

# StyleT2V0: Zero-Shot Short Video Generation Guided by Text and Style

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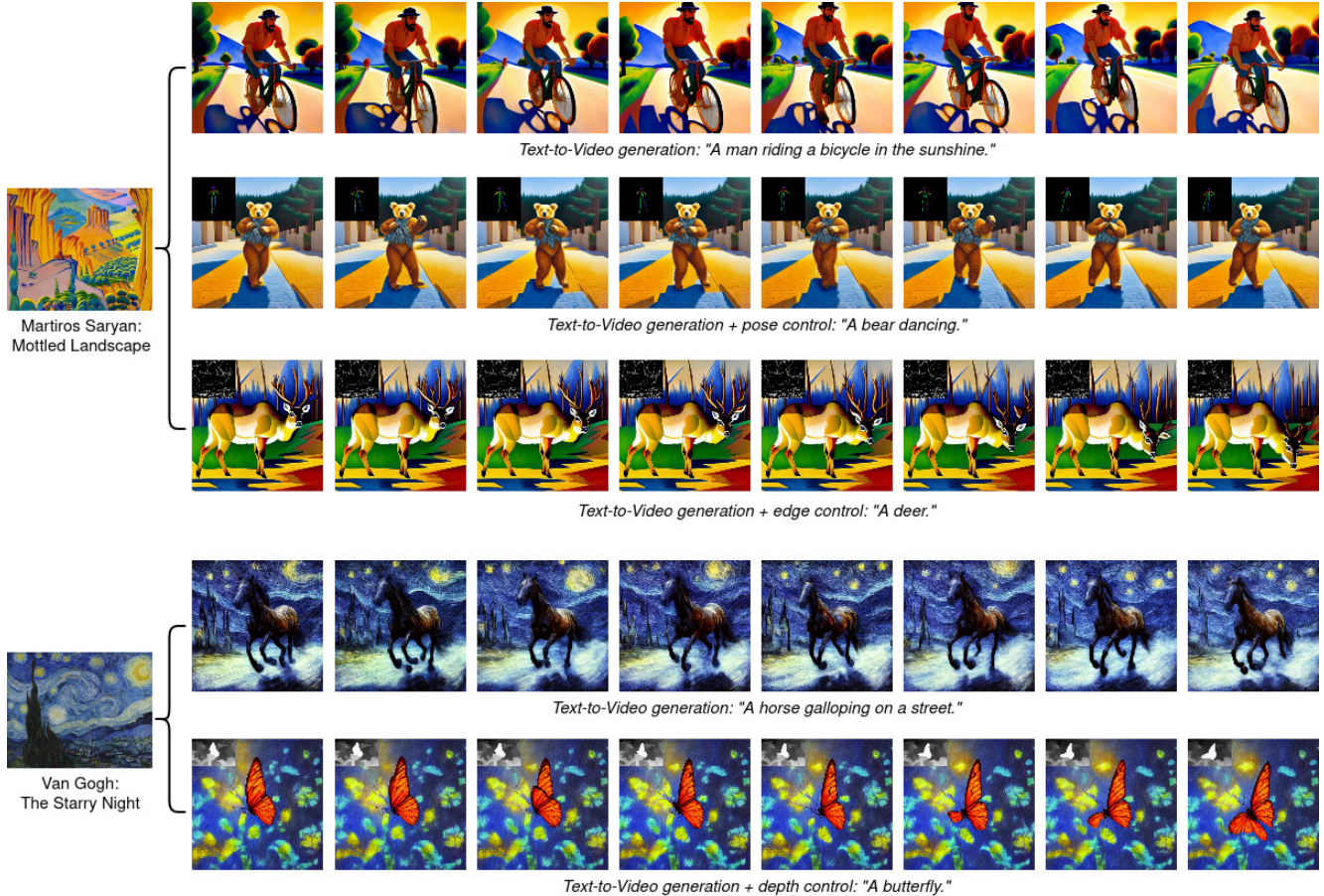


Figure 1. Our method **StyleT2V0** advances video generation capabilities by utilizing: (i) a textual prompt combined with a style image to capture the ensure stylistic and contextual accuracy across frames (see rows 1, 4), (ii) pose control to ensure dynamic accuracy, (iii) edge control for enhanced visual definition, and (iv) depth control to enrich spatial perception. The showcased frames illustrate the method’s robust ability to adaptively integrate enhanced controls while consistently reflecting the artistic intent and textual specifications.

## Abstract

Recent advancements in text-to-image diffusion models have markedly expanded the potential for generating images based on textual descriptions. Yet, the extension of these capabilities to video generation poses significant challenges, particularly in achieving stylistic consistency

across video frames. Traditional approaches often require extensive training processes and face difficulties in maintaining a coherent artistic style throughout the sequence, especially when attempting to emulate a **specific style referenced by an image**.

To bridge this gap, we introduce StyleT2V0, a zero-shot methodology that enables the generation of videos directly

*from textual prompts while ensuring stylistic alignment with a designated reference image. By integrating shared attention mechanisms within the generation process, our approach guarantees that each frame not only adheres contextually to the textual description but also maintains stylistic coherence with the reference image, all without the need for additional model training. This method represents a significant step forward in simplifying the video production process from textual descriptions, offering a streamlined and efficient solution for creating stylistically consistent video content.*

## 1. Introduction

The advent of generative artificial intelligence (AI) has marked a new era of digital content creation, enabling the synthesis of complex and creative outputs from simple inputs. However, the field of text-to-video generation, despite its potential, grapples with significant challenges due to the inherently complex nature of videos and the intricate relationship between text and visual content. Ensuring temporal and contextual consistency in video generation remains a major hurdle, as videos require each frame to be visually coherent with adjacent frames and aligned with the text prompt provided. Furthermore, maintaining style consistency across all frames introduces an additional layer of complexity.

Current text-to-video generation systems often struggle to capture the desired style effectively. This is partly because conveying style through textual descriptions lacks the explicitness and spatial specificity provided by visual references. The StyleCrafter [14] introduces several innovative strategies to overcome existing limitations in video generation. It leverages a style control adapter, trained with image datasets, to infuse style features from a reference image into video content. It incorporates a Scale-Adaptive Fusion Module, which balances text-based content features with image-based style features. Despite its advanced capabilities, StyleCrafter requires extensive training. It uses a style control adapter that must be trained with image datasets to effectively transfer style features from reference images into the generated video. The method also involves fine-tuning of the T2V models to improve their temporal dynamics and ensure consistent style throughout the video, enhancing the overall quality.

Other techniques in the field utilize various forms of style transfer technology, embedding style features directly within the generative process. These methods typically employ deep learning architectures that integrate style information at different layers of the network, allowing for the dynamic adjustment of style elements during the video generation process. However, these methods can suffer from high computational costs and may require significant tun-

ing to balance style preservation with the visual output.

Despite these challenges, there are promising opportunities for advancements in text-to-video generation technology. Leveraging zero-shot learning could potentially enable models to generate content without prior specific training on similar tasks, thus enhancing the model’s flexibility and applicability. Additionally, improving methods for consistent style information throughout the video could lead to more cohesive and aesthetically pleasing outputs. In response to these challenges, we introduce StyleT2V0, a zero-shot and versatile approach that facilitates the creation of videos in any desired style using a reference image (see Fig. ??? [this will be teaser figure] and further results). This method offers two key benefits: it enhances the styling capabilities of text-to-video (T2V) models without prior training and ensures a more precise representation of the intended style than text descriptions alone, leading to videos that are not only consistent but also uniform in stylistic execution.

The experiments demonstrate that StyleT2V0 successfully generates temporally consistent short videos from text while ensuring stylistic alignment with a specified reference image. To summarize, our contributions are three-fold:

- A novel attention sharing mechanism to simultaneously enforce style guidance and ensure frame consistency during generation in a zero-shot manner.
- A range of applications demonstrating the effectiveness of our approach, such as generating videos conditioned on a style image alongside edge, pose, or depth maps.
- Our approach demonstrates superior performance in terms of color transformation from style image while preserving the text alignment.

## 2. Related Work

### 2.1. Text-to-Image Generation

The field of text-to-image synthesis has witnessed transformative advancements with the development of sophisticated models that generate highly realistic and semantically accurate images directly from textual descriptions. Generative Adversarial Networks (GANs) [6] have been pivotal in pioneering early developments in text-to-image generation. Innovations such as StackGAN [30] and AttnGAN [29] have laid foundational work by using stacked generative models and attention mechanisms, respectively, to progressively refine images from coarse-to-fine details based on textual inputs. These models excelled in generating visually appealing images that closely align with the given text descriptions, setting a new benchmark for image quality and text alignment in early stages. The adoption of transformers, initially designed for text processing, into image synthesis marked a significant leap. Models like DALL-E [18] from OpenAI leveraged transformers to handle complex, abstract concepts and create images with unprecedented creativity

from textual prompts. The model’s ability to parse and interpret nuanced text has enabled it to generate images that are not only high in quality but also rich in context and imagination. Diffusion models have recently revolutionized the text-to-image landscape by offering another layer of sophistication. Unlike GANs, diffusion models convert noise into structured images gradually, controlled by the conditioning provided by text. This process has proven to produce images with exceptional detail and lower rates of artifacts. GLIDE [16] and Imagen [23] are notable examples that use diffusion techniques to achieve remarkable levels of realism and accuracy, outperforming previous methodologies in various benchmarks. Advancing further, Latent Diffusion Models (LDMs) [20] operate by encoding images in a compressed latent space before applying the diffusion process, significantly enhancing computational efficiency. This method maintains the high quality of generated images while reducing the resource intensity typically associated with such processes. The use of latent space allows these models to manage and manipulate higher-resolution images more effectively, facilitating broader applications.

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## 2.2. Text-to-Video Generation

Text-to-video generation continues to be a vibrant and challenging area of research within artificial intelligence, aiming to transform textual descriptions into dynamic visual narratives. This field has evolved significantly, initially leveraging generative adversarial networks (GANs) [6] and vector quantized variational autoencoders (VQ-VAE) [26] to more advanced methodologies employing transformer-based architectures and diffusion models.

The utilization of autoregressive transformers has profoundly impacted the development of this technology. Models such as NUWA [28] and Phenaki [27] are at the forefront, with NUWA implementing a 3D transformer encoder-decoder framework that supports both image and video generation from text. Phenaki advances this paradigm using a bidirectional masked transformer that employs a causal attention mechanism, enabling the synthesis of extended video sequences from concise text descriptions. These advancements underscore the versatility of transformers in handling the complexities of video content generation.

In parallel, diffusion models have brought substantial enhancements to the fidelity and coherence of video outputs. CogVideo [12] extends the principles of the CogView2 [3] model, incorporating structured multi-frame-rate hierarchical training strategies to refine the alignment of text and video. Large-scale diffusion models like those used in Video Diffusion Models (VDM) [11] and Imagen Video [10] apply cascading diffusion models tailored for video, achieving high-resolution outputs that maintain temporal consistency.

Further expanding the model ecosystem, approaches like Make-A-Video [24] and Gen-1 [5] explore unsupervised learning paradigms and content-guided video editing, respectively. Make-A-Video leverages existing text-to-image architectures, applying them to video synthesis in an unsupervised manner. Conversely, Gen-1 introduces a novel framework for video editing that is guided by structural and content-based descriptions, facilitating refined control over the generated video content.

The introduction of diffusion transformers has further revolutionized video generation, leading to solutions like Latte [15] and Sora [1], which can produce minute-long videos of high visual quality that faithfully follow human instructions. These models demonstrate remarkable capabilities in interpreting and visualizing complex scenarios presented in textual form.

Recent innovations have introduced zero-shot, training-free methodologies to this discipline. Text2Video-Zero [13] exemplifies this approach by adapting pre-existing text-to-image models to generate video sequences without extensive retraining, significantly reducing the computational footprint. This model enhances latent codes with motion dynamics and integrates cross-frame attention, ensuring temporal consistency across generated video frames from textual prompts. However, Text2Video-Zero lacks the capability for style image guidance, thus necessitating the training of a DreamBooth [22] model for each unique style to facilitate video generation with the desired stylistic characteristics. The StyleCrafter [14], on the other hand, introduced a style control adapter, trained with image datasets, to infuse style features from a reference image into video content. It incorporated a Scale-Adaptive Fusion Module, which bal-



anced text-based content features with image-based style features. Despite its advanced capabilities, StyleCrafter requires extensive training. It uses a style control adapter that must be trained with image datasets to effectively transfer style features from reference images into the generated video. The method also involves fine-tuning of the T2V models to improve their temporal dynamics and ensure consistent style throughout the video, enhancing the overall quality.

Unlike the aforementioned methods, StyleT2V0 is entirely training-free, eliminating the necessity for significant computing resources or multiple GPUs, all while allowing for style image guidance during the generation process.

### 3. Preliminaries

#### 3.1. Stable Diffusion (SD)

SD is a diffusion model operating in the latent space of an autoencoder  $\mathcal{D}(\mathcal{E}(\cdot))$ , namely VQ-GAN [4] or VQ-VAE [26], where  $\mathcal{E}$  and  $\mathcal{D}$  are the corresponding encoder and decoder, respectively. More precisely if  $x_0 \in \mathbb{R}^{h \times w \times c}$  is the latent tensor of an input image  $Im$  to the autoencoder, i.e.  $x_0 = \mathcal{E}(Im)$ , diffusion forward process iteratively adds Gaussian noise to the signal  $x_0$ :

$$q(x_t|x_{t-1}) = \mathcal{N}(x_t; \sqrt{1 - \beta_t}x_{t-1}, \beta_t I), \quad t = 1, \dots, T \quad (1)$$

where  $q(x_t|x_{t-1})$  is the conditional density of  $x_t$  given  $x_{t-1}$ , and  $\{\beta_t\}_{t=1}^T$  are hyperparameters.  $T$  is chosen to be as large that the forward process completely destroys the initial signal  $x_0$  resulting in  $x_T \sim \mathcal{N}(0, I)$ . The goal of SD is then to learn a backward process

$$p_\theta(x_{t-1}|x_t) = \mathcal{N}(x_{t-1}; \mu_\theta(x_t, t), \Sigma_\theta(x_t, t)) \quad (2)$$

for  $t = T, \dots, 1$ , which allows to generate a valid signal  $x_0$  from the standard Gaussian noise  $x_T$ . To get the final image generated from  $x_T$  it remains to pass  $x_0$  to the decoder of the initially chosen autoencoder:  $Im = \mathcal{D}(x_0)$ .

After learning the above mentioned backward diffusion process (see DDPM [9]) one can apply a deterministic sampling process, called DDIM [25]:

$$x_{t-1} = \sqrt{\alpha_{t-1}} \left( \frac{x_t - \sqrt{1 - \alpha_t} \epsilon_\theta^t(x_t, \tau)}{\sqrt{\alpha_t}} \right) + \sqrt{1 - \alpha_{t-1}} \epsilon_\theta^t(x_t, \tau), \quad t = T, \dots, 1, \quad (3)$$

where  $\alpha_t = \prod_{i=1}^t (1 - \beta_i)$  and

$$\epsilon_\theta^t(x_t) = \frac{\sqrt{1 - \alpha_t}}{\beta_t} x_t + \frac{(1 - \beta_t)(1 - \alpha_t)}{\beta_t} \mu_\theta(x_t, t). \quad (4)$$

To get a text-to-image synthesis framework, SD guides the diffusion processes with a textual prompt  $\tau$ . Particularly

for DDIM sampling, we get:

$$x_{t-1} = \sqrt{\alpha_{t-1}} \left( \frac{x_t - \sqrt{1 - \alpha_t} \epsilon_\theta^t(x_t, \tau)}{\sqrt{\alpha_t}} \right) + \sqrt{1 - \alpha_{t-1}} \epsilon_\theta^t(x_t, \tau), \quad t = T, \dots, 1. \quad (5)$$

It is worth noting that in SD, the function  $\epsilon_\theta^t(x_t, \tau)$  is modeled as a neural network with a UNet-like [21] architecture composed of convolutional and (self- and cross-) attentional blocks.  $x_T$  is called the latent code of the signal  $x_0$  and there is a method [2] to apply a deterministic forward process to reconstruct the latent code  $x_T$  given a signal  $x_0$ . This method is known as DDIM inversion. Sometimes for simplicity, we will call  $x_t, t = 1, \dots, T$  also the *latent codes* of the initial signal  $x_0$ .

#### 3.2. Text2Video-Zero

Our baseline model, which is termed *Text2Video-Zero* [13], is a state-of-the-art framework designed to generate videos directly from textual descriptions without requiring any fine-tuning or pre-training on video datasets. Text2Video-Zero harnesses the synergy of motion dynamics and cross-frame attention mechanisms to forge a path toward seamless text-to-video generation, ensuring both temporal consistency and fidelity to the input text description.

The process begins with a randomly sampled latent code  $x_1^T$  for the initial frame, which is then refined using  $\Delta t$  DDIM [25] backward steps to derive  $x_1^{T'}$ , utilizing a pre-trained Stable Diffusion model (SD) [19]. To incorporate motion dynamics, a specified motion field is applied across the video sequence, employing a warping function  $W_k$  for each frame  $k$ , which transforms  $x_1^{T'}$  into  $x_k^{T'}$ . This step is pivotal as it incorporates the latent codes with the requisite motion dynamics, thereby dictating the global scene and camera motion. This ensures temporal consistency across the background and the overall scene, enhancing the clarity of the generated video.

Subsequently, the denoised latent codes  $x_k^T$  for each frame  $k$  are obtained through a forward application of the DDPM [9] process. This probabilistic method allows for a greater degree of freedom in managing the motion of objects within the video, further contributing to a dynamic and realistic representation of the scene.

The core of Text2Video-Zero's capability to maintain visual continuity and object identity lies in its innovative use of cross-frame attention. By employing keys and values derived from the first frame's latent code, the model generates subsequent frames ( $k = 1, \dots, m$ ) that not only share stylistic and narrative elements but also preserve the appearance and identity of foreground objects across the sequence. This cross-frame attention is crucial for ensuring that each frame contributes to a cohesive video narrative, reflecting the text prompt accurately.

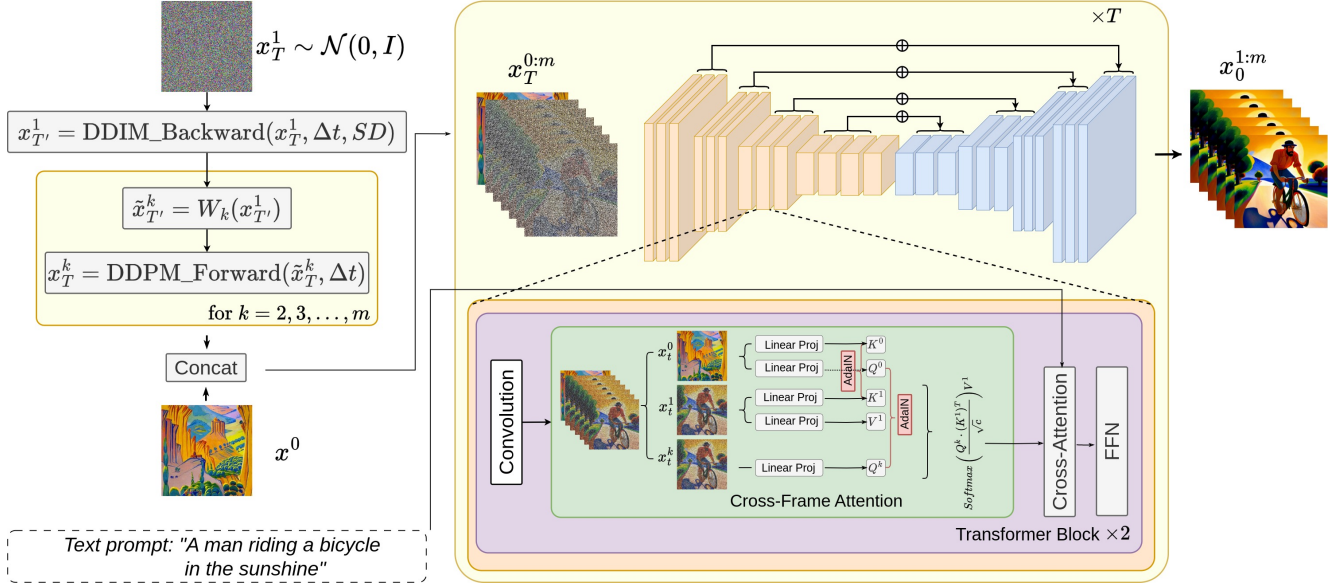


Figure 2. Architectural Overview of the StyleT2V0 Method: This diagram illustrates the integration of the shared attention mechanism within the Text2Video-Zero framework to enhance style consistency across video sequences. The method employs Adaptive Instance Normalization (AdaIN) to adjust the query ( $\hat{Q}_i$ ) for each frame  $F_i$  by blending style features from a designated reference image ( $Q_r$ ), ensuring stylistic uniformity and alignment with textual prompts. Keys ( $K$ ) and values ( $V$ ) are also adapted using AdaIN, aligning them with the style of the reference image while maintaining content continuity from the first frame ( $K_1$  and  $V_1$ ). The resultant attention mechanism applies these style-consistent queries, keys, and values to produce a video sequence where each frame is coherent in both style and content.

In addition to these mechanisms, Text2Video-Zero optionally applies background smoothing to enhance the visual quality of the video. Through salient object detection, a mask  $M_k$  is generated for each frame  $k$  to identify foreground pixels distinctly. A convex combination of the latent code  $x_1^t$  from the first frame warped to match the current frame  $k$  and the current frame’s latent code  $x_k^t$  is used to refine the background.

$$x_k^t = M_k \odot x_k^t + (1 - M_k) \odot (\alpha \hat{x}_k^t + (1 - \alpha) x_k^t), \quad (6)$$

This method ensures that the background remains consistent and unobtrusive, allowing foreground elements to stand out and maintain narrative focus.

#### 4. Method

In StyleT2V0, we extend the Text2Video-Zero framework by embedding a novel shared attention mechanism tailored to enhance style consistency across video sequences. This approach integrates style features from a designated reference image, ensuring uniform style across all frames while aligning with textual prompts. This is achieved through a specialized shared attention mechanism that utilizes Adaptive Instance Normalization (AdaIN) to blend style and content features across the sequence.

The style image is consistently referenced to dictate the

overarching style, while the content from the first frame serves as a baseline for content continuity.

For each frame  $F_i$  (where  $i = 1, \dots, n$ ), the query  $Q_i$  is adjusted using AdaIN to incorporate style features from the style image:

$$\hat{Q}_i = \text{AdaIN}(Q_i, Q_r), \quad (7)$$

where  $Q_r$  is derived from the style image, effectively embedding its stylistic attributes into the current frame’s query.

The keys  $K$  and values  $V$  are derived using AdaIN to align with the style properties of the style image, while maintaining content integrity from the first frame:

$$K = \text{AdaIN}(K_1, K_r), \quad V = V_1, \quad (8)$$

where  $K_1$  and  $V_1$  originate from the first frame, ensuring that the content features are consistently propagated through the video sequence.  $K_r$  is derived from the style image, aligning the keys with the desired style.

The attention for each frame  $F_i$  is computed to ensure that both style and content are coherently blended:

$$\text{Attention}(\hat{Q}_i, K, V) = \text{softmax} \left( \frac{\hat{Q}_i K^T}{\sqrt{d_k}} \right) V \quad (9)$$

where  $d_k$  is the dimensionality of the key vectors.

In addition to the shared attention mechanism, StyleT2V0 incorporates ControlNet [31] for detailed management of pose, edge, and depth attributes, enhancing the video’s motion dynamics and visual clarity. Pose control ensures accurate character movements through joint information, while edge control sharpens scene definitions, and depth control introduces a layered spatial perception for immersive depth effects.

## 5. Experiments

### 5.1. Implementation Details

Our implementation is based on the publicly available Text2Video-Zero codebase<sup>1</sup>, which we have adapted to integrate our methodological enhancements. These modifications include adjustments to the model’s handling of latent space dynamics and attention mechanisms. Importantly, our approach extends the input modality to not only include text prompts but also reference images, allowing for a richer context in video generation.

### 5.2. Qualitative Results

The StyleT2V0 framework showcases a significant advancement in generating stylistically consistent short videos directly from textual prompts. The framework’s flexibility was tested across a range of styles and contexts, showing its capability to adapt to various artistic and realistic styles. Whether translating a simple day-to-day activity or a complex, abstract concept into video format, StyleT2V0 maintained high fidelity to the original style and context of the reference image. Add references to Figs. In summary, StyleT2V0 not only advances the technical capabilities of text-to-video synthesis but also enhances the creative possibilities, allowing users to produce rich, stylistically consistent videos directly from textual descriptions. This marks a notable progress in the field of generative media, paving the way for more dynamic and aesthetically coherent video content generation.

### 5.3. Quantitative Results

**Testing Dataset.** To assess the effectiveness and scalability of our method, we assembled a test dataset comprising content prompts and style references. The content prompts consist of recognizable textual descriptions generated by GPT-4 [17], categorized into four meta-groups: human, animal, object, and landscape. For style references, we curated a collection of 10 diverse single-reference stylized images, encompassing both renowned artworks and generic style images. This structured dataset allows for a comprehensive evaluation of our video generation model across a variety of content and stylistic contexts.

<sup>1</sup><https://github.com/Picsart-AI-Research/Text2Video-Zero>

## 5.4. Ablation Study

## 6. Conclusion

In this study, we introduced StyleT2V0, a zero-shot framework for generating stylistically consistent videos from textual prompts, aligned with specified reference images without the need for retraining. Our approach leverages shared attention mechanisms, ensuring that each video frame not only adheres to the contextual details of the text but also harmonizes with the aesthetic qualities of the reference image. Furthermore, the integration of pose, edge, and depth controls allows for precise manipulation of dynamic movement, visual sharpness, and spatial depth, enhancing the expressiveness and realism of the generated videos.

The experimental results validate StyleT2V0’s ability to produce videos that are temporally consistent and stylistically faithful to both the textual prompts and the visual style of the reference images. Additionally, the adaptability to various control settings demonstrates the framework’s versatility and its potential application across various domains such as digital marketing and personalized media production, significantly reducing the need for computational resources typically associated with training advanced video generation models.

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