8), conditional on a draw of $F$ from this conditional distribution, the conditional distribution of $S_{m}$ satisfies

$$
S_{m} \mid F \sim \mathcal{N}\left(\Sigma_{S} \Sigma_{F}^{-1}\left(F-\mu_{F}\right), \Sigma_{S}-\Sigma_{S} \Sigma_{F}^{-1} \Sigma_{S}\right)
$$

where $\Sigma_{S}=\sigma_{m}^{2} \Sigma_{m}\left(\rho_{m}\right)$.
3. $\left\{G_{j}\right\}_{j=1}^{25}$

Let $\Sigma_{g, j}=\omega^{2} \kappa_{g, j}^{2}\left(1-\lambda_{g, j}^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)$ and $\Sigma_{c, i}=\omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)$. From (5) and (6), we have that the conditional distribution of $\left\{G_{j}\right\}_{j=1}^{25}$ is independent and satisfies

$$
G_{j} \sim \mathcal{N}\left(m_{G_{j}}, V_{G_{j}}\right)
$$

where $V_{G_{j}}=\left(\Sigma_{g, j}^{-1}+\sum_{i=1}^{n} \mathbf{1}[J(i)=j] \lambda_{c, i}^{2} \Sigma_{c, i}^{-1}\right)^{-1}$ and $m_{G_{j}}=V_{G_{j}}\left(\lambda_{g, j} \Sigma_{g, j}^{-1} H_{K(j)}+\right.$ $\left.\sum_{i=1}^{n} \mathbf{1}[J(i)=j] \lambda_{c, i} \Sigma_{c, i}^{-1}\left(C_{i}-\mu_{c} \iota_{1}\right)\right)$ with $C_{i}=X_{i}-F$.
4. $\left\{H_{k}\right\}_{k=1}^{10}$

Let $\Sigma_{h, k}=\omega^{2} \kappa_{h, k}^{2} \Sigma_{U}\left(\theta_{h, k}^{u}\right)$ and $\Sigma_{g, j}=\omega^{2} \kappa_{g, j}^{2}\left(1-\lambda_{g, j}^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)$. From (6) and (7), we have that the conditional distribution of $\left\{H_{k}\right\}_{j=1}^{25}$ is independent and satisfies

$$
H_{k} \sim \mathcal{N}\left(m_{H_{k}}, V_{H_{k}}\right)
$$

where $V_{H_{k}}=\left(\Sigma_{h, k}^{-1}+\sum_{j=1}^{25} \mathbf{1}[K(j)=k] \lambda_{g, j}^{2} \Sigma_{g, j}^{-1}\right)^{-1}$ and $m_{H_{k}}=V_{H_{k}}\left(\sum_{j=1}^{25} \mathbf{1}[K(j)=\right.$ $\left.k] \lambda_{g, j} \Sigma_{g, j}^{-1} G_{j}\right)$.
5. $\mu_{c}$

Under a flat prior for $\mu_{c}$, the conditional posterior distribution is

$$
\mu_{c} \sim \mathcal{N}\left(\frac{\sum_{i=1}^{n} \iota_{1}^{\prime} \Sigma_{c, i}^{-1}\left(C_{i}-\lambda_{c, i} G_{J(i)}\right)}{\sum_{i=1}^{n} \iota_{1}^{\prime} \Sigma_{c, i}^{-1} \iota_{1}}, \frac{1}{\sum_{i=1}^{n} \iota_{1}^{\prime} \Sigma_{c, i}^{-1} \iota_{1}}\right)
$$

where $\Sigma_{c, i}=\omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)$, and $C_{i}=X_{i}-F$.
6. $\left\{\lambda_{c, i}\right\}_{i=1}^{n}$

From (5) and (13), the conditional posterior weight of $\lambda_{c, i}$ on $\lambda^{l}, l=1, \ldots, 25$ is proportional to
$p_{c, l}^{\lambda} \exp \left[-\frac{1}{2}\left(C_{i}-\mu_{c} \iota_{1}-\lambda^{l} G_{J(i)}\right)^{\prime}\left(\omega^{2} \kappa_{c, i}^{2}\left(1-\left(\lambda^{l}\right)^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)\right)^{-1}\left(C_{i}-\mu_{c} \iota_{1}-\lambda^{l} G_{J(i)}\right)\right]\left(1-\left(\lambda^{l}\right)^{2}\right)^{-(\tilde{q}+1) / 2}$ with $C_{i}=X_{i}-F$.
7. $\left\{\lambda_{g, j}\right\}_{j=1}^{25}$

From (6) and (14), the conditional posterior weight of $\lambda_{g, j}$ on $\lambda^{l}, l=1, \ldots, 25$ is proportional to
$p_{g, l}^{\lambda} \exp \left[-\frac{1}{2}\left(G_{j}-\lambda^{l} H_{K(j)}\right)^{\prime}\left(\omega^{2} \kappa_{g, j}^{2}\left(1-\left(\lambda^{l}\right)^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)\right)^{-1}\left(G_{j}-\lambda^{l} H_{K(j)}\right)\right]\left(1-\left(\lambda^{l}\right)^{2}\right)^{-(\tilde{q}+1) / 2}$.
8. $\left\{\theta_{c, i}^{u}\right\}_{i=1}^{n}$

From (13), the conditional posterior weight of $\theta_{c, i}^{u}$ on $\theta^{l}, l=1, \ldots, 100$ is proportional to

$$
p_{c, l}^{\theta}\left|\operatorname{det} \Sigma_{U}\left(\theta^{l}\right)\right|^{-1 / 2} \exp \left[-\frac{1}{2} \frac{U_{c,,}^{\prime} \Sigma_{U}\left(\theta^{l}\right)^{-1} U_{c, i}}{\omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right)}\right]
$$

with $U_{c, i}=X_{i}-F-\mu_{c} \iota_{1}-\lambda_{c, i} G_{J(i)}$.
9. $\left\{\theta_{g, j}^{u}\right\}_{j=1}^{25}$

From (14), the conditional posterior weight of $\theta_{g, j}^{u}$ on $\theta^{l}, l=1, \ldots, 100$ is proportional to

$$
p_{g, l}^{\theta}\left|\operatorname{det} \Sigma_{U}\left(\theta^{l}\right)\right|^{-1 / 2} \exp \left[-\frac{1}{2} \frac{U_{g, j}^{\prime} \Sigma_{U}\left(\theta^{l}\right)^{-1} U_{g, j}}{\omega^{2} \kappa_{g, j}^{2}\left(1-\lambda_{g, j}^{2}\right)}\right]
$$

with $U_{g, j}=G_{j}-\lambda_{g, j} H_{K(j)}$.
10. $\left\{\theta_{h, k}^{u}\right\}_{k=1}^{10}$

From (7) and (15), the conditional posterior weight of $\theta_{h, k}^{u}$ on $\theta^{l}, l=1, \ldots, 100$ is proportional to

$$
p_{h, l}^{\theta}\left|\operatorname{det} \Sigma_{U}\left(\theta^{l}\right)\right|^{-1 / 2} \exp \left[-\frac{1}{2} \frac{H_{k}^{\prime} \Sigma_{U}\left(\theta^{l}\right)^{-1} H_{k}}{\omega^{2} \kappa_{h, k}^{2}}\right] .
$$

11. $\omega^{2}$

Under the conjugate inverse gamma prior with unit median, $(1 / 2.198) / \omega^{2} \sim \chi_{1}^{2}$. From (13)-(15), the conditional posterior distribution for $\omega^{2}$ then is again inverse-gamma and satisfies

$$
\frac{1 / 2.198+S_{c}^{2}+S_{g}^{2}+S_{h}^{2}}{\omega^{2}} \sim \chi_{1+(n+10+25)(\tilde{q}+1)}^{2}
$$

where $S_{c}^{2}=\sum_{i=1}^{n} U_{c, i}^{\prime}\left(\kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)\right)^{-1} U_{c, i}, \quad S_{g}^{2}=\sum_{j=1}^{25} U_{g, j}^{\prime}\left(\kappa_{g, j}^{2}(1-\right.$ $\left.\left.\lambda_{g, j}^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)\right)^{-1} U_{g, j}$ and $S_{h}^{2}=\sum_{k=1}^{10} U_{h, k}^{\prime}\left(\kappa_{h, k}^{2} \Sigma_{U}\left(\theta_{h, k}^{u}\right)\right)^{-1} U_{h, k}$ with $U_{c, i}=X_{i}-F-$ $\mu_{c} \iota_{1}-\lambda_{c, i} G_{J(i)}, U_{g, j}=G_{j}-\lambda_{g, j} H_{K(j)}$ and $U_{h, k}=H_{k}$.
12. $\left\{\kappa_{c, i}\right\}_{i=1}^{n}$

From (5) and (13), the conditional posterior weight of $\kappa_{c, i}$ on $\kappa^{l}, l=1, \ldots, 25$ is proportional to

$$
p_{c, l}^{\kappa} \exp \left[-\frac{1}{2} \frac{U_{c, i}^{\prime} \Sigma_{U}\left(\theta_{c, i}^{u}\right)^{-1} U_{c, i}}{\omega^{2}\left(\kappa^{l}\right)^{2}\left(1-\lambda_{c, i}^{2}\right)}\right]\left(\kappa^{l}\right)^{-(\tilde{q}+1)}
$$

with $U_{c, i}=X_{i}-F-\mu_{c} \iota_{1}-\lambda_{c, i} G_{J(i)}$.
13. $\left\{\kappa_{g, j}\right\}_{j=1}^{25}$

From (6) and (14), the conditional posterior weight of $\kappa_{g, j}$ on $\kappa^{l}, l=1, \ldots, 25$ is proportional to

$$
p_{g, l}^{\kappa} \exp \left[-\frac{1}{2} \frac{U_{g, j}^{\prime} \Sigma_{U}\left(\theta_{g, j}^{u}\right)^{-1} U_{g, j}}{\omega^{2}\left(\kappa^{l}\right)^{2}\left(1-\lambda_{g, j}^{2}\right)}\right]\left(\kappa^{l}\right)^{-(\tilde{q}+1)}
$$

with $U_{g, j}=G_{j}-\lambda_{g, j} H_{K(j)}$.
14. $\left\{\kappa_{h, k}\right\}_{k=1}^{10}$

From (7) and (15), the conditional posterior weight of $\kappa_{h, k}$ on $\kappa^{l}, l=1, \ldots, 25$ is proportional to

$$
p_{h, l}^{\kappa} \exp \left[-\frac{1}{2} \frac{H_{k}^{\prime} \Sigma_{U}\left(\theta_{h, k}^{u}\right)^{-1} H_{k}}{\omega^{2}\left(\kappa^{l}\right)^{2}}\right]\left(\kappa^{l}\right)^{-(\tilde{q}+1)} .
$$

15. $\{J(i)\}_{i=1}^{n}$

From (5), the conditional posterior weight of $J(i)$ on $j=1, \ldots, 25$ is proportional to

$$
\exp \left[-\frac{1}{2}\left(C_{i}-\mu_{c} \iota_{1}-\lambda_{c, i} G_{j}\right)^{\prime}\left(\omega^{2} \kappa_{c, i}^{2}\left(1-\lambda_{c, i}^{2}\right) \Sigma_{U}\left(\theta_{c, i}^{u}\right)\right)^{-1}\left(C_{i}-\mu_{c} \iota_{1}-\lambda_{c, i} G_{j}\right)\right] .
$$

16. $\{K(j)\}_{i=1}^{25}$

From (6), the conditional posterior weight of $K(j)$ on $k=1, \ldots, 10$ is proportional to

$$
\exp \left[-\frac{1}{2}\left(G_{j}-\lambda_{g, j} H_{k}\right)^{\prime}\left(\omega^{2} \kappa_{g, j}^{2}\left(1-\lambda_{g, j}^{2}\right) \Sigma_{U}\left(\theta_{g, j}^{u}\right)\right)^{-1}\left(G_{j}-\lambda_{g, j} H_{k}\right)\right]
$$

17. $\left(f_{0}, \mu_{m}\right)$

Let $\iota_{1: 2}$ be the first two columns of $I_{\tilde{q}+1}$, so that $\iota_{1} f_{0}+\iota_{2} \mu_{m}=\iota_{1: 2}\left(f_{0}, \mu_{m}\right)^{\prime}$. Then from (8), the conditional posterior of $\left(f_{0}, \mu_{m}\right)$ under the improper flat prior is

$$
\left(f_{0}, \mu_{m}\right)^{\prime} \sim \mathcal{N}\left(m_{f}, V_{f}\right)
$$

where $V_{f}=\left(\iota_{1: 2}^{\prime} \Sigma_{F}^{-1} \iota_{1: 2}\right)^{-1}$ and $m_{f}=V_{f} \iota_{1: 2}^{\prime} \Sigma_{F}^{-1} F$, with $\Sigma_{F}=\sigma_{m}^{2} \Sigma_{m}\left(\rho_{m}\right)+\sigma_{\Delta a}^{2} \Sigma_{a}$.
18. $\sigma_{\Delta a}^{2}$

Under the conjugate inverse gamma prior with median $0.03^{2},\left(0.03^{2} / 2.198\right) / \sigma_{\Delta a}^{2} \sim \chi_{1}^{2}$.
From (11) the posterior conditional distribution then satisfies

$$
\frac{\left(0.03^{2} / 2.198\right)+A^{\prime} \Sigma_{a}^{-1} A}{\sigma_{\Delta a}^{2}} \sim \chi_{1+\tilde{q}+1}^{2}
$$

with $A=F-\iota_{1} f_{0}+\iota_{2} \mu_{m}-S_{m}$.
19. $\sigma_{m}$

From (10), the conditional posterior weight of $\sigma_{m}$ on $\sigma^{l}, l=1, \ldots, 25$ is proportional to

$$
p_{l}^{\sigma_{m}} \exp \left[-\frac{1}{2}\left(\sigma^{l}\right)^{-2} S_{m}^{\prime} \Sigma_{m}\left(\rho_{m}\right)^{-1} S_{m}\right]\left(\sigma^{l}\right)^{-(\tilde{q}+1)}
$$

where $p_{l}^{\sigma_{m}}$ are the triangular prior weights.
20. $\rho_{m}$

From (10), the conditional posterior weight of $\rho_{m}$ on $\rho^{l}, l=1, \ldots, 25$ is proportional to

$$
\exp \left[-\frac{1}{2} S_{m}^{\prime}\left(\sigma_{m}^{2} \Sigma_{m}\left(\rho^{l}\right)\right)^{-1} S_{m}\right]\left|\operatorname{det} \Sigma_{m}\left(\rho^{l}\right)\right|^{-1 / 2}
$$

21. $p_{c}^{\lambda}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{c, l}^{\lambda}\right\}_{l=1}^{25}$ satisfies

$$
\left\{p_{c, l}^{\lambda}\right\}_{l=1}^{25} \sim\left\{W_{l} / \sum_{i=1}^{25} W_{i}\right\}_{l=1}^{25}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{i=1}^{n} \mathbf{1}\left[\lambda_{c, i}=\lambda^{l}\right]+20 / 25$ degrees of freedom.
22. $p_{g}^{\lambda}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{g, l}^{\lambda}\right\}_{l=1}^{25}$ satisfies

$$
\left\{p_{g_{l}}^{\lambda}\right\}_{l=1}^{25} \sim\left\{W_{l} / \sum_{i=1}^{25} W_{i}\right\}_{l=1}^{25}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{j=1}^{25} \mathbf{1}\left[\lambda_{g, j}=\lambda^{l}\right]+20 / 25$ degrees of freedom.
23. $p_{c}^{\theta}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{c, l}^{\theta}\right\}_{l=1}^{100}$ satisfies

$$
\left\{p_{c, l}^{\theta}\right\}_{l=1}^{100} \sim\left\{W_{l} / \sum_{i=1}^{100} W_{i}\right\}_{l=1}^{100}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{i=1}^{n} \mathbf{1}\left[\theta_{c, i}^{u}=\theta^{l}\right]+20 / 100$ degrees of freedom.
24. $p_{g}^{\theta}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{g, l}^{\theta}\right\}_{l=1}^{100}$ satisfies

$$
\left\{p_{g, l}^{\theta}\right\}_{l=1}^{100} \sim\left\{W_{l} / \sum_{i=1}^{100} W_{i}\right\}_{l=1}^{100}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{j=1}^{25} \mathbf{1}\left[\theta_{g, j}^{u}=\theta^{l}\right]+20 / 100$ degrees of freedom.
25. $p_{h}^{\theta}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{h, l}^{\theta}\right\}_{l=1}^{100}$ satisfies

$$
\left\{p_{h, l}^{\theta}\right\}_{l=1}^{100} \sim\left\{W_{l} / \sum_{i=1}^{100} W_{i}\right\}_{l=1}^{100}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{k=1}^{10} \mathbf{1}\left[\theta_{h, k}^{u}=\theta^{l}\right]+20 / 100$ degrees of freedom.
26. $p_{c}^{\kappa}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{c, l}^{\kappa}\right\}_{l=1}^{25}$ satisfies

$$
\left\{p_{c, l}^{\kappa}\right\}_{l=1}^{25} \sim\left\{W_{l} / \sum_{i=1}^{25} W_{i}\right\}_{l=1}^{25}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{i=1}^{n} \mathbf{1}\left[\kappa_{c, i}=\kappa^{l}\right]+20 / 25$ degrees of freedom.
27. $p_{g}^{\kappa}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{g, l}^{\kappa}\right\}_{l=1}^{25}$ satisfies

$$
\left\{p_{g, l}^{\kappa}\right\}_{l=1}^{25} \sim\left\{W_{l} / \sum_{i=1}^{25} W_{i}\right\}_{l=1}^{25}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{j=1}^{25} \mathbf{1}\left[\kappa_{g, j}=\kappa^{l}\right]+20 / 25$ degrees of freedom.
28. $p_{h}^{\kappa}$

From the conjugate nature of the Dirichlet-multinomial prior, the conditional posterior for $\left\{p_{h, l}^{\kappa}\right\}_{l=1}^{25}$ satisfies

$$
\left\{p_{h, l}^{\kappa}\right\}_{l=1}^{25} \sim\left\{W_{l} / \sum_{i=1}^{25} W_{i}\right\}_{l=1}^{25}
$$

where $W_{l}$ are independent $\chi^{2}$ random variables with $\sum_{k=1}^{10} \mathbf{1}\left[\kappa_{h, k}=\kappa^{l}\right]+20 / 25$ degrees of freedom.

## Appendix 3 Additional Results

- Country-specific results for the posterior summarized in Table 4 are given in the file Table4_Results_by_Country.xlsx.
- A summary of the 50- and 100-year ahead predictive distributions by country are given in the file Predictive_Distributions_by_Country.xlsx..
- Figure 10: countries are listed in order in the file Figure_10_Countries.xlsx.
- Figure 11: by country are collected in the file Figure11_by_Country.zip
- Figure 14: by country are collected in the file Figure14_by_Country.zip.


## Additional References

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