Trainable Videorealistic Speech Animation

Tony Ezzat*  Gadi Geiger†  Tomaso Poggio‡
Center for Biological and Computational Learning
Massachusetts Institute of Technology

Abstract

We describe how to create with machine learning techniques a generative, videorealistic, speech animation module. A human subject is first recorded using a videocamera as he/she utters a predetermined speech corpus. After processing the corpus automatically, a visual speech module is learned from the data that is capable of synthesizing the human subject’s mouth uttering entirely novel utterances that were not recorded in the original video. The synthesized utterance is re-composited onto a background sequence which contains natural head and eye movement. The final output is videorealistic in the sense that it looks like a video camera recording of the subject. At run time, the input to the system can be either real audio sequences or synthetic audio produced by a text-to-speech system, as long as they have been phonetically aligned.

The two key contributions of this paper are 1) a variant of the multidimensional morphable model (MMM) to synthesize new, previously unseen mouth configurations from a small set of mouth image prototypes; and 2) a trajectory synthesis technique based on regularization, which is automatically trained from the recorded video corpus, and which is capable of synthesizing trajectories in MMM space corresponding to any desired utterance.


Keywords: facial modeling, facial animation, morphing, optical flow, speech synthesis, lip synchronization.

1 Overview

Is it not possible to record a human subject with a video camera, process the recorded data automatically, and then re-animate that subject uttering entirely novel utterances which were not included in the original corpus? In this work, we present such a technique for achieving videorealistic speech animation.

We choose to focus our efforts in this work on the issues related to the synthesis of novel video, and not on novel audio synthesis. Thus, novel audio needs to be provided as input to our system. This audio can be either real human audio (from the same subject or a different subject), or synthetic audio produced by a text-to-speech system. All that is required by our system is that the audio be phonetically transcribed and aligned. In the case of synthetic audio from TTS systems, this phonetic alignment is readily available from the TTS system itself [Black and Taylor 1997]. In the case of real audio, publicly available phonetic alignment systems [Huang et al. 1993] may be used.

Our visual speech processing system is composed of two modules: The first module is the multidimensional morphable model (MMM), which is capable of morphing between a small set of prototype mouth images to synthesize new, previously unseen mouth configurations. The second component is a trajectory synthesis module, which uses regularization [Girosi et al. 1993] [Wahba 1990] to synthesize smooth trajectories in MMM space for any specified utterance. The parameters of the trajectory synthesis module are trained automatically from the recorded corpus using gradient descent learning.

Recording the video corpus takes on the order of 15 minutes. Processing of the corpus takes on the order of several days, but, apart from the specification of head and eye masks shown in Figure 3, is fully automatic, requiring no intervention on the part of the user. The final visual speech synthesis module consists of a small set of prototype images (46 images in the case presented here) extracted from the recorded corpus and used to synthesize all novel sequences.

Application scenarios for videorealistic speech animation include: user-interface agents for desktops, TVs, or cell-phones; digital actors in movies; virtual avatars in chatrooms; very low bitrate coding schemes (such as MPEG4); and studies of visual speech production and perception. The recorded subjects can be regular people, celebrities, ex-presidents, or infamous terrorists.

In the following section, we begin by first reviewing the relevant prior work and motivating our approach.
2 Background

2.1 Facial Modeling

One approach to model facial geometry is to use 3D methods. Parke [1974] was one of the earliest to adopt such an approach by creating a polygonal facial model. To increase the visual realism of the underlying facial model, the facial geometry is frequently scanned in using Cyberware laser scanners. Additionally, a texture-map of the face extracted by the Cyberware scanner may be mapped onto the three-dimensional geometry [Lee et al. 1995b]. Guenter [1998] demonstrated recent attempts at obtaining 3D face geometry from multiple photographs using photogrammetric techniques. Pighin et al. [1998] captured face geometry and textures by fitting a generic face model to a number of photographs. Blanz and Vetter [1999] demonstrated how a large database of Cyberware scans may be morphed to obtain face geometry from a single photograph.

An alternative to the 3D modeling approach is to model the talking face using image-based techniques, where the talking facial model is constructed using a collection of example images captured of the human subject. These methods have the potential of achieving very high levels of videorealism, and are inspired by the recent success of similar sample-based methods for audio speech synthesis [Moulines and Charpentier 1990].

Image-based facial animation techniques need to solve the video generation problem: How does one build a generative model of novel video that is simultaneously photorealistic, videorealistic, and parsimonious? Photorealism means that the novel generated images exhibit the correct visual structure of the lips, teeth, and tongue. Videorealism means that the generated sequences exhibit the correct motion, dynamics, and coarticulation effects [Cohen and Massaro 1993]. Parsimony means that the generative model is represented compactly using a few parameters.

Bregler, Covell, and Slaney [1997] describe an image-based facial animation system called Video Rewrite in which the video generation problem is addressed by breaking down the recorded video corpus into a set of smaller audiovisual basis units. Each one of these short sequences is a triphone segment, and a large database with all the acquired triphones is built. A new audiovisual sentence is constructed by concatenating the appropriate triphone sequences from the database together. Photorealism in Video Rewrite is addressed by only using recorded sequences to generate the novel video. Videorealism is achieved by using triphone contexts to model coarticulation effects. In order to handle all the possible triphone contexts, however, the system requires a library with tens and possibly hundreds of thousands of subsequences, which seems to be an overly-redundant and non-parsimonious sampling of human lip configurations. Parsimony is thus sacrificed in favor of videorealism.

Essentially, Video Rewrite adopts a decidedly agnostic approach to animation: since it does not have the capacity to generate novel lip imagery from a few recorded images, it relies on the resequencing of a vast amount of original video. Since it does not have the capacity to model how the mouth moves, it relies on sampling the dynamics of the mouth using triphone segments.

The approach used in this work presents another approach to solving the video generation problem which has the capacity to generate novel video from a small number of examples as well as the capacity to model how the mouth moves. This approach is based on the use of a multidimensional morphable model (MMM), which is capable of multidimensional morphing between various lip images to synthesize new, previously unseen lip configurations. MMM’s are powerful models of image appearance because they combine the power of vector space representations with the realism of morphing as a generative image technique. Prototype example images of the mouth are decomposed into pixel flow and pixel appearance axes that represent basis vectors of image variation. These basis vectors are combined in a multidimensional fashion to produce novel, realistic, previously unseen lip configurations.

As such, an MMM is more powerful than other vector space representations of images which do not model pixel flow explicitly. Cosatto and Graf [1998], for example, describe an approach which is similar to ours, except that their generative model involved simple pixel blending of images, which fails to produce realistic transitions between mouth configurations.

An MMM is also more powerful than simple 1-dimensional morphing between 2 image end-points [Beier and Neely 1992], as well as techniques such as those of Scott et al. [1994] [1997] and Ezzat and Poggio [2000], which morphed between several visemes in a pairwise fashion. By embedding the prototype images in a vector space, an MMM is capable of generating smooth curves through lip space which handle complex speech animation effects in a non-ad hoc manner.

2.2 Speech Animation

Speech animation techniques have traditionally included both keyframing methods and physics-based methods, and have been extended more recently to include machine learning methods. In keyframing, the animator specifies particular key-frames, and the system generates intermediate values [Parke 1974] [Pearce et al. 1986] [Cohen and Massaro 1993] [LeGoff and Benoit 1996]. In physics-based methods, the animator relies on the laws of physics to determine the mouth movement, given some initial conditions and a set of forces for all time. This technique, which requires modeling the underlying facial muscles and skin, was demonstrated quite effectively by [Waters 1987] [Lee et al. 1995b]. Finally, machine learning methods are a new class of animation tools which are trained from recorded data and then used to synthesize new motion. Examples include hidden markov models (HMMs), which were demonstrated effectively for speech animation by [Brand 1999] [Masuko et al. 1998] [Brooke and Scott 1994].

Speech animation needs to solve several problems simultaneously: firstly, the animation needs to have the correct motion, in the sense that the appropriate phonemic targets need to be realized by the moving mouth. Secondly, the animation needs to be smooth, not exhibiting any unnecessary jerks. Thirdly, it needs to display the correct dynamics: plosives such as b and p need to occur fast. Finally, speech animation needs to display the correct coarticulation effects, which determine the effects of neighboring phonemes on the current phoneme shape.

In this work, we present a trajectory synthesis module to address the issues of synthesizing mouth trajectories with correct motion, smoothness, dynamics, and coarticulation effects. This module maps from an input stream of phonemes (with their respective frame durations) to a trajectory of MMM shape-appearance parameters. This trajectory is then fed into the MMM to synthesize the final visual stream that represents the talking face.

Unlike Video Rewrite [Bregler et al. 1997], which relies on an exhaustive sampling of triphone segments to model phonetic contexts, coarticulation effects in our system emerge directly from our speech model. Each phoneme in our model is represented as a localized Gaussian target region in MMM space with a particular position and covariance. The covariance of each phoneme acts as a spring whose tension pulls the trajectory towards each phonetic region with a force proportional to observed coarticulation effects in the data.

However, unlike Cohen and Massaro [1993] (who also modeled...
coarticulation using localized Gaussian-like regions), our model of coarticulation is not hand-tuned, but rather trained from the recorded corpus itself using a gradient descent learning procedure. The training process determines the position and shape of the phonetic regions in MMM space in a manner which optimally reconstructs the recorded corpus data.

3 System Overview
An overview of our system is shown in Figure 2. After recording the corpus (Section 4), analysis is performed to produce the final visual speech module. Analysis itself consists of three sub-steps: First, the corpus is pre-processed (Section 5) to align the audio and normalize the images to remove head movement. Next, the MMM is created from the images in the corpus (Section 6.2). Finally, the corpus sequences are analyzed to produce the phonetic models used by the trajectory synthesis module (Sections 6.4 and 7.2).

Given a novel audio stream that is phonetically aligned, synthesis proceeds in three steps: First, the trajectory synthesis module is used to synthesize the trajectory in MMM space using the trained phonetic models (Section 7). Secondly, the MMM is used to synthesize the novel visual stream from the trajectory parameters (Section 6.3). Finally, the post-processing stage composites the novel mouth movement onto a background sequence containing natural eye and head movements (Section 8).

4 Corpus
An audiovisual corpus of a human subject uttering various utterances was recorded. Recording was performed at a TV studio against a blue “chroma-key” background with a standard Sony analog TV camera. The data was subsequently digitized at a 29.97 fps NTSC frame rate with an image resolution of 640 by 480 and an audio resolution of 44.1KHz. The final sequences were stored as Quicktime sequences compressed using a Sorenson coder. The recorded corpus lasts for 15 minutes, and is composed of approximately 30000 frames.

The recorded corpus consisted of 1-syllable and 2-syllable words, such as “bed” and “dagger”. A total of 152 1-syllable words and 156 2-syllable words were recorded. In addition, the corpus included 105 short sentences, such as “The statue was closed to tourists Sunday”. The subject was asked to utter all sentences in a neutral expression. In addition, the sentences themselves were designed to elicit no emotions from the subject.

5 Pre-Processing
The recorded corpus data needs to be pre-processed in several ways before it may be processed effectively for re-animation.

Firstly, the audio needs to be phonetically aligned in order to be able to associate a phoneme for each image in the corpus. We perform audio alignment on all the recorded sequences using the CMU Sphinx system [Huang et al. 1993], which is publicly available. Given an audio sequence and an associated text transcript of the speech being uttered, alignment systems use forced Viterbi search to find the optimal start and end of phonemes for the given audio sequence. The alignment task is easier than the speech recognition task because the text of the audio being uttered is known a priori.

Secondly, each image in the corpus needs to be normalized so that only movement occurring in the entire frame is the mouth movement associated with speech. Although the subject was instructed to keep her head steady during recording, residual head movement nevertheless still exists in the final recorded sequences. Since the head motion is small, we make the simplifying assumption that it can be approximated as the perspective motion of a plane lying on the surface of the face. Planar perspective deformations [Wolberg 1990] have 8 degrees of freedom, and can be inferred using 4 corresponding points between a reference frame and the current frame. We employ optical flow [Horn and Schunck 1981] [Barron et al. 1994] [Bergen et al. 1992] to extract correspondences for 640x480 pixels, and use least squares to solve the overdetermined system of equations to obtain the 8 parameters of the perspective warp. Among the 640x480 correspondences, only those lying within the head mask shown in Figure 3 are used. Pixels from the background area are not used because they do not exhibit any motion at all, and those from the mouth area exhibit non-rigid motion associated with speech.

The images in the corpus also exhibit residual eye movement and eye blinks which need to be removed. An eye mask is created (see Figure 3) which allows just the eyes from a single frame to be pasted onto the rest of the corpus imagery. The eye mask is blurred at the edges to allow a seamless blend between the pasted eyes and the rest of face.

6 Multidimensional Morphable Models
At the heart of our visual speech synthesis approach is the multidimensional morphable model representation, which is a generative
model of video capable of morphing between various lip images to synthesize new, previously unseen lip configurations.

The basic underlying assumption of the MMM is that the complete set of mouth images associated with human speech lies in a low-dimensional space whose axes represent mouth appearance variation and mouth shape variation. Mouth appearance is represented in the MMM as a set of prototype images extracted from the recorded corpus. Mouth shape is represented in the MMM as a set of optical flow vectors [Horn and Schunck 1981] computed automatically from the recorded corpus. In the work presented here, 46 images are extracted and 46 optical flow correspondences are computed. The low-dimensional MMM space is parameterized by shape parameters \( \alpha \) and appearance parameters \( \beta \).

The MMM may be viewed as a “black box” capable of performing two tasks: Firstly, given as input a set of parameters \( (\alpha, \beta) \), the MMM is capable of synthesizing an image of the subject’s face with that shape-appearance configuration. Synthesis is performed by morphing the various prototype images to produce novel, previously unseen mouth images which correspond to the input parameters \( (\alpha, \beta) \).

Conversely, the MMM can also perform analysis: given an input lip image, the MMM computes shape and appearance parameters \( (\alpha, \beta) \) that represent the position of that input image in MMM space. In this manner, it is possible to project the entire recorded corpus onto the constructed MMM, and produce a time series of \( (\alpha, \beta) \) parameters that represent trajectories of mouth motion in MMM space. We term this operation analyzing the recorded corpus.

In the following sections, we describe how a multidimensional morphable model is defined, how it may be acquired automatically from a recorded video corpus, how it may be used for synthesis, and, finally, how such a morphable model may be used for analysis.

6.1 Definition

An MMM consists of a set of prototype images \( \{I_1\}_{i=1}^N \) that represent the various lip textures that will be encapsulated by the MMM. One image is designated arbitrarily to be the reference image \( I_1 \).

Additionally, the MMM consists of a set of prototype flows \( \{C_i\}_{i=1}^N \) that represent the correspondences between the reference image \( I_1 \) and the other prototype images in the MMM. The correspondence from the reference image to itself, \( C_1 \), is designated to be an empty, zero, flow.

In this work, we choose to represent the correspondence maps using relative displacement vectors:

\[
C_i(p) = \{d_x^i(p), d_y^i(p)\}.
\]  

A pixel in image \( I_1 \) at position \( p = (x, y) \) corresponds to a pixel in image \( I_i \) at position \( (x + d_x^i(x, y), y + d_y^i(x, y)) \).

Previous methods for computing correspondence [Beier and Neely 1992] [Scott et al. 1994] [Lee et al. 1995a] adopted feature-based approaches, in which a set of high-level shape features common to both images is specified. When it is done by hand, however, this feature specification process can become quite tedious and complicated, especially in cases when a large amount of imagery is involved. In this work, we make use of optical flow [Horn and Schunck 1981] [Barron et al. 1994] [Bergen et al. 1992] algorithms to estimate this motion. This motion is captured as a two-dimensional array of displacement vectors, in the same exact format shown in Equation 1. In particular, we utilize the coarse-to-fine, gradient-based optical flow algorithms developed by [Bergen et al. 1992]. These algorithms compute the flow displacements using the spatial and temporal image derivatives. In addition, they embed the flow estimation procedure in a multiscale pyramidal framework [Burt and Adelson 1983], where initial displacement estimates are obtained at coarse resolutions, and then propagated to higher resolution levels of the pyramid.

6.2 Building an MMM

An MMM must be constructed automatically from a recorded corpus of \( \{I_1\}_{i=1}^S \) images. The two main tasks involved are to choose the image prototypes \( \{I_1\}_{i=1}^N \), and to compute the correspondence \( \{C_1\}_{i=1}^N \) between them. We discuss the steps to do this briefly below. Note that the following operations are performed on the entire face region, although they need only be performed on the region around the mouth.

6.2.1 PCA

For the purpose of more efficient processing, principal component analysis (PCA) is first performed on all the images of the recorded video corpus. PCA allows each image in the video corpus to be represented using a set of low-dimensional parameters. This set of low-dimensional parameters may thus be easily loaded into memory and processed efficiently in the subsequent clustering and Dijkstra steps.

Performing PCA using classical autocovariance methods [Bishop 1995], however, usually requires loading all the images and computing a very large autocovariance matrix, which requires a lot of memory. To avoid this, we adopt an on-line PCA method, termed EM-PCA [Roweis 1998] [Tipping and Bishop 1999], which allows us to perform PCA on the images in the corpus without loading them all into memory. EM-PCA is iterative, requiring several iterations, but is guaranteed to converge in the limit to the same principal components that would be extracted from the classical autocovariance method. The EM-PCA algorithm is typically run in this work for 10 iterations.
Performing EM-PCA produces a set of $D$ 624x472 principal components and a matrix $\Sigma$ of eigenvalues. In this work, $D = 15$ PCA bases are retained. The images in the video corpus are subsequently projected on the principal components, and each image $I_j$ is represented with a $D$-dimensional parameter vector $p_j$.

### 6.2.2 K-means Clustering

Selection of the prototype images is performed using $k$-means clustering [Bishop 1995]. The algorithm is applied directly on the $\{(p_1, \ldots, p_N)\}$ low dimensional PCA parameters, producing $N$ cluster centers. Typically the cluster centers extracted by k-means clustering do not coincide with actual image datapoints, so the nearest images in the dataset to the computed cluster centers are chosen to be the final image prototypes $\{I_i\}_{i=1}^N$ for use in our MMM.

It should be noted that $k$-means clustering requires the use of an internal distance metric with which to compare distances between datapoints and the chosen cluster centers. In our case, since the image parameters are themselves produced by PCA, the appropriate distance metric between two points $p_m$ and $p_n$ is the Mahalanobis distance metric:

$$d(p_m, p_n) = (p_m - p_n)^T \Sigma^{-1}(p_m - p_n)$$  \hspace{1cm} (2)

where $\Sigma$ is the afore-mentioned matrix of eigenvalues extracted by the EM-PCA procedure.

We selected $N = 46$ image prototypes in this work, which are partly shown in Figure 4. The top left image is the reference image $I_1$. There is nothing magical about our choice of 46 prototypes, which is in keeping with the typical number of visemes other researchers have used [Scott et al. 1994] [Ezzat and Poggio 2000]. It should be noted, however, that the 46 prototypes have no explicit relationship to visemes, and instead form a simple basis set of image textures.

### 6.2.3 Dijkstra

After the $N = 46$ image prototypes are chosen, the next step in building an MMM is to compute correspondence between the reference image $I_1$ and all the other prototypes. Although it is in principle possible to compute direct optical flow between the images, we have found that direct application of optical flow is not capable of estimating good correspondence when the underlying lip displacements between images are greater than 5 pixels.

It is possible to use flow concatenation to overcome this problem. Since the original corpus is digitized at 29.97 fps, there are many intermediate frames that lie between the chosen prototypes. A series of consecutive optical flow vectors between each intermediate image and its successor may be computed and concatenated into one large flow vector that defines the global transformation between the chosen prototypes (see Appendix A for details on flow concatenation).

Typically, however, prototype images are very far apart in the recorded visual corpus, so it is not practical to compute concatenated optical flow between them. The repeated concatenation that would be involved across the hundreds or thousands of intermediate frames leads to a considerably degraded final flow.

To compute good correspondence between prototypes, a method is needed to figure out how to compute the path from the reference example $I_1$ to the chosen image prototypes $I_i$ without repeated concatenation over hundreds or thousands of intermediate frames. We accomplish this by constructing the corpus graph representation of the corpus: A corpus graph is an S-by-S sparse adjacency graph matrix in which each frame in the corpus is represented as a node in a graph connected to $k$ nearest images. The $k$ nearest images are chosen using the $k$-nearest neighbors algorithm [Bishop 1995], and the distance metric used is the Mahalanobis distance in Equation 2 applied to the PCA parameters $p_i$. Thus, an image is connected in the graph to the $k$ other images that look most similar to it. The edge-weight between a frame and its neighbor is the value of the Mahalanobis distance. We set $k = 20$ in this work.

After the corpus graph is computed, the Dijkstra shortest path algorithm [Cormen et al. 1989] [Tenenbaum et al. 2000] is used to compute the shortest path between the reference example $I_1$ and the other chosen image prototypes $I_i$. Each shortest path produced by the Dijkstra algorithm is a list of images from the corpus that cumulatively represent the shortest deformation path from $I_1$ to $I_i$ as measured by the Mahalanobis distance. Concatenated flow from $I_1$ to $I_i$ is then computed along the intermediate images produced by the Dijkstra algorithm. Since there are 46 images, $N = 46$ correspondences $\{C_i\}_{i=1}^N$ are computed in this fashion from the reference image $I_1$ to the other image prototypes $\{I_i\}_{i=1}^N$.

### 6.3 Synthesis

The goal of synthesis is to map from the multidimensional parameter space $(\alpha, \beta)$ to an image which lies at that position in MMM space. Since there are 46 correspondences, $\alpha$ is a 46-dimensional parameter vector that controls mouth shape. Similarly, since there are 46 image prototypes, $\beta$ is a 46-dimensional parameter vector that controls mouth texture. The total dimensionality of $(\alpha, \beta)$ is 92.

Synthesis first proceeds by synthesizing a new correspondence $C_1^{\text{synth}}$ using linear combination of the prototype flows $C_i$:

$$C_1^{\text{synth}} = \sum_{i=1}^N \alpha_i C_i.$$  \hspace{1cm} (3)

The subscript 1 in Equation 3 above is used to emphasize that $C_1^{\text{synth}}$ originates from the reference image $I_1$, since all the prototype flows are taken with $I_1$ as reference.

Forward warping may be used to push the pixels of the reference image $I_1$ along the synthesized correspondence vector $C_1^{\text{synth}}$. Notationally, we denote the forward warping operation as an operator $W(I, C)$ that operates on an image $I$ and a correspondence map $C$ (see Appendix B for details on forward warping).

However, a single forward warp will not utilize the image texture from all the examples. In order to take into account all image texture, a correspondence re-orientation procedure first described...
Analysis-by-synthesis algorithm is used to synthesize an image. However, the blending of multiple images in the MMM using the shown in Figure 6 are some of the synthetic images produced by [Jones and Poggio 1998] [Blanz and Vetter 1999], the synthesis algorithm is used in various forms in [Beymer et al. 1993] is adopted that re-orient the synthesized correspondence vector \( C_i^{\text{synth}} \) so that it originates from each of the other example images \( I_i \). Reorientation of the synthesized flow \( C_i^{\text{synth}} \) proceeds in two steps, shown figuratively in Figure 5. First, \( C_i \) is subtracted from the synthesized flow \( C_i^{\text{synth}} \) to yield a flow that contains the correct flow geometry, but which originates from the reference example \( I_i \), rather than the desired example image \( I_j \). Secondly, to move the flow into the correct reference frame, this flow vector is itself warped along \( C_i \). The entire re-orientation process may be denoted as follows:

\[
C_i^{\text{synth}} = W(C_i^{\text{synth}} - C_i, C_i). \tag{4}
\]

Re-orientation is performed for all examples in the example set.

The third step in synthesis is to warp the prototype images \( I_i \) along the re-oriented flows \( C_i^{\text{synth}} \) to generate a set of \( N \) warped image textures \( I_i^{\text{warped}} \).

\[
I_i^{\text{warped}} = W(I_i, C_i^{\text{synth}}). \tag{5}
\]

The fourth and final step is to blend the warped images \( I_i^{\text{warped}} \) using the \( \beta \) parameters to yield the final morphed image:

\[
I_i^{\text{morph}} = \sum_{i=1}^{N} \beta_i I_i^{\text{warped}}. \tag{6}
\]

Combining Equations 3 through 6 together, our MMM synthesis may be written as follows:

\[
I_i^{\text{morph}}(\alpha, \beta) = \sum_{i=1}^{N} \beta_i W(I_i, W(\sum_{j=1}^{N} \alpha_j C_j - C_i, C_i)). \tag{7}
\]

Empirically we have found that the MMM synthesis technique is capable of surprisingly realistic re-synthesis of lips, teeth, and tongue. However, the blending of multiple images in the MMM for synthesis tends to blur out some of the finer details in the teeth and tongue (See Appendix C for a discussion of synthesis blur). Shown in Figure 6 are some of the synthetic images produced by our system, along with their real counterparts for comparison.

6.4 Analysis

The goal of analysis is to project the entire recorded corpus \( \{I_i\}_{i=1}^{N} \) onto the constructed MMM, and produce a time series of \( (\alpha, \beta) \) parameters that represent trajectories of the original mouth motion in MMM space.

One possible approach for analysis of images is to perform analysis-by-synthesis: In this approach, used in various forms in [Jones and Poggio 1998] [Blanz and Vetter 1999], the synthesis algorithm is used to synthesize an image \( I_i^{\text{synth}}(\alpha, \beta) \), which is then compared to the novel image using an error metric (ie, the L2 norm). Gradient-descent is then usually performed to change the parameters in order to minimize the error, and the synthesis process is repeated. The search ends when a local minimum is achieved. Analysis-by-synthesis, however, is very slow in the case when a large number of images are involved.

In this work we choose another method that is capable of extracting parameters \( (\alpha, \beta) \) in one iteration. In addition to the image \( I_i^{\text{novel}} \) to be analyzed, the method requires that the correspondence \( C_i^{\text{novel}} \) from the reference image \( I_i \) in the MMM to the novel image \( I_i^{\text{novel}} \) be computed beforehand. In our case, most of the novel imagery to be analyzed will be from the recorded video corpus itself, so we employ the Dijkstra approach discussed in Section 6.2.3 to compute good quality correspondences between the reference image \( I_i \) and \( I_i^{\text{novel}} \).

Given a novel image \( I_i^{\text{novel}} \) and its associated correspondence \( C_i^{\text{novel}} \), the first step of the analysis algorithm is to estimate the parameters \( \alpha \) which minimize

\[
\| C_i^{\text{novel}} - \sum_{i=1}^{N} \alpha_i C_i \|. \tag{8}
\]

This is solved using the pseudo-inverse:

\[
\alpha = (C^T C)^{-1} C^T C_i^{\text{novel}} \tag{9}
\]

where \( C \) above is a matrix containing all the prototype correspondences \( \{C_i\}_{i=1}^{N} \).

After the parameters \( \alpha \) are estimated, \( N \) image warps are synthesized in the same manner as described in Section 6.3 using flow-reorientation and warping:

\[
I_i^{\text{warped}} = W(I_i, W(\sum_{i=1}^{N} \alpha_i C_j - C_i, C_i)). \tag{10}
\]

The final step in analysis is to estimate the values of \( \beta \) as the values which minimize

\[
\| I_i^{\text{novel}} - \sum_{i=1}^{N} \beta_i I_i^{\text{warped}} \| \quad \text{subject to} \quad \beta_i > 0 \ \forall i \ \text{and} \ \sum_{i=1}^{N} \beta_i = 1. \tag{11}
\]
The non-negativity constraint above on the \( \beta_i \) parameters ensures that pixel values are not negated. The normalization constraint ensures that the \( \beta_i \) parameters are computed in a normalized manner for each frame, which prevents brightness flickering during synthesis. The form of the imposed constraints cause the computed \( \beta_i \) parameters to be sparse (see Figure 7), which enables efficient synthesis by requiring only a few image warps (instead of the complete set of 46 warps). Equation 11, which involves the minimization of a quadratic cost function subject to constraints, is solved using quadratic programming methods. In this work, we use the Matlab function \texttt{quadprog}.

Each utterance in the corpus is analyzed with respect to the 92-dimensional MMM created in Section 6.2, yielding a set of \( z_i = (\alpha_i, \beta_i) \) parameters for each utterance. Analysis takes on the order of 15 seconds per frame on a circa 1998 450 MHz Pentium II machine. Shown in Figure 9 in solid blue are example analyzed trajectories for \( \alpha_{12} \) and \( \beta_{28} \) computed for the word \texttt{tabloid}.

7 Trajectory Synthesis

7.1 Overview

The goal of trajectory synthesis is to map from an input phone stream \( \{p_i\} \) to a trajectory \( y_i = (\alpha_i, \beta_i) \) of parameters in MMM space. After the parameters are synthesized, Equation 7 from Section 6.3 is used to create the final visual stream that represents the talking face.

The phone stream is a stream of phonemes \( \{p_i\} \) representing that phonetic transcription of the utterance. For example, the word \texttt{one} may be represented by a phone stream \( \{p_i\}^{15} = \{w, w, w, w, \ddash, \ddash, u_h, n, n, \ddash, \ddash, \ddash\} \). Each element in the phone stream represents one image frame. We define \( T \) to be the length of the entire utterance in frames.

Since the audio is aligned, it is possible to examine all the flow and texture parameters for any particular phoneme. Shown in Figure 8 are histograms for the \( \alpha_i \) parameter for the \( \{w, \ddash, \ddash\} \) and \( \{a, w\} \) phones. Evaluation of the analyzed parameters from the corpus reveals that parameters representing the same phoneme tend to cluster in MMM space. We represent each phoneme \( p \) mathematically as a multidimensional Gaussian with mean \( \mu_p \) and diagonal covariance \( \Sigma_p \). Separate means and covariances are estimated for the flow and texture parameters $^{1}$.

The trajectory synthesis problem is framed mathematically as a regularization problem [Girois et al. 1993] [Wahba 1990]. The goal is to synthesize a trajectory \( y \) which minimizes an objective function \( E \) consisting of a target term and a smoothness term:

\[
E = (y - \mu)^T D^T \Sigma^{-1} D (y - \mu) + \lambda y^T W^T W y. \tag{12}
\]

The desired trajectory \( y \) is a vertical concatenation of the individual \( y_i = (\alpha_i, \beta_i) \) terms at each time step (or \( y_i = (\alpha_i, \beta_i) \), since we treat flow and texture parameters separately):

\[
y = \begin{bmatrix}
y_1 & \cdots & y_T
\end{bmatrix}
\]

The target term consists of the relevant means \( \mu \) and covariances \( \Sigma \) constructed from the phone stream:

\[
\mu = \begin{bmatrix} \mu_{p_1} & \cdots & \mu_{p_T} \end{bmatrix}, \Sigma = \begin{bmatrix} \Sigma_{p_1} & \cdots & \Sigma_{p_T} \end{bmatrix} \tag{14}
\]

The matrix \( D \) is a duration-weighting matrix which emphasizes the shorter phonemes and de-emphasizes the longer ones, so that the objective function is not heavily skewed by the phonemes of longer duration:

\[
D = \begin{bmatrix}
\sqrt{1 - \frac{D_{p_1}}{D}} & \cdots & \cdots & \cdots & \cdots & \sqrt{1 - \frac{D_{p_T}}{D}}
\end{bmatrix}
\tag{15}
\]

One possible smoothness term consists of the first order difference operator:

\[
W = \begin{bmatrix}
-1 & 1 & \cdots & \cdots & \cdots & -1 & 1
\end{bmatrix}
\tag{16}
\]

Higher orders of smoothness are formed by repeatedly multiplying \( W \) with itself: second order \( W^T W^T \), third order \( W^T W^T W^T \), and so on.

Finally, the regularizer \( \lambda \) determines the trade-off between both terms.

Taking the derivative of Equation 12 and minimizing yields the following equation for synthesis:

\[
(D^T \Sigma^{-1} D + \lambda W^T W)y = D^T \Sigma^{-1} D \mu. \tag{17}
\]

Given known means \( \mu \), covariances \( \Sigma \), and regularizer \( \lambda \), synthesis is simply a matter of plugging them into Equation 17 and solving for \( y \) using Gaussian elimination. This is done separately for the flow and the texture parameters. In our experiments a regularizer of degree four yielding multivariate additive septic splines [Wahba 1990] gave satisfactory results (see next subsection).

\footnote{Technically, since the texture parameters are non-negative, they are best modeled using Gamma distributions not Gaussians. In that case, Equation 12 needs to be re-written for Gamma distributions. In practice, however, we have found Gaussians to work well enough for texture parameters.}
Coarticulation effects in our system are modeled via the magnitude of the variance $\Sigma_p$ for each phoneme. Small variance means the trajectory must pass through that region in phoneme space, and hence neighboring phonemes have little coarticulatory effect. On the other hand, large variance means the trajectory has a lot of flexibility in choosing a path through a particular phonetic region, and hence it may choose to pass through regions which are closer to a phoneme’s neighbors. The phoneme will thus experience large coarticulatory effects.

There is no explicit model of phonetic dynamics in our system. Instead, phonetic dynamics emerge implicitly through the interplay between the magnitude of the variance $\Sigma_p$ for each phoneme (which determines the phoneme’s “spatial” extent), and the input phoneme stream (which determines the duration in time of each phoneme). Equation 12 then determines the speed through a phonetic region in a manner which balances nearness to the phoneme with smoothness of the overall trajectory. In general, we find the trajectories speed up in regions of small duration and small variance (e.g. plosives), while they slow down in regions of large duration and large variance (e.g. silences).

### 7.2 Training

The means $\mu_p$ and covariances $\Sigma_p$ for each phone $p$ are initialized directly from the data using sample means and covariances. However, the sample estimates tend to average out the mouth movement so that it looks under-articulated. As a consequence, there is a need to adjust the means and variances to better reflect the training data.

Gradient descent learning [Bishop 1995] is employed to adjust the mean and covariances. First, the Euclidean error metric is chosen to represent the error between the original utterance $z$ and the synthetic utterance $y$:

$$E = (z - y)^T(z - y).$$

The parameters $\{\mu_p, \Sigma_p\}$ need to be changed to minimize this objective function $E$. The chain rule may be used to derive the relationship between $E$ and the parameters:

$$\frac{\partial E}{\partial \mu_i} = \left(\frac{\partial E}{\partial y}\right)^T \frac{\partial y}{\partial \mu_i}$$

$$\frac{\partial E}{\partial \sigma_{ij}} = \left(\frac{\partial E}{\partial y}\right)^T \frac{\partial y}{\partial \sigma_{ij}}.$$  

Finally, gradient descent is performed by changing the previous values of the parameters according to the computed gradient:

$$\mu_i^{\text{new}} = \mu_i^{\text{old}} - \eta \frac{\partial E}{\partial \mu_i}$$

$$\Sigma_{ij}^{\text{new}} = \Sigma_{ij}^{\text{old}} - \eta \frac{\partial E}{\partial \Sigma_{ij}}.$$  

Cross-validation sessions were performed to evaluate the appropriate value of $\lambda$ and the correct level of smoothness $W$ to use. The learning rate $\eta$ was set to 0.00001 for all trials, and 10 iterations performed. Comparison between batch and online updates indicated that online updates perform better, so this method was used throughout training. Testing was performed on a set composed of 1-syllable words, 2-syllable words, and sentences not contained in the training set. The Euclidean norm between the synthesized trajectories and the original trajectories was used to measure error. The results showed that the optimal smoothness operator is fourth order and the optimal regularizer is $\lambda = 1000$. Figure 9 depicts synthesized trajectories for the $\alpha_{12}$ and $\beta_{28}$ parameters before training (in green dots) and after training (in red crosses) for these optimal values of $W$ and $\lambda$.

### 8 Post-Processing

Due to the head and eye normalization that was performed during the pre-processing stage, the final animations generated by our system exhibit movement only in the mouth region. This leads to an unnerving “zombie”-like quality to the final animations. As in [Cosatto and Graf 1998] [Bregler et al. 1997], we address this issue by compositing the synthesized mouth onto a background sequence which contains natural head and eye movement.

The first step in the composition process is to add Gaussian noise to the synthesized images to regain the camera image sensing noise that is lost as a result of blending multiple image prototypes in the MMM. We estimate means and variances for this noise by computing differences between original images and images synthesized by our system, and averaging over 200 images.

After noise is added, the synthesized sequences are composited onto the chosen background sequence with the help of the masks.
shown in Figure 3. The head mask is first forward warped using optical flow to fit across the head of each image of the background sequence. Next, optical flow is computed between each background image and its corresponding synthetic image. The synthetic image and the mouth mask from Figure 3 are then perspective-warped back onto the background image. The perspective warp is estimated using only the flow vectors lying within the background head mask. The final composite is made by pasting the warped mouth onto the background image using the warped mouth mask. The mouth mask is smoothed at the edges to perform a seamless blend between the background image and the synthesized mouth. The compositing process is depicted in Figure 10.

9 Computational Issues

To use our system, an animator first provides phonetically annotated audio. The annotation may be done automatically [Huang et al. 1993], semi-automatically using a text transcript [Huang et al. 1993], or manually [Sjlander and Beskow 2000].

Trajectory synthesis is performed by Equation 17 using the trained phonetic models. This is done separately for the flow and the texture parameters. After the parameters are synthesized, Equation 7 from Section 6.3 is used to create the visual stream with the desired mouth movement. Typically only the image prototypes \( I_i \), which are associated with top 10 values of \( \beta \), are warped, which yields a considerable savings in computation time. MMM synthesis takes on the order of about 7 seconds per frame for an image resolution of 624x472. The background compositing process adds on a few extra seconds of processing time. All times are computed on a 450 MHz Pentium II.

10 Evaluation

We have synthesized numerous examples using our system, spanning the entire range of 1-syllable words, 2-syllable words, short sentences, and long sentences. In addition, we have synthesized songs and foreign speech examples.

Experimentally we have found that reducing the number of prototypes below 30 degrades the quality of the final animations. An open question is whether increasing the number of prototypes significantly beyond 46 will lead to even higher levels of videorealism.

In terms of corpus size, it is possible to optimize the spoken corpus so that several words alone elicit the 46 prototypes. This would reduce the duration of the corpus from 15 minutes to a few seconds. However, this would degrade the quality of the correspondences computed by the Dijkstra algorithm. In addition, the phonetic training performed by our trajectory synthesis module would degrade as well. Further systematic experiments need to be made in order to evaluate how final performance changes with the size of the corpus.

We evaluated our results by performing three different visual “Turing tests” to see whether human subjects can distinguish between real sequences and synthetic ones. In the first experiment (“single presentation”), subjects were asked to view one visual sequence at a time, and identify whether it is real or synthetic. In a similar second experiment (“fast single presentation”), the subjects were asked to make the judgments in a fast manner while the utterances were being presented without pauses in between. In a third experiment (“double presentation”), the subjects were asked to view pairs of the same utterance, where one item in the pair is real and the other is synthetic (but randomly ordered). The subjects in this experiment were asked to identify which utterance in the pair is real, and which is synthetic. 16 or 18 utterances were presented to each subject, with half being real and half being synthetic. As seen from Table 1, performance in all three experiments was close to chance level (50%) and not significantly different from it.

Finally, we also evaluated our system by performing intelligibility tests in which subjects were asked to lip read a set of natural and

<table>
<thead>
<tr>
<th>Experiment</th>
<th># subjects</th>
<th>% correct</th>
<th>t</th>
<th>p&lt;</th>
</tr>
</thead>
<tbody>
<tr>
<td>Single pres.</td>
<td>22</td>
<td>54.3%</td>
<td>1.243</td>
<td>0.3</td>
</tr>
<tr>
<td>Fast single pres.</td>
<td>21</td>
<td>52.1%</td>
<td>0.619</td>
<td>0.5</td>
</tr>
<tr>
<td>Double pres.</td>
<td>22</td>
<td>46.6%</td>
<td>-0.75</td>
<td>0.5</td>
</tr>
</tbody>
</table>

Table 1: Levels of correct identification of real and synthetic sequences. “t” represents the value from a standard t-test with significance level indicated in the “p<” column.
Details on all experiments are forthcoming in a separate article.

11 Further Work

The main limitation of our technique is the difficulty of re-compositing synthesized mouth sequences into background sequences which involve 1) large changes in head pose, 2) changes in lighting conditions, and 3) changes in viewpoint. All these limitations can be alleviated by extending our approach from 2D to 3D. It is possible to envision a real-time 3D scanner that is capable of recording a 3D video corpus of speech. Alternatively, techniques such as those presented in [Guenter et al. 1998] [Pighin et al. 1998] [Blanz and Vetter 1999] can be used to map a 2D video corpus into 3D.

The geodesic trajectory synthesis equations described by Brand et al. [1999] [2000] are analogous (and more sophisticated than the trajectory synthesis techniques we use (Equations 12 and 17). Although those equations require considerably more training data, it is possible they could lead to higher levels of videorealism.

Clearly the face is used as a conduit to transmit emotion, so one possible avenue to explore is the synthesis of speech with various emotional states. It is possible to record various corpora under different emotional states and create MMMs for each state. During synthesis, the appropriate MMM is selected. An open question is whether they could lead to higher levels of videorealism.

The authors would like to dedicate this work in the memory of Christian Benoît [LeGoff and Benoît 1996] who was a pioneer in audiovisual speech research. The authors would like to thank Bereisdorf and Dynasty Models; Craig Milanesi, Dave Konstine, Christian Benoit [LeGoff and Benoit 1996] who was a pioneer in audiovisual speech research. The authors would like to thank Bereisdorf and Dynasty Models; Craig Milanesi, Dave Konstine, Christian Benoit from MIT Video Productions; Marypat Fitzgerald and Terzopoulos, Ryan Rifkin, and anonymous reviewers for countless helpful comments on the paper; David Beymer, Mike Jones, Vinay Kumar, Steve Lines, Roberto Brunelli and Steve Librande for laying the groundwork.

This work was partially funded by the Association Christian Benoît, NTT Japan, Office of Naval Research, DARPA, National Science Foundation (Adaptive Man-Machine Interfaces). Additional support was provided by: Central Research Institute of Electric Power Industry, Eastman Kodak Company, DaimlerChrysler AG, Honda R&D Co., Ltd., Komatsu Ltd., Toyota Motor Corporation and The Whitaker Foundation.

12 Acknowledgments

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References


Appendix: Flow Concatenation

Given a series of consecutive images $I_0, I_1, \ldots, I_n$, we would like to construct the correspondence map $C(0,n)$ relating $I_0$ to $I_n$. We focus on the case of the 3 images $I_{i-1}, I_i, I_{i+1}$ since the concatenation algorithm is simply an iterative application of this 3-frame base case. Optical flow is first computed between the consecutive frames to yield $C_{i-1,i}$. Note that it is not correct to construct $C_{i,i+1}$ as the simple addition of $C_{i-1,i} + C_{i+1,i}$ because the two flow fields are with respect to two different reference images. Vector addition needs to be performed with respect to a common origin.

Our concatenation thus proceeds in two steps: to place all vector fields in the same reference frame, the correspondence map $C_{i,i+1}$ itself is warped backwards [Wolberg 1990] along $C_{(i-1),i}$ to create $C_{warp}$. Now $C_{warp}$ and $C_{i,i+1}$ are both added to produce an approximation to the desired concatenated correspondence:

$$C_{i,(i-1)+1} = C_{i,(i-1)} + C_{warp,i,i+1}.$$  (26)

A procedural version of our backwarp warp is shown in Figure 11. $\text{BILINEAR}$ refers to bilinear interpolation of the 4 pixel values closest to the point $(x, y)$.

B Appendix: Forward Warping

Forward warping may be viewed as “pushing” the pixels of an image $I$ along the computed flow vectors $C$. We denote the forward warping operation as an operator $W(I, C)$ that operates on an image $I$ and a correspondence map $C$, producing a warped image $I_{warp}$ as final output. A procedural version of our forward warp is shown in Figure 12.

It is also possible to forward warp a correspondence map $C'$ along another correspondence map $C$, which we denote as $W(C', C)$. In this scenario, the $x$ and $y$ components of $C'(p) = \{d_x(p), d_y(p)\}$ are treated as separate images, and warped individually along $C$: $W(dx, C)$ and $W(dy, C)$.

C Appendix: Hole-Filling

Forward warping produces black holes which occur in cases where a destination pixel was not filled in with any source pixel value. This occurs due to inherent nonzero divergence in the optical flow, particularly around the region where the mouth is expanding. To remedy this, a hole-filling algorithm [Chen and Williams 1993] was adopted which pre-fills a destination image with a special reserved background color. After warping, the destination image is traversed in rasterized order and the holes are filled in by interpolating linearly between their non-hole endpoints.

In the context of our synthesis algorithm in Section 6.3, hole-filling can be performed before blending, or after blending. Throughout this paper, we assume hole-filling is performed before blending, which allows us to subsume the hole-filling procedure into our forward warping operator $W$ and simplify our notation. Consequently (as in Equation 6), the blending operation becomes a simple linear combination of the hole-filled warped intermediates $I_{i,warp}$.

In practice, however, we perform hole-filling after blending, which reduces the size of the holes that need to be filled, and leads to a considerable reduction in synthesis blur. Post-blending hole-filling requires a more complex blending algorithm than as noted in Equation 6 because the blending algorithm now needs to keep track of holes and non-holes in the warped intermediate images $I_{i,warp}$.

$$p_{morph}(x,y) = \sum_{x,y \notin \text{hole}} \beta_{x,y} I_{warp}(x,y)$$

Typically an accumulator array is used to keep track of the denominator term in Equation 27 above. The synthesized mouth images shown in Figure 6 were generated using post-blending hole-filling.