

A HIERARCHICAL ORNSTEIN–UHLENBECK MODEL FOR CONTINUOUS REPEATED MEASUREMENT DATA

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In this paper, we present a diffusion model for the analysis of continuous-time change in multivariate longitudinal data. The central idea is to model the data from a single person with an Ornstein–Uhlenbeck diffusion process. We extend it hierarchically by allowing the parameters of the diffusion process to vary randomly over different persons. With this approach, both intra and interindividual differences are analyzed simultaneously. Furthermore, the individual difference parameters can be regressed on covariates, thereby providing an explanation of between-person differences. Unstructured and unbalanced data pose no problem for the model to be applied. We demonstrate the method on data from an experience sampling study to investigate changes in the core affect. It can be concluded that different factors from the five factor model of personality are related to features of the trajectories in the core affect space, such as the cross-correlation and variability of the changes.

Key words: Ornstein–Uhlenbeck process, Hierarchical, Multivariate, Bayesian.

1. Introduction

Change over time is a central and nonnegligible concept that psychologists frequently encounter while studying different phenomena such as learning processes, developmental issues, or mood changes. Emotions and related phenomena are prime examples that are subject to change over time, but also measures that are intuitively believed to be very stable (such as personality characteristics, see, e.g., Borkenau & Ostendorf, 1998) reveal their changing nature when repeated measures are taken. Going a step further, following the arguments of van Montfort, Oud, and Satorra (2007), any serious causal analyses should be based on longitudinal data.

Many longitudinal studies rely on panel designs: a large number of subjects with typically fewer than 10 measurements per person. However, technical innovations such as palmtops, beepers, and online questionnaires have made it possible to measure variables of interest more “intensively” than in the typical panel designs. Such intensive longitudinal designs (see Walls & Schafer, 2006) have recently become quite popular. They typically consists of relatively long (e.g., more than 50 measurement occasions) data chains from different subjects. Intensive longitudinal data frequently stem from experience sampling techniques (Bolger, Davis, & Rafaeli, 2003; Csikszentmihalyi & Larson, 1987; Larson & Csikszentmihalyi, 1983; Russell & Feldman-Barrett, 1999), dyadic interaction studies (Ferrer & Nesselrode, 2003), or cognitive and sensorimotor performance research (Li, Huxhold, & Schmiedek, 2004). These methods often result in a highly unstructured longitudinal dataset, since the records may be taken at different time points for different individuals. Moreover, missing data often occur, which leads to unequal numbers of observations.

Apart from a few exceptions (e.g., Oud & Singer, 2008; Singer, 2008), time is generally handled in a discrete way in longitudinal models. This property makes modeling such data somewhat unrealistic, since the measured phenomena do not cease to exist between observations. Oud (2002) argues that most processes in behavioral sciences unfold in continuous time and should

be handled accordingly. Treating a theoretically continuous variable as discrete may lead to biased results, for a discussion see Delsing, Oud, and Bruyn (2005). In contrast, continuous-time modeling appears to be a more realistic solution.

In this paper, we focus on introducing a hierarchical model for analyzing intensively measured variables while allowing for change in continuous time. This way, intra and interindividual differences are studied simultaneously. The model incorporates two distinctive properties. First of all, it concentrates on the dynamic feature of the change process by investigating a mean-reverting tendency over time. Second, it explores interindividual variability from different perspectives, some of which have not been considered so far. Both aspects are summarized below.

Concerning the first special property, the paper introduces a stochastic process with Markovian properties, namely the Ornstein–Uhlenbeck (OU) process, to serve as the basic model for change within an individual. The OU process can be seen as a mathematical model of temporal change for phenomena with regulatory, mean-reverting, or centralizing mechanisms. This property makes the process especially useful for modeling moods and emotions (Larsen, 2000; Lykken & Tellegen, 1996), but it can also be applied to other change processes with a possible regulatory mechanism (e.g., balance control). When we use the OU process as an analytical tool, the focus is more on the dynamics of the process, and not so much on the systematic or structural changes with respect to the mean level, as is commonly investigated by mixed models (e.g., Diggle, Heagerty, Liang, & Zeger, 2002; Verbeke & Molenberghs, 2000). This way, our approach is closer to the area of time series analysis, but a distinctive aspect is that our main emphasis lies on studying interindividual differences in temporal change, while time series analysis mainly focuses on a single measurement chain.

Regarding the second characteristic, the proposed hierarchical model is suited for exploring interindividual differences from aspects which have been neglected so far. Other existing techniques, like structural equation modeling (SEM, Bollen, 1989), multilevel modeling (Goldstein, 2003; Raudenbush & Bryk, 2002), state-space modeling either combined with SEM or with Kalman filter estimation procedure (Oud, 2007; Oud & Singer, 2008) do not commonly deal with interindividual differences with respect to all model parameters. In the hierarchical model presented here, parameters like the serial- or the cross-correlation are considered to be person-specific, and in this way we allow new aspects of interindividual differences to be investigated.

Regarding statistical inference, we will make use of a Bayesian approach (Gelman, Carlin, Stern, & Rubin, 2004). The Bayesian methodology offers a sound way for statistical inference in models with a complex hierarchical structure. With respect to the present model, since all the OU parameters can be turned into random effects and can be regressed onto predictors or covariates (De Boeck & Wilson, 2004), the parameter estimation in the classical framework would involve a high-dimensional integration over the numerous random effect distributions. In contrast, the paper will demonstrate that the hierarchical OU model can be fitted in a straightforward way using the Bayesian framework. Also, the most commonly used statistical inference technique in continuous time modeling involves some approximation methods to estimate the parameters, while with the Bayesian approach this is no longer necessary.

The OU process in particular and some variants of it (e.g., the integrated OU process) have been proposed as models for the analysis of longitudinal profiles in several fields. For example, single time series from measurements of animal movement have been modeled as an OU process by Blackwell (1997, 2003), Brillinger, Preisler, Ager, and Kie (2004) and Dunn and Gipson (1977). A major difference between our psychological approach and this biological application is that in the latter there is no interest in describing and explaining differences between subjects. In the context of mixed models for longitudinal data analysis, the integrated OU process is proposed to model serial correlation between measurements (e.g., see De la Cruz-Mesía & Marshall, 2006; Sy, Taylor, & Cumberland, 1997; Taylor, Cumberland, & Sy, 1994), but no interindividual variation is allowed in the driving parameters of the process.

The structure for the remainder of the article is as follows. In the next section, we explain the OU diffusion process together with the interpretation of its parameters. Subsequently, we discuss a hierarchical extension. The following section summarizes the statistical inference and afterward an application to the modeling of core affect trajectories is presented. The last part presents our conclusion.

2. The Theory of the Ornstein–Uhlenbeck Diffusion Process

In this section, we give a nontechnical and self-contained account of the OU diffusion process. More detailed explanations can be found in Cox & Miller (1972), Dunn & Gipson (1977), Karlin & Taylor (1981), and Blackwell (2003). Usually, the model is presented as a solution to a first-order stochastic differential equation (SDE). However, since the details of the SDE may not be generally known, we will start with the solution and introduce the SDE only later. At this point, we do not yet present an application, so our explanation will be in general terms without reference to a substantive area.

Let us assume that the state of an individual at time t ($t \geq 0$) can be represented as a point $\mathbf{Y}(t) = (Y_1(t), Y_2(t), \dots, Y_q(t))^T$ in a continuous, q -dimensional space. In general, an OU diffusion process is a continuous-time Gaussian process $\{\mathbf{Y}(t) : t \geq 0\}$ defined on this q -dimensional space such that given that the process was in state $\mathbf{Y}(t)$ at time t , the conditional distribution of the position $\mathbf{Y}(t+d)$ d time units later is

$$\mathbf{Y}(t+d) | \mathbf{Y}(t) \sim N_q(\boldsymbol{\mu} + e^{-\mathbf{B}d}(\mathbf{Y}(t) - \boldsymbol{\mu}), \boldsymbol{\Gamma} - e^{-\mathbf{B}d}\boldsymbol{\Gamma}e^{-\mathbf{B}'d}), \quad (1)$$

where $\boldsymbol{\mu}$ is a q -dimensional vector and \mathbf{B} and $\boldsymbol{\Gamma}$ are $q \times q$ matrices. The function e^M (with M a square matrix) is the matrix exponential. The matrix exponential e^M is defined as follows:

$$e^M = I + \sum_{j=1}^{\infty} \frac{M^j}{j!}.$$

As a special case (see later), if M is a diagonal matrix with diagonal elements m_1, \dots, m_q , it is equal to

$$e^M = \begin{pmatrix} e^{m_1} & \dots & 0 \\ \vdots & \ddots & \vdots \\ 0 & \dots & e^{m_q} \end{pmatrix},$$

where e^{m_i} is the scalar exponential function value of m_i .

If \mathbf{B} is a positive definite matrix (and if $E(\mathbf{Y}(0)) = \boldsymbol{\mu}$ and $\text{var}(\mathbf{Y}(0)) = \boldsymbol{\Gamma}$), then the Ornstein–Uhlenbeck process is stationary. This can be seen intuitively by letting $d \rightarrow \infty$, so that the matrix exponential part $e^{-\mathbf{B}d}$ goes to zero and the process has the following equilibrium or stationary distribution:

$$\mathbf{Y}(t) \sim N_q(\boldsymbol{\mu}, \boldsymbol{\Gamma}). \quad (2)$$

From (2), it can be deduced that $\boldsymbol{\mu}$ is the mean of the equilibrium distribution and $\boldsymbol{\Gamma}$ is its covariance matrix (hence, $\boldsymbol{\Gamma}$ is positive definite). The assumption of stationarity implies that if the process runs for an infinitely long period of time, this equilibrium density is the density function of the visited points in the q -dimensional space. It will also be assumed in this paper that the distribution of $\mathbf{Y}(0)$ (the first measurement, when time t equals 0) is the equilibrium distribution (in which case the model is strictly stationary).

To explain the interpretation of the distinct parameter vectors and matrices in the model, it is best to start by examining the conditional mean vector $\boldsymbol{\mu} + e^{-\mathbf{B}d}(\mathbf{Y}(t) - \boldsymbol{\mu})$ from (1). It can be seen that this conditional mean depends on the previous position $\mathbf{Y}(t)$, on the time difference d between the two measurements and on the $\boldsymbol{\mu}$ and \mathbf{B} parameters of the process. The interpretation of $\boldsymbol{\mu}$ and \mathbf{B} is the following: The parameter $\boldsymbol{\mu}$ is the vector of expected values of the equilibrium distribution (let $d \rightarrow \infty$) and can thus be seen as a fixed point attractor in the q -dimensional space. Because of this property, $\boldsymbol{\mu}$ can be called the average position or the “homebase” of the process. The other parameter, the matrix \mathbf{B} controls the strength of the centralizing tendency, which keeps the process in the vicinity of the homebase. This matrix represents a mean reverting or dampening force, since it impedes the process to diffuse away from the homebase. To illustrate the role of the centralizing tendency, it is easier to see its function in a one-dimensional model (i.e., $q = 1$). In that case, the conditional mean of $Y(t + d)$ given $Y(t)$ is equal to $\mu + e^{-\beta d}(Y(t) - \mu)$ (with, as required, $\beta > 0$). From the latter equation, it can be easily derived that (1) if β is large, the conditional mean is very close to the homebase and (2) if β approaches zero, the homebase takes a value close to the previous position.

To simplify the interpretation in the general q -dimensional case, in this paper we only deal with the subset of isotropic \mathbf{B} matrices, that is, $\mathbf{B} = \beta \mathbf{I}$, where \mathbf{I} is the $q \times q$ identity matrix. The isotropic restriction is motivated by two main reasons. The first is a pragmatic one since making the matrix \mathbf{B} isotropic reduces the complexity of the model substantially. It is not unusual in the applied literature on the OU process to use such a constraint (see, e.g., Blackwell, 2003). For a general \mathbf{B} , not only does just \mathbf{B} have to be positive definite but the matrix $\mathbf{B}\boldsymbol{\Gamma} + \boldsymbol{\Gamma}\mathbf{B}'$ as well (Dunn & Gipson, 1977). Satisfying this constraint in the estimation process is quite cumbersome. A second and more substantive reason to prefer the isotropic parametrization is that it does not give any special importance to the chosen coordinate system. If the matrix \mathbf{B} is isotropic, the expected trajectories near the homebase are straight lines and, therefore, the centralizing tendency matrix is invariant under rotation and reflection (see also Blackwell, 1997). Accordingly, the centralizing tendency is controlled only by the distance from the homebase of the process and not by its direction. However, for nonisotropic centralizing tendency matrices, the expected trajectories close to the homebase are generally not straight lines, but are curved. In the latter case, special importance is given to the coordinate system at hand. However, in most psychological applications, the coordinate system is arbitrary (see also the application below).

If \mathbf{B} is restricted to be isotropic, then it holds that $e^{-\mathbf{B}d} = e^{-\beta d} \mathbf{I}$. Hence, if β is large (in which case there is a strong centralizing tendency), the exponential factor in the conditional mean $\boldsymbol{\mu} + e^{-\beta d}(\mathbf{Y}(t) - \boldsymbol{\mu})$ goes to 0, so that the next point is a draw from a normal distribution with the homebase $\boldsymbol{\mu}$ as a mean. Alternatively, if β is small (i.e., weak centralizing tendency), the exponential factor is close to 1 and the next position is a draw from a normal distribution with the previous position as its mean. Thus, as in the one-dimensional case, the matrix exponential $e^{-\mathbf{B}d}$ part behaves as a weighting function, taking values between 0 and 1, and adding a certain proportion to the homebase from the distance between the previous point and the homebase. For isotropic \mathbf{B} matrices, it holds that the conditional mean lies somewhere on the straight line connecting $\boldsymbol{\mu}$ and $\mathbf{Y}(t)$. This is further illustrated in Figure 1 where the 0.5 probability contour curves of three conditional distributions from (1) for a two-dimensional ($q = 2$) model are drawn. A cross \times denotes the center of each conditional distribution. In all three cases, all parameters are kept constant, except for the centralizing tendency which can take a small, medium, or large value (see the figure for the exact numerical values). With a small β value, the conditional mean of the distribution of the next point is close to the previous point but as β increases the conditional mean is moving closer to the homebase. Note that the conditional variance also depends on the centralizing tendency parameter: a larger β value implies a larger variability (see later).

The centralizing tendency controls the autocorrelation function. In Appendix A, it is shown that the (continuous) autocorrelation function $\rho(d)$ of an OU process equals $e^{-\beta d}$. Because the

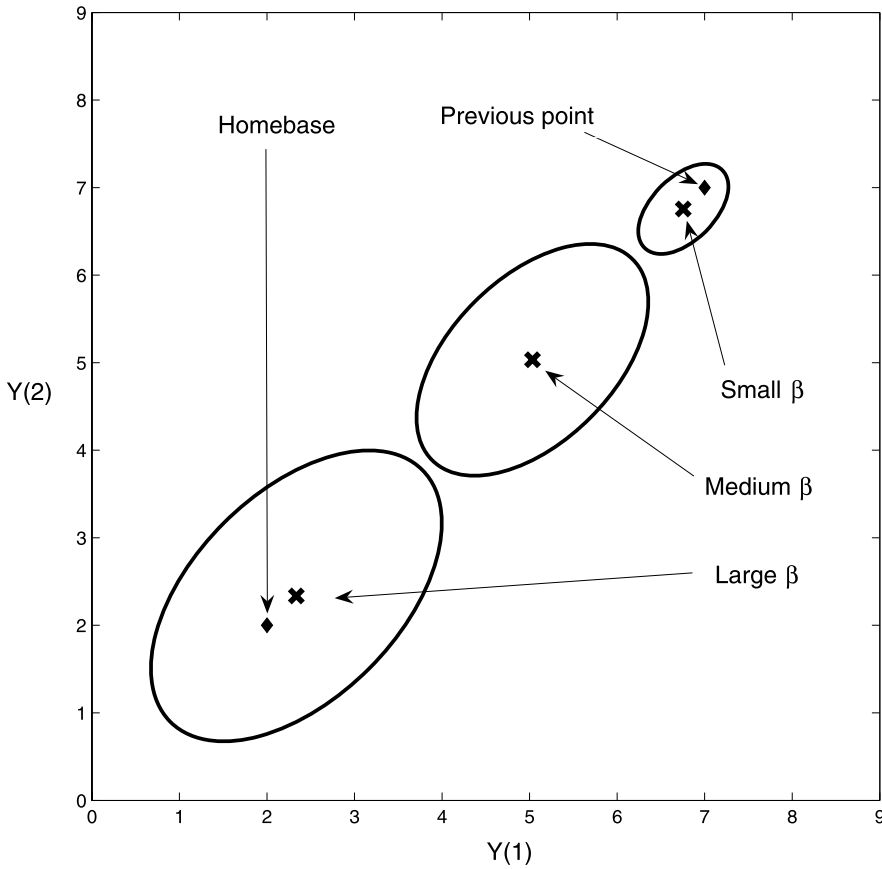


FIGURE 1.

The 0.5 probability contour curves of conditional distributions—the centers denoted by \times -s—with three different values for β .

autocorrelation function is an exponentially decaying function of (continuous) time, the OU process is the continuous time variant of an autoregressive process of order 1. Brockwell and Davis (2002) denote such a process as a CAR(1) process (where C stands for continuous-time). Figure 2 shows the change in the autocorrelation function as a function of time for four different β values. Although the range of the β parameter is relatively small, it has a remarkable effect on the slope of the autocorrelation function. Large (small) β values lead in general to small (large) autocorrelation because the serial correlation function decays more (less) steeply.

The matrix $\mathbf{\Gamma}$ is the covariance matrix of the stationary distribution in (2) and is part of the conditional covariance. It is a positive definite, symmetric $q \times q$ matrix containing the variances for each dimension on the diagonal (denoted with the corresponding row index) and the covariances as off-diagonal elements:

$$\mathbf{\Gamma} = \begin{pmatrix} \gamma_1 & \cdots & \gamma_{1q} \\ \vdots & \ddots & \vdots \\ \gamma_{q1} & \cdots & \gamma_q \end{pmatrix}.$$

Large variance values imply that the process can go through many changes (i.e., it is very volatile), while small variances lead to smoother trajectories. The covariances represent the ex-

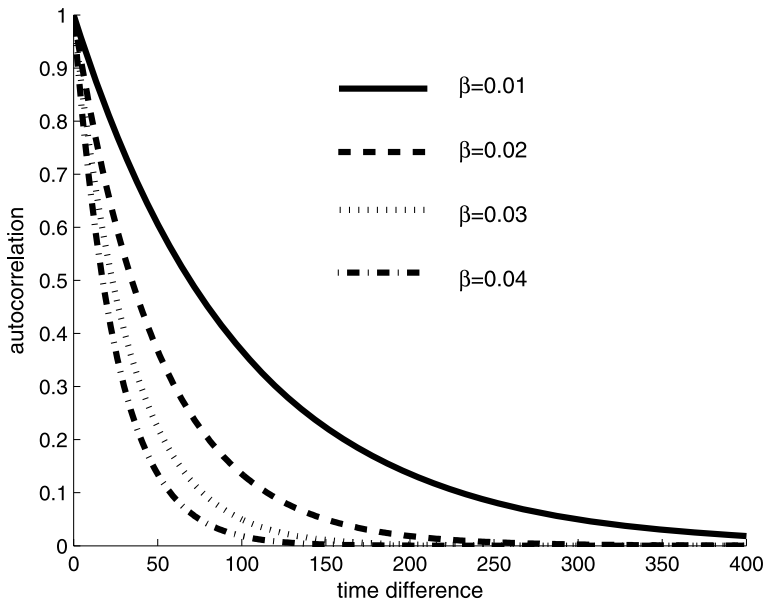


FIGURE 2.
Four different β values and their corresponding autocorrelation functions.

tent to which changes in one dimension tend to covary with changes in another dimension. In a unidimensional context, the single γ parameter is often referred to as the volatility parameter.

Examining the features of the instantaneous variance $\mathbf{\Gamma} - e^{-\mathbf{B}d}\mathbf{\Gamma}e^{-\mathbf{B}'d}$, we can see that as the exponential part goes to 0 (i.e., a large centralizing tendency and/or time difference), the instantaneous variance converges to the variance of the stationary distribution. As the exponential part goes to 1 (i.e., small centralizing tendency and/or time difference), the conditional variance becomes very small. To illustrate this, assume for simplicity that $q = 2$ (and \mathbf{B} is isotropic). Then $\mathbf{\Gamma} - e^{-\mathbf{B}d}\mathbf{\Gamma}e^{-\mathbf{B}'d}$ equals:

$$\begin{pmatrix} \gamma_1(1 - e^{-2\beta d}) & \gamma_{12}(1 - e^{-2\beta d}) \\ \gamma_{12}(1 - e^{-2\beta d}) & \gamma_2(1 - e^{-2\beta d}) \end{pmatrix}$$

from which we can see the effect of the centralizing tendency β for a higher dimensional case (given a constant time difference). On the one hand, when β is large, when γ_1 and γ_2 are multiplied by a number close to one, and hence the instantaneous variances are near to the variances of the stationary distribution. On the other hand, when β is small, the conditional variances are close to 0. The same reasoning applies to the covariances. However, when considering the instantaneous cross-correlations, it can easily be seen that they are independent of the centralizing tendency and the time difference. Moreover, they are equal to the correlations of the stationary distribution ($\rho = \frac{\gamma_{12}}{\sqrt{\gamma_1}\sqrt{\gamma_2}}$).

As already mentioned above, (1) is part of the solution of a stochastic differential equation (SDE). Without presenting too much technical detail, the SDE is a convenient and intuitively appealing way of representing the OU process because it provides a link with the more familiar field of deterministic differential equations. A full and rigorous treatment of SDEs can be found in Arnold (1974), Karlin and Taylor (1981), and Smith (2000). For demonstrative purposes, we simplify matters here, and we will assume that $q = 1$, such that the vector $\boldsymbol{\mu}$ and the matrices $\mathbf{\Gamma}$ and \mathbf{B} reduce to the scalars μ , γ and β , respectively. For the one-dimensional OU process as

considered in this paper, the corresponding SDE equals

$$dY(t) = \beta(\mu - Y(t))dt + \sqrt{2\beta\gamma} dW(t), \quad (3)$$

where $dY(t)$ is the (random) change in the process $Y(t)$ in a small time interval $(t, t + dt)$ and $W(t)$ represents a univariate standard (i.e., driftless and variance equal to one) Wiener process, which is a mathematical model for a continuous-state continuous-time process with independent increments (i.e., Brownian motion). The quantity $dW(t)$ is the increment of this process W in the small time interval $(t, t + dt)$. As can be seen from (3), the change in $Y(t)$ consists of two components: a deterministic part and a stochastic part (embodied by the Wiener process term). The solution of such a SDE requires a special stochastic calculus which we will not further discuss here (a short introduction can be found in Brockwell & Davis, 2002).

However, if we let $\gamma \rightarrow 0$, so that the stochastic part disappears, we are left with a simple first-order linear differential equation. One can also see that the magnitude of change in the deterministic part depends on the difference between the homebase μ and the current position $Y(t)$. If there were no stochastic disturbance in (3), the solution of the deterministic differential equation (given an arbitrary initial value at time 0 of $Y(0)$) would be equal to

$$Y(t) = \mu + (Y(0) - \mu)e^{-\beta t}.$$

This solution represents the scalar version of the mean of the conditional distribution in (1). Given the initial value $Y(0)$, the exact position of the deterministic process can be found for every time difference t . Moreover, if the time difference becomes large, the process converges to the homebase μ . However, when the stochastic disturbance term is added again, the OU process is retrieved and the exact position of the process is unpredictable because of the inherent stochastic nature of the process.

In the following section, we will discuss how to extend the basic OU process in order to make it appropriate for studying interindividual differences.

3. Hierarchical Extension of the OU Process

For the case where longitudinal data are collected for a random sample of persons, as is often the case in psychological research, it is natural to consider a hierarchical extension of the OU diffusion process in order to describe and explain interindividual differences. A hallmark of the presented model is that all parameters are allowed to vary over individuals (and not only the means as is commonly done).

Let us first fix some notation. A specific person p ($p = 1, \dots, P$) is measured n_p times at the following sequence of time points: $t_{p1}, t_{p2}, \dots, t_{ps}, \dots, t_{p,n_p}$. Note that we do not require that persons are measured at regular time intervals or that they are measured at exactly the same time points. The measured sequence of positions in the multidimensional space is denoted as $Y(t_{p1}), \dots, Y(t_{ps}), \dots, Y(t_{p,n_p})$. (Note that we will usually set t_{p1} equal to 0 for all persons to align the measurements of the different persons.) For all persons, the model for the first observation of the chain of measurements is the person-specific equilibrium distribution:

$$Y(t_{p1}) \sim N_q(\mu_p, \Gamma_p). \quad (4)$$

This assumption can be justified, since in many applications the process has been diffusing long enough to have converged to its stationary distribution and forgotten its initial position. For the subsequent points, we rely on (1) the conditional distribution of a single person p at time t_s

given its position at the previous measurement occasion $Y(t_{p,s-1})$ is normal with conditional mean vector δ_{ps} and conditional covariance matrix Λ_{ps}

$$Y(t_{ps}) | Y(t_{p,s-1}) \sim N_q(\delta_{ps}, \Lambda_{ps}), \quad (5)$$

where

$$\delta_{ps} = \mu_p + e^{-\mathbf{B}_p(t_{ps}-t_{p,s-1})} (Y(t_{p,s-1}) - \mu_p)$$

and

$$\Lambda_{ps} = \Gamma_p - e^{-\mathbf{B}_p(t_{ps}-t_{p,s-1})} \Gamma_p e^{-\mathbf{B}_p'(t_{ps}-t_{p,s-1})}.$$

Note that all parameter vectors and matrices carry an index p to denote that they are allowed to be person specific.

In a hierarchical model, the individual parameters are assumed to be drawn from a population distribution. Instead of merely listing these population distributions, we will also indicate how the model can be complemented with person-specific covariate information, with which we will attempt to explain interindividual differences in the basic OU parameters. If one only wants to describe the amount of between-person variation in some parameters, the covariates can be removed from the model so that only the intercept is left.

Let us suppose that k covariates are measured and x_{jp} denotes the score of person p on covariate j ($j = 1, \dots, k$). Then we can collect all covariate scores into a vector (together with a constant 1 for the intercept) $\mathbf{x}'_p = (1, x_{1p}, x_{2p}, \dots, x_{kp})$. Furthermore, let α_{μ_1} be the vector (of length $k+1$) with regression coefficients for the regression of the individual homebases onto the covariates. The vectors with regression coefficients for the other parameters are given names in a similar fashion (e.g., α_β for β , etc.).

At this point, we will introduce a simplification of our model in order to make the exposition not overly complex. Most potential applications we encountered for the hierarchical OU model are two-dimensional in nature and, therefore, we will assume in the remainder of the paper that $q = 2$. However, the extension to higher-dimensional cases is mostly evident, except for one part of the model (the variances), but it will be indicated explicitly how the general case can be handled.

The person-specific homebase μ_p (for $q = 2$) is assumed to be a draw from the following bivariate normal distribution:

$$\mu_p \sim N_2(\alpha_\mu \mathbf{x}'_p, \Sigma_\mu), \quad (6)$$

where

$$\alpha_\mu = (\alpha_{\mu_1}, \alpha_{\mu_2}),$$

so that α_μ is a $(k+1) \times 2$ matrix of regression coefficients. Furthermore, the matrix Σ_μ , which is defined as

$$\Sigma_\mu = \begin{pmatrix} \sigma_{\mu_1}^2 & \sigma_{\mu_1\mu_2} \\ \sigma_{\mu_1\mu_2} & \sigma_{\mu_2}^2 \end{pmatrix},$$

is the (residual) covariance matrix, representing the variations and associations that exist in the population between the individual means of the stationary distribution after taking into account the person covariates. As said above, if only the intercept is present in the covariate vector, then the model just describes the population mean vector of the homebases and the variability in the population. Note that the regression of μ_p onto covariates is in general a q -variate multiple regression problem since μ_p is of length q .

Not only the mean of the stationary distribution is assumed to be person-specific, but its covariance matrix as well:

$$\mathbf{\Gamma}_p = \begin{pmatrix} \gamma_{1p} & \gamma_{12p} \\ \gamma_{12p} & \gamma_{2p} \end{pmatrix}.$$

We need to propose a population distribution for these covariance matrices such that positive definiteness of $\mathbf{\Gamma}_p$ is ensured. Moreover, we would like to regress the variances and covariances on covariates. An obvious choice would be to assume that $\mathbf{\Gamma}_p$ is a draw from an inverse-Wishart distribution. However, such a distribution does not allow the regression of variances and covariances on covariates in a natural way. One possible alternative is to decompose the covariance matrix (Barnard, McCulloch, & Meng, 2000), usually into standard deviations and correlation matrices and assume proper distributions for these in order to ensure the positive definiteness. In this paper, we make use of the fact that $\mathbf{\Gamma}_p$ is a two-by-two covariance matrix such that it can be decomposed into two variances and a correlation. Next, the logarithms of the variances are assumed to be sampled from a normal distribution. After applying the Fisher- z transformation to the correlation coefficient, the transformed value is taken as a draw from a normal distribution, which also provides the possibility of regressing the mean of this distribution on covariates.

The two diagonal elements of the covariance matrix $\mathbf{\Gamma}_p$ are regressed on the covariates in the following way:

$$\begin{aligned} \log(\gamma_{1p}) &\sim N(\mathbf{x}'_p \boldsymbol{\alpha}_{\gamma_1}, \sigma_{\gamma_1}^2), \\ \log(\gamma_{2p}) &\sim N(\mathbf{x}'_p \boldsymbol{\alpha}_{\gamma_2}, \sigma_{\gamma_2}^2), \end{aligned}$$

with $\boldsymbol{\alpha}_{\gamma_1}$ and $\boldsymbol{\alpha}_{\gamma_2}$ being regression coefficient vectors with $k + 1$ components. Since it is assumed that the log-transformed γ -parameters are normally distributed, the original γ -parameters follow a lognormal distribution:

$$f(\gamma_{up}) = \frac{1}{\gamma_{up} \sqrt{2\pi \sigma_{\gamma_u}^2}} e^{-\frac{1}{2} \frac{(\log(\gamma_{up}) - \mathbf{x}'_p \boldsymbol{\alpha}_{\gamma_u})^2}{\sigma_{\gamma_u}^2}}, \quad (7)$$

for $u = 1, 2$ and where $f(\cdot)$ will be used in the remainder of the paper as the generic symbol to denote a probability density function.

The covariance parameter γ_{12p} of $\mathbf{\Gamma}_p$ can be expressed in terms of standard deviations and the correlation: $\gamma_{12p} = \sqrt{\gamma_{1p}} \times \sqrt{\gamma_{2p}} \times \rho_p$. Instead of proposing a population distribution for the covariance parameter, it will be assumed that the Fisher- z transformed (or z -transformed for short) individual-specific cross-correlation coefficient $F(\rho_p)$ is drawn from a normal population distribution whose mean depends on covariates

$$F(\rho_p) \sim N(\mathbf{x}'_p \boldsymbol{\alpha}_\rho, \sigma_\rho^2).$$

The parameter ρ_p is the cross-correlation for a person p and it indicates the extent to which changes in one dimension tend to correlate with changes in the other dimensions for person p . From the mean of the population distribution of the z -transformed ρ_p , it can be learned whether there is on average (i.e., in the population) a positive, negative, or zero correlation between the changes in two dimensions. The density of the original ρ_p then equals (applying the transforma-

tion of variables technique; see, e.g., Mood, Graybill, & Boes, 1974):

$$f(\rho_p) = \left| \frac{dF(\rho_p)}{d\rho_p} \right| \phi(F(\rho_p); \mathbf{x}'_p \boldsymbol{\alpha}_\rho, \sigma_\rho^2) \\ = \frac{1}{(1 - \rho_p)(1 + \rho_p)} \frac{1}{\sqrt{2\pi\sigma_\rho^2}} \exp\left(-\frac{1}{2} \frac{\left(\frac{1}{2} \log\left(\frac{1+\rho_p}{1-\rho_p}\right) - \mathbf{x}'_p \boldsymbol{\alpha}_\rho\right)^2}{\sigma_\rho^2}\right), \quad (8)$$

where $F(\cdot)$ is the Fisher- z transform and $\phi(x; \mu, \sigma^2)$ is the normal density evaluated at x with mean μ and variance σ^2 . Again, $\boldsymbol{\alpha}_\rho$ contains $k + 1$ regression coefficients.

It should be noted that the solution outlined here, where the elements of the matrix $\boldsymbol{\Gamma}_p$ are regressed onto covariates while still maintaining the positive definiteness of $\boldsymbol{\Gamma}_p$, is only valid in the two-dimensional case. For $q > 2$ and with a regression of the elements of the covariance matrix on predictors, we refer to Daniels and Pourahmadi (2002), whose approach is based on a Cholesky decomposition of the covariance matrix $\boldsymbol{\Gamma}_p$.

Finally, because the centralizing tendency matrix \mathbf{B}_p is isotropic, we need to assume a population distribution only for the single parameter β_p . Since β_p has to be positive, similarly to the variance parameters, it is assumed that the log-transformed β_p -values follow a normal distribution whose mean again depends on the covariates:

$$\log(\beta_p) \sim N(\mathbf{x}'_p \boldsymbol{\alpha}_\beta, \sigma_\beta^2),$$

where the length of vector $\boldsymbol{\alpha}_\beta$ is $k + 1$. As in the case of the variance parameters (γ_{1p} and γ_{2p}), the distribution of β_p is lognormal:

$$f(\beta_p) = \frac{1}{\beta_p \sqrt{2\pi\sigma_\beta^2}} e^{-\frac{1}{2} \frac{(\log(\beta_p) - \mathbf{x}'_p \boldsymbol{\alpha}_\beta)^2}{\sigma_\beta^2}}. \quad (9)$$

This completes the description of the model and puts us in a position to address issues of statistical inference.

4. Statistical Inference for the Ornstein–Uhlenbeck Model

The intrinsic complexity of the model motivated us to perform all statistical inferences in a Bayesian framework. The model complexity is mainly the result of the fact that all parameters are treated in a hierarchical sense and are allowed to differ over persons. The resulting high-dimensional integration over the many random effects distributions cannot be handled using brute force quadrature methods. Therefore, parameter estimation is done by sampling from the posterior density using a MCMC algorithm. For the model selection, we rely on the Deviance Information Criterion (DIC; Spiegelhalter, Best, Carlin, & van der Linde, 2002). The major part of this section is devoted to the estimation of the parameters. More details about the Bayesian methodology can be found in Gelman et al. (2004) and Robert and Casella (2004).

As a first step, we derive the likelihood. Let us denote the data from person p as follows: $\{\mathbf{Y}(t_{ps})\}_{s=1}^{n_p}$. Given the sequence of observations from person p , the likelihood contribution for

person p reads as

$$\begin{aligned}
 & f(Y(t_{p1}) \mid \boldsymbol{\mu}_p, \mathbf{B}_p, \boldsymbol{\Gamma}_p) \prod_{s=2}^{n_p} f(Y(t_{ps}) \mid Y(t_{p,s-1}), \boldsymbol{\mu}_p, \mathbf{B}_p, \boldsymbol{\Gamma}_p) \\
 & \propto \prod_{s=1}^{n_p} |\mathbf{V}_{ps}|^{-\frac{1}{2}} e^{-\frac{1}{2}(\boldsymbol{\lambda}_{ps}^T \mathbf{V}_{ps}^{-1} \boldsymbol{\lambda}_{ps})}
 \end{aligned} \tag{10}$$

with

$$\begin{aligned}
 \boldsymbol{\lambda}_{ps} &= \begin{cases} Y_{p1} - \boldsymbol{\mu}_p & \text{if } s = 1, \\ Y_{ps} - [\boldsymbol{\mu}_p + e^{-\mathbf{B}_p(t_{ps}-t_{p,s-1})}(Y(t_{p,s-1}) - \boldsymbol{\mu}_p)] & \text{if } s > 1, \end{cases} \\
 \mathbf{V}_{ps} &= \begin{cases} \boldsymbol{\Gamma}_p & \text{if } s = 1, \\ \boldsymbol{\Gamma}_p - e^{-\mathbf{B}_p(t_{ps}-t_{p,s-1})} \boldsymbol{\Gamma}_p e^{-\mathbf{B}_p'(t_{ps}-t_{p,s-1})} & \text{if } s > 1. \end{cases}
 \end{aligned} \tag{11}$$

The likelihood of all person-specific parameters (given the data from all persons $1, \dots, P$ and making use of the fact that the persons are independent) then equals

$$\begin{aligned}
 & f(\{Y(t_{1s})\}_{s=1}^{n_1}, \dots, \{Y(t_{Ps})\}_{s=1}^{n_P} \mid \boldsymbol{\mu}_1, \dots, \boldsymbol{\mu}_P, \mathbf{B}_1, \dots, \mathbf{B}_P, \boldsymbol{\Gamma}_1, \dots, \boldsymbol{\Gamma}_P) \\
 &= \prod_{p=1}^P f(\{Y(t_{ps})\}_{s=1}^{n_p} \mid \boldsymbol{\mu}_p, \mathbf{B}_p, \boldsymbol{\Gamma}_p) \propto \prod_{p=1}^P \prod_{s=1}^{n_p} |\mathbf{V}_{ps}|^{-\frac{1}{2}} e^{-\frac{1}{2}(\boldsymbol{\lambda}_{ps}^T \mathbf{V}_{ps}^{-1} \boldsymbol{\lambda}_{ps})},
 \end{aligned} \tag{12}$$

where $\boldsymbol{\lambda}_{ps}$ and \mathbf{V}_{ps} are defined as in (11).

To find the posterior distribution, let us first collect (for simplicity) all model parameters in a single parameter vector $\boldsymbol{\theta}$ which contains the (unique) elements of the person-specific vectors and matrices $\boldsymbol{\mu}_1, \dots, \boldsymbol{\mu}_P, \boldsymbol{\Gamma}_1, \dots, \boldsymbol{\Gamma}_P, \mathbf{B}_1, \dots, \mathbf{B}_P$, the regression coefficients $\boldsymbol{\alpha}_\mu, \boldsymbol{\alpha}_{\gamma_1}, \boldsymbol{\alpha}_{\gamma_2}, \boldsymbol{\alpha}_\rho, \boldsymbol{\alpha}_\beta$, the residual (co)variances $\boldsymbol{\Sigma}_\mu, \sigma_{\gamma_1}^2, \sigma_{\gamma_2}^2, \sigma_\rho^2, \sigma_\beta^2$. The joint posterior of all parameters can then be written as follows:

$$\begin{aligned}
 & f(\boldsymbol{\theta} \mid \{Y(t_{1s})\}_{s=1}^{n_1}, \dots, \{Y(t_{Ps})\}_{s=1}^{n_P}) \\
 & \propto f(\{Y(t_{1s})\}_{s=1}^{n_1}, \dots, \{Y(t_{Ps})\}_{s=1}^{n_P} \mid \boldsymbol{\mu}_1, \dots, \boldsymbol{\mu}_P, \mathbf{B}_1, \dots, \mathbf{B}_P, \boldsymbol{\Gamma}_1, \dots, \boldsymbol{\Gamma}_P) \\
 & \quad \times \prod_{p=1}^P f(\boldsymbol{\mu}_p \mid \boldsymbol{\alpha}_\mu, \boldsymbol{\Sigma}_\mu) \times \prod_{p=1}^P \prod_{u=1}^2 f(\gamma_{up} \mid \boldsymbol{\alpha}_{\gamma_u}, \sigma_{\gamma_u}^2) \\
 & \quad \times \prod_{p=1}^P f(\rho_p \mid \boldsymbol{\alpha}_\rho, \sigma_\rho^2) \times \prod_{p=1}^P f(\beta_p \mid \boldsymbol{\alpha}_\beta, \sigma_\beta^2) \\
 & \quad \times f(\boldsymbol{\alpha}_\mu) f(\boldsymbol{\alpha}_{\gamma_1}) f(\boldsymbol{\alpha}_{\gamma_2}) f(\boldsymbol{\alpha}_\rho) f(\boldsymbol{\alpha}_\beta) f(\boldsymbol{\Sigma}_\mu) f(\sigma_{\gamma_1}^2) f(\sigma_{\gamma_2}^2) f(\sigma_\rho^2) f(\sigma_\beta^2),
 \end{aligned} \tag{13}$$

where we have assumed independence between all sets of prior parameters. The distributions of the person-specific parameters (third line of (13)) are given in (6), (7), (8), and (9).

In (13), we have not yet specified the prior distributions (last line of the equation). For all regression coefficients $\boldsymbol{\alpha}$ we assume a uniform prior:

$$f(\boldsymbol{\alpha}_g) \propto 1,$$

where g can be equal to $\mu_1, \mu_2, \gamma_1, \gamma_2, \rho$ or β (for the two-dimensional case). The prior distribution of Σ_μ is assumed to be equal to

$$f(\Sigma_\mu) \propto |\Sigma_\mu|^{-(q+1)/2}$$

which is Jeffreys prior, and $q = 2$ in the two-dimensional case. We choose noninformative priors for all the other variance parameters σ_g^2 as well:

$$f(\sigma_g^2) \propto \sigma_g^{-2}$$

which is a uniform prior on $\log \sigma$. For the two-dimensional case g can be equal to $\mu_1, \mu_2, \gamma_1, \gamma_2, \rho$, or β .

To sample from the joint posterior, we make use of the Gibbs sampler (Gelman et al., 2004; Robert & Casella, 2004). For this, we need to derive the full conditionals, that is, the conditional distribution of each parameter given the other parameters and the data. In many cases, the full conditionals are known densities from which one can sample directly. If the full conditional is an unknown distribution, we make use of a Metropolis–Hastings step in the Gibbs sampler to obtain a draw from it. Because of reasons of efficiency, when deriving the full conditionals, we try to treat parameters that logically belong together as one block (e.g., Σ_μ or α_μ). We start with the full conditionals of the regression coefficients for the variances, the centralizing tendency and the cross-correlation and subsequently treat the residual variances of these parameters. Next, we treat the regression coefficients and residual covariance matrix for the mean positions and move then to the lowest level parameters (the individual homebases, variances, cross-correlations, and centralizing tendencies).

With regard to the unidimensional parameters (γ_1, γ_2, ρ , and β), the full conditionals of their regression coefficients ($\alpha_{\gamma_1}, \alpha_{\gamma_2}, \alpha_\rho$, and α_β) and their residual variances ($\sigma_{\gamma_1}^2, \sigma_{\gamma_2}^2, \sigma_\rho^2$, and σ_β^2) can be derived in a similar fashion. Here, we give the example of the full conditional of α_{γ_1} and $\sigma_{\gamma_1}^2$, but γ_1 could be substituted by γ_2, ρ , or β as well. The full conditional of α_{γ_1} reads as:

$$f(\alpha_{\gamma_1} | \gamma_{11}, \dots, \gamma_{1P}, \sigma_{\gamma_1}^2) \propto |V_g|^{-\frac{1}{2}} \exp \left(-\frac{1}{2} (\alpha_{\gamma_1} - X \hat{\alpha}_{\gamma_1})' V_g^{-1} (\alpha_{\gamma_1} - X \hat{\alpha}_{\gamma_1}) \right), \quad (14)$$

where X is a $P \times (k+1)$ matrix defined by stacking the person-specific covariate vectors \mathbf{x}'_p below each other. If we denote $\mathbf{g} = (\log(\gamma_{11}), \dots, \log(\gamma_{1P}))'$ such that $\hat{\alpha}_{\gamma_1} = (X'X)^{-1} X' \mathbf{g}$ and $V_g = \sigma_{\gamma_1}^2 (X'X)^{-1}$, it can be seen that the full conditional of α_{γ_1} is a normal density with mean $X \hat{\alpha}_{\gamma_1}$ and covariance matrix V_g .

The full conditional for $\sigma_{\gamma_1}^2$, the residual variance of α_{γ_1} , follows a scaled inverse- χ^2 distribution:

$$f(\sigma_{\gamma_1}^2 | \gamma_{11}, \dots, \gamma_{1P}) \propto (\sigma_{\gamma_1}^2)^{-\left(\frac{P-k-1}{2}+1\right)} e^{-\frac{(P-k-1)s^2}{2\sigma_{\gamma_1}^2}}$$

with

$$s^2 = \frac{1}{P-k-1} (\mathbf{g} - X \hat{\alpha}_{\gamma_1})' (\mathbf{g} - X \hat{\alpha}_{\gamma_1}),$$

where \mathbf{g} , X , and $\hat{\alpha}_{\gamma_1}$ are defined in the same way as in (14).

The full conditional of α_μ is also a known density, but it is somewhat harder to obtain since it involves a multivariate regression problem. The Bayesian treatment of multivariate regression is described in Zellner (1971). The solution lies in treating α_μ and Σ_μ together: First, we draw

Σ_μ given all other parameters (except α_μ) and subsequently we draw α_μ , conditional upon all other parameters and Σ_μ . To start, we define the matrix \mathbf{M} as the $P \times 2$ matrix of individual homebases, that is, $\mathbf{M} = (\mu_1, \dots, \mu_P)'$. Then the least squares regression coefficient matrix (of the regression of \mathbf{M} on \mathbf{X} , where the latter is defined in (14)), equals $\hat{\mathbf{A}}_\mu = (\mathbf{X}'\mathbf{X})^{-1}\mathbf{X}'\mathbf{M}$. Stacking the two columns of $\hat{\mathbf{A}}_\mu$ below each other results in $\hat{\alpha}_\mu = (\hat{\mathbf{A}}'_{\mu_1}, \hat{\mathbf{A}}'_{\mu_2})'$. In the same vein, stacking the two columns of α_μ below one another gives $\tilde{\alpha}_\mu = (\alpha'_{\mu_1}, \alpha'_{\mu_2})'$. Also, define $\mathbf{S} = (\mathbf{M} - \mathbf{X}\hat{\mathbf{A}}_\mu)'(\mathbf{M} - \mathbf{X}\hat{\mathbf{A}}_\mu)$. The full condition of Σ_μ then equals

$$f(\Sigma_\mu \mid \mu_1, \dots, \mu_P) \propto |\Sigma_\mu|^{-\frac{P-k+2}{2}} e^{-\frac{1}{2}\text{tr}\Sigma_\mu^{-1}\mathbf{S}} \quad (15)$$

which is an inverse-Wishart distribution with scale matrix \mathbf{S} and degrees of freedom $v = P - k - 1$. The full conditional of the matrix α_μ is

$$f(\alpha_\mu \mid \Sigma_\mu, \mu_p, \mathbf{x}'_p) \propto |\Sigma_\mu|^{-k/2} e^{-\frac{1}{2}(\tilde{\alpha}_\mu - \hat{\alpha}_\mu)' \Sigma_\mu^{-1} \otimes (\mathbf{x}_p \mathbf{x}'_p) (\tilde{\alpha}_\mu - \hat{\alpha}_\mu)} \quad (16)$$

with \otimes denoting the Kronecker product.

The full conditional of μ_p (with $p = 1, \dots, P$) is a bivariate normal distribution (because of conjugacy of the relevant parts of the likelihood and prior):

$$\mu_p \mid \{\mathbf{Y}_{ps}\}_{s=1}^{n_p}, \mathbf{B}_p, \Gamma_p, \alpha_\mu, \Sigma_\mu \sim N_2(\Omega_p, \Phi_p)$$

where

$$\begin{aligned} \Phi_p &= \left(\Sigma_\mu^{-1} + \Gamma_p^{-1} + \sum_{s=2}^{n_p} \mathbf{V}_{ps}^{-1} - \sum_{s=2}^{n_p} \mathbf{V}_{ps}^{-1} e^{-\mathbf{B}_p d_{ps}} \right. \\ &\quad \left. - \sum_{s=2}^{n_p} (e^{-\mathbf{B}_p d_{ps}})^T \mathbf{V}_{ps}^{-1} + \sum_{s=2}^{n_p} (e^{-\mathbf{B}_p d_{ps}})^T \mathbf{V}_{ps}^{-1} e^{-\mathbf{B}_p d_{ps}} \right)^{-1}, \\ \Omega_p &= \Phi_p \left(\Sigma_\mu^{-1} \mathbf{x}'_p \alpha_\mu + \Gamma_p^{-1} \mathbf{Y}_{p1} + \sum_{s=2}^{n_p} \mathbf{V}_{ps}^{-1} \mathbf{Y}_{ps} - \sum_{s=2}^{n_p} \mathbf{V}_{ps}^{-1} e^{-\mathbf{B}_p d_{ps}} \mathbf{Y}_{p,s-1} \right. \\ &\quad \left. - \sum_{s=2}^{n_p} (e^{-\mathbf{B}_p d_{ps}})^T \mathbf{V}_{ps}^{-1} \mathbf{Y}_{ps} + \sum_{s=2}^{n_p} (e^{-\mathbf{B}_p d_{ps}})^T \mathbf{V}_{ps}^{-1} e^{-\mathbf{B}_p d_{ps}} \mathbf{Y}_{p,s-1} \right) \end{aligned}$$

with $d_{ps} = t_{ps} - t_{p,s-1}$ and \mathbf{V}_{ps} is defined in (11).

Unfortunately, there is no closed form solution for the rest of the full conditional distributions of the person specific diffusion parameters. To calculate the posterior of the covariance matrix Γ_p , we use the decomposition method which has been described earlier. Consequently, we have to deal with the calculation of the conditional distributions of the variances and the correlation. The full conditional of the variance γ_{1p} is

$$f(\gamma_{1p} \mid \{\mathbf{Y}_{ps}\}_{s=1}^{n_p}, \mu_p, \gamma_{2p}, \rho_p, \mathbf{B}_p) \propto f(\gamma_{1p}) \prod_{s=1}^{n_p} |\mathbf{V}_{ps}|^{-\frac{1}{2}} e^{-\frac{1}{2}(\lambda_{ps}^T \mathbf{V}_{ps}^{-1} \lambda_{ps})}, \quad (17)$$

where λ_{ps} and \mathbf{V}_{ps} are defined as in (11). For the exact expression of $f(\gamma_{1p})$, please see (7). The formulation of the posterior distribution of γ_{2p} follows exactly the same principle.

The expression for the full conditional of ρ_p is very similar to that of the variances:

$$f(\rho_p \mid \{\mathbf{Y}_{ps}\}_{s=1}^{n_p}, \boldsymbol{\mu}_p, \gamma_{1p}, \gamma_{2p}, \mathbf{B}_p) \propto f(\rho_p) \prod_{s=1}^{n_p} |\mathbf{V}_{ps}|^{-\frac{1}{2}} e^{-\frac{1}{2}(\boldsymbol{\lambda}_{ps}^T \mathbf{V}_{ps}^{-1} \boldsymbol{\lambda}_{ps})} \quad (18)$$

since $\boldsymbol{\lambda}_{ps}$ and \mathbf{V}_{ps} are specified as in (11). For the formula of $f(\rho_p)$, see (8).

Since $\mathbf{B}_p = \beta_p \mathbf{I}$, we have to deal only with β_p . The full conditional for β_p equals (note that because this parameter does not play a role in the distribution of the first observation the product starts only at $s = 2$):

$$f(\beta_p \mid \{\mathbf{Y}_{ps}\}_{s=1}^{n_p}, \boldsymbol{\mu}_p, \boldsymbol{\Gamma}_p) \propto f(\beta_p) \prod_{s=2}^{n_p} |\mathbf{V}_{ps}|^{-\frac{1}{2}} e^{-\frac{1}{2}(\boldsymbol{\lambda}_{ps}^T \mathbf{V}_{ps}^{-1} \boldsymbol{\lambda}_{ps})}, \quad (19)$$

where $f(\beta_p)$ is as in (9) and $\boldsymbol{\lambda}_{ps}$ and \mathbf{V}_{ps} are shown in (11).

As has been discussed above, we use the Gibbs sampler for sampling for the posterior distribution. If a full conditional distribution is not a known distribution from which it is easy to sample, a Metropolis–Hastings step is used. For this procedure, reasonable candidate generating distributions have to be assigned. The types of these distributions were always chosen to be identical to the population distribution of these parameters with the previously accepted value as a mean and with a variance which ensured a reasonable acceptance ratio (around 0.44, see Gelman et al., 2004, p. 306). The acceptance ratio was monitored and updated during the burn-in. The evaluation of the convergence is based on a visual assessment of the trace plots and on the values of the \hat{R} diagnostic as it is described by Gelman et al. (2004).

A software program to sample from the joint posterior has been written in MATLAB. However, as can be seen from the equations of the full conditionals for all person-specific parameters (with the exception of $\boldsymbol{\mu}_p$), we have to calculate a product with n_p factors involved (or a sum of n_p terms on the logscale). Since this calculation has to be performed many times in an MCMC algorithm, the process is computationally very demanding. For that reason, we have written the most computationally intensive subroutines of the code—namely the above mentioned person-specific likelihood parts—in C++, which then can be called from MATLAB in a straightforward way. Consequently, the computation time is highly reduced. As an example, 10,000 iterations take 2 hours on a computing node with an AMD Opteron250 processor and 2 Gb of RAM.

To carry out subsequent model selection, we opted for the Deviance Information Criterion (DIC) statistic (Spiegelhalter et al., 2002). The DIC takes into account two important features of the model: the complexity (based on the number of the parameters) and the fit (typically measured by a deviance statistic). DIC examines the two features together and gives a measure which balances between the two. Its formula is the sum of the effective number of parameters and the posterior mean of the deviance (defined as -2 times loglikelihood). Theoretically, the model with smaller DIC would better predict a replicate dataset of the same structure.

5. Application to Core Affect Trajectories

The hierarchical OU model as described in the previous sections will serve as a model for the trajectories of individuals in the core affect space (Russell, 2003). According to Russell (2003), core affect lies at the heart of a person's emotional experience and can be characterized as a compound of hedonic (pleasure-displeasure) and arousal (deactivated-activated) values. This core affect is always part of the human psyche as a consciously accessible state, which changes continuously over time. The core affect space is defined by two dimensions: activation

(vs. deactivation) and pleasantness (vs. unpleasantness). Consequently, the emotional experience at a particular moment can be represented as a single point in the two-dimensional plane, and the itinerary of a person’s emotional experience is the core affect trajectory. Our goal with modeling the core affect variability with an OU model is twofold. First, we want to describe the individual and population characteristics of movement throughout the core affect space. By making use of a stochastic model approach such as the OU model, we are able to treat the elapsed time as continuous and model the two dependent variables (pleasantness and activation) simultaneously (together with their cross-correlation). Moreover, we are able to evaluate the strength of the centralizing tendency and the magnitude of individual difference in it. As a second goal, we want to explain the individual differences in the characteristics of individual trajectories. In doing this, we could try to answer such questions as: Can the average position (mood) of an individual be predicted from some of their major personality dimensions?

The most common method for collecting data about such trajectories is experience sampling (Bolger et al., 2003; Csikszentmihalyi & Larson, 1987; Larson & Csikszentmihalyi, 1983; Russell & Feldman-Barrett, 1999). Persons are surveyed repeatedly at randomly chosen time points with respect to their position in the core affect space and this assessment takes place in the natural environment of the participants. From data obtained by experience sampling, we are able to investigate the intra and interindividual variation in core affect position.

As an illustration of the hierarchical OU model for the core affect trajectories, we used a dataset of which a subset has been described in Kuppens, Van Mechelen, Nezlek, Dossche, and Timmermans (2007). In the present study, 80 students from the University of Leuven were paid to give systematic self-reports about their emotional state in the core affect space during one week. The participants were provided with a booklet in which they could indicate their positions on the Affect Grid (Russell, Weiss, & Mendelsohn, 1989; see the used format of the grid on Figure 3), when a preprogrammed wristwatch beeped (9 times a day).

Not surprisingly, some of the planned measurements were missing. If the participants missed a beep, they had to indicate why they missed it, and in almost all cases the reason was that they did not hear it. Therefore, the missing data were considered to be missing completely at random

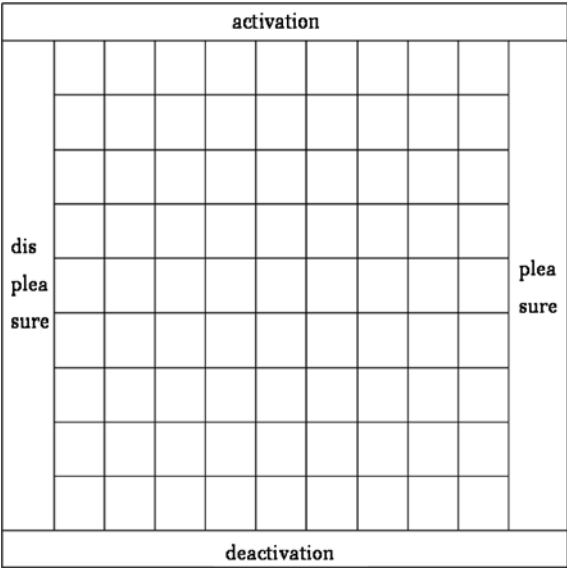


FIGURE 3.
The Affect Grid used in the application.

(MCAR; Little & Rubin, 2002) and it was assumed that there was no observation at that particular time (which makes the data unbalanced).

The average age of the participants was 21.7 years ($SD = 4.7$) and 60% of them were women. The maximum number of measurements for a single person was 63 and on average there were 60 measurements per person ($SD = 3.4$). The elapsed time interval between the measurements was semi-random. The participants were asked to give information about the time when they were awake, this interval was divided into equal periods, and a random beep was scheduled into each period. As a result of this procedure, we do not have measurements for the nights.

In addition to the experience sampling, in an introductory session the participants also completed the Dutch version of NEO-FFI (Hoekstra, Ormel, & de Fruyt, 1996), which is a questionnaire to measure the dimensions of the Five Factor model of personality (Big Five). The NEO-FFI consists of 60 items divided equally into five scales which assess Neuroticism, Extraversion, Openness to experience, Agreeableness, and Conscientiousness. All items are rated on a 5-point scale ranging from 1 (strongly disagree) to 5 (strongly agree). The five factors will be used as covariates to explain individual differences in the characteristics of core affect trajectories.

5.1. Exploratory Data Analysis

An exploratory data analysis was carried out to investigate the main characteristics of the measurements. First, we present two typical person profiles from the data set (Figure 4).

Subsequent measurements are connected with straight lines. We can see from these profiles that there is variability in the occupied positions in the core affect space but we also notice that the extent of this variability may differ somewhat between individuals.

Figure 5 shows a smoothed heat map of the visit frequencies in the core affect space, based on the aggregated data. We can clearly see a central area, where most of the visits are concentrated. We expect the population distribution of homebases to be located somewhere in that area.

Figure 6 shows the estimated vector-field of the core affect grid with the data pooled together for all individuals in the dataset. For each cell in the core affect grid, the length of the vector is proportional to the estimated escape velocity from that cell and its angle corresponds to the direction of escape. At the more central locations in the grid, the average velocity to move away is very small. However, as farther away from this central point, there is a tendency to be pulled toward the center: the directions of most of the vectors are more or less toward the central location

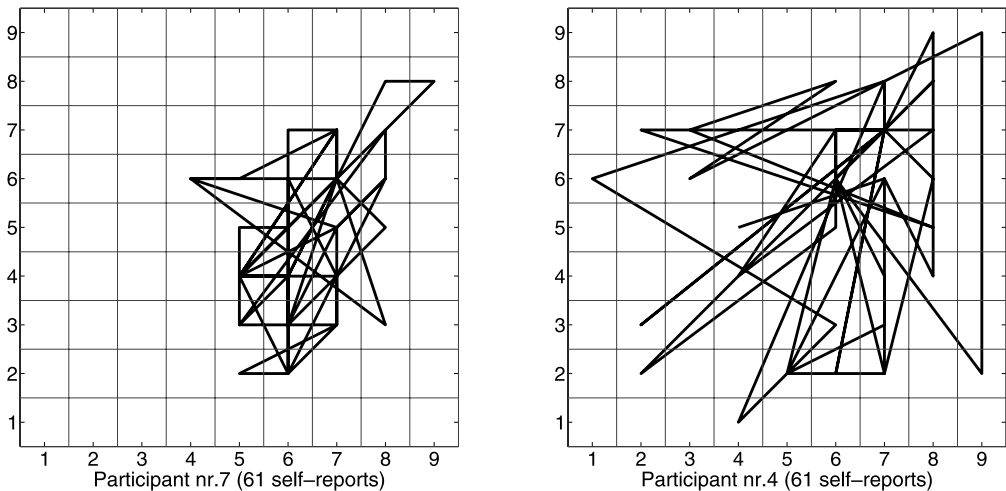


FIGURE 4.
Person profiles in the core affect grid.

and with increasing distance from the central point, the vector length tends to increase. (At the border cells, we notice more irregularity—these are due to sampling variability because there are much fewer visits to these outer cells; see also the heat map in Fig. 5.)

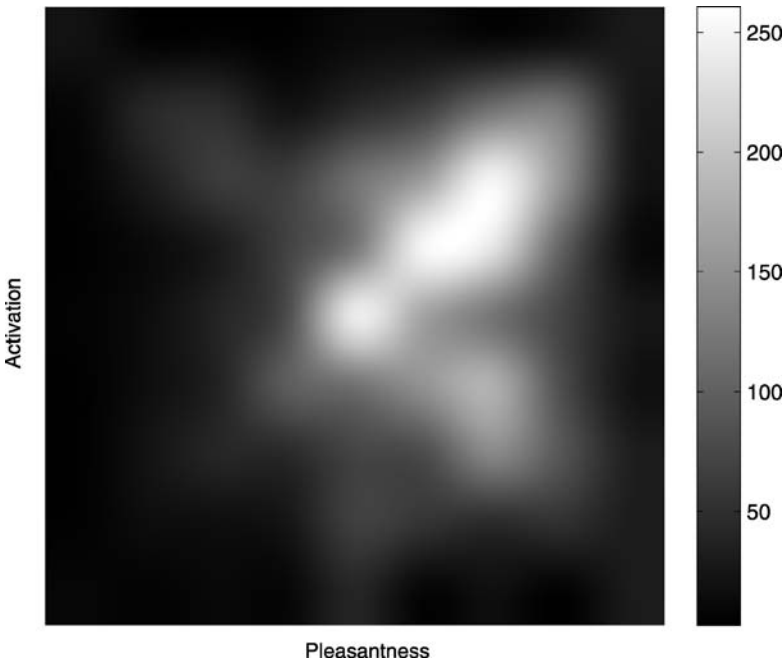


FIGURE 5.
Smoothed heat map of the visit frequencies in the core affect space.

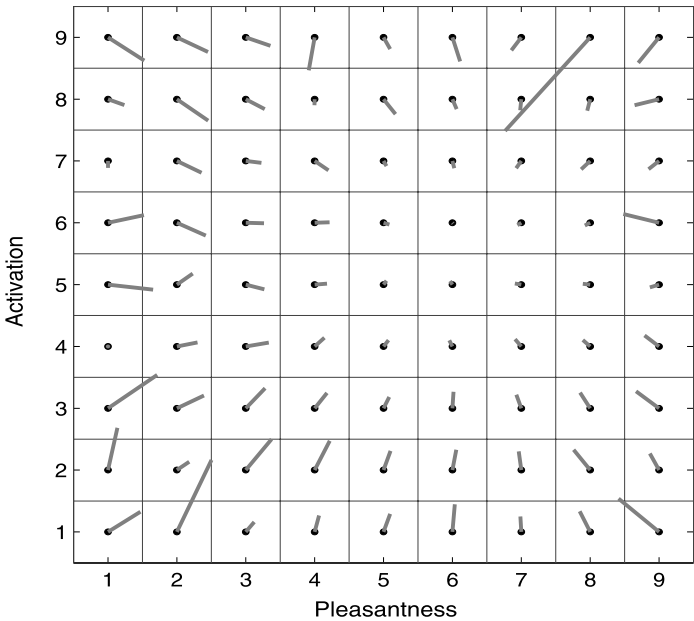


FIGURE 6.
Estimated vector field of the core affect grid.

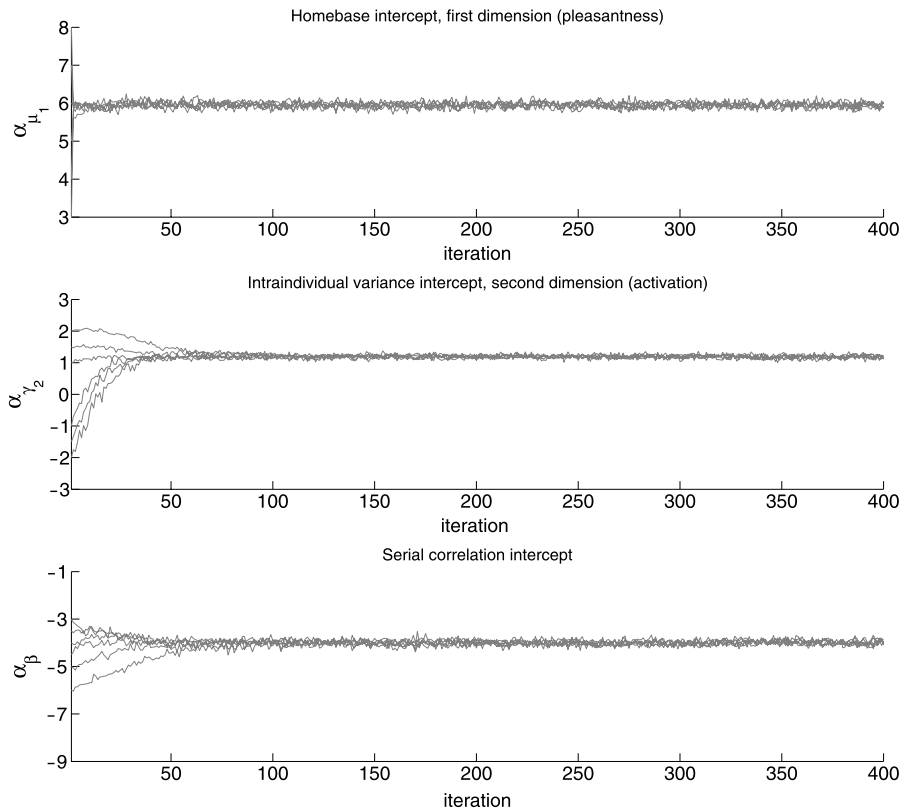


FIGURE 7.

Convergence of the core affect space data with respect to the population mean of the pleasantness homebase, intraindividual variance in activation and the serial-correlation.

5.2. Implementation of the MCMC Algorithm

The results presented below are based on 36,000 draws from the posterior distribution, which come from six chains with 6,000 iterations each. Each chain started with a burn-in period, to be discarded, of length 4,000. The initial values of the chains are randomly perturbed rational values derived from the data (e.g., the sample average for the homebases). Convergence checks showed no problems (all $\hat{R} < 1.1$). Generally, the convergence was fast for all parameters. Figure 7 shows the six iteration histories for three parameters (in each case starting from different initial values).

To demonstrate the efficiency, we ran some simulations for illustrative purposes. We simulated two data sets with the estimated population means, but with different numbers of people. Figure 8 shows the recovered population means with sample size 80, as in the application in the first column and with sample size 20 in the second column. The black line shows the simulated value. As can be seen, as the sample size increases, convergence is faster and the uncertainty with respect to the mean decreases, but the model does reasonably well with a relatively small sample size (20 subjects) as well.

5.3. Results

In a first step, we estimated the model without covariates to get a general description about the core affect space (technically by reducing the covariate vector \mathbf{x}'_p to the scalar value 1 for each person). Table 1 shows a summary of the results, containing the posterior mean (which is

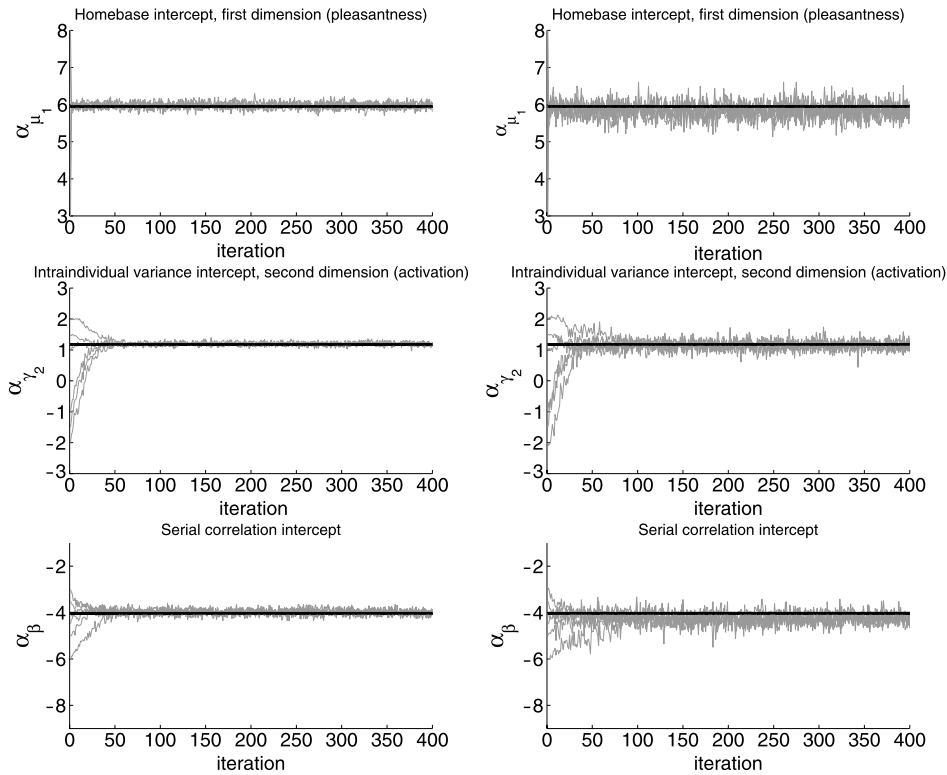


FIGURE 8. Convergence of the simulated data with respect to the population mean of the pleasantness homebase, intraindividual variance in activation and the serial-correlation for population size 80 (first column) and 20 (second column).

TABLE 1. Summary measures of the posterior distributions for the most important parameters of the OU model estimated without covariates.

Model parameter	Description	Posterior mean	95% posterior credibility interval		Posterior SD
	Pleasantness				
α_{μ_1}	Average homebase	5.95	5.79	6.10	0.08
σ_{μ_1}	SD of the average homebase	0.42	0.29	0.61	0.08
α_{γ_1}	Average log-variability	0.94	0.81	1.07	0.06
	Activation				
α_{μ_2}	Average homebase	5.23	5.09	5.37	0.07
σ_{μ_2}	SD of the average homebase	0.30	0.19	0.45	0.06
α_{γ_2}	Average log-variability	1.18	1.07	1.30	0.05
$\sigma_{\mu_1\mu_2}$	Covariance between the homebases	0.05	−0.04	0.16	0.05
α_{ρ}	Average Fisher transformed cross-correlation	0.02	−0.04	0.08	0.03
α_{β}	Average log-centralizing tendency	−4.03	−4.23	−3.82	0.10

technically the intercept) and standard deviation and the endpoints of the 95% posterior credibility interval. The estimated means of the homebase population distribution are (5.95, 5.23), which correspond to the findings of previous research (Russell et al., 1989). It shows that on average, the emotional state of persons is slightly pleasant and rather activated than deactivated.

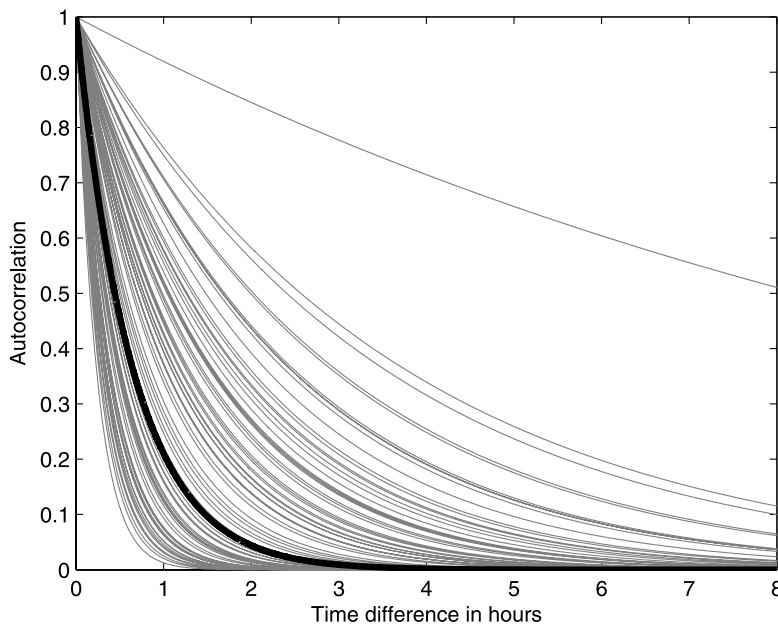


FIGURE 9.

Change in the autocorrelation function according to the estimated β_p parameters of the individuals.

If we look at the standard deviations of the homebases (0.42, 0.30) on the two dimensions, it appears that there is an approximately equal amount of variability in the homebases across persons. The covariance in the population between the homebases of the dimensions is estimated to be 0.05, which means that across persons, the homebases of the pleasantness and the activation dimensions are not related.

The average log-variance of the activation-deactivation dimension (1.18, which corresponds to an expected mean of 3.65 on the normal scale) is somewhat larger than the average log-variance of the pleasantness-unpleasantness dimension (0.94, i.e., 2.94 on the normal scale). Considering the posterior mean of the Fisher transformed correlation in the stationary distribution, it can be seen that the latter is rather small (i.e., 0.02), suggesting that on average, there is not much cross-correlation (i.e., on average, the bivariate person-specific stationary distributions have zero correlation). However, there is considerable variability in the estimated person-specific untransformed cross-correlation parameters: The average value is 0.02, while the between-person standard deviation is large (0.21), with the endpoints of the 95% credibility interval equal to -0.42 and 0.39 , respectively. (It must be emphasized here that there is a large conceptual difference between the population correlation of the homebase distribution as discussed above—equal to 0.05—and the average correlation of the stationary distribution as discussed in this paragraph; both concepts are unrelated to each other.)

The mean of the population distribution of the logarithm of the centralizing tendency is estimated to be -4.03 . If we convert it back to the original scale, the estimate for the mean of the population average centralizing tendency is 0.026. We may also look at the posterior estimates for the person-specific centralizing force (i.e., β_p). To illustrate the interpretation of this average and the individual variation in centralizing tendency graphically, we convert the posterior person-specific β_p estimates to the corresponding autocorrelation functions. Figure 9 shows the person-specific autocorrelation function together with the autocorrelation functions based on the average population value (thick line). For most participants, the autocorrelation between subsequent core affect positions separated 2 hours in time is fallen below 0.2.

TABLE 2.
Summary of the regression coefficients with a 95% posterior credibility not containing 0.

Model parameter	Description	Covariate	Posterior mean	95% posterior credibility interval		Posterior SD
	Pleasantness					
$\alpha_{\mu_{1N}}$	Homebase	Neuroticism	−0.38	−0.64	−0.11	0.13
$\alpha_{\gamma_{1N}}$	Variability	Neuroticism	0.23	0.01	0.45	0.11
$\alpha_{\gamma_{1A}}$	Variability	Agreeableness	−0.29	−0.58	−0.00	0.14
α_{ρ_C}	Cross-correlation	Conscientiousness	−0.17	−0.29	−0.06	0.06

Note. Model parameters refer to the regression weights. For example, $\alpha_{\mu_{1N}}$ is the regression weight for neuroticism relating to the homebase in the pleasantness dimension (μ_1).

In a next step, each of the six person-specific Ornstein–Uhlenbeck parameters (i.e., the two homebases μ_p and the log-variances $\log(\gamma_{up})$ of the pleasantness-unpleasantness and the activation-deactivation dimensions, the Fisher transformed cross-correlation between the dimensions $F(\rho_p)$ and the log-centralizing tendency $\log(\beta_p)$) were regressed onto the Big Five personality dimensions. Table 2 summarizes posterior means and standard deviations for the four regression coefficients for which the 95% posterior credibility intervals do not contain zero.

The current analysis shows that the neurotic individuals tend to have a lower homebase with respect to pleasantness. Also, they show higher intra-individual variance in this dimension. Based on these findings, it seems that they generally feel quite unpleasant but it changes dynamically, which might suggest emotional instability. In contrast, agreeable individuals tend to have lower variation with respect to pleasantness.

The cross-correlation parameter presents another interesting aspect of the core affect space. It shows how changes with respect to hedonic and arousal values coincide. Interestingly, for conscientious individuals, changes in one dimension are negatively correlated with the changes in the other one. In their case, high levels of pleasantness might often be accompanied with relatively low levels of activation and vice versa.

The presented results are based on the basic OU model in which every diffusion parameter was modeled as a random effect, that is, we assumed that people would differ from each other with respect to all of the modeled variables. However, we can construct simpler models, by removing individual difference dimensions from the model (i.e., restricting parameters to be equal across persons) or by dropping parameters altogether. Three constrained models were fitted and compared to the basic model using the DIC (estimated on 4,000 posterior draws). The DIC value of the basic model is 18,103. In the first constrained model, we did not allow β_p to vary across persons, but we still estimated a nonperson-specific β . The DIC of this model is higher, namely 18,366. Then we tested whether there is any need for autocorrelation in the model at all by fixing β to a large value. The resulting DIC of this constrained model was even higher, 19,377, indicating that we indeed need to take into account the autocorrelation in the model. Finally, we estimated a very simple constrained model where just the homebases were person-specific. This model showed the highest DIC value (i.e., 19,723). These outcomes tend to confirm our assumption that people do differ with respect to the modeled parameters.

6. Conclusion

In this paper, we have introduced a model for the analysis of multivariate longitudinal profiles with continuous and possibly unbalanced measurement times. The presented hierarchical diffusion modeling approach offers three distinctive characteristics.

First of all, in our approach a multivariate Ornstein–Uhlenbeck diffusion process serves as a model for individual behavior. The OU diffusion process is an intuitively appealing way of describing the continuous change of certain phenomena over time, certainly for constructs related to mood and emotion. Although we must mention that the current methodology is developed for quantitative variables, and not yet able to deal with categorical or mixed responses.

Second, we have succeeded in allowing for individual differences in all parameters of the model. Traditionally, one only looks for between-person differences in the mean structure. However, our application suggests that with respect to variation in core affect trajectories, there are also differences between people in their variabilities, cross-correlations between the dimensions and autocorrelations. The latter aspects are usually thought of as fixed over persons. That there is individual variation in these features has been confirmed by regressing the individual differences parameters on covariates and thereby explaining the existing variation. Although we have only a single application, it holds interesting suggestions for further substantive research in the core affect domain.

Third, we have fitted the model using a Bayesian approach. The parameter estimation procedure is exact (in the sense that it is not an approximated model that is fitted) and relatively fast (a single analysis takes only a few hours). Bayesian parameter estimation has the advantage that relatively complex models (as the one described in this paper) can be handled and that the uncertainty in the parameter estimates is easy to express (by means of credibility intervals or even visually by plotting marginal posterior distributions). It also provides a useful framework for model selection.

A major remaining challenge regarding the present modeling framework is the issue of systematic model testing. In this paper, we have concentrated on model selection (i.e., selecting the best fitting model among a predefined set of possible models) using the DIC. This is a common way of working in the domain on hierarchical or multilevel models where model selection is carried out using approximative methods such as DIC, AIC, or BIC (the latter two employed in a non-Bayesian context). In addition, absolute measures of model fit can also be considered, but this is a far more involved issue not in the last place because of the hierarchical nature of the model. Candidate procedures that could be considered involve, for instance, the posterior predictive check framework (e.g., Gelman et al., 2004; Gelman, Goegebeur, Tuerlinckx, & Van Mechelen, 2000).

Appendix A. Derivation of the Autocorrelation Function for a Multivariate OU Process

For simplicity and without loss of generality, let us assume that $\boldsymbol{\mu}$ is a $q \times 1$ vector of zeros. Then the covariance matrix function of the OU process with a general \mathbf{B} and $\boldsymbol{\Gamma}$ can be derived as follows (see also Schach, 1971):

$$\begin{aligned} E[\mathbf{Y}(t)\mathbf{Y}'(t+d)] &= E[E(\mathbf{Y}(t)\mathbf{Y}'(t+d) \mid \mathbf{Y}(t))] \\ &= E[\mathbf{Y}(t)E(\mathbf{Y}'(t+d) \mid \mathbf{Y}(t))] \\ &= E[\mathbf{Y}(t)\mathbf{Y}'(t)e^{-\mathbf{B}'d}] \\ &= E[\mathbf{Y}(t)\mathbf{Y}(t)']e^{-\mathbf{B}'d} \\ &= \boldsymbol{\Gamma}e^{-\mathbf{B}'d}. \end{aligned}$$

If \mathbf{B} equals $\beta\mathbf{I}$ then the covariance matrix function becomes $\boldsymbol{\Gamma}e^{-\beta d}$. Because $\text{var}[\mathbf{Y}(t)] = \text{var}[\mathbf{Y}(t+d)] = \boldsymbol{\Gamma}$, the autocorrelation function $\rho(d)$ is equal to $e^{-\beta d}$.

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